
Dennis S. Bernstein, Principal Investigator


$$
0=A_{s} Q+Q A_{s}^{T}+A Q_{A}^{T}+V_{1}+\left(A \cdot B R_{2 s}^{-1} P_{s}\right) \hat{Q}\left(A-B R_{2 s}^{-1} P_{s}\right)^{T}-Q_{s} \mathbf{v}_{2 s}^{-1} Q_{s}^{T}+r_{1} Q_{s} V_{2 s}^{-1} Q_{s}^{F_{1}^{T}} T_{1}^{T}
$$

$$
0=A_{s}^{\top} P+P A_{s}+A^{T} P_{A}+R_{1}+\left(A-Q_{s} V_{2 s}^{-1}\right)^{T} P^{\prime}\left(A-Q_{s} V_{2 s}^{-1}{ }^{C}\right)-P_{s}^{T_{2}} R_{2 s}^{-1} P_{s}+\tau_{1}^{\top} P_{s}^{T_{s}} R_{2 s}^{-1} P_{s} T_{1}
$$

$$
0=\left(A_{s}-B_{s} R_{2 s}^{-1} p_{s}\right) \hat{Q}+\hat{Q}\left(A_{s}-B_{s} R_{2 s}^{-1} p_{s}\right)^{T}+Q_{s} v_{2 s}^{-1} Q_{s}^{T} \tau_{1} Q_{s} v_{2 s}^{-1} Q_{s}^{T} T
$$

$$
0=\left(A_{s}-Q_{s} V_{2 s}^{-1} C_{s}\right)^{\top} \hat{P}+\hat{P}\left(A_{s}-Q_{s} V_{2 s}^{-1} C_{s}\right)+p_{s}^{\top} R_{2 s}^{-1} P_{s} \tau_{1}^{\tau_{P}} T_{s}^{T} R_{2 s}^{-1} P_{s}^{r}
$$

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| 2a. SECURITY CLASSIFICATION AUTHOAITY |  |  |  | Approved for public release; distribution unlimited. |  |  |  |
| 20. Declassification/oowngaliding Schedule |  |  |  |  |  |  |  |
| 4. PERFOAMING ORGANIZATION AEPORT Numberisi |  |  |  | 5. MONITOAING ORGANIZATION REPORT NUMBER(S) AFOSR.TR. \& 9-0066 |  |  |  |
| Ga nAme of performing organization <br> Harris Corporation |  |  | 3b. OFFICE SYMBOL (If applicable) | 7a. name of monitoring organization <br> AFOSR/NM |  |  |  |
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| 11. TITLE IInciude Security Clamificationt OPUS: Optimal Projection for Uncertain Systems. |  |  |  |  | 2304 | A1 |  |
| 12. PERSONAL AUTHOR(S) <br> Dennis S. Bernstein |  |  |  |  |  |  |  |
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This is the final report for the research project entitled OPUS: Optimal Projection for Uncertain Systems. OPUS is a unified approach to controlsystem design and analysis for high performance, multivariable applications such as large flexible space structures. In particular, OPUS yields loworder, robust controllers which meet both time- and frequency-domain
objectives. The present report is divided into three main research areas:

1) Fixed-Structure Design
2) Robust Analysis and Design
3) Further Extensions

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1) A unified approach to reduced-order, robust modeling, estimation, and control including singular problems and decentralized architectures
2) A computationally tractable approach to designing low-order, finitedimensional controllers for distributed parameter systems
3) A thorough development of quadratic Lyapunov bounds for robust stability and performance analysis
4) Complete unification of $L_{2}$ (time-domain) and $H_{m}$ (frequency-domain) design criteria for full-2 and reduced-order modeling, estimation, and control

The report includes reproductions of 41 research papers.

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# OPUS: Optimal Projection for Uncertain Systems 

For
Air Force Office of Scientific Research (AFSC)
Building 410
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Washington, DC 20332

Attention:
Major James Crowley

October 1988

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## Abstract

This is the final report for the research project entitled OPUS: Optimal Projection for Uncertain Systems OPUS is a unified approach to control-system design and analysis for high-performance, multivariable applications such as large flexible space structures. In particular, OPUS yields low-order, robust controllers which meet both time- and frequency-domain objectives. The present report is divided into three main research areas:

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4) Complete unification of $\mathrm{I}_{2}^{\mu}$ (time-domain) and $\mathrm{H}_{60}$ (frequency-domain) design criteria for full- and reduced-order modeling, estimation, and control. (id) =

The report includes reproductions of 41 research papers.


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SECTION 1.0
INTRODUCTION

### 1.0 INIRODUCTION

### 1.1 Overview

Over the past 10-15 years controls researchers have come to the realization that classical controls analysis and design techniques are inadequate in the face of modern large scale, high-performance applications. In particular, the principal motivation for OPUS is the problem of vibration suppression in large lightweight flexible space structures characterized by high-dimensional, highly uncertain models. In addition, stringent performance specifications in the face of high disturbance levels place severe demands on existing control-design techniques. Specifically, performance tradeoffs involving sensors, processors, actuators, and identification accuracy must be cut as tightly as possible to minimize hardware and testing costs. For feasibility and cost effectiveness, system design must also be performed efficiently with respect to human and computer resources.

The goal of this project has been to develop a mathematically rigorous control-design methodology which directly addresses these technology issues. In particular, optimal projection theory addresses the need for low-order, high-performance controllers which can be implemented on-board for real-time operation. Low-order controllers are necessitated by cost, weight, and reliability constraints associated with space-qualified processors. Furthermore, opus incorporates a fundamental theory of robust controller synthesis to account for unavoidable modeling uncertainties arising for reasons such as material and manufacturing variations, thermal and aging effects, as well as limits to identification accuracy. The principal contribution of OPUS is thus a unified theory which simultaneously accounts for both real-time processor constraints and modeling uncertainty. A high level overview of OPUS is given in [88] (Appendix A).

During the course of this project OPUS has, in addition, been extended to a large class of problems in systems and control theory. The current scope of the theory includes (see Figure 1-1):


Figure 1-1. Scope of OPUS

1. A unified treatment of reduced-order modeling, estimation, and control (Appendix B) ;
2. Robust estimation and control via quadratic Lyapunov functions including robust performance (Appendices G,H,I);
3. A unified approach to $\mathrm{I}_{2}$ and $H_{\infty}$ control including parametric robustness (Appendix J);
4. Decentralized, nonstrictly proper, and sampled-data control (Appendices D, E,I).

Of particular interest is the recent extension to $H_{\infty}$ control. As shown in [117] (Appendix J), we have developed a method for directly imbedding $H_{\infty}$ design constraints within OPUS theory and thus, in particular, within LQG. These results are given by a system of modified Riccati equations which directly generalize LQG theory and which have the potential for significant computational savings compared to existing $\mathrm{H}_{\infty}$ synthesis methods.

The underlying philosophy of OPUS is to capture as many design constraints as possible within a single system of design equations. This is demonstrated in [117] by the unification of time- and frequency-domain criteria addressed by the $L_{2} / H_{60}$ design equations. An additional example is provided by the results obtained in $[119,94]$ (Appendices $H$ and I) for robust stability and performance via fixed-order compensation in the presence of realvalued structured parameter uncertainty. In these algebraic design equations the projection matrix automatically enforces a constraint on controller order, while additional terms guarantee both robust stability and performance. Note that for full-order controllers in the absence of uncertainty, these four coupled equations reduce to the standard pair of separated Riccati equations of LQG theory. Versions of these equations have been developed for each of the problems shown in Figure 1-1. These results are discussed in more technical detail in the following sections.

The justification for this line of research is based upon several considerations. First, and most obvious, is the fact that our results show that numerous design constraints can be captured simultaneously within a constructive theory which directly generalizes IQG theory. Such an approach provides the capability for simultaneously performing multiple design tradeoffs
for multivariable systems with respect to competing constraints such as sensor noise, control authority, controller order, robustness, disturbance attenuation, mean-square error, sample rate, degree of decentralization, etc. Next we stress that rathe than being ad hoc constructions, these design equations follow directly from the optimality of well-defined performance objectives. Thus, these results are useful in assessing the suboptimality of alternative methods. For example, as shown in [20] several suboptimal approaches to reduced-order control design can be viewed as approximations to the optimal projection equations.

### 1.2 Status of Computational Results

Overall, OPUS can be viewed as a theory for characterizing solutions to constrained control-design jroblems. Transforming OPUS into a practical design methodology requires the development of effective computational algorithms. Such development has been carried out in related work by S. Richter at Harris Corporation. Using homotopic continuation methods, Richter has developed efficient algorithms which fully account for the structure of these modified Riccati equations and their coupling terms. Homotopy algorithms, as reviewed in
S. Richter and R. DeCarlo, "Continuation Methods: Theory and Applications," IEEE Trans. Autom. Contr., Vol. 28, pp. 660-665, 1983.
offer several advantages over both gradient-based and Newton-type methods. For example, homotopy methods have a strong theoretical foundation based upon differential topology, in particular, topological degree theory, while in practice these methods effectively address the key issues of startup, convergence, and global optimality. Homotopy algorithms have also reached a high degree of maturity and availability with the advent of HOMPACK described in

L. T. Watson, S. C. Billups, and A. P. Morgan, "A Suite of Codes for Globally Convergent Homotopy Algorithms," ACM Trans. Math. Software, Vol. 13, pp. 281-310, 1987.

The continuation algorithm developed for the optimal projection equations essentially follows a smooth path connecting an easily solvable version of the equations with the final, desired form. The algorithm utilizes the tensor derivatives of the terms in the optimal projection equations to integrate along the solution paths. To demonstrate the algorithm, an 8 th-order, nonminimum phase example originally due to
R. H. Cannon, Jr., and D. E. Rosenthal, "Experiments in Control of Flexible Structures with Noncolocated Sensors and Actuators," AIAA J. Guid. Contr. Dyn., Vol. 7, pp. 546-553, 1984.
was considered. This problem was used in
Y. Liu and B. D. O. Anderson, "Controller Reduction Via Stable
Factorization and Balancing," Int. J. Contr., Vol. 44, pp. 507 -
531 , 1986 .
to compare several reduced-order control-design methods. The comparisons performed by Liu and Anderson highlight the suboptimal nature of these methods. Specifically, several methods failed to yield stabilizing controllers for $10 \%$ of the cases while others failed for as many as $60 \%$. In contrast, as reported in [68,102], the optimal projection approach yielded stabilizing controllers for all cases considered. While the methods compared by Liu and Anderson were most prone to failure at high authority levels, the optimal projection results were within $20 \%$ of the LQG performance at $10^{2}-10^{3}$ higher authority levels. In addition, using topological degree theory, an upper bound has been obtained on the number of solutions of the design equations. Letting $n=$ plant dimension, $n_{u}=$ dimension of the unstable plant subspace, $n_{C}=$ compensator order, $\ell=$ number of measurements, and $m=$ number of controls, the number of solutions for the case $n_{C} \geq n_{u}$ is not greater than

$$
\begin{gathered}
\binom{\min (n, m, \ell)-n_{u}}{n_{C}-n_{u}}, n_{c} \leq \min (n, m, \ell) \\
1,
\end{gathered}
$$

Hence, for the case in which the controller order is greater than the number of inputs or outputs (so that the controller is not ill-conditioned), the equations possess at most one solution corresponding to the global minimum. Furthermore, since in many practical cases of interest this bound is small, it suffices to compute each such solution to determine the global optimum. These results along with suitable extensions to related problems have been used widely throughout this project. For example, recent results on fixed-order control of distributed parameter systems described in Section 2.3 were obtained using the homotopy algorithm.

### 1.3 Long-Range Goals of the Project

The long-range (5-10 year) goal of this project is the development of a truly effective computer-aided design methodology for multivariable control design. Numerical solution of the design equations would form the basis for such a design tool. This methodology would be appropriate in an engineering envirorment since the user need not be familiar with the mathematics of the design equations being solved. We envision a methodology similar to finite element modeling used routinely by structural analysts. An OPUS design package would go far beyond currently available packages whose multivariable design capabilities are based largely upon LQG theory.

### 1.4 Plan of This Report

Since this is the final report for this project our goal is to accomplish the following objectives:

1) Review the evolution and maturation of the research plan throughout the project;
2) Highlight the principal research accomplishments; and
3) Summarize open problems and point out future research directions.

Detailed technical discussion of results obtained will not appear in the main body of the report. Rather, the appendices contain a fairly complete (and lengthy) collection of the principal research results. We note that the ordering of the appendices is not chronological but instead reflects the most logical order according to subject matter.

SECTION 2.0
FIXED-STRUCTURE DESIGN

### 2.0 FIXED-STRUCIURE DESIGN

### 2.1 Motivation

While achieving the system design specifications (stability, performance, etc.) the control-design process must not lose sight of restrictions which arise in controller implementation. Indeed, control-design methods which focus primarily on performance specifications often pay a serious price by producing controllers which are difficult, if not impossible, to implement in practice. Hence our approach rests upon the notion of fixed-structure design. That is, we seek to meet design specifications within a framework which constrains the class of implementable designs. In this way the burden of hardware implementation (sensors, processors, and actuators) can be minimized to the greatest possible extent.

### 2.2 The Three Basic Problems

The most fundamental restriction arising in fixed-structure design is that of the order, or dimension, of the controller. In addressing this problem we have developed a unified treatment of three basic problems in reduced-order design, namely, modeling, estimation, and control. These three problems form a fundamental hierarchy of design problems in system theory, namely, to determine a system of fixed degree which, for a given system, approximates, estimates, or controls selected plant states. The solutions to these problems, given in [32,29,24] (Appendix B), reveal a surprising degree of common structure. Specifically, the solutions involve systems of 2,3 , and 4 modified algebraic Riccati and Lyapunov equations coupled by a projection matrix (the "optimal projection"). In addition, the estimation and control results provide transparent generalizations of steady-state Kalman filter and LQG theory.

Although the structure of these equations is aesthetically appealing by itself, the principal benefit for practical purposes is computational. That is, by exploiting the structure of these equations it is possible to significantly reduce the computational burden inherent in commonly used
gradient search techniques. This point has been amply demonstra+ $+\cdots$ in $[63,68$ ] as discussed in Section 1.

### 2.3 Finite-Dimensional Control of Distributed Parameter Systems

The problem of controller order becomes exacerbated when the plant is infinite dimensional since infinite-dimensional controllers cannot be implemented precisely, while finite-dimensional plant approximations may be of arbitrarily high order. To address this problem the fixed-structure controldesign results of [24] were generalized in [37] (Appendix C) to the case in which the plant is infinite dimensional. The resulting design equations now comprise a system of four operator equations coupled by a finite-rank nonselfadjoint projection operator. In spite of the infinite dimensionality of the plant, the design equations directly characterize fixed-order, finitedimensional dynamic compensator gains (Figure 2-1). Corresponding results for fixed-order finite-dimensional modeling and fixed-order finite-dimensional state estimation can also be obtained in an analogous manner.

Application of the operator-theoretic results of [37], however, requires finite-dimensional approximation of the design equations. In practice one could solve the design equations for a sequence of plant approximations of increasingly high order while keeping the controller order fixed. The limiting controller would then serve as a nearly optimal fixed-order finite-dimensional controller for the original distributed parameter system (Figure 2-2). This was investigated numerically in [122] in a collaborative project with Professor I. G. Rosen. In [122] two alternative approaches were considered for obtaining finite-dimensional controllers for infinite-dimensional systems. The first approach, which has been widely studied, involves computing a sequence of fullorder LQG controllers for a sequence of high-order plant approximations, while the second approach assumes a fixed order for the dynamic controller. To demonstrate these methods, two examples were considered, namely, a onedimensional parabolic (heat/diffusion) system and a hereditary (delay) system. For each example a sequence of spline-based, Ritz-Galerkin finite element approximations was derived for use in the control-design procedure. LQG theory and the optimal projection approach were then used to obtain full- and first-


Figure 2-1. The Optimal Projection Equations For Finite-Dimensional Fixed-Order Dynamic Compensation of Infinite-Dimensional Systems Provide a Direct Path to Optimal Physically Realizable Controllers for Distributed Parameter Systems

## CONVERGENCE OF SUBOPTIMAL REDUCED-ORDER COMPENSATORS

IDEA: DESIGN A SEQUENCE OF REDUCED-ORDER COMPENSATORS WHILE INCREASING THE ORDER OF THE APPROXIMATE MODEL AND KEEPING THE ORDER OF THE COMPENSATOR FIXED


Figure 2-2. Numerical Solution of the Optimal Projection Equations for Fixed-Order Dynamic Compensation Provides a Path to the optimal Fixed-Order Controller for an Infinite-Dimensional System
order controllers for each example with plant approximations up to 32nd order. For the parabolic system the perfomance degradation of the first-order controllers was only $2 \%$ compared to the full-order controllers (Figure 2-3), while for the hereditary system the degradation was less than $10 \%$. The difference in implementation requirements for a first-order versus a 32 nd-order controller is, of course, considerable.

### 2.4 Decentralized Control

In addition to incorporating constraints on the order of the feedback compensator, the fixed-structure approach allows additional constraints on the complexity of the feedback law. In particular, the results of [24] assumed a centralized structure for the dynamic compensator. In many applications, however, a decentralized controller architecture permits a simplified feedback communication structure and allows increased parallelism in the control law execution.

The fixed-structure approach is ideally suited to the decentralized design problem. For each fixed decentralized architecture, the design procedure can be performed to assess the ability to meet specifications for the given configuration. If specifications cannot be met, then the feedback architecture can be modified to improve performance, robustness, etc.

For the case of dynamic compensation, it was shown in [76] that the optimal projection technique provides a direct means for characterizing decentralized controllers. The key step is the realization that each subcontroller in the decentralized configuration must be an optimal centralized controller when viewed as a controller for the plant and remaining subcontrollers. This observation immediately suggests a sequential design algorithm in which individual subcontrollers are alternately refined until convergence is achieved. Because the method is based upon optimization principles, performance improvement is guaranteed at each step. This technique was demonstrated numerically in [76] (Appendix D) where a two-channel decentralized controller, fourth-order in each channel, was designed for a pair of interconnected simply supported beams. The algorithm demonstrated


Figure 2-3. The LQG Closed-Loop cost via Full-Order Controllers is Compared to First-Order Optimal Projection Designs for a Parabolic System
convergence to a decentralized controller whose performance was within $10 \%$ of the fully centralized controller.

For the case in which each subcontroller is a static (proportional) feedback law, it is possible to simultaneously characterize the optimal gains in each control channel without requiring a sequential approach. A thorough treatment of this case, including robust stability and performance, is given in [121] (Appendix D).

### 2.5 Singular Control

An important generalization of the results of [24] involves the case in which the controller includes a static feedthrough component. One technical issue which arises in the problem formulation is that the $I_{2}$ norm of a control signal corrupted by white noise (as a result of measurement feedthrough) is infinite. Hence the measurement feedthrough problem is only well-posed when either the measurement noise intensity or the control weighting matrix is singular. As is well known from the singular control literature, however, singular problem data often lead to complex behavior including impulsive controls and singular arcs. The imposition of a smooth controller structure via the fixed-structure approach thus precludes such complex behavior.

The fixed-order state estimation and dynamic compensation results of [29,24] were partially extended to the singular case in $[78,79]$. Even in the full-order case the singular control results are novel since standard LQG theory yields only strictly proper controllers. The results of $[78,79$ ] were incomplete, however, since the gains associated with certain estimation and feedback paths were not given explicitly. For the singular estimation problem this defect was remedied in
Y. Halevi, "The Optimal Reduced-Order Estimator for Systems with Singular Measurement Noise," IEEE Trans. Autom. Contr., Vol. 34, 1989.
where all feedback gains were explicitly characterized. In addition, this
solution was shown to agree completely with results obtained using standard limiting methods. For the corresponding dynamic-compensation problem the complete singular solution has been derived in joint research with Professor Y. Halevi and will be reported in $[130,138,139]$ (Appendix E).

SECTION 3.0

### 3.0 ROBUST ANALYSIS AND DESIGN

### 3.1 Motivation

The purpose of feedback control is to achieve desirable performance in the face of uncertainty. Although identification can reduce uncertainty to some extent, it is often impractical and residual modeling discrepancies always remain. For example, modeling uncertainty in flexible structures may arise in the mass, damping, and stiffness operators. Controllers must therefore be robust to achieve desired disturbance rejection in the presence of such modeling uncertainty.

### 3.2 Stochastic Modeling

Our approach to robust control was originally inspired by stochastic parameter modeling within a linear-quadratic optimization framework. In a series of early papers [1-16], D. C. Hyland explored the ramifications of a multiplicative white noise model as a consequence of the minimum information modeling technique based upon the Maximum Entropy Principle of Jaynes. The intent was not to view the white noise process as a literal model of parameter uncertainty, however, but rather to construct a tractable design model which captures the effects of parameter uncertainty upon system performance.

An interesting feature of the Maximm Entropy modeling approach was that the multiplicative white noise model was not to be rigorously interpreted as an Ito differential model, but rather in terms of the Stratonovich formulation. Recasting the Stratonovich model in terms of Ito differentials then led to additional "correction" terms. It is precisely these terms which were shown to play a crucial role in capturing the effects of parameter uncertainty. Such effects include decorrelation, i.e., the decrease in cross-correlation of system states due to parameter uncertainty, as well as equilibration, i.e., the tendency of state variances to equalize in the presence of high levels of uncertainty thus rendering different states indistinguishable. These effects of parameter uncertainty are fundamental features of high-order, lightly damped modal systems. An interesting treatment of these ideas for structural and
acoustic analysis can be found in

R. H. Lyon, Statistical Energy Analysis of Dynamical Systems: Theory and Application, MIT Press, Cambridge, MA, 1975.

For feedback design within fixed-structure design theory, the Stratonovich model produces controllers possessing intuitively appealing features. Specifically, such control laws exhibit high-authority control in the low-frequency, well-modeled portion of the structure along with lowauthority, rate dissipative action in the high-frequency region [35] (Appendix F). The ability to merge and unify these control regimes is a unique and significant contribution of the Maximum Entropy approach.

As a control-design methodology, however, it remained to validate the approach as a rigorous robust design technique. Optimal controllers designed in the presence of white noise disturbances, it was reasoned, are automatically desensitized to actual deteministic plant parameter variations. This idea was confirmed empirically by numerical studies in $[36,39]$ which showed an efficient design tradeoff between performance and robustness in the presence of structured real-valued parameter variations. Further robustness studies confirming these results were carried out in
A. Gruzen, "Robust Reduced Order Control of Flexible Structures," C. S. Draper Laboratory Report \#CSDI-T900, April 1986.
A. Gruzen and W. E. Vander Velde, "Robust Reduced-Onder Control of Flexible Structures Using the Optimal Projection/Maximum Entropy Design Methodology," AIAA Guid. Nav. Contr. Conf., Williamsburg, VA, August 1986.
M. Cheung and S. Yurkovich, "On the Robustness of MEOP Design versus Asymptotic LQG Synthesis," IEEE Trans. Autom. Contr., Vol. 33, 1988.

In spite of these results, it was clear that issues concerning stochastic modeling, such as stochastic stability and the physical interpretation of the model, tended to obscure the effectiveness of the technique for robust control. Thus a crucial step in the evolution of our approach was the ability
to show in [77] (Appendix F) that such controllers are guaranteed to be robust for all cases in which the design equations are solvable. In particular, it was shown that a second-moment stochastic stability condition in the presence of a time-exponential cost weighting induces a Lyapunov function which guarantees deterministic robust stability over a prescribed range of parameter variations. This result thus provided the bridge to cross over from the world of stochastic modeling (a statistical theory) to deterministic robustness theory (a theory of worst-case bounds).

### 3.3 Robust Analysis

For a given controller, it is often necessary to assess the stability and worst-case performance of the closed-loop system as parameters vary within a specified range of uncertainty. This is a problem of robust analysis, whose consideration precedes the more complex problem of robust controller synthesis.

Our principal matrematical technique in robustness analysis is Lyapunov stability theory. Here the idea is to determine a Lyapunov function which guarantees robust stability over a range of uncertain parameters. For linear systems we employ the quadratic Lyapunov function

$$
\begin{equation*}
V(x)=x^{T} P x \tag{1}
\end{equation*}
$$

or, equivalently, the Lyapunov equation

$$
\begin{equation*}
0=A^{T} P+P A+R \tag{2}
\end{equation*}
$$

for the linear system

$$
\begin{equation*}
\dot{x}=A x+W \tag{3}
\end{equation*}
$$

The dual equation

$$
\begin{equation*}
0=A Q+Q A^{T}+V \tag{4}
\end{equation*}
$$

is also useful for robust performance analysis since $V$ can be interpreted as the intensity of the additive white noise signal w . In robust analysis one typically replaces (4) by

$$
\begin{equation*}
0=A Q+Q A^{T}+\Omega+V \tag{5}
\end{equation*}
$$

where $\Omega$ is an additional nonnegative-definite matrix. Now robust stability of the perturbed system

$$
\begin{equation*}
\dot{x}=(A+\Delta A) x+w \tag{6}
\end{equation*}
$$

is assured so long as

$$
\begin{equation*}
\Delta A Q+g \Delta A^{T}<\Omega . \tag{7}
\end{equation*}
$$

This can be seen by rewriting (5) as

$$
\begin{equation*}
0=(A+\Delta A) Q+Q(A+\triangle A)^{T}+\left[\Omega-\left(\Delta A Q+Q \Delta A^{T}\right)\right]+V \tag{8}
\end{equation*}
$$

Furthermore, it is also possible to guarantee robust performance since the solution $Q_{\triangle A}$ of

$$
\begin{equation*}
0=(A+\Delta A) Q_{\triangle A}+Q_{\Delta A}(A+\Delta A)^{T}+V \tag{9}
\end{equation*}
$$

satisfies

$$
\begin{equation*}
Q_{\triangle A} \leq \Omega \tag{10}
\end{equation*}
$$

The above technique, developed in [115] (Appendix G), provides a simple
3-4
approach to robust stability and performance.

To develop a more sophisticated approach one can replace (5) by

$$
\begin{equation*}
0=A Q+Q A^{T}+\Omega(Q)+V \tag{11}
\end{equation*}
$$

where $\Omega(\cdot)$ is now a bounding operator which satisfies

$$
\begin{equation*}
\Delta A Q+\varrho \Delta A^{T}<\Omega(Q) \tag{12}
\end{equation*}
$$

for all variations $\Delta A$ in a specified uncertainty set and for all nonnegativedefinite matrices Q . This approach now guarantees the bounding a priori via (12) and the problem is to determine whether or not there exists a solution to (11).

The a priori bounding technique shown in (11), (12) has been given a fairly complete treatment in [123] (Appendix G). The goal in [123] was to systematically investigate candidate choices for the function $\Omega(\cdot)$. This investigation also provides a unified setting for particular bounds which have been used in various control-design contexts. For example, for $A=\sigma_{1} A_{1},\left|\sigma_{1}\right| \leq 1$, the absolute value bound

$$
\begin{equation*}
\Omega(Q)=\left|A_{1} Q+Q A_{1} T\right| \tag{13}
\end{equation*}
$$

where $|\cdot|$ replaces each eigenvalue by its absolute value, was proposed in
S. S. L. Chang and T. K. C. Peng, "Adaptive Guaranteed Cost Control of Systems with Uncertain Parameters", IEEE Trans. Autom. Contr., Vol. AC-17, pp. 474-483, 1972.

On the other hand, writing $A_{1}=D_{1} E_{1}$, the bound

$$
\begin{equation*}
\Omega(Q)=D+Q E Q, \tag{14}
\end{equation*}
$$

where $D=D_{1} D_{1} T$ and $E=E_{1} T_{E_{1}}$, was studied in
I. R. Petersen and C. V. Hollot, "A Riccati Equation Approach to the Stabilization of Uncertain Systems", Automatica, Vol. 22, pp. 433-448, 1986.
D. Hinrichsen and A. J. Pritchard, "Stability Radius for Structured Ferturbations and the Algebraic Riccati Equation", Sys. Contr. Lett., Vol. 8, pp. 105-113, 1986.

Finally, the choice

$$
\begin{equation*}
\Omega(Q)=\alpha Q+\alpha^{-1} A_{A_{1} Q A_{1}} T \tag{15}
\end{equation*}
$$

corresponds to the bound arising from a multiplicative white noise model as discussed in [77] and Section 3.2. We call (14) the quadratic bound (since it is quadratic in Q ) and (15) the linear bound (since it is linear in Q ).

### 3.4 Robust Synthesis

The principal payoff of our robust stability and performance technique is the ability to incorporate these bounds directly within the fixed structure design methodology. This can be done by bounding the cost over the class of parameter uncertainties prior to determining the feedback gains. The resulting bound is then treated as an auxiliary cost which can then be minimized by suitable feedback gains. The solution to this optimization problem is thus guaranteed to yield robust stability and performance.

To carry out this procedure it is essential that the bound $\Omega(\cdot)$ be differentiable with respect to $Q$. Furthenmore, $\Omega(\cdot)$ will be differentiable with respect to the feedback gains if it is differentiable with respect to $A_{1}$ (which involves gains in the control-design setting). These requirements thus suggest the linear bound (15) ard the quadratic bound (14) as the prime candidates for robust synthesis.

As discussed previously, the linear bound (15) was originally suggested by a multiplicative white noise model. By incorporating this bound within the design procedure, sufficient conditions for robust estimation and robust control were developed in [95,119] (Appendix H). In addition, a unified treatment of robust, reduced-order modeling, estimation, and control was given in [89] (Appendix H).

The quadratic bound (14) has also been developed extensively within a design context. In $[101,83,94]$ (Appendix I) the problems of reduced-order modeling, estimation, and control were considered via this bound. Finally, both the linear and quadratic bounds were considered simultaneously in [113] (Appendix I).

## $3.5 \mathrm{H}_{6}$ Theory

The robustness theory discussed in the previous subsections addresses the problem of real-valued structured parameter uncertainty. In many applications, however, uncertainty is present in the form of unstructured perturbations to the plant transfer function. A typical case is the presence of high-frequency, unmodeled dynamics.

A mathematically rigorous approach to this problem involves defining a suitable norm on the space of plant transfer functions to characterize uncertainty in terms of neighborhoods of the nominal plant. The resulting $H_{\infty}$ theory was pioneered by Zames in
G. Zames, "Feedback and Optimal Sensitivity: Model Reference Transformations, Multiplicative Seminorms, and Approximate Inverses," IEEE Trans. Autom. Contr., Vol. AC-26, pp. 301-320, 1981.
while recent overviews were given in
B. A. Francis and J. C. Doyle, "Linear Control Theory with an $H_{\infty}$ Optimality Criterion," SIAM J. Contr. Optim., Vol. 25, pp. 815844, 1987.
B. A. Francis, A Course in $\mathrm{H}_{6}$ Control Theory, Springer-Verlag, New York, 1987.

The most fundamental problem of $\mathrm{H}_{\infty}$ control design is the so-called Standard Problem considered by Francis: determine a feedback compensator which minimizes the peak (worst-case) disturbance attenuation of the closedloop system. By introducing suitable weighting matrices and problem transformations, solutions to the Standard Problem can be used to yield robust controllers for unstructured plant uncertainty.

Current $H_{\infty}$ synthesis methods, however, possess two principal drawbacks: they are computationally intensive and they often yield excessively high-order controllers. These difficulties have been removed with the advent of new state space solutions to the Standard Problem given in [117] (Appendix J) and
I. R. Petersen, "Disturbance Attenuation and $H^{\infty}$ Optimization: A Design Method Based on the Algebraic Riccati equation," IEEE Trans. Autom. Contr., Vol. AC-32, pp. 427-429, 1987.
P. P. Khargonekar, I. R. Petersen, and M. A. Rotea, " $\mathrm{H}^{\infty}$ optimal Control with State Feedback," IEEE Trans. Autom. Contr., Vol. 33, pp. 786-788, 1988.
J. C. Doyle, K. Glover, P. P. Khargonekar, and B. A. Francis, "State-Space Solutions to Standard $\mathrm{H}_{2}$ and $\mathrm{H}_{\infty}$ Control Problems," Proc. Amer. Contr. Conf., pp. 1691-1696, Atlanta, GA, June 1988.

These papers characterize solutions to the Standard Problem in terms of modified Riccati equations. The computational savings of this approach over earlier methods is considerable, possibly two orders of magnitude. In addition, the dynamic compensators obtained from these Riccati equatiows are of the same order as the plant model. This approach thus removes the principal drawbacks of earlier $\mathrm{H}_{\infty}$ synthesis methods.

By incorporating the fixed-structure approach we have, in addition, obtained the most general solution thus far available for the Standard Problem. Specifically, in [117] (Appendix J) we consicir the minimization of an $L_{2}$ performance criterion subject to a constraint on the $H_{\infty}$ closed-loop performance. This multi-norm problem formulation thus allows the designer to
perform tradeoffs between these competing performance measures. In addition we impose a constraint on the order of the dynamic compensator to obtain optimal low-order feedback controllers which satisfy the $H_{\infty}$ performance constraint. Utilizing an eighth-order nonminimm phase example given in
R. H. Cannon, Jr., and D. E. Rosenthal, "Experiments in Control of Flexible Structures with Noncolocated Sensors and Actuators," AIAA J. Guid. Contr. Dyn., Vol. 7, pp. 546-553, 1984.
we used these results to obtain 9 dB improvement over the corresponding $L_{Q} G$ design (Figure 3-1).

Inmediate spinoffs of these results include the problems of model reduction and state estimation. The $H_{\infty}$ model reduction problem [114] (Appendix $J$ ) addresses one of the most fundamental problems of linear system theory, namely, given a linear time-invariant system of degree $n$, find a linear time-invariant transfer function of degree $n_{m a}<n$ which minimizes the $H_{\infty}$ distance between the full- and reduced-order systems. Although the Hankel norm modelreduction problem has been widely studied as in

> K. Glover, "All Optimal Hankel-Norm Approximations of Linear Multivariable Systems and Their $L^{\infty}$-Error Bounds," Int. J. Contr., Vol. 39 , pp. $1115-1193,1984$.
the solution to the $H_{\infty}$ problem had not been given previousl. .

For state estimation the Kalman filter provides the least squares ( $L_{2}$ ) optimal solution. In certain applications, however, it may be desirable to minimize the worst-case frequency content of the error signal. This problem is addressed in [116] (Appendix J) where the standard steady-state Kalman filter is generalized to include a bound on the $H_{\infty}$ norm of the error signal.

Finally, it is reasonable to expect that in practice both structured and unstructured plant uncertainty will be present. This leads to consideration of the Standard Problem in the presence of parametric uncertainty. Thus it is of interest to design feedback controllers which are guaranteed to satisfy a

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Figure 3-1. The $L Q G / H_{\infty}$ Design Equations Yield 9 dB Improvement Over The Corresponding IQG Design for an 8 th-Order Nomminimum Phase Plant
specified $H_{\infty}$ disturbance attenuation constraint over a range of parametric uncertainty. This problem has been addressed in [105] (Appendix J) where the results of [117] on $H_{\infty}$ design have been merged with these of [94,119] on parametrically robust design. Again the development has been carried out in the context of fixed-order dynamic compensation for maximal design flexibility.

SECTION 4.0
FURIHER EXIENSIONS

### 4.0 FURIHER EXIENSIONS

### 4.1 Motivation

The previous sections have addressed two principal problems in control design, namely, fixed-structure design and robustness. Both of these problems concern fundamental issues in the practical implementation of feedback controllers. In this section we extend these results in two directions in order to address larger classes of design problems.

### 4.2 Tracking

All of the feedback control theory discussed in Sections 2 and 3 addresses the problem of feedback control for regulation in the presence of external disturbances. Many control problems, however, are of a tracking or servomechanism nature. While a limited class of such problems can be recast without loss of generality as regulation problems, many important ones cannot. For example, the standard transformations given in
H. Kwakernaak and R. Sivan, Linear Optimal Control Systems, Wiley, New York, 1972.
B. A. Francis, A Course in $\mathrm{H}_{6}$ Control Theory, Springer-Verlag, New York, 1987.
assume that the command signals can be represented as an augmentation of the plant dynamics. There are many important cases, such as the tracking of steps and ramps, which must be represented by uncontrollable, unstable dynamics, where this transformation cannot be applied. Furthermore, such transformations often ignore controller effort. To fill this gap we have undertaken a systematic program for developing a tracking control theory consistent with earlier developments. As a first step we have considered the problem of regulation about a prescribed nonzero set point, which corresponds to the step command tracking problem. Our work in this area was originally motivated by results obtained in
Z. Artstein and A. Leizarowitz, "Tracking Periodic Signals with the Overtaking Criterion," IEEE Trans. Autom. Contr., Vol. AC-30, pp. 1123-1126, 1985.
A. Leizarowitz, "Tracking Nomperiodic Trajectories with the Overtaking Criterion," Appl. Math. Optim., Vol. 14, pp. 155171, 1986.
A. Leizarowitz, "Infinite Horizon Stochastic Regulation and Tracking with the Overtaking Criterion," Stochastics, 1987.

References [67,103] (Appendix K) present general solutions to the nonzero set point problem for both static and dynamic controllers. The overall controller configurations for these problems are shown in Figures 4-1 and 4-2. Note that these controllers involve two components, namely, a closed-loop feedback component similar to a regulator and an open-loop feedforward component which has no counterpart in the standard theory and which cannot be obtained from standard transformations.

Recent activities have focused on extending the nonzero set point results to broader classes of command and disturbance signals. It turns out that the challenging case (as with steps and ramps) involves signals generated by unstable command or reference dynamics. As a critical first step in addressing this problem we have considered the problem of reduced-order steady-state observer design for unstable plants. These results appear in [125] (Appendix K). This optimal subspace observer problem gives rise to yet another projection which we denote by $\mu$. The most general estimation problem involving all three projections $\tau, \nu$, and $\mu$ has also been solved and will be reported in [134,139].

### 4.3 Sampled-Data Control

The discussion in the previous sections has focused on continuous-time systems subject to continuous-time (analog) controllers. In practice, however, controller implementation will almost invariably utilize digital controllers within the context of sampled-data control systems. Rigorous consideration of such systems is critical, particularly for distributed parameter systems which possess modal frequencies beyond the Nyquist rate of any digital controller.


Figure 4-1. The Static Controller for Nonzero Set Point Regulation Involves Both Feedforward and Feedback Gains


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Figure 4-2. The Dynamic Controller for Nonzero Set Point Regulation Involves Both an Open-Loop Gain and a Closed-Loop Dynamic Compensator

Hence, a rigorous theory of sampled-data control design must be developed which accounts precisely for all effects arising from analog-to-digital and digital-to-analog operations.

Optimal projection theory for discrete-time systems was developed in [41] and applied to sampled-data systems in [44] (Appendix L). As a next step it is desirable to obtain robust control results. To this end, the optimal projection equations for reduced-order discrete-time estimation and control in the presence of multiplicative white noise were obtained in [54,69] (Appendix L). After these results were obtained, it became clear that a true sampleddata robustness theory must account for the exponential matrix structure which arises from the sampling process. For example, if $A+\Delta A$ denotes the continuoustime dynamics matrix, where $A$ is the nominal matrix and $\Delta A$ denotes uncertainty, then the equivalent discrete-time dynamics matrix is given by $e^{(A+\Delta A) h}$, where $h$ is the sample interval. Because of the exponential function, however, this discrete-time dynamics matrix does not have the additive structure considered in the discrete-time theory in $[54,69]$. Moreover, a linear approximation for the exponential will not be valid in the presence of system time constants near or above the sample rate.

Although an attempt to bound this discrepancy resulted in new inequalities in [92] and questions of decomposition in [87] (Appendix L), this approach appears inadequate. The crucial clue to the most natural approach was ultimately found in
A. R. Tiedemann and W. L. DeKoning, "The Equivalent Discrete-Time optimal Control Problem for Continuous-Time Systems with Stochastic Parameters," Int. J. Cont., Vol. 40, pp. 449-466, 1984.
which studied the propagation of multiplicative white noise in the presence of $A / D$ and D/A interfaces. Motivated by these results, we have obtained results which extend the robust performance bounds obtained for continuous-time systems to the sampled-data problem. Specifically, by considering the evolution of the linear parameter uncertainty bound over the sample interval, a robust stability
condition was developed in [128] (Appendix L). This result is unique in that it accounts directly for the exponential structure of the parameter uncertainty.

SECTION 5.0
OPEN PROBLEMS AND FUIURE DIRECTIONS

### 5.0 OPEN PROBLEMS

### 5.1 Motivation

The value and importance of the results obtained under this project lie largely in the foundation they provide for future research. The purpose of this section is to collect together various questions and problems as a guide to future endeavors. The order of listing of these questions roughly parallels the order of the previous sections.

### 5.2 Fixed-Structure Design

Since the fixed-structure design approach involves a nonconvex optimization problem, there arise several questions concerning the structure of the space of solutions.

- Do there exist verifiable a priori conditions which guarantee stabilizability of a given linear time-invariant plant by fixed-order dynamic compensation? As in the fullorder case, one would expect such conditions to play a fundamental role in determining the existence of solutions to the design equations. Conversely, when the plant is known to be stabilizable by a controller of order $n_{C}$, does the underlying optimization problem always possess a solution? Will the design equations always yield at least one such stabilizing controller? How is the ability to find stabilizing controllers affected by the choice of weightings and noise intensities?
- Is it possible to design all subcontrollers of a decentralized dynamic compensator simultaneously without performing sequential iterations? If a sequential algorithm is used, then under what conditions is the algorithm guaranteed to find the global minimum?
- How can the fixed-structure approach be extended to address the simultaneous stabilization problem, i.e., the problem of finding a single controller which stabilizes several different plants simultaneously?
- The $L_{2}$ model reduction theory of [32] (Appendix B) can readily be extended to the problem of characterizing optimal finite-dimensional models for infinite-dimensional systems using the method of [37] (Appendix C). Can such finite-
dimensional models serve as useful lumped approximations to distributed parameter systems? Can the $\mathrm{L}_{2} / \mathrm{H}_{\infty}$ model reduction theory of [114] (Appendix J) be used similarly?
- How can the fixed-structure approach be used to design controllers with additional constraints on their internal structure, such as prespecified pole locations? This question is the basis for ongoing work in [131].


## 5. 2 Robust Analysis and Design

There exist a variety of open questions concerning the conservatism and effectiveness of the parametric robustness bounds and the $\mathrm{H}_{\infty}$ design equations.

- For which class of parameter uncertainty structures are the quadratic Lyapunov bounds nonconservative? How can the robustified design equations be used iteratively to reduce design conservatism?
- The multiplicative noise model was shown in [77] (Appendix F) to guarantee deterministic robustness. However, this result involved a uniform right shift rather than the variable left shift arising from the Stratonovich interpretation of the multiplicative noise. Can it be shown rigorously that the Stratonovich model yields robust controllers? Furthermore, can the relationship between Stratonovich design and positive real controllers for modal systens be made precise?
- The basis for the $\mathrm{H}_{\infty}$ design results obtained in [117] (Appendix $J$ ) is the quadratic bound developed for parametrically robust control in [94] (Appendix I). This raises the following question: Does there exist an alternative interpretation of the linear bound which can be used to guarantee disturbance attenuation for some specified class of disturbances?
- The $H_{\infty}$ control design results are virtually identical to the optimality conditions for the problem of minimizing an exponential-of-quadratic cost criterion as considered in
P. R. Kumar and J. H. van Schuppen, "On the optimal Control of Stochastic Systems with an Exponential-of-Integral Performance Index," J. Math. Anal. Appl., Vol. 80, pp. 312-332, 1981.
P. Whittle, "Risk-Sensitive Linear/Quadratic/Gaussian Control," Adv. Appl. Prob., Vol. 13, pp. 764-777, 1981.
A. Bensoussan and J. H. van Schuppen, "Optimal Control of Partially Observable Stochastic Systems With an Exponential-ofIntegral Performance Index," SIAM J. Contr. Optim., Vol. 23, pp. 599-613, 1985.
P. Whittle and J. Kuhn, "A Hamiltonian Formalation of RiskSensitive Linear/Quadratic/Gaussian Control," Int. J. Contr. Vol. 43, pp. 1-12, 1986.

Is it possible to directly extend these results using the fixed structure approach? Also, can the fixed-structure approach be used to extend the Maximum Entropy theory of
D. Mustafa and K. Glover, "Controllers Which Satisfy a ClosedLoop $H_{\infty}$ Norm Bound and Maximize an Entropy Integral," Proc. IEEE Conf. Dec. Contr., Austin, TX, December 1988.

- The $\mathrm{I}_{2} / \mathrm{H}_{\infty}$ model reduction theory given in [114] (Appendix J) minimizes an $I_{2}$ criterion subject to a constraint on the $H_{\infty}$ distance between the full- and reduced-order models. Can the $\mathrm{L}_{2}$ criterion be neutralized so as to obtain a "pure" $\mathrm{H}_{\infty}$ result as is done in [117] (Appendix J) for full-order control design? Can the resulting $\mathrm{H}_{\infty}$ solution be shown to actually characterize the $\mathrm{H}_{6}$ optimal reduced-order model by taking the $H_{\infty}$ constraint to be sufficiently small? Similar questions apply to fixed-order control design. For example, does there exist a "pure" $H_{b}$ reduced-order control design theory? Can these results be shown to be necessary as well as sufficient?
- What is the generalization of the $H_{\infty}$ control and estimation results to the singular problem? To the cross-weighting problem?
- Is it possible to extend the $\mathrm{L}_{2}$ and $\mathrm{L}_{2} / \mathrm{H}_{\infty}$ model reduction results to allow the reduced-order model to be nonstrictly proper?


### 5.4 Tracking and Sampled-Data Control

With regard to tracking and sampled-data theory a number of problems remain to be explored.

- Is it possible to develop a methodology for designing tracking controllers which applies to a broad range of signal models? For example, the command signal may be known exactly in advance (such as a specified square wave) while, at the other extreme, it may only be known to be an element of a large class of signals. For example, step commands are known to be steps but their exact level is not known until they actually occur during operation. Other command signals may only be known to be outputs of systems driven by random noise. A classification scheme based upon the degree and type of priori knowledge of the command signal should lead to a hierarchy of control designs ranging from poorly known to well-known command sigmals. In addition, it is important to distinguish between a priori command signal knowledge available during the design phase and command signal knowledge available during operation. The differences between these cases can be used to account for differing assumptions appearing in the literature. Relevant references include
B. D. O. Anderson and J. B. Moore, Linear Optimal Control, Prentice-Hall, Englewood Cliffs, NJ, 1970.
C. D. Johnson, "Accommodation of External Disturbances in Linear Regulator and Servomechanism Problems," IEEE Trans. Autom. Contr., Vol. AC-16, pp. 635-644, 1971.
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W. E. Schmitendorf and B. R. Barmish, "Robust Asymptotic Tracking for Linear Systems with Unknown Parameters," Automatica, Vol. 22, pp. 355-360, 1986.
- How can the new subspace projection $\mu$, which arises in the observer design problem in [125] (Appendix K), be used to design servocompensators? That is, can $\mu$ be used to design controllers which track the output of an unstable command model?
- Is it possible to develop a theory of robust sampled-data controller synthesis which accounts directly for the exponential structure of the equivalent discrete-time model? The results of [128] (Appendix L) provide a starting point in this regard.
- What is the form of the equations for the $H_{\infty}$-constrained discrete-time control-design problem?

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PROGRAM PERSONNEL

### 7.0 Program Personnel

### 7.1 Dr. Dennis S. Bernstein

The principal investigator for this project was Dr. D. S. Bernstein. Dr. Bernstein received the Sc.B. in applied mathematics from Brown University and both the M.S. and Ph.D. degrees from the University of Michigan in control engineering. From 1982 to 1984 he was a staff member at Lincoln Laboratory, MIT, before joining Harris in 1984. He is a member of IEEE and SIAM.

### 7.2 Professor Wassim M. Haddad

Throughout this project Dr. Bernstein was assisted by Professor W. M. Haddad. Professor Haddad received the Ph.D. degree in 1987 from Florida Institute of Technology in mechanical engineering with specialization in systems and control. His doctoral research on robust control theory contributed directly to many of the developments achieved under this project. He is currently a faculty member at FIT and is a member of IEEE.

### 7.3 Acknowledgements

Many persons contributed either directly or indirectly to the results obtained under this project. At Harris we specifically acknowledge numerous helpful conversations with Dr. D. Hyland, Dr. J. Shipley, Dr. E. Collins, Mr. S. Richter, Mr. S. Greeley, Mr. L. Davis, Mr. A. Daubendiek, Mr. D. Phillips, and Mr. A. Tellez. We are also grateful to Mssrs. Richter, Greeley, Daubendiek, and Tellez for performing the numerical calculations in support of the project.

Finally, we are particularly indebted to Ms. J. Straehla who, without assistance, mastered PC TEX and singlehandedly procuced every paper emanating from this project. Her contribution was exceptional and invaluable to the project. We also thank Ms. L. Ford for producing this final report.
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# Optimal Projection for Uncertain Systems (OPUS): A Unified Theory of Reduced-Order, Robust Control Design 

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This research was supported in part by the Air Force office ot Scientific Research under contracts E49620-86-C-0002 and F49620-86-C-0038.

## Abstract

OPUS (Optimal Projection for Uncertain Systems) provides new machinery for designing active controllers for suppressing vibration in flexible atructures. The purpose of this paper is to review this machinery and demonstrate its practical value in addressing the structural control problem.

## 1. Introduction

For many years it has been widely recognized that the desire to orbit large, lightveight space atructures possessing high-performance capabilities would require active feedback control techniques. More generally, the need for such techniques aay arise due to the combinations of either 1) moderate performance requirements for highly flexible structures with low-frequency modes or 2) stringent performance requirements for semi-rigid structures with relatively high-frequency modes (Figure 1). Applications include pointing. slewing, and aperture shape control for optical and RF syeteus.
"Small" structures

- Older generation of spacecratl
- Mosi civil engineering struclures.
(from stiength static loading point of view)

CONTROLLER AUTHOAITY AND OA RESPONSE AANGE OF INTEAEST
poini bi view

"Large" structures

- Highly fiexibie spacecratt. tall buildings. rapid transit structures. etc
And/or
- Stringent pointing accuracy and oplical quality requirements
- Noise abalemenl (acoustical) structural interaction)


Figure 1. The Need for Active Structural Control Arises From Stringent Performance Requirements or Low-Frequency Nodes


Figure 2. Vibration Control Systems Utilize Sensors. Processors and Actuators to Suppress Disturbances

The problew of active vibration suppression (Figure 2) entails the following considerations:

1. Multiple, highly coupled feedback loops. The potentially large number of sensors and actuators leads to a fully coupled multi-input, multi-output feedback contral system.
2. Limited actustor power. The control authority available from on-board actuators is limited by weight, size. cost and power considerations.
3. High-dimensional models. Large structures subjected to broadband disturbances are typically represented by high-order finite element arodels.
4. Limited processor capacity. Reliability and cost considerations limit the processor caparity available for on-board real-time implementation of the control system.
5. Highly uncertain models wirh atructured uncertainty. Finite element models often exhibit sienificant error particularly as modal frequency increases. Although modal testing and related identification methods may be used to improve modeling accuracy, residual uncertajnty always remains anc unpredictable on-orbit changes due to eging, thermal effects, etc., must be tolerated.
6. Stringent performance requirements. Since active space structure control is most relevant in precision applications, it can readily be expected that performance specifications will be particularly stringent.
7. Design efficiency. Because of implementation complexity due to the presence of multiple loops, high dimension, and high levels of uncerteinty, the control design approach should efficiently utilize both synthesis and analysis techniques (Figure 3).


Figure 3. Control-System Desien Must Efficiently Utilize Both Synthesis and Analysis Techniques

These considerations pose a congiderable challenge to the state-of-the-art in control-design methodologies. For example, the presence of multiple, coupled feedback paths essentially precludes the effectiveness of single-loop design techniques. The sheer number of loops, their interaction, and the need to address a host of other issues render such methods inefficient and unwieldy.

In addition to the presence of multiple loops, the high dimensionaliry of dynamic models places a severe burden on control-design methodologies. For example, although LQG (linear-quadratic-Gaussian) design is applicable to multi-loop problems. such controllers are of the same order as the structural model (Figures 4 and 5). Thus $L Q G$ and similar hjgh-order controllers can be expected to plece an unacceptable computational burden on the real-time pronessing capability. Hence it is not surprising that a variety of techniques have haer. $\overline{\mathrm{F}}$ : : used to reduce the order of LQG controllers, A comparison of several such met.lods is given in [1].

All of the above difficulties are severely exacerbated by the fact that the dynamic (i.e.. finite element) model upon which the control design is predicared may be highly inaccurate in spite of extensive modal identification. Kence, applicable control-design methodologies must account for modeling uncertainties by providing robust (i.e.. insensitive) controllers. Furthermore, because of stringent

HIGH-ORDER PLANT XGR ${ }^{n}$


Figure 4. LQG Theory Addresces the Problea of Designing a Quadratically Optimal, Eull-Order Dynamic Compencator

## FULL-ORDER CONTROLLER GAINS

$$
\begin{aligned}
& A_{c}=A-Q \bar{Y}-Y P \\
& B_{c}=O C^{T} V_{2}^{-1} \\
& C_{c}=-R_{2}^{-1} B^{T} P
\end{aligned}
$$

## SEPARATED RICCATI EQUATIONS

$$
\begin{aligned}
& O=A Q+O A^{\top}+V_{1}-Q E \quad \text { (Kalman Filler) } \\
& 0=A^{\top} P+P A+R_{1}-P Y P \quad \text { (Regulatior) } \\
& \sum \equiv B A_{2}^{-1} g^{\top} \quad \bar{\Xi} \equiv C^{\top} V_{2}^{-1} C
\end{aligned}
$$

Figure 5. The Optimal Full-Order (LQG) Controller Is Detemined by a Pait of Separsted Riccari Equations
performance requirements, robust control design must avoid conservatism with respect to modeling uncertainty which may unnecessarily degrade performance. A galient example of conservatism is illustrated in Figure 6. If uncertainty in the modal frequency is complexified in a transfer function setting, then the resulting pole location uncertainty has the form of a disk. This disk, however, intersects the right half plane in violation of energy dissipation. Hence one source of conservatism is the inability to differentiate between physically diatinct parameters auch as modal frequency and modal damping.


Figure 6. Complexification of Real Parameters May Lead ro Robustness Conservatiam

Although classical methods are inappropriate for vibration control, a wide variety of modern techniques are available. These include both multi-loop frequencydomain methods and fime-domain techniques. A comprehensive review of such methods will not be attempted here. Rather, we shall mereiy point out aspects of several methods which motivate the philosophy of OPUS development.

As is well known, dynamic models can be transformed (at least in theory) between the frequency and time domains. Significant differences arise, however, in artempting to represent modeling errors. Specifically, model-error characterization of a particular type. which is natural and tracrable in one domain, may become extremely cumbersome when transformed into the other domain. Fo: example, consider a stete space model with psrameter uncertainties arising in the syaten matrices $(A, B, C)$. Upon transforming ro a frequency domain model $G(s)=C(s I-A)^{-1} B$ the parametric uncertainties may perturb the transfer function coefficients in a
complicated manner. A more natural measure of uncertainty for transfer functions has been developed in [2] where system uncertainty in the frequency domain is modeled by meanc of normed neighborhoods in the tl-infinify ropology. There are limitationc with this approach, however. in designing controllers for vibration suppression. For example. as Ghown in Eigure 6. complexification of real-parameter uncertainties cuch as modal frequencies may yield unnecestary conservatism, while norm bounds often fail to preserve the physical structure of parameter variations. A case in point is the lightly damped oscillator. As shown in [A42], norm bounds predict stability over a frequency range on the order of the damping while in fact the oscillator is unconditionally stable. Furthemore, with regard to processor throughpur tradeoffs, modern frequency-domain methods typically yield high-order controllers.

Although LQG addregses performance/actuator and performance/sencor tradeoffs in a multi-loop setting, it fails to incorporate modeling uncertainty. Thus it is not surprising. as shown in [3]. that LQG designs fail to possess guaranteed gain margin. Since LQG designs lack such margins, attempts have been made to apply frequency-domein rechniques to improve their characreristics. One such method, known as LQG/LTR ( $[4,5]$ ) seeks to recover the gain margin of full-gtatefeedback controllers. Specifically, full-state-feedback $L Q R$ controllers are guarenteed to remain stable in the face of perturbations of the input marrix $B$ of the form $\alpha B$ where $\alpha \in[1 / 2, \infty)$. As shown in $[6,7]$, however, the full-state-feedbark gain margin fails to provide robustness with respect to perturbations which are not of this form. For instance, the example given in [6] with $B=[01]^{T}$ can be destabilized for auitable performance weightings with perturbation $B(E)=\left\{\varepsilon[1]^{T}\right.$ for arbitrarily small $E$ in spite of the 6 dB margin. Furthermore. since LQG/LTR loop shaping is based upon singular value norm bounds, trearment of physically meaningful real parameter variationg may lead to unnecessary conservatism. Several approaches have been proposed for circuaventing these difficulties (see.e.g., [8]).

The importance of addressing the problew of ctructured uncertainty in finite element modele cannor be overemphasized. Structural characteristics such as modal frequencies, damping ratios, and mode shapes appear explicitly in ( $A, B, C$ ) state-space modela as physically meaningful parameters. Uncertainty in mode shapes, for example, which appear as columns of the $B$ matrix, cannot in general be expected to be of multiplicative form in accordance with traditional gain-marein specifications. This is precisely the problem illustrated by the example of (6) discussed above. Furthermore, uncertainties in modal frequencies and damping rarios must be carefully differentiated since, roughly epeaking, modal frequency uncertainties affect only the imaginary part of the pole location while damping uncertainty affects the real part. Although these and related observations
concerning uncertainty in the dynamjc characteristics of lightly damped structures may be self evident, they have remained largely unexploited in standard controldesign methods.

## 2. OPUS: New Machinery for Control-System Design

In view of the ability of LQG theory to synthesize dynamic controllers for multi-input, multi-output controllers, it is nor surprising that LQG forms the basis for a variety of structural control methods. However. as discussed previously. LQG lacks the ability to address performance/procescor and performance/robustness tradeoffs. This situation has thus motivated the development of numerous variants of LQG which entail additional procedures which attempt to remedy these defects. OPUS. however. is distinctiy different. Rather than append additional procedures to LQG design. OPUS extends LQG theory itself by generalizing the basic underlying machinery.

As shown in Figure 5. the basic machinery of LQG consists of a pair of geparated Riccati equations whose solutions serve to directly and explicitly synthesize the gains of an optimal dynamic compensator. The contriburion of OPUS is to directly expand this machinery. The overall approach is illustrated in figure 7 which portrays two distinct seneralizations of the basjc LQG machinery. As Figure 7 illustrates, thece generalizations can be developed individually when ejther loworder or robust controllers are desired. The appealing aspect of opuS, however. is the ability to estend LQG to address both problems simultaneously in a unified manner.


Pigure 7. The Standard LQG Result Is Generalized by Both the Fixed-Order Coneraint and Modeling of Parameter Uncertainties

In the following sections the generalizations depicted in Figure 7 will be reviewed following the left branch. That is, the optianl projection approach to reduced-order controller design will first be discussed in Section 3 without introducing plant uncertainties. In Section 4 the reduced-order congtrajnt will be retajned while considering, in addition, uncertainties in the system model. In each case the discussion vill focus on the underlying ideas with a minimum of technical detail.

Clearly, in order for a novel design methodology ro be of pracrical value it mugt be computationally tractable. Hence Section 5 will present an overview of the current state of algorithm development for solving the OPUS design equations. Finally. Section 6 will briefly sumarize further OPUS generalizations of LQG theory which are relevant to atructural control.

## 3. Extensions of LQG to Reduced-Order Dynamic Compensation

The simplear, most direct way to obtain optimal reduced-order controllers is to redevelop the standard LQC result in the presence of a constraint on controller dimension (Eigure 8). The mathematical technique required to do this is remarkably strajghtforward. Specifically, the structure and order of the controller are fixed and the performance is optimized with respect to the controller gains. The recultinc necegsary conditions obrained using Lagrange multipliers thus characterize the opeimal gaing.


STEADY-STATE PERFORMANCE CRITERION

$$
J\left(A_{c}, B_{c}, C_{c}\right)=\lim _{1 \rightarrow \infty} E\left[x^{\top} R_{1} x+L^{\top} R_{2} u\right]
$$

Eigure 8. In Accordance With On-Board Processor Requirements, a Reduced-Order Constraint Is Imposed on the Diaension of the Dynamic Compensator


#### Abstract

This parameter optimization approach as such is not new and was invest igated extensively in the 1970's. Typically, however, the optimality conditions were found to be complex and unwieldy while offering little insight and requiting gradient search methods for numerical solution.


One curious aspect of the parameter optimization literature is that no attempt was made to actually use this direct method to rederive the LQG result itself. Such an exercise. it may be surmised, mjght revesi hidden structure within the optimality conditions which would shed light on the reduced-order case. Indeed. such an approach led to the realization that an oblique projection (idempotent matrix) js the key to unlocking the unvieldy oprimality conditions ([A7, Al7]). Although the result js mathemat ally straightforward, it is by no means obvious gince in the full-order (LQG) case the projection is the jdentity and hence not readily apparent.

By exploiting the presence of the projecrion, the necessary conditions can be transformed into coupled system of four algebrajc matrix equations consisting of a pair of modified Riccati equations and a pair of modified Lyapunov equations (Figure 9). The coupling is via the oblique projection $r$ which appears in all four equations and which is determined by the solutions $\dot{Q}$ and $\vec{P}$ of the modifjed Lyapunov equations. A satisfying feature of the optimality conditions is that in the fullorder case the projection becones the identity, the modified Lyapunov equations drop out, and. since $\boldsymbol{T}_{\perp}=0$, the modified Riccati equations specialize to the usual separated Riccati equations of LQG theory. Since, furthermore, $G=r=n \times n$ identity, the standard LQG gain expressions are recovered.

Alrhough the modified Riccati equationg specialize to the standard Riccari equations in the full-order case, the modified Lyapunov equations have no counterpart in the standard theory. The role of these equations can be understood by considering the problem of optimal model reduction alone. For this problem the optimal reducedorder model is characterized by a pair of coupled modified Lyapunov equarions (see [A22]). Thus the modified Lyapunov equarions arising in the reduced-order dynamic-compensation problem are direcely analogous to the modified Lyapunov equations arising in model reduction alone. The modified Lyapunov equations arising in the control problem, howeyer, are intimately coupled with rhe modified Ricati equations. Hence it cannot be expected that reduced-order control-design techniques besed upon LQG will generally yield optimal fixed-order controllere (Eigure 10). It is interesting to note that seversl such methods discussed in [1] are based upon balancing which was shown in [A22] to be suboptimal with respect to the quadratic (least Equares) oprimality criterion.

## REDUCED-ORDER CONTROLLER GAINS

$$
\begin{aligned}
& A_{c}=I\left(A-Q \Xi-\Sigma P / G^{\top}\right. \\
& B_{c}=V^{\top} O C^{\top} V_{2}^{-1} \\
& C_{c}=-A_{2}^{-1} B^{\top} \rho G^{\top}
\end{aligned}
$$

## COUPLED RICCATI/LYAPUNOV EQUATIONS

$$
\begin{aligned}
& 0=A Q+O A T+V_{1}-O \bar{X} Q+r_{1} Q \bar{E} Q r_{1} \\
& 0=A^{\top} P+P A+R_{1}-P \leq P+r_{1}^{T} P \leq P_{r_{1}}
\end{aligned}
$$

$$
\begin{aligned}
& 0=(A-Q E)^{r \hat{P}}+\hat{P}(A-Q \bar{S})+P \pm P-r_{t}^{\top} P \sum_{r_{1}} \\
& \operatorname{rank} \hat{O}=\operatorname{rank} \hat{P}=\operatorname{rank} \hat{\alpha} \hat{P}=n_{c} \\
& \hat{Q} \hat{P}=G^{\boldsymbol{T}} M r \quad \quad G^{\top}=I_{n_{c}} \\
& t=G^{\top} I^{\prime}=\hat{O} \hat{P}(\hat{O} \hat{P}) \| \quad T_{1}=I_{n}-T \\
& \Sigma=B R_{2}{ }^{-1} \mathrm{BT}^{\top} \quad \bar{I}=\mathrm{C}^{\top} \mathbf{V}_{2}{ }^{-1} \mathbf{C}
\end{aligned}
$$

Figure 9. The Oprimal Reduced-Order Compensator Is Determined by a Pair of Modified Riccati Equations and a Pair of Modified Lyapunov Equations Coupled by the oblique Projection $T$


Ejgure 10. The Optimal Projection Equations Provide a Direct Path to Optimal Reduced-Order Dynamic Compensators

In sumary, the optimal projection equations for reduced-order dynamic compensation comprise a direct extension of the bagic LQG machinery to the reducedorder control problem. The design equations, which reduce to the standard LQC result in the full-order case, provide direct synthesis of optimal reduced-order controllera in accordance with impleaentation constraints.

## 4. Exrensions of LQG to Uncerrain Podeling

Two fundanental sources of error in modeling flexible structures are truncated modes and parameter uncertainties. Since the optimal projection approach permita the utilization of the full dynamics model, modal truncation can be largely avoided. There remains, however, a tendency to truncate poorly known modes and thus jt is essential to incorporate a model of parameter uncertainty in both well-known and poorly known components of the system. Hence the problem formulation of $E$ gure 8 is now generalized in Eigure 11 to include uncertain parameters $\sigma_{i}$ appearing in the A. $B$ and $C$ matrices. The parameter $\sigma_{i}$ is assumed to lie within the interval $\left[-\delta_{i}, \delta_{i}\right]$ In accordance with identification accurscy. Clearly, when uncer:ainty is absent. i.e.. when $A_{i}, B_{i}, C_{i}=0_{\text {, }}$ the reduced-order design problem of $E$ igure 8 is recovered.

HIGH-ORDER, UNCERTAIN PLANT

- Stochastic disturbance model
- Deterministic parameter uncertainty model


Figure 11. Robust Optimal Projection Design Is Based Upon a Hybrid Uncertaincy Model Involving a Determiniatic Parameter Uncertainty Model and a Stochastic Disturbance Model

A salient feature of the design model is that uncertainty is modeled in two distinctly different ways. External uncertainty appearing as additive white noise is modeled stochastically. Such a model appears appropriate for dicturbances such as coolant flow for which only power spectral data are avajlable. On the other hand. internal uncertainty appearing as parameter variations is modeled deterministically. Such a model appears appropriate for uncertainty arising from directly measurable quantities such as mass and stiffness. Thus the overall uncertainty model is hybrid in the sense that it utilizes both deterministic and stochastic characterizations of uncertainty.

A natural performance meabure which accounts for borh types of uncertainty characterization involves the usual LQC quadratic criterion averaged over the disturbance statistics and then maximized over the uncertain parameters (Figure 12). Hence this performance measure incorporates on the average and worst case aspects in accordance with physjcal considerations.

PERFORMANCE CRITERION


ROBUST PERFORMANCE PROBLEM
Minimize $J\left(A_{c}, B_{c}, C_{c}\right)$ over the class of robustly stabilizing controllers ( $A_{c}, B_{c}, C_{c}$ )

Figure 12. Performance Is Defined To Be Worst Case Over the Uncertain Parameters and Average Over the Disturbance Statistics

The resuiting Robust Performance Problem thus involves determining the gaine ( $A_{c} \cdot B_{c} \cdot C_{c}$ ) to minimize the performance $J$. The static gain $D_{c}$ can also be included but will not be discussed here. Despite the apparent complexity of the problem, remarkably simple techniquec can be used, Specifically, first note that after taking the expected value the performance $J$ has the form

$$
\begin{equation*}
J\left(A_{c} \cdot{ }^{B} c_{c} \cdot C_{c}\right)=\sup _{\sigma_{i}} \lim _{t \rightarrow \infty} \sup \operatorname{tr} \dot{Q}(t) \tilde{R}, \tag{4.1}
\end{equation*}
$$

where "tr" denotes race of a matrix. $\bar{Q}(t)$ ia the covariance of the closed-loop ayster. and $\bar{R}$ is an augmented weighting matrix composed of $R_{1}, R_{12}$ and $R_{2}$. The covariance $\bar{Q}(t)$ satisfies the standard Lyapunov differential equation

$$
\begin{equation*}
\dot{\bar{Q}}=\left(\bar{A}+\sum \sigma_{i} \bar{A}_{i}\right) \bar{Q}+\bar{Q}\left(\tilde{A}+\sum \sigma_{i} \bar{A}_{i}\right)^{T}+\tilde{V}_{0} \tag{4.2}
\end{equation*}
$$

where $\bar{A}$ is the closed-loop dynamics, $\bar{A}_{i}$ is composed of $A_{i}, B_{i}$ and $C_{i}$, and $\vec{V}$ is the intensity of external disturbances for the closed-loop system including the plant and measurement noise.

Two distinct approaches to this problem will be considered. The first involves bounding the performance over the class of parameter uncertainties and then choosing the gaing to minimize the bound. Since bounding precedes control design this approach is known as robust design via a priori performance bounds. The second approach involves exploiting the nondestabilizing nature of structural systems via weak subsystem interaction.
4.1 Robust Design Via A Priori Performance Bounds

The key step in bounding the performance (4.1) is to replace (4.2) by a modified Lyapunov differential equation of the form

$$
\begin{equation*}
\dot{g}=\tilde{A} \tilde{g}+\tilde{g} A^{T}+\psi(\tilde{g})+\tilde{v} \tag{4.3}
\end{equation*}
$$

where the bound $\Psi$ satisfies the inequality

$$
\begin{equation*}
\sum \sigma_{i}\left(\bar{A}_{j} \tilde{\mathrm{O}}^{-}+\tilde{Q A}_{i}^{T}\right) \leq \Psi(\tilde{\mathrm{G}}) \tag{4.4}
\end{equation*}
$$

over the range of uncertain parameters $\sigma_{i}$ and for all candidate feedback gains. Note that the inequality ( 4.4 ) is defined in the sense of nonnegative-definite marrices. Now rewrite (4.3) by appropriate addition and subtraction as

$$
\begin{equation*}
\bar{Q}=\left(\bar{A}+\sum \sigma_{i} \bar{A}_{i}\right) \bar{Q}+\bar{Q}\left(\bar{A}+\sum \sigma_{i} \bar{A}_{i}\right)^{T}+\Psi(\bar{Q})-\sum \sigma_{i}\left(\bar{A}_{i} \bar{Q}_{\underline{Q}} \bar{Q}_{i}^{T}\right)+\overline{\bar{V}} \tag{4,5}
\end{equation*}
$$

Now subtract (4.2) from (4.5) to obtain

$$
\begin{equation*}
\overline{\bar{Q}}-\overline{\bar{Q}}=\left(\bar{A}+\sum \sigma_{i} \bar{A}_{i}\right)(\bar{Q}-\bar{Q})+(\bar{Q} \bar{Q})\left(\bar{A}+\sum \sigma_{i} \bar{A}_{i}\right)^{T}+\Psi(\bar{Q})-\sum \sigma_{i}\left(\bar{A} \bar{Q}_{i}+\underline{Q} \bar{A}_{i}^{T}\right) \tag{4.6}
\end{equation*}
$$

Since by (4.4) the term

$$
\begin{equation*}
\Psi(\bar{g})-\sum \sigma_{i}\left(\bar{A}_{i} \bar{S}+\overline{G A_{i}}\right) \tag{4.7}
\end{equation*}
$$

is nonnegative definite, it follows immediately that

$$
\begin{equation*}
\dot{Q} \leq \dot{0} \tag{4.8}
\end{equation*}
$$

over the class of uncertain parameters. Thus the performance (4.1) can be bounded by

$$
\begin{equation*}
J\left(A_{c}, B_{c}, C_{c}\right) \leq \underline{J}\left(A_{c}, B_{c} \cdot C_{c}\right) \triangleq \lim _{t \rightarrow \infty} \operatorname{tr} \ddot{Q R} \tag{4.9}
\end{equation*}
$$

The auxiliary cost $J$ is thus guaranteed to bound the actual cost J. This leads to the Auxiliary Minimization Problem: Minimize the auxiliary cost $J$ over the controller gains. The advantage of this approach is that necessary conditions for the Auxiliary Minimization Problem effectively serve as gufficient conditions for robust performance in the original problem. Since the bounding step precedes the optimization procedure, this approach is referred to as robust design via a priori performance bounds. This procedure is philosophically siailar ro guaranteed cost control ( 9,10$\}$ ). Nore rhar since bounding precedes optimizarion, the bound (4.4) must hold for all gains since the optimal gains are yer to be determined.

To obtain sufficient condirions for robust stability, the bounding function $\Psi$ must be specified. Since the ordering of nonnegative-definite marrices appearing in ( 4.4 ) is not a rotal ordering, a unique lowest bound should not be expected. Furthermore, each differentiable bound leads to a fundamentel extension of the optimal projection equarions and rhus of the basic LQG machinery. In work thus far. two bounds have been extensively investigated. Only one bound, the right shift/multiplicative white noise bound, will be discussed here. The strucrured stabiliry radius bound introduced in [11,12] is discussed in [A43].

The right shift/multiplicarive white nojse bound imestigated in [A29: A41] is given by

$$
\begin{equation*}
\psi(\vec{g})=\sum \delta_{j}\left(\alpha_{i} \tilde{Q}^{+} \alpha_{i}^{-1} \bar{A}_{i} \overline{Q A}_{i}^{T}\right) \tag{4,10}
\end{equation*}
$$

where $\alpha_{i}$, 0 are arbitrary scalars. Note that rhis bound consists of two distinct partis which must appear in an appropriate rario. The firsr rerm $\alpha_{i} \bar{Q}$ arises narurally when an exponential time weighting $e^{\alpha_{i}{ }^{t}}$ is included in the performance measure. As is well known ([13]) this leads to a prescribed uniform atability margin for the
closed-loop system (Figure 13). A uniform ctability margin. no marter how large, however, does not guarantee robustness with respect to arbitrary parameter variations. The complementary second term $\alpha_{i}^{-1} A_{j}^{-} \mathcal{Q A}_{i}^{T}$ is crucial in this regard.


Figure 23. Open-Loop Right-Shifted Dynamics Arising From Exponential Cost Weighting Lead to a Uniform Closed-Loop Stability Margin

Although terms of the form $\bar{A}_{i} \overline{Q A}_{i}^{-T}$ are unfamiliar in robust control design. they arise naturally in srochastic differential equarions with multiplicarive whire noise. That is, if the uncertain parameters $\sigma_{i}$ are repleced by white noise processes entering multiplicarively rather than additively, then che covariance equetion for $\bar{Q}$ automatically includes terms of the form $\tilde{u}_{i} \tilde{Q A}_{i}^{T}$. The literature on systems with multiplicative white noise is quite extensive; see [A38] for references. It should be stressed, however, rhar for our purposes the multiplicarive white noise model is not interpreted literally as having physical significance. Rather. aulriplicarive whire noise can be thought of as a useful degign model which correctiy captures the impact of uncertainty on the performance functionsl via the srate covariance. Furthermore, just as the right shifr term alone does not guarantee robustness. nejther does the multiplicative white noise term. Both terms must appear sfaltaneously. Roughly speaking, since mulriplicarive white noise disturbs the plant though uncertain parameters, the closed-loop system is automatically desensitized to actual parameter variations.

After incorporating the right chift/multiplicarive white noise bound (4.10) into (4.3) to obtain a bound $J$ for the performance, the optimel projection equations can be rederived following exactly the same parameter optimization procedure discuased in Section 3. Again. the mathematics required is but a straightforward applicarion of Lagrange multipliers. The addirional bounding rerms are carried through the derivation to yield a direct generalization of the optimal projection equations shown in Figure 14 with gains given in Eigure 15.

$$
\begin{aligned}
& 0=\left(A_{s}-B_{s} R_{2 s}^{-1} r_{s}\right) \hat{Q}+\hat{O}\left(A_{s}-B_{s} R_{2 s}^{-1} P_{s}\right)^{T}+\varrho_{s} V_{2 s}^{-1} V_{s}^{\top}-I_{1} N_{s} V_{2 s}^{-1} s_{s}^{\top} r_{1}^{\top}
\end{aligned}
$$

Figure 14. The Robustified Optimal Projection Design Equations Account for Both Reduced-Order Dynamic Compenaation and Parametric Uncertainty

## GAINS

$$
\begin{gathered}
A_{c}=I^{( }\left(A_{s}-B_{s} R_{2 s}^{-1} r_{s}-Q_{s} V_{2 s}^{-1} C_{s}\right) G^{\top} \\
B_{c}=V_{s} v_{2 s}^{-1} \\
C_{c}=-R_{2 s}^{-1}{ }_{s} G^{T} \\
\quad \text { NOTATION }
\end{gathered}
$$

$$
\begin{aligned}
& \hat{\mathbf{Q}} \hat{P}=G^{\top} M \Gamma, \quad \Gamma G^{\top}=I_{n_{c}}\left(\Rightarrow r=G^{\top} \Gamma=r^{2}\right) \\
& \lambda O A^{T}=\sum_{i=1}^{P} A_{i} O A_{i}, \quad \lambda Q A=\sum_{i=1}^{P} A_{i} Q B_{i} \text {, etc. } \\
& H_{2 s}=R_{2}+B^{\top}(\mathbf{P}+\hat{P}) B \quad V_{2 s}=V_{2}+c(\mathbb{Q}+\hat{Q}) c^{\top} \\
& Q_{s}=O C_{s}^{\top}+V_{12}+A(O+\hat{Q})^{\top} \quad r_{s}=B_{5}^{\top} P+R_{12}^{\top}+B^{\top}\left(P+\hat{P}_{)} A\right.
\end{aligned}
$$

Figure 15. The OPUS Controller Cains Are Explicitiy Characterized as a Direct Generalization of the Classical LQG Gains

The robustified optimal projection equarions comprise a system of four marrix equarions coupled by both the optimal projection and uncertainry rerms. When the uncertainty tems are absent, the optimal projection equations of figure 9 are immediately recovered. On the orher hand, if the order of the controller is set equal to the order of the plant, then all terac involving $T_{1}$ can be deleted. However, in this case the modified Lyapunov equations do nor drop out since $\vec{Q}$ and $\vec{P}$ still appeor in the modified Riccati equations. Nence rhe basic machinery of LQC is again exrended to include a pair of lyapunov equations coupled ro a generalizarion of the standard LQG equations. It is interesting to note thar a relared result in the context of multiplicative noise also appcared in the Sovier lirerature ([14]). It should also be pointed out that although the modified lyapunov equations ariaing in the reduced-order control-design problem have analogues in model reduction, the modifi :- Lyapunov equarions appearing in the full-order robuctified equarions represent new machinery not anticipated in robustness rheories. Hence using straightforward mathematical techniques, the basic LQC machinery has again been extended in novel directions.

Solving the design equarions shown in Figures 14 and 15 yields controllers with guaranteed levels of robustness. The actual robustress levels alay, however, be larger than specified by a friori bounds. Thus, to achipe desired robustification levels for the uncertainty structure specified by the a priori bounds, the design procedure may be utilized within an iterarive synthesis/analysis procedure (figure 16).


Figure 16. Optimal Projection/Guaranteed Cost Control Provides Direct Synthesis of Robust Dynamic Compencators

### 4.2 Robusr Design Via Weak Subeybrem Interaction

The mechanism by which LQC was robuctified in Section 4.1 avolved bounding the performance over the class of narameter uncertainties and rhen derermining optimal controller gains for the performance bound. As discussed in Section 2.
however, flexible structures possess special properties which may, in addition, be exploited ro achieve robustness. Specifically, aside from rigid-body modes, energy dissipation implies that mechanical structures are open-loop stable regardless of the level of uncertainty. That is, flexible structures possess only nondesrabilizing uncertainties. llence, in the closed loop, a given controller may or may not render a particular uncertainty destabilizing. A priori bounds on conrroller performance must, however, be valid for all gains since bounding precedes optimization. Hence. a priori bounding may in certain caces fail to exploit nondestabilizinc uncertainties.

A familiar example of a nondestabilizing uncertainty involves uncertain modal frequencies. Such an uncertainty will not, of course, destabilize an uncontrolled (open-1oop) structure. If particular modal frequencies are poorly known then it is clearly advisable to avoid applying high authority control. Hence, rather than the right-shift approach of Eigure 13, it appears advantageous (although, at first. counterintuitive) to utilize just the opposire. namely, a left shift (Figure 17). Furthermore, in view of the fact thar uncertainty usually increases with modal frequency (Eigure 18), a variable left shift appears to be more appropriate than a uniform left shift. By left-shifting high-frequency poorly known modes, the control-system design procedure applies correspondingly reduced authority ro modes "percejved" as highly damped. Hence the variable left shift can be roughly rhought of as a device for achieving suitable authority rolloff. As will be seen, however, the underlying robust ification mechanism. nemely, weak subsystem interaction. is far more subtle than the approach of classical rolloff rechniques. It is also interesting to note that the weak subsystem interaction approach ro robustness is entirely distinct from classical robustness approaches which urilize high loop gain to reduce sensitivity.

$$
\dot{x}=A x \Rightarrow \dot{x}=\left(A+\frac{1}{2} \sum_{i=1}^{p} A_{i}^{2}\right) x
$$



Figure 17. A Variable Left Shift Exploita Open-Loop Nondestebilizing Uncertaintas


Figure 18. Modal Uncertainty Generally Increases With Erequency

A variable left shift can readily be introduced into the robustified optimal projection design equarions by replacing $A$ by

$$
\begin{equation*}
A_{B}=A+\frac{l}{2} \sum A_{i}^{2} \tag{4.11}
\end{equation*}
$$

where $A_{j}$ denotes the structure of modal frequency uncertainty (Figure 19). Most interestingly. such a modification of the dynamics matrix arises narurally froma multiplicative white noise model defined not in the usual Ito sense bur rather in the sense of Stratonovich. Thus, as in the a priori bounding approach, a stochastic


Eigure 19. For Modal Syatems With Frequency Uncertainty the Stratonovich Correction Corresponds to a Variable Left Shift
model serves to suggest a mechanisn for robustification (Figure 20). Again it is important to stress that the multiplicative white noise model is not interpreted literally as having physical significance, but rather can be thought of as a useful design model which correctly captures the impact of uncertainty on the performance functional via the state covariance.

ROBUSTNESS BOUNDS


Figure 20. Stochastic Models and Robustness Bounds Are Fundamentally Related

In earlier work the stratonovich dynamics model was justified by means of the minimum information/maximum entropy approach ([A1-AlS]). A central result of the maximum entropy approach is that the high aurhority/low authority rransition of a vibration control system from well-known low-frequency modes to poorly known highfrequency modes (Figure 18) is directly reflected in the ctructure of the stare covariance matrix (Figure 21). A full-stare feedback design applied ro a simply


Figure 21. Frequency Uncertainties in the Stratonovich Model Lead to Suppressed Cross Correlation in the Steady-State Covariance
supported beas illuatrates this pojnt (Ficure 22). By assuming that uncertainty in modal frequencies increases linearly with frequency. the structure of the covariance marrix leads directiy ro the conrol goina illustrated in Figure 23. Note that in the high-frequency region rhe position gains are essentially zero and thus the control law approaches positivemeal energy dissipative rate feedback. This. of course, is precisely the type of atructural controller expected in the presence of poor modeling information. of course, any effective control-design theory for active vibration suppression in flexible structures should produce energy dissipative controllera when structural modeling information is highly uncertain.


- SImply. supporito blam with force actuator
- homdimensiomal eouations of monton $\left(\vec{u}_{n}=n^{2}\right)$
- "Emergr" state-wighnimg
- ymcentannties in open-loop fitoutncics
$T_{K}=-\bar{Z}^{\prime \prime}\left(\sigma_{K} \bar{\sigma}_{R}\right)^{-1}$
$\sigma_{K}=\begin{gathered}\text { stambaho devation } \\ \text { OF Kim mode preourmet }\end{gathered}$
- simple uncentaimty modet
- func-state fegodacx

Figure 22. The Effects of Erequency Uncertainties Can Be Illustrated for a One-Dimencional Beam With Idealized Full-State Feedback

To carry out robustified oprimal projection design in the presence of leftshifted open-1 oop dynamics, it is only necessary to utilize the left-shifred dynamics matrix (4.11) in place of the right-shifted matrix. All of the robustified optimal projection machinery, including gain expressions, can be utilized directly. It is also important to seress that the left shifr must be used in conjunction with terms of the form $\vec{A}_{i} \tilde{Q A}_{j}$.

One explanation for the mechanism by which robustificarion is achieved is illustrated in Eigure 24. By left shifting the open-loop dynamica within the degign process, the compensator poles are similarly left-shifted. Thus the compensaror polec are effectively moved further into the left half plane away from the actual plant poles. Since the interaction between compensator and plant poles is weakened, the closed-loop system is correspondingly robustified with respect to uncertajnties in the plant pole locarions. A sensitivity analysis of this mechanism urilizing a uniform left shift in the context of LQG design is given in [15].

a owestip onect mesuti or staatonovich monft
Figure 23. The Maxitum Entropy Controller Approaches Rate Feedback in the Limit of Poor Modeling Information (High Uncertainty)

$$
\underbrace{\left[\begin{array}{cc}
A & c B C_{c} \\
c B_{c} C & A_{c}
\end{array}\right]}_{c=1 \Rightarrow L O G}=\underbrace{\left[\begin{array}{cc}
A & 0 \\
0 & A_{c}
\end{array}\right]}_{\substack{\text { PlanUCompensator } \\
\text { Subsystems }}}+\underbrace{\left[\begin{array}{cc}
0 & c B C_{c} \\
t B_{c} C & 0
\end{array}\right]}_{\substack{\text { Subsystem interactions } \\
0 \leq c \leq 1}}
$$



> LOG puts compensator poles near plant poles $\Rightarrow$ O( $\kappa)$ shift (sirong inleracilon)

$\Rightarrow \underbrace{0\left(c^{2}\right) \text { thilt (weak Interection) }}_{\substack{\text { Wecreased sensilivity to } \\ \text { plant variations }}}$

Figure 24. The Stratonovich Variable Left-Shift Model Effectively Places the Compensator Poles Eurther Into the Left Half Plane Where Plant/Compensator Intersction Is Weakened

As discuesed above, the left-shift approach exploits open-loop
nondegtabilizing uncertainties and thus cannot operate through priori bounding. Thus the actual level of robustification achieved from the robustified optimal projection equations for a given level of uncertainty nodeling cannot be predicted a priori. i.e.. in advance of control design. Indeed, this situation is to be expected when nondestabilizing uncertainties are explojted in a nonconservative design theory. Thus a suitable robust analysis technique is required for nonconservatively determining the robustification of the closed-loop system with respect to open-loop nondestabilizing uncertainties.

A suitable robustness analysis technique, known as majorant Lyapunov analysis, has indeed been developed ([A42]). Essentially, this technique employs a new type of Lyapunov function for assessing robustness due to weak subsyatem interaction. The underlying machinery consists of the block-norm matrix which is a nonnegative matrix each of whose elements is the norm of a block of a suitably partitioned matrix (Figure 25). A matrix which bounds the block-norm matix in the sense of nonnegative matrices, i.e.. element by elsment, ia known as amajorant. Majoranta were introduced in [16] and were applied to atability analysie of integration algorithms for ODE's in [17].
(Ostrowski, 1961; Dahlquist, 1983)


NONNEGATIVE CONE ORDERING
$m \leq \leq \hat{m}$

[^0]To apply majorants to dynamical systems, the model is written in the form shown in Figure 26. The matrix A is block diagonal and consjsts of subsystem dynamics. The subsystem interactions represented by the partitioned matrix $G$ are assumed to be uncertain. By suitable manipulation. uncertainties in the diagonal blacks of A can also be captured by $G$. By assuming that the spectral norm (largest aingular value) of the blocks of $G$ satisfy given bounds. the covariance block-norm inequality is obtajned (Figure 27). This inequality is interpreted in the sense of nonnegative matrices, i.e., element-by-element, and $*$ denotes the Hadamard (element-by-element) product.

$$
\dot{i}=(A+G) x+w \quad \dot{Q}=(A+G) Q+a(A+G)^{\top}+v
$$




Nolse intensify


Figure 26. The Large-Scale System Model Involves Known Local Dynamics and Uncertain Inceractions

$$
\begin{array}{ll}
\dot{x}=(A+G) x+w & J=E\left(x^{\top} R x\right]=\operatorname{tr} O R \\
0=(A+G) O+O(A+G)^{\top}+v & J=\left\{\Omega\left(A_{i} \oplus A_{j}\right) \mid\right.
\end{array}
$$

$$
\nu=\left[\begin{array}{cc}
\left\|v_{1}\right\|_{F} & \left\|v_{12}\right\|_{F}-- \\
\left\|v_{21}\right\|_{F} & \left\|v_{2}\right\|_{F} \\
\vdots &
\end{array}\right] \quad Q=\left[\begin{array}{cc}
\left\|a_{1}\right\|_{F} & \left\|a_{12}\right\|_{F}-- \\
\left\|a_{21}\right\|_{F} & \left\|a_{2}\right\|_{F} \\
1 &
\end{array}\right]
$$

$$
\left[\begin{array}{ccc}
0 & \bar{\sigma}\left(G_{12}\right) & --- \\
\bar{\sigma}\left(G_{21}\right) & 0 & \\
\vdots & & \ddots \\
1 & &
\end{array}\right] \leq \leq!i
$$

$\Omega$

$$
A \cdot Q \leq \leq G Q+Q G T+D
$$

Figure 27. The Block-Norm Matrix of the State Covariance Satisfies a Lyapunov-Type Inequality Involving Nonnegative Matrices

To achieve robustness, the covariance block-norm inequality is replaced by the moran Lyapunov equation (Figure 28). The solution of the majorant Lyapunov equation provides bound (majorant) for the block norm of the covariance thereby guaranteeing both robust stability and performance.

## MAJORANT LYAPUNOV EQUATION

$$
\left.\begin{array}{c}
A \cdot \widehat{Q}=S \widehat{Q}+\widehat{Q} S_{T}+\mathcal{D} \\
\bar{\sigma}\left(G_{i j}\right) \leq \mathcal{S i j} \\
\Omega \\
Q \leq \widehat{Q} \\
\Omega
\end{array}\right]=\text { Robust Stability }
$$

Figure 28. The Corresponding Nonnegative Matrix Equation Yields a Majorant for the State Covariance and Hence Robust Stability and Performance

It ie interesting to note that numerical solution of the majorant Lyapunov equation requires no new techniques. Utilizing properties of matrices, the solution can be obtained monotonically by means of a straightforward iterative technique (Figure 29).

MLE has a unique solution ill $\left.\{\widehat{\{ })_{K}, K=0,1, \ldots, \infty\right\}$ where:

$$
\begin{aligned}
& \hat{Q}_{0}=0 \\
& \widehat{\hat{Q}}_{K+1}=1 \|^{\prime \prime} \cdot\left(i ;{\hat{V_{K}}}_{K}+\widehat{( }\right)_{K}(; \boldsymbol{T}+1) \\
& \left(11 I_{m n} \triangleq 11-1 m n\right)
\end{aligned}
$$

converges. If so, then:

$$
\widehat{\hat{V}}=\lim _{K \rightarrow \infty} \widehat{G}_{K}
$$

$$
\begin{aligned}
& J-J_{0} \leq 2 \sum_{K=1}^{r}\left(\operatorname{tr} \hat{P}_{K}\right)\left(\left(\hat{\beta} \hat{S}_{K K}\right.\right. \\
& \left(0=A_{K}^{\top} \hat{P}_{K}+\hat{P}_{K} A_{K}+R_{K}\right)
\end{aligned}
$$

Figure 29. By Exploiting the Properties of M-Matrices. the Majorant Lyapunov Equation Can Be Solved Monotonically by Means of a Simple Iterative Technique

An jllustrative application of the majorant Lyapunov equation involves lightly damped aubsyens (Eigure 30). As shown in [A42] (and expected intuitively). robustnesa with respect to uncertain subsystem interaction is proportional to the frequency separation between the subsysteme. The ability to capture thia robustifjcation mechaniam ia a unique feature of the majorant lyapunov function not available from quadratic (i.e., scalar) Lyapunov functions or vector Lyapunov functions ([18.19]).


Majorant Lyapunov Equation Bound $\sim \sqrt{(2 v)^{2}+\left(\omega_{1}-\omega_{2}\right)^{2}}$

Eigure 30. Robustneas Bounds for Uncertain Coupling in Modal Systema Are Proportional to the Frequency Separation Between Subsystems

The next atep in the majorant development involves a hierarchy of finer and finer robugnese bounds which account for higher order aubsyctem interactions, e.g., the interaction between the ith and jth subsystems via the kth aubsystem. The second member of the hierarchy (Figure 31) provides robustness guarantees with respect to frequency uncertainties. The interesting aspect of this robustness test is the fact that the performance bound is characterized precisely by a Stratonovich model. Hence the Stratunovich model can be viewed as an approximation to a robustness bound, while exploiting the Stratonovich/majorant relationshjp leads to natural aynthesia/analysis acheme (Figure 32) which nonconservatively exploits open-loop nondestabilizing uncertainties.


ANALYSIS
Stratonovich synthesis = approximation to majorant analysis

Eigure 31. The Stratonovich Synthesis Model Provides a Ejrst Approximation to the Majorant Analysie Bounds

Second member of the hierarchy:

$$
\begin{aligned}
& n \cdot \hat{(Q)}+\hat{M}[\hat{()}]=(; \hat{\hat{Q}})\rangle+\langle\hat{(\imath)}\rangle(; \boldsymbol{j}+1) \\
& J-\operatorname{tr}[\hat{Q} R] \leq 2 \sum_{K=1}^{r}\left(\operatorname{tr} \hat{P}_{K}\right)\left(\dot{j}\langle\hat{Q})_{K K}\right. \\
& 0=A \hat{Q}+\hat{Q} A^{T}+1 \|[\hat{Q}]+V \\
& 0=A T \hat{P}+\hat{P} A+l l[[\hat{P}]+R
\end{aligned}
$$

where:

$$
\begin{aligned}
& \widehat{\hat{Q}\rangle}>\triangleq \text { off-diagonal part of } \widehat{\mathrm{Q}} \\
& \|[.]=\text { Stratonovich model operator }
\end{aligned}
$$

- Tighter bound-incorporates more intormation on $\mathbf{A}$ and $\mathbf{G}$
- Predicts stability when $\left(A+A^{T}\right)$ stable, $\mathbf{G}=-\mathbf{G}^{\boldsymbol{T}}$
- "Nominal" performance, $\operatorname{tr}$ [ $\hat{Q} \mathrm{R}$ ], given by Stratonovich model

Pigure 32. The Refined Majorant Bound Incorporates atratonovich Covariance Model

## 5. Numerical Algorithms and Examples

Practical design of controllers je only possible when efficient, reliable algorithma are avajlable. Indeed, the optimal projection equations are readily solvable and have been applied to a wide variety of examples. Numerical results appear in [A3-A6, A8, A11, A12, A14-A16, A18, A19, A21-A24, A26-A2B, A30-A33, A39, A42, A44, A46]. Two distinctly different algorithma have been developed thus far, namely, an iterat jve method and a homotopy algorithm.

The iterative method, developed in [A14, A16, A44] and further studied in [20,21], is outlined in Figure 33. The nice feature of this approach is that only a atandard LQG software package is required for ils implementation. The basic motivation for the wethod is the observation that the main source of coupling is via the terma involving $T_{1}$. The coupling is absent, of course, when $\tau$ is the identity, i.e.. LQG. Note also that the terms involving $\tau_{\perp}$ are small when $R_{2}$ and $V_{2}$ are large. i.e.. when control cost is high and the measurement noise is significant. This case, which yields low-authority controllers. is approximately characterized by decoupled control-deaign and controller-reduction operations. Thus it is not surprising that LQG reduction techniques are most successful when controller authority is low.

Since the $\tau_{L}$ terwi occasion the greatest difficulty, it appears advantageous to bring them into play sradually. This can be accomplished by fixing $T$ after each iteration to yield updated values of $Q$. $P, \dot{Q}$ and $\hat{P}$. Furthermore, $T$ is introduced gradually by means of $\alpha$ to reduce its rank.

The crucial atep of the algorithm concerns the construction of the projection ftrom the pseudogramians $\vec{Q}$ and $\vec{P}$. Specifically, $T$ can be characterized (aee (A22]) aa the sum of eigenprojections of $\bar{Q}$, where each choice of eigenprojections may correspond to a local extremal. However. the necessary conditions do not specify which eigenprojections are to be selected for obtaining a particular local solution. Nevertheless, there do exist useful methods for constructing 7 . For example, component-cost decomposition methods ([22]) when applied within the optimal projection framework often permit efficient identification of the global optimum.

Although the iterative method is convenient to use because it utilizes readily evailable software, it is aboptimal in the menae that it does not fully exploit the structure of the equatione. Specifically, while the iterative method addresses a syetem of four non matrix equations, careful analysis reveals that because of the rank deficiency of the projection the froblem can be recast as four $n_{c} x n^{2}$ equations. Hence, when $n_{c}$ is much amaller than $n$. which is clearly the most


Figure 33. The Iterative Method far Solving the Robustified OPUS Design Equations Requires Only an LQG Software Package and Involves Refinement of the Optimal Projectiont
desirable case for practical implementation, there exists considerable opportunity for increased computational efficiency. Furthermore, and most satisfying, the computational complexity decreases with $n_{c}$ as is intuitively expected below that required by LQG design. Hence the optimal projection approach has computational complexity less then LQG reduction aethods for which LQG is but the first step.
S. Richter ([23,A46]) has developed a howotopy algorithm which fully exploits this crucial structure. Numerical experiments thus far have shown that considerable computational savings can be achieved over the iterative method. Eurthermore, by applying topological degree theory to investigare the branches and character of the local extremals. it can be shown that the maximum number of possible extremals is

$$
\binom{\sin (n, m, 2)}{n_{c}}
$$

if $n_{c} \leq \min (n, m, l)$ or 1 otherwise. Hence in most practical cases the equations support relatively small number of solutions.

Both the iterative method and the homotopy algorithm have been applied to a design problem involving an 8 th-order flexible structure originally due to D. Enns and concidered in [1]. Specifically, a variety of LQG reduction methods are compared in [1] for a range of controller authorities. These methods include:

1. Enns: This method is e frequency-weighted, balanced realization technique appicable to either model or controller reduction.
2. Glover: This method utilizes the theory of Hankel norm optimal approximation for controller reduction.
3. Davis and Skelton: This is a modification of compensator reduction via balancing which addresses the case of unstable controllers.
4. Yousuff and Skelton: This is a further modification of balancing for handling stable or unstable controllers.
5. Lju and Anderson: In place of using a balanced approximation of the compensator transfer functic: directly, this method approximates the component parts of fractional representarion of the compensator.

All of the above methods proceed by first obtaining the full-order LQG compencaror design for a high-order atate-space model and then reducing the dimenajon of the resulting LQG compensator.

Figure 34 summarizes the resulrs reported in [1] for the above LQG reduction methode along with results obtained using the iterative method for colving the optimal projection equations. Here $q_{2}$ is a scale factor for the plant disturbance noise affecting controller authority. Clearly, LQG reduction methods experience increasing difficulry as authority increases. i.e.. as the $\boldsymbol{r}_{\perp}$ terms become incressingly more important in coupling the control and reduction operations. Eor the low authority cases, the optimal projection calculations, which were performed on a Harrie 4800 minicomputer, appeared to incur roughly the same computarional burden as the LQG reduction methods. Although the optimal projection computational burden increasea with authority, comparison with the LQG reduction methods is not meaningful because of the difficulty experienced by these methods in schieving closed-loop etability. See [A44] for further detalla and for comparisonc involving transient response.

The homotopy algorithe was also applied to the example considered in [1]. One of the main goals of the development effort was to extend the range of disturbance ineansity or. equivalently, observer bandwideh, out beyond $q_{2}=2000$. To this end, second-order ( $n_{c}=2$ ) controllers were obrained with relatively little compuration for $q_{2}=10,000,100,000$ and $1,000,000$. In addition. the performance of each reduced-order controller was within 255 of LQG. These cases can curely be expected to present a nontrivial challenge to both the LQG reduction methoda and the jeerative optimal projection method.

Numerjcal solution of the robustified optimel projection equations hes been carried out for eeveral examples. For illuarrative purposes a $2 \times 2$ example was considered in [A26] and the resuita illuatrated in Eigure 35 indicate performancel robuseness tradeoffa possible. The varieble left-shift technique was applied in [A19] to the NASA SCOLE problem with frequency uncertainties. The robustnesa of LQG and two robustified designs is shown in Eigure 36. The plots illustrate the degradation in performance due to sjmultaneous perturbation of all modal frequencies. Note that LQG is rendered unatable by $+5 \%$ frequency perturbation while a highauthority robustified design improves this region ro $+8 \%$. The low-authority design increatee thia region significantly while acrificing 67 nominal performance.


S - The cloced-loop syatem is atable
U - The closed-loop syetem is unstable

Figure 34. The Optinal Projection Approach Was Compared to Several LQG Reduction Techniques Over a Range of Controller Authorities for en Example of Enns


Figure 35. The Robutified Oprimal Projection Equations Provide Robustnesa/Performance Tradeoffa for a Highly Sensitive Nowinal LQG Desigu

RMS LOS ERROR DEGRADATION (X)


Figure 36. The Stratonovich Model Robugtifies the LQG Design for the NASA SCOLE Model with Uncertsin Model Erequencies

## 6. Additionel Extencions

The robustified optimal projection design machinery has been further extended to encompase a larger number of design cases arising in practical application. Here ve merely list the extensions:

1. Discrete-tite and sampled-deta controllers ([A28, A30, A34.A35]).
2. Decentralized contrallers ([A39]).
3. Nonstrictly proper controllers ([A37]).
4. Diatributed paraneter systems ([A25]).

## 7. Concluding Remarks

The machinery provided by OPUS for deaigning active controllers for flexible structures has been reviewed. The basic machinery is a system of coupled Riccati and Lyapunov equations which directly generalize the classical LQG result to include both a conatraine on controller order and a model of parameter uncertainty. The overall approsch thus encompasses all major design tradeoffs arising in vibration-suppresaion applications. Substantial numerical experience has been gained through an iterative method requiring only an LQG software package and. more recently, by means of highly efficjent homotopy algorithm developed by S. Richter. The overall approach opens the door for effective design of implementable controllers for large precision space atructures.

Acknowledgment. We wish to thank th. Jill M. Straehla for the excellent preparation of this paper.

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## APPENDIX B: Fixed-Structure Design

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# The Optimal Projection Equations for Model Reduction and the Relationships Among the Methods of Wilson, Skelton, and Moore 

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#### Abstract

First-order mecessary conditions for quadratically optimal reduced-order modeling of linear time-iavariant systems are derived in the form.of a pair of modified Lyapunov equations coupled by an oblique projection which determines the optimal reduced-order model. This form of the necessary conditions considerably simplifies previous results of Wilson [1] and clearly demonstrates the quadratic exiremality and monoptinallity of the balancing method of Moore [2]. The possible existence of multiple solutions of the optimal projection equations is demonstrated and a relaxation-type algorithm is proposed for computing these local extreme. A component-cost analysis of the modelerror criterion similar to the approach of Skelton [3] is utilized at each iteration to direet the algoritim to the global minimum.


## I. Introduction

THE problem of approximating a high-order linear dynamical system with a relatively simpler system, i.e., the model-reduction problem, has received considerable attention in recent years. Among the my riad papers devoted to this problem are the notable contributions of Wilson [1]. Moore [2], and Skelton [3] with which the present paper is concerned. In his 1970 paper, Wilson proposed an optimality-based approach to model reduction which invoives minimizing the steady-state, quadratically weighted' output error when the original system and reduced-order model are subjected to white-noise inputs. For the resulting parameter-optimization problem, he obtained first-order

[^1]necessary conditions which have the form of an aggregation (as, e.g., [4]) and which involve the solution of two Lyapunov equations each of order $n+n_{m}$, where $n$ and $n_{m}$ are the orders of the original and reduced-order models, respectively [5]. [6].
Some time later, Moore proposed a quite different approach to model reduction based upon system-theoretic arguments as opposed to optimality criteria. Using the eigenvalues of the product of the controllability and observability gramians (which satisfy $n \times n$ Lyapunov equations), his method identifies subsystems which contribute little to the impulse response of the overal: system. Such "weak" subsystems are thus eliminated to obtain a reduced-order model. This technique, known as balancing, has been vigorously developed in the recent literature [7][11]. Since this approach is completely independent of optimality considerations, there is, of course. no expectation that such reduced-order models are in any sense optimal.
A third approach to model reduction, proposed by Skelton [3], [12], also utilizes a quadratic optimality criterion as in [1]. However, rather than proceeding from necessary conditions as does Wilson, Skelton determines for a given basis the contribution (cost) of each state in a decomposition of the error criterion and truncates those with the least value. Although this approach is guided by optimality considerations, no rigorous guarantee of optimality is possible because of dependence on the choice of state space basis.
The present paper has five main objectives, the first of which is to show how the complex optimality conditions of Wilson can be transformed without loss of generality into much simpler and more tractable forms. The transformation is facilitated by exploiting the presence of an oblique (i.e., nonorthogonal) projection which was not recognized in $[1]^{2}$ and which arises as a direct consequence of optimality. The resulting "optimal projection equations" constitute a coupled system of two $n \times n$

[^2]modified Lyapunov equations (see (2.13), (2.14) or (2.21), (2.22)] whose solutions are given by a pair of rank $-n_{m}$ controllability and observability pseudogramians. The highly structured form of these equations gives crucial insight into the set of local extrema satisfying the first-order necessary conditions.

The second objective of the paper is to show how the optimal projection equations provide a rigorous extremality context for Moore's balancing method and to clearly demonstrate its quadratic nonoptimality. Although for some problems the "weak subsystem' hypothesis leads to a nearly optimal reduced-order model, we construct examples for which the reduced-order model obtained from the balancing method is much worse with respect to the least-squares criterion than the quadratically optimal reducedorder model. In general, all that can be said is that the presence of a weak subsystem indicates that the reduced-order model obtained by truncation in the balanced basis may be in the proximity of an extremal of the quadratically optimal model-reduction problem; however, this extremal may very well be a global maximum. It should be noted that in a recent paper [13] Kabamba has used bounds on the model error to demonstrate the quadratic nonoptimality of the balancing method.

The third objective of the paper is to demonstrate via an example the mechanism responsible for the existence of multiple extrema of the optimal model-reduction problem. By characterizing the optimal projection as a sum of rank-1 eigenprojections of the product of the rank-deficient pseudogramians. it is immediately clear that the first-order necessary conditions of the problem are ambiguous in the sense that they fail to specify which $n_{m}$ eigenprojections comprise the optimal projection corresponding to a solution (i.e., global minimum) of the optimal model-reduction problem. Specifically, since the pseudogramians can be rank deficient in $\left(n_{m}\right)=n!/ n_{m}!\left(n-n_{m}\right)!$ ways. there may be precisely this many extremal projections corresponding to an identical number of local extrema.

The fourth objective of the paper is to propose a numerical algorithm for solving the optimal projection equations by exploiting their structure and taking advantage of the available insights. By expressing the modified Lyapunov equations in the form of "standard" Lyapunov equations, an iterative relaxation-type algorithm is developed. The crucial aspect of the proposed algorithm involves extracting an oblique projection at each step from the product of the solutions of the Lyapunov equations. Since $\left(n_{m}^{n}\right)$ rank $-n_{m}$ projections can be extracted from the product of two $n \times n$ positive-definite matrices, it is quickly evident that the criterion by which the $n_{m}$ eigenprojections are chosen determines which of the numerous local extrema will be reached. If. for example, the projection is chosen in accordance with the $n_{m}$ largest eigenvalues of the product of the solutions of the Lyapunov equations, then it should not be surprising in view of the previous discussion that a global maximum may very well be reached. In this case, the first iteration of this algorithm involves Lyapunov equations whose solutions are the controllability and observability gramians and the eigenvalues in question are precisely the squares of the second-order modes [2, p. 24]. Thus, the first iteration coincides with the (nonoptimal) balancing approach of [2].

Since the optimal projection equations are a consequence of iifferential (local) properties, it should not be expected that they alone would possess the inherent ability to identify the global minimum. Moreover, because of the number of local extrema. second-order necessary conditions appear to be useless. Instead, we investigate an approach which chooses the eigenprojections according to a component-cost analysis of the model-error criterion. This technique can lead to a global minimum by effectively eliminating the local extrema which have considerably greater cost than the global minimum. This approach is philosophically identical to the component cost analysis of Skelton [3], [12]. Essentialiy, then, component cost analysis is utilized at each iteration to direct the algorithm to the global minimum. Although our application of this technique is admittedly heuristic, it should be noted that it is essentially proposed as a device for efficiently
"sorting out" the local extrema which satisfy the otherwise mathematically rigorous necessary conditions. Hence, we propose component cost analysis as a crucial step in bridging the gap between local extremality and global optimality.

It should be pointed out that neither the numerical algorithm proposed in this paper nor the iterative algorithm developed in [4] and [5] has been proven to be convergent. The principal contribution of the present paper, however, is not a particular proposed algorithm but rather the revelations concerning the structure of the first-order necessary conditions. The proposed numerical algorithm should be considered but a prelude to a full investigation into numerical algorithms for the optimal projection equations. It should also be noted that the presence of the optimal projection was not exploited in developing the iterative algorithms in [4] and [5] (in fact, it did not even appear in [1]) and hence crucial insight into local extrema was lacking.
The fifth and last objective of the paper is to point out the connection berween the optimal projection equations for model reduction obtained herein and the first-order necessary conditions obtained recently for two closely related problems, namely, reduced-order state estimation and fixed-order dynamic compensation.
The plan of the paper is as follows. Section II begins with general notation and definitions followed by the model-reduction problem statement and the main theorem which presents the optimal projection equations for model reduction. A series of remarks considers various aspects of the main theorem and sets the stage for discussing connections with [1] and [2]. Section III contains a detailed discussion of the sense in which the optimal projection equations simplify the necessary conditions given in [1], and Section IV shows how the approach of [2] is approximately extremal. Section V presents a simple example which clearly displays the possible existence of multiple extrema satisfying the optimal projection equations. This example shows that the balancing method of [2] may lead to a nonoptimal reduced-order model and suggests a heuristic procedure for selecting the eigenprojections comprising the projection corresponding to the global minimum. i.e., the optimal projection. In Section VI, a numerical algorithm for solving the optimal projection equations is proposed and applied to an illustrative example considered previously in [1] and [2] as well as to some interesting examples considered recently by Kabamba in [13]. Related results on reduced-order dynamic compensation and state estimation are briefly reviewed in Section VII and suggestions for further research are given in Section VIII. The proof of the main theorem appears in the Appendix.

## II. Problem Statement and Main Result

The following notation and definitions will be used throughout the paper:

| ${ }^{\prime}{ }_{Z}{ }^{T}$ | $r \times r$ identity matrix transpose of vector or matrix $Z$ |
| :---: | :---: |
| $Z^{-T}$ | $\left(Z^{T}\right)^{-1}$ or $\left(Z^{-1}\right)^{T}$ |
| $\rho(Z)$ | rank of matrix $Z$ |
| $\begin{aligned} & \operatorname{tr} Z \\ & \\|Z\\| \end{aligned}$ | trace of square matrix $Z$ $\left[\operatorname{tr} Z Z^{T}\right]^{1 / 2}$ |
| $Z_{i j}$ | (i, j)-element of matrix $Z$ |
| diag ( $\alpha_{1}, \cdots, \alpha_{r}$ ) | $r \times r$ diagonal matrix with listed diagonal elements |
| $E_{1}$ | matrix with unity in the ( $i, i$ ) position and zeros elsewhere |
| 5 | expected value |
| fin, iffes | real numbers, $r \times s$ real matrices |
|  | left half plane |
| nonnegative-definite matrix | symmetric matrix with nonnegative eigenvalues |
| positive-definite | symmetric matrix with positive |
| matrix | eigenvalues |

semisimple matrix
nonnegative
semisimple matrix
positive-semisimple matrix
positive-diagonal
matrix
$n, m, \ell, n_{m}$
$x, u, y, x_{m}, y_{m}$
$A, B, C$
$A_{m}, B_{m}, C_{m}$
$R, V$
matrix similar to a diagonal matrix [14, p. 10]
matrix similar to a nonnegativedefinite matrix matrix similar to a positivedefinite matrix
diagonal matrix with positive diagonal elements
positive integers, $1 \leq n_{m} \leq n$
$n, m, l, n_{m}, \ell$-dimensional vectors
$n \times n, n \times m, \ell \times n$ matrices
$n_{m} \times n_{m}, n_{m} \times m, \ell \times n_{m}$ matrices
$\ell \times \ell, m \times m$ positive-definite matrices

We consider the following problem.
Optimal Model-Reduction Problem: Given the controllable and observable system

$$
\begin{gather*}
\dot{x}=A x+B u,  \tag{2.1}\\
y=C x \tag{2.2}
\end{gather*}
$$

find a reduced-order model

$$
\begin{gather*}
\dot{x}_{m}=A_{m} x_{m}+B_{m} u,  \tag{2.3}\\
y_{m}=C_{m} x_{m} \tag{2.4}
\end{gather*}
$$

which minimizes the quadratic model-reduction criterion ${ }^{3}$

$$
J\left(A_{m}, B_{m}, C_{m}\right) \triangleq \lim _{t \rightarrow \infty} \mathbb{E}\left[\left(y-y_{m}\right)^{T_{R}} R\left(y-y_{m}\right)\right]
$$

where the input $u(t)$ is white noise with positive-definite intensity $V$. To guarantee that $J$ is finite, it is assumed that $A$ is stable and we restrict our attention to the set of admissible reduced-order models

$$
Q \triangleq\left\{\left(A_{m}, B_{m}, C_{m}\right): \quad A_{m} \text { is stable }\right\}
$$

Since the value of $J$ is independent of the internal realization of the transfer function corresponding to (2.3) and (2.4), we further restrict our attention to the set
$a_{+} \triangleq\left\{\left(A_{m}, B_{m}, C_{m}\right) \in Q:\right.$

- $\left(A_{m}, B_{m}\right)$ is controllable and $\left(A_{m}, C_{m}\right)$ is observable $\}$. The following lemma is needed for the statement of the main result.

Lemma 2.1: Suppose $Q, P \in \mathbb{R}^{n \times n}$ are nonnegative definite. Then $Q P$ is nonnegative semisimple. Furthermore, if $\rho(Q P)=$ $n_{m}$ then there exist $G, \Gamma \in \cap^{n_{m} \times n}$ and positive-semisimple $M \in$ $\beta_{\beta^{n} m} \times n_{m}$ such that

$$
\begin{gather*}
\phi P=G^{\top} M \Gamma  \tag{2.5}\\
\Gamma G^{T}=I_{r_{m}} . \tag{2.6}
\end{gather*}
$$

Proof: By [14, Theorem 6.2.5, p. 123], there exists $n \times n$ invertible $\Phi$ such that the nonnegative-definite matrices $D_{0} \triangleq$ $\Phi Q^{r}$ and $D_{\beta} \hat{=} \phi^{-r} \beta \dot{\phi}-1$ are both diagonal. Hence, $D_{Q} D_{\beta}$ is nonnegative definite and $Q \dot{P}=\Phi^{-1} D_{Q} D_{\rho} \phi$ is nonnegative semisimple. Next introduce $n \times n$ orthogonal $U$ to effect a rearrangement of basis if necessary so that

$$
Q P=\Phi\left[\begin{array}{ll}
\Lambda & 0 \\
0 & 0
\end{array}\right] \Phi^{-1}
$$

[^3]where $\Phi \triangleq \$ U$ and $n_{m} \times n_{m} \Lambda$ is positive diagonal. Hence, for all $\boldsymbol{n}_{\boldsymbol{m}} \times \boldsymbol{n}_{\boldsymbol{m}}$ invertible $\boldsymbol{S}$.
\[

\phi P=\Phi\left[$$
\begin{array}{l}
S \\
0
\end{array}
$$\right]\left(S^{-1} \Lambda S\right)\left[S^{-1} \quad 0\right] \Phi^{-1}
\]

and thus, (2.5) and (2.6) hold with $G=\left[S^{T} 0\right] \Phi^{T}, M=S^{-1} \Lambda S$ and $r=\left[S^{-1} 0\right] \Phi^{-1}$.
For convenience in stating the main theorem, we shall refer to $G, \Gamma \in\left[\beta^{n} m_{m} \times n\right.$ and positive-semisimple $M \in \beta^{n}{ }^{n} \times n_{m}$ satisfying (2.5) and (2.6) as a ( $G, M, \Gamma$ )-foctorization of $Q P$. Also, define the positive-definite controllability and observability gramians

$$
\begin{aligned}
& W_{c} \triangleq \int_{0}^{\infty} e^{A^{\prime} B V B^{T_{e}} e^{T_{t}} d t,} \\
& W_{0} \triangleq \int_{0}^{\infty} e^{A^{T}} C^{\top} R C e^{A t} d t,
\end{aligned}
$$

which satisfy the dual Lyapunov equations

$$
\begin{align*}
& 0=A W_{c}+W_{c} A^{T}+B^{T} B^{T},  \tag{2.7}\\
& 0=A^{T} W_{o}+W_{0} A+C^{T} R C . \tag{2.8}
\end{align*}
$$

Main Theorem: Suppose $\left(A_{m}, B_{m}, C_{m}\right) \in Q$, solves the optimal model-reduction problem. Then there exist nonnegativedefinite matrices $\hat{O}, \mathcal{P} \in \mathbb{B}^{n \times n}$ such that, for some $(G, M, \Gamma)$ factorization of $Q P, A_{m}, B_{m}$, and $C_{m}$ are given by

$$
\begin{gather*}
A_{m}=\Gamma A G^{T}  \tag{2.9}\\
B_{m}=\Gamma B  \tag{2.10}\\
C_{m}=C G^{r} \tag{2.11}
\end{gather*}
$$

and such that, with $\tau \triangleq G^{\top} \Gamma$, the following conditions are sati. fied:

$$
\begin{align*}
& \rho(\phi)=\rho(P)=\rho(\emptyset P)=n_{m},  \tag{2:2}\\
& 0=\tau\left[A Q+Q A^{\top}+B V B^{\top}\right],  \tag{2.15}\\
& 0=\left[A^{\top} P+P A+C^{\top} R C\right] \tau . \tag{2.1.}
\end{align*}
$$

Several comments are in order. First, note that the main theorem consists of necessary conditions in the form of two modified Lyapunov equations (2.13) and (2.14) plus rank conditions (2.12) which must possess nonnegative-definite solutions $\phi_{2}$ $\beta$ when an optimal reduced-order model exists. We shall call $\boldsymbol{Q}$ and $\tilde{P}$ the controllability and observability pseudogramians, respectively, since they are analogous to $W_{c}^{\prime}$ and $W_{o}$ and yet have rank deficiency. The modified Lyapunov equations are coupled by the $n \times n$ matrix $T$ which is a projection (idempotent matrix) since

$$
r^{2}=C^{T} T C^{T} T=C^{T} I_{n_{m}} \Gamma=\tau
$$

Note that, in general, $T$ is an oblique projection and not necessarily an orthogonal projection since it may not be symmetric. We shall refer to a projection $\tau$ corresponding to a solution (i.e., global minimum) of the optimal model-reduction problem as an "optimal projection." It should be stressed that the form of the optimal reduced-order model (2.7)-(2.9) is a direct consequence of optimality and not the result of an a priori assumption on the structure of the reduced-order model.

Since the optimal projection equations are first-order necessary conditions for optimality, they ma;' possess multiple solutions corresponding to various local extrema such as local maxima. local minima, saddle points, etc. The following definition will prove useful.

Definition 2.1: Nonnegative-definite $Q, \beta \in\left\{\int^{n \times n}\right.$ are extremal if (2.12)-(2.14) are satisfied. $\left(A_{m}, B_{m}, C_{m}\right) \in Q_{+}$is
extremal if there exist extremal $\boldsymbol{Q}, \boldsymbol{P}$ such that $\left(A_{m}, B_{m}, C_{m}\right)$ is given by (2.9)-(2.11) for some ( $G, M, \Gamma$ )-factorization of $Q P$. The corresponding projection $\tau$ is an extremal projection.

Proposition 2.1: Suppose $\left(A_{m}, B_{m}, C_{m}\right)$ is extremal. Then the model-reduction error is given by ${ }^{4}$

$$
\begin{equation*}
J\left(A_{m}, B_{m}, C_{m}\right)=2 \operatorname{tr}\left[\left(\phi P-W_{r} W_{0}\right) A\right] \tag{2.15}
\end{equation*}
$$

Proof: The proof is given at the end of Appendix A.
Remark 2.1: Noting the identities

$$
\begin{equation*}
-2 \operatorname{\pi r}\left[W_{c} W_{0} A\right]=\operatorname{tr}\left[C^{\top} R C W_{c}\right]=\operatorname{tr}\left[B V B^{\top} W_{0}\right], \tag{2.16}
\end{equation*}
$$

which follow from (2.7) and (2.8), (2.15) can be written for extremal $\left(A_{m}, B_{m}, C_{m}\right)$ as

$$
\begin{align*}
& J\left(A_{m}, B_{m}, C_{m}\right)=2 \operatorname{tr}[Q P A]+\operatorname{tr}\left[C^{\top} R C W_{c}\right] \\
&=2 \operatorname{tr}[\oint P A]+\operatorname{tr}\left[B V B^{\top} W_{o}\right] . \tag{2.17}
\end{align*}
$$

For convenience in the following discussion, let $Q, P, G, M$, $\Gamma$, and $\tau$ correspond to some extremal $\left(A_{m}, B_{m}, C_{m}\right)$. Now observe that if $x_{m}$ is replaced by $S x_{m}$, where $S$ is an arbitrary nonsingular matrix, then an "equivalent" reduced-order model is obtained with ( $A_{m}, B_{m}, C_{m}$ ) replaced by $\left(S A_{m} S^{-1}, S B_{m}\right.$, $C_{m} S^{-1}$. Since $J\left(A_{m}, B_{m}, C_{m}\right)=J\left(S A_{m} S^{-1}, S B_{m}, C_{m} S^{-1}\right)$, one would expect the main theorem to apply also to ( $S A_{m} S^{-1}, S B_{m}$, $C_{m} S^{-1}$ ). Indeed, the following result shows that this transformation corresponds to the alternative factorization $O P=$ $\left(S^{-T} G\right)^{T}\left(S M S^{-1}\right)(S \Gamma)$ and, moreover, that all $(G, M, \Gamma)$ factorizations of $\boldsymbol{Q} \boldsymbol{P}$ are related by an invertible transformation. Proposition 2.2: If $S \in \mathbb{R}^{n_{m} \times n_{m}}$ is invertible, then $\mathcal{G}=$ $S^{-T} G, \Gamma=S \Gamma$ and $\vec{M}=S M S^{-1}$ satisfy

$$
\begin{align*}
& \phi P=G^{T} \bar{M} \Gamma  \tag{2.5}\\
& \Gamma G^{r}=I_{n_{m}} \tag{2.6}
\end{align*}
$$

Conversely, if $\bar{G}, \Gamma \in\left[\beta^{n_{m} \times n}\right.$ and invertible $\bar{M} \in \mathbb{R}^{n_{m} \times n_{m}}$ satisfy (2.5)' and (2.6)', then there exists invertible $S \in \mathbb{R}^{n_{m} \times n_{m}}$ such that $\bar{G}=S^{-T} G, \hat{\Gamma}=S \Gamma$ and $\bar{M}=S M S^{-1}$.

Proof: The first part is immediate. The second part follows by taking $S \triangleq \bar{M}^{-1} \Gamma G^{T} M$, noting $S^{-1}=M \Gamma G^{T} M^{-1}$ and using the identities $\Gamma G^{T} M \Gamma \bar{G}^{T}=\bar{M}$ and $M \Gamma^{C^{r}}=\Gamma \bar{G}^{T} \bar{M}^{T}$.

The next result shows that there exists a similarity transformation which simultaneously diagonalizes $Q P$ and $\tau$.

Proposition 2.3: There exists invertible $\Phi \in \mathbb{B}^{n \times n}$ such that

$$
\begin{align*}
& Q=\Phi^{-1}\left[\begin{array}{cc}
\Lambda_{Q} & 0 \\
0 & 0
\end{array}\right] \Phi^{-T}, \quad \hat{P}=\Phi^{T}\left[\begin{array}{cc}
\Lambda_{P} & 0 \\
0 & 0
\end{array}\right] \Phi,  \tag{2.18}\\
& \phi P=\Phi^{-1}\left[\begin{array}{cc}
\Lambda & 0 \\
0 & 0
\end{array}\right] \Phi, \quad T=\Phi^{-1}\left[\begin{array}{cc}
I_{n_{m}} & 0 \\
0 & 0
\end{array}\right] \Phi, \quad \text { (2.19 } \tag{2.19a,b}
\end{align*}
$$

where $\Lambda_{Q}, \Lambda_{\beta} \in \mathbb{R}^{n_{m} \times n_{m}}$ are positive diagonal, $\Lambda \triangleq \Lambda_{Q} \Lambda_{\beta}$ and the diagonal elements of $\Lambda$ are the eigenvalues of $M$. Consequently,

$$
\begin{equation*}
Q=\tau \emptyset, P=P_{\tau} \tag{2.20}
\end{equation*}
$$

Proof: By [14. Theorem 6.2.5, p. 123], and by (2.12), there exists $n \times n$ invertible $\Phi$ such that (2.18) holds and thus (2.19a) also holds. Define

$$
C=\left[\begin{array}{ll}
I_{n_{m}} & 0
\end{array}\right]^{-T}, \bar{M}=\Lambda \quad \text { and } \Gamma=\left[\begin{array}{ll}
I_{n_{m}} & 0
\end{array}\right] \Phi
$$

so that (2.5)' and (2.6)' are satisfied. By the second part of Proposition 2.2 there exists invertible $S \in \|^{n_{m}} \times{ }^{n}$ such that $G=$

4 The expressions (2.15)-(2.17) and (2.23)-12.24) will be used in Sections
$S^{T} \bar{G}, M=S^{-1} \bar{M} S$ and $\Gamma=S^{-1} \Gamma$. Now (2.19b) follows from

$$
\tau=G^{T} \Gamma=G^{T} \Gamma=\Phi^{-1}\left[\begin{array}{cc}
I_{n_{m}} & 0 \\
0 & 0
\end{array}\right] \Phi .
$$

It is useful to present an alternative form of the optimal modelreduction equations (2.13) and (2.14). For convenience, define the notation

$$
\tau_{\perp} \triangleq I_{n}-\tau
$$

Proposition 2.4: Equations (2.13) and (2.14) are equivalent, respectively, to

$$
\begin{align*}
& 0=A Q+Q A^{T}+B V B^{T}-\tau_{1} B V B^{T_{T}} T_{1},  \tag{2.21}\\
& 0=A^{T} S_{+}+P A+C^{\top} R C-\tau_{1}^{T} C^{T} R C T_{1} . \tag{2.22}
\end{align*}
$$

Proof: By (2.20), (2.21) $=(2.13)+(2.13)^{\tau}+(2.13) \tau$ and (2.13) $=\boldsymbol{T}(2.21)$. Similarly, (2.14) and (2.22) are equivalent.

Remark 2.2: Noting the identities

$$
\begin{equation*}
-2 \operatorname{tr}[\oint P A]=\operatorname{tr}\left[C^{\top} R C Q\right]=\operatorname{tr}\left[B V B^{\top} P\right] \tag{2.23}
\end{equation*}
$$

which follow from (2.20)-(2.22), (2.17) can be written for extremal $\left(A_{m}, B_{m}, C_{m}\right)$ as

$$
\begin{equation*}
J\left(A_{m}, B_{m}, C_{m}\right)=\operatorname{tr}\left[C^{\top} R C\left(W_{c}-\phi\right)\right]=\operatorname{tr}\left[B V B^{T}\left(W_{0}-P_{0}\right)\right] \tag{2.24}
\end{equation*}
$$

To facilitate the discussion in the following sections, we consider the change of basis $\hat{x} \triangleq \Phi x$, where $\Phi$ is given by Proposition 2.3. Writing (2.1) and (2.2) as

$$
\begin{gather*}
\dot{\vec{x}}=\tilde{A} \hat{x}+\hat{B} u,  \tag{2.25}\\
y=\hat{C} \hat{x}, \tag{2.26}
\end{gather*}
$$

where

$$
\hat{A} \leqq \Phi A \Phi^{-1}, \quad E \leqq \Phi B, \quad C \cong C \Phi^{-1}
$$

(2.9)-(2.11) become

$$
\begin{gather*}
A_{m}=\hat{A} \bar{A} \mathcal{C}^{T}  \tag{2.27}\\
B_{m}=\hat{f} B  \tag{2.28}\\
C_{m}=\hat{C} \mathcal{G}^{T} \tag{2.29}
\end{gather*}
$$

where

$$
\mathrm{f} \triangleq \Gamma \Phi^{-1}, \quad G \triangleq G \Phi^{T}
$$

satisfy

$$
G^{T}{ }^{\top}=\left[\begin{array}{cc}
I_{n_{m}} & 0  \tag{2.30}\\
0 & 0
\end{array}\right], \quad \Gamma G^{T}=I_{n_{m}}
$$

Note that (2.30) implies

$$
f=\left[\begin{array}{ll}
S & 0
\end{array}\right], \quad G=\left[\begin{array}{ll}
S^{-T} & 0 \tag{2.31}
\end{array}\right]
$$

for some $n_{m} \times n_{m}$ invertible $S$. Partitioning

$$
\begin{aligned}
\hat{x}=\left[\begin{array}{c}
\hat{X}_{m} \\
\hat{X}_{2}
\end{array}\right], \quad \hat{A}=\left[\begin{array}{cc}
\hat{A}_{m} & \hat{A}_{m 2} \\
\hat{A}_{2 m} & \hat{A}_{22}
\end{array}\right], \\
\hat{B}=\left[\begin{array}{c}
B_{m} \\
\hat{B}_{2}
\end{array}\right], \quad C=\left[\begin{array}{ll}
C_{m} & C_{2}
\end{array}\right],
\end{aligned}
$$

where $\hat{X}_{m} \in \mathbb{X}^{n_{m}}$ and $\hat{A}_{m}, \mathcal{B}_{m}$ and $\hat{C}_{m}$ are $n_{m} \times n_{m}, n_{m} \times m$ and
$\ell \times n_{m}$, respectively, (2.27)-(2.29) and (2.31) yield

$$
A_{m}=S A_{m} S^{-1}, \quad B_{m}=S B_{m}, \quad C_{m}=C_{m} S^{-1}
$$

This shows that the optimal reduced-order model (modulo a state transiormation) can be obtained by truncating the last $n-n_{m}$ states of the original system when it is expressed in the basis with respect to which $\boldsymbol{Q}$ and $\boldsymbol{\beta}$ have the diagonal forms

$$
\left[\begin{array}{cc}
\Lambda_{0} & 0 \\
0 & 0
\end{array}\right] \text { and }\left[\begin{array}{cc}
\Lambda_{P} & 0 \\
0 & 0
\end{array}\right] .
$$

Since the optimal projection $\tau$ has the simple form

$$
\left[\begin{array}{cc}
I_{I_{m}} & 0 \\
0 & 0
\end{array}\right]
$$

in this basis, we shall refer to (2.25) and (2.26) as an optimal projection realization of (2.1) and (2.2). Note that when (2.21) and (2.22) are expanded in an optimal projection basis (i.e., a basis corresponding to an optimal projection realization) they assume the form

$$
\begin{align*}
0= & \hat{A}_{m} \Lambda_{Q}+\Lambda_{Q} \hat{A}_{m}^{T}+\hat{B}_{m} V B_{m}^{T},  \tag{2.32}\\
& 0=\hat{A}_{2 m} \Lambda_{Q}+\hat{B}_{2} V B_{m}^{T},  \tag{2.33}\\
0= & \hat{A}_{m}^{T} \Lambda_{P}+\Lambda_{P} \hat{A}_{m}+C_{m}^{T} R C_{m},  \tag{2.34}\\
& 0=\Lambda_{P} \hat{A}_{m 2}+C_{m}^{T} R C_{2} . \tag{2.35}
\end{align*}
$$

If $\Phi$ in Proposition 2.23 is replaced by

$$
\left[\begin{array}{cc}
\left(\Lambda_{Q}^{-1} \Lambda_{P}\right)^{1 / 4} & 0 \\
0 & I_{n-n_{m}}
\end{array}\right] \Phi,
$$

which corresponds to a change of basis for the reduced-order model obtained by truncation, then $\Lambda_{Q}$ and $\Lambda_{\rho}$ are both replaced by $\left(\Lambda_{Q} \Lambda_{\beta}\right)^{1 / 2}$ and hence this can be called a balanced optimal projection basis, utilizing the terminology of [2]. Thus, in a balanced optimal projection realization, $\Lambda_{\varrho}$ and $\Lambda_{\rho}$ appearing in (2.32)-(2.35) are equal.

The nexi result provides an interesting closed-form characterization of an extremal projection in terms of the Drazin generalized inverse of $\phi \hat{P}$. Since $(\hat{Q} \hat{P})^{2}=G^{\top} M^{2} \Gamma$, and hence $\rho(\phi \hat{P})^{2}=$ $\rho(Q P)$, the "index" of $\phi P$ (see [15, p. 121]) is 1. In this case, the Drazin inverse is traditionally called the group inverse and is denored by ( $Q P)^{\prime}\left[15, \mathrm{p}\right.$. 124]. Since, as is easily verified, $(Q P)^{\prime \prime}$ $=G^{\top} M^{-1} \Gamma$, (2.6) leads to the following result.
Proposition 2.5: An extremal projection $\tau$ is given by

$$
\begin{equation*}
\tau=Q P(Q P)^{\wedge} . \tag{2.36}
\end{equation*}
$$

An alternative representation for an extremal projection will prove useful for developing a numerical algorithm for solving (2.21) and (2.22). If $Q, P \in G^{r \times r}$ are nonnegative definite then by Lemma 2.1 QP is nonnegative semisimple and thus there exists invertible $\Psi \in \mathbb{E g}^{\prime \times r}$ such that

$$
Q P=\Psi^{-1} \Omega \Psi
$$

where $\Omega=\operatorname{diag}\left(\omega_{1}, \cdots, \omega_{r}\right)$ and $\omega_{i} \geq 0$ are the eigenvalues of $Q P$. Now define the ith eigenprojection [16, p. 41]

$$
\Pi_{i}[Q P] \leqq \Psi^{-1} E_{i} \Psi
$$

which is a rank-1 oblique projection. Note that $Q P$ has the decomposition

$$
Q P=\sum_{i=1}^{r} \omega_{i} \Pi_{i}(Q P]
$$

Proposition 2.6: An extremal projection T is given by

$$
\begin{equation*}
\left.T=\sum_{i=1}^{n_{n}} \Pi_{i} Q Q P\right] \tag{2.37}
\end{equation*}
$$

where the ith eigenprojection $\Pi_{i}[Q P[$ corresponds to the $i$ th nonzero eigenvalue $\lambda$; of $\emptyset \boldsymbol{P}$.

## III. Relationship to Wilson's Form of the Necessary Conditions

The optimal model-reduction problem considered in the previous section is identical to the problem considered by Wilson in [1] with the minor exception that he sets $R=I_{1}$. In [1] $G$ and $\Gamma$ are denoted by $\theta_{2}^{T}$ and $\theta_{1}$, (2.6) appears as (15), and (2.9)-(2.11) are given by (14a, b). Note that in (1], $\theta_{1}$ and $\theta_{2}$ depend upon the solutions of a pair of $\left(n+n_{m}\right) \times\left(n+n_{m}\right)$ Lyapunov equations [see (7), (9) of [1] or (A.2), (A.3) of the present paper] whose coefficients and nonhomogeneous terms depend in turn on $A_{m}$, $B_{m}$, and $C_{m}$ [see (A.10)-(A.15)]. The advantage of the $n \times n$ optimal projection equations (2.21) and (2.22) over the form of the necessary conditions given in [1] [see (A.10)-(A.15)] is that the optimal projection equations are independent of $A_{m}, B_{m}$, and $C_{m}$. Hence, this permits the development of numerical algorithms which avoid the need to choose starting values for $A_{m}, B_{m}$, and $C_{m}$. To see this, note that although the unknowns $A_{m}, B_{m}$, and $C_{m}$ appear explicity in (A.10)-(A.15), all data in the optimal projection equations (2.13) and (2.14) are known except for the solutions $Q$ and $\bar{P}$. Moreover, the optimal projection $\tau$, which was not recognized in [1], can be seen to play a fundamental role by coupling the modified Lyapunov equations (2.21) and (2.22) and determining (since $\tau=G^{T} \Gamma$ ) $A_{m}, B_{m}$, and $C_{m}$ in (2.7)-(2.9).

## IV. Relationship to Moore's Balancing Method

In contrast to Wilson's method for model reduction which is based on optimality principles, the approach due to Moore [2] relies on system-theoretic ideas. The main thrust of this approach "is to eliminate any weak subsystem which contributes little to the impulse response matrix" [2, p. 26]. The concept of a "weak subsystem" is defined by means of a dominance relation [2, p. 28] involving similarity invariants called second-order modes. Moore evaluates reduced-order models obtained in this way by computing the relative error in the impulse response given for MIMO systems by [2, p. 29]

$$
\epsilon\left(A_{m}, B_{m}, C_{m}\right) \triangleq\left[\int_{0}^{\infty}\left\|H_{e}(t)\right\|^{2} d t / \int_{0}^{\infty}\|H(t)\|^{2} d t\right]^{1 / 2},
$$

where $H_{e}(t) \triangleq H(t)-H_{m}(t), H(t) \triangleq R^{1 / 2} C e^{\wedge}\left(B V^{1 / 2}\right.$ and $H_{m}(t)$ $\triangleq R^{1 / 2} C_{m} e^{A} m^{1} B_{m} V^{1 / 2}$. To discuss this approach in the context of the optimal model-reduction problem, we assume that $V=I_{m}$ and $R=$ 1.

Proposition 4.1: Suppose $\left(A_{m}, B_{m}, C_{m}\right) \in Q$. Then

$$
\begin{align*}
\epsilon\left(A_{m}, B_{m}, C_{m}\right) & =\left[-\frac{1}{2} J\left(A_{m}, B_{m}, C_{m}\right) / \operatorname{tr}\left(W_{c} W_{0} A\right)\right]^{1 / 2} \\
& =\left[J\left(A_{m}, B_{m}, C_{m}\right) / \operatorname{tr}\left(C^{\top} R C W_{c}\right)\right]^{1 / 2} \\
& =\left[J\left(A_{m}, B_{m}, C_{m}\right) / \operatorname{tr}\left(B V B^{\top} W_{0}\right)\right]^{1 / 2} \tag{4.1}
\end{align*}
$$

Proof: The result follows from (A.1). (A.8), and (A.9) which hold without regard to either optimality or extremality.
Note that Proposition 4.1 shows that the relative error in the impuise response is minimized precisely when $J\left(A_{m}, B_{m}, C_{m}\right)$ is minimized. Actually, this result is to be expected since, as shown in [1], $J$ can be obtained alternatively by taking $u(t)$ to be an impulse at $t=0$.

To draw interesting comparisons with the results of [2], choose
$n \times n$ invertible $\Psi$ such that $\Psi W_{c} \Psi^{r}$ and $\Psi^{-T} W_{0} \Psi^{-1}$ are both diagonal and hence

$$
\begin{equation*}
W_{c} W_{o}=\boldsymbol{\Psi}^{-1} \Sigma^{2} \Psi \tag{4.2}
\end{equation*}
$$

where $\Sigma \triangleq \operatorname{diag}\left(\sigma_{1}, \cdots, \sigma_{n}\right)$ and the second-order modes $\sigma_{i}$ (i.e., the positive square roots of the eigenvalues of $W_{c} W_{0}$ ) satisfy $\sigma_{1} \geq \sigma_{2}$ $\geq \cdots \geq \sigma_{n}>0$. This transformation corresponds to replacing (2.1), (2.2) by

$$
\begin{gather*}
\dot{x}=\bar{A} \bar{x}+B u,  \tag{4.3}\\
y=\bar{C} \bar{x}, \tag{4.4}
\end{gather*}
$$

where

$$
\begin{equation*}
\bar{x} \triangleq \Psi x, \quad \bar{A} \triangleq \Psi A \Psi^{-1}, \quad B \triangleq \Psi B, \quad C \triangleq C \Psi^{-1} \tag{4.5}
\end{equation*}
$$

The transformed system (4.3), (4.4), called a principal axis realization [17], can further be chosen so that

$$
\begin{equation*}
\Psi W_{c} \Psi^{T_{ \pm}} \boldsymbol{\Psi}^{-r_{W_{0}} \Psi^{-1}=\Sigma, ~} \tag{4.6}
\end{equation*}
$$

i.e., the balanced realization. Using (4.5), (2.7) and (2.8) become

$$
\begin{align*}
& 0=\bar{A} \Sigma+\Sigma \bar{A}^{T}+B V B^{T}  \tag{4.7}\\
& 0=\bar{A}^{T} \Sigma+\Sigma \bar{A}+C^{T} R C \tag{4.8}
\end{align*}
$$

The model-reduction procedure suggested in [2] involves partitioning

$$
\begin{aligned}
& \vec{x}=\left[\begin{array}{l}
x_{m} \\
X_{2}
\end{array}\right], \quad \vec{A}=\left[\begin{array}{cc}
A_{m} & A_{m 2} \\
A_{2 m} & A_{22}
\end{array}\right], \\
& \vec{B}=\left[\begin{array}{l}
B_{m} \\
B_{2}
\end{array}\right], \quad C=\left[\begin{array}{ll}
C_{m} & C_{2}
\end{array}\right],
\end{aligned}
$$

where $\bar{X}_{m} \in \mathbb{F}_{J^{n}}$ and $\bar{A}_{m}, \bar{B}_{m}$, and $\bar{C}_{m}$ have corresponding dimension, and extracting the reduced-order model ( $\bar{A}_{m}, \bar{B}_{m}$, $\boldsymbol{C}_{m}$ ). Hence, the reduced-order model ( $\bar{A}_{m}, \boldsymbol{B}_{m}, \boldsymbol{C}_{m}$ ) is extracted from (4.3), (4.4) in essentially the same way the optimal reducedorder model ( $A_{m}, B_{m}, C_{m}$ ) is extracted from (2.25), (2.26). To see how the optimal-projection realization compares to a princi-pal-axis realization, first note that (2.13) and (2.14) are satisfied by $\boldsymbol{Q}=W_{c}$ and $\boldsymbol{P}=W_{0}$ when the rank conditions (2.10) are ignored. Indeed, since $W_{c}$ and $W_{o}$ are positive definite, the rank conditions (2.12) do not hold. If, however, the system (2.1), (2.2) is expressed in the balanced coordinate system (4.3), (4.4) (so that $W_{c}=W_{o}=\Sigma$ ), then the assumption $\sigma_{n_{m}} \rightarrow \sigma_{n_{m}+1}$ implies that $\rho\left(W_{c}\right), \rho\left(W_{o}\right)$ and $\rho\left(W_{c} W_{0}\right)$ are "approximately" equal to $n_{m}$ and thus, in this sense, condition (2.10) is satisfied. This observation leads to the suggestion that when $\sigma_{n_{m}}>\sigma_{n_{m}+1}, W_{c}$ and $W_{o}$ are approximations to solutions $\boldsymbol{Q}$ and $\bar{P}$ of the optimal projection equations and the reduced-order model ( $A_{m}, B_{m}, C_{m}$ ) of Moore is an approximation to some extremal $\left(A_{m}, B_{m}, C_{m}\right)$. There is no guarantee, of course, that any particular extremum corresponds to the global minimum, or even to a local minimum.

## V. Existence of Multiple Extrema and Component-Cost Ranking

In this section, we show by means of a simple example that the optimal projection equations may possess nonunique solutions corresponding to multiple extrema, e.g., local minima or maxima. We also show how decomposing the cost can identify the global minimum from among the numerous extrema. To begin, let $m=\ell=n, R=V=I_{n}$,

$$
A \leq \operatorname{diag}\left(-\alpha_{1}, \cdots,-\alpha_{n}\right)
$$

where $\alpha_{i}>0, i=1, \cdots, n$, and suppose $B$ and $C$ are such that

$$
B B^{T}=\operatorname{diag}\left(\beta_{1}, \cdots, \beta_{n}\right), \quad C^{\top} C=\operatorname{diag}\left(\gamma_{1}, \cdots, \gamma_{n}\right),
$$

where $\beta_{i}>0, \gamma_{i}>0, i=1, \cdots, n$. Hypothesizing diagonal solutions $\boldsymbol{Q}$ and $\mathcal{P}$ of (2.21) and (2.22) leads to

$$
\phi_{i i}=\frac{\beta_{i}}{2 \alpha_{i}} \delta_{i}, \quad P_{i i}=\frac{\gamma_{i}}{2 \alpha_{i}} \delta_{i},
$$

where each $\delta_{l}, i=1, \cdots, n$ is either zero or one and exactly $n_{m}$ of the $\delta_{i}$ 's are equal to one. Hence $r=\operatorname{diag}\left(\delta_{1}, \cdots, \delta_{n}\right)$. Note that there are $\binom{n}{n_{m}}$ such solutions of the optimal projection equations corresponding to $\binom{n}{n_{m}}$ local extrema.

Since

$$
W_{c}=-\frac{1}{2} A^{-1} B B^{T}, \quad W_{0}=-\frac{1}{2} A^{-1} C^{\tau} C, \quad \phi=\tau W_{c}, \quad P=\tau W_{0}
$$

and $A, W_{c}$, and $W_{o}$ commute, (2.15) becomes

$$
J\left(A_{m}, B_{m}, C_{m}\right)=-\frac{1}{2} \operatorname{tr} \tau_{1} A^{-1} B B^{T} C^{\tau} C
$$

Hence,

$$
\begin{equation*}
J\left(A_{m}, B_{m}, C_{m}\right)=\sum_{i=1}^{n} \zeta_{i}\left(1-\delta_{i}\right) . \tag{5.1}
\end{equation*}
$$

where

$$
\zeta_{i} \triangleq \beta_{i} \gamma_{i} / 2 \alpha_{i} .
$$

To minimize $J$, it is clear that $\delta_{i}$ should be chosen to be unity for the largest $n_{m}$ elements of the set $\left\{\zeta_{i}\right\}_{i=1}^{n}$ and zero otherwise. Although this choice is not necessarily unique, it does yield a global minimum. Note that choosing $\delta_{i}=1$ is equivalent to selecting a particular eigenprojection $\Pi_{i}\left[W_{c} W_{o}\right]$ corresponding to the eigenvalue $\beta_{i} \gamma_{i} / 4 \alpha_{i}^{2}$.

Remark 5.1: The expression in (5.1) can be regarded as a decomposition of the cost in terms of the state variables. The idea of deleting states based on their "component costs" is precisely the "component cost analysis" approach of Skelton [3], [12].

Using the example, it is easy to see that the balancing method of [2], which selects eigenprojections based upon the magnitude of the eigenvalues of $W_{c} W_{o}$, i.e., the (squares of the) second-order modes, may yield a grossly suboptimal reduced-order model. To this end, let
$\alpha_{1}=1, \quad \alpha_{2}=10^{6}, \quad \beta_{1}=1, \quad \beta_{2}=10^{6}, \quad \gamma_{1}=1, \gamma_{2}=10^{3}$
so that

$$
\zeta_{1}=0.5, \quad \zeta_{2}=500 .
$$

Clearly, $J$ is minimized $\left(J=\zeta_{1}\right)$ by choosing $\delta_{1}=0, \delta_{2}=1$, which corresponds to truncating the first state variable. If, however, the method of [2] is utilized, then judging by the secondorder modes

$$
\sigma_{1}=0.5, \quad \sigma_{2}=(2.5)^{1 / 2} \cdot 10^{-2}=0.012
$$

the second state variable should be deleted. This, however, corresponds to choosing $\delta_{1}=1, \delta_{2}=0$ with the higher $\operatorname{cost} J=$ $\zeta_{2}$. The fact that the balancing approach of [2] fails to determine a solution of the optimal model-reduction problem should not be surprising in view of the fact that the error criterion plays no role in the balancing technique.

Although the above solution exploted the simple structure of this example, it is clear that choosing the global minimum from among the local extrema involves an eigenprojection decomposi-
tion of the cost $J$. To extend this idea to more general systems, e invoke the following heuristic approximation.
Approximation 5.I: Let $\Psi$ define the balanced basis as in (4.6). Then also approximately defines a balanced optimal projection basis, i.e.,

$$
\begin{equation*}
\psi Q \psi T_{=\Psi} \psi-T P_{\psi}-1=T \Sigma^{2} . \tag{5.2}
\end{equation*}
$$

where extremal

$$
\begin{equation*}
\bar{f} \triangleq \Psi \Psi^{-1}=\operatorname{diag}\left(\delta_{1}, \cdots, \delta_{n}\right) \tag{5.3}
\end{equation*}
$$

and

$$
\delta_{i} \in\{0,1\}, \quad \sum_{i=1}^{n} \delta_{i}=n_{m} .
$$

Proposition S.I: If Approximation 5.1 holds for extremal ( $A_{m}, B_{m}, C_{m}$ ) then, with $\bar{f}_{+} \triangleq I_{n}-\bar{j}$,

$$
\begin{align*}
J\left(A_{m}, B_{m}, C_{m}\right) & =-2 \operatorname{tr}\left[\overline{\bar{I}}_{\perp} \Sigma^{2} A\right] \\
& =2 \sum_{i=1}^{n}-\sigma_{i}^{2} A_{i i}\left(1-\delta_{i}\right) . \tag{5.4}
\end{align*}
$$

Remark 5.2: From (4.7) and (4.8), it follows that (5.4) can be written either as

$$
\begin{align*}
J\left(A_{m}, B_{m}, C_{m}\right) & =\operatorname{tr}\left[\tilde{r}_{\perp} \Sigma B V B^{r}\right] \\
& =\sum_{i=1}^{n} \sigma_{i}\left(B V B^{T}\right)_{m}\left(1-\delta_{i}\right) \tag{5.5}
\end{align*}
$$

or

$$
\begin{align*}
J\left(A_{m}, B_{m}, C_{m}\right) & =\operatorname{tr}\left[\bar{F}_{\perp} \Sigma C^{\top} R C\right] \\
& =\sum_{i=1}^{n} \sigma_{i}\left(C^{\top} R C\right)_{i}\left(1-\delta_{i}\right) . \tag{5.6}
\end{align*}
$$

Hence, Approximation 5.1 leads to the following component-cost ranking (again. in the sense of Skelton [3], [12]) of the ( $n_{m}^{\prime}$ ) exirema satisfying the optimal projection equations.
Component-Cost Ranking: Assume Approximation 5.1 is valid and choose the eigenprojections comprising extremal $\bar{j}$ such that
$\delta_{1}=1$, if $-\sigma_{i}$ is among the $n_{m}$
largest elements of the set $\left\{-\sigma_{r}^{2} \bar{A}_{m}\right\}_{-1}$;

$$
\delta_{i}=0, \quad \text { otherwise. }
$$

For comparison purposes, we shall also consider the following ranking of the eigenprojections based upon the eigenvalues of $W_{C} W_{0}$ (i.e., second-order modes).
Eigenvalue Ranking: Choose the eigenprojections comprising extremal $\bar{f}$ such that
$\delta_{i}=1$, if $-\sigma_{i}^{2} \bar{A}_{i n}$ is among the $n_{m}$
largest elements of the set $\{-\sigma,\}^{n}=1$;

$$
\delta_{i}=0, \quad \text { otherwise. }
$$

Remark 5.3: The observation that the second-order modes alone may be a poor guide to determining an optimal reducedorder model has recently been made in [13] where bounds on the model-error criterion were given involving both the second-order modes and suitable weights called balanced gains. II can be seen
.. m Proposition 5.1 that the role of balanced gains in our approach is played by the elements $-\sigma_{i} A_{i /}$ when Approximation 5.1 holds. It can also be seen that the balanced gains of Kabamba yield bounds on the component costs of Skelton.

## VI. Numercal Solution of the Optimal Prcjection Equations

Insofar as the ultimate aim of any model-reduction technique is to permit the development of numerical procedures for reducing high-order models, the optimal projection equations, comprising a coupled system of modified Lyapunov equations, appear promising in this regard. Therefore, we present an iterative computational algorithm that exploits the structure of these equations and the available insights. The reader is strongly reminded that the proposed algorithm is but a first attempt at solving these new equations and alternative algorithms may yet be devised. The basis of this algorithm is the ability to write the modified Lyapunov equations (2.21), (2.22) in the form of "standard" Lyapunov equations (6.1), (6.2) such that the pseudogramians $Q$ and $\mathcal{F}$ are extracted at the final step (6.6). It follows from (2.32)(2.35) that (2.21), (2.22) are indeed equivalent to (6.1), (6.2) (with $k=\infty$ ) and (6.6).
Algorithm:
Step 1) Initialize $\tau^{(0)}=I n$.
Step 2) Solve for $\tilde{\underline{Q}}^{(1)}, p^{n}(k)$

$$
\begin{align*}
& 0=\left(A-\tau^{(k)} A \tau_{\perp}^{(k)}\right) \dot{Q}^{(k)}+\dot{Q}^{\prime(k)}\left(A-\tau^{(k)} A \tau_{\perp}^{(k)}\right)^{T}+B V B^{T},  \tag{6.1}\\
& 0=\left(A-\tau_{\perp}^{(k)} A \tau^{(k)}\right)^{\top(k)}+\dot{P}^{(k)}\left(A-\tau_{\perp}^{(k)} A \tau^{(k)}\right)+C^{\top} R C . \tag{6.2}
\end{align*}
$$

Step 3) Balance

$$
\begin{gather*}
\Phi^{(k)} \dot{\Phi}^{(k)}\left(\Phi^{(k)}\right)^{\tau}=\left(\Phi^{(k)}\right)^{\tau} \bar{P}^{(k)}\left(\Phi^{(k)}\right)^{-1}=\Sigma^{(k)},  \tag{6.3}\\
\Sigma^{(k)}=\operatorname{diag}\left(\sigma_{1}^{(k)}, \cdots, \sigma_{n}^{(k)}\right), \quad \sigma_{1}^{(k)} \geq \sigma_{2}^{(k)} \geq \cdots \geq \sigma_{n}^{(k)} \geq 0 .
\end{gather*}
$$

Step 4) If $k>1$ check for convergence

$$
\begin{equation*}
e_{k} \triangleq\left[\frac{\operatorname{tr}\left(C^{\top} R C W_{c}\right)-\operatorname{tr}\left(C^{\top} R C \tau^{(k)} \underline{Q}^{(k)}\left(\tau^{(k)}\right)^{T}\right)}{\operatorname{tr}\left(C^{\top} R C W_{c}\right)}\right]^{1 / 2} \tag{6.4}
\end{equation*}
$$

If $\left|e_{k}-e_{k-1}\right|<$ tolerance then go to step 8); else continue; Step 5) Select $n_{m}$ eigenprojections

$$
\begin{aligned}
& \Pi_{i, 1}\left[\delta^{(k)} \dot{p}^{(k)}\right], \cdots, \Pi_{i_{-1}}\left[\delta^{(k)} \dot{p}^{(k)}\right], \\
& \left.\Pi_{i} \dot{\mathscr{O}}^{(k) \dot{\boldsymbol{P}}(k)}\right] \triangleq \Phi^{(k)} E_{i}\left(\Phi^{(k)}\right)^{-1} .
\end{aligned}
$$

Step 6) Update

$$
\begin{equation*}
\tau^{(k+1)}=\sum_{i=1}^{n_{i}} \Pi_{i,}\left[Q^{(k)} \hat{p}^{(k)}\right] . \tag{6.5}
\end{equation*}
$$

Step 7) Check for convergence; if not, increment $k$ and return to Step 2).

Step 8) Set

$$
\begin{equation*}
\phi=\gamma^{(\infty)} \bar{Q}\left(\tau^{(\infty)}\right)^{\tau}, \quad \rho_{=}=\left(\tau^{(\infty)}\right)^{\tau} \hat{\mathcal{F}}^{(\infty)} . \tag{6.6}
\end{equation*}
$$

For convenience, we shall adopt the notation $\left(A_{m}^{(k)}, B_{m}^{(k)}, C_{m}^{(k)}\right)$. where $k>0$, to denote the reduced-order model obtained as a result of applying the projection $\tau^{(k)}$, and we define (see Section IV)

$$
e_{k} \triangleq \in\left(A_{m}^{(k)}, B_{m}^{(k)}, C_{m}^{(k)}\right)
$$

i.c., the relative error associated with $\left(A_{m}^{(k)}, B_{m}^{(k)}, C_{m}^{(k)}\right)$. Note that. in general, $\epsilon_{k} \neq e_{k}$ since $e_{k}$ denotes the relative error only for an extremum, i.e., when convergence has been reached.

It should be clear from the discussion in the previous section that the crucial step of the algorithm is Step 5)-the choice of the eigenprojections. For the examples which follow, we shall invoke consistently at Step 5) either the component-cost ranking based upon Approximation 5.1 or the eigenvalue ranking.

Remark 6.1: Note that in the special case $R=I_{m}$ and $V=I_{\mathrm{f}}$. the first iteration of the algorithm yields $\boldsymbol{Q}^{(0)}=W_{c}, \boldsymbol{\beta}^{(0)}=W_{0}$. If, at Step 5), we choose $i_{r}=r, r=1, \cdots, n_{m}$, i.e., the eigenprojections are selected according to the eigenvalue ranking, then ( $\left.A^{(1)}, B^{(1)}, C_{p}^{(1)}\right)$ is precisely the reduced-order model obtained from balancing.
We shall first consider the following example which was treated by both Wilson and Moore. In this example, and those that follow, assume $R=I_{m}, V=I_{\mathrm{f}}$.

Example 6. $i:$ :
$A=\left[\begin{array}{rrrr}0 & 0 & 0 & -150 \\ 1 & 0 & 0 & -245 \\ 0 & 1 & 0 & -113 \\ 0 & 0 & 1 & -19\end{array}\right], \quad B=\left[\begin{array}{l}4 \\ 1 \\ 0 \\ 0\end{array}\right], C=\left[\begin{array}{llll}0 & 0 & 0 & 1\end{array}\right]$.
Table I summarizes the results obtained for the three cases $n_{m}$ $=3,2,1$ utilizing the eigenvalue ranking. In each case, the proposed algorithm converged linearly in less than eight iterations and, in each case, improvement is evident over previously published results. As pointed out in [2], Wilson's result seems to imply a lack of final convergence. For this example, the balancing approach yields a reduced-order model close to the globa! minimum.

We now turn to a pair of interesting examples considered in [13].

Example 6.2:

$$
A=\left[\begin{array}{cc}
-0.005 & -0.99 \\
-0.99 & -5000
\end{array}\right], \quad B=\left[\begin{array}{c}
1 \\
100
\end{array}\right], \quad C=B^{T} .
$$

Table II summarizes the results obtained using the eigenvalue ranking and Table III gives the results when the component-cost ranking is used. It is clear that the former method directs the algorithm to the global maximum whereas the latter approach yields the global minimum.

Example 6.3:

$$
A=\left[\begin{array}{cc}
-0.25 & -0.4 \\
-0.4 & -0.72
\end{array}\right], \quad B=\left[\begin{array}{c}
1 \\
1.2
\end{array}\right], \quad C=B^{\tau} .
$$

Table IV reports the results obtained using either the compo-nent-cost ranking or the eigenvalue ranking which agree for this example. If the alternative eigenprojection is selected then, as expected, the algorithm converges to a global maximum (see Table V). The interesting aspect of this example, as discussed in [13], is that the error $\epsilon_{1}=0.5245$ (see [13]) for the reduced-order model obtained by either eigenprojection ranking is actually greater than $\epsilon_{1}=0.3849$ obtained by choosing the alternative reduced-order model. This situation seems to indicate that proper eigenprojection selection based upon a cost decomposition is able to direct the algorithm to the global minimum in cases for which the starting values are not nearby.

## VII. The Optimal Projection Equations for Fixed-Order Dynamic Compensation and Reduced-Order State Estimation

We briefly discuss the relationship between the optimal projection equations for model reduction and analogous results for reduced-order control and estimation: problems.
Fixed-Order Dynamic-Compensation Problem: Given the controlled system

$$
\begin{gather*}
\dot{x}=A x+B u+w_{1}  \tag{7.1}\\
y=C x+w_{2} \tag{7.2}
\end{gather*}
$$

table I
RELATIVE ERROR $e_{-}=\epsilon_{-}$

|  | Wilson [1] | Moore [2] | Optimal Projection <br> Equations |
| :---: | :---: | :--- | :--- |
| 3 | - | 0.001311 | 0.001306 |
| 2 | 0.04097 | 0.03938 | 0.03929 |
| 1 | - | 0.4321 | 0.4268 |

TABLE II
example 6.2 with eigenvalue ranking

| $k$ | $e_{k}$ |
| :---: | :---: |
| 1 | 0.9950371897 |
| 2 | 0.9950371691 |
|  | 0.0950371690 |

TABLE II
EXAMPLE 6.2 WITH COMPONENT-COST RANKING

| $k$ | $e_{k}$ |
| :---: | :---: |
| 1 | 0.0995037 |
| 2 | 0.0995449 |
| 3 | 0.0995924 |
| 4 | 0.0996520 |
| 5 | 0.0997346 |
| 6 | 0.0998648 |
| 7 | 0.1001125 |
| 8 | 0.1007724 |
| 9 | 0.1054569 |
| 10 | 0.0982006 |
| 11 | 0.0975409 |
| 12 | 0.0975342 |
| 13 | 0.0975330 |
| 14 | 0.0975329 |
| table iv <br> EXAMPLE 6.3 Using Either Ranking |  |
| $k$ | $e_{k}$ |
| 1 | 0.646996 |
| 2 | 0.418341 |
| 3 | 0.220994 |
| 4 | 0.177276 |
| 5 | 0.176576 |

table V
EXample 6.3 With the Opposite ranking

| $k$ | $e_{k}$ |
| :---: | :---: |
| 1 | 0.7624928516 |
| 2 | 0.9999999961 |
| 3 | 0.9999999975 |
| $\cdot$ | $\cdot$ |
| $\cdot$ | $\cdot$ |
| 29 | 0.9999999999 |

design a fixed-order dynamic compensator

$$
\begin{gather*}
\dot{x}_{c}=A_{d} x_{c}+B_{c} y_{1}  \tag{7.3}\\
u=C_{c} x_{c} \tag{7.4}
\end{gather*}
$$

which minimizes the performance criterion

$$
\begin{equation*}
J\left(A_{c}, B_{c}, C_{c}\right) \triangleq \lim _{r \rightarrow \infty}\left[x^{\top} R_{1} x+u^{\top} R_{2} u\right], \tag{7.5}
\end{equation*}
$$

where $u \in \mathbb{R}^{m}, x_{c} \in \mathbb{R}^{n} c, n_{c} \leq n_{1} w_{1}$ is white disturbance noise, $w_{2}$ is nonsingular white observation noise, $\boldsymbol{R}_{1}$ is nonnegative definite, and $\boldsymbol{R}_{\mathbf{2}}$ is positive definite.

Necessary conditions characterizing optimal ( $A_{c}, B_{c}, C_{c}$ ) have been developed in [18]-[22] along the same lines as the main theorem. These conditions, called the optimal projection equations for fixed-order dynamic compensation, consist of four matrix equations (two modified Riccati equations and two modified Lyapunov equations) coupled by a projection. The modified Riccati equations, not surprisingly, are similar in form to the covariance and cost Riccati equations of LQG theory and the modified Lyapunov equations are similar to the optimal model-reduction equations (2.13) and (2.14). Hence, while the modified Riccati equations govern optimal estimation and optimal control, the additional modified Lyapunov equations characterize "optimal reduction." The important fact that all four equations are coupled supports the view that optimal fixed-order dynamic compensators cannot, in general, be designed by means of a stepwise procedure, e.g., by either open-loop model reduction followed by LQG or LQG followed by closed-loop model reduction.

Midway between the model-reduction and fixed-order dy-namic-compensation problems lies the following problem.
Reduced-Order State-Estimation Problem: Given the observed system

$$
\begin{align*}
& \dot{x}=A x+w_{1},  \tag{7.6}\\
& y=C x+w_{2}, \tag{7.7}
\end{align*}
$$

design a reduced-order state estimator

$$
\begin{gather*}
\dot{x}_{\mathrm{e}}=A_{e} x_{e}+B_{e} y,  \tag{7.8}\\
y_{e}=C_{r} x_{e}, \tag{7.9}
\end{gather*}
$$

which minimizes the estimation criterion

$$
J\left(A_{e}, B_{e}, C_{e}\right) \triangleq \lim _{i \rightarrow \infty} E\left[\left(L x-y_{e}\right)^{\tau} R\left(L x-y_{e}\right)\right],
$$

where $x_{e} \in \mathbb{R}^{n_{e}}, L \in \mathbb{B}^{p \times n_{e}}$ and $L$ identifies the states, or linear combinations of states, whose estimates are desired. The order $n_{e}$ of the estimator state $x_{e}$ is determined by implementation constraints, i.e., by the computing capability available for realizing (7.8) and (7.9) in real time.

In view of the results already given, it should not be surprising (see [23]) that the optimal projection equations for reduced-order state estimation form a system of three matrix equations (a pair of modified Lyapunov equations along with a single modified Riccati equation) coupled by a projection which determines the gains of the optimal reduced-order estimator. This intrinsic coupling between the "operations" of optimal estimation (the modified Riccati equation) and optimal model reduction (the pair of modified Lyapunov equations) stresses the fact that reduced-order estimators designed by means of either model reduction followed by "full-order"' state estimation or full-order estimation followed by estimator reduction will generally not be optimal for the given order.

## Vill. Directions for Further Research

The most important area of research invoives the further development of algorithms for solving the optimal projection
equations. Although proving local convergence of the proposed algorithm appears possible, the more important problem is achieving global optimality via the component cost approach. Although the global minimum was attained for all examples attempted by the authors, it remains to treat considerably more complex systems.
An interesting extension of the main theorem involves the case in' which the original system (2.1), (2.2) is a distributed parameter system, e.g., a pertial differential equation or a functional differential equation. This generalization, which has been referred to as the "ultimate reduced-order problem" [24], may lead to the efficient generation of high-order discretizations for such systems. All of the mathematical machinery required to generalize the main theorem to this case has already been applied to fixed-order dynamic compensation in [25].

## IX. Conclusion

First-order necessary conditions for quadratically optimal reduced-order modeling of a linear time-invariant plant are expressed in the form of a pair of $n \times n$ modified Lyapunov equations coupled by an oblique projection. This form of the necessary conditions considerably simplifies the original form given by Wilson in [1] and clearly reveals the possible presence of numerous extrema. The balancing method of Moore given in [2] is shown to yield a reduced-order model that is "close" to an extremal given by the necessary conditions. A numerical example shows, however, that this extremal may very well be the global maximum rather than the desired global minimum. An algorithm is proposed which exploits the presence of the optimal projection and computes the various local extrema by the choice of eigenprojections comprising the projection. A component-cost ranking of the eigenprojections, which is very much in the spirit of Skelton's method in [3] and [12], is used to direct the algorithm to the global optimum.

It should be pointed out that Moore's balancing appears to have strong ties with the $L_{\infty}$ reduction problem via the Hankel norm [29]. Alternative settings for the Hankel operator, however, seem to indicate connections to the quadratic problem [30]. Finally, the robustness problem for reduced-order modeling, estimation, and control in a quadratic setting is discussed in [31].

## Appendix

## Proof of the Main Theorem

Introducing the augmented system

$$
\begin{gathered}
\dot{x}=\bar{A} \bar{x}+\overline{B u}, \\
\bar{y}=C \bar{x},
\end{gathered}
$$

where

$$
\begin{gathered}
\bar{x} \leqq\left[\begin{array}{c}
x \\
x_{m}
\end{array}\right], \quad \bar{y} \grave{y}-y_{m}, \\
A \triangleq\left[\begin{array}{cc}
A & 0 \\
0 & A_{m}
\end{array}\right], \quad B \triangleq\left[\begin{array}{c}
B \\
B_{m}
\end{array}\right], \quad C \leqq\left[\begin{array}{ll}
C & -C_{m}
\end{array}\right] .
\end{gathered}
$$

leads to the expression

$$
\begin{equation*}
J\left(A_{m}, B_{m}, C_{m}\right)=\operatorname{tr} Q \hat{R} \tag{A.I}
\end{equation*}
$$

where $\boldsymbol{R} \triangleq \mathcal{C}^{\top} R \mathbb{R}$ and the nonnegative-definite steady-state covariance $\boldsymbol{Q}$ of $\boldsymbol{x}$ is given by the (unique) solution of

$$
\begin{equation*}
0=A Q+Q A^{T}+V \tag{A.2}
\end{equation*}
$$

with $\bar{D} \leqq \bar{B} V \bar{B}^{T}$. To minimize (A.1) subject to the constraint (A.2), form the Lagrangian

$$
L\left(A_{m}, B_{m}, C_{m}, Q\right) \triangleq \operatorname{tr}\left|\lambda Q R+\left(A \bar{Q}+\emptyset A^{T}+\bar{V}\right) \mathscr{F}\right|
$$

with multipliers $\lambda \geq 0$ and $\dot{P} \in \mathbb{R}^{\left(n+n_{m}\right) \times\left(n+n_{m}\right)}$. Since $Q_{+}$is an open set, the standard Lagrange multiplier rule can be applied.

Using formulas for computing partial derivatives [26], it follows that

$$
0=L_{Q}=A^{T} \bar{P}+\bar{P} \bar{A}+\lambda R .
$$

Since $\lambda=0$ implies $\bar{P}=0$ (recall $A$ is stable), we can take $\lambda=1$ without loss of generality. Hence, $\dot{P}$ is the (unique nonnegativedefinite) solution of

$$
\begin{equation*}
0=A^{T} \tilde{P}+\bar{P} A+\tilde{R} \tag{A.3}
\end{equation*}
$$

Again using formulas from [26] and performing some manipulation, it follows that

$$
\begin{gather*}
0=L_{A_{m}}=Q_{12}^{T} P_{12}+Q_{2} P_{2}  \tag{A.4}\\
0=L_{B_{m}}=2\left(P_{12}^{T} B+P_{2} B_{m}\right) V  \tag{A.S}\\
0=L_{C_{m}}=2 R\left(C_{m} Q_{2}-C Q_{12}\right) \tag{A.6}
\end{gather*}
$$

where $\boldsymbol{Q}$ and $\tilde{P}$ have been partitioned as

$$
\boldsymbol{Q}=\left[\begin{array}{ll}
Q_{1} & Q_{12}  \tag{A.7}\\
Q_{12}^{T} & Q_{2}
\end{array}\right], \quad \dot{P}=\left[\begin{array}{cc}
P_{1} & P_{12} \\
P_{12}^{T} & P_{2}
\end{array}\right] .
$$

Since (as will be seen shortly) $Q_{2}$ and $P_{2}$ are positive definite, define

$$
\begin{equation*}
C \triangleq Q_{2}^{-1} Q_{12}^{T}, \quad \Gamma \triangleq-P_{2}^{-1} P_{12}^{T} \tag{A.8}
\end{equation*}
$$

so that (A.4)-(A.6) become (2.6), (2.10) and (2.11), respectively. Next, define the nonnegative-definite matrices

$$
\begin{equation*}
\phi \triangleq Q_{12} Q_{2}^{-1} Q_{12}^{\top}, \quad P \leqq P_{12} P^{-1} P_{12}^{\top} \tag{A.9}
\end{equation*}
$$

and note that (A.4) implies that (2.5) holds with $M \subseteq Q_{2} P_{2}$. Since $Q_{2} P_{2}=P_{2}^{-1 / 2}\left(P_{2}^{1 / 2} Q_{2} P_{2}^{1 / 2}\right) P_{2}^{1 / 2}, M$ is positive semisimple. The rank conditions (2.12) follow from Sylvester's inequality. Expanding (A.2) and (A.3) yields

$$
\begin{gather*}
0=A Q_{1}+Q_{1} A^{T}+B V B^{T},  \tag{A.10}\\
0=A Q_{12}+Q_{12} A_{m}^{T}+B V B_{m}^{T},  \tag{A.11}\\
0=A_{m} Q_{2}+Q_{2} A_{m}^{T}+B_{m} V B_{m}^{T}  \tag{A.12}\\
0=A^{\tau} P_{1}+P_{1} A+C^{\top} R C  \tag{A.13}\\
0=A^{T} P_{12}+P_{12} A_{m}-C^{\top} R C_{m},  \tag{A.14}\\
0=A_{m}^{T} P_{2}+P_{2} A_{m}+C_{m}^{\top} R C_{m} . \tag{A.15}
\end{gather*}
$$

Since $A_{m}$ is stable and $\left(A_{m}, B_{m}\right)$ is controllable, standard results (e.g., [27, p. 277]) imply that $Q_{2}$ is positive definite. Similarly, $P_{2}$ is positive definite.

It is easy to see at this point that $A_{m}, B_{m}$, and $C_{m}$ are independent of $Q_{1}$ and $P_{1}$ and thus (A.10) and (A.13) can be ignored. Now, substituting (2.10), (2.11) and the identities

$$
\begin{array}{ll}
Q_{12}=\phi \Gamma^{r}, & P_{12}=-P G^{\top} \\
Q_{2}=\Gamma Q \Gamma^{T}, & P_{2}=G F G^{T} \tag{A.17}
\end{array}
$$

into (A.11), (A.12), (A.14), and (A.15) yields

$$
\begin{gather*}
0=A Q \Gamma^{T}+Q \Gamma^{T} A_{m}^{T}+B V B^{T} \Gamma^{T},  \tag{A.18}\\
0=A_{m} \Gamma Q \Gamma^{T}+\Gamma \phi \Gamma^{T} A_{m}^{T}+\Gamma B V B^{T} \Gamma^{T},  \tag{A.19}\\
0=A^{T} \beta G^{T}+P G^{T} A_{m}+C^{\top} R C G^{T}  \tag{A.20}\\
0=A_{m}^{T} G \beta G^{r}+G P G^{T} A_{m}+G C^{\top} R C G^{\Gamma} . \tag{A.21}
\end{gather*}
$$

Computing (A.19)-Г(A.18) implies

$$
A_{m}=\Gamma A Q \Gamma^{T}\left(\Gamma Q \Gamma^{T}\right)^{-1}
$$

which, since $\Gamma \dot{Q} \Gamma^{T}=Q_{2}$, yields (2.9). Alternatively, (2.9) can be obtained from (A.21)-G(A.20).

If we now substitute (2.9) into (A.18)-(A.21) and use the easily verified relations (2.20), it follows that (A.19) $=\Gamma(A .18)$ and $(A .22)=G(A .21)$, and thus (A.19) and (A.21), are redundant. Finally, $G^{T}(A .18)^{T}$ and (A.20) C yield (2.13) and (2.14), respectively. Note that these last multiplications entail no loss of generality since $\rho(G)=\rho(\Gamma)=n_{m}$.

To show that the optimal projection equations entail no loss of generality over (A.2)-(A.6), let $\hat{Q}, \hat{P}$ be extremal and define $Q_{12}$, $Q_{2}, P_{12}, P_{2}$ by (A.16) and (A.17) for some ( $G, M, \Gamma$ )factorization of $Q \mathscr{P}$, and let $Q_{1}, P_{1}$ satisfy (A.10) and (A.13). Then it is straightforward to reverse the steps taken in the proof to arrive at (A.2)-(A.6).

Proof of Proposition 2.1: Extremal $Q, \beta$ leads to $Q, \bar{P}$ as in (A.7) satisfying (A.2)-(A.6). Computing

$$
\begin{aligned}
J\left(A_{m}, B_{m}, C_{m}\right)=\operatorname{tr}\left(Q_{1} C^{\top} R C-2 Q_{12} C_{m}^{\tau} R C\right) & +\operatorname{tr}\left(Q_{2} C_{m}^{r} R C_{m}\right) \\
& =\operatorname{tr}\left[C^{\tau} R C\left(W_{c}-Q\right)\right]
\end{aligned}
$$

noting that (2.13), (2.14) are equivalent to (2.21), (2.22) because of (2.20) and using (2.23), leads to (2.15).

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Demais S. Bernstein (M'82), for a photograph and biography, see p. 376 of the April 1985 issue of this TRANSACTIONS.

# The Optimal Projection Equations for Reduced-Order State Estimation 

## DENNIS S. BERNSTEIN AND DAVID C. HYLAND


#### Abstract

Absfract-First-order mecessary conditions for optimal, steady-state, reduced-order state estimation for a limear, time-invariant plant in the presence of correlated disturbance and nonsiagular mensurement nolse are derived in a mew and bighly simpllfied form. In comtran to the lose matrix Riccati equation arisiag in the fuli-order (Kalman filter) case, the optimal steady-atate reduced-order eatinator th charncterized by throe matrix equations (one modified Ricentl equation and two modified Lyapumov equations) conpled by a projection whose rank is precisely equal to the order of the extimator and which determines the optimal estimator gains. This compliag is a graphic reminder of the anbopimality of proposed approsches involvias cither model reduction followed by "full-order" extimator design or fuili-order extimator design followed by entimator-reduction techniques. The results givea mere complement recentiy obtained resalts which characterize the optimal redaced-order model by means of a pair of coupled modified Lyapunov equations [7] and the optimal fixed-order dyammic compensator by menas of a coupled syatem of two modified Riccati equations and two modified Lyapunov equations [6].


## I. INTRODUCTION

It has recently been shown (see [1]-[7]) that the first-order necessary conditions for the problems of optimal model reduction and optimal fixedorder dymamic compensation can be formulated in terms of an "optimal
Manuscrip received June 28, 1984; revised Sepvember 14, 1984. This work was
mpported is pan by Lincoln Laborstory. M.I.T.
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projection" matrix which arises as a direct consequence of optimality. These necessary conditions, by virtue of their remarkable simplicity, -yield insight into the structure of the optimal design and permit the development of aiternative numerical algorithms [2], [4], [7]. The purpose of this note is to develop analogous first-order necessary conditions for the reduced-order state-estimation problem. Since this problem falls midway between the problems of open-loop model reduction and closed-loop fixed-order dynamic compensation, it is not surprising that the necessary conditions for these problems are correspondingly related. Specifically, while the optimal projection equations for model reduction consist of a system of two matrix equations (a pair of modified Lyapunov equations) and the optimal projection equations for fixed-order dynamic compensation comprise a system of four matrix equations (a pair of modified Lyapunov equations plus a pair of modified Riccati equations), the optimal projection equations for reduced-order state estimation form a system of three matrix equations (a pair of modified Lyapunov equations along with a single modified Riccati equation). In each case the system of matrix equations is coupled by an - oblique projection (idempotent matrix) which determines the gains of the optimal reduced-order system, whether it be model, estimator, or compensator.

The need for designing an optimal reduced-order state estimator for a high-order dynamic system follows directly from real-world constraints on computing capability. A further motivation is the fact that although a system may have many degrees of freedom, it is often the case that estimates of only a small number of state variables are actually required. In the face of these practical motivations, numerous approaches to designing reduced-order state estimators have been proposed. See [8] for a recent review of previous results.

An important fact pointed out in [8] and [9] is that reduced-order estimators designed by means of either model reduction followed by "full-order" state estimation or full-order estimation followed by estimator reduction will nor be optimal for the given order. In the present paper this point is graphically confirmed by the fact that the three matrix equations characterizing the optimal reduced-order state estimator reveal intrinsic coupling (via the optimal projection) between the "operations" of optimal estimation (the modified Riccati equation) and optimal model reduction (the pair of modified Lyapunov equations).

## 11. Problem Statement and Main Result

The following notation and definitions will be used throughout the paper:

$$
\begin{aligned}
& n, l, n_{e, p} p \\
& x_{1}, y, x_{e}, y_{\text {e }} \\
& A, C, L \\
& A_{e}, B_{e}, C_{e} \\
& w_{1}(t), t \geq 0 \\
& w_{2}(t), t \geq 0 \\
& V_{12} \\
& R \\
& J_{r} \\
& \begin{array}{l}
Z^{T} \\
7^{-T}
\end{array} \\
& \mathfrak{N}(Z), \mathcal{R}(Z), p(Z) \\
& \text { 逪 }
\end{aligned}
$$

$$
\begin{aligned}
& \text { stable matrix }
\end{aligned}
$$

nonnegative-definite matrix
positive-definite matrix
nonnegative-semisimple matrix
positive integers, $1 \leq n_{e} \leq n$
$n, l . n_{e}, p$-dimensional vectors
$n \times n, l \times n, p \times n$ real matrices
$n_{f} \times n_{e} n_{e} \times l, p \times n_{e}$ real matrices $n$-dimensional white noise with nonne-gative-definite intensity $V_{1}$
$l$-dimensional white noise with posi-tive-definite intensity $V_{2}$
$n \times /$ matrix satisfying $\left[w_{1}(t) w_{2}(s)^{T}\right]$
$=V_{12} \delta(t-s)$
$p \times p$ positive-definite matrix
$r \times r$ indentity matrix
transpose of vector or matrix $Z$
$\left(Z^{T}\right)^{-1}$ or $\left(Z^{-1}\right)^{T}$
null space, range, rank of matrix $Z$

## expected value

real numbers, $r \times s$ real matrices matrix with eigenvalues in open left half plane
symmetric matrix with nonnegative eigenvalues
symmetric matrix with positive eigenvalues
matrix similar to nonnegative-definite matrix
positive-semisimple matrix
positive-diagonal matrix
matrix similar to a positive-definite matrix
diagonal matrix with positive diagonal elements

We consider the following optimal reduced-order state-estimation problem. Given the system

$$
\begin{align*}
& x=A x+w_{1}  \tag{2.1}\\
& y=C x+w_{2} \tag{2.2}
\end{align*}
$$

design a reduced-order state estimator

$$
\begin{align*}
& \dot{x}_{e}=A_{e} x_{e}+B_{u} y  \tag{2.3}\\
& y_{e}=C_{e} x_{e} \tag{2.4}
\end{align*}
$$

which minimizes the error criterion

$$
J\left(A_{e}, B_{e}, C_{e}\right) \triangleq \lim _{r \rightarrow-}\left\{\left[\left(L x-y_{e}\right)^{\top} R\left(L x-y_{e}\right)\right]\right.
$$

In this formulation the matrix $L$ identifies the states, or linear combinations of states, whose estimates are desired. The order $n_{e}$ of the estimator state $x_{e}$ is determined by implementation constraims, i.e., by the computing capability available for realizing (2.3), (2.4) in real time. Hence, $n_{e}$ is considered to be fixed in what follows and the problem is concerned with determining $A_{e}, B_{r}$, and $C_{r}$.

To guarantee that $J$ is finite it is assumed that $A$ is stable and we restrict our attention to the set of stable reduced-order estimators

$$
Q \triangleq\left\{\left(A_{e}, B_{e}, C_{e}\right): A_{e} \text { is stable }\right\}
$$

Since the value of $J$ is independent of the internal realization of the transfer function corresponding to (2.3) and (2.4), without loss of generality we further restrict our attention to the set of admissible estimators

$$
\cdots\left\{\left(A_{e}, B_{e}, C_{e}\right) \in Q:\right.
$$

( $A_{e}, B_{e}$ ) is controllable and ( $A_{e}, C_{e}$ ) is observable \}.
The following lemma, whose proof is given in [7], is needed for the statement of the main result.

Lemma 2,1: Suppose $Q, P \in\}^{n \times n}$ are nonnegative definite. Then $Q \dot{P}$ is nonnegative semisimple. Furthermore, if $\rho(\phi \dot{P})=n_{e}$, then there exist $G, \Gamma \in\{ \}^{n_{e} \times n}$ and positive-semisimple $M \in \mathbb{R}^{n_{e} \times n_{e}}$ such that

$$
\begin{gather*}
Q P=G^{r} M T  \tag{2.5}\\
\Gamma G^{r}=I_{n_{e}} \tag{2.6}
\end{gather*}
$$

For convenience in stating the Main Theorem we shall refer to $G, \Gamma \in$ [gn $\Omega^{\times n}$ and positive-semisimple $M \in H^{n_{e} \times n_{e}}$ satisfying (2.5) and (2.6) as a $(G, M, \Gamma)$-factorization of $\phi \mathcal{P}$. Funthermore, define the notation

$$
T \triangleq G T, r_{1} \triangleq t_{n}-T
$$

and

$$
Q \leqq Q C^{T}+V_{12}
$$

where $Q \in \beta_{3^{n \times n}}$.
Main Theorem: Suppose $\left(A_{e}, B_{e}, C_{f}\right) \in Q_{+}$solves the optimal reduced-order state-estimation problem. Then there exist nonnegativedefinite matrices $Q, Q, \mathcal{P} \in \mathcal{B}^{n \times n}$ such that, for some ( $G, M, \Gamma$ )factorization of $Q \stackrel{F}{P}, A_{e}, B_{e}$, and $C_{e}$ are given by

$$
\begin{align*}
& A_{e}=\Gamma\left(A-Q V_{2}^{-} C\right) G^{T}  \tag{2.7}\\
& B_{e}=\Gamma Q V_{2}^{-1}  \tag{2.8}\\
& C_{e}=L G^{r} \tag{2.9}
\end{align*}
$$

and such that the following conditions are satisfied.

$$
\begin{equation*}
0=A Q+Q A^{r}+V_{1}-Q V_{2}^{-1} Q^{r}+r_{1} Q V_{2}^{-} Q^{r_{T}^{\prime}} \tag{2.10}
\end{equation*}
$$

$$
\begin{align*}
& 0=A Q+Q A^{T}+Q V_{2}^{-1} Q^{T}-T_{1} Q V_{2}^{-1} Q^{T} r_{\perp}^{T}  \tag{2.11}\\
& 0=\left(A-Q V_{2}^{-1} C\right)^{T} P+P\left(A-Q V_{2}^{-} C\right)+L^{T} R L-r_{\perp}^{T} L^{T} R L T_{\perp}  \tag{2.12}\\
& \rho(Q)=\rho(F)=\rho(Q P)=n_{C} . \tag{2.13}
\end{align*}
$$

Remork 2.1: It is useful to note that (2.7) can be replaced by

$$
\begin{equation*}
A_{\mathrm{f}}=\Gamma A G^{T}-B_{\mathrm{f}} C G^{T} \tag{2.7}
\end{equation*}
$$

Remark 2.2: Because of (2.6) the $n \times n$ matrix $\tau$ which couples the three equations (2.10)-(2.12) is idempotent, i.e., $T^{2}=r$. In general, this "optimal projection" is an oblique projection (as opposed to an orthogonal projection) since it is not necessarily symmetric. Note that from Sylvester's inequality and (2.6) it follows that $\rho(\tau)=n_{e}$. It should be stressed that the form of the optimal reduced-order estimator (2.7)(2.9) is a direct consequence of optimality and not the result of an a priori assumption on the structure of the reduced-order estimator.
Remark 2.3: To obrain the standard steady-state Kalman filter result for the full-order case, set $p=n_{e}=n$ and $L=I_{n}$. Then $t=G=\Gamma=$ $I_{n}$ and thus (2.10) reduces to the standard observer Riccati equation [10, p. 367] and (2.7) and (2.8) yield the usual expressions. Furthermore, it follows from (2.7)' [11, Lemma 2.1] and standard results that (2.11)(2.13) are equivalent to the assumption that $\left(A_{e}, B_{e}, C_{f}\right)$ is controllable and observable.

Remark 2.4: Since $Q P$ is nonnegative semisimple it has a group generalized inverse ( $Q \rho$ ) ${ }^{\prime \prime}$ given by $G^{T} M^{-1} \Gamma$ (see, e.g., [12, p. 124]). Hence. by (2.6) the optimal projection $\tau$ is given by

$$
\begin{equation*}
t=Q P(\phi P)^{\prime} \tag{2.14}
\end{equation*}
$$

Remark 2.5: Replacing $x_{e}$ by $S x_{p}$, where $S$ is invertible, yields the "equivalent" estimator ( $S A_{1} S^{-1}, S B_{e}, C S^{-1}$ ). Since $J\left(A_{e}, B_{e}, C_{e}\right)=$ $J\left(S A_{c} S^{-1}, S B_{e}, C_{r} S^{-1}\right)$, one would expect the Main Theorem to apply also to $\left(S A S^{-1}, S B_{e}, C S^{-1}\right)$. This is indeed the case since transformation of the estimator state basis corresponds to the alternative factorization $Q^{\boldsymbol{F}}=\left(S^{-T} G\right)^{T}\left(S M S^{-i}\right)(S T)$.

Remark 2.6: Note that, for the optimal values of $A_{e}, B_{r}$, and $C_{e}$, (2.3) assumes the observer form

$$
\begin{equation*}
\dot{x}_{p}=\Gamma A G^{\top} x_{p}+\Gamma Q V_{i}^{-1}\left(y-C G^{{ }_{x}} x_{p}\right) \tag{2.15}
\end{equation*}
$$

By introducing the quasi-full-state estimate $\hat{x} \hat{\underline{x}} G^{T_{x_{p}}} \in \beta^{n}$ so that $\tau \hat{x}=$ $\hat{x}$ and $x_{e}=\Gamma \hat{x} \in \mathscr{S}^{n}{ }^{n}$, (2.15) can be written as

$$
\begin{equation*}
\dot{\hat{x}}=\tau A \tau \dot{x}+\tau Q V_{2}^{-1}(y-C \hat{x}) \tag{2.16}
\end{equation*}
$$

Note that although the implemented estimator (2.15) has the state $x_{e} \in$ $\mathfrak{B}^{n} e,(2.15)$ can be viewed as a quasi-full-order estimator whose geometric structure is entirely dictated by the projection $r$. Specifically, error inputs $Q V_{2}^{-1}(y-C \hat{x})$ are annihilated uniess they are contained in $[\mathscr{V}(r)]^{2}=$ $Q\left(\tau^{\eta}\right)$. Hence, the observation subspace of the estimator is precisely $Q\left(T T^{7}\right)$.

Remark 2.7: Although the form of (2.16) would lead one to surmise that the optimal reduced-order estimator is a projection of the optimal fuliorder estimator, this is not generally the case for the following simple reason. In the full-order case $Q$ (which appears in $Q$ ) is determined by solving a single Riccati equation, whereas in the reduced-order case $Q$ must be found in conjunction with $\boldsymbol{Q}$ and $\boldsymbol{\beta}$ to satisfy all three matrix equations (2.10)-(2.12). Hence, the value of $Q$ in the reduced-order case may be different from the value of $Q$ in the full-order case. Thus, (2.16) may not be obtainable by simply projecting the full-order result.

To further clarify the relationship between $\boldsymbol{Q}, \vec{P}$. and 9 , we now show that there exists similarity transformation which simultaneously diagonalizes $\boldsymbol{Q} \boldsymbol{P}$ and $\tau$.

Proposition 2.1: There exists invertible $\phi \in \int^{n \times \pi}$ such that

$$
\begin{aligned}
& Q=\phi^{-1}\left[\begin{array}{cc}
\Lambda_{Q} & 0 \\
0 & 0
\end{array}\right] \phi-T, \quad \beta_{=}=\phi^{T}\left[\begin{array}{cc}
\Lambda_{\beta} & 0 \\
0 & 0
\end{array}\right] \phi, \\
& Q \beta=\phi^{-1}\left[\begin{array}{ll}
\Lambda & 0 \\
0 & 0
\end{array}\right] \phi, \quad \tau=\phi \cdot\left[\begin{array}{cc}
A_{n_{p}} & 0 \\
0 & 0
\end{array}\right] \phi, \quad(2.18 \mathrm{a}, \mathrm{~b})
\end{aligned}
$$

where $\Lambda_{Q} \Lambda_{p} \in \sin _{\mathrm{n}} \times \mathrm{n}_{\mathrm{f}}$ are positive diagonal, $\Lambda \triangleq \Lambda_{Q} \Lambda_{\rho}$, and the diagonal elements of $\Lambda$ are the eigenvalues of $M$. Consequently,

$$
\begin{equation*}
Q=r Q . \quad P=\beta r_{r} . \tag{2.19}
\end{equation*}
$$

## III. Proof of the Man theorem

The proof proceeds exactly as in $[6]$. Using the fact that $Q_{+}$is open, the Frizz John version of the Lagrange multiplier theorem can be used to rigorously derive the first-order necessary conditions

$$
\begin{gather*}
0=A Q+Q A^{T}+D_{1}  \tag{3.1}\\
0=A^{T} P_{+}+P A+R_{1},  \tag{3.2}\\
0=P_{12}^{\top} Q_{12}+P_{2} Q_{2},  \tag{3.3}\\
B_{2}=-\left[\left(P_{2}^{-1} P_{12}^{T} Q_{1}+Q_{12}^{\top}\right) C^{T}+P_{2}^{-1} P_{12}^{T} V_{12} V_{2}^{-1},\right.  \tag{3.4}\\
C_{8}=L Q_{12} Q_{2}^{-1}, \tag{3.5}
\end{gather*}
$$

where

$$
\begin{gathered}
A=\left[\begin{array}{cc}
A & 0 \\
B_{c} C & A_{e}
\end{array}\right], \quad D=\left[\begin{array}{cc}
V_{1} & V_{12} B_{!}^{r} \\
B_{i} V_{12} & B_{i} V_{2} B_{!}^{T}
\end{array}\right], \\
R=\left[\begin{array}{cc}
L^{T} R L & -L^{T} R C_{i} \\
-C_{!}^{T} R L & C_{!}^{T} R_{2} C_{e}
\end{array}\right]
\end{gathered}
$$

and $\left(n+n_{f}\right) \times\left(n+n_{f}\right) \underline{Q} . \dot{P}$ are partitioned into $n \times n, n \times n_{e}$ and $n_{e}$ $\times n_{f}$ subblocks as

$$
Q=\left[\begin{array}{ll}
Q_{1} & Q_{12} \\
Q_{12}^{\tau} & Q_{2}
\end{array}\right], \quad P=\left[\begin{array}{cc}
P_{1} & P_{12} \\
P_{12}^{T} & P_{2}
\end{array}\right] .
$$

Expanding (3.1) and (3.2) yields

$$
\begin{align*}
& 0=A Q_{1}+Q_{1} A^{T}+V_{1} .  \tag{3.6}\\
& 0=A Q_{12}+Q_{12} A T+Q_{1}\left(B_{C} C\right)^{T}+V_{12} B!\text {. }  \tag{3.7}\\
& 0=A_{e} Q_{2}+Q_{2} A_{r}^{T}+B_{2} C Q_{12}+Q_{12}^{T}\left(B_{1} C^{r}+B_{e} V_{2} B_{6}^{T},\right.  \tag{3.8}\\
& 0=A^{T} P_{1}+P_{1} A+\left(B_{e} C\right)^{T} P_{12}^{T}+P_{12} B_{e} C+L^{T} R L \text {, }  \tag{3.9}\\
& 0=P_{12} A_{1}+A^{\top} P_{12}+\left(B_{1} C\right)^{T} P_{2}-L^{\top} R C_{e},  \tag{3.10}\\
& 0=A!P_{2}+P_{2} A_{1}+C_{!}^{T} R_{2} C_{6} . \tag{3.11}
\end{align*}
$$

Note that (3.9) is superfluous and can be omitted. Writing (3.8) as (see [13]. [14])

$$
0=\left(A_{e}+B_{f} C Q_{12} Q_{i}^{*}\right) Q_{2}+Q_{2}\left(A_{f}+B_{f} C Q_{12} Q_{2}^{+}\right)^{\gamma}+B_{f} V_{2} B_{f}^{r},
$$

where $Q_{i} \dot{i}$ is the Moore-Penrose or Drazin generalized inverse of $Q_{2}$, it follows from [11. Lemmas 2.1 and 12.2] that $Q_{2}$ is positive definite. Similarly. (3.11) implies that $P_{\mathbf{2}}$ is positive definite. This justifies (3.4) and (3.5).

Now define the $n \times n$ nonnegative-definite matrices (see [13]. [14])

$$
Q=Q_{1}-Q_{12} Q^{-1} Q_{12}^{\top}, \quad Q=Q_{12}^{\top} Q_{-}^{-1} Q_{12}^{\tau}, \quad \beta=P_{12} P_{2}^{-1} P_{12}^{\tau}
$$

and note that (3.3) implies (2.5) and (2.6) with

$$
G=Q_{i}^{-1} Q_{12}^{\top}, \quad M=Q_{2} P_{2}, \quad \Gamma=-P_{2}^{-1} P_{12}^{\top} .
$$

Since $Q_{2} P_{2}=P_{2}^{1 / 2}\left(P_{1}^{1 / 2} Q_{2} P l^{1 / 2}\right) P_{1^{\prime}} . M$ is positive semisimple. Sylvester's inequality yields (2.13). Nore (2.19) and the identities

$$
\begin{gather*}
Q_{1}=Q+Q .  \tag{3.12}\\
Q_{12}=Q r^{r} . P_{12}=-P G^{r} .  \tag{3.13}\\
Q_{2}=r Q r^{r} . P_{2}=G P G^{r} . \tag{3.14}
\end{gather*}
$$

Using (3.12)-(3.14), (3.4) and (3.5) yield (2.8) and (2.9). Also, the righ-hand sides of (3.8) and (3.7) yield (2.7). Substituting (2.7)-(2.9) imo (3.6)-(3.8), (3.10) and (3.11), it can be seen that (3.8) and (3.11) are also superfluous. Finally, linear combinations of the remaining three equations (3.6), (3.7), and (3.10) yield (2.10)-(2.12).

## IV. Concluding Remarks

The question of multiple local minima satisfying the optimal projection equations for reduced-order state estimation and the problem of constructing numerical methods for solving these equations are beyond the scope of this note. It should be pointed out, however, that promising numerical results for the model-reduction and fixed-order dynamic-compensation problems have been obtained by means of iterative algorithms that take full advantage of the presence and structure of the optimal projection [2]. [4]. [7].
Finally, the results of this paper can be extended to include the following related problems: 1) discrete-time system/discrete-time estimator: 2) infinite-dimensional system/finite-dimensional estimator [5]: and 3) parameter uncertainties [1], [15], [16].

## ACKNOWLEDGMENT

The authors wish to thank Dr. F. M. Ham for directing their attention to the reduced-order state-estimation problem as a fruitful application of the optimal projection approach.

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# The Optimal Projection Equations for Fixed-Order Dynamic Compensation 

## DAVID C. HYLAND and DENNIS S. BERNSTEIN

Abstract-First-order mecessary conditions for quadratically optimal, steady-state, fixed-order dynamic compensation of a linear, time-invariant plant in the presence of disturbance and observation moise are derived in a new and highly simplified form. In contrast to the pair of matrix Riccati equations for the full-order LQG case, the optimal steady-state fixed-order dynamic compeasator is characterized by four matrix equations (two modified Riccati equations and two modified Lyapunov equations) coupled by a projection whose rank is precisely equal to the order of the compensator and which determines the optimal compensator gaius. The coupling represents a graphic portrayal of the demise of the classical separation principle for the reduced-order controller case.

## I. INTRODUCTION

Because of constraints imposed by on-line computations, dynamic controllers for high-order systems such as flexible spacecrafi must be of relatively modest order. Hence, this paper is concerned with the design of quadratically optimal, fixed-order (i.e., reduced-order) dynamic compensation for a plant subject to stochastic disturbances and nonsingular measurement noise. Since white noise in all measurement channels precludes direct output feedback (see Section II), only purely dynamic controllers are considered. The requirements for resolution of this optimization problem include the following.

1) Conditions for the existence of an optimal, stabilizing compensator of the prescribed order. (In the full-order case these are the usual stabilizability and detectability conditions of LQG theory.)
2) Stationary conditions, i.e.. first-order necessary conditions, rendered in a tractable form to facilitate developments in items 3) and 4) below. (In the full-order case these conditions are precisely the LQG gain relations together with the regulator and observer Riccati equations.)
3) Sufficiency conditions. i.e., additional restrictions on solutions of the first-order necessary conditions which characterize local minima and single out the global minimum. (In the full-order case the global minimum is distinguished by the unique nonnegative-definite solutions to the LQG Riccati equations.)
4) Convergent numerical algorithms for simultaneous satisfaction of the necessary and sufficient conditions. In the full-order case numerical algorithms have been devised which take full advantage of the highly structured form of the Riccati equations.)
[^4]The present paper deals exclusively with item 2). Although the stationary conditions for the fixed-order compensation problem have been written down (see |||-|12|. for example). full expioitation has undoubtedly been impeded by their extreme complexity [see (3.3)-(3.11)]. What has been lacking. to quote the insightful remarks of $|9|$, "is a deeper understanding of the structural coherence of thesc equations." The contribution of the present paper is to show how the originally very complex stationary conditions can be transformed without loss of gencrality to much simpler and more tractable forms. The resulting equations (2.10)-(2.17) preserve the simple form of LQG relations for the gains in terms of covariance and cost matrices which, in turn, are determined by a coupled system of two modified Riccati equations and two modified Lyapunov equations. This coupling, by means of a projection (idempotent matrix) whose rank is pecisely equal to the order of the compensator, represents a graphic portrayal of the demise of the classical separation principle for the reduced-order controller case. When, as a special case, the order of the compensator is required to be equal to the order of the plant, the modified Riccati equations reduce to the standard LQG Riccati equations and the modified Lyapunov equations express the proviso that the compensator be minimal, i.e., controllable and observable. Since the LQG Riccati equations as such are nothing more than the necessary conditions for full-order compensation, we believe that the "optimal projection equations" provide a clear and simple generalization of standard LQG theory.

Since we are concerned with optimal fixed-order compensator design, our approach does not represent yet another model- or controllerreduction scheme along the lines of [13]-[17]. Indeed, the optimal projection equations, by virtue of their relatively transparent structure, can reveal the extent to which the design equations of a given ad hoc reduction scheme conform to the necessary conditions for optimality. For example, the oblique projection which arises in the present formulation may not be of the form $\left[\begin{array}{cc}l & 0 \\ 0 & 0\end{array}\right]$ even in the basis corresponding to the "balanced" realization [13]-\{16]. These issues are discussed in [18] where the results of [19] are simplified by means of the approach of the present paper and where the balancing method of [13] is reinterpreted in the context of optimality theory.

The fact that the optimal projection equations consist of four coupled matrix equations, i.e., two modified Riccati equations and two modified Lyapunov equations, should not be at all surprising for the following simple reason. Reduced-order control-design methods often involve either LQG applied to a reduced-order model or model reduction applied to a full-order LQG design. Both approaches, then, involve the solution of precisely four equations: two Riccati equations (for LQG) plus two Lyapunov equations (for model reduction via balancing, as in [13]). The coupled form of the optimal projection equations is thus a strong reminder that the LQG and order-reduction operations cannot be iterated but must, in a certain sense, be performed simultaneously.

## II. Problem Statement and the Maln Theorem

Given the control system

$$
\begin{gather*}
x(t)=A x(t)+B u(t)+w_{1}(t),  \tag{2.1}\\
y(t)=C x(t)+w_{2}(t) \tag{2.2}
\end{gather*}
$$

design a fixed-order dynamic compensator

$$
\begin{gather*}
\dot{x}_{c}(t)=A_{c} x_{r}(t)+B_{i} y(t),  \tag{2.3}\\
u(t)=C_{r} x_{r}(t) \tag{2.4}
\end{gather*}
$$

which minimizes the steady-state performance criterion

$$
\begin{equation*}
J\left(A_{c}, B_{c}, C_{c}\right) \triangleq \lim _{\rightarrow \rightarrow \infty} E\left(x(0)^{T} R_{1} x(t)+u(t)^{\tau} R_{2} u(t)\right] \tag{2.5}
\end{equation*}
$$

where: $x \in\left\|^{n}, u \in \pi^{m}, y \in\right\|^{\prime}, x_{r} \in \|^{n} c, n_{c} \leqslant n, A, B, C, A_{c}, B_{c}$, $C_{c}, R_{1}$, and $R_{2}$ are matrices of appropriate dimension with $R_{1}$ (symmetric) nonnegative definite and $R_{2}$ (symmetric) positive definite; $w_{1}$ is white disturbance noise with $n \times n$ nonnegative-definite intensity $V_{1}$ and $w_{2}$ is
white observation noise with $/ \times 1$ positive-definite intensity $V_{2} ; w_{1}$ and $w_{2}$ are uncorrelated and have zero mean. We note that the assumptions of nonsingular control weighting and nonsingular observation noise preclude the use of direct output feedback as in

$$
\begin{equation*}
u(t)=C_{r} x_{s}(t)+D_{c} y(t) \tag{2.6}
\end{equation*}
$$

since $J$ is undefined unless (see [7])

$$
\begin{equation*}
\operatorname{tr}\left[D_{c}^{I} R_{2} D_{c} V_{2}\right]=0 \quad\left(\Leftrightarrow R_{2} D_{1} V_{2}=0\right) \tag{2.7}
\end{equation*}
$$

To guarantee that $J$ is finite and independent of initial conditions we restrict our attention to the set of admissible stabilizing compensators

$$
Q \triangleq\left\{\left(A_{r}, B_{c}, C_{c}\right): A \triangleq\left[\begin{array}{cc}
A & B C_{r} \\
B_{c} C & A_{c}
\end{array}\right] \text { is asymptotically stable }\right\}
$$

where $\bar{A}$ is the closed-loop dynamics matrix. Since the value of $J$ is independent of the internal realization of the compensator, we can further restrict our attention to

$$
a, \triangleq\left\{\left(A_{c}, B_{c}, C_{c}\right) \in a\right.
$$

$$
\left.\left(A_{s}, B_{c}\right) \text { is controllable and }\left(C_{s}, A_{c}\right) \text { is observable }\right\} .
$$

For the following lemma call a square matrix nonnegative (respectively, positive) semisimple if it has a diagonal Jordan form and nonnegative (respectively, positive) eigenvalues. Let $I$, denote the $r \times r$ identity matrix.
Lemma 2.1: Suppose $Q, P \in \Omega^{n \times n}$ are nonnegative definite. Then $\phi \mathcal{P}$ is nonnegative semisimple. Furthermore, if $\operatorname{rank} \phi P=n_{c}$ then there exist $G, \Gamma \in \mathbb{R}^{n} c^{\times n}$ and positive-semisimple $M \in \mathbb{R}^{n^{*}{ }^{\times n} c}$ such that

$$
\begin{align*}
& Q P=G^{r} M \mathrm{~T}  \tag{2.8}\\
& \Gamma G^{r}=I_{n_{c}} \tag{2.9}
\end{align*}
$$

Proof. The result is an immediate consequence of [20, Theorem 6.2.5, p. 1231.

For convenience in stating the Main Theorem, define

$$
\Sigma \triangleq B R_{2}^{-1} B^{T}, \Sigma \triangleq C^{T} V_{2}^{-1} C
$$

and call $G, M$, and $\Gamma$ satisfying (2.8) and (2.9) a ( $G, M, \Gamma$ )-factorization of $Q P$.
Main Theorem: Suppose $\left(A_{c}, B_{c}, C_{c}\right) \in \mathcal{Q}_{+}$solves the steady-state fixed-order dynamic-compensation problem. Then there exist $n \times n$ nonnegative-definite matrices $Q, P, Q$, and $P^{\text {such that }} A_{c}, B_{c}$, and $C_{c}$ are given by

$$
\begin{align*}
& A_{c}=\Gamma(A-Q \Sigma-\Sigma P) G^{\top}  \tag{2.10}\\
& B_{c}=\Gamma Q C^{T} V_{2}^{-1}  \tag{2.11}\\
& C_{c}=-R_{2}^{-1} B^{T} P G^{T} \tag{2.12}
\end{align*}
$$

for some $(G, M, \Gamma)$-factorization of $Q P$. and such that with $\tau \underline{\underline{\Delta}} G T$ the following conditions are satisfied:

$$
\begin{align*}
0= & (A-\tau Q \Sigma) Q+Q(A-\tau Q S)^{r}+V_{1}+\tau Q \Sigma Q \tau^{T},  \tag{2.13}\\
0= & (A-\Sigma P \tau)^{T} P+P(A-\Sigma P \tau)+R_{1}+\tau^{T} P \Sigma P \tau,  \tag{2.14}\\
0= & \left.\tau(A-\Sigma P) Q+Q(A-\Sigma P)^{T}+Q \Sigma Q\right],  \tag{2.15}\\
0= & {\left[(A-Q \Sigma)^{r} P+P(A-Q \Sigma)+P \Sigma P\right] r, }  \tag{2.16}\\
& \operatorname{rank} Q=\operatorname{rank} P=\operatorname{rank} Q P=n_{c} . \tag{2.17}
\end{align*}
$$

Remark 2.1: Because of (2.9) the $n \times n$ matrix $\tau$ which couples the four equations (2.13)-(2.16) is idempotent, i.e., $r^{2}=\tau$. In general this "optimal projection" is an oblique projection (as opposed to an orthogonal projection) since it is not necessarily symmetric. Note that Sylvester's inequality and (2.9) imply that rank $t=n_{c}$.

Remark 2.2: Using the relations $Q=\tau Q$ and $\boldsymbol{P}=\boldsymbol{P}_{r}$ [see (3.12)],
the optimal projection equations (2.13)-(2.16) can be written in the equivalent form

$$
\begin{gather*}
0=A Q+Q A^{T}+V_{1}-Q \Sigma Q+\tau_{1} Q \Sigma Q Q_{1}^{T},  \tag{2.18}\\
0=A^{T} P+P A+R_{1}-P \Sigma P+\tau_{1}^{T} P \Sigma P_{\tau_{1}},  \tag{2.19}\\
0=(A-\Sigma P) Q+Q(A-\Sigma P)^{T}+Q \Sigma Q-\tau_{1} Q \Sigma Q \tau_{1}^{T},  \tag{2.20}\\
0=  \tag{2.21}\\
(A-Q \Sigma)^{T} P+P(A-Q \Sigma)+P \Sigma P-\tau_{1}^{T} P \Sigma P \tau_{1}
\end{gather*}
$$

where $r_{1} \triangleq I_{n}-T_{\text {. Note that in the full-order case } n_{c}=n, T=G=\Gamma}$ $=I_{n}$ and thus (2.18) and (2.19) reduce to the standard observer and regulator Riccati equations and (2.10)-(2.12) yield the usual LQG expressions. Furthermore, it can be shown that (2.20), (2.21), and (2.17) are equivalent to the assumption that $\left(A_{c}, B_{c}, C_{c}\right)$ is controllable and observable.

Remark 2.3: Since $\phi^{\prime} P$ is nonnegative semisimple it has a group generalized inverse ( $\varnothing$ $\boldsymbol{\rho}^{\prime \prime}$ given by $G^{T} M^{-1 \Gamma}$ (see e.g., [21, p. 124]). Hence, hy (2.9) the optimal projection $T$ is given by

$$
\begin{equation*}
\tau=\phi P(\phi P)^{\prime} \tag{2.22}
\end{equation*}
$$

Remark 2.4: The modified Riccati equations (2.13) and (2.14) are similar to the (single) "extended algebraic Riccati equation" which arises in the static output feedback problem (see, e.g., [22]).

Remark 2.5: Replacing $x_{c}$ by $S x_{c}$, where $S$ is invertible, yields the "equivalent"' compensator ( $S A_{C} S^{-1}, S B_{c}, C_{c} S^{-1}$ ). Since $J\left(A_{c}, B_{c}, C_{c}\right)$ $=J\left(S A_{S} S^{-1}, S B_{c}, C S^{-1}\right)$ one would expect the Main Theorem to apply also to $\left(S A_{C} S^{-1}, S B_{c}, C_{c} S^{-1}\right)$. This is indeed the case since transformation of the compensator state basis corresponds to the alternative factorization $Q P^{=}=\left(S^{-T} G\right)^{T}\left(S M S^{-1}\right)(S T)$. See [10] for related remarks.
Remark 2.6: By introducing the quasi-full-state estimate $f \triangleq G^{\boldsymbol{T}_{\boldsymbol{x}_{c}} \in}$ $\hat{\beta}^{n}$ so that $\tau \hat{x}=\hat{x}$ and $x_{c}=\Gamma \hat{x} \in \tilde{J}^{n} c,(2.1)-(2.4)$ can be written as

$$
\begin{gathered}
\dot{x}=A x+B C_{c} \tau \tilde{x}+w_{1}, \\
\hat{x}=\pi\left(A-\hat{B}_{r} C+B C_{c}\right) \tau \mathscr{X}+\tau \dot{B}_{c}\left(C x+w_{2}\right)
\end{gathered}
$$

where $B_{c} \Delta \underline{\underline{1}} Q C^{T} V_{2}^{-1}$ and $C_{c} \underline{\underline{\Delta}}-R_{2}^{-1} B^{T} P$. Although the implemented compensator has the state $x_{c} \in B^{n} c$, it can be viewed as a quasi-full-order compensator whose geometric structure is entirely dictated by the projection 7 . Sensor inputs $\mathcal{B}_{c} y$ are annihilated unless they are contained in $[\mathfrak{R}(\tau)]^{+}=\mathbb{R}\left(\tau^{\tau}\right)$, where $\mathfrak{N}$ and $R$ denote null space and range. Furthermore, the quasi-full-order state estimate $r \hat{X}$ employed in the control input is contained in $Q(r)$. Thus, $R(\tau)$ and $Q\left(\tau^{\tau}\right)$ are the control and observation subspaces of the compensator.

## IIL. Proof of the Main theorem

The proof given here considerably simplifies the original derivation given in [23] and [24]. Using the fact that $Q_{+}$is open, the Fritz John version of the Lagrange multiplier theorem can be used to rigorously derive the first-order necessary conditions ([7], see also [25])

$$
\begin{align*}
& 0=\bar{A} \bar{Q}+\bar{Q} \bar{A}^{T}+\bar{V} .  \tag{3.1}\\
& 0=\bar{A} \bar{T} \bar{f}+\bar{P} \bar{A}+\bar{R},  \tag{3.2}\\
& 0=P_{12}^{\tau} Q_{12}+P_{1} Q_{2},  \tag{3.3}\\
& B_{\mathrm{c}}=-\left(P_{2}^{-1} P_{12}^{\top} Q_{1}+Q_{12}^{\top}\right) C^{\top} V_{2}^{-1},  \tag{3.4}\\
& C_{i}=-R_{2}^{-1} B^{\top}\left(P_{1} Q_{12} Q_{2}^{-1}+P_{12}\right), \tag{3.5}
\end{align*}
$$

where

$$
\hat{V} \Delta\left[\begin{array}{cc}
V_{1} & 0 \\
0 & B_{i} V_{2} B_{r}^{T}
\end{array}\right], \quad R_{\Delta} \Delta\left[\begin{array}{cc}
R_{1} & 0 \\
0 & C_{!}^{\top} R_{2} C_{C}
\end{array}\right]
$$

and $\left(n+n_{c}\right) \times\left(n+n_{c}\right) Q$. $\mathcal{P}$ are pantitioned into $n \times n, n \times n_{c}$, and $n_{r} \times n_{r}$ subblocks as

$$
Q=\left[\begin{array}{ll}
Q_{1} & Q_{12} \\
Q_{1}^{r} & Q_{:}
\end{array}\right], \quad \beta=\left[\begin{array}{cc}
P_{1} & P_{12} \\
P_{1:}^{1} & P_{:}
\end{array}\right]
$$

Expanding (3.1) and (3.2) yields

$$
\begin{align*}
& 0=A Q_{1}+Q_{1} A^{T}+B C_{c} Q_{12}^{T}+Q_{12}\left(B C_{c}\right)^{r}+V_{1},  \tag{3.6}\\
& 0=A Q_{12}+Q_{12} A_{c}^{T}+B C_{c} Q_{2}+Q_{1}\left(B_{c} C\right)^{r}  \tag{3.7}\\
& 0=A_{c} Q_{2}+Q_{2} A T+B_{c} C Q_{12}+Q_{12}^{T}\left(B_{c} C\right)^{T}+B_{c} V_{2} B_{c}^{T},  \tag{3.8}\\
& 0=A^{T} P_{1}+P_{1} A+\left(B_{c} C\right)^{T} P_{12}^{T}+P_{12} B_{c} C+R_{1,}  \tag{3.9}\\
& 0=P_{12} A_{c}+A^{T} P_{12}+\left(B_{c} C\right)^{T} P_{2}+P_{1} B C_{c},  \tag{3.i0}\\
& 0=A_{c}^{T} P_{2}+P_{2} A_{c}+\left(B C_{c}\right)^{T} P_{12}+P_{12}^{T} B C_{c}+C_{c}^{T} R_{2} C_{r} \tag{3.11}
\end{align*}
$$

Writing (3.8) as (see [26], [27])

$$
0=\left(A_{c}+B_{\mathrm{r}} C Q_{12} Q_{2}^{*}\right) Q_{2}+Q_{2}\left(A_{\mathrm{c}}+B_{c} C Q_{12} Q_{:}^{*}\right)^{r}+B_{c} V_{2} B_{c}^{T}
$$

where $Q_{2}^{+}$is the Moore-Penrose or Drazin generalized inverse of $Q_{2}$, it follows from [28, Lemmas 2.1 and 12.2] that $Q_{2}$ is positive definite. Similarly, (3.11) implies that $P_{2}$ is positive definite. This justifies (3.4) and (3.5).

Now define the $n \times n$ ponnegative-definite matrices (see [26], [27])

$$
\begin{array}{ll}
Q \triangleq Q_{1}-Q_{12} Q_{2}^{-1} Q_{12}^{r}, & P \triangleq P_{1}-P_{12} P_{2}^{-1} P_{12}^{T}, \\
Q \triangleq Q_{1:} Q_{2}^{-} Q_{12}^{T}, & P \triangleq P_{12} P_{2}^{-1} P_{12}^{T}
\end{array}
$$

and note that (3.3) implies (2.8) and (2.9) with

$$
G \Delta Q_{2}^{-1} Q_{12}^{\tau}, M \triangleq Q_{2} P_{2}, \Gamma \Delta-P_{2}^{-1} P_{12}^{\tau}
$$

Since $Q_{2} P_{2}=P_{2}^{-1 / 2}\left(P_{2}^{1 / 2} Q_{2} P_{1}^{1 / 2}\right) P_{2}^{1 / 2}, M$ is positive semisimple. Sylvester's inequality yields (2.17). Note also that

$$
\begin{equation*}
Q=\tau Q, P=P_{7} . \tag{3.12}
\end{equation*}
$$

Next (2.11) and (2.12) follow from (3.4) and (3.5) by using the identities

$$
\begin{align*}
& Q_{1}=Q+Q, P_{1}=P+P  \tag{3.13}\\
& Q_{12}=Q \Gamma^{T}, P_{12}=-P G^{T}  \tag{3.14}\\
& Q_{2}=\Gamma Q \Gamma^{T}, P_{2}=G P G^{r} . \tag{3.15}
\end{align*}
$$

Now substitute (2.11), (2.12), and (3.13)-(3.15) into (3.6)-(3.11) and use the relations

$$
\begin{gathered}
B_{c} C=\Gamma Q \Sigma, B C_{c}=-\Sigma P G^{\top}, \\
B_{c} V_{2} B_{c}^{\top}=\Gamma Q \Sigma Q \Gamma^{\top}, C_{c}^{\top} R_{2} C_{c}=G P \Sigma P G^{\top} .
\end{gathered}
$$

Then (2.10) follows from (3.8)- $\Gamma(3.7)$. Substiruting (2.10) into (3.7), (3.8), (3.10), and (3.11) shows that ( 3.7$\left.)_{G}\right)^{r}$ and $-(3.10) \Gamma$ are precisely (2.15) and (216). Since $G^{T}(3.8) G=(2.15) r$ and $\Gamma^{T}(3.11) \Gamma=$ 7(2.16), (3.8) and (3.11) can be omitted. Finally, using (3.12) it follows that $(2.13)=(3.6)+(2.15) \tau-(2.15)-(2.15)^{\tau}$ and similarly for (2.14).

## IV. Directions for further Research

With regard to the existence of a stabilizing compensator, known results (e.g., [28]-[34]) can be exploited to a great extent. A numerical algorithm for solving the optimal projection equations has been developed in [24] and [35]. The proposed computational scheme is philosophically quite different from gradient search algorithms [2], [3], [6], (7). [9], [11]. [36], [37] in that it operates through dir ct solution of the optimal projection equations by iterative refinement of the optimal projection. Methods for eliminating local extrema are being investigated by applying component cost analysis [17]. Generalizations of the optimal projection equations can arise by considering the following extensions of the fixedorder dynamic-compensation protiem.

1) Discrete-Time System/Discrete-Time Compensafor: Digital implementation can be modeled by a discrete-time compensator with control of a continuous-time system facilitated by sampling and reconstruction devices.
2) Cross Weighsing/Correlated Disturbance and Observation Noise: This extension is straighforward and entircly analogous to the LQG cave (sec, e.g., (3. p. 351]).
3) Singular Observation Noise/Singular Control Weighting: With due attention to (2.7), direct oulput feedback can be used in the singular case. The nature of the problem forebodes all of the difficulties associated with the singular LQG problem. Note that the output feedback problem [22], [38], when viewed in this context, is highly singular.
4) Infinite-Dimensional Systems: The optimal projection equations have been extended in [39] and [40] to the case in which (2.1) is a distributed parameter system. for example, a parial or functional differential equation.
5) Decentralized Fixed-Order Controller: The optimal projection equations can be derived for the case in which the dynamic controlter has a fixed decentralized structure.
6) Parameter Uncertainties: The original derivation in [23] treated a Stratonovich state-dependent noise model representing parameter uncertainties in the plant. Further consideration of control- and measurementdependent noise raises the possibility of directly including the impact of parameter uncerainties in the design of robust, implementable compensation for large-order systems.

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# THE OPTIMAL PROJECTION EQUATIONS FOR FINITE-DIMENSIONAL FIXED-ORDER DYNAMIC COMPENSATION OF INFINITE-DIMENSIONAL SYSTEMS* 

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#### Abstract

Abatract. One of the major difficulties in designing implementable finite-dimensional controllers for distributed parameter systems is that such systems are inherently infinite dimensional while controller dimension is severely constrained by on-line computing capability. While some approaches to this probiem initially seek a correspondingly infinite-dimensional control law whose finite-dimensional approximation may be of impractically high order, the usual engineering approach involves first approximating the distributed parameter system with a high-order discretized model followed by design of a relatively low-order dynamic controller. Among the numerous approaches suggested for the latter step are model/controller reduction techniques used in conjunction with the standard LQG result. An alternative approach, developed in [36], relies upon the discovery in [31] that the necessary conditions for optimal fixed-order dynamic compensation can be transformed into a set of equations possessing remarkable structural coherence. The present paper generalizes this result to apply directly to the distributed parameter system itself. In contrast to the pair of operator Riccati equations for the "full-order" LQG case, the optimal finite-dimensional fixed-order dynamic compensator is characterized by four operator equations (two modified Riccati equations and two modified Lyapunov equations) coupled by an oblique projection whose rank is precisely equal to the order of the compensator and which determines the optimal compensator gains. This "optimal projection" is obtained by a full-rank factorization of the product of the finite-rank nonnegative-definite Hilber-space operators which satisfy the pair of modified Lyapunov equations. The coupling represents a graphic portrayal of the demise of the classical separation principle for the finite-dimensional reduced-order controller case. The results obtained apply to a semigroup formulation in Hilber space and thus are applicable to control problems involving a broad range of specific partial and functional differential equations.


Key words. optimality conditions, finite-dimensional fixed-order dynamic compensator, infinitedimensional system, distributed parameter system, semisimple operator, oblique projection, Drazin generalized inverse

1. Introduction. One of the major difficulties in designing active controllers for distributed parameter systems is that such systems are inherently infinite dimensional while implementable controllers are necessarily finite dimensional with controller dimension severely constrained by on-line computing capability. As pointed out by Balas ([1], see also [2]), control design for distributed parameter systems entails the practical constraints of 1) finitely many sensors and actuators, 2) a finite-dimensional controller and 3 ) natural system dissipation. The validity of 2 ) is apparent from the fact that processing and transmitting electrical signals by conventional analog or digital components constitutes finite-dimensional action. Although distributed parameter devices can also be utilized, their fabrication and implementation can incorporate at most a finite number of design specifications.' Hence, although distributed parameter systems are most accurately represented by infinite-dimensional models, real-world

[^5]constraints require that implementable controllers be modelled as lumped parameter systems.

Clearly, the above observations effectively preclude the possibility of realizing infinite-dimensional controllers that involve full-state feedback or full-state estimation (see, e.g., [4]-[6] and the numerous references therein). Although finite-dimensional approximation schemes have been applied to optimal infinite-dimensional control laws ([7]-[9]), these results only guarantee optimality in the limit, i.e., as the order of the approximating controller increases without bound. Hence, there is no guarantee that a particular approximate (i.e., discretized) controller is actually optimal over the class of approximate controllers of a given order dictated by implementation constraints. Moreover, even if an optimal approximate finite-dimensional controller could be obtained, it would almost certainly be suboptimal in the class of all controllers of the given order.

Although the usual engineering approach to this problem is to replace the distributed parameter system with a high-order finite-dimensional model, analogous, fundamental difficulties remain since application of $L Q G$ leads to a controller whose order is identical to that of the high-order approximate model. Attempts to remedy this problem usually rely upon some method of open-loop model reduction or closedloop controller reduction (see, e.g., [10]-[15]). Most of these techniques (with the exception of [11]) are ad hoc in nature, however, and hence guarantees of optimality and stability may be lacking.

A more direct approach that avoids both model and controller reduction is to fix the controller structure and optimize the performance criterion with respect to the controller parameters. Although much effort was devoted to this approach (see, e.g., [16]-[30]), progress in this direction was impeded by the extreme complexity of the nonlinear matrix equations arising from the first-order necessary conditions. What was lacking, to quote the insightful remarks of [24], was a "deeper understanding of the structural coherence of these equations." The key to unlocking these unwieldy equations was subsequently discovered by Hyland in [31] and developed in [32]-[36]. Specifically, it was found that these equations harbored the definition of an oblique projection (i.e., idempotent matrix) which is a consequence of optimality and not the result of an ad hoc assumption. By exploiting the presence of this "optimal projection," the originally very complex stationary conditions can be transformed without loss of generality into much simpler and more tractable forms. The resulting equations (see [36, (2.10)-(2.17)]) preserve the simple form of $L Q G$ relations for the gains in terms of covariance and cost matrices which, in turn, are determined by a coupled system of two modified Riccati equations and two modified Lyapunov equations. This coupling, by means of the optimal projection, represents a graphic portrayal of the demise of the classical separation principle for the reduced-order controller case. When, as a special case, the order of the compensator is required to be equal to the order of the plant, the modified Riccati equations immediately reduce to the standard LQG Riccati equations and the modified Lyapunov equations express the proviso that the compensator be minimal, i.e., controllable and observable. Since the LQG Riccati equations as such are nothing more than the necessary conditions for full-order compensation, the "optimal projection equations" appear to provide a clear and simple generalization of standard LQG theory.

The fact that the optimal projection equations consist of four coupled matrix equations, i.e., two modified Riccati equations and two modified Lyapunov equations, can readily be explained by the following simple reason. Reduced-order control-design methods often involve either $L Q G$ applied to a reduced-order model or model reduction
applied to a full-order LQG design, and hence both approaches require the solution of precisely four equations: two Riccati equations (for LQG) plus two Lyapunov equations (for system reduction via balancing, as in [12], [14]). The coupled form of the optimal projection equations is thus a strong reminder that the $L Q G$ and orderreduction operations cannot be iterated but must, in a precise sense, be performed simulraneously. This situation is partly due to the fact that the optimal projection matrix may not be of the form $\left[\begin{array}{ll}1 & 0 \\ 0 & 0\end{array}\right]$ even in the basis corresponding to the "balanced" realization [12], [14]. This point is explored in [37], [37a] where the solution to the optimal model-reduction problem is characterized by a pair of modified Lyapunov equations which are also coupled by an oblique projection.

Returning now to the distributed parameter problem, it should be mentioned that notable exceptions to the previously mentioned work on distributed parameter controllers are the contributions of Johnson [38] and Pearson [39], [40] who suggest fixing the order of the finite-dimensional compensator while retaining the distributed parameter model. Progress in this direction, however, was impeded not only by the intractability of the optimality conditions that were available for the finite-dimensional problem (as in [16]-[30]), but also by the lack of a suitable generalization of these conditions to the infinite-dimensional case. The purpose of the present paper is to make significant progress in filling these gaps, i.e., by deriving explicit optimality conditions which directly characterize the optimal finite-dimensional fixed-order dynamic compensator for an infinite-dimensional system and which are exactly analogous to the highly simplified optimal projection equations obtained in [31]-[34], [36] for the finite-dimensional case. Specifically, instead of a system for four matrix equations we obtain a system of four operator equations whose solutions characterize the optimal finite-dimensional fixed-order dynamic compensator. Moreover, the optimal projection now becomes a bounded idempotent Hilbert-space operator whose rank is precisely equal to the order of the compensator.

The mathematical setting we use is standard: a linear time-invariant differential system in Hilbert space with additive white noise, finitely many controls and finitely many noisy measurements (thus satisfying the first practical constraint mentioned above). The input and output maps are assumed to be bounded. Since the only explicit assumption on the unbounded dynamics operator is that it generate a strongly continuous semigroup, the results are potentially applicable to a broad range of specific partial and functional differential equations. The actual applicability of our results is essentially limited by practical constraint 3 ). Since we are concerned with the steadystate problem, we implicitly assume that the distributed parameter system is stabilizable, i.e., that there exists a dynamic compensator of a given order such that the closed-loop system is uniformly stable. We note that stabilizing compensators do exist for the wide class of problems considered in [41] and [42] which includes delay, parabolic and damped hyperbolic systems. The question of how much damping is required for stabilizability of hyperbolic systems is a crucial issue in designing controllers for large flexible space structures [7], [43]-[49a].

It is important to point out that the results of this paper can immediately be specialized to finite-dimensional systems by requiring that the Hilbert space characterizing the dynamical system be finite-dimensional. Then all unboundedness considerations can be ignored, adjoints can be interpreted as transposes and other obvious simplifications can be invoked. The only mathematical aspect requiring attention is the treatment of white noise which, for general handling of the infinite-dimensional case, is interpreted according to [6]. ${ }^{2}$ For the finite-dimensional case, however, the standard
classical notions suffice and the results go through with virtually no modifications.
The contents of the paper are as follows. Section 2 contains preliminary notation in addition to particular results for use later in the paper. Section 3 presents the optimal steady-state finite-dimensional fixed-order dynamic-compensation problem and the Main Theorem gives the necessary conditions in the form of the optimal projection equations (3.15)-(3.18). We then develop a series of results which serve to elucidate several aspects of the Main Theorem. Section 4 is devoted to the proof of the Main Theorem. The reader is alerted to the two crucial steps required. The first step involves generalizing to the infinite-dimensional case the derivation of the necessary conditions in their "primitive" form (see (4.27)-(4.29) and (4.48)-(4.53)). The derivation in [31]-[33], [36] involving Lagrange multipliers is invalid in the infinite-dimensional case due to the presence of the unbounded system-dynamics operator. Instead, we use the gramian form of the closed-loop covariance operator to obtain a dual problem formulation and then proceed to derive the primitive necessary conditions by means of a lengthy, but direct, computation (Lemma 4.7). The second crucial step involves transforming the primitive form of the necessary conditions to the final form given in the Main Theorem. This laborious computation was first carried out in [31], [32] and was subsequently facilitated in [33], [36] by means of a judicious change of variables (see (4.32), (4.33)). Finally, some concluding remarks are given in § 5 .
2. Preliminaries. In this section we introduce general notation along with basic definitions and results for use in later sections. Our principal references are [6], [50] and [51].

Throughout this section let $\mathscr{K}, \mathscr{K}^{\prime}$ and $\mathscr{K}^{\prime \prime}$ denote real separable Hilbert spaces with norm $\|\cdot\|$ and inner product $\langle\cdot, \cdot\rangle$ and let $\mathscr{B}\left(\mathscr{X}, \mathscr{X}^{\prime \prime}\right)$ denote the space of bounded linear operators from $\mathscr{H}$ into $\mathscr{H}^{\prime}$. For $L \in \mathscr{B}(\mathscr{H}, \mathscr{H}),\|L\|$ is the norm of $L, \mathscr{R}(L)$ is the range of $L, \mathcal{N}(L)$ is the null space of $L, \rho(L)$ is the rank of $L$ (set $\rho(L)=\infty$ if $L$ does not have finite rank), $L^{-1}$ is the inverse of $L$ when $L$ is invertible, i.e., when $L$ has a bounded inverse, $L^{*}$ is the adjoint of $L$ and $L^{-*} A\left(L^{*}\right)^{-1}$. Recall that $\|L\|=\left\|L^{*}\right\|$ and that $\rho(L)=\rho\left(L^{*}\right)\left[50\right.$, p. 161]. Now suppose that $\mathscr{K}=\mathscr{K}^{\prime}$ so that $L \in \mathscr{B}(\mathscr{H}) \wedge \mathscr{B}(\mathscr{X}, \mathscr{H})$. If $L L^{*}=L^{*} L$ then $L$ is normal and if $L=L^{*}$ then $L$ is selfadjoint. If $L$ is selfadjoint and $\langle L x, x\rangle \geqq 0, x \in \mathscr{H}$, then $L$ is now regative definite. Note that the selfadjointness assumption is included in the definitio since the Hilbert spaces are assumed real. If $L$ is nonnegative definite then $L^{1 / 2}$ denotes the (unique) nonnegative-definite square root of $L$. Call $L$ semisimple (resp., real semisimple, nonnegative semisimple) if there exists invertible $S \in \mathscr{B}(\mathscr{H})$ such that $S L S^{-1}$ is normal (resp., selfadjoint, nonnegative definite). This implies that $S L S^{-1}$ has a complete set of orthonormal eigenvectors and, in the real-semisimple or nonnegative-semisimple cases, has real or nonnegative eigenvalues.

Recall that if $L \in \mathscr{B}(\mathscr{X})$ is compact then $L$ has at most a countable number of eigenvalues and all nonzero eigenvalues have finite multiplicity. Hence, for $L \in$ $\mathscr{B}(\mathscr{X}, \mathscr{H})$ compact, let $\left\{\alpha_{i}\right\}$ be the (at most countable) sequence of eigenvalues of $\left(L L^{*}\right)^{1 / 2}$ with appropriate multiplicity and $\alpha_{1} \geqq \alpha_{2} \geqq \cdots>0$ [50, p. 261]. Then $\mathscr{B}_{1}\left(\mathscr{K}, \mathscr{K}^{\prime \prime}\right)$ denotes the set of race class (or nuclear) operators, i.e., the set of compact

[^6]$L \in \mathscr{B}\left(\mathscr{X}, \mathscr{X}^{\prime}\right)$ for which $\sum_{i} \alpha_{i}<\infty[50$, p. 521$] . \mathscr{B}_{1}\left(\mathscr{X}, \mathscr{X}^{\prime}\right)$ is a Banach space with norm
$$
\|L\|_{i} \triangleq \sum_{i} a_{i}
$$

If $\Sigma_{i} \alpha_{i}^{2}<\infty$ then $L \in \mathscr{B}_{2}\left(\mathscr{X}, \mathscr{X}^{\prime}\right)$, the set of Hilbert-Schmidt operators, which is a Banach space with norm

$$
\|L\|_{2} \underline{\underline{\Delta}}\left[\sum_{i} \alpha_{i}^{2}\right]^{1 / 2}
$$

Note that $\|L\| \leq\|L\|_{2} x\|L\|_{1},\|L\|=\left\|L^{*}\right\|,\|L\|_{1}=\left\|L^{*}\right\|_{1}$ and $\|L\|_{2}=\left\|L^{*}\right\|_{2}$. If $\mathscr{K}=\mathscr{H}^{\prime}$, then we write $\mathscr{B}_{1}(\mathscr{K})$ and $\mathscr{B}_{2}(\mathscr{X})$ for $\mathscr{B}_{1}(\mathscr{H}, \mathscr{X})$ and $\mathscr{B}_{2}(\mathscr{H}, \mathscr{H})$, respectively. Note that if nonnegative-definite $L \in \mathscr{B}_{1}(\mathscr{X})$ then $L^{1 / 2} \in \mathscr{B}_{2}(\mathscr{H})$.

If $L \in \mathscr{B}_{1}\left(\mathscr{H}, \mathscr{H}^{\prime}\right)$ and $S \in \mathscr{B}\left(\mathscr{H}^{\prime}, \mathscr{H}^{\prime \prime}\right)$ then

$$
\|S L\|_{1} \leqq\|S\|\|L\|_{1}
$$

and hence $S L \in \mathscr{B}_{1}\left(\mathscr{H}, \mathscr{K}^{\prime}\right)$. Similarly, under suitable hypotheses,

$$
\|L S\|_{1} \equiv\|S\|\|L\|_{1}
$$

and

$$
\|S L\|_{1} \leqq\|S\|_{2}\|L\|_{2} .
$$

Lemma 2.1. Suppose $L \in \mathscr{F}_{1}(\mathscr{H})$ and let $\left\{\lambda_{i}\right\}$ denote the nonzero eigenvalues of $L$ with appropriate multiplicity. Then [51, p. 89]

$$
\sum_{i}\left|\lambda_{i}\right| \leqq\|L\|_{1} .
$$

If $L$ is selfadjoint then [50, p. 522]

$$
\sum_{i}\left|\lambda_{i}\right|=\|L\|_{1} .
$$

If $L$ is nonnegative definite then

$$
\sum_{i} \lambda_{i}=\|L\|_{i} .
$$

Let $L \in \mathscr{B}_{1}(\mathscr{H})$. Then define $\left[50, p\right.$. 523] the trace functional tr: $\mathscr{B}_{1}(\mathscr{H}) \rightarrow \mathbb{R}$ by

$$
\operatorname{tr} L \triangleq \sum_{i}\left\langle L \phi_{i}, \phi_{i}\right\rangle
$$

where the summation is independent of the choice of orthonormal basis $\left\{\phi_{i}\right\}$. The trace satisfies $\operatorname{tr} L=\operatorname{tr} L^{*}, \operatorname{tr} S L=\operatorname{tr} L S$ for all $S \in \mathscr{B}(\mathscr{H}), \operatorname{tr} S T=\operatorname{tr} T S$ for all $S, T \in \mathscr{B}_{2}(\mathscr{H})$ and $\operatorname{tr}(\alpha T+\beta S)=\alpha(\operatorname{tr} T)+\beta(\operatorname{tr} S)$ for all $\alpha, \beta \in R$ and $S, T \in \mathscr{B}_{1}(\mathscr{H})$.

Lemma 2.2. Suppose $L \in \mathscr{B}_{1}(\mathscr{H})$ and let $\left\{\lambda_{i}\right\}$ denote the nonzero eigenvalues of $L$ with appropriate multiplicity. Then [51, p. 139]

$$
\operatorname{tr} L=\sum_{1} \lambda_{i}
$$

and hence (by Lemma 2.1)

$$
|\operatorname{tr} L| \leqq\|L\|_{1} .
$$

If $L$ is nonnegative definite then

$$
\operatorname{tr} L=\|L\|_{1}
$$

Corollary 2.1. For each $S \in \mathscr{F}(\mathscr{H})$ the linear functionals

$$
\begin{aligned}
& L \rightarrow \operatorname{tr} S L: \mathscr{B}_{1}(\mathscr{H}) \rightarrow R, \\
& L \rightarrow \operatorname{tr} L S: \mathscr{B}_{1}(\mathscr{H}) \rightarrow R
\end{aligned}
$$

are continuous. For each $L \in \mathscr{B}_{1}(\mathscr{H})$ the linear funcrionals

$$
\begin{aligned}
& S \rightarrow \operatorname{tr} L S: \mathscr{B}(\mathscr{X}) \rightarrow R, \\
& S \rightarrow \operatorname{tr} S L: \mathscr{B}(\mathscr{H}) \rightarrow R
\end{aligned}
$$

are continuous.
Although showing that a bounded linear operator is trace class is slightly more involved than the above characterizations of $\mathscr{B}_{1}(\mathscr{K})$, the following result will suffice for our purposes (see [52, p. 96], or [52a, p. 171]).

Lemma 2.3. Let $L \in \mathscr{B}(\mathscr{H})$ be nonnegative definite. Then

$$
\sum_{i}\left\langle L \phi_{i} \phi_{i}\right\rangle,
$$

whether finite or infinite, is independent of the orthonormal basis $\left\{\phi_{i}\right\}$. The summation is finite if and only if $L \in \mathscr{B}_{1}(\mathscr{H})$.

Many of the operators introduced in the following section have finite-dimensional domain or range space and hence are degenerate, i.e., have finite rank. Recall that degenerate operators are necessarily trace class. The following result, which generalizes [53, Thm. 2.1, p. 240] in certain respects, will be fundamental in decomposing finite-rank operators.

Lemma 2.4. Suppose $L_{1}, \cdots, L_{r} \in B\left(\mathscr{H}, \mathscr{H}^{\prime}\right)$ have finite rank. Then there exists a finite-dimensional subspace $\mathcal{M} \subset \mathscr{H}$ such that $L_{i} \mathcal{M}^{\perp}=0, i=1, \cdots, r$. Furthermore, if $\mathscr{H}=\mathscr{H}^{\prime}$ then $\mathscr{M}$ can be chosen such that $L_{i} \mathcal{M} \subset \mathcal{M}, i=1, \cdots, r$.

Proof. It suffices to consider the case $r=1$. Writing $L$ for $L_{1}$, note that since $\rho\left(L^{*}\right)<\infty, \mathcal{N}(L)^{\perp}=\mathscr{R}\left(L^{*}\right)[50$, p. 155] and $\mathcal{N}(L)$ is closed, the first statement holds with $\mathscr{M}=\mathcal{N}(L)^{\perp}$. When $\mathscr{H}=\mathscr{K}^{\prime}$ set $\mathscr{H}=\mathcal{N}(L)^{\perp}+\mathscr{R}(L)$ and note that $\mathcal{M}^{\perp}=$ $\mathcal{N}(L) \cap \mathscr{R}(L)^{\perp} \subset \mathcal{N}(L)$ and $L \mathcal{H} \subset \mathscr{R}(L) \subset \mu$.

The following generalization of Sylvester's inequality [54, p. 66] will be used repeatedly in handling finite-rank operators.

Lemma 2.5. Let $L \in \mathscr{B}\left(\mathscr{K}, \mathscr{H}^{\prime}\right)$ and $S \in \mathscr{B}\left(\mathscr{H}^{\prime}, \mathscr{H}^{\prime \prime \prime}\right)$. Then

$$
\begin{equation*}
\rho(S L) \leqq \min \{\rho(S), \rho(L)\} \tag{2.1}
\end{equation*}
$$

If $\operatorname{dim} \mathscr{X}^{\prime}=\nu<\infty$, then

$$
\begin{equation*}
\rho(S)+\rho(L)-\nu \leqq \rho(S L) \tag{2.2}
\end{equation*}
$$

Proof. If either $S$ or $L$ does not have finite rank then (2.1) is immediate. If both $S$ and $L$ have finite rank then the standard arguments [54] used to prove the finitedimensional version of (2.1) remain valid. To prove (2.2), note that Lemma 2.4 implies that there exist orthonormal bases for $\mathscr{H}$ and $\mathscr{K}^{\prime}$ with respect to which $L$ has the matrix representation $\left[\tilde{L} 0\right.$ ], where $\tilde{L} \in \mathbf{R}^{\wedge \times p}$. Similarly, there exist orthonormal bases for $\mathscr{H} \mathscr{K}^{\prime}$
 Since the two cited bases for $\mathscr{K}^{\mu}$ may be different, let orthogonal $U \in R^{\nu \times \nu}$ be the matrix representation (with respect to either basis for $\mathscr{H}^{\prime}$ ) for the change in orthonormal basis [6, p. 100]. Hence $S L$ has the matrix representation [ $\left.\begin{array}{cc}s U L & 0 \\ 0 & 0\end{array}\right]$ and (2.2) follows from the known result [54, p. 66].

As in the proof of Lemma 2.5, we shall utilize the infinite-matrix representation of an operator with respect to an orthonormal basis. All matrix representations given
here will consist of real entries since the Hilbert spaces involved are real. When the orthonormal bases are specified and no confusion can arise, we shall not differentiate between an operator and its matrix representation. We shall use the infinite identity matrix $I_{\infty}$ interchangeably with the identity $I_{x}$ on $\mathscr{K}$.

When dealing with finite-dimensional Euclidean spaces the notation and terminology introduced above will be utilized with only minor changes. For example, bounded linear operators will be represented by matrices whose elements are determined according to fixed orthonormal bases and hence we identify $\mathbf{R}^{m \times n}=\mathscr{B}\left(\mathbf{R}^{n}, \mathbf{R}^{m}\right)$. Note that if $L \in \mathscr{B}\left(\mathbf{R}^{n}, \mathscr{H}\right)$ and $S \in \mathscr{B}\left(\mathscr{H}, \mathbf{R}^{m}\right)$ then $S L$ is an $m \times n$ matrix which is independent of any particular orthonormal basis for $\mathscr{H}$. The transposes of $x \in \mathbf{R}^{n} \neq \mathbf{R}^{n \times 1}$ and $M \in R^{m \times n}$ are denoted by $x^{T}$ and $M^{T}$ and $M^{-T} \&\left(M^{\top}\right)^{-1}$. Let $I_{n}$ denote the $n \times n$ identity matrix.

To specialize some of the above operator terminology to matrices, let $M \in R^{n \times n}$. We shall say $M$ is nonnegative (resp., positive) diagonal if $M$ is diagonal with nonnegative (resp., positive) diagonal elements. $M$ is nonnegative (resp., positive) definite if $M$ is symmetric and $x^{T} M x \geqq 0$ (resp., $x^{T} M x>0$ ), $x \in R^{n}$. Recall that $M$ is symmetric (resp., nonnegative definite, positive definite) if and only if there exists orthogonal $U \in \mathbf{R}^{n \times n}$ such that $U M U^{T}$ is diagonal (resp., nonnegative diagonal, positive diagonal). $M$ is semisimple [55, p. 13], or nondefective [56, p. 375], if $M$ has $n$ linearly independent eigenvectors, i.e., $M$ has a diagonal Jordan canonical form over the complex field. $M$ is real (resp., nonnegative, positive) semisimple if $M$ is semisimple with real (resp., nonnegative, positive) eigenvalues. Note that $M$ is real (resp., nonnegative, positive) semisimple if and only if there exists invertible $S \in \mathbf{R}^{n \times n}$ such that $S M S^{-1}$ is diagonal (resp., nonnegative diagonal, positive diagonal). Alternatively, $M$ is real (resp., nonnegative, positive) semisimple if and only if there exists invertible $S \in \mathbf{R}^{n \times n}$ such that $S M S^{-1}$ is symmetric (resp., nonnegative definite, positive definite).

Lemma 2.6. The product of two nonnegative- (resp., positive-) definite matrices is nonnegative (resp., positive) semisimple.

Proof. If $S, L \in \mathbb{R}^{n \times n}$ are both nonnegative (resp., positive) definite then by [ 55 ,
 $D_{L} A^{A} \phi^{T} L \phi$ are nonnegative (resp., positive) diagonal. Hence, $S L=\phi D_{S} D_{L} \phi^{-1}$ is nonnegative (resp., positive) semisimple, as desired. Alternatively, if either $S$ or $L$ is positive definite, then the result follows from $S L=L^{-1 / 2}\left(L^{1 / 2} S L^{1 / 2}\right) L^{1 / 2}$ if $L$ is positive definite or $S L=S^{1 / 2}\left(S^{1 / 2} L S^{1 / 2}\right) S^{-1 / 2}$ if $S$ is positive definite.
3. Problem statement and the Main Theorem. We consider the following steadystate fixed-order dynamic-compensation problem. Given the dynamical system on $[0, \infty)$

$$
\begin{align*}
& \dot{x}(t)=A x(t)+B u(t)+H_{1} w(t),  \tag{3.1}\\
& y(t)=C x(t)+H_{2} w(t) \tag{3.2}
\end{align*}
$$

design a finite-dimensional fixed-order dynamic compensator

$$
\begin{align*}
& \dot{x}_{\tau}(t)=A_{c} x_{c}(t)+B_{c} y(t),  \tag{3.3}\\
& u(t)=C_{c} x_{c}(t) \tag{3.4}
\end{align*}
$$

which minimizes the steady-state performance criterion

$$
\begin{equation*}
J\left(A_{c} B_{c} C_{c}\right) \triangleq \lim _{i \rightarrow \infty} E\left[\left\langle R_{1} x(t), x(t)\right\rangle+u(t)^{T} R_{2} u(t)\right] \tag{3.5}
\end{equation*}
$$

The following data are assumed. The state $x(t)$ is an eiement of a real separable Hilbert space $\mathscr{X}$ and the state differential equation is interpreted in the weak sense (see, e.g., [6, pp. 229, 317]). The closed, densely defined operator $A: \mathscr{D}(A) \subset \mathscr{K} \rightarrow \mathscr{H}$ generates a strongly continuous semigroup $e^{A t}, t \geq 0$. The control $u(t) \in \mathbf{R}^{m}, B \in$ $\mathscr{B}\left(\mathbf{R}^{m}, \mathscr{X}\right)$ and the operator $R_{1} \in \mathscr{B}_{1}(\mathscr{H})$ and the matrix $R_{2} \in \mathbf{R}^{m \times m}$ are nonnegative definite and positive definite, respectively. $w(\cdot)$ is a zero-mean Gaussian "standard white noise process" in $L_{2}\left((0, \infty), \mathscr{H}^{\prime}\right)$ (see [6, p. 314]), where $\mathscr{H}^{\prime}$ is a real separable Hilbert space, $H_{1} \in \mathscr{B}_{2}\left(\mathscr{K}^{\prime}, \mathscr{H}\right), H_{2} \in \mathscr{B}\left(\mathscr{K}^{\prime}, \mathbf{R}^{\prime}\right)$ and "E" denotes expectation. We assume that $H_{1} H_{2}^{*}=0$, i.e., the disturbance and measurement noises are independent, ${ }^{3}$ and that $V_{2} \underline{L}_{H_{2}} H_{2}^{*} \in R^{\prime}$ is positive definite, i.e., all measurements are noisy. Note that $V_{1} \not H_{1} H_{1}^{*} \in \mathscr{F}_{1}(\mathscr{H})$ is nonnegative definite and trace class. ${ }^{4}$ The initial state $x(0)$ is Gaussian and independent of $w(\cdot)$. The observation $y(t) \in \mathbf{R}^{\prime}$ and $C \in \mathscr{B}\left(\mathscr{H}, \mathbf{R}^{\prime}\right)$. The dimension of the compensator state $x_{c}(t)$ is of fixed, finite order $n_{c} \leqq \operatorname{dim} \mathscr{H}$ and the optimization is performed over $A_{c} \in R^{n_{c} \times n_{c}}, B_{c} \in R^{n_{c} \times 1}$ and $C_{c} \in R^{m \times n_{c}}$.

To handle the closed-loop system (3.1)-(3.4), we introduce the augmented state space $\mathscr{H} \triangleq \mathscr{H} \oplus \mathbf{R}^{n} c$ which is a real separable Hilbert space with inner product $\left\langle\bar{x}_{1}, \bar{x}_{2}\right\rangle \geqslant$ $\left\langle x_{1}, x_{2}\right\rangle+x_{c 1}^{T} x_{c 2}, \dot{x}_{i} \xlongequal{\mu}\left(x_{i}, x_{c i}\right)$. An operator $L \in \mathscr{B}(\dot{\mathscr{H}})$ has a "decomposition" into operators $L_{1} \in \mathscr{B}(\mathscr{H}), L_{12} \in \mathscr{F}\left(\mathbb{R}^{n_{c}}, \mathscr{H}\right), L_{21} \in \mathscr{B}\left(\mathscr{H}, \mathbf{R}^{n_{r}}\right)$ and $L_{2} \in \mathbf{R}^{n_{c}{ }^{n_{c}}}$ in the sense that for $\dot{x} \neq\left(x, x_{c}\right) \in \dot{\mathscr{H}}, L \tilde{x}=\left(L_{1} x+L_{12} x_{c}, L_{21} x+L_{2} x_{c}\right)$, or, in "block" form,

$$
L=\left[\begin{array}{ll}
L_{1} & L_{12} \\
L_{21} & L_{2}
\end{array}\right]
$$

For later use note that

$$
\|L\| \leqq\left\|L_{1}\right\|+\left\|L_{12}\right\|+\left\|L_{21}\right\|+\left\|L_{22}\right\|
$$

and

$$
L^{*}=\left[\begin{array}{ll}
L_{1}^{*} & L_{21}^{*} \\
L_{12}^{*} & L_{2}^{T}
\end{array}\right]
$$

We can similarly construct unbounded operators in $\mathscr{H}$. Hence, define the closedloop dynamics operator $\dot{A}: \mathscr{H}(\tilde{A}) \subset \dot{\mathscr{H}} \rightarrow \dot{\mathscr{H}}$ on the dense domain $\mathscr{D}(\dot{A}) \xlongequal[1]{\underline{D}}(A) \times \mathbf{R}^{n}$. by $\tilde{A} \tilde{x}=\left(A x+B C_{c} x_{c}, B_{c} C x+A_{c} x_{c}\right)$. Since $\tilde{A}$ can be represented by

$$
\tilde{A}=\left[\begin{array}{cc}
A & B C_{r} \\
B_{c} C & A_{c}
\end{array}\right]=\left[\begin{array}{ll}
A & 0 \\
0 & 0
\end{array}\right]+\left[\begin{array}{cc}
0 & B C_{c} \\
B_{c} C & A_{c}
\end{array}\right]
$$

and since the closed-loop operator

$$
\left[\begin{array}{ll}
A & 0 \\
0 & 0
\end{array}\right]: \mathscr{D}(\tilde{A}) \rightarrow \tilde{H}
$$

generates the strongly continuous semigroup

$$
\left[\begin{array}{cc}
e^{A t} & 0 \\
0 & I_{n c}
\end{array}\right], \quad t \geq 0
$$

it follows from [50, Thm., p. 497] that $\bar{A}$ is also closed and generates a strongly continuous semigroup $e^{\lambda_{1}} \in \mathscr{B}(\dot{\mathscr{X}}), t \geq 0$. To guarantee that $J$ is finite and independent

[^7]of initial conditions we restrict our attention to the set of admissible stabilizing compensators
$$
\mathscr{A} \mathscr{\perp}\left\{\left(A_{o} B_{o}, C_{c}\right): e^{\lambda_{t}} \text { is exponentially stable }\right\}
$$

Hence if $\left(A_{c} B_{c} C_{c}\right) \in \mathscr{A}$ then there exist $\alpha>0$ and $\beta>0$ such that

$$
\begin{equation*}
\left\|e^{\lambda_{1}}\right\| \leqq \alpha e^{-\beta_{1}}, \quad t \geqq 0 \tag{3.6}
\end{equation*}
$$

Since the value of $J$ is independent of the internal realization of the compensator, we can further restrict our attention to

$$
\mathscr{A}_{+} A\left\{\left(A_{c}, B_{c} C_{c}\right) \in \mathscr{A}:\left(A_{c}, B_{c}\right) \text { is controllable and }\left(C_{c}, A_{c}\right) \text { is observable }\right\} .
$$

The following lemma is required for the statement of the Main Theorem.
Lemma 3.1. Suppose $\hat{Q}, \hat{P} \in \mathscr{B}(\mathscr{H})$ have finite rank and are nonnegative definite. Then $\hat{Q} \hat{P}$ is nonnegative semisimple. Furthermore, if $\rho(\hat{Q} \hat{P})=n_{c}$ then there exist $G$, $\Gamma \in \mathscr{B}\left(\mathscr{H}, \mathbf{R}^{n_{c}}\right)$ and positive-semisimple $M \in \mathbf{R}^{n_{c} \times n_{c}}$ such that

$$
\begin{align*}
& \hat{Q} \hat{P}=G^{*} M \Gamma  \tag{3.7}\\
& \Gamma G^{*}=I_{n_{c}}
\end{align*}
$$

Proof. By Lemma 2.4 there exists a finite-dimensional subspace $\mathscr{M} \subset \mathscr{X}$ such that $\hat{Q} \mathcal{M} \subset \mathcal{M}, \hat{Q} \mathcal{M}^{\perp}=0, \hat{P} \mathcal{M} \subset \mathcal{M}$ and $\hat{P} \mathcal{M}^{\perp}=0$. Hence there exists an orthonormal basis for $\mathscr{H}$ with respect to which $\hat{Q}$ and $\hat{P}$ have the infinite-matrix representations

$$
\hat{Q}=\left[\begin{array}{cc}
\hat{Q}_{1} & 0 \\
0 & 0
\end{array}\right], \quad \hat{P}=\left[\begin{array}{cc}
\hat{P}_{1} & 0 \\
0 & 0
\end{array}\right]
$$

where $\hat{Q}_{1}, \hat{P}_{1} \in R^{r \times r}$ are nonnegative definite and $r \triangleq \operatorname{dim} \mathcal{H}$. Since by Lemma 2.6 there exists invertible $\Psi \in R^{r \times r}$ such that $\bar{\Lambda}=\Psi^{-1} \hat{Q}_{1} \hat{P}_{1} \Psi$ is nonnegative diagonal, we have

$$
\hat{Q} \hat{P}=\left[\begin{array}{cc}
\Psi & 0 \\
0 & I_{\infty}
\end{array}\right]\left[\begin{array}{cc}
\tilde{\Lambda} & 0 \\
0 & 0
\end{array}\right]\left[\begin{array}{cc}
\Psi^{-1} & 0 \\
0 & I_{\infty}
\end{array}\right]
$$

which shows that $\hat{Q} \hat{P}$ is nonnegative semisimple. If, furthermore, $\rho(\hat{Q} \hat{P})=n_{c}$ then it is clear that $\Psi$ can be chosen (i.e., modified by an orthogonal matrix) so that

$$
\tilde{\Lambda}=\left[\begin{array}{ll}
\Lambda & 0 \\
0 & 0
\end{array}\right]
$$

where $\Lambda \in R^{n_{c} \times n_{n}}$ is positive diagonal. Hence,

$$
\left.\hat{Q} \hat{P}=\left[\begin{array}{cc}
\Psi & 0 \\
0 & I_{\infty}
\end{array}\right]\left[\begin{array}{c}
I_{n_{c}} \\
0
\end{array}\right]\right] \Lambda\left[\left[\begin{array}{lll}
I_{n_{c}} & 0
\end{array}\right] \quad 0\right]\left[\begin{array}{cc}
\Psi^{-1} & 0 \\
0 & I_{\infty}
\end{array}\right]
$$

which shows that (3.7) and (3.8) are satisfied with

$$
\left.G=\left[\begin{array}{ll}
{\left[\begin{array}{ll}
S^{T} & 0
\end{array}\right]} & 0
\end{array}\right]\left[\begin{array}{cc}
\Psi^{\top} & 0 \\
0 & I_{\infty}
\end{array}\right], \quad M=S^{-1} \Lambda S, \Gamma=\left[\begin{array}{lll}
{\left[S^{-1}\right.} & 0
\end{array}\right] \quad 0\right]\left[\begin{array}{cc}
\Psi^{-1} & 0 \\
0 & I_{\infty}
\end{array}\right]
$$

for all invertible $S \in \mathbf{R}^{n_{e} \times n_{r}} . \quad \square$
We shall refer to $G, \Gamma \in \mathscr{B}\left(\mathscr{H}, \mathbf{R}^{n_{c}}\right)$ and positive-semisimple $M \in R^{n_{c} \times n_{r}}$ satisfying (3.7) and (3.8) as a ( $G, M, \Gamma$ )-factorization of $\hat{Q} \hat{P}$. For convenience in stating the Main Theorem define

$$
\Sigma^{\wedge} B R_{2}^{-1} B^{*}, \quad \bar{\Sigma} \triangleq C^{*} V_{2}^{-1} C .
$$

Main Theorem. Suppose $\left(A_{c} B_{\boldsymbol{c}} C_{c}\right) \in \mathscr{A}_{+}$solves the steady-state fixed-order dynamic-compensation problem. Then there exist nonnegative-definite $Q, P, \hat{Q}, \hat{P} \in \mathscr{B}_{1}(\mathscr{H})$ such that $A_{c} B_{c}$ and $C_{c}$ are given by

$$
\begin{align*}
& A_{\mathrm{c}}=\Gamma(A-Q \bar{\Sigma}-\Sigma P) G^{*},  \tag{3.9}\\
& B_{c}=\Gamma Q C^{*} V_{2}^{-1},  \tag{3.10}\\
& C_{c}=-R_{2}^{-1} B^{*} P G^{*}, \tag{3.11}
\end{align*}
$$

for some ( $G, M, \Gamma$-factorization of $\hat{Q} \hat{P}$, and such that, with $\tau \hat{\#} G^{*} \Gamma$, the following conditions are satisfied:

$$
\begin{equation*}
Q: \mathscr{D}\left(A^{*}\right) \rightarrow \mathscr{X}(A), \quad P: \mathscr{I}(A) \rightarrow \mathscr{\mathscr { F }}\left(A^{*}\right), \tag{3.12a,b}
\end{equation*}
$$

$$
\begin{equation*}
\hat{Q}: \mathscr{H} \rightarrow \mathscr{D}(A), \quad \hat{P}: \mathscr{H} \rightarrow \mathscr{X}\left(A^{*}\right), \tag{3.13a,b}
\end{equation*}
$$

$(3.14 a, b, c)^{5}$

$$
\rho(\hat{Q})=\rho(\hat{P})=\rho(\hat{Q} \hat{P})=n_{0}
$$

$$
\begin{align*}
& 0=(A-\tau Q \bar{\Sigma}) Q+Q(A-\tau Q \bar{\Sigma})^{*}+V_{1}+\tau Q \bar{\Sigma} Q \tau^{*},  \tag{3.15}\\
& 0=\left(A-\Sigma P_{\tau}\right)^{*} P+P\left(A-\Sigma P_{\tau}\right)+R_{1}+\tau^{*} P \Sigma P_{\tau},  \tag{3.16}\\
& 0=\left[(A-\Sigma P) \hat{Q}+\hat{Q}(A-\Sigma P)^{*}+Q \bar{\Sigma} Q\right] \tau^{*},  \tag{3.17}\\
& 0=\left[(A-Q \Sigma)^{*} \hat{P}+\hat{P}(A-Q \bar{\Sigma})+P \Sigma P\right] \tau . \tag{3.18}
\end{align*}
$$

The content of the Main Theorem is clearly a set of necessary conditions which characterize the optimal steady-state fixed-order dynamic compensator when it exists. These necessary conditions consist of a system of four operator equations including a pair of modified Riccati equations (3.15) and (3.16) and a pair of modified Lyapunov equations (3.17) and (3.18). The salient feature of these four equations is the coupling by the operator $\tau \in \mathscr{B}(\mathscr{H})$ which, because of (3.8), is idempotent, i.e., $\tau^{2}=\tau$. In general, $\tau$ is an oblique projection and not an orthogonal projection since there is no requirement that $\tau$ be selfadjoint. Additional features of the Main Theorem will be discussed in the remainder of this section. For convenience, let $G, M, \Gamma, \tau, Q, P, \hat{Q}$ and $\hat{P}$ be as given by the Main Theorem and define $\Lambda \triangleq \operatorname{diag}\left(\lambda_{1}, \cdots, \lambda_{n_{c}}\right)$, where $\lambda_{1} \geqq \lambda_{2} \geqq \cdots \geqq$ $\lambda_{n}>0$ are the eigenvalues of $M$.

We begin by noting that if $x_{c}$ is replaced by $S x_{c}$ where $S \in \mathbb{R}^{n_{r} \times n_{r}}$ is invertible, then an "equivalent" compensator is obtained with $\left(A_{c}, B_{c}, C_{r}\right)$ replaced by $\left(S A_{c} S^{-1}, S B_{c} C_{c} S^{-1}\right.$ ).

Proposition 3.1. Let $\left(A_{c}, B_{c} C_{c}\right) \in \mathscr{A}+$. If $S \in \mathbb{R}^{n_{r} \times n_{c}}$ is invertible then $\left(S A_{c} S^{-1}, S B_{c} C_{c} S^{-1}\right) \in A_{+}$and

$$
\begin{equation*}
J\left(A_{c} B_{c}, C_{c}\right)=J\left(S A_{c} S^{-1}, S B_{c} C_{c} S^{-1}\right) \tag{3.19}
\end{equation*}
$$

Proof. Although the result is obvious from system-theoretic arguments, we shall prove it analytically by utilizing elements of the development in $\S 4$. Define

$$
\tilde{S} \pm\left[\begin{array}{cc}
I_{\infty} & 0 \\
0 & S
\end{array}\right] \in \mathscr{R}(\hat{\mathscr{R}})
$$

and note that replacing $\left(A_{c} B_{c} C_{c}\right)$ by $\left(S A_{S} S^{-1}, S B_{c} C_{c} S^{-1}\right.$ ) is equivalent to replacing $\tilde{\tilde{A}}, \tilde{V}$ and $\tilde{R}$ by $\tilde{S} \bar{A} \tilde{S}^{-1}, \tilde{S} \tilde{\tilde{V}} \tilde{S}^{*}$ and $\tilde{S}^{-*} \tilde{R} \tilde{S}^{-1}$, respectively. If $\alpha, \beta>0$ satisfy (3.6) then a straightforward application of the Hille-Yosida theorem [57, pp. 153-5] shows that

[^8]the strongly continuous semigroup generated by $\bar{S} \bar{A} \bar{S}^{-1}$ satisfies $\left\|e^{\boldsymbol{j} \boldsymbol{j}-1 /}\right\| \leqq$ $\|\tilde{S}\|\left\|\tilde{S}^{-1}\right\| \alpha e^{-\beta!}$, which proves the first assertion. Since $\tilde{S} e^{\boldsymbol{\lambda} / \tilde{S}^{-1}, t \geqq 0 \text {, is also a strongly }}$ continuous semigroup with generator $\tilde{S} \tilde{A} \tilde{S}^{-1}$, it follows that $\tilde{S} \boldsymbol{e}^{\boldsymbol{A} t} \tilde{S}^{-i}=e^{\tilde{S}_{\dot{S}}{ }^{-1} t}$. Hence
and (3.19) follows from $\operatorname{tr} \tilde{Q} \tilde{R}=\operatorname{tr}\left(\tilde{S} \tilde{Q}^{*} \tilde{S}^{*}\right)\left(\tilde{S}^{-*} \tilde{R} \tilde{S}^{-1}\right)$.
In view of Proposition 3.1 one would expect the Main Theorem to apply also to ( $S A_{c} S^{-1}, S B_{c} C_{\uparrow} S^{-1}$ ). Indeed, it may be noted that no claim was made as to the uniqueness of the $(G, M, \Gamma)$-factorization of $\hat{Q} \hat{P}$ used to determine $A_{c} B_{c}$ and $C_{c}$ in (3.9)-(3.11). These observations are reconciled by the following result which shows that a transformation of the compensator state basis corresponds to the alternative factorization $\hat{Q} \hat{P}=\left(S^{-T} G\right)^{T}\left(S M S^{-1}\right)(S \Gamma)$ and, moreover, that all ( $G, M, \Gamma$ )-factorizations of $\hat{Q} \hat{P}$ are related by a nonsingular transformation. Note that $\tau$ remains invariant over the class of factorizations.

Proposition 3.2. If $S \in \mathbb{R}^{n_{c} \times n_{c}}$ is invertible then $\bar{G} \triangleq S^{-T} G, \bar{\Gamma} \triangleq S \Gamma$ and $\bar{M} \triangleq S M S^{-1}$ satisfy
(3.7) ${ }^{\prime}$

$$
\begin{align*}
& \hat{Q} \hat{P}=\bar{G}^{*} \bar{M} \bar{\Gamma} \\
& \bar{\Gamma} \bar{G}^{*}=I_{n_{c}} \tag{3.8}
\end{align*}
$$

Conversely, if $\bar{G}, \bar{\Gamma} \in \mathscr{B}\left(\mathscr{H}, \mathbf{R}^{n_{c}}\right)$ and invertible $\bar{M} \in \mathbf{R}^{n_{c} \times n_{c}}$ satisfy, (3.7)' and (3.8)', then there exists invertible $S \in \mathbb{R}^{n_{c} \times n_{c}}$ such that $\bar{G}=S^{-T} G, \bar{\Gamma}=S \Gamma$ and $\bar{M}=S M S^{-1}$.

Proof. The first part of the proposition is immediate. The second part follows by taking $S \triangleq \bar{M}^{-1} \Gamma G^{*} M^{-1}$, noting $S^{-1}=M \bar{\Gamma} \bar{G}^{*} \bar{M}^{-1}$ and using the identities $\bar{I} G^{*} M \Gamma \bar{G}^{*}=\bar{M}$ and $M \Gamma \bar{G}^{*}=\Gamma \bar{G}^{*} \bar{M}$.

The next result shows that there exists a similarity transformation which simultaneously diagonalizes $\hat{Q} \hat{P}$ and $\tau$.

Proposition 3.3. There exists invertible $\Phi \in \mathscr{B}(\mathscr{H})$ such that

$$
\begin{array}{ll}
\hat{Q}=\Phi^{-1}\left[\begin{array}{cc}
\Lambda_{Q} & 0 \\
0 & 0
\end{array}\right] \Phi^{-*}, & \hat{P}=\Phi^{*}\left[\begin{array}{cc}
\Lambda_{\tilde{P}} & 0 \\
0 & 0
\end{array}\right] \Phi \\
\hat{Q} \hat{P}=\Phi^{-1}\left[\begin{array}{cc}
\Lambda & 0 \\
0 & 0
\end{array}\right] \Phi, & \tau=\Phi^{-1}\left[\begin{array}{cc}
I_{n_{e}} & 0 \\
0 & 0
\end{array}\right] \Phi \tag{3.21a,b}
\end{array}
$$

where $\Lambda_{\hat{Q}}, \Lambda_{\mathcal{F}} \in R^{n_{c} \times n_{s}}$ are positive diagonal and $\Lambda_{\hat{Q}} \Lambda_{\hat{P}}=\Lambda$. Consequently,

$$
\begin{equation*}
\hat{Q}=\tau \hat{Q}, \quad \hat{P}=\hat{P}_{\tau} . \tag{3.22a,b}
\end{equation*}
$$

Proof. Proceeding as in the proof of Lemma 3.1, choose an orthonormal basis for $\mathscr{H}$ with respect to which

$$
\hat{Q}=\left[\begin{array}{cc}
\hat{Q}_{1} & 0 \\
0 & 0
\end{array}\right] \text { and } \hat{P}=\left[\begin{array}{cc}
\hat{P}_{1} & 0 \\
0 & 0
\end{array}\right]
$$

where $\hat{Q}_{1}, \hat{P}_{1} \in \mathbf{R}^{\prime \times r}$ are nonnegative definite. By [55, Thm. 6.2.5, p. 123], there exists invertible $\Psi \in R^{r \times r}$ such that $\tilde{\Lambda}_{Q} \underline{\underline{\Delta}} \Psi \hat{Q}_{1} \Psi^{T}$ and $\hat{\Lambda}_{P}=\Psi^{-T} \hat{P}_{1} \Psi^{-1}$ are nonnegative diagonal. Because of (3.14), it is clear that $\Psi$ can be chosen so that

$$
\tilde{\Lambda}_{Q}=\left[\begin{array}{cc}
\Lambda_{Q} & 0 \\
0 & 0
\end{array}\right] \quad \text { and } \quad \tilde{\Lambda}_{\beta}=\left[\begin{array}{cc}
\Lambda_{\beta} & 0 \\
0 & 0
\end{array}\right]
$$

where $\Lambda_{\delta}, \Lambda_{\beta} \in R^{n_{c} \times n_{c}}$ are positive diagonal. Thus (3.20) holds with

$$
\Phi \triangleq\left[\begin{array}{cc}
\Psi & 0 \\
0 & I_{\infty}
\end{array}\right]
$$

From (3.20) it follows that

$$
\hat{Q} \hat{P}=\Phi^{-1}\left[\begin{array}{cc}
\Lambda_{\mathbf{Q}} \Lambda_{P} & 0 \\
0 & 0
\end{array}\right] \Phi
$$

Now define $\bar{G}=\left[\begin{array}{ll}I_{n_{c}} & 0\end{array}\right] \Phi^{-*}, \bar{M}=\Lambda_{\varrho} \Lambda_{\hat{\rho}}$ and $\bar{\Gamma}=\left[\begin{array}{ll}I_{n_{r}} & 0\end{array}\right] \Phi$ so that (3.7)' and (3.8)' are satisfied. By the second part of Proposition 3.2 there exists invertible $S \in \mathbf{R}^{n_{c} \times n_{c}}$ such that $G=S^{\top} \bar{G}, M=S^{-1} \bar{M} S$ and $\Gamma=S^{-1} \bar{\Gamma}$. Since $M$ and $\bar{M}$ have the same eigenvalues, $\bar{M}=\Lambda$ (modulo an ordering of the diagonal elements) and thus (3.21a) holds. Finally, (3.21b) follows from

$$
\tau=G^{*} \Gamma=\vec{G}^{*} \bar{\Gamma}=\Phi^{-1}\left[\begin{array}{cc}
I_{n_{c}} & 0 \\
0 & 0
\end{array}\right] \Phi
$$

Remark 3.1. Proposition 3.3 shows that $\lambda_{1}, \cdots, \lambda_{n_{c}}$ are the positive eigenvalues of $\hat{Q} \hat{P}$.

Remark 3.2. The simultaneous diagonalization in (3.20) has been effected by a contragredient transformation [55], [58]. For applications of this type of transformation to model reduction and realization problems see [12], [59]-[61]. Simultaneous diagonalization of operators is discussed in [53, p. 181].

The following result validates the precise handling of the unbounded operator $A$ in (3.9), (3.17) and (3.18).

Proposition 3.4. The following relations hold:
(3.23a, b, c)
(3.25a, b)

$$
\begin{aligned}
& \rho(G)=\rho(\Gamma)=\rho(\tau)=n_{\infty} \\
& \tau: \mathscr{H} \rightarrow \mathscr{H}(A), \quad \tau^{*}: \mathscr{H} \rightarrow \mathscr{D}\left(A^{*}\right), \\
& G^{*}: \mathbf{R}^{n_{c}} \rightarrow \mathscr{X}(A), \quad \Gamma^{*}: \mathbb{R}^{n_{c}} \rightarrow \mathscr{D}\left(A^{*}\right) .
\end{aligned}
$$

Proof. From (3.8) and (2.1) it follows that $n_{c}=\rho\left(\Gamma G^{*}\right) \leqq \min \left\{\rho(\Gamma), \rho\left(G^{*}\right)\right\}$. Since $\rho(\Gamma) \leqq n_{\mathrm{c}} \rho(G)=\rho\left(G^{*}\right)$ and $\rho(G) \leqq n_{c}$ (3.23a) and (3.23b) hold. To show (3.23c) either note (3.21b) or use (3.14a) and (3.22) to obtain

$$
n_{c}=\rho(\hat{Q})=\rho(\tau \hat{Q}) \leqq \rho(\tau)=\rho\left(G^{*} \Gamma\right) \leqq \rho(\Gamma)=n_{c}
$$

To prove (3.24a) note that (3.22a) implies $\mathscr{R}(\hat{Q}) \subset \mathscr{R}(\tau)$ and thus $\rho(\hat{Q})=\rho(\tau)$ implies $\mathscr{R}(\hat{Q})=\mathscr{R}(\tau)$, and similarly for (3.24b). Finally, (3.25) follows from (3.23), (3.24), the definition $\tau=G^{*} \Gamma$ and the fact that $\tau^{*}=\Gamma^{*} G$.

Since the domain of $A$ may not be all of $\mathscr{H}$, expressions involving $A$ require special interpretation. First note that because of the range condition (3.25a), the expression (3.9) indeed represents an $n_{c} \times n_{c}$ matrix (see, e.g., [6, p. 80]). Similarly, because of (3.25b), $A_{c}^{T}$ is given by

$$
\begin{equation*}
A_{c}^{T}=G\left(A^{*}-\Sigma Q-P \Sigma\right) \Gamma^{*} \tag{3.26}
\end{equation*}
$$

With regard to (3.15), note that because of (3.12a), the right-hand side of (3.15) is a linear operator with domain $\mathscr{D}\left(A^{*}\right)$. Since $\Theta \stackrel{A}{A}-\tau Q \Sigma Q-Q \bar{\Sigma} Q \tau^{*}+V_{1}+\tau Q \mathcal{E} Q \tau^{*}$ is continuous on $\mathscr{D}\left(A^{*}\right), A Q+Q A^{*}$ has a continuous extension on $\mathscr{H}$ given precisely by - $\Theta$. Similar remarks apply to (3.16). Analogous domain conditions were obtained in [5] for a deterministic infinite-dimensional linear-quadratic control problem with
full-state feedback. Finally, because of (3.24) the right-hand sides of (3.17) and (3.18) denote bounded linear operators on all of $\mathscr{H}$.

It is useful to present an alternative form of the optimal projection equations (3.15)-(3.18). For convenience define the notation

$$
\tau_{\perp} \triangleq I_{x}-\tau
$$

Proposition 3.5. Equations (3.15)-(3.18) are equivalent, respectively, to

$$
\begin{align*}
& 0=A Q+Q A^{*}+V_{1}-Q \bar{\Sigma} Q+\tau_{\perp} Q \bar{\Sigma} Q \tau_{\perp}^{*}  \tag{3.27}\\
& 0=A^{*} P+P A+R_{1}-P \Sigma P+\tau_{\perp}^{*} P \Sigma P \tau_{\perp}  \tag{3.28}\\
& 0=(A-\Sigma P) \hat{Q}+\hat{Q}(A-\Sigma P)^{*}+Q \bar{\Sigma} Q-\tau_{\perp} Q \bar{\Sigma} Q \tau_{\perp}^{*}  \tag{3.29}\\
& 0=(A-Q \Sigma)^{*} \hat{P}+\hat{P}(A-Q \bar{\Sigma})+P \Sigma P-\tau_{\perp}^{*} P \Sigma P \tau_{\perp} \tag{3.30}
\end{align*}
$$

Proof. The equivalence of (3.27) and (3.28) to (3.15) and (3.16) is immediate. Using (3.22a) in the form $\hat{Q}=\hat{Q} \tau^{*}$, we obtain (3.17) = (3.29) $\tau^{*}$. Conversely, from (3.22a) and $[(A-\Sigma P) \hat{Q}]^{*}=\hat{Q}(A-\Sigma P)^{*}$ (see, e.g., $[6$, p. 80]) it follows that $(3.29)=$ $(3.17)+(3.17)^{*}-\tau(3.17)$. Similarly, (3.18) and (3.30) are equivalent.

The form of the optimal projection equations (3.27)-(3.30) helps demonstrate the relationship between the Main Theorem and the classical LQG result when $\operatorname{dim} \mathscr{H}=n<$ $\infty$. In this case we need only note that the $(G, M, \Gamma)$-factorization of $\hat{Q} \hat{P}$ in the "full-order" case $n_{c}=n$ is given by $G=\Gamma=I_{n}$ and $M=\hat{Q} \hat{P}$. Since $\tau=I_{m}$, and thus $\tau_{\perp}=0$, (3.27) and (3.28) reduce to the standard observer and regulator Riccati equations and (3.9)-(3.11) yield the usual LQG expressions. Furthermore, note that in the full-order case

$$
\begin{equation*}
A_{c}=A+B C_{c}-B_{c} C \tag{3.31}
\end{equation*}
$$

and (3.29) and (3.31) can be written as

$$
\begin{align*}
& 0=\left(A_{c}+B_{c} C\right) \hat{Q}+\hat{Q}\left(A_{c}+B_{c} C\right)^{T}+B_{c} V_{2} B_{c}^{T},  \tag{3.32}\\
& 0=\left(A_{c}-B C_{c}\right)^{T} \hat{P}+\hat{P}\left(A_{c}-B C_{c}\right)+C_{c}^{T} R_{2} C_{c} \tag{3.33}
\end{align*}
$$

Since, as is well known, the stability of $\tilde{A}$ corresponds to the stability of $A+B C_{c}=$ $A_{c}+B_{c} C$ and $A-B_{c} C=A_{c}-B C_{c}$ it follows from standard results (e.g., [62, pp. 48, 277]) that the positive-definiteness conditions ( $3.14 a, b$ ) are equivalent to the assumption that ( $A_{c}, B_{c}, C_{c}$ ) is controllable and observable.

To obtain a geometric intepretation of the optimal projection we introduce the quasi-full-state estimate

$$
\hat{x}(t) \triangleq G^{*} x_{c}(t) \in \mathscr{H}
$$

so that $\tau \hat{x}(t)=\hat{x}(t)$ and $x_{c}(t)=\Gamma \hat{x}(t)$. Now, the closed-loop system (3.1)-(3.4) can be written as

$$
\begin{align*}
& \dot{x}(t)=A x(t)-B \hat{C}_{c} \tau \hat{x}(t)+H_{1} w(t)  \tag{3.34}\\
& \dot{\hat{x}}(t)=\tau\left(A+B \hat{C}_{c}-\hat{B}_{c} C\right) \tau \hat{x}(t)+\tau \hat{B}_{c}\left(C x(t)+H_{2} w(t)\right), \tag{3.35}
\end{align*}
$$

where (3.35) is interpreted in the sense of (3.34) since $\hat{x}(t) \in \mathscr{H}$ and where

$$
\hat{B}_{c} \triangleq Q C^{*} V_{2}^{-1}, \quad \hat{C}_{c} \triangleq-R_{2}^{-1} B^{*} P .
$$

It can thus be seen that the geometric structure of the quasi-full-order compensator is entirely dictated by the projection $\boldsymbol{\tau}$. In particular, control inputs $\boldsymbol{\tau} \hat{\boldsymbol{x}}(\boldsymbol{t})$ determined by
(3.35) are contained in $\mathscr{R}(\tau)$ and sensor inputs $\tau \hat{B}_{c} y(t)$ are annihilated unless they are contained in $[\mathcal{N}(\tau)]^{\perp}=\mathscr{R}\left(\tau^{*}\right)$. Consequently, $\mathscr{R}(\tau)$ and $\mathscr{R}\left(\tau^{*}\right)$ are the control and observation subspaces, respectively, of the compensator. Since $\tau$ is not necessarily an orthogonal projection, these (finite-dimensional) subspaces may be different.

From the form of (3.35) it is tempting to suggest that the optimal fixed-order dynamic compensator can be obtained by projecting the full-order (infinitedimensional) $L Q G$ compensator. However, this is generally impossible for the following simple reason. Although the expressions for $A_{c}, B_{c}$ and $C_{c}$ in (3.9)-(3.11) have the form of a projection of the full-order $L Q G$ compensator, the operators $Q$ and $P$ in (3.9)-(3.11) are not the solutions of the usual LQG Riccati equations but instead must be obtained by simultaneously solving all four coupled equations (3.15)-(3.18). This observation reinforces the statement made in § 1 that the optimal fixed-order dynamic compensator cannot in general be obtained by LQG followed by closed-loop controller reduction as in [14] and [15].

We now give an explicit characterization of the optimal projection in terms of $\hat{Q}$ and $\hat{P}$. Since $\hat{Q} \hat{P}$ has finite rank, its Drazin inverse $(\hat{Q} \hat{P})^{D}$ exists (see [63, Thm. 6, p. 108]) and, since $(\hat{Q} \hat{P})^{2}=G^{*} M^{2} \Gamma$, and hence $\rho(\hat{Q} \hat{P})^{2}=\rho(\hat{Q} \hat{P})$, the "index" of $\hat{Q} \hat{P}$ (see [63], [64]) is 1 . In this case the Drazin inverse is traditionally called the group inverse and is denoted by ( $\hat{Q} \hat{P}$ )* (see, e.g., [64, p. 124] or [65]).

Proposition 3.6. The optimal projection $\tau$ is given by

$$
\begin{equation*}
\tau=\hat{Q} \hat{P}(\hat{Q} \hat{P}) \tag{3.36}
\end{equation*}
$$

Proof. It is easy to verify that the conditions characterizing the Drazin inverse [63] for the case that $\hat{Q} \hat{P}$ has index 1 are satisfied by $G^{*} M^{-1} \Gamma$. Hence ( $\left.\hat{Q} \hat{P}\right)^{* *}=G^{*} M^{-1} \Gamma$ and (3.8) implies (3.36).

We now give an alternative characterization of the optimal projection by introducing the following notation from [51, p. 73]. For $\phi, \psi \in \mathscr{H}$ define the operator $\phi \otimes \psi \in$ $\mathscr{B}(\mathscr{H})$ by

$$
(\phi \otimes \psi) x \triangleq\langle x, \phi\rangle \psi, \quad x \in \mathscr{H},
$$

and note that $\rho(\phi \otimes \psi)=1$ if $\phi$ and $\psi$ are both nonzero and $(\phi \otimes \psi)^{*}=\psi \otimes \phi$. Using this notation, (3.21a) can be written as

$$
\begin{equation*}
\Phi \hat{Q} \hat{P} \Phi^{-1}=\sum_{i=1}^{n} \lambda_{i} \xi_{i} \otimes \xi_{i} \tag{3.37}
\end{equation*}
$$

where $\left\{\xi_{i}\right\}_{i=1}^{\infty}$ is an orthonormal basis for $\mathscr{H}$. In terms of the Riesz bases (see e.g., [52, p. 309])

$$
\phi_{i} \triangleq \Phi^{*} \xi_{i}, \quad \psi_{i} \triangleq \Phi^{-1} \xi_{i}, \quad i=1,2, \cdots,
$$

(3.37) is equivalent to

$$
\begin{equation*}
\hat{Q} \hat{P}=\sum_{i=1}^{n_{s}} \lambda_{i} \phi_{i} \otimes \psi_{i} \tag{3.38}
\end{equation*}
$$

which can be regarded as a specialized spectral decomposition of a semisimple operator. We emphasize that, in contrast to the singular value decomposition for compact nonnormal operators (see, e.g., [50, p. 261]), the $\lambda_{i}$ in (3.38) are eigenvalues of $\hat{Q} \hat{P}$ (see Remark 3.1), not singular values. Moreover, although $\left\{\phi_{i}\right\}_{i=1}^{\infty}$ and $\left\{\psi_{i}\right\}_{i=1}^{\infty}$ are bases for $\mathscr{H}$, they are not necessarily orthogonal. They are, however, biorthonormal, i.e., $\left(\phi_{i} \psi_{j}\right)=\delta_{i j}$, and hence $\phi_{i} \otimes \psi_{i}$ is a rank-one projection and $\left(\phi_{i} \otimes \psi_{i}\right)\left(\phi_{j} \otimes \psi_{j}\right)=0, i \neq j$.

Since $\boldsymbol{\tau}$ is a rank- $\boldsymbol{n}_{\boldsymbol{c}}$ projection, it is not surprising that $\boldsymbol{\tau}$ is given precisely by

$$
\begin{equation*}
\tau=\sum_{i=1}^{n} \phi_{i} \otimes \psi_{i} \tag{3.39}
\end{equation*}
$$

The following result summarizes the above observations.
Proposition 3.7. There exist biorthonormal linearly independent sets $\left\{\psi_{1}\right\}_{i=1}^{n} \subset$ $\mathscr{D}(A)$ and $\left\{\phi_{i}\right\}_{i=1}^{n_{i}} \subset \mathscr{D}\left(A^{*}\right)$ such that (3.38) and (3.39) hold. Furthermore, if the $(G, M, \Gamma)$-factorization of $\hat{Q} \hat{P}$ is chosen such that $M=\Lambda$, then, for all $x \in \mathscr{X}$,

$$
\begin{aligned}
& G x=\left(\left\langle x, \psi_{1}\right\rangle, \cdots,\left(x, \psi_{n_{c}}\right\rangle\right)^{T}, \\
& \Gamma x=\left(\left\langle x, \phi_{1}\right\rangle, \cdots,\left(x, \phi_{n_{c}}\right\rangle\right)^{T} .
\end{aligned}
$$

Remark 3.3. Note that $\hat{P} \hat{Q}$ and $\tau^{*}$ are given by

$$
\hat{P} \hat{Q}=\sum_{i=1}^{m_{n}} \lambda_{i} \psi_{i} \otimes \phi_{i} \quad \tau^{*}=\sum_{i=1}^{n_{n}} \psi_{i} \otimes \phi_{i}
$$

and, for all $y^{\perp}\left(y_{1}, \cdots, y_{n_{c}}\right)^{T} \in R^{n_{t}}, G^{*}$ and $\Gamma^{*}$ satisfy

$$
G^{*} y=\sum_{i=1}^{n_{i}} y_{i} \psi_{i} \quad \Gamma^{*} y=\sum_{i=1}^{n_{i}} y_{i} \phi_{i}
$$

4. Proof of the Main Theorem. We state and prove a series of lemmas which allow us to compute the Frechet derivatives of $J$ with respect to $A_{c} B_{c}$ and $C_{c}$ Requiring that these derivatives vanish leads to the necessary conditions in their "primitive" form. A transformation of variables then leads to the form of the necessary conditions (3.9)-(3.18).

Let "u-lim" denote the uniform limit (i.e., limit in operator norm) for bounded linear operators [ 50, p. 150] and, for strongly continuous $S(t) \in \mathscr{B}(\mathscr{H}), t \geqq 0$, interpret the strong integral $\int_{t_{1}^{2}}^{t_{2}} S(t) d t$ according to $\int_{i_{2}}^{t_{2}} S(t) z d t, z \in \mathscr{H}$ [50, p. 152]. Also recall the standard fact $\left[6, p\right.$. 186] that $\left(e^{\mathcal{A} \ell}\right)^{*}=e^{\boldsymbol{\lambda} \cdot t}$ and similarly for $\tilde{\mathcal{A}}$. Throughout this section let $\left(A_{c}, B_{c} C_{c}\right) \in \mathscr{A}+$ and let $\alpha, \beta>0$ satisfy (3.6).

To begin, note that the closed-loop system (3.1)-(3.4) can be written as

$$
\begin{equation*}
\dot{x}(t)=\tilde{A} \tilde{x}(t)+\tilde{H} w(t) \tag{4.1}
\end{equation*}
$$

where

$$
\tilde{H} \pm\left[\begin{array}{c}
H_{1} \\
B_{c} H_{2}
\end{array}\right] \in \mathscr{B}_{2}\left(\mathscr{H}^{\prime} \oplus R^{\prime}\right)
$$

For convenience define the nonnegative-definite operator

$$
\tilde{V} \triangleq \dot{H} \tilde{H}^{*}=\left[\begin{array}{cc}
V_{1} & 0 \\
0 & B_{c} V_{2} B_{c}^{T}
\end{array}\right] \in \mathscr{B}_{1}(\tilde{H})
$$

In terms of the augmented state $\bar{x}(t)$, the performance criterion (3.5) becomes

$$
\begin{equation*}
J\left(A_{\square} B_{o} C_{c}\right)=\lim _{t \rightarrow \infty} \mathrm{E}\langle\tilde{R} \tilde{x}(t), \tilde{x}(t)\rangle, \tag{4.2}
\end{equation*}
$$

where the nonnegative-definite operator $\tilde{R}$ is defined by

$$
\tilde{R} \triangleq\left[\begin{array}{cc}
R_{1} & 0 \\
0 & C_{c}^{r} R_{2} C_{c}
\end{array}\right] \in \mathscr{F}_{1}(\tilde{H})
$$

To write (4.2) in terms of the covariance of $\tilde{x}(t)$, recall [6, p. 308] that the covariance " $E\left[(\xi-E \xi)(\xi-E \xi)^{*}\right]$ " of a Hilbert-space-valued weak random variable $\xi$ is defined to be the nonnegative-definite operator $S$ which satisfies

$$
\langle S y ; z\rangle=E\langle\xi-E \xi, y\rangle\langle\xi-E \xi, z\rangle
$$

for all $y, z$ in the Hilbert space. Hence define [6, p. 317]

$$
\tilde{Q}(t) \triangleq E\left[(\tilde{x}(t)-E \tilde{x}(t))(\tilde{x}(t)-E \tilde{x}(t))^{*}\right] .
$$

Lemma 4.1. $\tilde{Q}^{\underline{\Delta} u-\lim _{t \rightarrow \infty}} \tilde{Q}(t)$ exists and is given by

$$
\begin{equation*}
\tilde{Q}=\int_{0}^{\infty} e^{\lambda_{1}} \tilde{V} e^{\tilde{\lambda}^{*} t} d t \tag{4.3}
\end{equation*}
$$

Furthermore,

$$
\begin{equation*}
J\left(A_{c}, B_{c} C_{c}\right)=\operatorname{tr} \tilde{Q} \hat{R} \tag{4.4}
\end{equation*}
$$

Proof. First compute (as in [6, p. 317])

$$
\begin{aligned}
& \langle\tilde{Q}(t) \tilde{y}, \bar{z}\rangle=E\left(\bar{x}(t)-e^{\tilde{A} t} E \tilde{x}(0), \tilde{y}\right\rangle\left\langle\tilde{x}(t)-e^{\tilde{A} / E \tilde{x}}(0), \tilde{z}\right\rangle \\
& =E\left\langle\int_{0}^{1} e^{\tilde{A}(1-s)} \tilde{H} \tilde{w}(s) d s, \tilde{y}\right\rangle\left\langle\int_{0}^{t} e^{\tilde{\lambda}(t-\sigma)} \tilde{H} \tilde{w}(\sigma) d \sigma, \tilde{z}\right\rangle \\
& +\left\langle\tilde{Q}(0) e^{\boldsymbol{i} \cdot t} \tilde{y}, e^{\boldsymbol{\lambda} \cdot \tau} \bar{z}\right\rangle \\
& =\mathbf{E} \int_{0}^{1} \int_{0}^{1}\left\langle\tilde{w}(s), \tilde{H}^{*} e^{\boldsymbol{\lambda} \cdot(1-s)} \tilde{y}\right\rangle\left(\tilde{w}(\sigma), \tilde{H}^{*} e^{\tilde{\lambda}(t-\sigma)} \tilde{z}\right) d s d \sigma \\
& +\left\langle e^{\bar{A}_{1}} \bar{Q}(0) e^{\boldsymbol{A}^{*} \cdot} \hat{y}, \bar{z}\right\rangle \\
& =\int_{0}^{1}\left\langle e^{\tilde{\boldsymbol{A}}(1-s)} \tilde{\boldsymbol{V}} e^{\tilde{\boldsymbol{A}}(1-s)} \tilde{y}, \tilde{z}\right\rangle d s+\left\langle e^{\tilde{A}_{1}} \tilde{Q}(0) e^{\left.\tilde{A^{\cdot} \cdot} \dot{\boldsymbol{y}}, \tilde{z}\right\rangle,}\right.
\end{aligned}
$$

which shows that $\bar{Q}(t)$ is given by

$$
\tilde{Q}(t)=e^{\bar{\lambda} t} \tilde{Q}(0) e^{\bar{\lambda}^{-1} t}+\int_{0}^{1} e^{\bar{A}_{s}} \tilde{V} e^{\bar{\lambda}^{\pi} s} d s
$$

Clearly, (4.3) makes sense as a strong integral since

$$
\|\tilde{Q}\| \leqq \int_{0}^{\infty}\left\|e^{\tilde{A_{t}}} \hat{V} e^{\tilde{A}^{\cdot} \cdot}\right\| d t \leqq \alpha^{2}\|\tilde{V}\| \int_{0}^{\infty} e^{-2 \beta_{t}} d t<\infty
$$

To demonstrate uniform convergence it need only be noted that

$$
\begin{aligned}
& \|\tilde{Q}-\tilde{Q}(t)\|=\sup _{\|\dot{y}\|=1}\|(\tilde{Q}-\tilde{Q}(t)) \dot{y}\| \\
& =\sup _{\|\hat{y}\|=1} \| \int_{1}^{\infty} e^{\lambda_{s}} \tilde{V} e^{\lambda_{i} \cdot \tilde{y} d s-e^{\lambda_{1}} \tilde{Q}(0) e^{\lambda_{0} \cdot} \tilde{y} \|} \\
& \leq \int_{1}^{\infty}\left\|e^{\lambda_{1}} \hat{V} e^{\lambda^{\prime} s}\right\| d s+\left\|e^{\lambda_{1}} \tilde{Q}(0) e^{\lambda^{-1} \cdot}\right\| \\
& \leq \frac{1}{2} \alpha^{2}\|\tilde{V}\| \beta^{-1} e^{-2 \beta I}+\|\tilde{Q}(0)\| e^{-2 \beta t} \text {. }
\end{aligned}
$$

Next, let $\left\{\phi_{i}\right\}_{i=1}^{\infty}$ be an orthonormal basis for $\tilde{\mathscr{H}}$ and use Parseval's equality to obtain

$$
J\left(A_{c} B_{\infty} C_{c}\right)=\lim _{t \rightarrow \infty} E\left\|\tilde{R}^{1 / 2} \tilde{x}(t)\right\|^{2}=\lim _{t \rightarrow \infty} E \sum_{i=1}^{\infty}\left(\tilde{R}^{1 / 2} \tilde{x}(t), \phi_{i}\right\rangle^{2}
$$

Since

$$
f_{n}(t) \triangleq \sum_{i=1}^{n}\left\langle\tilde{R}^{1 / 2} \tilde{x}(t), \phi_{i}\right\rangle^{2}, \quad t \geqq 0
$$

is nonnegative for each $n$ and is increasing in $n$ for each $t$ with limit $\langle\tilde{R} \tilde{x}(t), \tilde{x}(t)\rangle$, monotone convergence permits expectation-limit interchange. Hence using $\mathbb{E} \bar{x}(r)=$ $e^{\boldsymbol{\lambda}} \mathbf{E x}(0)$ we have

$$
\begin{aligned}
J\left(A_{\sigma} B_{o} C_{c}\right) & =\lim _{t \rightarrow \infty} \sum_{i=1}^{\infty} E\left(\tilde{x}(t), \tilde{R}^{1 / 2} \phi_{i}\right\rangle^{2} \\
& =\lim _{i \rightarrow \infty} \sum_{i=1}^{\infty}\left[\left\langle\tilde{Q}(t) \tilde{R}^{1 / 2} \phi_{i}, \tilde{R}^{1 / 2} \phi_{i}\right\rangle+\left\langle e^{\tilde{A}} \mathbf{E} \tilde{x}(0), \tilde{R}^{1 / 2} \phi_{i}\right\rangle^{2}\right] \\
& =\lim _{i \rightarrow \infty}\left\{\operatorname{tr}\left[\tilde{R}^{1 / 2} \tilde{Q}(t) \tilde{R}^{1 / 2}\right]+\left\|\tilde{R}^{1 / 2} e^{\pi / E} \tilde{x}(0)\right\|^{2}\right\}
\end{aligned}
$$

which by Corollary 2.1 yields (4.4).
We shall also require the "dual" of $\tilde{Q}$ given by

$$
\begin{equation*}
\tilde{P}=\int_{0}^{\infty} e^{\overline{A^{*}} t} \tilde{R} e^{\bar{A} t} d t \tag{4.5}
\end{equation*}
$$

Since $\tilde{\boldsymbol{V}}$ and $\tilde{\boldsymbol{R}}$ are nonnegative definite it is readily seen that $\tilde{Q}$ and $\tilde{P}$ are also nonnegative definite.

Lemma 4.2. $\tilde{Q}, \tilde{P} \in \mathscr{B}_{1}(\tilde{\mathscr{K}})$.
Proof. It suffices to consider $\tilde{Q}$ only since the situation for $\tilde{P}$ is exactly analogous. Since $\tilde{Q}$ is nonnegative definite, Lemma 2.3 can be used. Letting $\left\{\phi_{i}\right\}_{i=1}^{x}$ be an orthonormal basis for $\overline{\mathscr{H}}$, we have

$$
\begin{aligned}
\operatorname{tr} \tilde{Q} & =\sum_{i=1}^{\infty}\left\langle\tilde{Q} \phi_{i}, \phi_{i}\right\rangle=\sum_{i=1}^{\infty}\left\langle\int_{0}^{\infty} e^{\bar{i} \tau} \tilde{V} e^{\bar{i} \cdot} \phi_{i} d t, \phi_{i}\right\rangle \\
& =\lim _{n \rightarrow \infty} \int_{0}^{\infty} \sum_{i=1}^{n}\left\langle\tilde{V} e^{\bar{i} \cdot i} \phi_{i}, e^{\bar{A} \cdot t} \phi_{i}\right\rangle d t
\end{aligned}
$$

Let $f_{n}(t)$ denote the above integrand. Since $\tilde{\boldsymbol{V}}$ is nonnegative definite, $\left\{f_{n}(\cdot)\right\}$ is a monotonically increasing sequence of nonnegative functions such that $f_{n}(t) \rightarrow$ $\operatorname{tr} e^{\lambda_{1}} \dot{V} e^{\dot{\lambda}^{-1}}, t \geqq 0$. Hence, by monotone convergence and Lemma 2.2,

$$
\begin{aligned}
\operatorname{tr} \tilde{Q} & =\int_{0}^{\infty} \operatorname{tr}\left[e^{\lambda_{1}} \tilde{V} e^{\tilde{\lambda}^{-} t}\right] d t \\
& =\int_{0}^{\infty}\left\|e^{\lambda_{t}} \bar{V} e^{\lambda^{-} \cdot}\right\|_{1} d t \leqq \alpha^{2}\|\tilde{V}\|_{1} \int_{0}^{\infty} e^{-2 \beta t} d t<\infty
\end{aligned}
$$

Lemma 4.3. With $\tilde{Q}$ and $\tilde{P}$ given by (4.3) and (4.5) it follows that

$$
\begin{equation*}
\operatorname{tr} \tilde{Q} \tilde{R}=\operatorname{tr} \tilde{V} \tilde{P} \tag{4.6}
\end{equation*}
$$

Proof. For any orthonormal basis $\left\{\phi_{i}\right\}_{i=1}^{\infty}$, of $\hat{\mathscr{H}}$ we have

$$
\begin{aligned}
\operatorname{tr} \tilde{Q} \tilde{R} & =\operatorname{tr} \tilde{R} \tilde{Q}=\sum_{i=1}^{\infty}\left\langle\tilde{R} \int_{0}^{\infty} e^{\lambda_{1}} \tilde{V} e^{\lambda^{-i} t} \phi_{i} d t, \phi_{i}\right\rangle \\
& =\lim _{n \rightarrow \infty} \int_{0}^{\infty} \sum_{i=1}^{n}\left\langle\tilde{R} e^{\left.\tilde{\lambda_{i}} \tilde{V} e^{i^{-i} t} \phi_{i}, \phi_{i}\right\rangle d t}\right.
\end{aligned}
$$

Letting $f_{n}(t)$ denote the above integrand it follows that $f_{n}(t) \rightarrow \operatorname{tr} \tilde{R} e^{\tilde{\lambda} t} \tilde{V} e^{\tilde{A}^{\cdot} t}, t \geqq 0$, and

$$
\left|f_{n}(t)\right| \leqq \sum_{i=1}^{\infty}\left|\left\langle e^{\tilde{\lambda} t} \tilde{\mathbf{V}} e^{\tilde{A}^{+r} t} \phi_{i}, \tilde{R} \phi_{i}\right\rangle\right| \leqq \alpha^{2}\|\tilde{\mathrm{~V}}\| e^{-2 \beta t} \sum_{i=1}^{\infty}\left\|\tilde{R} \phi_{i}\right\| .
$$

If $\left\{\phi_{i}\right\}_{i=1}^{\infty}$ is chosen to be the set of orthonormal eigenvectors of $\tilde{R}$ then Lemma 2.1 implies $\sum_{i=1}^{\infty}\left\|\tilde{R} \phi_{i}\right\|=\|\tilde{R}\|_{1}$ and thus $\left|f_{n}(t)\right|$ is bounded on $[0, \infty)$ by an integrable function. Hence by dominated convergence,

$$
\operatorname{tr} \tilde{Q} \tilde{R}=\int_{0}^{\infty} \operatorname{tr}\left[\tilde{R} e^{\tilde{A_{1}} \tilde{V}} e^{\tilde{A}^{\cdot} t}\right] d t=\int_{0}^{\infty} \operatorname{tr}\left[e^{\tilde{A}^{\bullet} t} \tilde{R} e^{\tilde{A_{1}} \tilde{V}}\right] d t=\int_{0}^{\infty} \sum_{i=1}^{\infty}\left\langle\tilde{V} \phi_{i}, e^{\tilde{A}^{*} t} \tilde{R} e^{\tilde{i} t}\right\rangle d t
$$

And again using dominated convergence,

$$
\operatorname{tr} \tilde{Q} \tilde{R}=\sum_{i=1}^{\infty} \int_{0}^{\infty}\left\langle\tilde{V} \phi_{i} e^{\tilde{A}^{-t}} \tilde{R} e^{\tilde{A_{t}}} \phi_{i}\right\rangle d t=\sum_{i=1}^{\infty}\left\langle\tilde{V} \phi_{i}, \int_{0}^{\infty} e^{\tilde{A^{*}} \cdot \tilde{R}} e^{\tilde{A_{t}}} \phi_{i} d t\right\rangle=\operatorname{tr} \tilde{V} \tilde{P}
$$

The next result is important in that it allows us to treat $\tilde{Q}$ and $\tilde{P}$ as solutions of dual algebraic Lyapunov equations. For a similar result involving groups rather than semigroups see [50, pp. 555-557].

Lemma 4.4. $\hat{Q}$ is given by (4.3) if and only if $\tilde{Q} \in \mathscr{B}(\hat{\mathscr{H}})$ satisfies

$$
\begin{align*}
& \tilde{Q}: \mathscr{D}\left(\tilde{A}^{*}\right) \rightarrow \mathscr{D}(\tilde{A}),  \tag{4.7}\\
& 0=\tilde{A} \tilde{Q}+\tilde{Q} \tilde{A}^{*}+\tilde{V}, \tag{4.8}
\end{align*}
$$

where (4.8) holds in the sense discussed in § 3. Furthermore, $\tilde{P}$ is given by (4.5) if and only if $\tilde{P} \in \mathscr{B}(\tilde{\mathscr{H}})$ satisfies

$$
\begin{align*}
& \tilde{P}: \mathscr{D}(\tilde{A}) \rightarrow \mathscr{B}\left(\tilde{A}^{*}\right),  \tag{4.9}\\
& 0=\tilde{A}^{*} \tilde{P}+\tilde{P} \tilde{A}+\tilde{R} .
\end{align*}
$$

Proof. We consider $\tilde{Q}$ only. To prove necessity let $t^{\prime}>0$. Then for all $t \in\left[0, t^{\prime}\right)$ and $\tilde{x} \in \mathscr{D}\left(\tilde{A}^{*}\right)$ we can write

$$
\begin{aligned}
& e^{\boldsymbol{\lambda} t} \tilde{Q} e^{\boldsymbol{\lambda} \cdot t^{\prime}} \tilde{x}=\int_{0}^{\infty} e^{\boldsymbol{\lambda}(t+s)} \hat{V} e^{\boldsymbol{\lambda} \cdot\left(t^{\prime}+s\right)} \tilde{x} d s \\
& =\int_{1}^{\infty} \cdot e^{\tilde{\lambda} \sigma} \tilde{V} e^{\tilde{A} \cdot \sigma} e^{\tilde{\lambda}^{*}\left(t^{\prime}-1\right)} \tilde{x} d \sigma
\end{aligned}
$$

Hence,
which shows that $e^{\lambda_{1}} \tilde{Q} e^{\lambda^{\boldsymbol{H}} r^{\prime}}$ is strongly differentiable with respect to $t$ for all $t \in\left[0, t^{\prime}\right)$. In particular, setting $t=0$ it follows that $\tilde{Q} e^{\tilde{A}^{\boldsymbol{*}}} r^{\prime} \tilde{x} \in \mathscr{D}(\tilde{\boldsymbol{A}})$ for all $\tilde{x} \in \mathscr{D}\left(\tilde{\mathcal{A}}^{*}\right)$ (see, e.g., [6, p. 173] or [50, p. 485]). Performing the differentiation on the left-hand side of
(4.11) and setting $t=0$ yields

Now fix $\bar{x} \in \mathscr{S}\left(\tilde{A}^{*}\right)$. Then for $\left\{t_{i}\right\}_{i=1}^{\infty}, t_{i}>0, t_{i} \rightarrow 0$, we have

$$
\begin{aligned}
& \tilde{Q} e^{\tilde{i} \cdot} \iota_{i} \in \mathscr{X}(\tilde{A}), \quad i=1,2,3, \cdots, \\
& \tilde{Q} e^{\lambda^{*} \cdot} \cdot \bar{x} \underset{i \rightarrow \infty}{ } \tilde{Q} \tilde{x} .
\end{aligned}
$$

Now consider the sequence $\left\{\tilde{A} \bar{Q} e^{\tilde{X} \cdot t_{i}} \tilde{x}\right\}_{i=1}^{\infty}$. Letting $t^{\prime}=t_{i}$ in (4.12) and using dominated convergence to interchange limit and integration ( $\tilde{A}^{*} \bar{x}$ is a fixed element of $\tilde{\mathscr{C}}$ ), it follows that

$$
\begin{equation*}
\lim _{i \rightarrow \infty} \tilde{A} \tilde{Q} e^{\tilde{\lambda} \cdot \epsilon_{1}} \tilde{x}=-\int_{0}^{\infty} e^{\boldsymbol{\lambda} \sigma} \tilde{V} e^{\dot{\hat{A}} \cdot \sigma} \tilde{A}^{*} \tilde{x} d \sigma-\tilde{V} \tilde{x} \tag{4.13}
\end{equation*}
$$

Since $\tilde{\boldsymbol{A}}$ is closed, $\tilde{Q} \tilde{x} \in \mathscr{D}(\overline{\boldsymbol{A}})$. This proves (4.7). Also, since $\tilde{A}$ is closed we have

$$
\lim _{i \rightarrow \infty} \tilde{A} \tilde{Q} e^{\tilde{i}^{\cdot} \cdot} \tilde{x}=\tilde{A} \tilde{Q} \tilde{x}
$$

which with (4.13) implies

$$
\tilde{A} \tilde{Q} \tilde{x}=-\tilde{Q} \tilde{A}^{*} \tilde{x}-\tilde{y} \tilde{x},
$$

and hence

$$
\left(\tilde{A} \tilde{Q}+\tilde{Q} \tilde{A}^{*}+\tilde{V}\right) \tilde{x}=0, \quad \tilde{x} \in \mathscr{D}\left(\tilde{A}^{*}\right)
$$

as desired.
To prove sufficiency let $\tilde{x} \in \mathscr{O}(\tilde{A})$. Then $e^{\tilde{A}^{*} ‘} \tilde{x} \in \mathscr{D}\left(\tilde{A}^{*}\right), t \geq 0$, and hence

$$
\frac{d}{d t} e^{X_{1}} \tilde{Q} e^{\tilde{A}^{t_{t}}} \bar{x}=e^{\lambda_{1}}\left(\tilde{A} \bar{Q}+\tilde{Q} \tilde{A}^{*}\right) e^{\tilde{X}^{\cdot} t} \tilde{x}
$$

Thus

$$
e^{\tilde{\lambda} t} \tilde{Q} e^{\tilde{\lambda}^{\prime} t} \tilde{x}-\tilde{Q} \tilde{x}=\int_{0}^{1} e^{\lambda_{s}}\left(\tilde{A} \tilde{Q}+\tilde{Q} \tilde{A}^{*}\right) e^{\tilde{\lambda}^{\cdot} s} \tilde{x} d s, \quad \tilde{x} \in \mathscr{D}\left(\tilde{A}^{*}\right)
$$

Extending $\tilde{A} \bar{Q}+\tilde{Q} \hat{A}^{*}$ to all of $\hat{\mathscr{H}}$ we obtain

$$
e^{\lambda_{1}} \tilde{Q} e^{\lambda_{-1}} \tilde{x}-\tilde{Q} \bar{x}=-\int_{0}^{1} e^{\overline{A_{1}}} \tilde{V} e^{\lambda^{-s}} \tilde{x} d s, \quad \tilde{x} \in \hat{\mathscr{H}}
$$

Letting $t \rightarrow \infty$ yields (4.3).
We now introduce some notation which will prove to be most convenient in the following results. For ( $\left.A_{c}^{\prime}, B_{c}^{\prime}, C_{c}^{\prime}\right) \in \mathbf{R}^{n_{c} \times n_{c}} \times \mathbf{R}^{n_{c} \times 1} \times \mathbf{R}^{m \times n_{c}}$ define

$$
\delta_{A_{c}} \triangleq A_{c}^{\prime}-A_{\sigma} \quad \delta_{B_{c}} \triangleq B_{c}^{\prime}-B_{c} \quad \delta_{C_{c}} \triangleq C_{c}^{\prime}-C_{c}
$$

and

$$
\left\|\left(\delta_{A_{c}}, \delta_{B_{c}}, \delta_{C_{c}}\right)\right\| \wedge\left\|\delta_{A_{c}}\right\|+\left\|\delta_{B_{c}}\right\|+\left\|\delta_{C_{c}}\right\|
$$

Furthermore, let $\tilde{A}^{\prime}, \tilde{V}^{\prime}$ and $\tilde{R}^{\prime}$ denote $\tilde{A}, \bar{V}$ and $\tilde{R}$ with ( $A_{c}, B_{c} C_{c}$ ) replaced by
( $A_{c}^{\prime}, B_{c}^{\prime}, C_{c}^{\prime}$ ) and define

$$
\begin{aligned}
& \delta_{\lambda} \wedge \tilde{A}^{\prime}-\tilde{A}=\left[\begin{array}{cc}
0 & B \delta_{C_{c}} \\
\delta_{B_{c}} C & \delta_{A_{c}}
\end{array}\right], \\
& \delta_{\bar{V}} \triangleq \tilde{V}^{\prime}-\tilde{V}=\left[\begin{array}{cc}
0 & 0 \\
0 & B_{c} V_{2} \delta_{B_{c}}^{T}+\delta_{B_{c}} V_{2} B_{c}^{T}+\delta_{B_{c}} V_{2} \delta_{B_{c}}^{T}
\end{array}\right], \\
& \delta_{K} \triangleq \tilde{R}^{\prime}-\tilde{R}=\left[\begin{array}{ll}
0 & 0 \\
0 & C_{c}^{T} R_{2} \delta_{C_{c}}+\delta_{C_{c}^{T}}^{T} R_{2} C_{c}+\delta_{C_{c}}^{T} R_{2} \delta_{C_{c}}
\end{array}\right] .
\end{aligned}
$$

We shall also write $\tilde{Q}^{\prime}, \tilde{P}^{\prime}$ for $\tilde{Q}, \tilde{P}$ as given by (4.3) and (4.5) with $\tilde{A}, \tilde{\mathrm{~V}}, \tilde{R}$ replaced by $\tilde{A}^{\prime}, \hat{V}^{\prime}, \tilde{R}^{\prime}$ and define

$$
\delta_{\delta} \triangleq \tilde{Q^{\prime}}-\tilde{Q}, \quad \delta_{\mathcal{F}} \hat{A}^{\tilde{P}^{\prime}}-\tilde{P} .
$$

Lemma 4.5. $A$ is open.
Proof. Let $\left(A_{c} B_{c} C_{c}\right) \in \mathscr{A}$ be arbitrary and consider the open set

$$
\begin{equation*}
N \perp\left\{\left(A_{c}^{\prime}, B_{c}^{\prime}, C_{c}^{\prime}\right) \in \mathbf{R}_{c}^{n_{c} \times n_{c}} \times \mathbf{R}_{c}^{n_{c} \times 1} \times \mathbf{R}^{m \times n_{c}}:\left\|\left(\delta_{A_{c}}, \delta_{B_{c}} \delta_{C_{c}}\right)\right\|<\beta / 2 \alpha \gamma\right\}, \tag{4.14}
\end{equation*}
$$

where $\gamma^{\wedge} \max \{1,\|B\|,\|C\|\}$. Then, since $\tilde{A}^{\prime}=\tilde{A}+\delta_{\lambda}$ and $\delta_{\tilde{A}} \in \mathscr{B}(\tilde{\mathscr{X}})$ it follows from Theorem 2.1, p. 497 of [50], that for all ( $A_{c}^{\prime}, B_{c}^{\prime}, C_{c}^{\prime}$ ) $\in N$ and $t \geq 0$,

$$
\left\|e^{\lambda^{\prime} t}\right\| \leq \alpha e^{\left(-\beta+a \mid \delta_{A} \|\right)!} \leq \alpha e^{-\beta_{1} / 2}
$$

Hence, $N \subset \mathscr{A}$, as desired.
Lemma 4.6. There exists $c>0$ such that

$$
\begin{align*}
& \left\|\delta_{\delta}\right\| \leqq c\left\|\left(\delta_{A_{i}}, \delta_{B_{c}}, \delta_{c_{c}}\right)\right\|,  \tag{4.15}\\
& \left\|\delta_{\boldsymbol{F}}\right\| \leqq c\left\|\left(\delta_{A_{c}}, \delta_{\mathcal{B}_{c},} \delta_{c_{c}}\right)\right\|, \tag{4.16}
\end{align*}
$$

for all $\left(A_{c}^{\prime}, B_{c}^{\prime}, C_{c}^{\prime}\right) \in N$, where $N \subset \mathscr{A}$ is the open neighborhood of $\left(A_{c} B_{c} C_{c}\right)$ defined by (4.14).

Proof. We consider (4.15) only. Since $\left\|e^{\lambda_{t}^{\prime}}\right\| \leqslant \alpha e^{-\beta t / 2}, t \geqq 0,\left(A_{c}^{\prime}, B_{c}^{\prime}, C_{c}^{\prime}\right) \in N$, it follows that

$$
\begin{aligned}
& \text { (4.17) } \leq \alpha\left(\|\tilde{V}\|+\left\|\delta_{\tilde{v}}\right\|\right) \int_{0}^{\infty}\left\|e^{\left(\lambda^{*}+\delta^{*}{ }_{\beta}^{\prime \prime}\right.}-e^{\lambda \cdot t}\right\| e^{-\beta_{s} / 2} d t \\
& +\alpha^{2}\left\|\delta_{\tilde{v}}\right\| \int_{0}^{\infty} e^{-3 \beta_{t} / 2} d t+\alpha\|\tilde{V}\| \int_{0}^{\infty}\left\|e^{\left(\hat{\lambda}+\delta_{t}\right) t}-e^{i \lambda}\right\| e^{-\beta / 2} d t \\
& =\alpha\left(2\|\tilde{V}\|+\left\|\delta_{\dot{v}}\right\|\right) \int_{0}^{\infty}\left\|e^{\left(\lambda+\delta_{j}\right) t}-e^{\lambda_{t}}\right\| e^{-\beta_{1} / 2} d t+\frac{2 \alpha^{2}}{3 \beta}\left\|\delta_{\dot{v}}\right\| .
\end{aligned}
$$

From [50, p. 497], it follows that the perturbed semigroup $e^{\left(\lambda+s_{1}\right) t}$ has an expansion

$$
e^{\left(X+s_{1}\right)^{\prime}}=e^{A_{t}}+\sum_{i=1}^{\infty} U_{1}(t), \quad t \geq 0,
$$

where $U_{i}(t) \in \mathscr{F}(\mathscr{\mathscr { H }}), t \geq 0$, satisfy the estimates

$$
\left\|U_{i}(t)\right\| \leqq \alpha^{i+1}\left\|\delta_{A}\right\|^{i} e^{-\beta_{1}} t^{i} / i!
$$

Hence, for all $\left(A_{c}^{\prime}, B_{c}^{\prime}, C_{c}^{\prime}\right) \in N$,

$$
\begin{equation*}
\left\|e^{\left(\lambda+\delta_{i}\right)^{\prime}}-e^{\tilde{A^{\prime}}}\right\| \leqq \sum_{i=1}^{\infty}\left\|U_{i}(t)\right\| \leqq \alpha e^{-\beta \prime}\left[e^{a\left\|\delta_{i}\right\|^{\prime}}-1\right] \tag{4.18}
\end{equation*}
$$

From (4.17), (4.18) and the relations $\left\|\delta_{A}\right\| \leqq \gamma\left\|\left(\delta_{A_{c}}, \delta_{B_{c}}, \delta_{C_{c}}\right)\right\|<\beta / 2 \alpha$ and

$$
\int_{0}^{\infty}\left[e^{a\left\|\delta_{i}\right\|^{\prime}}-1\right] e^{-3 \beta t / 2} d t<\frac{\alpha \gamma}{3 \beta^{2}}\left\|\left(\delta_{A_{c}}, \delta_{B_{c}}, \delta_{C_{c}}\right)\right\|
$$

it follows that

$$
\begin{aligned}
\left\|\delta_{\tilde{Q}}\right\| \leqq & \frac{2 \alpha^{3} \gamma}{3 \beta^{2}}\left(2\|\tilde{V}\|+\left\|\delta_{\tilde{v}}\right\|\right)\left\|\left(\delta_{A_{c}}, \delta_{B_{c}}, \delta_{C_{c}}\right)\right\| \\
& +\frac{2 \alpha^{2}}{3 \beta}\left(2\left\|B_{c} V_{2}\right\|\left\|\delta_{B_{c}}\right\|+\left\|V_{2}\right\|\left\|\delta_{B_{c}}\right\|^{2}\right)
\end{aligned}
$$

which yields (4.15). $\square$
Since $\tilde{Q}, \tilde{P} \in \mathscr{B}(\hat{\mathscr{H}})$ we can write

$$
\tilde{Q}=\left[\begin{array}{ll}
Q_{1} & Q_{12} \\
Q_{12}^{*} & Q_{2}
\end{array}\right], \quad \tilde{P}=\left[\begin{array}{ll}
P_{1} & P_{12} \\
P_{12}^{*} & P_{2}
\end{array}\right]
$$

where $Q_{1} \in \mathscr{B}(\mathscr{H}), Q_{12} \in \mathscr{B}\left(R^{n_{c}}, \mathscr{H}\right), Q_{2} \in R^{n_{c} \times n_{c}}$ and similarly for $P_{1}, P_{12}$ and $P_{2}$. Note that $Q_{1}, Q_{2}, P_{1}$ and $P_{2}$ are nonnegative definite. Also, define the notation

$$
\hat{P} \hat{Q}=\left[\begin{array}{ll}
Z_{1} & Z_{12} \\
Z_{21} & Z_{2}
\end{array}\right]
$$

where

$$
\begin{array}{ll}
Z_{1} \triangleq P_{1} Q_{1}+P_{12} Q_{12}^{*}, & Z_{12} \triangleq P_{1} Q_{12}+P_{12} Q_{2} \\
Z_{21} \triangleq P_{12}^{*} Q_{1}+P_{2} Q_{12}^{*}, & Z_{2} \triangleq P_{12}^{*} Q_{12}+P_{2} Q_{2}
\end{array}
$$

and, for $\left(A_{c}^{\prime}, B_{c}^{\prime}, C_{c}^{\prime}\right) \in \mathscr{A}$, let

$$
\delta_{J}\left(\delta_{A_{c}}, \delta_{B_{c}}, \delta_{C_{c}}\right) \triangleq J\left(A_{c}^{\prime}, B_{c}^{\prime}, C_{c}^{\prime}\right)-J\left(A_{c} B_{o} C_{c}\right)
$$

Lemma 4.7. Let $\left(A_{c}^{\prime}, B_{c}^{\prime}, C_{c}^{\prime}\right) \in \mathscr{A}$. Then

$$
\begin{equation*}
\delta_{J}\left(\delta_{A_{c},}, \delta_{B_{c}}, \delta_{C_{c}}\right)=\mathscr{L}\left(\delta_{A_{c^{\prime}}} \delta_{B_{c},}, \delta_{C_{c}}\right)+o\left(\left\|\left(\delta_{A_{t}}, \delta_{B_{c}}, \delta_{C_{c}}\right)\right\|\right) \tag{4.19}
\end{equation*}
$$

where

$$
\begin{align*}
\mathscr{L}\left(\delta_{A_{c}}, \delta_{B_{c}}, \delta_{C_{c}}\right) \triangleq & 2 \operatorname{tr}\left[Z_{2} \delta_{A_{c}}\right]+2 \operatorname{tr}\left[\left(V_{2} B_{c}^{T} P_{2}+C Z_{21}^{*}\right) \delta_{B_{c}}\right]  \tag{4.20}\\
& \left.+2 \operatorname{tr}\left[Q_{2} C_{c}^{T} R_{2}+Z_{12}^{*} B\right) \delta_{C_{c}}\right]
\end{align*}
$$

and

$$
\begin{equation*}
\lim _{\left(\delta_{A_{c}}, \delta_{c}, b_{C_{c}}\right) \rightarrow 0}\left\|\left(\delta_{A_{c}}, \delta_{B_{c}}, \delta_{C_{c}}\right)\right\|^{-1} O\left(\left\|\left(\delta_{A_{c}}, \delta_{B_{c}}, \delta_{C_{c}}\right)\right\|\right)=0 \tag{4.21}
\end{equation*}
$$

Proof. Combining (4.8) and (4.10) with (4.6), J can be written as

$$
J\left(A_{c} B_{n} C_{c}\right)=\operatorname{tr}[\tilde{Q} \tilde{R}+\tilde{P} \tilde{V}]+\frac{1}{2} \operatorname{tr}\left[\tilde{Q} \mathrm{cl}\left(\tilde{A}^{*} \tilde{P}+\tilde{P} \tilde{A}\right)+\tilde{P} \operatorname{cl}\left(\tilde{A} \tilde{Q}+\tilde{Q} \tilde{A}^{*}\right)\right]
$$

and likewise for ( $A_{c}^{\prime}, B_{c}^{\prime}, C_{c}^{\prime}$ ), where "cl" denotes closure (i.e., extension) of a bounded operator to all of $\mathscr{\mathscr { K }}$. Now using the identity

$$
\operatorname{tr}\left[\tilde{Q}^{\prime} \tilde{R}^{\prime}+\tilde{P}^{\prime} \tilde{V}^{\prime}\right]-\operatorname{tr}[\tilde{Q} \tilde{R}+\tilde{P} \tilde{V}]=\operatorname{tr}\left[\tilde{Q} \delta_{\tilde{R}}+\tilde{P} \delta_{\tilde{V}}\right]+\operatorname{tr}\left[\delta_{\tilde{Q}} \tilde{R}^{\prime}+\delta_{\tilde{F}} \tilde{V}^{\prime}\right]
$$

we can compute

$$
\begin{aligned}
& \delta_{J}\left(\delta_{A_{c}}, \delta_{B_{c}}, \delta_{C_{c}}\right)=\operatorname{tr}\left[\bar{Q}_{\tilde{R}}+\tilde{P} \delta_{\tilde{v}}\right]+\frac{1}{\frac{1}{2}} \operatorname{tr}\left[\tilde{Q} \mathrm{Cl}\left(\tilde{\mathcal{A}}^{*}\left(\tilde{P}+\delta_{\boldsymbol{\beta}}\right)+\left(\tilde{P}+\delta_{\tilde{P}}\right) \tilde{A}^{\prime}\right)\right] \\
& +\frac{1}{2} \operatorname{tr}\left[\delta_{\tilde{S}} \operatorname{cl}\left(\tilde{A}^{\prime *} \tilde{\boldsymbol{P}}^{\prime}+\tilde{\boldsymbol{P}}^{\prime} \overline{\boldsymbol{A}}^{\prime}\right)\right] \\
& +\frac{1}{2} \operatorname{tr}\left[\tilde{P} \mathrm{cl}\left(\tilde{A}^{\prime}\left(\tilde{Q}+\delta_{\tilde{Q}}\right)+\left(\tilde{Q}+\delta_{\tilde{Q}}\right) \tilde{A}^{\prime *}\right)\right] \\
& +\frac{1}{2} \operatorname{tr}\left[\delta_{\tilde{P}} \mathrm{Cl}\left(\tilde{A}^{\prime} \tilde{Q}^{\prime}+\tilde{Q}^{\prime} \tilde{A}^{\prime *}\right)\right] \\
& -\frac{1}{2} \operatorname{tr}\left[\tilde{Q} \mathrm{cl}\left(\tilde{A}^{*} \tilde{P}+\tilde{P} \tilde{A}\right)+\tilde{P} \mathrm{cl}\left(\tilde{A} \tilde{Q}+\tilde{Q} \tilde{A}^{*}\right)\right] \\
& +\operatorname{tr}\left[\delta_{\widehat{Q}} \tilde{R}^{\prime}+\delta_{\beta} \tilde{V}^{\prime}\right] .
\end{aligned}
$$

Using $\tilde{A^{\prime}}=\tilde{A}+\delta_{\bar{A}}$ and combining the second, fourth and sixth terms yields

$$
\delta_{j}\left(\delta_{A_{c}} \delta_{B_{c}}, \delta_{C_{c}}\right)=\Lambda+\Omega,
$$

where

$$
\begin{aligned}
& \Lambda \hat{A} \operatorname{tr}\left[\tilde{Q} \delta_{\tilde{R}}+\tilde{P} \delta_{\dot{v}}\right]+\frac{1}{2} \operatorname{tr}\left[\tilde{Q}\left(\delta_{\dot{A}}^{*} \tilde{P}+\tilde{P} \delta_{\tilde{X}}\right)+\tilde{P}\left(\delta_{\mathcal{A}} \tilde{Q}+\tilde{Q} \delta_{\hat{A}}^{*}\right)\right] \\
& \quad=\operatorname{tr}\left[\tilde{Q} \delta_{\tilde{R}}+\tilde{P} \delta_{\tilde{i}}\right]+2 \operatorname{tr}\left[\delta_{i} \tilde{Q} \tilde{P}\right]
\end{aligned}
$$

and

$$
\begin{aligned}
& +\frac{1}{2} \operatorname{tr}\left[\delta_{\delta} \mathrm{cl}\left(\tilde{\mathcal{A}^{\prime *}} \tilde{P}^{\prime}+\tilde{P}^{\prime} \tilde{A}^{\prime}\right)+\delta_{\tilde{\beta}} \mathrm{cl}\left(\tilde{A}^{\prime} \tilde{Q}^{\prime}+\tilde{Q}^{\prime} \tilde{A}^{\prime * *}\right)\right]+\operatorname{tr}\left[\delta_{\tilde{O}} \tilde{R}^{\prime}+\delta_{\tilde{\mathcal{P}}} \tilde{V}^{\prime}\right] .
\end{aligned}
$$

Computing

$$
\operatorname{tr}\left[\tilde{Q} \delta_{\tilde{R}}+\tilde{P} \delta_{\tilde{V}}\right]=2 \operatorname{tr}\left[V_{2} B_{c}^{T} P_{2} \delta_{B_{c}}\right]+2 \operatorname{tr}\left[Q_{2} C_{c}^{T} R_{2} \delta_{C_{c}}\right]+\operatorname{tr}\left[P_{2} \delta_{B_{c}} V_{2} \delta_{B_{c}}^{T}+Q_{2} \delta_{c_{c}}^{T} R_{2} \delta_{C_{c}}\right]
$$

and

$$
2 \operatorname{tr}\left[\delta_{A} \tilde{Q} \tilde{P}\right]=2 \operatorname{tr}\left[Z_{2} \delta_{A_{\mathrm{C}}}\right]+2 \operatorname{tr}\left[C Z_{21}^{*} \delta_{B_{\mathrm{B}}}\right]+2 \operatorname{tr}\left[Z_{12}^{*} B \delta_{C_{\mathrm{C}}}\right]
$$

and retaining first-order terms, we obtain (4.20).
To evaluate $\Omega$, use (4.8) and (4.10) to replace $\tilde{R}^{\prime}$ and $\hat{V}^{\prime}$ iat the last term in $\Omega$ and write $\tilde{A}^{\prime}=\tilde{A}+\delta_{\bar{A}}$, to obtain

$$
\begin{align*}
& \Omega=\frac{1}{2} \operatorname{tr}\left[\tilde{Q} \mathrm{cl}\left(\tilde{A}^{*} \delta_{\tilde{\mathcal{F}}}+\delta_{\tilde{\mathcal{F}}} \tilde{A}\right)+\tilde{P} \mathrm{cl}\left(\tilde{A} \delta_{\tilde{\rho}}+\delta_{\boldsymbol{Q}} \tilde{A}^{*}\right)\right] \\
& +\frac{1}{2} \operatorname{tr}\left[\tilde{Q}\left(\delta_{\hat{A}}^{*} \delta_{\tilde{P}}+\delta_{\beta} \delta_{\bar{A}}\right)+\tilde{P}\left(\delta_{\bar{A}} \delta_{\tilde{Q}}+\delta_{\delta} \delta_{\hat{A}}^{*}\right)\right]  \tag{4.22}\\
& -\frac{1}{2} \operatorname{tr}\left[\delta_{\tilde{Q}} \mathrm{cl}\left(\tilde{A}^{\prime *} \tilde{P}^{\prime}+\tilde{P}^{\prime} \tilde{A}^{\prime}\right)+\delta_{\tilde{B}} \mathrm{cl}\left(\tilde{A}^{\prime} \tilde{Q}^{\prime}+\tilde{Q}^{\prime} \tilde{A}^{\prime *}\right)\right] .
\end{align*}
$$

Next we note that

$$
\begin{equation*}
\operatorname{tr}\left[\tilde{Q} \operatorname{cl}\left(\tilde{A}^{*} \delta_{\mathcal{\beta}}+\delta_{\beta} \tilde{A}^{*}\right)\right]=\operatorname{tr}\left[\delta_{\bar{\beta}} \operatorname{cl}\left(\tilde{A} \tilde{Q}+\tilde{Q} \tilde{A}^{*}\right)\right] . \tag{4.23}
\end{equation*}
$$

To see this we observe that by arguments similar to those used in the proof of Lemma 4.4 and the fact that $\delta_{\beta}: \mathscr{D}(\bar{A}) \rightarrow \mathscr{D}\left(\bar{A}^{*}\right)$ it follows that

$$
\delta_{\beta}=-\int_{0}^{\infty} e^{\boldsymbol{A} \cdot t} \operatorname{cl}\left(\tilde{A}^{*} \delta_{\beta}+\delta_{\beta} \tilde{A}\right) e^{\tilde{X}_{1}} d t .
$$

Now, using the technique of Lemma 4.3 with the role of $\tilde{R}$ played by $-\mathrm{cl}\left(\bar{A}^{*} \delta_{\beta}+\delta_{\beta} \tilde{A}\right)$,
we see that

$$
\operatorname{tr}\left[\tilde{Q} \operatorname{cl}\left(\tilde{A}^{*} \delta_{\mathcal{P}}+\delta_{\beta} \tilde{A}\right)\right]=-\operatorname{tr}\left[\delta_{\tilde{\beta}} \tilde{V}\right]=\operatorname{tr}\left[\delta_{\tilde{\beta}} \operatorname{cl}\left(\tilde{A} \tilde{Q}+\tilde{Q}_{\dot{A}} \tilde{A}^{*}\right)\right]
$$

Similarly, it can be shown that

$$
\begin{equation*}
\operatorname{tr}\left[\tilde{P} \mathrm{cl}\left(\tilde{A} \delta_{\delta}+\delta_{\S} \tilde{A}^{*}\right)\right]=\operatorname{tr}\left[\delta_{\dot{Q}} \mathrm{cl}\left(\tilde{A^{*}} \tilde{P}+\tilde{P} \tilde{A}\right)\right] \tag{4.24}
\end{equation*}
$$

Now substitute (4.23) and (4.24) into (4.22) and rearrange the second term in (4.22) so that

$$
\begin{aligned}
& \boldsymbol{\Omega}=\frac{1}{2} \operatorname{tr}\left[\delta_{\delta} \operatorname{cl}\left(\tilde{A}^{*} \tilde{P}+\tilde{P} \tilde{A}\right)+\delta_{\tilde{P}} \operatorname{cl}\left(\tilde{A} \tilde{Q}+\tilde{Q} \tilde{A}^{*}\right)\right] \\
& +\frac{1}{2} \operatorname{tr}\left[\delta_{\boldsymbol{O}}\left(\delta_{\dot{A}}^{*} \tilde{P}+\tilde{P} \delta_{\boldsymbol{A}}\right)+\delta_{\mathcal{F}}\left(\delta_{\boldsymbol{A}} \bar{Q}+\overline{\mathrm{Q}} \delta_{\dot{A}}^{*}\right)\right] \\
& -\frac{1}{2} \operatorname{tr}\left[\delta_{\varrho} \mathrm{cl}\left(\tilde{A}^{\prime *} \tilde{P}^{\prime}+\tilde{P}^{\prime} \tilde{A}^{\prime}\right)+\delta_{\beta} \mathrm{cl}\left(\tilde{A}^{\prime} \tilde{Q}^{\prime}+\tilde{Q}^{\prime} \tilde{A}^{\prime *}\right)\right] \\
& =-\frac{1}{2} \operatorname{tr}\left[\delta_{\delta} \mathrm{cl}\left(\tilde{A}^{\prime *} \delta_{\tilde{P}}+\delta_{\bar{P}} \tilde{A}^{\prime}\right)+\delta_{\bar{\beta}} \mathrm{cl}\left(\tilde{A}^{\prime} \delta_{\boldsymbol{Q}}+\delta_{\bar{\Phi}} \bar{A}^{\prime *}\right)\right] .
\end{aligned}
$$

Using (4.8) to obtain

$$
0=\tilde{A}^{\prime} \delta_{\dot{Q}}+\delta_{\bar{Q}} \tilde{A}^{\prime *}+\delta_{A} \tilde{Q}+\tilde{Q} \delta_{\dot{\lambda}}^{*}+\delta_{\tilde{v}}
$$

and (4.10) to obtain a similar relation involving $\tilde{P}$, we have

$$
\left.\Omega=\operatorname{tr}\left[\delta_{\bar{Q}}\left(\delta_{\hat{\lambda}}^{*} \tilde{P}+\tilde{P} \delta_{\hat{A}}+\delta_{\dot{R}}\right)\right]+\operatorname{tr}\left[\delta_{\beta} \tilde{Q}+\tilde{Q} \delta_{\hat{\lambda}}^{*}+\delta_{\tilde{v}}\right)\right] .
$$

Restricting ( $A_{c}^{\prime}, B_{c}^{\prime}, C_{c}^{\prime}$ ) to $N$ (see (4.14)), using Lemma 4.6 and noting that $\delta_{\bar{A}}$ and $\delta_{\mathcal{R}}$ have finite rank, it follows that there exists $c_{1}>0$ such that

$$
\begin{equation*}
\|\Omega\| \leqq c_{1}\left\|\left(\delta_{A_{c}}, \delta_{B_{c}}, \delta_{C_{c}}\right)\right\|^{2} \tag{4.25}
\end{equation*}
$$

Combining $\Omega$ with the second-order terms in $\Lambda$ yields the desired result.
Lemma 4.8. $\mathscr{A}_{+}$is open.
Proof. From the "generic" property of controllability and observability [62, p. 44] there exists an open neighborhood of ( $A_{c}, B_{c} C_{c}$ ) each of whose elements is minimal. Combining this fact with Lemma 4.5 yields the desired result.

Lemma 4.9. $Q_{2}$ and $P_{2}$ are positive definite.
Proof. First note that expanding the $R^{n_{c} \times n_{c}}$-component of the Lyapunov equation (4.8) yields (4.50) below. By a minor extension of results from [66] or [67], (4.50) can be rewritten as

$$
0=\left(A_{c}+B_{c} C Q_{12} Q_{2}^{+}\right) Q_{2}+Q_{2}\left(A_{c}+B_{c} C Q_{12} Q_{2}^{+}\right)^{T}+B_{c} V_{2} B_{c}^{T}
$$

where $Q_{2}^{+}$is the Moore-Penrose or Drazin generalized inverse of $Q_{2}$. Next note that since $\left(A_{c}, B_{c}\right)$ is controllable then so is $\left(A_{c}+B_{c} C Q_{12} Q_{2}^{+}, B_{c} V_{2}^{1 / 2}\right)$. Now, since $Q_{2}$ and $B_{c} V_{2} B_{c}^{T}$ are nonnegative definite, it follows from [62, Lemma 12.2] that $Q_{2}$ is positive definite. Similar arguments show that $P_{2}$ is positive definite.

Having established Lemmas 4.1-4.9, we can now proceed with the proof of the Main Theorem. Let $\left(A_{c}, B_{c}, C_{c}\right) \in \mathscr{A}_{+}$be as in the Main Theorem and consider (4.19)
 linear functional and $\mathscr{A}_{+}$is open, the convergence in (4.21) implies that $\mathscr{L}$ is precisely the Frechet derivative of $J$ with respect to $\left(A_{c} B_{c} C_{c}\right)$. Since $\mathscr{A}_{+}$is open, the optimality of ( $A_{a} B_{c} C_{c}$ ) implies

$$
\begin{equation*}
\mathscr{L}\left(\delta_{A_{c}}, \delta_{\mathbf{D}_{c}}, \delta_{C_{c}}\right)=0 \tag{4.26}
\end{equation*}
$$

for all ( $\delta_{A_{c}}, \delta_{g_{c}}, \delta_{C_{c}}$ ). Clearly, (4.26) is equivalent to

$$
\begin{equation*}
Z_{2}=0 \tag{4.27}
\end{equation*}
$$

$$
\begin{align*}
& V_{2} B_{c}^{T} P_{2}+C Z_{21}^{*}=0,  \tag{4.28}\\
& Q_{2} C_{c}^{\tau} R_{2}+Z_{12}^{*} B=0 . \tag{4.29}
\end{align*}
$$

Thus, $B_{c}$ and $C_{s}$ are given by

$$
\begin{align*}
& B_{\mathrm{c}}=-P_{2}^{-1} Z_{21} C^{*} V_{2}^{-1},  \tag{4.30}\\
& C_{c}=-R_{2}^{-1} B^{*} Z_{12} Q_{2}^{-1}
\end{align*}
$$

Although $B_{c}$ and $C_{c}$ are now determined in terms of $\tilde{Q}$ and $\tilde{P}, A_{c}$ remains to be found. Moreover, $\tilde{Q}$ and $\tilde{P}$ themselves depend (via (4.8) and (4.10)) on $B_{c}$ and $C_{c}$ Hence our task now is to consolidate and simplify (4.7)-(4.10), (4.27), (4.30) and (4.31) to obtain the more tractable conditions (3.9)-(3.18). To this end let us define new variables
(4.32a, b)

$$
\begin{array}{ll}
Q \triangleq Q_{1}-Q_{12} Q_{2}^{-1} Q_{12}^{*}, & P \underline{\underline{\Delta}} P_{1}-P_{12} P_{2}^{-1} P_{12}^{*}, \\
Q^{\wedge} Q_{12} Q_{2}^{-1} Q_{12}^{*}, & \hat{P} \triangleq \underline{\perp} P_{12} P_{2}^{-1} P_{12}^{*}
\end{array}
$$

(4.33a, b)

Clearly, $\hat{Q}$ and $\hat{P}$ are nonnegative definite and have finite rank. Since by Lemma 4.2 $\tilde{Q}, \tilde{P} \in \mathscr{B}_{1}(\mathscr{H})$, it can be seen that $Q_{1}, P_{1} \in \mathscr{B}_{1}(\mathscr{H})$, which implies $Q, P \in \mathscr{B}_{1}(\mathscr{H})$. To show that $Q$ and $P$ are nonnegative definite, note that $Q$ is the $\mathscr{B}(\mathscr{H})$-component of the nonnegative-definite operator $\mathscr{Q} \bar{Q} \mathscr{Q}^{*} \in \mathscr{B}(\mathscr{\mathscr { H }})$, where

$$
2 \perp\left[\begin{array}{cc}
I_{x} & -Q_{12} Q_{2}^{-1} \\
0 & -I_{n_{r}}
\end{array}\right]
$$

Similarly, $P$ is nonnegative definite.
From the domain conditions (4.7) and (4.9) it follows that

$$
\begin{equation*}
Q_{1}: \mathscr{D}\left(A^{*}\right) \rightarrow \mathscr{S}(A), \tag{4.34a,b}
\end{equation*}
$$

$P_{1}: \mathscr{D}(A) \rightarrow \mathscr{D}\left(A^{*}\right)$,
(4.35a, b)
$Q_{12}: \mathbf{R}^{n_{c} \rightarrow \mathscr{D}(A), ~}$
$P_{12}: \mathbb{R}^{n} \rightarrow \mathscr{D}\left(A^{*}\right)$,
which lead to (3.12) and (3.13).
Next note that (4.27) is equivalent to (3.8) with

$$
\begin{equation*}
G \triangleq Q_{2}^{-1} Q_{12}^{*}, \quad \Gamma \triangleq-P_{2}^{-1} P_{12}^{*} \tag{4.36a,b}
\end{equation*}
$$

and that (3.7) holds with

$$
\begin{equation*}
M \stackrel{\Delta}{\triangle} Q_{2} P_{2} \tag{4.37}
\end{equation*}
$$

Since $Q_{2}$ and $P_{2}$ are positive definite, Lemma 2.6 implies that $M$ is positive semisimple. We can also define $\tau=G^{*} \Gamma$ which, by (3.8) satisfies $\tau^{2}=\tau$. It is helpful to note the identities

$$
\begin{equation*}
\hat{Q}=Q_{12} G=G^{*} Q_{12}^{*}, \quad \hat{P}=-P_{12} \Gamma=-\Gamma^{*} P_{12}^{*}, \tag{4.38a,b}
\end{equation*}
$$

$$
\begin{equation*}
\hat{Q}=G^{*} Q_{2} G, \quad \hat{P}=\Gamma^{*} P_{2} \Gamma, \tag{4.39a,b}
\end{equation*}
$$

(4.40a, b)

$$
\tau G^{*}=G^{*}, \quad \Gamma \tau=\Gamma
$$

$$
\begin{equation*}
\hat{Q}=\tau \hat{Q}, \quad \hat{P}=\hat{P}_{\tau}, \tag{4.41a,b}
\end{equation*}
$$

$$
\begin{equation*}
\hat{Q} \hat{P}=-Q_{12} P_{12}^{*} . \tag{4.42}
\end{equation*}
$$

From (3.8) and (2.1) it follows that
(4.43a, b)
(4.44a, b)

$$
\begin{aligned}
& \rho(G)=\rho(\Gamma)=n_{c} \\
& \rho\left(Q_{12}\right)=\rho\left(P_{12}\right)=n_{c}
\end{aligned}
$$

Hence, (2.2) and (4.38) imply $n_{c}=\rho\left(Q_{12}\right)+\rho(G)-n_{c} \leqq \rho(\hat{Q}) \leqq \rho\left(Q_{12}\right)=n_{o}$ which yields (3.14a). Similarly, (3.14b) holds and (3.14c) follows from (2.2) and (4.42).

Using (4.38) and (4.39), the components of $\tilde{Q}$ and $\tilde{P}$ can be written in terms of $G, \Gamma, Q, P, \hat{Q}$ and $\hat{P}$ as

$$
\begin{array}{ll}
Q_{1}=Q+\hat{Q}, & P_{1}=P+\hat{P}, \\
Q_{12}=\hat{Q} \Gamma^{*}, & P_{12}=-\hat{P} G^{*}, \\
Q_{2}=\Gamma \hat{Q} \Gamma^{*}, & P_{2}=G \hat{P} G^{*} . \tag{4.47}
\end{array}
$$

Now (3.10) and (3.11) can be obtained by substituting (4.45)-(4.47) into (4.30) and (4.31).

Expanding the $\mathscr{B}(\mathscr{H}), \mathscr{B}\left(\mathbb{R}^{n_{c}}, \mathscr{H}\right)$ and $\mathbb{R}^{n_{،} \times n_{c}}$ components of (4.8) and (4.10) yields

$$
\begin{align*}
& 0=A Q_{1}+Q_{1} A^{*}+B C_{c} Q_{12}^{*}+Q_{12}\left(B C_{c}\right)^{*}+V_{1}  \tag{4.48}\\
& 0=A Q_{12}+Q_{12} A_{c}^{T}+B C_{c} Q_{2}+Q_{1}\left(B_{c} C\right)^{*}  \tag{4.49}\\
& 0=A_{c} Q_{2}+Q_{2} A_{c}^{T}+B_{c} C Q_{12}+Q_{12}^{*}\left(B_{c} C\right)^{*}+B_{c} V_{2} B_{c}^{T}  \tag{4.50}\\
& 0=A^{*} P_{1}+P_{1} A+\left(B_{c} C\right)^{*} P_{12}^{*}+P_{12} B_{c} C+R_{1}  \tag{4.51}\\
& 0=P_{12} A_{c}+A^{*} P_{12}+\left(B_{c} C\right)^{*} P_{2}+P_{1} B C_{c}  \tag{4.52}\\
& 0=A_{c}^{T} P_{2}+P_{2} A_{c}+\left(B C_{c}\right)^{*} P_{12}+P_{12}^{*} B C_{c}+C_{c}^{T} R_{2} C_{c} \tag{4.53}
\end{align*}
$$

Substituting (4.45)-(4.47) into (4.48)-(4.53), using the identities

$$
\begin{aligned}
& B_{c} C=\Gamma Q \Sigma, \quad B C_{c}=-\Sigma P G^{*} \\
& B_{c} V_{2} B_{c}^{T}=\Gamma Q \bar{\Sigma} Q \Gamma^{*}, \quad C_{c}^{T} R_{2} C_{c}=G P \Sigma P G^{*}
\end{aligned}
$$

and defining

$$
A_{Q} \pm A-Q \bar{\Sigma}, \quad A_{P} \underline{\underline{\Delta}} A-\Sigma P
$$

we obtain

$$
\begin{align*}
& 0=A Q+Q A^{*}+A_{P} \hat{Q}+\hat{Q} A_{P}+V_{1}  \tag{4.54}\\
& 0=\left[A_{P} \hat{Q}+Q \bar{\Sigma} Q+\hat{Q}\left(\Gamma^{*} A_{c}^{T} G+\bar{\Sigma} Q\right)\right] \Gamma^{*},  \tag{4.55}\\
& 0=\Gamma\left[G^{*} A_{c} \Gamma \hat{Q}+Q \bar{\Sigma} \hat{Q}+Q \bar{\Sigma} Q+\hat{Q}\left(\Gamma^{*} A_{c}^{T} G+\bar{\Sigma} Q\right)\right] \Gamma^{*},  \tag{4.56}\\
& 0=A^{*} P+P A+A_{Q}^{*} \hat{P}+\hat{P} A_{Q}+R_{1}  \tag{4.57}\\
& 0=-\left[A_{Q}^{*} \hat{P}+P \Sigma P+\hat{P}\left(G^{*} A_{c} \Gamma+\Sigma P\right)\right] G^{*}  \tag{4.58}\\
& 0=G\left[\Gamma^{*} A_{c}^{T} G \hat{P}+P \Sigma \hat{P}+P \Sigma P+\hat{P}\left(G^{*} A_{c} \Gamma+\Sigma P\right)\right] G^{*} \tag{4.59}
\end{align*}
$$

We are now in a position to determine $A_{c}$ by computing (4.56) $-\Gamma(4.55)$ which yields (3.9). Alternatively, $A_{c}$ can be obtained by computing (4.59) $+G(4.58)$. As mentioned in $\S 3$, (3.9) is valid since $G^{*}: R^{n_{c}} \rightarrow \mathscr{\mathscr { F }}(\mathrm{~A})$ and $A_{c}^{T}$ is given by (3.26).

Next we substitute the expressions for $A_{c}$ and $A_{c}^{\top}$ into (4.55), (4.56), (4.58) and (4.59) and compute the relations (4.55) $G, G^{*}(4.56) G,-(4.58) \Gamma$ and $\Gamma^{*}(4.59) \Gamma$ to obtain, respectively,

$$
\begin{align*}
& 0=\left[A_{P} \hat{Q}+\hat{Q} A_{P}^{*}+Q \bar{\Sigma} Q\right]^{*},  \tag{4.60}\\
& 0=\tau\left[A_{P} \hat{Q}+\hat{Q} A_{P}^{*}+Q \bar{\Sigma} Q\right] \tau^{*},  \tag{4.61}\\
& 0=\left[A_{Q}^{*} \hat{P}+\hat{P} A_{Q}+P \Sigma P\right] \tau,  \tag{4.62}\\
& 0=\tau^{*}\left[A_{\hat{Q}}^{*} \hat{P}+\hat{P} A_{Q}+P \Sigma P\right] \tau . \tag{4.63}
\end{align*}
$$

Note that (4.60)-(4.63) are equivalent to (4.55), (4.56), (4.58) and (4.59) since $G$ and $\Gamma$ have full rank. Since $(4.61)=\tau(4.60)$ and $(4.63)=\tau^{*}(4.62),(4.61)$ and (4.63) are superfluous and can be omitted. Thus we have derived (3.17) and (3.18).

To obtain (3.15) and (3.16) we need only compute the relations (4.54) $+\tau(4.60)-$ (4.60) - (4.60)* and (4.57) $+\tau^{*}(4.62)-(4.62)-(4.62)^{*}$ and use (4.41).

Finally, to show that the preceding development entails no loss of generality in the optimality conditions we now use (3.9)-(3.18) to obtain (4.7)-(4.10) and (4.27)(4.29). Let $A_{n} B_{c}, C_{\Omega} G, \Gamma, \tau, Q, P, \hat{Q}, \hat{P}$ be as in the theorem statement and define $Q_{1}, Q_{12}, Q_{2}, P_{1}, P_{12}, P_{2}$ by (4.45)-(4.47). Note that (3.12) and (3.13) imply (4.34) and (4.35) and hence (4.7) and (4.9). Using (3.8), (3.10), (3.11) and (3.22) it is easy to verify (4.27)-(4.29). Finally, substitute (4.32), (4.33) and (4.36) into (3.15)-(3.18), reverse the steps taken earlier in the proof and use (3.9)-(3.11) to obtain (4.8) and (4.10), which completes the proof.
5. Concluding remarks. This paper has considered the problem of quadratically optimal, steady-state, fixed-order dynamic compensation for linear infinite-dimensional systems. The Main Theorem presents the stationarity conditions of the optimization problem in a highly simplified and rigorous form. The "optimal projection equations" (3.15)-(3.18) (or, equivalently, (3.27)-(3.30)) of the Main Theorem reveal the essential structure of the first-order necessary conditions and display the central role played by the optimal projection $\tau$. The relationship of the Main Theorem to the standard finite-dimensional steady-state $L Q G$ problem can be demonstrated by replacing $\tau$ with the identity matrix and noting that (3.27) and (3.28) reduce immediately to the familiar pair of operator Riccati equations and that (3.29) and (3.30) yield the controllability and observability gramians of the controller.

Inasmuch as the Main Theorem is a fundamental generalization of classical steady-state $L Q G$ theory, a number of issues must be reexamined. Hence, in conclusion we should like to point out some possible extensions of the Main Theorem along with directions for further research.

1. Sufficiency theory. Although sufficient conditions for the existence of an optimal compensator were not investigated in this paper, auxiliary conditions based upon the structure of (3.15)-(3.18) could perhaps be imposed upon $Q, P, \hat{Q}$ and $\hat{P}$ to single out the global optimum from amongst the local minima. This would be similar to the situation in $L Q G$ theory where, under stabilizability and detectability hypotheses, optimal stabilizing $Q$ and $P$ are identified as the unique nonnegative-definite solutions of the pair of algebraic Riccati equations.
2. Stabilizability. Just as in the full-order $L Q G$ problem, one would expect a natural relationship between the structure of the optimal solution and stabilizability/detectability hypotheses. The results of [41], [42] and [68] could serve as a starting point in this regard.
3. Numerical algorithms. In practical situations, the distributed parameter system would be replaced by a high-order discretized model for which the matrix version (rather than the operator version) of the optimal projection equations could be solved numerically. A numerical algorithm for solving the matrix version of the optimal projection equations has been developed in [32] and [34]. The proposed computational scheme is fundamentally quite different from gradient search algorithms [17], [18], [21], [22], [24], [25], [28], [30] in that it operates through direct solution of the optimal projection equations by iterative refinement of the optimal projection.
4. Convergence. One of the principal uses for the optimal projection equations will be to understand the relationship between fixed-order dynamic-compensator
designs which are optimal with respect to approximate models and the optimal fixed-order dynamic compensator for the distributed parameter system itself. By considering a sequence of $n$ th-order approximate models which converge to the distributed parameter system, conditions would be sought guaranteeing that the sequence of fixed-order compensators based on each approximate model approach the optimal dynamic compensator based upon the distributed parameter system (see [38]-[40]). This approach is analogous to the convergence results obtained in [7], [8] with the major difference being that the optimal projection equations permit the order of the compensator to remain fixed in accordance with real-world implementation constraints whereas in [7]-[9] the order of the compensator increases without bound.
5. Unbounded control and observation. An important generalization of the problem considered in this paper involves the case in which the input and output operators $B$ and $C$ are unbounded. The mathematical details for this problem are considerably more complex (see, e.g., [69]).
6. Singular observation noise/singular control weighting. As pointed out in [22], [33], [36] the assumptions of nonsingular control weighting and nonsingular observation noise preclude the use of direct output feedback as in

$$
\begin{equation*}
u(t)=C_{c} x_{c}(t)+D_{c} y(t) \tag{5.1}
\end{equation*}
$$

since $J$ is undefined unless

$$
\operatorname{tr}\left[D_{c}^{T} R_{2} D_{c} V_{2}\right]=0\left(\Leftrightarrow R_{2} D_{c} V_{2}=0\right)
$$

Although with due attention to (5.1) direct output feedback can be used in the singular case, the nature of the problem forebodes all of the difficulties associated with the singular LQG problem. Note that the deterministic output feedback problem [70], when viewed in this context, is highly singular.
7. Discrete-time system/discrete-time compensator. Digital implementation can be modelled by a discrete-time compensator with control of a continuous-time system facilitated by sampling and reconstruction devices. See [71], [73] for results in this direction.
8. Cross weighting/correlated disturbance and observation noise. This extension is straightforward and entirely analogous to the LQG case (see, e.g., [18, p. 351]).

Acknowledgments. We wish to thank Ardeth P. Grant for excellent and careful typing of the original manuscript version of this paper. Word processing and revision support was provided by Harris Corporation, GASD.

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# Finite-Dimensional Approximation for <br> Optimal Fixed-Order Compensation of <br> Distributed Parameter Systems 

by

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#### Abstract

In controlling distributed parameter systems it is often desirable to obtain low-order, finitedimensional controllers in order to minimize real-time computational requirements. Standard approaches to this problem employ model/controller reduction techniques in conjunction with LQG theory. In this paper we consider the finite-dimensional approximation of the infinite-dimensional Bernstein/Hyland optimal projection theory. Our approach yields fixed-finite-order controllers which are optimal with respect to high-order, approximating, finite-dimensional plant models. We illustrate the technique by computing a sequence of first-order controllers for one-dimensional, single-input/single-output, parabolic (heat/diffusion) and hereditary systems using spline-based, Ritz- Galerkin, finite element approximation. Our numerical studies indicate convergence of the feedback gains with less than $2 \%$ performance degradation over full-order LQG controllers for the parabolic system and $10 \%$ degradation for the hereditary system.


[^9]
## 1. Introduction

Approximation methods for the optimal control of distributed parameter systems have been widely studied. In particular, the approach taken in [1-12] involves approximating the original distributed parameter system by a sequence of finite-dimensional systems and then using finitedimensional control-design techniques to obtain a sequence of approximating, sub-optimal control laws, observers, or compensators. Furthermore, in these treatments it was demonstrated that if the open-loop system is approximated appropriately, then it is possible to guarantee convergence of the sequence of sub-optimal controllers, observers, or compensators, respectively, to the optimal controller, observer, or compensator for the original infinite-dimensional system. In addition, it can also be shown that when the approximating sub-optimal control laws or estimators are applied to the original system, near-optimal performance can frequently be obtained. These ideas were pursued in the context of both open- and closed-loop control, in both continuous and discrete-time, and for both full-state-feedback control and LQG (i.e., Kalman-filter-based) state estimation and compensation.

In practical situations, however, it is often of interest to obtain the simplest (i.e., the lowest order) controller which provides a given, desired feedback performance. This is usually achieved in one of two ways. Either the plant approximation order is reduced prior to controller design, or, alternatively, reduction techniques are applied to a given high-order control law. Unfortunately, the former approach may result in undesirable spillover effects while the latter may yield low-order controllers of low authority which perform unacceptably. In fact, with the second approach, this may occur even when a suitable controller is known to exist. For example, as is shown in [13], controller reduction techniques may even destabilize the closed-loop system.

A third, more direct approach involves fixing the controller order a priori, and then optimizing a performance criterion over the class of fixed-order controllers. In a finite-dimensional setting, a set of necessary conditions in the form of four coupled matrix equations (as a direct extension of the pair of the separated Riccati equations of LQG theory) which characterize the optimal fixedorder compensator was derived in [14]. These four equations are coupled via an oblique projection (idempotent) matrix. In the full-order case, this projection becomes the identity thus effectively eliminating the additional two equations, and the necessary conditions reduce to the standard LQG Riccati equations.

The notion that this direct (i.e., fixed-finite-order) approach can be applied to distributed parameter systems was first suggested by Johnson in [15] and further developed in [16] and [17]. To realize such an approach, however, would require a suitable generalization of the optimality conditions for the finite-dimensional fixed-order theory. This result was subsequently obtained in [18] where the matrix optimal projection equations obtained in [14] for finite-dimensional systems were extended to a set of four coupled operator Riccati and Lyapunov equations characterizing optimal fixed-finite-order controllers for infinite-dimensional systems.

In developing numerical schemes to actually compute fixed-finite-order compensators for infinite-dimensional systems, one might consider an approach wherein LQG reduction procedures are applied to a sequence of controllers obtained by using finite-dimensional full-order design techniques in conjunction with high-order finite-dimensional plant approximations. However, such an approach is unappealing for two reasons. First, since such methods are not predicated on the minimization of a performance index, prospects for convergence are slim. And, second, controllerreduction methods have not proven to be reliable in producing stabilizing compensators (see, for example, [13]).

Hence, on the other hand, we develop an abstract approximation framework (and ultimately computational schemes) which combine the infinite-dimensional optimal projection theory of [18] with the approximation ideas developed in [9-12] for infinite-dimensional LQG problems. More precisely, our approach involves constructing a sequence of approximating finite-dimensional subspaces of the original, underlying, infinite-dimensional Hilbert state space along with corresponding sequences of bounded linear operators which approximate the given input, output, and system operators. Then, by choosing bases for these approximating subspaces and applying the finitedimensional optimal projection theory, a sequence of matrix equations characterizing a sequence of approximating optimal, fixed-finite-order compensators for the distributed system is obtained. Finally, numerical techniques for solving the matrix optimal projection equations (for example, the homotopic continuation algorithm described in [19] and [20]) can be used to compute the sequence of approximating gains.

Our primary aim in this paper is to describe the general approach we are proposing, to discuss its implementation, and to demonstrate its feasibility and practicality. We offer no convergence arguments here, but rather reserve them for a more theoretical paper to follow. Instead, we consider the application of our technique to two examples. One involves the control of a one-dimensional,
single-input, single-output parabolic (heat/diffusion) system while the other involves a single-input single-output one-dimensional hereditary control system. These relatively simple examples have been used throughout the distributed parameter control literature to illustrate the application of new theories and techniques. A detailed discussion of the application of our ideas to more complex control systems, for instance, the vibration control of flexible structures, will also appear elsewhere. We use spline- based Ritz-Galerkin finite element schemes to approximate the open-loop systems (one for which convergence can be demonstrated in the LQG case) and present and discuss some of the numerical results which we have obtained using our general approximation framework.

We now outline the remainder of the paper. In Section 2 we briefly review the infinitedimensional optimal projection theory from [18], describe the approximation framework, and derive the corresponding equivalent matrix equations and feedback gains which characterize the approximating fixed-finite-order compensator. In Section 3 we consider the examples, construct the approximation schemes, and discuss our numerical findings. Section 4 contains a summary and some concluding remarks.

## 2. Optimal Projection Theory and Finite-Dimensional Approximation

We consider the following fixed-finite-order dynamic-compensation problem. Given the infinitedimensional control system

$$
\begin{equation*}
\dot{x}(t)=A x(t)+B u(t)+H_{1} w(t) \tag{2.1}
\end{equation*}
$$

with measurements

$$
\begin{equation*}
y(t)=C x(t)+H_{2} w(t) \tag{2.2}
\end{equation*}
$$

where $t \in[0, \infty)$, design a finite-dimensional, $n_{c}$ th-order dynamic compensator

$$
\begin{gather*}
\dot{x}_{c}(t)=A_{c} x_{c}(t)+B_{c} y(t),  \tag{2.3}\\
u(t)=C_{c} x_{c}(t) \tag{2.4}
\end{gather*}
$$

which minimizes the steady-state performance criterion

$$
\begin{equation*}
J\left(A_{c}, B_{c}, C_{c}\right) \triangleq \lim _{t \rightarrow \infty} \mathbb{E}\left[\left\langle R_{1} x(t), x(t)\right\rangle+u(t)^{\mathrm{T}} R_{2} u(t)\right] \tag{2.5}
\end{equation*}
$$

For convenience we denote the infinite-dimensional plant by $\Pi$; that is,

$$
\Pi \triangleq\left\{A, B, C, R_{1}, R_{2}, V_{1}, V_{2}\right\}
$$

Here $x(t)$ lies in a real, separable Hilbert space $\mathcal{X}$ with inner product $\langle\cdot, \cdot\rangle, A: \operatorname{Dom}(A) \subset X \rightarrow X$ is a closed, densely defined operator which generates a $C_{0}$ semigroup $\{T(t): t \geq 0\}$ of bounded linear operators on $\mathscr{X}, B \in \mathcal{L}\left(\mathbb{R}^{m}, X\right)$, and $C \in \mathcal{L}\left(X, \mathbb{R}^{\ell}\right)$. We assume that the state and measurement are corrupted by a white noise signal $w(t)$ in a real, separable Hilbert space $\hat{X}$ (see [21] or [22]), that $H_{1} \in \mathcal{L}(\hat{X}, X), H_{2} \in \mathcal{L}\left(\hat{X}, \mathbb{R}^{\ell}\right), R_{1} \in \mathcal{L}(X)$ is (self-adjoint) nonnegative definite, and that $R_{2}$ is an $m \times m$ (symmetric) positive-definite matrix. We define $V_{1}=H_{1} H_{i}^{*}$ and $V_{2}=H_{2} H_{2}^{*}$, where ( ) ${ }^{*}$ denotes adjoint, and assume for convenience that $H_{1} H_{2}^{*}=0$ and that $V_{2}$ is positive definite. The compensator is assumed to be of fixed, finite order $n_{c}$ (i.e., $x_{c}(t) \in \mathbb{R}^{n_{c}}$ ) and that $A_{c}, B_{c}$, and $C_{c}$ are matrices of appropriate dimension. For further details and discussion on the problem statement and the above assumptions, see [18].

We summarize here the primary result from [18] characterizing optimal fixed-finite-order controllers. For convenience define $\Sigma \triangleq B R_{2}^{-1} B^{*}$ and $\bar{\Sigma} \triangleq C^{*} V_{2}^{-1} C$. Also let $I_{n}$, and $I_{\Sigma}$ denote respectively the $n_{c} \times n_{c}$ identity matrix and the identity operator on $Y$.

Theorem 2.1. Let $n_{c}$ be given and suppose there exists a controllable and observable $n_{c}$ thorder dynamic compensator ( $A_{c}, B_{c}, C_{c}$ ) which minimizes $J$ given by (2.5) and for which the closedloop semigroup generated by

$$
A \triangleq\left[\begin{array}{cc}
A & B C_{c} \\
B_{c} C & A_{c}
\end{array}\right]
$$

is uniformly exponentially stable. Then there exist nonnegative-definite operators $Q, P, \hat{Q}, \hat{P}$ on $\mathcal{X}$ such that $A_{c}, B_{c}$, and $C_{c}$ are given by

$$
\begin{align*}
& A_{c}=\Gamma(A-Q \bar{\Sigma}-\Sigma P) G^{*}  \tag{2.6}\\
& B_{c}=\Gamma Q C^{*} V_{2}^{-1}  \tag{2.7}\\
& C_{c}=-R_{2}^{-1} B^{*} P G^{*} \tag{2.8}
\end{align*}
$$

where $Q: \operatorname{Dom}\left(A^{*}\right) \rightarrow \operatorname{Dom}(A), P: \operatorname{Dom}(A)^{\bullet} \rightarrow \operatorname{Dom}\left(A^{*}\right), \hat{Q}: \mathcal{X} \rightarrow \operatorname{Dom}(A), \hat{P}: \mathcal{X} \rightarrow$ $\operatorname{Dom}\left(A^{*}\right)$, and $G, \Gamma \in \mathcal{L}\left(X, \mathbb{R}^{n_{0}}\right)$, and such that the following conditions are satisfied:

$$
\begin{align*}
\operatorname{rank} \hat{Q} & =\operatorname{rank} \hat{P}=\operatorname{rank} \hat{Q} \hat{P}=n_{c} \\
\hat{Q} \hat{P} & =G^{*} M \Gamma, \quad \Gamma G^{*}=I_{n_{e}} \tag{2.10}
\end{align*}
$$

for some $M \in \mathbb{R}^{n_{0} \times n_{e}}$,

$$
\begin{align*}
& 0=A Q+Q A^{*}+V_{1}-Q \bar{\Sigma} Q+\tau_{\perp} Q \Sigma Q \tau_{\perp}^{*},  \tag{2.11}\\
& 0=A^{*} P+P A+R_{1}-P \Sigma P+\tau_{\perp}^{*} P \Sigma P r_{\perp},  \tag{2.12}\\
& 0=(A-\Sigma P) \hat{Q}+\hat{Q}(A-\Sigma P)^{*}+Q \Sigma Q-\tau_{\perp} Q \Sigma Q r_{\perp}^{*},  \tag{2.13}\\
& 0=(A-Q \bar{\Sigma})^{*} \hat{P}+\hat{P}(A-Q \bar{\Sigma})+P \Sigma P-\tau_{\perp}^{*} P \Sigma P r_{\perp}, \tag{2.14}
\end{align*}
$$

where

$$
\tau \triangleq G^{*} \Gamma, \quad \tau_{\perp} \triangleq I_{x}-\tau
$$

It is shown in [18] that the factorization (2.10) for the nonnegative-definite operators $\hat{Q}$ and $\hat{P}$ satisfying rank $\hat{Q} \hat{P}=n_{c}$ always exists and is unique except for a change of basis in $\mathbb{R}^{n_{0}}$. It is also shown in [18] that $\Gamma^{*}: \mathbb{R}^{n_{0}} \rightarrow \operatorname{Dom}\left(A^{*}\right)$ so that the expression (2.6) is well defined.

Equations (2.11)-(2.14) are, in general, infinite-dimensional operator equations. To actually use them to compute the optimal fixed-order compensator, a finite-dimensional plant approximation is required. For each $N=1,2, \ldots$, let $\mathcal{X}^{N}$ denote a finite-dimensional subspace of $X$ and let $\rho^{N}$ : $\mathcal{X} \rightarrow \mathcal{X}^{N}$ denote the corresponding orthogonal projection of $\mathcal{X}$ onto $X^{N}$. Let $A^{N} \in \mathcal{L}\left(X^{N}\right), B^{N} \in$ $\mathcal{L}\left(\mathbb{R}^{m}, \chi^{N}\right), C^{N} \in \mathcal{L}\left(X^{N}, \mathbb{R}^{\ell}\right), R_{\mathrm{i}}^{N} \in \mathcal{L}\left(X^{N}\right)$, and $V_{1}^{N} \in \mathcal{L}\left(Y^{N}\right)$. We consider the system (2.6)(2.14) with the plant $\Pi$ replaced by the plant $\Pi^{N}$ given by

$$
\Pi^{N} \triangleq\left\{A^{N}, B^{N}, C^{N}, R_{1}^{N}, R_{2}, V_{1}^{N}, V_{2}\right\} .
$$

Typically, the operators $B^{N}, C^{N}, R_{1}^{N}$ and $V_{1}^{N}$ are chosen as $B^{N}=\rho^{N} B, C^{N}=C \rho^{N}, R_{1}^{N}=\rho^{N} R_{1}$ and $V_{1}^{N}=P^{N} V_{1}$ with the requirement that $P^{N}$ converge strongly to the identity $I_{X}$ as $N \rightarrow \infty$. The operator $A^{N}$ is chosen so that it and its adjoint satisfy the hypotheses of the Trotter-Kato semigroup approximation theorem (i.e., stability and consistency, see, for example, [23]). That is, $A^{N}$ is chosen so that $\lim _{N \rightarrow \infty} T^{N}(t) P^{N} \phi=T(t) \phi$, and $\lim _{N \rightarrow \infty} T^{N}(t)^{*} P^{N} \phi=T(t)^{*} \phi$, uniformly in $t$ for $t$ in bounded intervals, for each $\phi \in X$, where $T^{N}(t)=\exp \left(t A^{N}\right), t \geq 0$. We shall say more about these choices for $A^{N}, B^{N}, C^{N}, R_{1}^{N}$, and $V_{1}^{N}$ when we address convergence questions below.

Although with the plant $\Pi^{N}$ equations (2.11)-(2.14) are finite dimensional, they are still operator equations. It is their matrix equivalents which are used in computations. Unless orthonormal bases are chosen for the subspaces $X^{N}$ (which is typically not the case in practice) some care must be taken to obtain the appropriate matrix system.

For each $N=1,2, \ldots$, let $\left\{\phi_{j}^{N}\right\}_{j=1}^{k_{j}^{N}}$ be a basis for $X^{N}$ and choose the standard bases for all Euclidean spaces. For a linear operator $L$ with domain and range $X^{N}$ or any Euclidean space, let $[L]$ denote its matrix representation with respect to the bases chosen above. Also, let $\Phi^{N}$ denote the $k^{N}$-square Gram matrix corresponding to the basis $\left\{\phi_{j}^{N}\right\}_{j=1}^{k_{j}^{N}}$; that is, $\Phi_{i j}^{N}=\left\langle\phi_{i}^{N}, \phi_{j}^{N}\right\rangle$, $i, j=1,2, \ldots, k^{N}$. Noting that

$$
\begin{array}{lll}
{\left[\left(A^{N}\right)^{*}\right]=\left(\Phi^{N}\right)^{-1}\left[A^{N}\right]^{\mathrm{T}} \Phi^{N},} & {\left[\left(B^{N}\right)^{*}\right]=\left[B^{N}\right]^{\mathrm{T}} \Phi^{N},} & {\left[\left(C^{N}\right)\right]=\left(\Phi^{N}\right)^{-1}\left[C^{N}\right]^{\mathrm{T}}} \\
{\left[\left(\tau_{\perp}^{N}\right)^{*}\right]=\left(\Phi^{N}\right)^{-1}\left[\tau_{\perp}^{N}\right]^{\mathrm{T}} \Phi^{N},} & {\left[\Sigma^{N}\right]=\left[B^{N}\right] R_{2}^{-1}\left[B^{N}\right]^{\mathrm{T}} \Phi^{N},} & {\left[\Sigma^{N}\right]=\left(\Phi^{N}\right)^{-1}\left[C^{N}\right]^{\mathrm{T}} V_{2}^{-1}\left[C^{N}\right]}
\end{array}
$$

the matrix equivalents of the operator equations (2.11)-(2.14) become

$$
\begin{align*}
0= & {\left[A^{N}\right]\left[Q^{N}\right]+\left[Q^{N}\right]\left(\Phi^{N}\right)^{-1}\left[A^{N}\right]^{T} \Phi^{N}+\left[V_{1}^{N}\right]-\left[Q^{N}\right]\left[\bar{\Sigma}^{N}\right]\left[Q^{N}\right] } \\
& +\left[\tau_{\perp}^{N}\right]\left[Q^{N}\right]\left[\bar{\Sigma}^{N}\right]\left[Q^{N}\right]-\left(\Phi^{N}\right)^{-1}\left[\tau_{\perp}^{N}\right]^{T} \Phi^{N},  \tag{2.15}\\
0= & \left(\Phi^{N}\right)^{-1}\left[A^{N}\right]^{T} \Phi^{N}\left[P^{N}\right]+\left[P^{N}\right]-\left[A^{N}\right]+\left[R_{1}^{N}\right]-\left[P^{N}\right]\left[\Sigma^{N}\right]\left[P^{N}\right] \\
& +\left(\Phi^{N}\right)^{-1}\left[\tau_{\perp}^{N}\right]^{T} \Phi^{N}\left[P^{N}\right]-\left[\Sigma^{N}\right]\left[P^{N}\right]\left[\tau_{\perp}^{N}\right],  \tag{2.16}\\
0= & \left(\left[A^{N}\right]-\left[\Sigma^{N}\right]\left[P^{N}\right]\right)\left[\hat{Q}^{N}\right]+\left[\hat{Q}^{N}\right]\left(\Phi^{N}\right)^{-1}\left(\left[A^{N}\right]-\left[\Sigma^{N}\right]\left[P^{N}\right]\right)^{T} \Phi^{N} \\
& +\left[Q^{N}\right]\left[\bar{\Sigma}^{N}\right]\left[Q^{N}\right]-\left[\tau_{\perp}^{N}\right]\left[Q^{N}\right]\left[\Sigma^{N}\right]\left[Q^{N}\right]\left(\Phi^{N}\right)^{-1}-\left[\tau_{\perp}^{N}\right]^{\mathrm{T}} \Phi^{N},  \tag{2.17}\\
0= & \left(\Phi^{N}\right)^{-1}\left(\left[A^{N}\right]-\left[Q^{N}\right]\left[\Sigma^{N}\right]\right)^{T^{N}}\left[\hat{P}^{N}\right]+\left[\hat{P}^{N}\right]\left(\left[A^{N}\right]-\left[Q^{N}\right]\left[\Sigma^{N}\right]\right) \\
& +\left[P^{N}\right]\left[\Sigma^{N}\right]\left[P^{N}\right]-\left(\Phi^{N}\right)^{-1}\left[\tau_{\perp}^{N}\right]^{T} \Phi^{N}\left[P^{N}\right]\left[\Sigma^{N}\right]\left[P^{N}\right]\left[\tau_{\perp}^{N}\right] . \tag{2.18}
\end{align*}
$$

Therefore, if we define the $k^{N} \times k^{N}$ nonnegative-definite matrices

$$
\begin{array}{ll}
Q_{0}^{N} \triangleq\left[Q^{N}\right]\left(\Phi^{N}\right)^{-1}, & P_{0}^{N} \triangleq \Phi^{N}\left[P^{N}\right] \\
\hat{Q}_{0}^{N} \triangleq\left[\hat{Q}^{N}\right]\left(\Phi^{N}\right)^{-1}, & \hat{P}_{0}^{N} \triangleq \Phi^{N}\left[\hat{P}^{N}\right] \\
V_{0}^{N} \triangleq\left[V_{1}^{N}\right]\left(\Phi^{N}\right)^{-1}, & R_{0}^{N} \triangleq \Phi^{N}\left[R_{1}^{N}\right] \\
\Sigma_{0}^{N} \triangleq\left[B^{N}\right] R_{3}^{-1}\left[B^{N}\right]^{\mathbf{T}}, & \bar{\Sigma}_{0}^{N} \triangleq\left[C^{N}\right]^{\mathbf{T}} V_{2}^{-1}\left[C^{N}\right]
\end{array}
$$

we can solve the matrix optimal projection equations given in [14] corresponding to the matrix plant model

$$
\Pi_{0}^{N} \triangleq\left\{\left[A^{N}\right],\left[B^{N}\right],\left[C^{N}\right], R_{0}^{N}, R_{2}, V_{0}^{N}, V_{2}\right\}
$$

to obtain the matrices $Q_{0}^{N}, P_{0}^{N}, \hat{Q}_{0}^{N}$ and $\hat{P}_{0}^{N}$. The approximating optimal $n_{c}$ th-order dynamic compensator $\left\{A_{c}^{N}, B_{c}^{N}, C_{c}^{N}\right\}$ is then given by

$$
\begin{aligned}
& A_{c}^{N}=\Gamma_{0}^{N}\left(\left[A^{N}\right]-Q_{0}^{N} \Sigma_{0}^{N}-\Sigma_{0}^{N} P_{0}^{N}\right)\left(G_{0}^{N}\right)^{\mathrm{T}} \\
& B_{c}^{N}=\Gamma_{0}^{N} Q_{0}^{N}\left[C^{N}\right]^{\mathrm{T}} V_{2}^{-1} \\
& C_{c}^{N}=-R_{2}^{-1}\left[B^{N}\right]^{\mathrm{T}} P_{0}^{N}\left(G_{0}^{N}\right)^{\mathrm{T}}
\end{aligned}
$$

where $\Gamma_{0}^{N}, G_{0}^{N} \in \mathbb{R}^{n_{\bullet} \times k^{N}}$, and $M^{N} \in \mathbb{R}^{n_{0} \times n_{0}}$ satisfy

$$
\begin{gathered}
\hat{Q}_{0}^{N} \hat{P}_{0}^{N}=G_{0}^{N} M^{N} \Gamma_{0}^{N}, \quad \Gamma_{0}^{N}\left(G_{0}^{N}\right)^{\mathrm{T}}=I_{n_{a}} \\
{\left[\tau^{N}\right]=\left(G_{0}^{N}\right)^{\mathrm{T}} \Gamma_{0}^{N}, \quad\left[\tau_{\perp}^{N}\right]=I_{k^{N}}-\left[\tau^{N}\right] .}
\end{gathered}
$$

When an infinite-dimensional controller will suffice, $C_{c}=-R_{2}^{-1} B^{*} P \in \mathcal{L}\left(X, \mathbb{R}^{m}\right)$ and $B_{c}=$ $Q C^{*} V_{2}^{-1} \in \mathcal{L}\left(R^{\ell}, X\right)$ are the usual infinite-dimensional LQG controller and observer gains (see [9]). The operators $P, Q \in \mathcal{L}(X)$ are the nonnegative-definite solutions to the two decoupled operator algebraic Riccati equations (2.11) and (2.12) with $\tau$ and $\tau_{\perp}$ formally set to $I_{\Upsilon}$ and 0 , respectively. Since $C_{c}$ has range in $\mathbb{R}^{m}$ and $B_{c}$ has domain $\mathbb{R}^{\ell}$, there exist vectors $c_{c}=\left(c_{c}^{1}, \ldots, c_{c}^{m}\right)^{\mathbf{T}} \in \times_{j=1}^{m} X$ and $b_{c}=\left(b_{c}^{1}, \ldots, b_{c}^{l}\right)^{\mathrm{T}} \in x_{j=1}^{l} X$ such that

$$
\left[C_{c} x\right]_{i}=\left\langle c_{c}^{i}, x\right\rangle, \quad i=1,2, \ldots, m, \quad x \in X
$$

and

$$
B_{c} y=b_{c}^{\mathrm{T}} y=\sum_{i=1}^{\ell} y_{i} b_{c}^{i}, \quad y=\left(y_{1}, \ldots, y_{\ell}\right) \in \mathbb{R}^{\ell}
$$

The vectors $c_{c}$ and $b_{c}$ are referred to as the optimal LQG controller and observer functional gains respectively.

With regard to approximation for the full-order LQG problem, for each $N=1,2, \ldots$ we take $n_{c}=k^{N}$. Then it is not difficult to show that

$$
C_{c}^{N}\left[P^{N} x\right]=\left\langle c_{c}^{N}, x\right\rangle, \quad x \in X,
$$

and

$$
B_{c}^{N} y=\left(b_{c}^{N}\right)^{\mathbf{T}} y, \quad y \in \mathbb{R}^{\ell}
$$

where $c_{c}^{N} \in \times_{j=1}^{m} \chi^{N}$ and $b_{c}^{N} \in x_{j=1}^{\ell} \chi^{N}$ are given by $c_{c}^{N}=C_{c}^{N}\left(\Phi^{N}\right)^{-1} \phi^{N}$ and $b_{c}^{N}=\left(B_{c}^{N}\right)^{\mathrm{T}} \phi^{N}$ respectively with $\phi^{N}=\left(\phi_{1}^{N}, \ldots, \phi_{k^{N}}^{N}\right) \in \times_{j=1}^{k^{N}} \chi^{N}$. The vectors $c_{c}^{N}$ and $b_{c}^{N}$ are referred to as the approximating optimal LQG controller and observer functional gains. To compute them we need only solve two standard decoupled matrix algebraic Riccati equations for the $k^{N} \times k^{N}$ nonnegativedefinite matrices $Q_{0}^{N}$ and $P_{0}^{N}$.

A rather complete convergence theory for LQG approximation can be found in [9]. Essentially, it is shown there that if the approximating subspaces ${X^{N}}^{N}$ are chosen so that the projections $P^{N}$
converge strongly to the identity as $N \rightarrow \infty$, the operators $A^{N}, B^{N}, C^{N}, R_{1}^{N}$, and $V_{1}^{N}$ are chosen as was described above, and the operators $Q^{N}$ and $P^{N}$ are uniformly bounded in $N$, then $Q^{N}$ and $P^{N}$ converge weakly to $Q$ and $P$, respectively as $N \rightarrow \infty$. This in turn implies that $C_{c}^{N} \rightarrow C_{c}$, strongly, $B_{c}^{N} \rightarrow B_{c}$, weakly, $c_{c}^{N} \rightarrow c_{c}$ and $b_{c}^{N} \rightarrow b_{c}$, weakly, and the closed-loop semigroup for the approximating optimal LQG compensator converges weakly to the closed-loop semigroup for the optimal infinite-dimensional LQG compensator, as $N \rightarrow \infty$. If, in addition, the operators $S^{N}(t)=T^{N}(t)+B^{N} C_{c}^{N}$ and $\hat{S}^{N}(t)=T^{N}(t)-B_{c}^{N} C^{N}$ are uniformly exponentially stable, uniformly in $N$, then $Q^{N} \rightarrow Q$ and $P^{N} \rightarrow P$, strongly, $C_{c}^{N} \rightarrow C_{c}$ and $B_{c}^{N} \rightarrow B_{c}$, in norm, $c_{c}^{N} \rightarrow c_{c}$ and $b_{c}^{N} \rightarrow b_{c}$, strongly, and the closed-loop semigroups converge strongly, as $N \rightarrow \infty$. If $R_{1}^{N}$ and $V_{1}^{N}$ are coercive, uniformly in $N$, then $S^{N}(t)$ and $\hat{S}^{N}(t)$ will be uniformly exponentially stable. If it is also true that $R_{1}$ and $V_{1}$ are trace class and $R_{1}^{N} P^{N} \rightarrow R_{1}$ and $V_{1}^{N} P^{N} \rightarrow V_{1}$ in trace norm then $Q$ and $P$ are trace class and $Q^{N} P^{N} \rightarrow Q$ and $P^{N} P^{N} \rightarrow P$ in trace norm as $N \rightarrow \infty$.

Returning to the fixed-finite-order case, we note that in general the approximating optimal projection equations may not possess a unique solution. However, in [19] it is shown for the finitedimensional case that it is possible to obtain an upper bound for the number of stabilizing solutions. Using topological degree theory, the following result was obtained in [19].

Theorem 2.2. Consider the equations (2.11)-(2.14) with the infinite-dimensional plant $\Pi$ replaced by the finite-dimensional plant $\Pi^{N}$. Let $n_{u}$ denote the dimension of the unstable subspace of $A^{N}$ and assume that $n_{c} \geq n_{u}$. Then in the class of nonnegative-definite operators $Q^{N}, P^{N}, \hat{Q}^{N}, \hat{P}^{N}$ on $\mathcal{X}^{N}$ satisfying rank $\hat{Q}^{N}=\operatorname{rank} \hat{P}^{N}=\operatorname{rank} \hat{Q}^{N} \hat{P}^{N}=n_{c}$, there exist at most

$$
\left(\begin{array}{cl}
\min \left(k^{N}, m, \ell\right)-n_{u} \\
n_{c}-n_{u} & n_{c} \leq \min \left(k^{N}, m, \ell\right) \\
1, & \text { otherwise }
\end{array}\right.
$$

solutions of (2.11)-(2.14), each of which is stabilizing. If, in addition, the plant ( $A^{N}, B^{N}, C^{N}$ ) is stabilizable by an $n_{c}{ }^{+}+$-order $\sim$ ratroller, then there exists at least one stabilizing solution of (2.9)-(2.14).

Theorem 2.2 shows that while there may exist multiple solutions to the finite-dimensional optimal projection equations, in practice this number can be quite small. For example, if $n_{c} \geq n_{u}$ and the system is either single input ( $m=1$ ) or single output ( $\ell=1$ ) then there exists at most one solution to (2.9)-(2.14) for the plant $\Pi^{N}$. The existence of at least one stabilizing solution of course depends upon whet, $r$ or not the plant is stabilizable by an $n_{c}$ th-order controller
(for relevant results, see [24]). Finally, while it may be possible to stabilize a plant with $n_{c}<n_{u}$, this case lies outside the scope of the analysis given in [19].

## 3. Examples and Numerical Results

We first consider the one-dimensional, single-input/single-output, parabolic (heat/diffusion) control system with Dirichlet boundary conditions given by

$$
\begin{gather*}
\frac{\partial v}{\partial t}(t, \eta)=a \frac{\partial^{2} v}{\partial \eta^{2}}(t, \eta)+b(\eta) u(t)+h_{1}(\eta) w_{1}(t, \eta), \quad 0<\eta<1, \quad t>0  \tag{3.1}\\
v(t, 0)=0, \quad v(t, 1)=0, \quad t>0  \tag{3.2}\\
y(t)=\int_{0}^{1} c(\eta) v(t, \eta) d \eta+h_{2} w_{2}(t), \quad t>0 \tag{3.3}
\end{gather*}
$$

where $a>0$, and $b(\cdot)$ and $c(\cdot)$ are given by

$$
b(\eta)= \begin{cases}\frac{1}{\beta_{2}-\beta_{1}}, & \beta_{1} \leq \eta \leq \beta_{2} \\ 0, & \text { elsewhere }\end{cases}
$$

and

$$
c(\eta)= \begin{cases}\frac{1}{\gamma_{2}-\gamma_{1}}, & \gamma_{1} \leq \eta \leq \gamma_{2} \\ 0, & \text { elsewhere }\end{cases}
$$

with $0 \leq \beta_{1}<\beta_{2} \leq 1$ and $0 \leq \gamma_{1}<\gamma_{2} \leq 1$. In (3.1) and (3.3), $h(\cdot) \in L_{\infty}(0,1), w_{1}(t, \cdot) \in$ $L_{2}(0,1)$, a.a. $t \in[0, \infty)$, (see $[22]$, p. 314), $h_{2}$ is a nonzero constant and $w_{2}(\cdot)$ is unit-intensity white noise.

To rewrite (3.1)-(3.3) in the form (2.1), (2.2), in the usual way we take $X=L_{2}(0,1)$ endowed with the standard $L_{2}$ inner product, let $x(t)=v(t, \cdot), t \geq 0$, define $A: \operatorname{Dom}(A) \subset \mathcal{X} \rightarrow \mathcal{X}$ by $A \phi=a D^{2} \phi$ for $\phi \in \operatorname{Dom} A \triangleq H^{2}(0,1) \cap H_{0}^{1}(0,1)$, and define $B \in \mathcal{L}\left(\mathbb{R}^{1}, X\right)$ and $C \in \mathcal{L}\left(X, \mathbb{R}^{1}\right)$ by $B u=b(\cdot) u$ for $u \in \mathbb{R}^{1}$, and $C \phi=\int_{0}^{1} c(\eta) \phi(\eta) d \eta$, for $\phi \in L_{2}(0,1)$, respectively. Furthermore, let $\hat{X} \triangleq L_{2}(0,1) \times \mathbb{R}$, set $w(t) \triangleq\left(w_{1}(t, \cdot), w_{2}(t)\right) \in \hat{X}$, and define $H_{1} \in \mathcal{L}(\hat{X}, \underline{X})$ and $H_{2} \in \mathcal{L}\left(\hat{X}, \mathbb{R}^{1}\right)$ by $H_{1} z=h_{1}(\cdot) z_{1}$ and $H_{2} z=h_{2} z_{2}$ for $z=\left(z_{1}, z_{2}\right) \in \hat{X}$.

It is well known (see, for example, [23]) that $A$ is closed, densely defined, and negative definite. Furthermore, $A$ is the infinitesimal generator of a uniformly exponentially stable, analytic (abstract parabolic) semigroup $\{T(t): t \geq 0\}$ of bounded, self-adjoint linear operators on $\mathcal{X}$.

We consider linear spline-based Ritz-Galerkin approximation for the open-loop system. For each $N=2,3, \ldots$, let $\left\{\phi_{j}^{N}\right\}_{j=1}^{N-1}$ be the linear spline ("hat") functions defined on the interval $[0,1]$
with respect to the uniform partition $\left\{0, \frac{1}{N}, \frac{2}{N}, \ldots, 1\right\}$, i.e.,

$$
\phi_{j}^{N}(\eta)= \begin{cases}N \eta-j+1, & \eta \in\left[\frac{j-1}{N}, \frac{j}{N}\right), \\ j+1-N \eta, & \eta \in\left[\frac{i}{N}, \frac{i+1}{N}\right), \\ 0, & \text { elsewhere on }[0,1],\end{cases}
$$

$j=1,2, \ldots, N-1$. Set $X^{N}=\operatorname{span}\left\{\phi_{j}^{N}\right\}_{j=1}^{N-1}$ and note that $k^{N}=\operatorname{dim} X^{N}=N-1$, and $\mathcal{X}^{N} \subset H_{0}^{1}(0,1)$ for all $N$. If $\mathcal{P}^{N}: X \rightarrow X^{N}$ denotes the orthogonal projection of $\mathscr{X}$ onto $\mathcal{X}^{N}$, then standard convergence estimates for interpolatory splines (see [25]) can be used to show that $\lim _{N \rightarrow \infty} \rho^{N} \phi=\phi$ in $L_{2}(0,1)$ for $\phi \in L_{2}(0,1)$.

There are two equivalent ways to obtain an operator representation for the usual Ritz-Galerkin approximation to $A$. First, $A$ can be extended to a bounded linear operator from $H_{0}^{1}(0,1)$ onto its dual, $H^{-1}(0,1)$, via

$$
\begin{equation*}
(A \phi)(\psi)=-a\langle D \phi, D \psi\rangle, \quad \phi, \psi \in H_{0}^{1}(0,1) . \tag{3.4}
\end{equation*}
$$

Since $\mathcal{X}^{N} \subset H_{0}^{1}(0,1)$ for all $N=2,3, \ldots$, we define $A^{N} \in \mathcal{L}\left(X^{N}\right)$ by $A^{N} \phi^{N}=A \phi^{N}, \phi^{N} \in X^{N}$, with $A \phi^{N} \in H^{-1}(0,1)$ considered to be a linear functional on $\mathcal{X}^{N}$. From the Riesz Representation theorem we obtain $A^{N} \phi^{N}=\psi^{N}$ where $\psi^{N}$ is that element in $\chi^{N}$ which satisfies $\left(A^{N} \phi^{N}\right)\left(\chi^{N}\right)=$ $-a\left\langle D \phi^{N}, D \chi^{N}\right\rangle=\left\langle\psi^{N}, \chi^{N}\right\rangle$.

Alternatively and equivalently, by using the fact that $A$ is self-adjoint, we can define $A^{N}$ as follows. Let $\mathcal{P}_{1}^{N}: H_{0}^{1}(0,1) \rightarrow \mathcal{X}^{N}$ denote the orthogonal projection of the Hilbert space $H_{0}^{1}(0,1)$ onto $\mathscr{X}^{N}$. Using the definition (3.4), it is not difficult to show that $-A \in \mathcal{L}\left(H_{0}^{1}(0,1), H^{-1}(0,1)\right)$ is coercive and, therefore, that $A^{-1}: H^{-1}(0,1) \rightarrow H_{0}^{1}(0,1)$ exists and is bounded. We then define $A^{N} \in \mathcal{L}\left(X^{N}\right)$ to be the inverse of the operator given by $\left(A^{N}\right)^{-1}=\left.P_{1}^{N} A^{-1}\right|_{I^{N}}$.

Using either definition, it is easily argued that $A^{N}$ is well defined, self-adjoint, and is the infinitesimal generator of a uniformly exponentially stable (uniformly in $N$ ) semigroup, $T^{N}(t)=$ $\exp \left(t A^{N}\right), t \geq 0$, of bounded linear operators on $X^{N}$. Also, using the approximation properties of splines, it is not difficult to show that $\lim _{N \rightarrow \infty}\left(A^{N}\right)^{-1} P^{N} \phi=A^{-1} \phi, \phi \in X$. Consequently, the hypotheses of the Trotter-Kato theorem (see [23]) are satisfied and we have $\lim _{N \rightarrow \infty} T^{N}(t) P^{N} \phi=$ $T(t) \phi$ and $\lim _{N \rightarrow \infty} T^{N}(t)^{*} p^{N} \phi=T(t)^{*} \phi, \phi \in X$, uniformly in $t$ for $t$ in bounded intervals. A detailed discussion of the results just outlined can be found in [8].

We define $B^{N}=P^{N} B$ and $C^{N}=C P^{N}$, from which it immediately follows that $\lim _{N \rightarrow \infty} B^{N}=$ $B$ and $\lim _{N \rightarrow \infty} C^{N}=C$ in norm and similarly for their adjoints. For the example we shall consider
here, we have chosen $R_{1}=r_{1} I_{x}, R_{2}=r_{2} I_{m}$, with $r_{1}, r_{2}>0$. Setting $h_{1}(\eta)=v_{1}^{\frac{1}{2}}, 0<\eta<1$, and $h_{2}=v_{2}^{\frac{1}{2}}$ with $v_{1}, v_{2}>0$, we obtain $V_{1}=v_{1} I_{X}$ and $V_{2}=v_{2}$. We then take $R_{1}^{N}=\rho^{N} R_{1}$ and $V_{1}^{N}=\rho^{N} V_{1}$. For the LQG problem, the open-loop uniform exponential stability of both the infinite-dimensional system and the approximating systems is sufficient to conclude the strong convergence of the approximating Riccati operators to the solutions of the infinite-dimensional Riccati equations, the uniform norm convergence of the approximating controller and observer gains, and the strong convergence of the functional gains, as $N \rightarrow \infty$.

Since the basis elements $\left\{\phi_{j}^{N}\right\}_{j=1}^{N-1}$ are piecewise linear with respect to the uniform mesh $\left\{0, \frac{1}{N}, \frac{2}{N}, \ldots, 1\right\}$ on $[0,1]$, the equivalent matrix representations for the operators defined above can be computed directly and in closed form. The Gram matrix $\Phi_{i j}^{N}=\left\langle\phi_{i}^{N}, \phi_{j}^{N}\right\rangle, i, j=1,2, \ldots, N-1$ is given by $\Phi^{N}=\frac{1}{N} \operatorname{Tridiag}\left\{\frac{1}{6}, \frac{2}{3}, \frac{1}{6}\right\}$, and if we define the generalized atifness matrix $\Psi^{N}$ by $\Psi_{i j}^{N}=-a\left\langle D \phi_{i}^{N}, D \phi_{j}^{N}\right\rangle, i, j=1,2, \ldots, N-1$, then $\Psi^{N}=a N$ Tridiag $\{1,-2,1\}$. It follows that $\left[A^{N}\right]=\left(\Phi^{N}\right)^{-1} \Psi^{N},\left[B^{N}\right]=\left(\Phi^{N}\right)^{-1} b^{N},\left[C^{N}\right]=c^{N}$, with $b_{i}^{N}=\left\langle b, \phi_{i}^{N}\right\rangle=\frac{1}{\beta_{2}-\beta_{1}} \int_{\beta_{1}}^{\beta_{2}} \phi_{i}^{N}(\eta) d \eta$, and $c_{i}^{N}=\left\langle c, \phi_{i}^{N}\right\rangle=\frac{1}{\gamma_{2}-\gamma_{1}} \int_{\gamma_{2}}^{\gamma_{2}} \phi_{i}^{N}(\eta) d \eta, i=1,2, \ldots, N-1$, and that $R_{0}^{N}=r_{1} \Phi^{N}$ and $V_{0}^{N}=v_{1}\left(\Phi^{N}\right)^{-1}$.

For our numerical study we set $a=1, \beta_{1}=.75-.03 \sqrt{2}, \beta_{2}=75+.04 \sqrt{2}, \gamma_{1}=.25-.04 \sqrt{2}, \gamma_{2}=$ $.25+.03 \sqrt{2}, r_{1}=v_{1}=1, r_{2}=v_{2}=10^{-4}, h_{1}(\eta) \equiv 1$, and used our technique to compute approximating optimal LQG (i.e., $n_{c}=N-1$ ) and 1 st order (i.e., $n_{c}=1$ ) compensators for various values of $N$. The open-loop stability of system (3.1)-(3.3) and the approximating systems imply that the finite-dimensional approximating optimal projection equations have a solution. Theorem 2.2 on the other hand, with $n_{u}=0$ and $n_{c}=1$ or $n_{c}=N-1$, implies that they have at most one solution. Consequently, the system of equations (2.11)-(2.14) with the plants $\Pi^{N}$ admits a unique solution.

The optimal projection equations (2.11)-(2.14) were solved using the homotopic continuation algorithm described in [19]. It is shown in [19] that the operation count for the algorithm is proportional to $p\left(2 n^{3}+(m+\ell) n^{2}+(m+\ell)^{3} n_{c}^{3}\right)$ where $p$ is the number of integration steps and $n$ is the dimension of the finite-dimensional plant. This is competitive with the operation count for the Hamiltonian solution of the standard Riccati equations which is $O\left(16 n^{3}\right)$ for LQG. Also, note that the computational burden for the solution of the optimal projection equations decreases with $n_{c}$.

Since $m=\ell=1$ in the LQG case, the optimal functional observer and feedback control gains $b_{c}$ and $c_{c}$ and the approximating gains $b_{c}^{N}$ and $c_{c}^{N}$, are all simply $L_{2}$ functions with $b_{c}^{N}$ and $c_{c}^{N}$ elements in $X^{N}$. We plot the functions $b_{c}^{N}$ and $c_{c}^{N}$ we obtained for various values of $N$ respectively
in Figures 3.1 and 3.2 below. That convergence is indeed achieved can immediately be observed in the figures.



Figure 3.2

In the fixed-order case with $n_{c}=1$, the compensator gains $A_{c}, B_{c}$, and $C_{c}$ are all scalars. Also, for a first-order controller there are only two independent parameters, $A_{c}$ and $B_{c} C_{c}$. In Table 3.1 below we give the values we obtained for $A_{c}^{N}$ and $B_{c}^{N} C_{c}^{N}$ for various values of $N$. Once again, it is clear that the gains are converging as $N$ increases. In addition, in Table 3.1 we provide the closed-loop costs $J_{\text {LQG }}^{N}$ and $J_{F O}^{N}$ for the LQG and first-order controllers. These closed-loop costs were evaluated using a 64th-order modal approximation to the infinite-dimensional system. For all values of $N$ the performance of the fixed-order compensator was within $2 \%$ of the corresponding LQG controller. Thus, for example, the replacement of a 32nd-order approximating optimal LQG controller by an approximating optimal first-order controller will yield considerable implementation simplification with only minor performance degradation. Note that for the example we consider here, it is possible to compute the open-loop cost for the infinite-dimensional system in closed form. We have

$$
\begin{aligned}
J_{O L} & =\operatorname{tr} \int_{0}^{\infty} V_{1} T^{*}(t) R_{1} T(t) d t=v_{1} r_{1} \operatorname{tr} \int_{0}^{\infty} T(t)^{2} d t \\
& =v_{1} r_{1} \sum_{n=1}^{\infty} \int_{0}^{\infty} e^{-2 n^{2} \pi^{2} a t} d t=\frac{v_{1} r_{1}}{2 \pi^{2} a} \sum_{n=1}^{\infty} \frac{1}{n^{2}} \\
& =\frac{v_{1} r_{2}}{12 a}=\frac{1}{12} \approx .08333 .
\end{aligned}
$$

Finally, for comparison purposes, we tried applying balancing techniques to the LQG controllers to reduce their order. However, with $n_{c}=1$, such controllers were found to be destabilizing. Based upon the results in [13], this was not unexpected.

| $N$ | $A_{c}^{N}$ | $B_{c}^{N} C_{c}^{N}$ | $J_{L Q G}^{N}$ | $J_{F O}^{N}$ |
| :---: | :---: | :---: | :---: | :---: |
| 4 | -687.6 | 5470 | .06999 | .07014 |
| 8 | -720.9 | 5231 | .06870 | .06993 |
| 12 | -730.9 | 5182 | .06872 | .06991 |
| 16 | -734.3 | 5145 | .06874 | .06990 |
| 20 | -738.0 | 5127 | .06875 | .06990 |
| 24 | -737.6 | 5108 | .06876 | .06990 |
| 28 | -739.8 | 5109 | .06876 | .06990 |
| 32 | -738.7 | 5099 | .06877 |  |

Table 3.1

As a second example we consider the one-dimensional, single-input, single-output hereditary control system given by

$$
\begin{gather*}
\dot{v}(t)=a_{0} v(t)+a_{1} v(t-\rho)+b_{0} u(t)+h_{1} w(t), \quad t>0,  \tag{3.5}\\
y(t)=c_{0} v(t)+h_{2} w(t), \quad t>0, \tag{3.6}
\end{gather*}
$$

where $a_{0}, a_{1}, b_{0}, c_{0}, h_{1}, h_{2}, \rho \in \mathbb{R}^{1}$ with $h_{2} \neq 0$, and $w$ is a unit-intensity white noise process. To rewrite (3.5), (3.6) in the form (2.1), (2.2), we take $\mathcal{X}=\mathbb{R}^{1} \times L_{2}(-\rho, 0)$ endowed with the usual product space inner product, $\langle(\eta, \phi),(\xi, \psi)\rangle=\eta \xi+\int_{-\rho}^{0} \phi \psi$, and let $x(t)=\left(v(t), v_{t}\right), t \geq 0$, where for $t \geq 0, v_{t} \in L_{2}(-\rho, 0)$ is given by $v_{t}(\theta)=v(t+\theta),-\rho \leq \theta \leq 0$. Define $A: \operatorname{Dom}(A) \subset X \rightarrow X$ by $A \hat{\phi}=\left(a_{0} \phi(0)+a_{1} \phi(-\rho), D \phi\right)$ for $\hat{\phi}=(\phi(0), \phi) \in \operatorname{Dom}(A) \triangleq\left\{(\xi, \psi) \in \mathcal{X}: \psi \in H^{1}(-\rho, 0), \psi(0)=\right.$ $\xi\}$, and let $B \in \mathcal{L}\left(\mathbb{R}^{1}, X\right)$ and $C \in \mathcal{L}\left(X, \mathbb{R}^{1}\right)$ be given by $B u=\left(b_{0} u, 0\right)$ and $C(\eta, \phi)=c_{0} \eta$, respectively. Let $\hat{X}=\mathbb{R}^{1}$ and define $H_{1} \in \mathcal{L}(\hat{X}, X)$ and $H_{2} \in \mathcal{L}\left(\hat{X}, \mathbb{R}^{1}\right)$ by $H_{1} z=\left(h_{1} z, 0\right)$ and $H_{2} z=h_{2} z$, for $z \in \mathbb{R}^{1}$.

The operator $A$ is densely defined and is the infinitesimal generator of a $C_{0}$ semigroup $\{T(t): t \geq$ $0\}$ of bounded linear operators on $\mathcal{X}$ with $T(t)(\eta, \phi)=\left(v(t ; \eta, \phi), v_{t}(\eta, \phi)\right), t \geq 0$, where $\bar{v}(\cdot ; \eta, \phi)$ is the unique solution to (3.5) with $b_{0}=h_{1}=0$, and initial conditions $v(0)=\eta, v_{0}=\phi$. We take $R_{1} \in \mathcal{L}(X)$ and $R_{2} \in \mathcal{L}\left(\mathbb{R}^{1}\right)$ to be $R_{1}(\eta, \phi)=\left(r_{1} \eta, 0\right)$ and $R_{2} u=r_{2} u$, respectively, with $r_{1}, r_{2}>0$. The definitions of $H_{1}$ and $H_{2}$ given above imply that $V_{1} \in \mathcal{L}(\mathcal{X})$ and $V_{2} \in \mathcal{L}\left(\mathbb{R}^{1}\right)$ are given by $V_{1}(\eta, \phi)=\left(h_{1}^{2} \eta, 0\right)$ and $V_{2} z=h_{2}^{2} z$, for $(\eta, \phi) \in X$ and $z \in \mathbb{R}^{1}$.

We employ an approximation scheme recently proposed by Ito and Kappel in [26]. We briefly outline it here; a more detailed discussion can be found in [26]. For each $N=1,2, \ldots$ let $\chi_{j}^{N} \in$ $L_{2}(-\rho, 0)$ denote the characteristic function for the interval $[-j \rho / N,-(j-1) \rho / N), j=1,2, \ldots, N$, and let $X^{N}$ be the ( $N+1$ )-dimensional subspace of $X$ defined by

$$
x^{N}=\operatorname{span}\left\{(1,0),\left(0, \chi_{1}^{N}\right), \ldots,\left(0, \chi_{N}^{N}\right)\right\} .
$$

Let $\mathcal{P}^{N}: \mathcal{X} \rightarrow \mathcal{X}^{n}$ denote the orthogonal projection of $X$ onto $\mathscr{X}^{N}$. Let $\left\{\phi_{j}^{N}\right\}_{j=0}^{N}$ denote the linear $B$ spline functions defined on the interval $[-\rho, 0]$ with respect to the uniform mesh $[-\rho, \ldots,-\rho / N, 0\}$, and set $X_{1}^{N}=\operatorname{span}\left\{\left(\phi_{j}^{N}(0), \phi_{j}^{N}\right)\right\}_{j=0}^{N}$. Then $X_{1}^{N}$ is an $(N+1)$-dimensional subspace of $\operatorname{Dom}(A)$ and it is not difficult to demonstrate that the restriction of $P^{N}$ to $X_{1}^{N}$ is a bijection onto $X^{N}$. Using the fact that $A$ restricted to $X_{1}^{N}$ has range in $X^{N}$, we define $A^{N} \in \mathcal{L}\left(X^{N}\right)$ by $A^{N}=A\left(P^{N}\right)^{-1}$, and set $T^{N}(t)=\exp \left(A^{N} t\right), t \geq 0$. Noting that $R(B) \subset \mathcal{X}^{N}$, we take $B^{N} \in \mathcal{L}\left(\mathbb{R}^{1}, X^{N}\right)$ to be given by $B^{N}=B$. Similarly, we take $R_{1}^{N}=R_{1}$ and $V_{1}^{N}=V_{1}$. We set $C^{N}=C$.

It is shown in [26] that $P^{N}(\eta, \phi) \rightarrow(\eta, \phi), T^{N}(t) P^{N}(\eta, \phi) \rightarrow T(t)(\eta, \phi)$, and $T^{N}(t)^{\bullet} P^{N}(\eta, \phi) \rightarrow T(t)^{*}(\eta, \phi)$ for $(\eta, \phi) \in X$ as $N \rightarrow \infty$, uniformly in $t$, for $t$ in bounded subsets of $[0, \infty)$. It then follows that $\lim _{N \rightarrow \infty} B^{N}=B$ and $\lim _{N \rightarrow \infty} C^{N} \rho^{N}=C$, in norm.

For the LQG (full-order) problem, the optimal functional observer and feedback control gains $b_{c}$ and $c_{c}$ are of the form $b_{c}=\left(\beta_{0}, \beta_{1}\right)$ and $c_{c}=\left(\gamma_{0}, \gamma_{1}\right)$ with $\beta_{0}, \gamma_{0} \in \mathbb{R}^{1}$, and $\beta_{1}, \gamma_{1} \in L_{3}(-\rho, 0)$. The approximating gains are of the form $b_{c}^{N}=\left(\beta_{0}^{N}, \beta_{1}^{N}\right)$ and $c_{c}^{N}=\left(\gamma_{0}^{N}, \gamma_{1}^{N}\right)$ with $\beta_{0}^{N}, \gamma_{0}^{N} \in \mathbb{R}^{1}$ and $\beta_{1}^{N}, \gamma_{1}^{N} \in \operatorname{span}\left\{\chi_{j}^{N}\right\}_{j=1}^{N}$. Since we are treating a one-dimensional example, if $b_{0} \neq 0$, the theory in [26] implies that $\beta_{0}^{N} \rightarrow \beta_{0}$ and $\gamma_{0}^{N} \rightarrow \gamma_{0}$ in $\mathbb{R}^{2}$, and $\beta_{1}^{N} \rightarrow \beta_{1}$, and $\gamma_{1}^{N} \rightarrow \gamma_{1}$ in $L_{2}(-\rho, 0)$, as $N \rightarrow \infty$.

Once again, as in the first example, matrix representations for the operators $A^{N}, B^{N}, C^{N}, R_{1}^{N}$, and $V_{1}^{N}$ are not difficult to compute in closed form. Indeed, the $(N+1) \times(N+1)$ matrix representation for the bijection $P^{N}: \mathscr{X}_{1}^{N} \rightarrow X_{1}$ is given by

$$
\left[P^{N}\right]=\left[\begin{array}{cccccccc}
1 & 0 & & & & & 0 & \\
\frac{1}{2} & \frac{1}{2} & 0 & & & & & \\
0 & \frac{1}{2} & \frac{1}{2} & 0 & & & & \\
& & & & \ddots & & & \\
& 0 & & & 0 & \frac{1}{2} & \frac{1}{2} & 0 \\
& 0 & & & & 0 & \frac{1}{2} & \frac{1}{2}
\end{array}\right] .
$$

Then $\left[A^{N}\right]=\left[K^{N}\right]\left[P^{N}\right]^{-1}$, where

$$
\left[K^{N}\right]=\left[\begin{array}{cccccccc}
a_{0} & 0 & & & & 0 & & a_{1} \\
\frac{N}{\rho} & -\frac{N}{\rho} & 0 & & & & & \\
0 & \frac{N}{\rho} & -\frac{N}{\rho} & 0 & & & & \\
\\
& & & & \ddots & & & \\
& 0 & & & & & \frac{N}{\rho} & -\frac{N}{\rho} \\
& 0 & & & 0 & \frac{N}{\rho} & -\frac{N}{\rho}
\end{array}\right] .
$$

We have the $(N+1) \times 1$ matrix $\left[B^{N}\right]=\left[b_{0} 0 \ldots 0\right]^{T}$ and the $1 \times(N+1)$ matrix $\left[C^{N}\right]=\left[c_{0} 0 \ldots 0\right]$, while $\left[R_{1}^{N}\right]=r_{1}\left[\mathcal{M}^{N}\right]$ and $\left[V_{1}^{N}\right]=h_{1}^{2}\left[\mathcal{M}^{N}\right]$ where the $(N+1) \times(N+1)$ matrix $\left[\mathcal{M}^{N}\right]$ is given by

$$
\left[M^{N}\right]=\left[\begin{array}{llll}
1 & 0 & & 0 \\
0 & & & 0 \\
& & \ddots & \\
0 & & 0
\end{array}\right]
$$

We set $a_{0}=a_{1}=b_{0}=c_{0}=r_{1}=h_{1}=\rho=1, r_{2}=.1$, and $h_{2}=\sqrt{.1}$ and computed approximating optimal LQG (i.e., $n_{c}=N+1$ ) and first-order (i.e., $n_{c}=1$ ) compensators for
$N=8,16,24$, and 32. The optimal LQG observer gains are given in Table 3.3 and Figure 3.3; the control gains are given in Table 3.4 and Figure 3.4. The first 23 open-loop poles of the system (see [27]) are given in Table 3.2. The approximating first-order compensator gains along with the corresponding and LQG closed-loop costs are given in Table 3.5 below. These costs were computed using an evaluation model obtained by setting $N=64$.. Note that the performance of the firstorder controllers is within $10 \%$ of the performance of the LQG controllers. Once again it is clear that convergence is achieved.

$$
\begin{aligned}
1.278465 & \\
-1.588317 & \pm 4.155305 \mathrm{i} \\
-2.417631 & \pm 10.68603 \mathrm{i} \\
-2.861502 & \pm 17.05611 \mathrm{i} \\
-3.167754 & \pm 23.38558 \mathrm{i} \\
-3.401945 & \pm 29.69798 \mathrm{i} \\
-3.591627 & \pm 36.00146 \mathrm{i} \\
-3.751047 & \pm 42.29965 \mathrm{i} \\
-3.888543 & \pm 48.59442 \mathrm{i} \\
-4.009422 & \pm 54.88686 \mathrm{i} \\
-4.117267 & \pm 61.17761 \mathrm{i} \\
-4.214618 & \pm 67.46710 \mathrm{i}
\end{aligned}
$$

Table 3.2


Figure 3.3

| $N$ | 8 | 16 | 24 | 32 |
| :---: | :---: | :---: | :---: | :---: |
| $\beta_{0}^{N}$ | 4.4213 | 4.4229 |  | 4.4233 |

Table 3.3


Figure 3.4

| $N$ | 8 | 16 | 24 | 32 |
| :---: | :---: | :---: | :---: | :---: |
| $\gamma_{0}^{N}$ | -4.4213 | -4.4229 | -4.4233 | -4.4234 |

Table 3.4

| $N$ | $A_{c}^{N}$ | $B_{c}^{N} C_{c}^{N}$ | $J_{L Q G}^{N}$ | $J_{F O}^{N}$ |
| :---: | :---: | :---: | :---: | :---: |
| 8 | -4.835 | -16.057 | 1.4042 | 1.5221 |
| 16 | -4.936 | -16.343 | 1.403877 | 1.5298 |
| 24 | -4.959 | -16.378 | 1.403856 | 1.5309 |
| 32 | -4.962 | -16.404 | 1.403852 | 1.5317 |

Table 3.5

## 4. Summary and Concluding Remarks

We have proposed an approximation technique for computing optimal fixed-order compensators for distributed parameter systems. Our approach involves using the optimal projection theory for infinite-dimensional systems (which characterizes the optimal fixed-order compensator) developed in [18] in conjunction with finite-dimensional approximation of the infinite-dimensional plant. We demonstrated the feasibility of our approach with two examples wherein we used spline-based RitzGalerkin finite element schemes to compute approximating optimal first-order controllers for onedimensional, singe-input/output, parabolic (heat/diffusion) and hereditary control systems. Our numerical studies indicate that convergence of the compensator gains is achieved and that using the first-order controller would lead to only minimal performance degradation over a standard LQG compensator while yielding significant implementation simplification.

At this point one is led naturally to ask the question of whether or not a satisfactory convergence theory could be developed. We are working on this at present and expect that such a theory would conform closely in form and spirit to the convergence results for LQG approximation found in [9] and [10] and outlined in Section 2 above. We also intend to consider our approximation ideas in the context of discrete-time or sampled-data systems, and for continuous-time systems involving unbounded input and/or output (for example, boundary control systems), and systems with control or measurement delays, see [11],[12]). Finally, we intend to investigate the application of our approximation framework to other infinite-dimensional control systems, in particular the vibration control of flexible structures (i.e., second-order systems such as wave, beam, or plate equations).

Acknowledgment. The authors gratefully acknowledge Mr. S. W. Greeley, Mr. S. Richter, and Mr. A. Daubendiek of Harris Corp. for carrying out the numerical computations reported in Section 3. We also wish to thank Ms. J. M. Straehla, also of Harris Corp., for preparing the original manuscript in $T_{E} X$.

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## APPENOIX D: Decentralized Control

[^10]
# Sequential design of decentralized dynamic compensators using the optimal projection equations 

DENNIS S. BERNSTEIN $\dagger$


#### Abstract

The optimal projection equations for quadratically optimal centralized fixed-order dynamic compensation are generalized to the case in which the dynamic compensator has, in addition, a fixed decentralized structure. Under a stabilizability assumption for the particular feedback configuration, the resulting optimality conditions explicitly characterize each subcontroller in terms of the plant and remaining subcontrollers. This characterization associates an oblique projection with each subcontroller and suggests an iterative sequential design algorithm. The reselts are applied to an interconnected flexible beam example.


## 1. Introduction

The purpose of this note is to consider the problem of designing decentralized dynamic feedback controllers using recently obtained results on quadratically optimal fixed-order dynamic compensation (Hyland and Bernstein 1984). As in Bernussou and Titli (1982). Looze et al. (1978), and Singh (1981), the overall approach is to fix the structure (information pattern and order) of the linear controller and optimize the steady-state regulation cost with respect to the controller parameters. The underlying philosophy is that the ability to carry out such an optimization procedure permits the evaluation of a particular decentralized configuration which may be dictated by implementation constraints. If there is some flexibility in designing the decentralized architecture, then these results can be used to evaluate the optimal performance of each permissible configuration, and hence to determine preferable structures. Since the present paper is confined to the question of optimal regulation, trade-offs with It ard to robustness in the presence of plant variations are not considered. Such trade-offs can be included, however, by utilizing the Stratonovich multiplicative white noise approach developed by Bernstein and Hyland (1985).

To further motivate our approach, consider the problem of controlling an $n$ thorder plant $\mathscr{P}$ by means of a decentralized dynamic compensator consisting of subcontrollers $\mathscr{C}_{1}$ and $\mathscr{C}_{2}$. A straightforward design technique that immediately comes to mind is that of sequential optimization (Davison and Gesing 1979, Jamshidi 1983). To begin, ignore $\mathscr{C}_{2}$ and design $\mathscr{C}_{1}$ as a centralized controller for $\mathscr{P}$. Next, regard the closed-loop system consisting of $\mathscr{P}$ and $\mathscr{C}_{1}$ as an augmented system $\mathscr{P}^{\prime}$ and design $\mathscr{C}_{2}$ as a centralized controller for $\mathscr{P}^{\prime}$. Now redesign $\mathscr{C}_{1}$ to be a centralized controller for the augmented closed-loop system composed of $\mathscr{P}$ and $\mathscr{C}_{2}$, and so forth. One difficulty with this scheme, however, is that of dimension. If, for example, one were to employ LQG at each step of this algorithm, then on the first iteration $\mathscr{C}_{1}$ would have dimension $n$ and thus $\mathscr{C}_{2}$ would have dimension $2 n$. On the second iteration, $\mathscr{C}_{1}$ would require dimension $3 n$ and $\mathscr{C}_{2}$ would have order $4 n$, and so forth. Such

[^11]difficulties can be avoided by setting $n=0$, which essentially corresponds to static output feedback. Although easier to implement, static output feedback lacks filtering abilities such as are inherent in LQG controllers, which are purely dynamic (i.e. strictly proper).

As discussed by Sandell et al. (1978), p. 119, the explanation for this difficulty is provided by the 'second-guessing' phenomenon: when LQG is used, each subcontroller must consist of linear feedback, not only of estimates of the plant states but also of estimates of the other subcontrollers' estimates. Hence the 'optimal' controller is given by an irrational transfer function, i.e. a distributed parameter (infinitedimensional) system. Such controllers, of course, must be ruled out since their design and implementation (except in special cases) violate physical realizability (see, for example, Bernstein and Hyland 1986).

Having thus ruled out zeroth-order and infinite-order decentralized controllers, we focus on the problem of designing purely dynamic decentralized compensators. Moreover, by invoking the constraint of fixed subcontroller order, we overcome the second-guessing phenomenon. Utilizing the parameter optimization approach thus leads to a generalization of the result obtained by Hyland and Bernstein (1984) for centralized control. In brief, it was shown in Hyland and Bernstein (1984) that the unwieldy first-order necessary conditions for fixed-order dynamic compensation can be simplified by exploiting the presence of a previously unrecognized oblique projection. The resulting optimal projection equations, which consist of a pair of modified Riccati equations and a pair of modified Lyapunov equations coupled by the optimal projection, yield insight into the structure of the optimal dynamic compensator and emphasize the breakdown of the separation principle for reduced-order controller design. For example, the optimal compensator is the projection of a fullorder dynamic controller which is generally different from the LQG design. Furthermore, this full-order controller and the oblique projection are intricately related since they are simultaneously determined by the coupled design equations. An immediate consequence is the observation that stepwise schemes employing either model reduction followed by LQG or LQG followed by model reduction are generally suboptimal. For computational purposes, the optimal projection equations permit the development of novel numerical methods which operate through successive iteration of the oblique projectic . (Hyland and Bernstein 1985). Such algorithms are thus philosophically and operationally distinct from gradient search methods.

The generalization of the optimal projection equations to the decentralized case is straightforward and immediate. In the optimization process each subcontroller is viewed as a centralized controller for an augmented 'plant' consisting of the actual plant and all other subcontrollers. It need only be observed that the necessary conditions for optimality for the decentralized problem must consist of the collection of necessary conditions obtained by optimizing over each subcontroller separately while keeping the other subcontrollers fixed. More precisely, this statement corresponds to the fact that setting the Frechet derivative to zero is equivalent to setting the individual partial derivatives to zero. Hence it is not surprising that the optimal projection ${ }_{-}$- uations for the decentralized problem involve multiple oblique projections, one associated with each subcontroller. Furthermore, each subcontroller incorporates an internal model (in the sense of an oblique projection of full-order dynamics) not only of the plant but also of all other subcontrollers. The structure of the equations suggests a sequential design algorithm such as that proposed in this work.

The simplicity with which this result is obtained should not belie its relevance to the decentralized control problem. Specifically, our approach is distinct from sub-system-decomposition techniques (Ikeda and Siljak 1980, 1981, Ikeda et al. 1981, 1984, Lindner 1985, Linnemann 1984, Ozguner 1979, Ramakrishna and Viswanadham 1982, Saeks 1979, Sezer and Huseyin 1984, Silkak 1978, 1983) and model-reduction methods since the optimal projection equations retain the full, interconnected plant at all times. For the proposed algorithm, decomposition techniques which exploit subsystem-interconnection data can play a role by providing a starting point for subsequent iterative refinement and optimization. Decomposition methods may also play a role when very high dimensionality precludes direct solution of the optimal projection equations. These are areas for future research.

With regard to the role of the oblique projection, it should be noted that such transformations do not, in general, preserve plant characteristics such as poles, zeros, subspaces, etc. Indeed, since the oblique projection arises as a consequence of optimality, approaches that seek to retain system invariants (e.g. Uskokovic and Medanic 1985) are generally suboptimal. In addition, the complex coupling among the plant and subcontrollers via multiple oblique projections provides an additional measure for evaluating the suboptimality of the methods proposed.

The plan of the paper is as follows. The fixed-structure decentralized dynamiccompensation problem is stated in $\S 2$ along with the generalization of the optimal projection equations. In § 3 we propose a sequential design algorithm for solving these equations and state conditions under which convergence is guaranteed. Finally, in § 4 the algorithm is applied to the 8th-order model of a pair of simply supported beams connected by a spring. For this example, we obtain a two-channel decentralized design which is 4th-order in each channel and compare its performance with the (8thorder) centralized LQG design.

## 2. Problem statement and main theorem

Given the controlled system

$$
\begin{align*}
& \dot{x}(t)=A x(t)+\sum_{i=1}^{p} B_{i} u_{i}(t)+w_{0}(t)  \tag{2.1}\\
& y_{i}(t)=C_{i} x(t)+w_{i}(t), \quad i=1, \ldots, p \tag{2.2}
\end{align*}
$$

design a fixed-structure decentralized dynamic compensator

$$
\begin{gather*}
\dot{x}_{\mathrm{c} i}(t)=A_{\mathrm{c} i} x_{\mathrm{c} i}(t)+B_{\mathrm{c} i} y_{j}(t), \quad i=1, \ldots, p  \tag{2.3}\\
u_{i}(t)=C_{\mathrm{c} i} x_{\mathrm{c} i}(t), \quad i=1, \ldots, p \tag{2.4}
\end{gather*}
$$

which minimizes the steady-state performance criterion

$$
\begin{equation*}
J\left(A_{\mathrm{c} 1}, B_{\mathrm{c} 1}, C_{\mathrm{c} 1}, \ldots, A_{\mathrm{cp}}, B_{c p}, C_{\mathrm{c} p}\right) \triangleq \lim _{t \rightarrow \infty} \mathbb{E}\left[x(t)^{\mathrm{T}} R_{0} x(t)+\sum_{i=1}^{p} u_{i}(t)^{\mathrm{T}} R_{i} u_{i}(t)\right] \tag{2.5}
\end{equation*}
$$

where, for $i=1, \ldots, p: x \in \mathbb{R}^{n}, u_{i} \in \mathbb{R}^{m_{i}}, y_{i} \in \mathbb{R}^{l_{i}}, c_{c i} \in \mathbb{R}^{n_{c i}}, n_{c} \triangleq \sum_{i=1}^{p} n_{c i}, n_{c i} \leqslant n+n_{c}-n_{c i}$. $A, B_{i}, C_{i}, A_{c i}, B_{c i}, C_{c i}, R_{0}$ and $R_{i}$ are matrices of appropriate dimension with $R_{0}$ (symmetric) non-negative definite and $R_{i}$ (symmetric) positive definite; $w_{0}$ is white disturbance noise with $n \times n$ non-negative-definite intensity $V_{0}$, and $w_{i}$ is white
observation noise with $l_{i} \times l_{i}$ positive-definite intensity $V_{i}$, where $w_{0}, w_{1}, \ldots, w_{p}$ are mutually uncorrelated and have zero mean. $\mathbb{E}$ denotes expectation and superscript $T$ indicates transpose.

To guarantee that $J$ is finite and independent of initial conditions we restrict our attention to the set of admissible stabilizing compensators

$$
\mathscr{A} \triangleq\left\{\left(A_{\mathrm{c} 1}, B_{\mathrm{c} 1}, C_{\mathrm{c} 1}, \ldots, A_{\mathrm{cp}}, B_{\mathrm{cp}}, C_{\mathrm{cp}}\right): \tilde{A} \text { is asymptotically stable }\right\}
$$

where the closed-loop dynamics matrix $\tilde{A}$ is given by

$$
\tilde{A} \triangleq\left[\begin{array}{cc}
A & \tilde{B} C_{\mathrm{c}} \\
B_{\mathrm{c}} \tilde{C} & A_{\mathrm{c}}
\end{array}\right]
$$

where

$$
\begin{aligned}
& \tilde{B} \triangleq\left[\begin{array}{lll}
B_{1} & \ldots & B_{p}
\end{array}\right], \quad \tilde{C} \triangleq\left[\begin{array}{c}
C_{1} \\
\vdots \\
C_{p}
\end{array}\right] \\
& A_{c} \triangleq \text { block-diagonal }\left(A_{\mathrm{cl}}, \ldots, A_{\mathrm{cp}}\right) \\
& B_{c} \triangleq \text { block-diagonal }\left(B_{\mathrm{cl}}, \ldots, B_{\mathrm{cp}}\right) \\
& C_{c} \triangleq \text { block-diagonal }\left(C_{\mathrm{c} 1}, \ldots, C_{c p}\right)
\end{aligned}
$$

(For possibly non-square matrices $S_{1}, S_{2}$, block-diagonal ( $S_{1}, S_{2}$ ) denotes the

$$
\left.\operatorname{matrix}\left[\begin{array}{cc}
S_{1} & 0 \\
0 & S_{2}
\end{array}\right]\right)
$$

It is possible that for certain decentralized structures the system is nut stabilizable, i.e. $\mathscr{A}$ is empty (Wang and Davison 1973, Seraji 1982, Sezer and Siljak 1981). Our approach, however, is to assume that $\mathscr{A}$ is not empty and characterize the optimal decentralized controller over the stabilizing class. Since the value of $J$ is independent of the internal realization of each subcompensator, without loss of generality we can further restrict our attention to

$$
\mathscr{A}_{+} \triangleq\left\{\left(A_{\mathrm{c} 1}, B_{\mathrm{c} 1}, C_{\mathrm{c} 1}, \ldots, A_{\mathrm{cp}}, B_{\mathrm{c} p}, C_{\mathrm{c} p}\right) \in \mathscr{A}:\left(A_{\mathrm{c} i}, B_{\mathrm{ci}}\right)\right. \text { is controllable and }
$$

$$
\left.\left(C_{\mathrm{ci}}, A_{\mathrm{c} i}\right) \text { is observable, } i=1, \ldots, p\right\}
$$

The following lemma is an immediate consequence of Theorem 6.2.5, p. 123 of Rao and Mitra (1971). Let $I_{r}$ denote the $r \times r$ identity matrix.

## Lemma 2.1

Suppose $\hat{Q}, \hat{P} \in \mathbb{R}^{9 \times q}$ are non-negative definite and rank $\hat{Q} \hat{P}=r$. Then there exist $G, \Gamma \in \mathbb{R}^{\times \times q}$ and invertible $M \in \mathbb{R}^{r \times r}$ such that

$$
\begin{gather*}
\hat{Q} \hat{P}=G^{\mathrm{T}} M \Gamma  \tag{2.6}\\
\Gamma G^{\mathrm{T}}=I_{r} \tag{2.7}
\end{gather*}
$$

For convenience in stating the main theorem, call ( $G, M, \Gamma$ ) satisfying (2.6), (2.7) a
projective factorization of $\hat{Q} \hat{P}$. Such a factorization is unique modulo an arbitrary change in basis in $\mathbb{R}^{\boldsymbol{r}}$, which corresponds to nothing more than a change of basis for the internal representation of the compensator (or subcompensators in the present context).

We shall also require the following notation. Let $\tilde{A}_{i}$ denote $\tilde{A}$ with the rows and columns containing $A_{\mathrm{c} i}$ deleted. Similarly, let $\tilde{R}_{i}$ be obtained by deieting the rows and columns corresponding to $C_{\mathrm{ci}}^{\mathrm{T}} R_{i} C_{\mathrm{ci}}$ in the matrix

$$
\tilde{R} \triangleq \text { block-diagonal }\left(R_{0}, C_{\mathrm{c} 1}^{\mathrm{T}} R_{1} C_{\mathrm{c} 1}, \ldots, C_{\mathrm{cp}}^{\mathrm{T}} R_{p} C_{\mathrm{cp}}\right)
$$

And furthermore, $\widetilde{V}_{i}$ is obtained by deleting the rows and columns containing $B_{\mathrm{ci}} V_{i} B_{\mathrm{ci}}^{\mathrm{T}}$ in

$$
\tilde{V} \triangleq \text { block-diagonal }\left(V_{0}, B_{c 1} V_{1} B_{c 1}^{\mathrm{T}}, \ldots, B_{c p} V_{p} B_{c p}^{\mathrm{T}}\right)
$$

Also define

$$
\tilde{B}_{i} \triangleq\left[\begin{array}{c}
B_{i} \\
0_{\left(n_{\mathrm{c}}-n_{\mathrm{c} i}\right) \times m_{i}}
\end{array}\right], \quad \tilde{C}_{i} \triangleq\left[\begin{array}{ll}
C_{i} & 0_{l_{i} \times\left(n_{c}-n_{c i}\right)}
\end{array}\right]
$$

where $0_{r \times s}$ denotes the $r \times s$ zero matrix. Note that $\tilde{A}_{i}, \tilde{B}_{i}, \tilde{C}_{i}, \tilde{R}_{i}$ and $\tilde{V}_{i}$ essentially represent the closed-loop system minus the $i$ th subcontroller as controlled by the latter. Finally, define

$$
\Sigma_{i} \triangleq \tilde{B}_{i} R_{i}^{-1} \tilde{B}_{i}^{\mathrm{T}}, \quad \Sigma_{i} \triangleq \tilde{C}_{i}^{\mathrm{T}} V_{i}^{-1} \tilde{C}_{i}
$$

and, for $\tau \in \mathbb{R}^{r \times r}$, let

$$
\tau_{\perp} \triangleq I_{r}-\tau
$$

## Main theorem

Suppose $\left(A_{\mathrm{c} 1}, B_{c 1}, C_{\mathrm{c} 1}, \ldots, A_{c p}, B_{c p}, C_{c p}\right) \in \mathscr{A}+$ solves the steady-state fixedstructure decentralized dynamic-compensation problem. Then for $i=1, \ldots, p$ there exist $\left(n+n_{\mathrm{c}}-n_{\mathrm{c} i}\right) \times\left(n+n_{\mathrm{c}}-n_{\mathrm{c} i}\right)$ non-negative-definite matrices $Q_{i}, P_{i}, \hat{Q}_{i}$ and $\hat{P}_{i}$ such that $A_{\mathrm{c} i}, B_{\mathrm{c} i}$ and $C_{\mathrm{c} i}$ are given by

$$
\begin{gather*}
A_{\mathrm{c} i}=\Gamma_{i}\left(\tilde{A}_{i}-Q_{i} \Sigma_{i}-\Sigma_{i} P_{i}\right) G_{i}^{\mathrm{T}}  \tag{2.8}\\
B_{\mathrm{c} i}=\Gamma_{i} Q_{i} \tilde{C}_{i}^{\mathrm{T}} V_{i}^{-1}  \tag{2.9}\\
C_{\mathrm{ci} i}=-R_{i}^{-1} \tilde{B}_{i}^{\mathrm{T}} P_{i} G_{i}^{\mathrm{T}} \tag{2.10}
\end{gather*}
$$

for some projective factorization $G_{i}, M_{i}, \Gamma_{i}$ of $\hat{Q}_{i} \hat{P}_{i}$, and such that, with $\tau_{i}=G_{i}^{\top} \Gamma_{i}$, the following conditions are satisfied:

$$
\begin{gather*}
0=\tilde{A}_{i} Q_{i}+Q_{i} \tilde{A}_{i}^{\mathrm{T}}+\tilde{V}_{i}-Q_{i} \Sigma_{i} Q_{i}+\tau_{i 1} Q_{i} \Sigma_{i} Q_{i} \tau_{i \perp}^{\mathrm{T}}  \tag{2.11}\\
0=\tilde{A}_{i}^{\mathrm{T}} P_{i}+P_{i} \tilde{A}_{i}+\tilde{R}_{i}-P_{i} \Sigma_{i} P_{i}+\tau_{i \perp}^{\mathrm{T}} P_{i} \Sigma_{i} P_{i} \tau_{i \perp}  \tag{2.12}\\
0=\left(\tilde{A}_{i}-\Sigma_{i} P_{i}\right) \hat{Q}_{i}+\hat{Q}_{i}\left(\tilde{A}_{i}-\Sigma_{i} P_{i}\right)^{\mathrm{T}}+Q_{i} \Sigma_{i} Q_{i}-\tau_{i \perp} Q_{i} \Sigma_{i} Q_{i} \tau_{i \perp}^{\mathrm{T}}  \tag{2.13}\\
0=\left(\tilde{A}_{i}-Q_{i} \Sigma_{i}\right)^{\mathrm{T}} \hat{P}_{i}+\hat{P}_{i}\left(\tilde{A}_{i}-Q_{i} \Sigma_{i}\right)+P_{i} \Sigma_{i} P_{i}-\tau_{i \perp}^{\mathrm{T}} P_{i} \Sigma_{i} P_{i} \tau_{i \perp}  \tag{2.14}\\
\operatorname{rank} \hat{Q}_{i}=\operatorname{rank} \hat{P}_{i}=\operatorname{rank} \hat{Q}_{i} \hat{P}_{i}=n_{\mathrm{c} i} \tag{2.15}
\end{gather*}
$$

## Remark 2.1

Because of (2.7) the matrix $\tau_{i}$ is idempotent, i.e. $\tau_{i}^{2}=\tau_{i}$. This projection corresponding to the ith subcontroller is an oblique projection (as opposed to an orthogonal projection) since it is not necessarily symmetric. Furthermore, $\tau_{i}$ is given in closed form by

$$
\tau_{i}=\hat{Q}_{i} \hat{P}_{i}\left(\hat{Q}_{i} \hat{P}_{i}\right)^{*}
$$

where ( ) ${ }^{*}$ denotes the (Drazin) group generalized inverse (see, for example, Campbell and Meyer, 1979, p. 124).

## 3. Proposed algarithm

## Sequential design algorithm

Step 1. Choose a starting point consisting of initial subcontroller designs;
Step 2. For a sequence $\left\{i_{k}\right\}_{k=1}^{\infty}$, where $i_{k} \in\{1, \ldots, p\}, k=1,2, \ldots$, redesign subcontroller $i_{k}$ as an optimal fixed-order centralized controller for the plant and remaining subcontrollers;
Step 3. Compute the cost $J_{k}$ of the current design and check $J_{k}-J_{k-1}$ for convergence.
Note that the first two steps of the algorithm consist of (i) bringing suboptimal subcontrollers 'on line' and (ii) iteratively refining each subcontroller. As discussed in $\S 1$, the choice of a starting design for Step 1 can be obtained by a variety of existing methods such as subsystem decomposition. As for subcontroller refinement, note that each subcontroller redesign procedure is equivalent to replacing a suboptimal subcontroller with a subcontroller which is optimal with respect to the plant and remaining subcontrollers.

## Proposition 3.1

For a given starting design and redesign sequence $\left\{i_{k}\right\}_{k=1}^{\infty}$ suppose that the optimal projection equations can be solved for each $k$ to yield the global minimum. Then $\left\{J_{k}\right\}_{k=1}^{\infty}$ is monotonically non-increasing and hence convergent.

Determining both a suitable starting point and redesign sequence for solvability and attaining the decentralized global minimum remain areas for future research. With regard to algorithms for solving the optimal projection equations for each subcontroller redesign procedure, details of proposed algorithms can be found in the works of Hyland (1983, 1984) and Hyland and Bernstein (1985).

## 4. Application to interconnected flexible beams

To demonstrate the applicability of the main theorem and the sequential design algorithm, we consider a pair of simply supported Euler-Bernoulli fiexible beams interconnected by a spring (see the Figure). Each beam possesses one rate sensor and one force actuator. Retaining two vibrational modes in each beam, we obtain the 8thorder interconnected model

$$
\begin{gathered}
A=\left[\begin{array}{ll}
A_{11} & A_{12} \\
A_{21} & A_{22}
\end{array}\right], \quad B_{1}=\left[\begin{array}{c}
B_{11} \\
0_{4 \times 1}
\end{array}\right], \quad B_{2}=\left[\begin{array}{l}
0_{4 \times 1} \\
B_{22}
\end{array}\right] \\
C_{1}=\left[\begin{array}{lll}
C_{11} & 0_{1 \times 4}
\end{array}\right], \quad C_{2}=\left[\begin{array}{ll}
0_{1 \times 4} & C_{22}
\end{array}\right]
\end{gathered}
$$


where
$A_{i i}=\left[\begin{array}{cccc}0 & \omega_{1 i} & 0 & 0 \\ -\omega_{1 i}-\left(k / \omega_{1 i}\right)\left(\sin \pi c_{i}\right)^{2} & -2 \zeta_{i} \omega_{1 i} & -\left(k / \omega_{2 i}\right)\left(\sin \pi c_{i}\right)\left(\sin 2 \pi c_{i}\right) & 0 \\ 0 & 0 & 0 & \omega_{2 i} \\ -\left(k / \omega_{1 i}\right)\left(\sin \pi c_{i}\right)\left(\sin 2 \pi c_{i}\right) & 0 & -\omega_{2 i}-\left(k / \omega_{2 i}\right)\left(\sin 2 \pi c_{i}\right)^{2} & -2 \zeta_{i} \omega_{2 i}\end{array}\right]$
$A_{i j}=\left[\begin{array}{cccc}0 & 0 & 0 & 0 \\ \left(k / \omega_{1 j}\right)\left(\sin \pi c_{i}\right)\left(\sin \pi c_{j}\right) & 0 & \left(k / \omega_{2 j}\right)\left(\sin \pi c_{i}\right)\left(\sin 2 \pi c_{j}\right) & 0 \\ 0 & 0 & 0 & 0 \\ \left(k / \omega_{1 j}\right)\left(\sin \pi c_{j}\right)\left(\sin 2 \pi c_{i}\right) & 0 & \left(k / \omega_{2 j}\right)\left(\sin 2 \pi c_{i}\right)\left(\sin 2 \pi c_{j}\right) & 0\end{array}\right]$

$$
\begin{gathered}
B_{i i}=\left[\begin{array}{c}
0 \\
-\sin \pi a_{i} \\
0 \\
-\sin 2 \pi a_{i}
\end{array}\right], \quad C_{i i}=\left[\begin{array}{llll}
0 & \sin \pi s_{i} & 0 & \sin 2 \pi s_{i}
\end{array}\right] \\
a_{i}=\hat{a}_{i} / L_{i}, \quad s_{i}=\hat{s}_{l} / L_{i}, \quad c_{i}=\hat{c}_{i} / L_{i}
\end{gathered}
$$

In the above definitions, $k$ is the spring constant, $\omega_{j i}$ is the $j$ th modal frequency of the $i$ th beam, $\zeta_{i}$ is the damping ratio of the $i$ th beam, $L_{i}$ is the length of the $i$ th beam, and $\hat{a}_{i}, \hat{s}_{i}$ and $\hat{c}_{i}$ are, respectively, the actuator, sensor and spring-connection coordinates as measured from the left in the Figure. The chosen values are

$$
\begin{array}{ccc}
k=10 \\
\omega_{1 i}=1, & \omega_{2 i}=4, & \zeta_{i}=0.005, \quad L_{i}=1, \quad i=1,2 \\
\hat{a}_{1}=0.3, & \hat{s}_{1}=0.65, & \hat{c}_{1}=0.6 \\
\hat{a}_{2}=0.8, & \hat{s}_{2}=0.2, \quad \hat{c}_{2}=0.4
\end{array}
$$

In addition, weighting and intensity matrices are chosen to be

$$
\begin{aligned}
& R_{1}=\text { block-diagonal }\left(\left[\begin{array}{cc}
1 & 0 \\
0 & 1 / \omega_{11}
\end{array}\right],\left[\begin{array}{cc}
1 & 0 \\
0 & 1 / \omega_{21}
\end{array}\right],\left[\begin{array}{cc}
1 & 0 \\
0 & 1 / \omega_{12}
\end{array}\right],\left[\begin{array}{cc}
1 & 0 \\
0 & 1 / \omega_{22}
\end{array}\right]\right) \\
& R_{2}=R_{3}=0 \cdot 1 I_{2} \\
& V_{0}=\text { block-diagonal }\left(\left[\begin{array}{ll}
0 & 0 \\
0 & 1
\end{array}\right],\left[\begin{array}{ll}
0 & 0 \\
0 & 1
\end{array}\right],\left[\begin{array}{ll}
0 & 0 \\
0 & 1
\end{array}\right],\left[\begin{array}{ll}
0 & 0 \\
0 & 1
\end{array}\right]\right) \\
& V_{1}=V_{2}=0.1 I_{2}
\end{aligned}
$$

For this problem the open-loop cost was evaluated and the centralized 8th-order LQG design was obtained to provide a baseline. To provide a starting point for the sequential design algorithm, a pair of 4th-order LQG controllers were designed for each beam separately ignoring the interconnection, i.e. setting $k=0$. The optimal projection equations were then utilized to iteratively refine each subcontroller. The results are summarized in the Table.

| Design | Cost |
| :---: | :---: |
| Open loop | $163 \cdot 5$ |
| Centralized LQG |  |
| $n_{c}=8$ | 19.99 |
| Suboptimal decentralized |  |
| $n_{\mathrm{c} 1}=n_{\mathrm{c} 2}=4$ | 59.43 |
| Redesign subcontroiler 2 | $28 \cdot 19$ |
| Redesign subcontroller 1 | $23 \cdot 29$ |
| Redesign subcontroller 2 | 23.04 |
| Redesign subcontroller 1 | 22.25 |
| Redesign subcontroller 2 | 21.94 |
| Redesign subcontroller 1 | 21.86 |
| Redesign subcontroller 2 | 21.81 |
| Redesign subcontroller 1 | 21.79 |

## Acknowledgments

This research was supported in part by the Air Force Office of Scientific Research under contracts $\mathrm{F} 49620-86-\mathrm{C}-0002$ and $\mathrm{F} 49620-86-\mathrm{C}-0038$.

The author wishes to thank Scott $W$. Greeley for providing the beam model in § 4 and for carrying out the design computations.

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## Robust Decentralized Output Feedback: The Static Controller Case

by

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Keywords: robust, guaranteed bounds, optimal, structured uncertainty, real parameters


#### Abstract

Sufficient conditions are developed for designing robust decentralized static output feedback controllers. The approach involves deriving necessary conditions for minimizing a bound on closedloop performance over a specified range of uncertain parameters. The effect of plant parameter variations on the closed-loop covariance is overbounded by means of a modified Lyapunov equation whose solutions are guaranteed to provide robust stability and performance.


Supported in part by the Air Force Office of Scientific Research under contract F49620-86-C0038.

## 1. Introduction

Because of implementation constraints, cost, and reliability considerations, a decentralized controller architecture is often required for controlling large scale systems. Furthermore, such controllers must be robust to variations in plant parameters. The present paper addresses both of these concerns within the context of a robust decentralized theory for continuous-time static controllers.

The approach to controller design considered herein invoives optimizing closed-loop performance with respect to the feedback gains. This approach to output feedback was studied for centralized controllers in $[8,9]$ and for decentralized controllers in [10]. An interesting feature of $[9,10]$ is the recognition of an oblique projection (idempotent matrix) which allows the necessary conditions to be written in terms of a modified Riccati equation. When the problem is specialized to full-state feedback, the projection becomes the identity and the modified Riccati equation coincides with the standard Riccati equation of LQR theory. It should be pointed out that this oblique projection is distinct from the oblique projection arising in dynamic compensation ([7]). A unified treatment of the static/dynamic (nonstrictly proper) centralized control problem involving both projections is given in [2].

The present paper goes beyond earlier work by deriving sufficient conditions for robust stability and performance with respect to variations in the plant parameters. Although plant disturbances are represented in the usual stochastic manner by means of additive white noise, uncertainty in the plant dynamics is modeled deterministically by means of constant structured parameter variations within bounded sets. Thus, for example, the dynamics matrix $A$ is replaced by $A+\sum_{k=1}^{p} \sigma_{k} A_{k}$, where $\sigma_{k}$ is a constant uncertain parameter assumed only to lie within the interval $\left[-\alpha_{k}, \alpha_{k}\right]$ but otherwise unknown, and $A_{k}$ is a fixed matrix denoting the structure of the uncertain parameter $\sigma_{k}$ as it appears in the nominal dynamics matrix $A$. The system performance is defined to be the worst-case value over the class of parameter uncertainties of a quadratic criterion averaged over the disturbance statistics.

Since the closed-loop performance can be written in terms of the second-moment matrix, a performance bound over the class of uncertain parameters can be obtained by bounding the state covariance. The key to bounding the state covariance is to replace the usual Lyapunov equation for the second-moment matrix by a modified Lyapunov equation. In the present paper the modified Lyapunov equation is constructed by adding two additional terms. The first term corresponds to a
uniform right shift of the open-loop dynamics. As is well known ([1]), such a shift may arise from an exponential performance weighting and leads to a uniform stability margin for the closed-loop system. In order to guarantee robustness with respect to specified structured parameter variations, however, an additional term of the form $A_{k} Q A_{k}^{T}$ is required. Such terms arise naturally in systems with multiplicative white noise; see $[3,4]$ and the references therein for further details. The exponential cost weighting and multiplicative noise interpretations for the uncertainty bound have no bearing in the present paper, however, since parameter variations are modeled deterministically as constant variations within bounded sets.

Having bounded the state covariance over the class of parameter uncertainties, the worst-case performance can thus be bounded in terms of the solution of the modified Lyapunov equation. The performance bound can be viewed as an auxiliary cost and thus leads to the Auxiliary Minimization Problem: Minimize the performance bound while satisfying the modified Lyapunov equation. The nice feature of the auxiliary problem is that necessary conditions for optimality of the performance bound now serve as sufficient conditions for robust performance in the original problem. Thus our approach seeks to rectify one of the principal drawbacks of necessity theory, namely, guarantees of stability and performance. Furthermore, it should be noted that if numerical solution of the optimality conditions yields a local extremal which is not the global optimum, then robust stability and performance are still guaranteed, although the performance of the extremal may not be as good as the performance provided by the global minimum. Philosophically, the overall approach of control design for a performance bound is related to guaranteed cost control ([6]). We note, however, that the bound utilized in [6] is nondifferentiable, which precludes the approach of the present paper.

A further extension of previous approaches considered in the present paper involves the types of feedback loops considered. Specifically, the usual approach to static output feedback involves nonnoisy measurements and weighted controls, while the dual problem involves feeding back noisy measurements to unweighted controls. This situation leads to an additional projection ([2]) which is dual to the projection discussed in $[9,10]$. The inclusion of the dual case now leads to a pair of modified Riccati equations coupled by both the uncertainty bounds and the oblique projections.

In addition to the two types of loops discussed above, one may wish to consider the two remaining cases, namely, feeding back noisy measurements to weighted controls and feeding back nonnoisy measurements to unweighted controls. It is easy to show, however, that the former case
leads to an undefined (i.e., infinite) value for the performance while the latter case is highly singular and thus will not be treated here.

Finally, the scope of the present paper is limited to the development of sufficient conditions for robust decentralized output feedback. Numerical solution of these equations can be carried out by extending available algorithms for centralized output feedback. Numerical algorithms for solving a single modified Riccati equation in the absence of uncertainty bounds are discussed in [10].

## 2. Notation and Definitions

| $\mathbb{R}, \mathbb{R}^{\text {r }}$ • $, \mathbb{R}^{r}, \underline{\mathbb{E}}$ | real numbers, $r \times s$ real numbers, $\mathbb{R}^{r \times 1}$, expectation |
| :---: | :---: |
| $I_{r},()^{T}$ | $r \times r$ identity, transpose |
| $\oplus, \otimes$ | Kronecker sum, Kronecker product ([5]) |
| $\mathbb{S}^{r}$ | $r \times r$ symmetric matrices |
| $\mathrm{N}^{\mathbf{r}}$ | $r \times r$ symmetric nonnegative-definite matrices |
| IP ${ }^{\text {r }}$ | $r \times r$ symmetric positive-definite matrices |
| $Z_{1} \leq Z_{2}$ | $Z_{2}-Z_{1} \in \mathbb{N}^{r}, \quad Z_{1}, Z_{2} \in \mathbb{S}^{r}$ |
| $Z_{1}<Z_{2}$ | $Z_{2}-Z_{1} \in \mathbb{P}^{r}, \quad Z_{1}, Z_{2} \in \mathbb{S}^{r}$ |
| asymptotically stable matrix | matrix with eigenvalues in open left half plane |
| $n, r, s, p$ | positive integers |
| $i, j, k$ | indices, $i=1, \ldots, r, \quad j=1, \ldots, s, \quad k=1, \ldots, p$ |
| $m_{i}, \hat{l}_{i}$ | positive integers, $i=1, \ldots, r$ |
| $\hat{m}_{j}, \ell_{j}$ | positive integers, $j=1, \ldots, s$ |
| $x$ | $n$-dimensional vector |
| $u_{i}, \hat{y}_{i}$ | $m_{i}, \hat{l}_{i}$-dimensional vectors, $i=1, \ldots, r$ |
| $\hat{u}_{j}, y_{j}$ | $\hat{m}_{j}, \ell_{j}$-dimensional vectors, $j=1, \ldots, s$ |
| $A, \Delta A$ | $n \times n$ matrices |
| $B_{i}, \Delta B_{i} ; \hat{C}_{i}$ | $n \times m_{i}$ matrices; $\quad \hat{l}_{i} \times n$ matrices, $i=1, \ldots, r$ |
| $\hat{B}_{j} ; C_{j}, \Delta C_{j}$ | $n \times \hat{m}_{j}$ matrices; $\quad \ell_{j} \times n$ matrices, $j=1, \ldots, s$ |
| $A_{k}$ | $n \times n$ matrices, $k=1, \ldots p$ |
| $B_{i k}$ | $n \times m_{i}$ matrices, $\quad i=1, \ldots, r, \quad k=1, \ldots, p$ |
| $C_{j k}$ | $\ell_{j} \times n$ matrices, $\quad j=1, \ldots, s, \quad k=1, \ldots, p$ |
| $D_{c i}$ | $m_{i} \times \hat{l}_{i}$ matrices, $i=1, \ldots, r$ |
| $E_{c j}$ | $\hat{m}_{j} \times \ell_{j}$ matrices, $j=1, \ldots, s$ |
| $\alpha$ | positive number |
| $\boldsymbol{A}_{\alpha}$ | $A+\frac{9}{3} I_{n}$ |
| $\alpha_{k}$ | positive number, $k=1, \ldots, p$ |
| $\gamma_{k}$ | $\alpha_{k}^{2} / \alpha, \quad k=1, \ldots, p$ |
| $\sigma_{k}$ | real number, $k=1, \ldots, p$ |
| $w_{0}(t), w_{j}(t)$ | $n$-dimensional, $\ell_{j}$-dimensional white noise, $j=1, \ldots$ |

$V_{0}, V_{j}$
$V_{0 j}$
$\boldsymbol{R}_{0}, \boldsymbol{R}_{\mathrm{i}}$
$R_{0 i}$
$\tilde{A}, \tilde{A}_{a}$
$\Delta \tilde{A}$
$\tilde{w}(t)$ $\tilde{R}$
$\tilde{V}$
intensities of $w_{0}, w_{j} ; \quad V_{0} \in \mathbb{N}^{n}, V_{j} \in \mathbb{P}^{\ell_{j}}, j=1, \ldots, s$ $n \times \ell_{j}$ cross intensity of $w_{0}, w_{j}, \quad j=1, \ldots, s$
state and control weightings; $\quad R_{0} \in \mathbb{N}^{n}, R_{i} \in \mathbb{P}^{m_{i}}, \quad i=1, \ldots, r$
$n \times m_{i}$ cross weighting; $\quad R_{0}-R_{0 i} R_{i}^{-1} R_{0 i}^{T} \geq 0, \quad i=1, \ldots, r$
$A+\sum_{i=1}^{r} B_{i} D_{c i} \hat{C}_{i}+\sum_{j=1}^{e} \hat{B}_{j} E_{c j} C_{j} \quad \tilde{A}+\frac{\alpha}{2} I_{n}$
$\Delta A+\sum_{i=1}^{r} \Delta B_{i} D_{c i} \hat{C}_{i}+\sum_{j=1}^{d} \hat{B}_{j} E_{c j} \Delta C_{j}$
$w_{0}(t)+\sum_{j=1}^{e} \hat{B}_{j} E_{c j} w_{j}(t)$
$R_{0}+\sum_{i=1}^{r}\left[R_{0 i} D_{c i} \hat{C}_{i}+\hat{C}_{i}^{T} D_{c i}^{T} R_{0 i}^{T}+\hat{C}_{i}^{T} D_{c i}^{T} R_{i} D_{c i} \hat{C}_{i}\right]$
$V_{0}+\sum_{j=1}^{e}\left[V_{0 j} E_{c j}^{T} \hat{B}_{j}^{T}+\hat{B}_{j} E_{c j} V_{0 j}^{T}+\hat{B}_{j} E_{c j} V_{j} E_{c j}^{T} \hat{B}_{j}^{T}\right]$

For arbitrary $n \times n Q, P$ define:

$$
\begin{aligned}
& R_{a i} \triangleq R_{i}+\sum_{k=1}^{p} \gamma_{k} B_{i k}^{T} P B_{k}, \quad P_{a i} \triangleq B_{i}^{T} P+R_{0 i}^{T}+\sum_{k=1}^{p} \gamma_{k} B_{i k}^{T} P A_{k}, \quad i=1, \ldots, r, \\
& V_{a j} \triangleq V_{j}+\sum_{k=1}^{p} \gamma_{k} C_{j k} Q C_{j k}^{T}, \quad Q_{a j} \triangleq Q C_{j}^{T}+V_{0 j}+\sum_{k=1}^{p} \gamma_{k} A_{k} Q C_{j k}^{T}, \quad j=1, \ldots, s .
\end{aligned}
$$

## 3. Robust Stability and Performance Problem

In this section we state the Robust Stability and Performance Problem along with related notation for later use. Let

$$
U \subset \mathbb{R}^{n \times n} \times \mathbb{R}^{n \times m_{1}} \times \ldots \times \mathbb{R}^{n \times m_{r}} \times \mathbb{R}^{\ell_{1} \times n} \times \ldots \times \mathbb{R}^{L_{0} \times n}
$$

denote the set of uncertain perturbations $\left(\Delta A, \Delta B_{1}, \ldots, \Delta B_{r}, \Delta C_{1}, \ldots, \Delta C_{s}\right)$ of the nominal system matrices $A, B_{1}, \ldots, B_{r}, C_{1}, \ldots, C_{s}$.

Robust Stability and Performance Problem. Determine ( $D_{c 1}, \ldots, D_{c r}, E_{c 1}, \ldots, E_{c s}$ ) such that the closed-loop system consisting of the $n$ th-order controlled and disturbed plant

$$
\begin{equation*}
\dot{x}(t)=(A+\Delta A) x(t)+\sum_{i=1}^{r}\left(B_{i}+\Delta B_{i}\right) u_{i}(t)+\sum_{j=1}^{n} \hat{B}_{j} \hat{u}_{j}(t)+w_{0}(t), \quad t \in[0, \infty) \tag{3.1}
\end{equation*}
$$

nonnoisy and noisy measurements

$$
\begin{gather*}
\hat{y}_{i}(t)=\hat{C}_{i} x(t), \quad i=1, \ldots, r  \tag{3.2}\\
y_{j}(t)=\left(C_{j}+\Delta C_{j}\right) x(t)+w_{j}(t), \quad j=1, \ldots, s, \tag{3.3}
\end{gather*}
$$

and static output feedback controller

$$
\begin{equation*}
u_{i}(t)=D_{\mathrm{c} i} \hat{y}_{i}(t), \quad i=1, \ldots, r \tag{3.4}
\end{equation*}
$$

$$
\begin{equation*}
\hat{u}_{j}(t)=E_{c j} y_{j}(t), \quad j=1, \ldots, s \tag{3.5}
\end{equation*}
$$

is asymptotically stable for all variations in $U$ and the performance criterion

$$
\begin{align*}
& J\left(D_{c 1}, \ldots, D_{c r}, E_{c 1}, \ldots, E_{\mathrm{ct}}\right) \triangleq \\
& \quad \sup _{u} \limsup _{t \rightarrow \infty} \mathbb{E}\left[x^{T}(t) R_{0} x(t)+2 \sum_{i=1}^{r} x^{T}(t) R_{0 i} u_{i}(t)+\sum_{i=1}^{r} u_{i}^{T}(t) R_{i} u_{i}(t)\right] \tag{3.6}
\end{align*}
$$

is minimized.
For each controller ( $D_{c 1}, \ldots, D_{c r}, E_{c 1}, \ldots, E_{c \varepsilon}$ ) and variation in $U$, the closed-loop system (3.1)(3.5) is given by

$$
\begin{equation*}
\dot{x}(t)=(\tilde{A}+\Delta \tilde{A}) x(t)+\tilde{w}(t), \quad t \in[0, \infty) \tag{3.7}
\end{equation*}
$$

where $\tilde{w}(t)$ is white noise with intensity $\tilde{V} \in \mathbb{N}^{n}$.
Remark 3.1. In the case $\Delta A, \Delta B_{i}, \Delta C_{j}=0$ it is well known that stabilizability is related to the existence of fixed modes ([11]). When plant uncertainties are present the problem is, of course, far more complex. In the present paper sufficient conditions for robust stability are obtained as a consequence of the existence of robust performance bounds.

Remark 3.2. Note that the controller architecture is quite general in that it includes two distinctly different types of decentralized loops. The first type, indexed by $i=1, \ldots, r$, involves feeding back nonnoisy measurements to weighted controls. This is the standard setting in the optimal output-feedback literature ( $[8-10]$ ). In addition, we include the dual situation, indexed by $j=1, \ldots, s$, which involves feeding back noisy measurements to unweighted controls. The case in which only one type of loop is present can be formally recovered from our results by ignoring $B_{i}$ and $\hat{C}_{i}$ or $\hat{B}_{j}$ and $C_{j}$ as required. As noted in Section 1 , noisy measurements cannot be fed back to weighted controls via static control, while feeding back nonnoisy measurements to unweighted controls is a singular problem.

Remark 3.3. Note that the problem statement is restrictive in the sense that uncertainties in both the control and observation matrices are not permitted within the same feedback loop. Although it is indeed possible to permit such simultaneous uncertainties, the development is considerably more complex and hence is not treated here.

Remark 3.4. The cost functional (3.6) is identical to the LQG criterion (usually stated in terms of an averaged integral) with the exception of the supremum for evaluating worst case over $u$.

## 4. Sufficient Conditions for Robust Stability and Performance

In practice, steady-state performance is only of interest when the closed-loop system (3.7) is stable over $U$. The following result, which expresses the performance in terms of the state covariance, is immediate.

Lemma 4.1. Let ( $D_{c 1}, \ldots, D_{c r}, E_{c 1}, \ldots, E_{c s}$ ) be given and suppose the system (3.7) is stable for all variations in $U$. Then

$$
\begin{equation*}
J\left(D_{c 1}, \ldots, D_{c r}, E_{c 1}, \ldots, E_{c s}\right)=\sup _{u} \operatorname{tr} Q_{\Delta \tilde{A}} \tilde{R} \tag{4.1}
\end{equation*}
$$

where $Q_{\Delta \hat{A}} \triangleq \lim _{t \rightarrow \infty} \mathbb{E}\left[x(t) x^{T}(t)\right] \in \mathbb{N}^{n}$ is the unique solution to

$$
\begin{equation*}
0=(\tilde{A}+\Delta \tilde{A}) Q_{\Delta \tilde{A}}+Q_{\Delta \tilde{A}}(\tilde{A}+\Delta \tilde{A})^{T}+\tilde{V} \tag{4.2}
\end{equation*}
$$

Remark 4.1. When $U$ is compact, "sup" in (4.1) can be replaced by "max".
We now seek upper bounds for $J\left(D_{c 1}, \ldots, D_{c r}, E_{c 1}, \ldots, E_{c \iota}\right)$. Our assumptions allow us to obtain robust stability as a consequence of robust performance.

Theorem 4.1. Let $\Omega: \mathbb{N}^{n} \times \mathbb{R}^{m_{1} \times \hat{\ell}_{1}} \times \ldots \mathbb{R}^{m_{r} \times \hat{L}_{1}} \times \mathbb{R}^{\boldsymbol{m}_{1} \times \boldsymbol{\ell}_{1}} \times \ldots \times \mathbb{R}^{\boldsymbol{m}_{1} \times \boldsymbol{C}_{0}} \rightarrow \mathbb{S}^{n}$ be such that

$$
\begin{gather*}
\Delta \tilde{A} Q+Q \Delta \tilde{A}^{T} \leq \Omega\left(Q, D_{c 1}, \ldots, D_{c r}, E_{c 1}, \ldots, E_{c \epsilon}\right), \\
\left(\Delta A, \Delta B_{1}, \ldots, \Delta B_{r}, \Delta C_{1}, \ldots, \Delta C_{\bullet}\right) \in U \\
\left(Q, D_{c 1}, \ldots, D_{c r}, E_{c 1}, \ldots, E_{c c}\right) \in \mathbb{N}^{n} \times \mathbb{R}^{m_{1} \times \mathcal{I}_{1}} \times \ldots \times \mathbb{R}^{m_{r} \times \boldsymbol{亡}} \times \mathbb{R}^{m_{1} \times \ell_{1}} \times \ldots \times \mathbb{R}^{m_{0} \times \ell_{\bullet}} . \tag{4.3}
\end{gather*}
$$

Furthermore, for given ( $D_{c 1}, \ldots, D_{c r}, E_{c 1}, \ldots, E_{c s}$ ) suppose there exists $Q \in \mathbb{N}^{n}$ satisfying

$$
\begin{equation*}
0=\tilde{A} Q+Q \tilde{A}^{T}+\Omega\left(Q, D_{c 1}, \ldots, D_{\mathrm{cr}}, E_{\mathrm{c} 1}, \ldots, E_{\mathrm{c} \boldsymbol{r}}\right)+\tilde{V} . \tag{4.4}
\end{equation*}
$$

Then the pair $\left(\tilde{A}+\Delta \tilde{A}, \tilde{V}^{\frac{1}{2}}\right)$ is stabilizable for all variations in $U$ if and only if $\tilde{A}+\Delta \tilde{A}$ is asymptotically stable for all variations in $U$. In this case,

$$
\begin{equation*}
Q_{\Delta i} \leq Q, \tag{4.5}
\end{equation*}
$$

where $Q_{\Delta A^{j}}$ satisfies (4.2), and

$$
\begin{equation*}
J\left(D_{c 1}, \ldots, D_{c r}, E_{c 1}, \ldots, E_{c t}\right) \leq \operatorname{tr} Q \tilde{R} . \tag{4.6}
\end{equation*}
$$

Proof. For all variations in $U$, (4.4) is equivalent to

$$
\begin{equation*}
0=(\tilde{A}+\Delta \tilde{A}) Q+Q(\tilde{A}+\Delta \tilde{A})^{T}+\Phi\left(Q, D_{c 1}, \ldots, D_{c r}, E_{c 1}, \ldots, E_{c \theta}, \Delta \tilde{A}\right)+\tilde{V} \tag{4.7}
\end{equation*}
$$

where

$$
\Phi\left(Q, D_{c 1}, \ldots, D_{c r}, E_{c 1}, \ldots, E_{c \theta}, \Delta \tilde{A}\right) \triangleq \Omega\left(Q, D_{c 1}, \ldots, D_{c r}, E_{c 1}, \ldots, E_{c я}\right)-\left(\Delta \tilde{A} Q+Q \Delta \tilde{A}^{T}\right)
$$

Note that by (4.3), $\Phi(\cdot) \geq 0$ for all variations in $U$. If $\left(\tilde{A}+\Delta \tilde{A}, \tilde{V}^{\frac{1}{2}}\right)$ is stabilizable for all variations in $U$, it follows from Theorem 3.6 of $[12]$ that $\left(\tilde{A}+\Delta \tilde{A},\left[\tilde{V}+\Phi\left(Q, D_{c 1}, \ldots, D_{c r}, E_{c 1}, \ldots, E_{c s}, \Delta \tilde{A}\right)\right]^{\frac{1}{2}}\right)$ is stabilizable for all variations in $U$. Hence Lemma 12.2 of [12] implies $\bar{A}+\Delta \tilde{A}$ is asymptotically stable for all variations in $U$. The converse is immediate. Next, substracting (4.2) from (4.7) yields

$$
0=(\tilde{A}+\Delta \tilde{A})\left(Q-Q_{\Delta \tilde{A}}\right)+\left(Q-Q_{\Delta \tilde{A}}\right)(\tilde{A}+\Delta \tilde{A})^{T}+\Phi\left(Q, D_{c 1}, \ldots, D_{c r}, E_{c 1}, \ldots, E_{c e}, \Delta \tilde{A}\right)
$$

or, equivalently, (since $\tilde{A}+\Delta \tilde{A}$ is asymptotically stable)

$$
Q-Q_{\Delta \tilde{A}}=\int_{0}^{\infty} e^{(\tilde{A}+\Delta \tilde{A}) t} \Phi\left(Q, D_{c 1}, \ldots, D_{c r}, E_{c 1}, \ldots, E_{c \iota}, \Delta \tilde{A}\right) e^{(\tilde{A}+\Delta \tilde{A})^{T_{t}}} d t \geq 0
$$

which implies (4.5). Finally, (4.5) and (4.1) yield (4.6).
Remark 4.2. If $\tilde{V}$ is positive definite then the stabilizability hypothesis of Theorem 4.1 is automatically satisfied for all variations in $U$.

## 5. Uncertainty Structure and the Quadratic Lyapunov Bound

The uncertainty set $U$ is assumed to be of the form

$$
\begin{align*}
& U=\left\{\left(\Delta A, \Delta B_{1}, \ldots, \Delta B_{r}, \Delta C_{1}, \ldots, \Delta C_{s}\right):\right. \\
& \Delta A=\sum_{k=1}^{p} \sigma_{k} A_{k}, \Delta B_{i}=\sum_{k=1}^{p} \sigma_{k} B_{i k}, \quad i=1, \ldots, r,  \tag{5.1}\\
& \left.\Delta C_{j}=\sum_{k=1}^{p} \sigma_{k} C_{j k}, \quad j=1, \ldots, s, \quad \sum_{k=1}^{p} \sigma_{k}^{2} / \alpha_{k}^{2} \leq 1\right\},
\end{align*}
$$

where, for $k=1, \ldots, p:\left(A_{k}, B_{1 k}, \ldots, B_{r k}, C_{1 k}, \ldots, C_{s k}\right)$ are fixed matrices denoting the structure of the parametric uncertainty; $\alpha_{k}$ is a given uncertainty bound; and $\sigma_{k}$ is an uncertain parameter. Note that the uncertain parameters $\sigma_{k}$ are assumed to lie in a specified ellipsoidal region in $\mathbb{R}^{p}$. The closed-loop system thus has structured uncertainty of the form

$$
\begin{equation*}
\Delta \tilde{A}=\sum_{k=1}^{p} \sigma_{k} \tilde{A}_{k}, \tag{5.2}
\end{equation*}
$$

where

$$
\begin{equation*}
\tilde{A}_{k} \triangleq A_{k}+\sum_{i=1}^{r} B_{i k} D_{c i} \hat{C}_{i}+\sum_{j=1}^{\dot{B}} \hat{B}_{j} E_{c j} C_{j k}, \quad k=1, \ldots, p \tag{5.3}
\end{equation*}
$$

To obtain explicit gain expressions for ( $D_{c 1}, \ldots, D_{\mathrm{cr}}, E_{\mathrm{c} 1}, \ldots, E_{\mathrm{ca}}$ ) we assume that, for each $k \in\{1, \ldots, p\}$, at most one of the matrices $B_{1 k}, \ldots, B_{r k}, C_{1 k}, \ldots, C_{s k}$ is nonzero. Note that this assumption does not preclude the treatment of uncertainties in the input and output matrices. It requires only that such uncertainties be modeled as uncorrelated.

Given the structure of $U$ defined by (5.1), the bound $\Omega$ satisfying (4.3) can now be specified. In the following result $Q$ denotes an arbitrary element of $\mathbb{N}^{n}$, not necessarily a solution of (4.4).

Proposition 5.1. Let $\alpha$ be an arbitrary positive scalar. Then the function

$$
\begin{equation*}
\Omega\left(Q, D_{\mathrm{c} 1}, \ldots, D_{\mathrm{cr}}, E_{\mathrm{c} 1}, \ldots, E_{\mathrm{c} \boldsymbol{\prime}}\right)=\alpha Q+\alpha^{-1} \sum_{k=1}^{p} \alpha_{k}^{2} \tilde{A}_{k} Q \tilde{A}_{k}^{T} \tag{5.4}
\end{equation*}
$$

satisfies (4.3) with $U$ given by (5.1).
Proof. Note that

$$
\begin{aligned}
0 & \leq \sum_{k=1}^{p}\left[\left(\alpha^{\frac{1}{2}} \sigma_{k} / \alpha_{k}\right) I_{n}-\left(\alpha_{k} / \alpha^{\frac{1}{2}}\right) \tilde{A}_{k}\right] Q\left[\left(\alpha^{\frac{1}{2}} \sigma_{k} / \alpha_{k}\right) I_{n}-\left(\alpha_{k} / \alpha^{\frac{1}{2}}\right) \tilde{A}_{k}\right]^{T} \\
& =\alpha \sum_{k=1}^{p}\left(\sigma_{k}^{2} / \alpha_{k}^{2}\right) Q+\alpha^{-1} \sum_{k=1}^{p} \alpha_{k}^{2} \tilde{A}_{k} Q \tilde{A}_{k}^{T}-\sum_{k=1}^{p} \sigma_{k}\left(\tilde{A}_{k} Q+Q \tilde{A}_{k}^{T}\right)
\end{aligned}
$$

which yields (4.3).
Remark 5.1. Note that the bound $\Omega$ given by (5.4) consists of two distinct terms. The first term $\alpha Q$ can be thought of as arising from an exponential time weighting of the cost, or, equivalently, from a uniform right shift of the open-loop dynamics ([1]). The second term $\alpha^{-1} \sum_{k=1}^{p} \alpha_{k}^{2} \tilde{A}_{k} Q \tilde{A}_{k}^{T}$ arises naturally from a multiplicative white noise model ([3,4]). Such interpretations have no bearing on the results obtained here since only the bound $\Omega$ defined by (5.4) is required. Note that the bound is valid for all positive $\alpha$. .

Remark 5.2. The conservatism of the bound (5.4) is difficult to predict for two reasons. First, the overbounding (4.3) holds with respect to the partial ordering of the nonnegative-definite matrices for which no scalar measure of conservatism is available. And, second, the bound (4.3) is required to hold for all nonnegative-definite matrices $Q$ and feedback gains ( $D_{c 1}, \ldots, D_{c r}, E_{c 1}, \ldots, E_{c g}$ ). The conservatism will thus depend upon the actual values of $Q, D_{c 1}, \ldots, D_{c r}, E_{c 1}, \ldots, E_{c,}$ determined by solving (4.4).

## 6. The Auxiliary Minimization Problem and Necessary Conditions for Optimality

Rather than minimizing the actual cost (3.6), we shall consider the upper bound (4.6). This leads to the following problem.

Auxiliary Minimization Problem. Determine ( $Q, D_{c 1}, \ldots, D_{c r}, E_{c 1}, \ldots, E_{c s}$ ) which minimizes

$$
\begin{equation*}
J\left(Q, D_{c 1}, \ldots, D_{c r}, E_{c 1}, \ldots, E_{c \theta}\right) \triangleq \operatorname{tr} Q \tilde{R} \tag{6.1}
\end{equation*}
$$

subject to

$$
\begin{gather*}
Q \in \mathbb{N}^{n},  \tag{6.2}\\
0=\tilde{A}_{\alpha} Q+Q \tilde{A}_{\alpha}^{T}+\sum_{k=1}^{p} \gamma_{k} \tilde{A}_{k} Q \tilde{A}_{k}^{T}+\tilde{V} \tag{6.3}
\end{gather*}
$$

The relationship between the Auxiliary Minimization Problem and the Robust Stability and Performance Problem is straightforward as shown by the following observation.

Proposition 6.1. Suppose ( $Q, D_{c 1}, \ldots, D_{c r}, E_{c 1}, \ldots, E_{c s}$ ) satisfies (6.2)-(6.4). Then

$$
\begin{equation*}
\left(\tilde{A}+\Delta \tilde{A}, \tilde{V}^{\frac{1}{2}}\right) \text { is stabilizable for all variations in } U \tag{6.4}
\end{equation*}
$$

if and only if $\tilde{A}+\Delta \tilde{A}$ is asymptotically stable for all variations in $U$. In this case,

$$
\begin{equation*}
J\left(D_{c 1}, \ldots, D_{c r}, E_{c 1}, \ldots, E_{c \rho}\right) \leq J\left(Q, D_{c 1}, \ldots, D_{c r}, E_{c 1}, \ldots, E_{c s}\right) \tag{6.5}
\end{equation*}
$$

Proof. With $\Omega$ given by (5.4), Proposition 5.1 implies that (4.3) is satisfied. Since the hypotheses of Theorem 4.1 are satisfied, robust stability with performance bound (4.6) is guaranteed. Note that with definition (6.1), (6.5) is merely a restatement of (4.6).

The derivation of the necessary conditions for the Auxiliary Minimization Problem is based upon the Fritz John form of the Lagrange multiplier theorem.* Rigorous application of this technique requires that $\left(Q, D_{c 1}, \ldots, D_{c r}, E_{c 1}, \ldots, E_{c a}\right)$ be restricted to the open set

$$
S \triangleq\left\{\left(Q, D_{c 1}, \ldots, D_{c r}, E_{c 1}, \ldots, E_{c a}\right): \quad Q \in \mathbb{P}^{n} \text { and } \ell \text { is asymptotically stable }\right\}
$$

[^12]where
$$
A \triangleq \tilde{A}_{\alpha} \oplus \tilde{A}_{\alpha}+\sum_{k=1}^{p} \gamma_{k} \tilde{A}_{k} \otimes \tilde{A}_{k}
$$

The requirement $\left(Q, D_{c 1}, \ldots, D_{c r}, E_{c 1}, \ldots, E_{c s}\right) \in S$ implies that $Q$ and its nonnegative-definite dual $P$ are unique solutions to the modified Lyapunov equations (6.3) and

$$
\begin{equation*}
0=\tilde{A}_{\alpha}^{T} P+P \tilde{A}_{\alpha}+\sum_{k=1}^{p} \gamma_{k} \tilde{A}_{k}^{T} P \tilde{A}_{k}+\tilde{R} . \tag{6.6}
\end{equation*}
$$

An additional technical requirement is that $\left(Q, D_{c 1}, \ldots, D_{c r}, E_{c 1}, \ldots, E_{c a}\right)$ be confined to the set

$$
\begin{aligned}
& S^{+} \triangleq\left\{\left(Q, D_{\mathrm{c} 1}, \ldots, D_{c r}, E_{c 1}, \ldots, E_{c \beta}\right) \in S: \quad \hat{C}_{i} Q \hat{C}_{i}^{T}>0, \quad i=1, \ldots, r\right. \\
&\text { and } \left.\hat{B}_{j}^{T} P \hat{B}_{j}>0, \quad j=1, \ldots, s\right\} .
\end{aligned}
$$

The positive definiteness conditions in the definition of $S^{+}$hold when $\hat{C}_{i}$ and $\hat{B}_{j}$ have full row and column rank, respectively, and $Q$ and $P$ are positive definite. As can be seen from the proof of Theorem 6.1 these conditions imply the existence of the projections $\nu_{i}$ and $\hat{\nu}_{j}$ corresponding to the two distinct types of feedback loops. Note that $S^{+}$is open.

Remark 6.1. As pointed out in Remark 3.1, the set $S$ may be empty in which case, of course, our results do not apply. As will be seen, however, our approach does not require explicit verification that $S$ be nonempty since robust stability is obtained as a consequence of robust performance.

Remark 6.2. As will be seen, the constraint $\left(Q, D_{c 1}, \ldots, D_{c r}, E_{c 1}, \ldots, E_{c a}\right) \in S$ need not be verified in practice and is not required for either robust stability or robust performance since Proposition 6.1 shows that only (6.2)-(6.4) are needed. Rather, the set $S$ constitutes sufficient conditions under which the Lagrange multiplier technique is applicable to the Auxiliary Minimization Problem. Specifically, the condition $Q>0$ replaces (6.2) by an open set constraint, while the asymptotic stability of $\ell$ serves as a normality condition which further implies that the dual $P$ of $Q$ is nonnegative definite.

Necessary conditions for the Auxiliary Minimization Problem can now be obtained.
Theorem 6.1. If $\left(Q, D_{c 1}, \ldots, D_{c r}, E_{c 1}, \ldots, E_{c \varepsilon}\right) \in S^{+}$solves the Auxiliary Minimization Problem with $U$ given by (5.1), then there exist $Q, P \in \mathbb{N}^{n}$ such that $D_{c 1}, \ldots, D_{c r}, E_{c 1}, \ldots, E_{c s}$ are given by

$$
\begin{align*}
& D_{c i}=-R_{a i}^{-1} P_{a i} Q \hat{C}_{i}^{T}\left(\hat{C}_{i} Q \hat{C}_{i}^{T}\right)^{-1}, \quad i=1, \ldots, r,  \tag{6.7}\\
& E_{c j}=-\left(\hat{B}_{j}^{T} P \hat{B}_{j}\right)^{-1} \hat{B}_{j}^{T} P Q_{a j} V_{a j}^{-1}, \quad j=1, \ldots, s, \tag{6.8}
\end{align*}
$$

and such that $Q, P$ satisfy

$$
\begin{align*}
& 0=\left(A_{\alpha}-\sum_{i=1}^{r} B_{i} R_{a i}^{-1} P_{a i} \nu_{i}\right) Q+Q\left(A_{\alpha}-\sum_{i=1}^{r} B_{i} R_{a i}^{-1} P_{a i} \nu_{i}\right)^{T}+V_{0} \\
&+\sum_{k=1}^{p} \gamma_{k}\left(A_{k}-\sum_{i=1}^{r} B_{i k} R_{a i}^{-1} P_{a i} \nu_{i}\right) Q\left(A_{k}-\sum_{i=1}^{r} B_{i k} R_{a i}^{-1} P_{a i} \nu_{i}\right)^{T} \\
&-\sum_{j=1}^{\dot{L}} Q_{a j} V_{a j}^{-1} Q_{a j}^{T}+\sum_{j=1}^{\dot{\nu}} \hat{\nu}_{j \perp} Q_{a j} V_{a j}^{-1} Q_{a j}^{T} \hat{\nu}_{j \perp}^{T},  \tag{6.9}\\
& 0=\left(A_{\alpha}-\sum_{j=1}^{\dot{L}} \hat{\nu}_{j} Q_{a j} V_{a j}^{-1} C_{j}\right)^{T} P+P\left(A_{\alpha}-\sum_{j=1}^{\dot{\nu}} \hat{\nu}_{j} Q_{a j} V_{a j}^{-1} C_{j}\right)+R_{0} \\
&+\sum_{k=1}^{p} \gamma_{k}\left(A_{k}-\sum_{j=1}^{\dot{\nu}} \hat{\nu}_{j} Q_{a j} V_{a j}^{-1} C_{j k}\right)^{T} P\left(A_{k}-\sum_{j=1}^{\dot{\nu}} \hat{\nu}_{j} Q_{a j} V_{a j}^{-1} C_{j k}\right) \\
&-\sum_{i=1}^{r} P_{a i}^{T} R_{a i}^{-1} P_{a i}+\sum_{i=1}^{r} \nu_{i \perp}^{T} P_{a i}^{T} R_{a i}^{-1} P_{a i} \nu_{i \perp},  \tag{6.10}\\
&  \tag{6.11}\\
& \nu_{i} \triangleq Q \hat{C}_{i}^{T}\left(\hat{C}_{i} Q \hat{C}_{i}^{T}\right)^{-1} \hat{C}_{i}, \quad \nu_{i \perp} \triangleq I_{n}-\nu_{i}, \quad i=1, \ldots, r,  \tag{6.12}\\
& \hat{\nu}_{j} \triangleq \hat{B}_{j}\left(\hat{B}_{j}^{T} P \hat{B}_{j}\right)^{-1} \hat{B}_{j}^{T} P, \quad \hat{\nu}_{j \perp} \triangleq I_{n}-\hat{\nu}_{j}, \quad j=1, \ldots, s .
\end{align*}
$$

Furthermore, the auxiliary cost is given by

$$
\begin{align*}
& J\left(Q, D_{c 1}, \ldots, D_{c r}, E_{c 1}, \ldots, E_{c i}\right)= \\
& \quad \operatorname{tr}\left[Q\left(R_{0}+\sum_{i=1}^{p} \nu_{i}^{T} P_{a i}^{T} R_{a i}^{-1} R_{i} R_{a i}^{-1} P_{a i} \nu_{i}-2 R_{0 i} R_{a i}^{-1} P_{a i} \nu_{i}\right)\right] . \tag{6.13}
\end{align*}
$$

Conversely, if there exist $Q, P \in \mathbb{N}^{n}$ satisfying (6.9) and (6.10) then $Q$ satisfies (6.3) with ( $D_{c 1}, \ldots, D_{c r}, E_{c 1}, \ldots, E_{c a}$ ) given by (6.7) and (6.8), and $J\left(Q, D_{c 1}, \ldots, D_{c r}, E_{c 1}, \ldots, E_{c 0}\right)$ is given by (6.13).

Proof. To optimize (6.1) over the open set $S^{+}$, subject to the constraint (6.3), form the Lagrangian

$$
\mathcal{L}\left(Q, D_{c 1}, \ldots, D_{c r}, E_{c 1}, \ldots, E_{c e}\right) \triangleq \operatorname{tr}\left[\lambda Q \tilde{R}+\left(\tilde{A} Q+Q \tilde{A}^{T}+\sum_{k=1}^{p} \gamma_{k} \tilde{A}_{k} Q \tilde{A}_{k}^{T}+\tilde{V}\right) P\right]
$$

where the Lagrange multipliers $\lambda \geq 0$ and $P \in \mathbb{R}^{n \times n}$ are not both zero. Setting $\partial \mathcal{L} / \partial Q=0, \lambda=0$ implies $P=0$ since $\Omega$ is asymptotically stable. Hence, without loss of generality set $\lambda=1$. Thus the stationarity conditions are given by

$$
\begin{equation*}
\frac{\partial \mathcal{L}}{\partial Q}=\tilde{A}_{\alpha}^{T} P+P \tilde{A}_{\alpha}+\sum_{k=1}^{p} \gamma_{k} \tilde{A}_{k}^{T} P \tilde{A}_{k}+\tilde{R}=0 \tag{6.14}
\end{equation*}
$$

$$
\begin{array}{ll}
\frac{\partial \mathcal{L}}{\partial D_{c i}}=R_{a i} D_{c i} \hat{C}_{i} Q \hat{C}_{i}^{T}+P_{a i} Q \hat{C}_{i}^{T}=0, & i=1, \ldots, r, \\
\frac{\partial \mathcal{L}}{\partial E_{c j}}=\hat{B}_{j}^{T} P \hat{B}_{j} E_{c j} V_{a j}+\hat{B}_{j}^{T} P Q_{a j}=0, & j=1, \ldots, s . \tag{6.16}
\end{array}
$$

Since $\left(Q, D_{c 1}, \ldots, D_{c r}, E_{c 1}, \ldots, E_{c a}\right) \in S^{+}, \hat{C}_{i} Q \hat{C}_{i}^{T}$ and $\hat{B}_{j}^{T} P \hat{B}_{j}$ are invertible and hence (6.15) and (6.16) imply (6.7) and (6.8). Finally, (6.9) and (6.10) are equivalent to (6.3) and (6.6).

Remark 6.3. Several special cases can be recovered formally from Theorem 6.1. For example, when the control weighting is nonsingular and the measurement noise is zero, i.e., when $\hat{u}_{i}$ and $y_{i}$ are absent for $i=1, \ldots, r$, delete (6.8) and set $\hat{\nu}_{j}=0$ in (6.9). In this case the last two terms in (6.9) can be deleted. Deleting also the uncertainty terms $A_{k}, B_{i k}, C_{j k}$ yields the results of [10] with the added features of correlated plant/measurement noise ( $V_{0 j}$ ) and cross weighting ( $R_{0 i}$ ). Furthermore, assuming a centralized structure for the static controller, i.e., $r=1$, yields the usual static output feedback result ( $[8,9]$ ).

## 7. Sufficient Conditions for Robust Stability and Performance

We now combine Proposition 6.1 and Theorem 6.1 to obtain sufficient conditions for robust stability and performance.

Theorem 7.1. Suppose there exist $Q, P \in \mathbb{N}^{n}$ satisfying (6.9) and (6.10). Then with $\left(D_{c 1}, \ldots, D_{c r}, E_{c 1}, \ldots, E_{c a}\right)$ given by (6.6) and (6.7), ( $\left.\tilde{A}+\Delta \tilde{A}, \tilde{V}^{\frac{1}{2}},\right)$ is stabilizable for all variations in $U$ if and only if $\tilde{A}+\Delta \tilde{A}$ is asymptotically stable for all variations in $U$. In this case the performance of the closed-loop system satisfies the bound

$$
\begin{equation*}
J\left(D_{c 1}, \ldots, D_{c r}, E_{c 1}, \ldots, E_{c a}\right) \leq \operatorname{tr}\left[Q\left(R_{0}+\sum_{i=1}^{p} \nu_{i}^{T} P_{g i}^{T} R_{a i}^{-1} R_{i} R_{a i}^{-1} P_{a i} \nu_{i}-2 R_{0 i} R_{a i}^{-1} P_{a i} \nu_{i}\right)\right] \tag{7.1}
\end{equation*}
$$

Proof. The converse of Theorem 6.1 shows that $Q$ satisfies (6.3) with ( $D_{c 1}, \ldots, D_{c r}$, $E_{c 1}, \ldots, E_{c e}$ ) given by (6.7) and (6.8). Hence, with the stabilizability assumption (6.4), Proposition 6.1 implies robust stability and performance.

Remark 7.1. The application of Theorem 7.1 in practice requires 1 ) numerical solution of (6.9) and (6.10), and 2) verification of the stabilizability hypothesis. No other assumptions need be verified in applying this result.

## 8. Concluding Remarks

We have developed a theory of robust decentralized output feedback via static control. The development permits the treatment of noisy and nonnoisy measurements, weighted and unweighted controls, and structured real-valued parameter uncertainties in the plant matrices. The theory provides a robustification of results given in [8-10] for both centralized and decentralized optimal output feedback. The theory is constructive in nature rather than existential. Specifically, the main result, Theorem 7.1, involves a coupled pair of modified Riccati equations (6.9), (6.10) whose solutions, when they exist, are used to explicitly construct feedback gains (6.7), (6.8) which are guaranteed to provide both robust stability and performance. Future research is required for evaluating the conservativeness of the theory. The numerical algorithms developed in [10] provide a starting point in this regard.

## References

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inverse of the noise intensity matrix. Hence, it is not surprising that a sizable body of literature has been devurad to the singular measurennent noise problem in both continuous and discrete time [2]-[14]. For an overview of stochastic observer theory, see [15].

Much of the continuous-time singular extimation literature ettermpts to overcome the noise singularity by introducing new measurements obtained by differentiating noise-free measurements. The present note complements these results in the following way. For the available noisy and noise-free measurements we simultancously design a reduced-order dynamic estimator for the noisy measurements and a static estimator for the noise-free measurements. We are not concerned here with the question of how the measurements are generated (e.g., vin successive differentiation). Rather, cur goal is to develop a unified dynamic/static estimation design theory which permits full utilization of both noisy and noise-free measurements. Application of these results to previously proposed approaches to singular estimation involving differentiation and transformation should be an interesting area for future reseasch.

The results given herein directly generalize the results obtained in [1]. Specifically, the modified Riccati/Lyapunov equations are now coupled by a pair of oblique projections. As in [1] the reguirement for reduced estimator order gives rise to the projection

$$
\begin{equation*}
T_{2}=Q P(Q P)^{0} \tag{1.1}
\end{equation*}
$$

where ()' denotes group generalized inverse and $\boldsymbol{Q}$ and $\boldsymbol{P}$ are rankdeficient nonnegative-definite matrices analogous to the controllability and observability Gramians of the estimator. In addition, the presence of noise-free measurements

$$
\begin{equation*}
y(t)=C_{2} x(t) \tag{1.2}
\end{equation*}
$$

leads to the projection

$$
\begin{equation*}
T_{1}=Q C_{2}^{\top}\left(C_{2} Q C_{2}^{\top}\right)^{-1} C_{2} \tag{1.3}
\end{equation*}
$$

where $Q$ is the steady-state error covariance. The contribution of the present note is a concise, unified statement of the optimality conditions in a form which clearly displays the role of the oblique projections $r_{1}$ and $r_{2}$ in explicitly characterizing oprimal static/dynamic (nonstrictly proper) estimators. An additional feature of the present note is the presence of state- and measurement-dependent white noise in the plant model. This model has been studied in a state-estimator context in [16]-[18] and has been justified as an approach to robustness in [19]-[22].

In Section III of the note, we consider the case in which the noisy and noise-free measurements are fed to the dynamic and static estimators, respectively. In Section IV, we note that feeding the noisy measurements to the static estimator results in an ill-posed problem, and we consider the general case in which the noise-free measurements are fed to both the static and dynamic estimators. Optimality conditions now lead to the interesting disjointness condition

$$
\begin{equation*}
0=r_{2} r_{1} \tag{1.4}
\end{equation*}
$$

concerning the relationship between the static and dynamic estimators. The meaning of (1.4) for proposed singular estimation schemes will be explored in future papers.

The goal of this note is confined to a rigorous development of necessary conditions for the optimal estimation problem. In support of this aim it should be noted that the usefulness of necessary conditions in optimization and optimal control has been amply demonstrated by classical resules such as the maximum principle and Euler-Lagrange theory. For practical purposes, necessary conditions are largely free from restrictive apecial assumptions which invariably accompeny sufficiency theory. Moet importantly, success in addressing the problems of existence, sufficiency and global optimality is far more likely after the full elucidation of the necessary conditions has been echieved. Indeed, sufficiency conditions are often obtained by strengthening neceseary conditions by means of additional restrictive assumptions.

Even without a complete resolution of quemions pertaining to eximence
and sufficiency, the necessury conditions fulfill several immedinte needs. Specifically, the atructure of these conditions provides insight into the properties of the solution arising from optimality considerations. This has been demonstrated for the closely related problem of reduced-order modeling for which local minima are characterized in terms of an eigensystem decomposition [23]. Potentially more useful than insight for practical applications are prospects for constructing novel computational algorithms which avoid traditional gradient search methods. Thus far, two distinct algorithms have been developed, namely, an iterative method which exploits the structure of the oblique projection [23] and a homotopy algorithm which eliminates the peed for eigensystem calculations and provides the means for attaining global optimality [24]. For computational purposes it should also be noted that under an existence assumption the necessary conditions are guaranteed to possess a solution to the problem, while sufficient conditions may fail in this regard.

## it. Problem statement and Mann theorem

## Reduced-Order State-Estimation Problem

Given the nth-order observed system

$$
\begin{gather*}
x(t)=\left(A+\sum_{i=1}^{f} w_{1}(t) A_{i}\right) x(t)+w_{0}(t)  \tag{3.1}\\
y_{1}(t)=\left(C_{1}+\sum_{i=1}^{t} v_{i}(t) C_{11}\right) x(t)+w_{1}(t),  \tag{3.2}\\
y_{2}(t)=C_{2} x(t) \tag{3.3}
\end{gather*}
$$

|  |  |
| :---: | :---: |
|  | $I_{n},()^{T},()^{\prime \prime}$ |
|  | $\oplus, \otimes$ |
|  | tr 2 |
|  | ${ }^{+}$ |
|  | $n_{1} l_{1}, l_{2}, n_{e}, p, q$ |
|  | n |
|  | $x, x_{0}$ |
|  | $y_{1}, y_{2}, y_{0}$ |
|  | $A_{1}, A_{i} ; C_{1}, C_{1 i}$ |
|  | $C_{2}$ |
|  | $\dot{A}_{e}, B_{e}, C_{e}, D_{e}$ |
|  | $v_{i}(l)$ |
|  | $w_{0}(t), w_{1}(t)$ |
|  | $V_{0}$ |
|  | $V_{1}$ |
|  | $V_{01}$ |
|  | $\boldsymbol{R}$ |
|  | $L$ |
|  | $\boldsymbol{A}, \boldsymbol{A}_{\boldsymbol{i}}$ |
|  | $\boldsymbol{K}$ |
|  | $\underset{\boldsymbol{A}}{\underline{A}}$ |

$\bar{W}(1)$
$D$
$R$
asymptotically stable matrix nonnegative-semisimple matrix nonnegative-definite matrix positive-definite matrix

For qubitrary $\boldsymbol{n} \times \boldsymbol{n} \boldsymbol{Q}, \boldsymbol{Q}$ define:

$$
\begin{gathered}
V_{11} \triangleq V_{1}+\sum_{i=1}^{\infty} C_{11}(Q+Q) C_{11}^{T} \\
Q \triangleq Q C_{1}^{T}+V_{01}+\sum_{i=1}^{\ell} A_{1}(Q+Q) C_{1}^{r} \\
A_{Q} \triangleq A-Q_{1} V_{13}^{-1} C_{1}
\end{gathered}
$$

II. Notation and Definitions
real numbers, $r \times s$ real matrices, $\mathbb{R}^{r \times 1}$, expectation
$n \times n$ idencity, transpose, group generalized inverse [25, p. 124]
Kronecker sum, Kronecker product [26]
trace of a square matrix $Z$
$I_{n}-T_{,} T \in \mathbb{R}^{n \times n}$
positive integers, $1 \leq n_{p} \leq n$
$n+n_{0}$
$n, n_{e}$-dimensional vectors
$I_{1}, l_{2}, q$-dimensional vectors
$n \times n$ matrices; $l_{1} \times n$ matrices, $l=1, \cdots, p$
$l_{2} \times n$ matrix
$n_{e} \times n_{e}, n_{e} \times l_{1}, q \times n_{e}, q \times l_{2}$ matrices
unit variance white noise, $i=1, \cdots, p$
$n$-dimensional, $I_{1}$-dimensional white noise processes
$n \times n$ nonnegative-definite intensity of $w_{0}(t)$
$l_{1} \times l_{1}$ positive-definite intensity of $w_{1}(t)$
$n \times l_{1}$ cross intensity of $w_{0}(t), w_{1}(t)$
$q \times q$ positive-definite matrix
$9 \times n$ matrix
$\left[\begin{array}{cc}A_{i} & 0 \\ B_{e} C_{1} & A_{e}\end{array}\right],\left[\begin{array}{cc}A_{i} & 0 \\ B_{e} C_{1} & 0\end{array}\right], i=1, \cdots, p$
$n_{4} \times l_{2}$ matrix
$\left[\begin{array}{cc}A_{1} & 0 \\ B_{r} C_{1}+K C_{2} & A_{f}\end{array}\right]$
$\left[\begin{array}{c}w_{0}(t) \\ B_{0} w_{1}(t)\end{array}\right]$
$\left[\begin{array}{cc}V_{0} & V_{01} B_{!}^{r} \\ B_{e} V_{01}^{r} & B_{e} V_{1} \dot{B}_{!}^{r}\end{array}\right]$

matrix with eigenvalues in open left-half plane
semisimple (nondefective) matrix with nonnegative eigenvalues symmetric matrix with nonnegative eigenvalues symmetric matrix with positive eigerivalues.
where $t \in[0, \infty)$, design an $n_{\text {eth }}$-order state extimator

$$
\begin{align*}
& x_{t}(t)=A_{0} x_{t}(t)+B_{d} y_{1}(t)  \tag{3.4}\\
& y_{d}(t)=C_{r} x_{t}(t)+D_{1} y_{2}(t) \tag{3.5}
\end{align*}
$$

which minimizes the then-astimation error crterion

$$
\begin{equation*}
J\left(A_{e}, D_{e}, C_{0}, D_{f}\right) \Delta \lim _{t \rightarrow \infty} \in\left[L x(t)-y_{f}(t)\right]^{r} R\left[L x(t)-y_{0}(t)\right] . \tag{3.0}
\end{equation*}
$$

To guarantec that $J$ is finite, assume that $A$ is asymptotically stable and consider the set of asymptotically stable reduced-order (i.e., fixed-order) estimators
$A \triangleq\left\{\left(A_{e}, B_{r}, C_{n}, D_{r}\right): A_{8}\right.$ is asymprotically stable $\}$.
Since the value of $J$ is independent of the internal realization of the transfer function corresponding to (3.4) and (3.5), without loss of genernlity we further restrict our attention to the set of admissible estimators

$$
\begin{gathered}
A^{+} \triangleq\left\{\left(A_{e}, B_{e}, C_{e}, D_{e}\right) \in A:\left(A_{e}, B_{e}\right)\right. \text { is controilable } \\
\text { and } \left.\left(A_{e}, C_{e}\right) \text { is observable }\right\} .
\end{gathered}
$$

An additional technical requirement is that ( $A_{e}, B_{e}, C_{e}, D_{e}$ ) be confined to the set

$$
\boldsymbol{A}^{+} \triangleq\left(\left(A_{e}, B_{e}, C_{e}, D_{e}\right) \in A^{+}:\right.
$$

$$
\left.C_{2}\left(Q_{1}-Q_{12} Q_{2}^{-1} Q_{12}^{\top}\right) C_{2}^{T} \text { is positive defi.ite }\right\}
$$

where

$$
\bar{Q} \triangleq\left[\begin{array}{ll}
Q_{1} & Q_{12} \\
Q_{12}^{T} & Q_{2}
\end{array}\right] \in \mathbb{W}^{\prime \times N}
$$

satisfies

$$
0=A Q+Q A^{T}+\sum_{i=1}^{T} A_{i} Q A_{i}^{T}+D
$$

and $Q_{2}$ is invertible since $\left(A_{e}, B_{e}\right)$ is controllable. The positive definiteness condition holds when $C_{2}$ has full row rank and $\bar{Q}$ is positive definite. As can be seen from the proof of Theorem 3.1, this condition implies the existence of the projection $\tau_{1}$ defined below.
The following factorization lemma is needed for the statement of the main result.
Lemma 3.1: Suppose $n \times n \boldsymbol{Q}, \dot{P}$ are nonnegative definite. Then $\boldsymbol{Q P}$ is nonnegative semisimple. If, in addition, rank $\varnothing \boldsymbol{O}=n_{e}$, then there exist $n_{e} \times n G, \Gamma$ and $n_{e} \times n_{e}$ invertible $M$ such that

$$
\begin{gather*}
\phi^{\beta}=G^{\top} M \Gamma,  \tag{3.78}\\
\Gamma G^{T}=I_{n_{\theta}} . \tag{3.7b}
\end{gather*}
$$

Furthermore, $G, M$, and $\Gamma$ are unique modulo a change of basis in isfe.

Proof: The result follows from [27, Theorem 6.2.5].
Since $Q \mathcal{Q}$ is semisimple (diagonalizable) it has a group generalized inverse $\left(Q^{\top}\right)^{\prime \prime}=G^{\top} M^{-\Gamma} \Gamma$ and

$$
\begin{equation*}
T_{2} \cong Q P(Q P)^{\prime}=G^{r} \Gamma \tag{3.8}
\end{equation*}
$$

is an oblique projection.
Theorem 3.1: Suppose $A$ is asymptotically stable and ( $\left.A_{e}, B_{e}, C_{f}, D_{f}\right)$ $\in A^{+}$solves the reduced-order state-estimation problem. Then there exist $n \times n$ nonnegative-definite matrices $Q, Q$, and $\boldsymbol{P}^{\boldsymbol{P}}$ such that $A_{e}, B_{e}$, $C_{e}$, and $D_{e}$ are given by

$$
\begin{gather*}
A_{d}=\Gamma\left(A-Q, V_{1 s}^{-1} C_{1}\right) C^{T},  \tag{3.9}\\
B_{\mathrm{e}}=\Gamma Q V_{12}^{-1},  \tag{3.10}\\
C_{e}=L T_{1 A} G_{1}^{r}  \tag{3.11}\\
D_{d}=L Q C_{2}^{r}\left(C_{2} Q C_{2}^{r}\right)^{-1} \tag{3.12}
\end{gather*}
$$

and such that $\boldsymbol{Q}, \boldsymbol{Q}$, and $\boldsymbol{\beta}$ amisfy

$$
0=A Q+Q A^{r}+\sum_{1=1} A_{1}(Q+Q) A_{1}^{T}+V_{0}-Q V_{14}^{-1} Q_{1}^{T}+r_{21} Q V_{10}^{-1} Q_{1}^{r_{21}},
$$

$$
\begin{align*}
& 0=A Q+Q A^{r}+Q V_{14}^{-1} Q_{3}^{r}-\tau_{2_{A}} Q V_{14}^{-1} Q_{1}^{r} \tau_{i_{1}}^{T},  \tag{3.14}\\
& 0=A_{Q}^{T} \beta+\beta A_{Q}+T_{1 /}^{T} L^{r} R L T_{1+}-T_{21}^{T} T_{1 L} L^{T} R L T_{1+} T_{21},  \tag{3.15}\\
& \operatorname{rank} \phi=\operatorname{rank} P=\operatorname{rank} \phi P=n_{\boldsymbol{q}} \tag{3.16}
\end{align*}
$$

where

$$
\begin{equation*}
T_{1} \triangleq Q C_{2}^{\top}\left(C_{2} Q C_{2}^{r}\right)^{-1} C_{2} \tag{3.17}
\end{equation*}
$$

Remark 3.1: Several special cases can be recovered from Theorem 3.1. For example, when the observation noise is nonsingular, i.e., when $y_{2}$ is absent, delect (3.12) and set $\tau_{1}=0$ [22]. Deleting also the multiplicative noise terms yields the Main Theorem of [1].

Specializing Theorem 3.1 to the full-order case $n_{e}=n$ reveals that the Lyapunov equation for $\beta$ is superfluous. In this case $G=\Gamma^{-1}$ and thus $G$ $=\Gamma=I_{n}$ without loss of generality.
Corollary 3.1: Assume $n_{e}=n, A$ is asymptotically stable and ( $A_{e}, B_{e}$, $\left.C_{e}, D_{e}\right) \in A^{+}$solves the full-order state-estimation problem. Then there exist $n \times n$ nonnegative-definite matrices $Q$ and $Q$ such that $A_{e}, B_{e}, C_{e}$, $D_{e}$ are given by

$$
\begin{gather*}
A_{e}=A-Q V_{1 s}^{-1} C_{1},  \tag{3.18}\\
B_{c}=Q V_{1 s}^{-1},  \tag{3.19}\\
C_{8}=L T_{1+}  \tag{3.20}\\
D_{8}=L Q C_{2}^{T}\left(C_{2} Q C_{2}^{T}\right)^{-1} \tag{3.21}
\end{gather*}
$$

and such that $Q$ and $Q$ satisfy

$$
\begin{gather*}
0=A Q+Q A^{T}+\sum_{i=1}^{Q} A_{1}(Q+Q) A_{1}^{T}+V_{0}-Q V_{1,}^{-1} Q,  \tag{3.22}\\
0=A Q+Q A^{T}+Q V_{1,}^{-1} Q ; . \tag{3.23}
\end{gather*}
$$

Remark 3.2: Note that by setting $A_{i}=0, C_{i}=0, i=1, \cdots, p$, it follows that (3.22) and (3.23) are decoupled and (3.23) is superfluous. To recover the standard Kalman filter which involves nonsingular noise, set $C_{2}=0$, delete (3.21) and define $\boldsymbol{T}_{1}=0$.

## IV. Additional Estimator Paths

We now consider the more general estimator

$$
\begin{gather*}
x_{\mathrm{f}}(t)=A_{\mathrm{e}} x_{\mathrm{r}}(t)+B_{d} y_{1}(t)+K y_{2}(t),  \tag{4.1}\\
y_{0}(t)=C_{r} x_{r}+D_{e} y_{2}(t)+R y_{1}(t) \tag{4.2}
\end{gather*}
$$

involving the additional gains $K$ and $R$.
Note that the additional path introduced in (4.2) implies that $J$ is infinite and thus the problem is meaningless. Hence, set $R=0$, and consider the additional path introduced by (4.1), i.e., filtering the noise-free measurement.

Replacing (3.4) by (4.1) and optimizing with respect to $K$ yields

$$
\begin{equation*}
0=G P Q C_{2} \quad, \tag{4.3}
\end{equation*}
$$

which implies

$$
\begin{equation*}
0=\tau_{2} \tau_{1} . \tag{4.4}
\end{equation*}
$$

Using (4.3), $Q=\tau_{2} Q$ and $\beta=\mathcal{P}_{2}$ [see (5.17)], the fitter gains (3.9) (3.15) become

$$
\begin{gather*}
A_{i}=\Gamma\left(A-Q V_{4}^{-1} C_{1}\right) G^{r}-K C_{2} G^{r}  \tag{4.5}\\
B_{\mathrm{r}}=\Gamma Q_{4} V_{\mathrm{k}}^{-1} \tag{4.6}
\end{gather*}
$$

$$
\begin{align*}
& C_{d}=L_{T_{1 \alpha}} G^{T}  \tag{4.7}\\
& D_{f}=L Q C_{2}^{\top}\left(C_{2} Q C_{2}^{\prime}\right)^{-1} \text {. }  \tag{4.8}\\
& 0=A Q+Q A^{T}+\sum_{i=1}^{n} A_{1}(Q+Q) A_{i}^{T}
\end{align*}
$$

$$
\begin{align*}
& 0=A Q+Q A^{T}+Q V_{i s}^{-1} Q_{s}^{r}-\tau_{2} Q V_{u}^{-1} Q_{r}^{T} r_{2 L}^{T}-Q C_{2}^{T} K^{T} G-G^{T} K C_{2} Q,  \tag{4.10}\\
& 0=\lambda_{Q}^{r} \beta+\beta A_{Q}+\tau_{1_{\perp}}^{T} L^{r} R L T_{1+}-\tau_{2 \perp}^{T} r_{1 \perp}^{T} L^{r} R L_{\tau_{1 \perp}} \tau_{2 \perp}, \tag{4.11}
\end{align*}
$$

where

$$
\lambda \triangleq A-G^{T} K C_{2}, A_{Q} \triangleq A_{Q}-G^{T} K C_{2}
$$

## V. Proof of the Man theorem

Using the notation of Section II the augmented system (3.1)-(3.4) can be written as

$$
\begin{equation*}
\dot{x}(t)=\left(A+\sum_{i=1} v_{i}(t) A_{i}\right) \hat{x}(t)+\tilde{w}(t) \tag{5.1}
\end{equation*}
$$

where

$$
f(t) \triangleq\left[x^{\top}(t), x_{f}^{\top}(t)\right]^{\top}
$$

To analyze (5.1) define the second-moment matrix

$$
\begin{equation*}
\phi(t)=\mathbb{R}\left[z(t) R^{T}(t)\right] . \tag{5.2}
\end{equation*}
$$

It follows from [28, Theorem 8.5.5, p. 142] thet $\boldsymbol{Q}(t)$ satisfies

$$
\begin{equation*}
\dot{Q}(t)=A \underline{Q}(t)+Q(t) A^{T}+\sum_{i=1}^{\ell} A_{i} \hat{Q}(t) A_{i}^{T}+D, \quad t \geq 0 \tag{5.3}
\end{equation*}
$$

Lemma 5.1: $A_{\varepsilon} \in A$ if and only if

$$
Q \triangleq A \oplus A+\sum_{i=1}^{n} A_{i} \otimes A_{i}
$$

is esymptotically stable.
Proof: The result follows from properties of the Kronecker product applied to partitioned matrices. See [22], [26] for details.
Hence, A stable assures

$$
\phi \triangleq \lim _{t \rightarrow \infty}\left\{\left[\tilde{x}(t) x^{T}(t)\right]\right.
$$

exists. Furchermore, $\boldsymbol{Q}$ and its nonnegative-definite dual $\boldsymbol{P}$ are unique solutions of the modified Lyapunov equations

$$
\begin{align*}
& 0=A Q+Q A^{T}+\sum_{i=1}^{n} A_{i} Q A_{i}^{T}+D,  \tag{5.4}\\
& 0=A^{T} \beta+P A+\sum_{i=1}^{N} A_{i}^{T} \beta A_{i}+R . \tag{5.5}
\end{align*}
$$

Parition $n \times \boldsymbol{N}, \boldsymbol{P}$ into $n \times n_{0} n \times n_{e}$ and $n_{c} \times n_{p}$ subblocks as

$$
\boldsymbol{\theta}=\left[\begin{array}{ll}
Q_{1} & Q_{12} \\
Q_{12}^{F} & Q_{2}
\end{array}\right], \beta=\left[\begin{array}{ll}
P_{1} & P_{12} \\
P_{12}^{r} & P_{2}
\end{array}\right]
$$

and define the $n \times n$ nonnegativedefinite matrices

$$
\begin{aligned}
& Q Q_{1}-Q_{12} Q_{2}^{-1} Q_{12}^{T}, \quad P \triangleq P_{1}-P_{12} P_{2}^{-1} P_{12}^{T} \\
& Q \triangleq Q_{12} Q_{2}^{-1} Q_{12}^{T}, \quad P \triangleq P_{12} P_{2}^{-1} P_{12}^{T}
\end{aligned}
$$

and the $n_{e} \times n_{f} n_{f} \times n_{e}, n_{f} \times n$ matrices

$$
G \triangleq Q_{2}^{-1} Q_{12}^{T}, M \triangleq Q_{2} P_{2}, \Gamma \triangleq-P_{2}^{-1} P_{12}^{r}
$$

To optimize (3.6) subject to the constraint (5.4) over $A^{+}$, form the Lagrangian
$\mathcal{L}\left(A_{e}, B_{f}, C_{E}, D_{f}, \boldsymbol{Q}, \boldsymbol{P}_{!} \lambda\right)$

$$
\hat{\operatorname{str}}\left[\lambda \phi R+\left(\lambda Q+Q A^{T}+\sum_{i=1} A_{i} Q A_{i}^{r}+\nabla\right) P\right]
$$

where the Lagrange multipliers $\lambda \geq 0$ and $\bar{P} \in \mathbb{R}^{\boldsymbol{N} \times{ }^{*}}$ are nor both zero and $Q$ and $\boldsymbol{P}$ are viewed as arbitrary $\tilde{n} \times \tilde{n}$ matrix variables. Setting $\partial L /$ $\partial Q=0, \lambda=0$ implies $P=0$ since $\left(A_{e}, B_{e}, C_{e}, D_{e}\right) \in A^{+}$. Hence, without loss of generality, set $\lambda=1$. Thus, the stationarity conditions are given by'

$$
\begin{align*}
& \frac{\partial L}{\partial P}=A \phi+Q A^{r}+\sum_{i=1} A_{i} Q A_{i} r^{2} P=0,  \tag{5.6}\\
& \frac{\partial L}{\partial Q}=A^{\top} P+P A+\sum_{i=1}^{P} A_{i}^{T} F A_{i}+\bar{R}=0, \\
& \frac{\partial L}{\partial A_{e}}=P P_{12} Q_{12}+P_{2} Q_{2}=0, \\
& \frac{\partial L}{\partial B_{e}}=\left(P_{12}^{r} Q_{1}+P_{2} Q_{12}^{r}\right) C_{i}^{r}+\sum_{i=1}^{n}\left(P_{12}^{r} A_{i} Q_{1} C_{1 i}^{r}\right. \\
& \left.+P_{2} B_{8} C_{11} Q_{1} C_{1 i}^{T}\right)+P_{12}^{T} V_{01}+P_{2} B_{e} V_{1}=0,  \tag{5.9}\\
& \frac{\partial L}{\partial C_{4}}=-R L Q_{12}+R D_{4} C_{2} Q_{12}+R C_{e} Q_{2}=0,  \tag{5.10}\\
& \frac{\partial L}{\partial D_{1}}=-R L Q_{1} C_{2}^{T}+R D_{5} C_{2} Q_{1} C_{2}^{T}+R C_{4} Q_{12}^{J} C_{2}^{T}=0 . \\
& \text { (5.7) } \\
& \text { (5.8) } \\
& \text { (5.11) } \\
& \text { Expanding (5.6) and (5.7) yields } \\
& 0=A Q_{1}+Q_{1} A^{T}+\sum_{i=1} A_{1} Q_{1} A_{i}^{T}+V_{0} .  \tag{5.12}\\
& 0=A Q_{12}+Q_{1} C_{i}^{T} B_{e}^{T}+Q_{12} A_{e}^{T}+\sum_{i=1}^{P} A_{1} Q_{1} C_{11}^{T} B_{e}^{T}+V_{01} B_{e}^{T},  \tag{5.13}\\
& 0=B_{e} C_{1} Q_{12}^{T}+A_{i} Q_{2}+Q_{12}^{T} C_{1}^{T} B_{e}^{T}+Q_{2} A_{e}^{T}+\sum_{i=1}^{r} B_{e} C_{11} Q_{1} C_{i i}^{T} B_{e}^{T}+B_{e} V_{1} B_{e}^{T},  \tag{5.14}\\
& 0=A^{T} P_{12}+C_{1}^{T} B^{T} P_{2}+P_{12} A_{f}-L^{T} R C_{e}+C_{2}^{T} D_{f}^{T} R C_{e} \text {, }  \tag{5.15}\\
& 0=A_{e}^{r} P_{2}+P_{2} A_{e}+C_{e}^{T} R C_{e} . \tag{5.16}
\end{align*}
$$

Note that the $(1,1)$ subblock of $(5.7)$ characterizing $P_{1}$ has been omitted from the above equations since the extimator gains are independent of $P_{1}$. Note that (5.8) implies (3.7a) and (3.7b). Since

$$
Q_{2} P_{2}=P_{2}^{-1 / 2}\left(P_{2}^{1 / 2} Q_{2} P_{2}^{1 / 2}\right) P_{2}^{1 / 2}
$$

$\boldsymbol{M}$ is positive semisimple. Sylvester's inequality yields (3.16). Note also that

$$
\begin{equation*}
Q=r_{2} Q, \quad \beta=\beta_{\tau_{2}} \tag{5.17}
\end{equation*}
$$

Next (3.10), (3.11), and (3.12) follow from (5.9), (5.10), and (3.11) by

[^13]using the identities
\[

$$
\begin{array}{ll}
Q_{1}=Q+Q . & P_{1}=P+P, \\
Q_{12}=Q \Gamma^{r}, & P_{12}=-P G^{r}, \\
Q_{2}=\Gamma Q \Gamma^{r}, & P_{2}=G A G^{r} . \tag{3.20}
\end{array}
$$
\]

Substituting (3.10), (3.11), (3.12) and (5.18)-(5.20) into (5.12)-(5.16) and using (5.12) $+G^{T} \Gamma(5.13) G-(5.13) G-(5.13 G)^{T}$ and $G^{T} \Gamma(5.13) G-(5.13) G-(5.13 G)^{T}$ yields (3.13) and (3.14). Using $\Gamma^{T} G(5.15) \Gamma$ - (5.15) $\Gamma$ - (5.15 $)^{T}$ yields (3.15). Finally, $\Gamma(5.13)$ (5.14) or $\boldsymbol{G}(5.15)-(5.16)$ yields (3.9).

Remark 5.1. Equations (4.5)-(4.11) are derived in a similar manner with $A$ replaced by $A$ in (5.1).

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# The Optimal Projection Equations for Static and Dynamic Output Feedback: The Singular Case 

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Dedicated to the memory of Professor Violet B. Haas November 23, 1926-January 21, 1986


#### Abstract

Oblique projections have been shown to arise naturally in both static and dyamic optimal design problems. For static controllers an oblique projection was inherent in the early work of Levine and Athans, while for dynamic controllers an oblique projection was developed by Hyland and Bernstein. This mote is motivated by the following antural question: What is the relationship berween the oblique projection arising in optimal static outpat feedback and the oblique projection arising in optimal fixed-order dyanmic compensation? We show that in monstictly proper optimal output feedback there are, indeed, three disthnct oblique projections corresponding to slagular measurement noise, singular control weighting, and reduced compensator order. Moreover, we uaify the Levine-Athans and Hyland-Berastein appronches by rederiviug the optimal projection equations for combined static/dyammic (monstrictly proper) output feedback in a form which clearly Ifrastrates the role of the thrse projections in characterixing the optimal feedback galins. Even when the dynamic componeut of the monstrictly proper controlier is of full order, the controller is characterbed by four matrix equations which generalize the staadard LQG result.


## I. Introduction

The optimal static output-feedback problem [1]. [2] and the optimal fixed-order dynamic-compensation problem [3], 14] have been extensively investigated. A salient feature of the necessary conditions for each of these problems is the presence of an oblique projection (idempotent matrix) which arises as a direct consequence of optimality. For the static problem with noise-free measurements (i.c., singular measurement noise) the necessary conditions involve the projection [2]

$$
r_{1}=Q C^{T}(C Q C)^{-1} C
$$

where $Q$ is the steady-state closed-loop state covariance. The dual projection

$$
T_{2}=B\left(B^{T} P B\right)^{-1} B^{T} P
$$

arises analogousiy in the corresponding problem involving singular control weighting. Furthermore, for fixed-order dynamic compensation with noisy measurements, it has recently been shown [4] that the necessary conditions give rise to the projection

$$
r_{3}=Q P\left(Q^{P}\right)^{\prime}
$$

where () ${ }^{\prime \prime}$ denotes group generalized inverse and $\boldsymbol{Q}$ and $\boldsymbol{P}$ are rankdeficient nonnegative-definite matrices analogous to the controllability

[^14]and observability Gramians of the compensator. To understand the reletionships among $\tau_{1}, \tau_{2}$, and $\tau_{3}$, the contribution of the present note is a unified treatment of the necessary conditions for optimal static/ dynamic feedbock compensation which clearly illustrates the role of the three projections in characterizing the optimal feedback gains. Even in the full-order case in which $\tau_{3}$ is the identity, the result provides a generalization of the standard LQG result to nonstrictly proper conarollers in which case the separation principle does not hold.

To clarify the ramifications of noise and weighting singularities in optimal output feedback, consider the problem of minimizing

$$
\begin{equation*}
J=\lim _{r \rightarrow \infty} T\left[x^{\top} R_{0} x+u^{\top} R_{1} u\right] \tag{1.1}
\end{equation*}
$$

with plank dynamics

$$
\begin{gather*}
x=A x+B u+w_{0}  \tag{1.2}\\
y=C x+w_{1} \tag{1.3}
\end{gather*}
$$

and nonstrictly proper feedback compensator

$$
\begin{align*}
& x_{c}=A_{c} x_{c}+B_{c} y  \tag{1.4}\\
& u=C_{c} x_{c}+D_{c} y \tag{1.5}
\end{align*}
$$

As pointed out in [3], $J$ is finite only if

$$
\begin{equation*}
0=\sigma\left[D_{c}^{T} R_{1} D_{c} V_{1}\right] \propto 0=R_{1} D_{c} V_{1} \tag{1.6}
\end{equation*}
$$

where $V_{1}$ denotes the intensity of $w_{1}$. Clearly, when $R_{1}$ and $V_{1}$ are nonsingular (1.6) implies $D_{c}=0$, and hence direct feedthrough is not pernitted, i.e., the compensator must be strictly proper. Conversely, to utilize a static gain $D_{c}$, either $R_{1}$ or $V_{1}$ must be singular. By writing singular $R_{1}$ and $V_{1}$ without loss of generality as

$$
R_{1}=\left[\begin{array}{cc}
R_{1} & 0  \tag{1.7}\\
0 & 0
\end{array}\right], V_{1}=\left[\begin{array}{ll}
D_{1} & 0 \\
0 & 0
\end{array}\right]
$$

it follows that che static transmission between noisy measurements and weighted controls must be zero (see Fig. 1).

The reader will observe that three feedback paths which are not ruled


Fig. 1.
out by (1.6) do not appear in Fig. 1. Specifically: 1) nonnoisy measurements can be fed back to unweighted controls: 2) dynamiccompensator cutputs can be fed back to unweighted controls; and 3) nonnoisy measurements can serve as inputs to the dynamic compensator. The reason for considering the more limited configuration shown in Fig. 1 is that only these paths are explicitly characterized by the necessary conditions. Hence, for simplicity we first consider only the scheme of Fig. 1, and later introduce the remaining permissible pachs. Interestingly, while these additional gains are not completely determined by the necessary conditions, they appear to play an important role in governing geometric interrelationships among the three projections.

Two final comments are in order. First, since our results are carried out in a multiplicative noise setting, we generalize previous results on state feedbeck [15]-[18] and dynamic compensation [9]-[11]. The mocivation for using a multiplicative white noise model is to represent piant parameter uncertainties and thereby obtain robust controllers [12]. Also, the derivations of the necessary conditions are straightforward extensions of the Lagrange multiplier technique used in [4] and hence have been omitted.

## II. Notation and Definitions

S, $\Omega^{r \times s}, \mathbb{R}^{\prime}$, ⿷
I. ( $)^{T}$ ( $)^{\prime \prime}$
$\oplus, \otimes$
${ }^{1}+$
asymptoxically stable matrix
nonnegative-semisimple matrix
$n, m_{1}, m_{2}, l_{1}, l_{2}, n_{c}, p$
$x_{1} u_{1}, u_{2,} y_{1}, y_{2}, x_{c}$
$A_{1} A_{i} ; B_{1}, B_{1 i} ; C_{1}, C_{11}$
$B_{2} . C_{2}$
$A_{c}, B_{c}, C_{c}, D_{c}, E_{c}$
u, (I)
$w_{0}(t), w_{1}(t)$
$V_{0}, V_{1}$
$V_{01}$
$R_{0}, R_{1}$
$R_{01}$
$\boldsymbol{A}_{1}, \boldsymbol{A}_{1}$
$D . R$
$\boldsymbol{A}_{\boldsymbol{A}} \boldsymbol{A}_{i}$
real numbers, $r \times s$ real matrices, $J_{\beta} \times 1$, expectation
$r \times r$ identity, transpose, group generalized inverse [13, p. 124]
Kronecker sum, Kronecker product
$I_{n}-T_{1}+\in \mathbb{T}^{n \times n}$
matrux with eigenvalues in open lefi-half plane
semisimple (nondefective) matrix with nonnegative eigenvalues positive integers
$n, m_{1}, m_{2}, l_{1}, l_{2}, n_{c}$-dimensional vectors
$n \times n$ matrices, $n \times m_{1}$ matrices, $I_{\mathrm{i}} \times n$ matrices, $i=1, \cdots, p$
$n \times m_{2}$ matrix, $l_{2} \times n$ matrix
$n_{c} \times n_{c}, n_{c} \times l_{1}, m_{1} \times n_{c}, m_{1} \times l_{2}, m_{2} \times l_{1}$ matrices
unit variance white noise, $i=1, \cdots, p$
$n$-dimensional, $l_{1}$-dimensional white noise
intensities of $w_{0}, w_{1} ; V_{0} \geq 0, V_{1}>0$
$n \times l_{1}$ cross intensity of $w_{0}, w_{1}$
state and control weightings; $R_{0} \geq 0, R_{1}>0$
$n \times m_{1}$ cross weighting: $R_{0}-R_{01} R_{1}^{-1} R_{01}^{T_{1}} \geq 0$
$A+B_{1} D_{c} C_{2}+B_{2} E_{r} C_{1}, A_{1}+B_{11} D_{c} C_{2}+B_{2} E_{c} C_{1,} i=1, \cdots, p$
$\left[\begin{array}{c}\boldsymbol{I}_{n} \\ \left(B_{2} E_{c}\right)^{r}\end{array}\right] r\left[\begin{array}{cc}V_{0} & V_{01} \\ V_{01}^{\tau} & V_{1}\end{array}\right]\left[\begin{array}{c}I_{n} \\ \left(B_{2} E_{c}\right)^{T}\end{array}\right] \cdot\left[\begin{array}{c}I_{n} \\ D_{c} C_{2}\end{array}\right] r\left[\begin{array}{ll}R_{0} & R_{01} \\ R_{01}^{\tau} & R_{1}\end{array}\right]\left[\begin{array}{c}I_{n} \\ D_{r} C_{2}\end{array}\right]$
$\left[\begin{array}{cc}\boldsymbol{A} & B_{1} C_{c} \\ B_{c} C_{1} & A_{c}\end{array}\right] \cdot\left[\begin{array}{cc}\boldsymbol{A}_{1} & B_{1} C_{c} \\ B_{c} C_{13} & 0\end{array}\right]$

$$
\begin{aligned}
& \text { V } \quad\left[\begin{array}{cc}
D & V_{01} B_{c}^{T}+B_{2} E_{c} V_{1} B_{c}^{\tau} \\
B_{c} V_{01}^{T}+B_{c} V_{1}\left(B_{2} E_{c}\right)^{\top} & B_{c} V_{1} B_{c}^{T}
\end{array}\right] \\
& \boldsymbol{R} \quad\left[\begin{array}{cc}
\underset{R}{R} \\
C_{c}^{r} R_{01}^{r}+C_{c}^{r} R_{1} D_{c} C_{2} & R_{01} C_{c}+\left(D_{c} C_{2}\right)_{R_{1}} C_{c} \\
C_{c}^{r} R_{1} C_{c}
\end{array}\right]
\end{aligned}
$$

For arbitrary $\boldsymbol{n} \times \boldsymbol{n} \boldsymbol{Q}, \mathbf{P}, \mathbf{Q}, \boldsymbol{P}_{,}, \boldsymbol{\tau}_{1}, \boldsymbol{\tau}_{2}$ define:

$$
\begin{aligned}
& R_{18} \triangleq R_{i}+\sum_{i=1}^{N} B_{i 1}^{r} P B_{1 /}, \quad V_{18} \text { 全 } V_{1}+\sum_{i=1}^{R} C_{11} Q C_{11}^{T}, \\
& \text { Q. } \triangleq Q C_{1}^{T}+V_{01}+\sum_{i=1}^{P} A_{i} Q C_{i, 1}^{T} \quad \mathcal{O}_{3} \triangleq B_{i}^{T} P+R_{01}^{T}+\sum_{i=1}^{\mathcal{E}} B_{1 i}^{T} P A_{i} \text {, } \\
& R_{11} \triangleq R_{1}+\sum_{i=1}^{n} B_{1 i}^{T}(P+\beta) B_{1 i}, \quad \gamma_{1} \triangleq V_{1}+\sum_{i=1}^{p} C_{11}(Q+Q) C_{11}^{r}
\end{aligned}
$$

Q. $\triangleq Q C_{1}^{T}+V_{01}+\sum_{i=1}^{P} A_{i}(Q+Q) C_{i=}^{T}$

$$
\theta_{s} \triangleq B_{i}^{T} P+R_{i 1}^{T}+\sum_{i=1} B_{i 1}^{T}(P+P) A_{i}
$$

$A \triangleq A-B_{1} R_{n}^{-1}\left(\theta_{s}+B_{i}^{T} P_{1} r_{1}-r_{2}\left(\Theta_{1}+Q C_{1}^{\top}\right) P_{18}^{-1} C_{1}\right.$.

$$
A_{i} \triangleq A_{i}-B_{11} R_{1}^{-1}\left(Q_{1}+B_{i}^{T} P\right) \tau_{1}-T_{2}\left(\mathcal{Q}_{1}+Q C_{1}\right) D_{18}^{-1} C_{11},
$$



$D_{0} \triangleq V_{0}-V_{01} D_{1 s}^{-1}\left(\Theta_{0}+Q C_{1}^{\top}\right)^{T_{2}^{T}} \boldsymbol{T}_{2} T_{2}\left(\Theta_{0}+Q C_{1}^{T}\right) V_{18}^{-1} V_{\theta 1}^{T}$

$$
+\tau_{2}\left(Q_{C}+Q C_{1}^{T}\right) P_{1:}^{-1} V_{1} \nabla_{1}^{-1}\left(Q_{0}+\phi C_{1}^{T}\right)^{\tau_{2} T} .
$$

## III. Static Output Feedback

## Static Output Feedback Problem

Given the controlled system
$\dot{x}(t)=\left(A+\sum_{i=1}^{p} v_{i}(t) A_{i}\right) x(t)$

$$
\begin{gather*}
+\left(B_{1}+\sum_{i=1}^{n} v_{1}(t) B_{11}\right) u_{1}(t)+B_{2} w_{2}(t)+w_{0}(t)  \tag{3.1}\\
y_{1}(t)=\left(C_{1}+\sum_{i=1}^{\infty} v_{i}(t) C_{1 i}\right) x(t)+w_{1}(t)  \tag{3.2}\\
y_{2}(t)=C_{2} x(t) \tag{3.3}
\end{gather*}
$$

where $t \in[0, \infty)$, determine $D_{c}$ and $E_{c}$ wuch that the static output feedback'law

$$
\begin{align*}
& \omega_{1}(t)=D_{c} y_{2}(t)  \tag{3.4}\\
& \omega_{2}(t)=E_{c} y_{1}(t) \tag{3.5}
\end{align*}
$$

minimizes the performance criterion

$$
\begin{equation*}
J \Delta \lim _{t \rightarrow \infty}\left[x(t)^{r} R_{0} x(t)+2 x(t)^{r} R_{01} w_{1}(t)+w_{1}(t)^{r} R_{1} w_{1}(t)\right] . \tag{3.6}
\end{equation*}
$$

To develop neceasary conditions for this problem, $D_{c}$ and $E_{c}$ must be restricted to the set of second-moment-stabilizing gains
$8 \Delta\left\{\left(D_{c}, E_{c}\right): A \oplus A+\sum_{i=1}^{n} A_{i} \otimes A_{i}\right.$ is asymptotically stable $\}$
The requiremem $\left(D_{c}, E_{c}\right) \in \delta$ implies the existence of the steady-state closed-loop atate covariance $Q \triangleq \lim _{\rightarrow \rightarrow \infty}\left[\left[x(t) x(t)^{T}\right]\right.$. Furthermore, $Q$ and its nomnegative-definite dual $P$ are the unique solutions of the modified Lyapunov equations

$$
\begin{align*}
& 0=A Q+Q A^{T}+\sum_{i=1}^{n} A_{i} Q A_{i}^{T}+D  \tag{3.7}\\
& 0=A^{T} P+P A+\sum_{i=1}^{N} A_{i}^{T} P A_{i}+R . \tag{3.8}
\end{align*}
$$

An additional technical assumption is that $\left(D_{c}, E_{c}\right)$ be confined to the set

$$
S+\Lambda\left\{\left(D_{s}, E_{c}\right) \in \delta: C_{2} Q C_{2}^{T}>0 \text { and } B_{2}^{T} P B_{2}>0\right\}
$$

In order to obtain closed-form expressions for the feedback gains we make the additional assumption here and in Section IV that

$$
\begin{equation*}
\left[B_{1} \neq 0=C_{11}=0\right], i=1, \cdots, p \tag{3.9}
\end{equation*}
$$

i.e., for each $i, B_{1 /}$ and $C_{1 /}$ are not both noazero. By optimizing (3.6) with respect to $D_{c}$ and $E_{c}$ and manipulating (3.7) and (3.8), we obtain the following result.

Theorem 3.1: Suppose $\left(D_{c}, E_{c}\right) \in \$+$ solves the static output feedback problem. Then there exist $n \times n$ nonnegative-definite $Q, P$ such that

$$
\begin{align*}
& D_{c}=-R_{1 s}^{-1} P_{t} Q C_{2}^{T}\left(C_{2} Q C_{2}^{\eta}\right)^{-1},  \tag{3.10}\\
& E_{f}=-\left(B_{2}^{T} P B_{2}\right)^{-1} B_{2}^{T} P Q_{c} V_{12}^{-1}, \tag{3.11}
\end{align*}
$$

and such that $Q$ and $P$ satisfy

$$
\begin{align*}
& 0=\left(A-B_{3} R_{i s}^{-1} P_{3} r_{1}\right) Q+Q\left(A-B_{3} R_{13}^{-1} P_{3} r_{1}\right)^{T}+V_{0} \\
& +\sum_{i=1}^{\infty}\left(A_{i}-B_{11} R_{i=}^{-1} \theta_{3} r_{1}\right) Q\left(A_{j}-B_{1 j} R_{j g}^{-1} \theta_{s} r_{1}\right)^{T} \\
& -Q_{i} V_{18}^{-1} Q_{j}^{\top}+T_{21} Q_{1} V_{i s}^{-1} Q_{i}^{T} T_{21}^{T},  \tag{3.12}\\
& 0=\left(A-\tau_{2} Q_{0} V_{1 .}^{-1} C_{1}\right)^{T} P+P\left(A-\tau_{2} Q_{1} V_{1,}^{-1} C_{1}\right)+R_{0} \\
& +\sum_{i=1}^{n}\left(A_{i}-\tau_{2} Q_{0} V_{11}^{-1} C_{11}\right)^{\top} P\left(A_{i}-\tau_{2} Q V_{13}^{-1} C_{11}\right)
\end{align*}
$$

where

$$
\begin{align*}
& \tau_{1} \otimes Q C_{2}^{T}\left(C_{2} Q C_{2}^{T}\right)^{-1} C_{2}  \tag{3.14}\\
& \eta_{2} \otimes B_{2}\left(B_{2}^{T} P B_{3}\right)^{-1} B_{2}^{T} P . \tag{3.15}
\end{align*}
$$

Remork 3.1: Several epecial cases can be recovered formally from Theorem 3.1. For example, when the coatrol weighting is nonaingular
and the measurement noise is zero, i.c., when $u_{2}$ and $y_{1}$ are absent, delete (3.11) and set $\mathrm{T}_{2}=0$. Deleting also the multiplicative noise terms yields the usual static output feedback result [1], [2].

## IV. Dynamic Output Feedback

We now expand the formulation of the static problem to include a purely dynamic (strictly proper) dynamic compensator.

## Dynamic Output Feedback Problem

Given (3.1)-(3.3), determine $A_{c}, B_{c}, C_{c}, D_{c}, E_{c}$ such that the static and dynamic output feedback law

$$
\begin{gather*}
\dot{x}_{c}(t)=A_{c} x_{c}(t)+B_{c} y_{1}(t),  \tag{4.1}\\
u_{1}(t)=C_{c} x_{c}(t)+D_{c} y_{2}(t),  \tag{4.2}\\
u_{2}(t)=E_{c} y_{1}(t) \tag{4.3}
\end{gather*}
$$

minimizes the performance criterion (3.6).
We restrict our attention to second-moment-stabilizing controllers
$D \hat{\rho}\left\{\left(A_{c}, B_{c}, C_{c}, D_{c}, E_{c}\right): \hat{A} \oplus \hat{A}+\sum_{i=1} \tilde{A}_{i} \otimes \tilde{A}_{i}\right.$ is

$$
\text { asymptotically stable and } \left.\left(A_{c}, B_{c}, C_{c}\right) \text { is minimal }\right\} \text {. }
$$

which implies the existence of $\bar{Q} \lim _{t \rightarrow \infty} E[\bar{x}(r) \vec{x}(t)]^{T}$, where $\tilde{x}(t) \triangleq$ $\left(x(t)^{T}, x_{S}(t)^{\eta}\right)^{r}$. Furthermore, $\mathcal{Q}$ and its dual $P$ are the unique solutions of the modified Lyapunov equations

$$
\begin{align*}
& 0=\vec{A} \bar{Q}+\bar{Q}^{T}+\sum_{i=1}^{p} \vec{A}_{i} \bar{Q}_{i}^{T}+V  \tag{4.4}\\
& 0=\vec{A}^{T} \vec{P}+\vec{P} \vec{A}^{2}+\sum_{i=1}^{\infty} \bar{A}_{i}^{T} \vec{P} \vec{A}_{i}+\vec{R} \tag{4.5}
\end{align*}
$$

Partitioning

$$
\phi=\left[\begin{array}{cc}
Q_{1} & Q_{12} \\
Q_{12}^{T} & Q_{2}
\end{array}\right], P=\left[\begin{array}{ll}
P_{1} & P_{12} \\
P_{12}^{T} & P_{2}
\end{array}\right]
$$

where $Q_{12}$ and $P_{12}$ are $n \times n_{5}$, we also require

$$
\begin{aligned}
D^{\cdot} \stackrel{\vdots}{=}\left\{\left(A_{c}, B_{c}, C_{r}, D_{c}, E_{c}\right) \in D: C_{2}\left(Q_{1}-Q_{12} Q_{2}^{-1} Q_{12}^{r}\right) C_{2}^{T}>0\right. \\
\left.\quad \text { and } B_{2}^{T}\left(P_{1}-P_{12} P_{2}^{-1} P_{12}^{T}\right) B_{2}>0\right\} .
\end{aligned}
$$

Optimizing (3.6) over $\mathscr{D}^{+}$, introducing new variables

$$
\begin{gather*}
Q \triangleq Q_{1}-Q_{12} Q_{2}^{-1} Q_{12}^{T}, \quad P \triangleq P_{1}-P_{12} P_{2}^{-1} P_{12}^{T}  \tag{4.6}\\
Q \leqq Q_{12} Q_{2}^{-1} Q_{12}^{T}, \quad P \leqq P_{12} P_{2}^{-1} P_{12}^{T} \tag{4.7}
\end{gather*}
$$

and manipulating (4.4) and (4.5). we obtain the dynamic extension of Theorem 3.1. The following lemma is required for the statement of the result.

Lemma 4.1: Suppose $n \times n \boldsymbol{Q}, P$ are nonnegative definite. Then $Q P$ is nonnegative semisimple. If, in addition, rank $\phi^{\beta}=n_{c}$, then there exist $n_{c} \times n G, \Gamma$ and $n_{c} \times n_{c}$ invertible $M$ such that

$$
\begin{equation*}
\phi \beta=G^{T} M T, \Gamma G^{T}=I_{n_{r}} \tag{4.8a,b}
\end{equation*}
$$

Proof: The result follows from [14, Theorem 6.2.5].
Since $Q \boldsymbol{P}$ is semisimple it has a group inverse $\left(Q F^{\prime}=G^{\top} M^{-1} \Gamma\right.$ and

$$
\begin{equation*}
T_{3} \triangleq Q P(\phi P)^{\prime}=G^{T} T \tag{4.9}
\end{equation*}
$$

is an oblique projection.
Theorem 4.1: Suppose $\left(A_{c}, B_{c}, C_{c}, D_{c}, E_{c}\right) \in D^{+}$solves the dynamic curput feedback problem. Then there exist $n \times n$ nonnegative-definite $Q$,

P, Q. Psuch that
where $\tau_{1}$ and $\tau_{2}$ are given by (3.14) and (3.15), $G, \Gamma$ satisfy (4.8a), (4.8b), and such that, with $\tau_{3}$ given by (4.9), Q, P, Q, and $P$ satisfy
$0=A^{T} P+P \hat{A}+\hat{R}_{0}+\sum_{i=1}^{R}\left[\hat{A}_{i}^{T} P \hat{A}_{i}+\left(\hat{A}_{1}-\left[r_{2 L} \Theta_{1}-\tau_{2} Q C_{i}^{T} D_{18}^{-1} C_{11}\right)^{T}\right.\right.$

$$
\left.\cdot P\left(A_{1}-\left[r_{2 \perp} Q_{1}-r_{2} \phi C_{1}^{T}\right] V_{1 s}^{-1} C_{1 i}\right)\right]-\left(\theta_{s} \tau_{1 \perp}-B_{1}^{T} \beta_{r_{1}}\right)^{T} R_{1 s}^{-1}
$$



Remark 4.1: Setting $\tau_{1}=\tau_{2}=0, D_{c}=0, E_{c}=0$ yields the results of [4], [11].

Remark 4.2: Suppose. $n_{c}=n$ so that $\tau_{3}=I_{n}$. Then the resulting fullorder nonstrictly proper controller is characterized by four matrix equations which generalize the standard LQG result. In this case the separation principle is no longer valid.

## V. Additional Feedback Paths

We now introduce the feedback paths not shown in Fig. 1. For the static problem replace (3.5) by

$$
\begin{equation*}
u_{2}(t)=E_{c} y_{1}(t)+K_{1} y_{2}(t) \tag{5.1}
\end{equation*}
$$

Opaimizing with respect to $\boldsymbol{K}_{1}$ yields the additional condition

$$
\begin{equation*}
0=C_{1} Q P B_{2} \tag{5.2}
\end{equation*}
$$

which implies

$$
\begin{equation*}
0=r_{2} r_{1} \tag{5.3}
\end{equation*}
$$

This geometric condition holds when $\boldsymbol{K}_{1}$ is optimally chosen. Although $\boldsymbol{K}_{1}$

$$
\begin{aligned}
& 0=A_{Q_{s}}^{T} P+P A_{Q_{T}}+\left(\phi_{s} \tau_{1 \perp}-B_{1}^{T} \phi_{T_{1}}\right)^{T} R_{1 s}^{-1}\left(\beta_{1} \tau_{1+}-B_{1}^{7} \beta_{\tau_{1}}\right)
\end{aligned}
$$

$$
\begin{align*}
& \operatorname{rank} Q=\operatorname{rank} P=\operatorname{rank} \phi \beta=n_{c} .
\end{align*}
$$

$$
\begin{aligned}
& 0=\hat{A} Q+Q \hat{A}^{T}+D_{0}+\sum_{i=1}^{R}\left[\hat{A}_{i} Q \hat{A}_{i}^{T}+\left(\hat{A}_{i}-B_{1}, R_{i s}^{-1}\left(\otimes_{s} r_{14}-B_{i}^{Y} \hat{P}_{\tau_{1}}\right]\right)\right. \\
& \left.\cdot \boldsymbol{Q}\left(\hat{A}_{1}-B_{11} \mathcal{R}_{18}^{-1}\left[\phi_{1} \tau_{11}-B_{1}^{T} P_{r_{1}}\right]\right)^{r}\right]-\left(r_{2 \perp} Q-r_{2} \phi C_{1}^{T}\right) \hat{1}_{18}^{-1}
\end{aligned}
$$

$$
\begin{align*}
& A_{c}=\Gamma\left(A-B_{1} R_{11}^{-1} \theta_{i}-Q_{1} D_{1:}^{-1} C_{1}\right) G^{T}  \tag{4.10}\\
& B_{c}=\Gamma\left(r_{2} \theta_{-}-r_{2} \phi C_{i}^{\top} P_{1:}^{-1},\right.  \tag{4.11}\\
& C_{e}=-R_{1,}^{-1}\left(\mathcal{P}_{s} \tau_{1_{1}}-B_{1}^{r} \boldsymbol{P}_{T_{1}}\right) G^{r},  \tag{4.12}\\
& D_{\mathrm{f}}=-R_{1:}^{-1}\left(\mathcal{S}_{s}+B_{1}^{T}\right)^{\prime} Q C_{2}^{T}\left(C_{2} Q C_{2}^{T}\right)^{-1},  \tag{4.13}\\
& E_{c}=-\left(B_{2}^{T} P B_{2}\right)^{-1} B_{2}^{T} P\left(Q_{1}+Q C_{1}^{T}\right) D_{1 s}^{-1} \tag{4.14}
\end{align*}
$$

is nor given explicitly, it does play a role in the necessary conditions since $A$ is replaced by $A+B_{2} K_{1} C_{2}$.
For the dynamic problem replace (4.1) and (4.3) by

$$
\begin{align*}
& x_{c}(t)=A_{c} x_{c}(t)+B_{c} y_{1}(t)+K_{3} y_{2}(t)  \tag{5.4}\\
& u_{2}(t)=E_{c} y_{1}(t)+K_{2} x_{c}(t)+K_{1} y_{2}(t) . \tag{5.5}
\end{align*}
$$

Optimizing with respect to $K_{1}, K_{2}, K_{3}$ yields

$$
\begin{gather*}
0=C_{2}(Q P+Q P+Q P) B_{2},  \tag{5.6a}\\
0=Q P B_{2},  \tag{5.6b}\\
0=C_{2} Q P, \tag{5.6}
\end{gather*}
$$

which imply

$$
\begin{align*}
& 0=\tau_{2} \tau_{1},  \tag{5.7a}\\
& 0=\tau_{2} \tau_{3},  \tag{5.7b}\\
& 0=\tau_{3} \tau_{1} . \tag{5.7c}
\end{align*}
$$

Note that (5.7b) and (5.7c) imply

$$
\begin{equation*}
0=\tau_{2} \tau_{3} \tau_{1} . \tag{5.8}
\end{equation*}
$$

Using (5.7b), (5.7c), $Q=r_{3} Q, P=\mathcal{P}_{T_{3}}($ see [4]), (4.10)-(4.18) become

$$
\begin{equation*}
A_{\mathrm{c}}=\Gamma\left(A-B_{1} R_{1 s}^{-1} \vec{W}_{s}-2 V_{15}^{-1} C_{1}+B_{2} K_{1} C_{2}\right) G^{\top}+\Gamma B_{2} K_{2}-K_{3} C_{2} G^{r} . \tag{5.9}
\end{equation*}
$$

$$
\begin{align*}
& B_{\mathrm{c}}=\Gamma r_{2} \mathrm{Q}_{\mathrm{s}} \hat{\mathrm{~V}}_{\mathrm{ts}}^{-1} \text {, }  \tag{5.10}\\
& C_{\mathrm{r}}=-\boldsymbol{R}_{{ }_{s}}^{-1} \hat{\sigma}_{s} \tau_{1} G^{T},  \tag{5.11}\\
& D_{c}=-R_{1,}^{-1} \Phi_{s} Q C_{2}^{\top}\left(C_{2} Q C_{2}^{T}\right)^{-1} .  \tag{5.12}\\
& E_{\mathrm{t}}=-\left(B_{2}^{\top} P B_{2}\right)^{-1} B_{2}^{T} P \hat{Q} \hat{V}_{1 s}^{-1},  \tag{5.13}\\
& 0=\bar{A}_{1} Q+Q \bar{A}_{3}^{T}+\hat{V}_{0}+\sum_{i=1}^{n}\left[\hat{A}, Q \hat{A}_{1}^{T}+\left(\hat{A}_{1}-B_{1}, \hat{R}_{13}^{-1} \bar{D}_{2} \tau_{12}\right)\right.
\end{align*}
$$

$$
\begin{align*}
& 0=\hat{A}_{2}^{T} P+P \hat{A}_{2}+\hat{R}_{0}+\sum_{i=1}^{D}\left(\hat{A}_{i}^{I} P \hat{A}_{1}+\left(\hat{A}_{1}+\tau_{2} \hat{Q}_{2} \hat{V}_{15}^{-1} C_{1 i}\right)^{r}\right. \tag{5.14}
\end{align*}
$$

$$
\begin{align*}
& +\tau_{3_{1}}^{\top} \tau_{1}^{\top} \hat{Q}_{s}^{T} R_{11}^{-1} \hat{\phi}_{s} \tau_{1} \tau_{3_{1}}, \tag{5.15}
\end{align*}
$$

where

$$
\hat{A}_{3}=\hat{A}+B_{2} K_{1} C_{2}-G^{r} K_{3} C_{2}, \quad \hat{A_{2}}=\hat{A}+B_{2} K_{1} C_{2}+B_{2} K_{2} \Gamma
$$

$$
\hat{A}_{Q_{1}} \triangleq A_{Q}+B_{2} K_{1} C_{2}-G^{T} K, C_{2}, \dot{A}_{A_{7}} \triangleq A_{P_{3}}+B_{2} K_{1} C_{2}+B_{2} K_{2} \Gamma
$$

## VI. Directions for Further Research

More general solutions can be obtained by incorporating singular estimation techniques [15] where noise-free measurements are repeatedly differentiated to enlarge the class of available outputs.

## ACKNOWLEDGMENT

The author wishes to thank D. C. Hyland for numerous helpful discussions.

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$$
\begin{align*}
& 0=A_{A} Q+Q A_{A}^{T}+r_{2 A} Q D_{1,}^{-1} Q_{s} T_{2 L} T_{1} \\
& -\tau_{31} \tau_{21} Q D_{11}^{-1} Q_{1}^{r} \tau_{2+}^{r} \tau_{3_{1}}^{r}-G^{r} K_{1} C_{2} Q-Q\left(K_{3} C_{2}\right)^{r} G, \tag{5.16}
\end{align*}
$$

$$
\begin{align*}
& \left.-\tau_{3_{1}}^{T} \tau_{14}^{T} \otimes_{2}^{T} R_{18}^{-1} \phi_{s}\right) \tau_{14} T_{3_{1}}-\Gamma^{\top} K_{2} B_{2} P-P\left(K_{2} B_{2}\right) T . \tag{5.17}
\end{align*}
$$

## A. Riccati Equation Approach to the Singular LQG Problem

(ABSTRACT)
by

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## 1. Introduction

The singular LQG control problem has been of considerable interest for almost two decades ([1-15]). Such problems arise when some of the measurements are noise free or when some of the control signals are unweighted. This will be the case, for example, if the sensor noise is colored or if actuator dynamics are included. Augmentation of the plant dynamics by means of noise filters or actuator dynamics thus leads directly to the singular problem formulation.

Most of the literature on the singular LQG problem is based upon limiting procedures in which suitable weighting matrices and noise intensities approach zero. These results demonstrate the types of behavior which can arise in the limiting solution including impulsive controls and singular arcs.

The available literature is concerned, of course, with determining the optimal limiting (i.e., singular) control. In practical applications, however, it is often of interest to determine the optimal controller within a prespecified class of controllers. In particular, we consider the singular LQG problem in which the controller is preconstrained to possess a fixed dynamic feedback structure. One benefit of this approach is that the fixed structure constraint eliminates the possibility of impulsive controls and other complex behavior.

Preliminary results for the singular LQG problem were obtained in [15] using the fixed structure approach. For generality, the problem considered in [15] permits the design of fixed-order, i.e., reduced-order, dynamic compensators. As in [16], the solution is given by a system of coupled algebraic Riccati and Lyapunov equations whose solutions (denoted by $Q, P, \hat{Q}, \hat{P}$ ) are used to explicitly characterize the optimal feedback gains. The coupling is due to a pair of oblique projections (i.e., idempotent matrices) which arise as a direct consequence of the fixed structure constraint. The order-reduction projection $r$ defined by

$$
\tau \triangleq \hat{Q} \hat{P}(\hat{Q} \hat{P})^{\#},
$$

where ( )* denotes group generalized inverse, appeared originally in [16], while the static projection $\nu$ given by

$$
\nu \triangleq Q C^{T}\left(C Q C^{T}\right)^{-1} C
$$

is familiar from least squares analysis.
The results of [15] are incomplete, however, in that the gains associated with certain feedback paths were not given explicitly. For the corresponding singular estimation problem ([17]) this de-
fect was remedied in [18] where all feedback gains were explicitly characterized. In addition, the solution obtained in [18] was shown to agree completely with results obtained using standard limiting methods when the (unconstrained) optimal singular estimator does not possess differentiators ([19]). The results of [18] thus provide an alternative approach to the singular estimation problem considered in $[20,21,22]$ and the numerous references therein.

The contribution of the present paper is thus to complete the development of [15] by incorporating the methods used in [18]. Accordingly, we derive a coupled system of modified Riccati and Lyapunov equations which explicitly characterize the feedback gains of the fixed-structure singular LQG controller. For generality we consider partial or total singularity in both the control weighting and measurement noise intensity matrices, and we allow the dynamic compensator to be of arbitrary dimension less than or equal to the number of plant states minus the number of noisefree measurements. In the special case in which the order of the dynamic compensator is equal to the number of plant states minus the number of noise-free measurements (i.e., the quasi full-order case), then we show that the optimal solution decomposes (separates) into a reduced-order observer followed by state feedback.

An additional benefit of our approach is the ability to impose an upper bound on the number of differentiators to be included in the feedback controller. That is, while certain measurement signals may be noise free and hence differentiable, it may be undesirable in practice to implement more than one level of differentiation or, perhaps any differentiation at all. Furthermore, as in [18] we demonstrate connections with earlier results by showing that the fixed structure solution agrees with the standard limiting solution when the latter possesses the same number of differentiators as are included in the prespecified controller structure.

To illustrate the solution we consider a numerical example of fourth order with two noisefree measurements and one noisy measurement. (Numerical results for the singular control case are immediate from duality). By solving the coupled systems of modified Riccati and Lyapunov equations by means of a homotopy algorithm ([23]), we obtain the quasi full-order solution (secondorder controller) as well as an optimal first-order controller.

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## APPENDIX F: Stochastic Modeling

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# THE OPTIMAL PROJECTION/MAXIMUM ENTROPY APPRRACH <br> TO DESIGNING LOW-ORDER, ROBUST CONTROLLERS FOR FLEXIBLE STRUCTURES 

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Abstract. The Optimal Projection/Maximum Entropy approach to designing low-order controllers for high-order systems with parameter uncertainties is reviewed.
The philosophy of representing uricertain parameters by means of Stratonovich multiplicative white noise is motivated by means of the Maximum Entropy Principle of Jaynes and statistical analysis of modal systems. The main result, the optimal projection equations for fixed-order dynamic compensation in the presence of state-, control-and measurement-dependent noise, represents a fundamental generalization of classical LQG theory.

## 1. Overview

Optimal Projection/Maximum Entropy Stochastic Modelling and Reduced-Order Design Synthesis is a rigorous new approach to designing robust, implementabie feedback controllers. Inspired by Statistical Energy Analysis [1], a branch of dynamic modal analysis developed for analyzing acoustic vibrations, its present stage of development [2-22], embodies a mathematically rigorous, fundamental generalization of classical steady-state Kalman filter and linear-quadraticGaussian (LQG) optimal control theory. Al though LQG theory is an effective tool for optimally quantifying performance/sensor-resolution and performance/actuation-level tradeoffs, it suffers from two fundamental defects which severely 11 mit its usefulness in practice.

1. Whereas the dimension of an LQG controller must equal that of the controlled plant, optimal projection design characterizes the quadratically optimal controller of fixed dimension less than that of the plant in accordance wth implementation constraints (e.g., reliability, complexity or real-time computing capability).
2. Whereas LQG presumes exact knowledge of each and every parameter appearing in the statespace plant description, maximum entropy modelling provides a stochastic plant model which admits ignorance with regard to parameter values In accordance with unavoidable plant modelling errors.
with regard to the latter item, it should be stressed that one of the major problems in designing high-performance control systems is that of robustness, i.e., the ability of the controller to tolerate errors in the plant model upon which its design is predicated. Maximum entropy modelling directly addresses this problem by incorporating into the dynamic model a representation of ignorance (i.e., uncertainty) regarding physical parameters. Roughly speaking, the idea behind the approach is to use a probabi: listic representation of each imperfectly known plant parameter so that the quadratically optimal control system designed under this probabilistic model is automatically desensitized to actual parameter variations when the control system is implemented. The overall control-design procedure thus avoids laborious trial and error post-design "tweaking."

## 2. Motivation

The inherent time- and frequency-domain duality in representing innear dynamic systems (i.e.. state space versus transfer functions) provides control-system designers with complementary methodologies for assessing tradeoffs between performance objectives and the design constraints of sensor resolution, actuation levels, plant modelling accuracy and controller complexity. In spite of the ability of LQG to optimally quantify performance/sensor-resolution and performance/ actuation-level tradeoffs in a state-space setting, its enormous sensitivity to plant modelling errors has forced practitioners to seek generalizations of classical frequency domain methods. In numerous practical situations, however, input/output techniques possess fundamental 1 imitations. For example, representing modelling uncertainty in a frequency-domain plant model $G(s)$ by means of

$$
G(s)+\Delta G(s)
$$

where $\Delta G$ remains in a normed neighborhood of $G$, is essentially a black -box (nonparametric) approach: By failing to exploit physical laws (such as conservation of energy), systems represented by $G+\Delta G$ may actually be physically impossible, resulting in unwarranted design conservatism at the expense of system performance. Hence, when some knowledge of internal mechanisms is avallable (1.e., the "grey-box" situation), state-space representations may provide greater modelling fidelity. These observations are motivated by the probiem of controliing vibration in flexible structures
where internal energy dissipation preciudes right-half-plane poles and where nigh-order finite -1 ement models have highly structured dynamics but possess numerous uncertain parameters. More specifically, frequency uncertainties for higher order modes are much larger in magnitude than damping uncertainties. Hence, the inability to differentiate between these physical parameters in an input-output representation leads to severe performance consequences.

The optimal projection/maximum entropy approach generalizes LQG theory in two fundamental respects: design of reduced-order controllers plus accommodation of a priori parameter uncertainties. For clarity, we discuss these generalizations separately following the left branch of Fig. 1. Optimal projection design is discussed in Section 3, followed by maximum entropy modelling in Sections 4, 5, and 6.

## OPTIMAL PROJECTION/MAXIMUM ENTROPY APPROACH TO LOW-ORDER, ROBUST CONTROLLER DESIGN



Fig. 1

## 3. Review of the Optimal Projection Approach

Most research into the design of reduced-order controllers involves one of two sequential procedures: model reduction followed by controlier design, or controller design followed by controller reduction. The optimal projection equations represent a radical departure from both of these approaches by directly characterizing the quadratically optimal reduced-order controller for a high-order model. Assuming a purely dynamic linear structure for the desired compensator, whose order is determined by implementation constraints, parameter optinization approach is taken. There is, of course, nothing novel about this approach per se and it has been widely studied in the control ifterature (see. -iferences listed in [18]). This approach, however, fell into disrepute because of the extreme complexity of the grossly unwieldy first-order necessary conditions with afforded ifttle insight and engendered brute-force gradient search techniques. The crucial discovery occurred [7] where it was revealed that the necessary conditions for the dynamiccomoensation problem give rise to the definition
of an optimal projection as a rigorous, unassailable consequence of quadratic optimality without recourse to ad hoc methods. Exploitation of this projection leads to immense simplification of the "primitive" form of the necessary conditions which now provide a transparent generalization of the pair of separated Riccati equations of standard LQG theory. In particular, the optimal projection equations comprise a system of four matrix equations coupled by an oblique projection which determines the optimal controller gains. The system of matrix equations includes a pair of modified Riccati equations which are analogous to the standard Riccati equations, along with a pair of modified Lyapunov equations which arise separately in the model reduction problem [19]. The coupling by means of the projection reveals the inherent
inseparability of these operations in the reduced-order case since optimality considerations demand that, in a very precise sense, "reduction" and "control design" be performed simultaneousiy. Hence the full-arder model is retained throughout the control design process and there is no need to truncate the plant model.

## 4. Maximum Entropy Modelling

Although optimal projection design deals directly and rigorousiy with the question of system dimension by trading order off against performance, it is, nevertheless, predicated upon the avallability of a completely accurate plant and disturbance model. Maximum entropy modelifing, however, addresses the robustness problem by directiy including parameter uncertainties in the plant and disturbance models so that optimal projection design plus maximum entropy modelling automatically yíelds control designs that trade performance off against modelifing uncertainties. In order to review the maximum entropy approach it is important to discuss the class of problems that motivated this work, namely, control of flexible structures. A finite element model of a large fiexible structure is, generally, an extremely high-order system. For example, a version of the widely studied Draper Model $\$ 2$ includes 3 rigid body modes, 147 elastic modes and 6 disturbance states, i.e., a total of 306 states, along with 9 sensors and 9 actuators. Besides the high order of these systems, finite element modelling is known to have poor accuracy, particularly for the high-order modes.
Reasonable and not overly conservative uncertainty estimates predict 30-50 percent error in modal frequencies after the first 10 modes, with the situation considerably more complex (and pessimistic) for damping estimates.

Maximum entropy modelling is a form of stochastic modelling. Al though external disturbances are traditionally modelled as random processes, the use of stochastic theory to model plant parameter uncertainty has seen relatively 1 imited application. To dispel all objections to a stochastic parameter-uncertainty model, we invoke the modern information-theoretic interpretation of probability theory. Rather than regard the probability of an event as an objective quantity such as the 1 imiting frequency of outcomes of .umerous repetitions (as, e.g., the number of heads in 1,000 coin tosses), we adopt the view that the probability of an event is a subjective quantity which reflects the observer's certainty as to
whether a particular event will or will not occur. This quantity is nothing more than a measure of the information (including, e.g., all theoretical analysis and empirical data) avallable to the observer. In this sense the validity of a stochastic model of a flexible structure, for example, does not rely upon the existence of a fleet of such objects (substitute "ensemble" for "fleet" in the classical terminology) but rather resides in the interpretation that it expresses the engineer's certainty or uncertainty regarding the values of physical parameters such as stiffnesses of structural components. This view of probability theory has its roots in Shannon's information theory but was first articulated unambiguously by Jaynes [23-26].

The preeminent problem in modelling the real world is thus the following: Given limited (incomplete) a priori data, how does one construct a well-defined (complete) probability model which is consistent with the avaliable data but which avoids inventing data which does not exist? To this end we invoke Jaynes Maximum Entropy Principle: First, define a measure of ignorance in terms of the information-theoretic entropy, and then determine the probability distribution which maximizes this measure subject to agreement with the available data. The reasoning behind this principle is that the probability distribution which maximizes the a priori ignorance must be the least presumptive (i.e., least likely to invent data) on the average since the corresponding amount of a posteriori learned information (should all uncertainty suddenly disappear) would necessarily be maximized. If, for some probability distribution, the a priori ignorance and hence the a posteriort learning were less than their potentially maximum value, then this distribution must be based upon invented and hence generally incorrect data. The Maximum Entropy Principle is clearly desirable for control-system design where the introduction of false data is to be assiduously avoided.

## 5. Minimum-Information Modelling of Parameter Uncertaintles

For dynamical-system modelling it was first shown by Hyland [2] that for structural systems the minimum information linear stochastic dynamic model induced by the Maximum Entropy Principle of Jaynes is a Stratonovich multiplicative white noise model. In the present paper we adopt this model and explore its ramifications for general systems. The basic model is given by

$$
\begin{equation*}
\dot{\tilde{x}}(t)=\left(\tilde{A}+\sum_{i=1}^{p} v_{i}(t) \tilde{A}_{i}\right) \tilde{x}(t)+\tilde{w}(t), \tag{5.1}
\end{equation*}
$$

where $\tilde{x}(t) \in \underline{R}^{\tilde{n}}, \tilde{A}, \tilde{A_{1}} \in \underline{R}^{\tilde{n} \times \tilde{n}}, \tilde{W}(t)$ is zeromean Gaussian white disturbance goise with non-negative-definite intensity $\widehat{\nabla}$, and $v_{i}(t)$ are zero-mean, unit-intensity Gaussian white noise processes which are mutually uncorrelated and uncorrelated with $\bar{W}(t)$. The multiplicative white nolse model (5.1) can be regarded as a parameter uncertainty model where each $v_{i}(t)$ corresponds to a single uncertain parameter whose pattern and magnitude are given by $\hat{X}_{1} /\left\|A_{i}\right\|$ and $\left\|\widetilde{A}_{1}\right\|$. respectively.

To see why (5.1) is a minimum information model of parameter uncertainty, note that when the pattern $\mathbb{X}_{1} /\left\|\mathbb{K}_{1}\right\|$ of an uncertain parameter
is known all available data (theoretical and empirical) can be used to determine a suitable value for the magnitude $\left\|\tilde{A}_{i}\right\|$ to reflect the corresponding level of uncertainty. Clearly, the collection of magnitudes constitutes the minimum data set needed to render (5.1) well defined. For the harmonic oscillator with uncertain natural frequency, the uncertainty magnitude is given by the reciprocal of the decorrelation time (Fig. 2). Note that the uncertainty representation (5.1) is a minimum information model in the sense that it eschews detailed descriptions of joint probability statistics of unknown parameters.

## MINIMUM-INFORMATION MODELLING DECORRELATION TIME





Fig. 2

To eliminate the white noise formalism, the model (5.1) is usually rigorized by the Ito differential equation

$$
\begin{equation*}
d \tilde{x}_{t}=\left(\tilde{A} d t+\sum_{i=1}^{p} d v_{i t} \tilde{A}_{i}\right) \tilde{x}_{t}+d \tilde{w}_{t} \tag{5.2}
\end{equation*}
$$

where dvit and dwt are Brownian motions, 1.e., Wiener processes. Although such models were studied extensively for control design, this approach fell into disrepute with the pubiication of $[27,28]$ where it was shown for discrete-time systems that sufficiently high uncertainty levels (i.e., magnitudes $\| \widetilde{A}_{j}| |$ above a "threshold")
lead to the nonexistence of a steady-state solution. Although it was purported that this "phenomenon" was an "obvious" consequence of high uncertainty levels, these conclusions failed to take into account (possibly because of the discrete-time setting) the subtle relationship between the ordinary differential equation (5.1) and the stochastic differential equation (5.2). Indeed, it was shown in [29] that if a stochastic differential equation is regarded as the 1 imit of a sequence of approximating ordinary differential equations, then (5.2) is not the correct version of (5.1).

Instead, the ordinary differential equation (5.1) with multiplicative white noise corresponds to the corrected Ito differential equation

$$
\begin{equation*}
d \tilde{x}_{t}=\left(\tilde{A}_{s} d t+\sum_{i=1}^{p} d v_{i t} \tilde{A}_{i}\right) \tilde{x}_{t}+d \tilde{X}_{t} \tag{5.3}
\end{equation*}
$$

where

$$
\begin{equation*}
\tilde{A}_{s}=\tilde{A}+\frac{1}{2} \sum_{i=1}^{p} \tilde{A}_{i}^{2} \tag{5.4}
\end{equation*}
$$

which differs from the "naive" equation (5.2) by a systematic drift term. The form of (5.3) was corroborated completely independently by Stratonovich [30], whose results actually appeared in the Russian literature prior to 1965. His approach is based upon an alternative definition of the stochastic integral which differs from the lto definition by a mathematical technicality.

In spite of the glaring technicality of the Stratonovich correction, almost all research on the estimation and control of such systems failed to perceive its physical significance. Specifically, the stratonovich correction neutralizes the "threshold uncertainty principle": For systems which are inherently stable under particular parameter variations (e.g., structures with uncertain stiffness matrices), the Stratonovich formulation correctly predicts unconditional second-moment stability in contrast to the Ito formulation within which a stringent uncertainty threshold is encountered. We shall now proceed to demonstrate this fact by means of a compelling example relevant to the modeling and control of flexible structures, in particular, and hyperbolic systems, in general.
First, suppose that zero-point deviations of $\widetilde{x}(t)$ are of interest and are evaluated according to

$$
\begin{equation*}
J=\lim _{t \rightarrow \infty} E \bar{X}(t)^{T} \bar{R} \bar{X}(t)=\lim _{t \rightarrow \infty} \operatorname{tr} \bar{Q}(t) \bar{R}, \tag{5.5}
\end{equation*}
$$

where $\tilde{R} \in \mathbb{R}^{\tilde{\pi} \times \tilde{f}}$ and the second moment of the state is

$$
\begin{equation*}
\tilde{o}(t) \stackrel{\hat{N}}{E}\left[\tilde{x}(t) \tilde{x}(t)^{\top}\right] . \tag{5.6}
\end{equation*}
$$

The obvious fact cannot be overemphasized that the primary state statistic of design interest in linear-quadratic optimization is the state covariance (5.6). From Ito calculus it follows that $\bar{Q}(t)$ is given for the naive model (5.2) by

$$
\begin{equation*}
\dot{\tilde{Q}}(t)=\tilde{A} \bar{O}(t)+\tilde{Q}(t) \tilde{A}^{\top}+\sum_{i=1}^{p} \tilde{A}_{i} \tilde{Q}(t) \tilde{A}_{i}^{\top}+\tilde{v} \tag{5.7}
\end{equation*}
$$

and for the corrected model (5.3) by

$$
\begin{equation*}
\dot{\tilde{Q}}(t)=\tilde{x}_{s} \tilde{Q}(t)+\widetilde{Q}(t) \tilde{A}_{s}^{\top}+\sum_{i=1}^{p} \tilde{A}_{i} \tilde{Q}(t) \tilde{\lambda}_{i}^{\top}+\tilde{v} . \tag{5.8}
\end{equation*}
$$

Each of these "stochastic" Lyapunov differential equations, wich govern the evolution of the second moment, should be regarded as $n(n+1) / 2$ ordinary differential equations. Hence we wish
to address the following question: How do the solutions of the stochastic Lyapunov equations (5.7) and (5.8) differ from each other and from the "deterministic" Lyapunov equation

$$
\begin{equation*}
\dot{\tilde{\sigma}}(t)=\tilde{A} \tilde{Q}(t)+\tilde{Q}(t) \tilde{R}+\tilde{\nabla} \tag{5.9}
\end{equation*}
$$

particularly in the presence of high uncertainty levels? The answer to this question of course depends upon the stochastic modification terms which for the naive model are given by

$$
\begin{equation*}
\underline{M}_{1}[\widetilde{Q}(t)] \triangleq \sum_{i=1}^{p} \tilde{A}_{i} \widetilde{Q}(t) \tilde{A}_{i}^{\top} \tag{5.10}
\end{equation*}
$$

and for the corrected model by
$M_{s}[\tilde{Q}(t)]=\sum_{i=1}^{P}\left[\frac{1}{2} \tilde{A}_{i}^{2} \tilde{Q}(t)+\frac{1}{2} \tilde{Q}(t) \tilde{A}_{i}^{2 T}+\tilde{A}_{i} \tilde{Q}(t) \tilde{A}_{i}^{T}\right]$.
Consider a system consisting of a pair of lightly damped modes so that

$$
\tilde{A} \equiv\left[\begin{array}{cccc}
0 & -\omega_{1} & 0 & 0 \\
\omega_{1} & -2 \eta_{1} & 0 & 0 \\
0 & 0 & 0 & -\omega_{2} \\
0 & 0 & \omega_{2} & -2 \eta_{2}
\end{array}\right] \text {. }
$$

where $\eta_{i}=\boldsymbol{s}_{1} \omega_{1}$, and to represent frequency uncertainties let

$$
\begin{aligned}
& \tilde{A}_{1}=\gamma_{1}\left[\begin{array}{rrrr}
0 & -1 & 0 & 0 \\
1 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0
\end{array}\right], \\
& \tilde{A}_{2}=\gamma_{2}\left[\begin{array}{rrrr}
0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 \\
0 & 0 & 0 & -1 \\
0 & 0 & 1 & 0
\end{array}\right],
\end{aligned}
$$

where for simplicity we have 1 gnored the effects of frequency uncertainties on the effective decay rate $\eta_{i}$. The magnitudes of the uncertainties are scaled by means of $\gamma_{1}$ and $\gamma_{2}$. For this example the Ito stochastic modification $M_{I}[\widetilde{Q}(t)]$ has the form

Although the off-diagonal terms have a stabilizing effect, it is clear that the diagonal elements destabilize the state variances. Hence, it is not surpising that for sufficiently high uncertainty levels, i.e., $\gamma_{i} \gg 0$, the [to mode]
is second-moment unstable. These observations are completely in accordance with the threshold uncertainty principle. The Stratonovich corrected stochastic modification Ms[ $\mathrm{M}(t)]$. however, has the form
which also has stabilizing off-diagonal elements but has fundamentally different diagonal elements: Rather than destabilizing the state variances, the diagonal elements of the corrected stochastic modification are equilibrating. This effect is even more striking when $M_{I}$ and $M_{S}$ are transformed into the basis with respect to which
$\tilde{A}=\left[\begin{array}{cccc}-j \omega_{1}-\eta_{1} & 0 & 0 & 0 \\ 0 & j \omega_{1}-\eta_{1} & 0 & 0 \\ 0 & 0 & -j \omega_{2}-\eta_{2} & 0 \\ 0 & 0 & 0 & j \omega_{2}-\eta_{2}\end{array}\right]$,
where higher order terms in $\eta$ have been ignored. In this basis, the diagonal terms of $M_{I}[\phi(t)]$ are destablifing whereas the diagonal terms of $\underline{M}_{S}[\widetilde{\sigma}(t)]$ exactly vanish.

The negative coefficients in the off-diagonal terms imply progressive decorrelation between pairs of dynamical states. This informational or statistical damping phenomenon is a direct resuit of parameter uncertainties captured by the mulitplicative white noise model. The Stratonovich correction, moreover, is crucial: By neutralizing the threshold unceriainty principle, it permits the consideration of long-term effects for arbitrary uncertainty levels.

As an example of the ramifications of these observations, assume las is usually the case in practice) that uncertainties in modal frequency obtained from finite-element analysis of a flexible structure increase with mode number. From the form of $M_{s}[\hat{O}(t)]$ it is easy to deduce that the steady-state covariance

$$
\tilde{\alpha^{2}} \lim _{t \rightarrow \infty} \delta(t)
$$

satisfying

$$
\begin{equation*}
0=\lambda_{s} \widetilde{Z}+\partial \tilde{A}_{s}^{T}+\sum_{i=1}^{p} \chi_{i} \tilde{\alpha}_{i}^{T}+\tilde{\gamma} \tag{5.12}
\end{equation*}
$$

becomes increasingly diagonally dominant with increasing frequency and thus assumes the qualitative form given in Fig. 3. The benefits of this sparse form are important: The computational effort required to determine the steadystate covariance (and thus to design a closedloop controller, for example) is directly proportional to the amount of information reposed in the model or, equivalently, inversely proportional to the level of modelled parameter uncertainty. This casts new light on the computational design burden vis-a-vis the modelling question: The computational burden depends only upon the information actually available. A simple control-design exercise involving full-state feedback illustrates this point. The gains for the higher order modes of the beam in Fig. 4, whose frequency uncertainties increase linearly with frequency, were obtained with modest computational effort in spite of $\pi=100$ (see Fig. 5). Another important ramification of the qualitative form of $\mathcal{Q}$ is the automatic generation of a high-/low-authority control law. Note that for the higher order and hence highly uncertain modes the control gains Indicate an inherently stable, low-performance rate-feedback control law, whereas for the lowest-order modes the control law is high authority, i.e.. "LQ" in character.

## EFFECT OF FREQUENCY UNCERTAINTIES ON THE QUALITATIVE STRUCTURE OF THE STEADY-STATE COVARIANCE $0=\lim _{t \rightarrow \infty} E\left[x(t) x(t)^{T}\right]$



INFORMATION REGIMES
Fig. 3

FULL-STATE FEEDBACK CONTROL SCHEME


Fig. 4

STOCHASTIC BEAM PROBLEM-GAIN MAGNITUDES


Fig. 5
6. Optimal Projection Design with Stratonovich Multiplicative White Noise
To state the Reduced-Order Dynamic Compensation Problem, we require the following notation. Let
$x: \underline{R}^{n}, y \in \underline{R}^{2}, u \in \mathbb{R}^{m}, A, A, \ldots, A_{p} \in \mathbb{R}^{n \times n}$,
$B, B_{1}, \ldots, B_{p} \in \mathbb{R}^{n \times m}, C, C_{1}, \ldots, C_{p} \in \mathbb{R}^{l \times n}, n_{c} \leq n$,
$x_{c} \in R^{n^{n}}{ }^{c}, A_{c} \in R^{n^{n} c^{x n}} c, B_{c} \in R^{n^{n} c^{x \ell}}, C_{c} \in R^{m \times n} c$,
$R_{1} \in R^{n \times n}, R_{1}>0, R_{2} \in R^{m \times m}, R_{2}>0, R_{12} \in R^{n \times m}$.

Furthermore, let $v_{1}, \ldots, v_{p}$ be unit-intensity, zero-mean and mutually uncorrelated white notse processes and let $w_{1} \in \mathbb{R}^{n}$ and $w_{2} \in R^{P}$ be zero-mean white noise processes with intensities $v_{1} \geq 0$ and $v_{2}>0$, respectively, and crossintensity $V_{12} \in R^{n x l}$. Assume that $V_{i}, w_{i}$ and $x(0)$ are uncorrelated. We require the technical assumption that, for each $1, B_{i} \neq 0$ fmplies $C_{i}=0$, i.e., the control - and measurement-dependent noises are uncorrelated.

Optimal Reduced-Order Dynamic-Compensation Problem. Given the controlled system

$$
\begin{align*}
& \dot{x}=\left(A+\sum_{i=1}^{p} v_{i} A_{i}\right) x+\left(B+\sum_{i=1}^{p} v_{i} B_{i}\right) u+w_{1},  \tag{6.1}\\
& y=\left(C+\sum_{i=1}^{p} v_{i} C_{i}\right) x+w_{2} \tag{6.2}
\end{align*}
$$

design an $n_{c}$ th-order dynamic compensator

$$
\begin{align*}
& \dot{x}_{c}=A_{c} x_{c}+B_{c} y_{1}  \tag{6.3}\\
& u=C x \tag{6.4}
\end{align*}
$$

which minimizes the performance criterion

To guarantee that $J$ is finite and independent of initial conditions, we restrict ( $A_{c}, B_{c}, C_{c}$ )
to the (open) set of second-moment-stabilizing triples
$\underline{s} \hat{\sum}\left\{\left(A_{c}, B_{c}, C_{c}\right): \tilde{A}_{s} \oplus \tilde{A}_{s}+\sum_{i=1}^{p} \bar{A}_{i} \otimes \tilde{A}_{i}\right.$ is stable and $\left(A_{c}, B_{c}, C_{c}\right)$ is controllable and observable $\}$, where $\oplus$ and denote. Kronecker sum and product and $\left.\tilde{A}_{s} \hat{\hat{e}^{A_{s}}} \begin{array}{ll}B_{s} C_{c} \\ B_{c} C_{s} & A_{c}\end{array}\right], \quad \tilde{A}_{i} \hat{\hat{人}}\left[\begin{array}{ll}A_{i} & B_{i} C_{c} \\ B_{c} C_{i} & 0\end{array}\right]$, $A_{s}{ }_{s} A+\frac{1}{2} \sum_{i=1}^{p} A_{i}^{2}, B_{s} \hat{S} B+\frac{1}{2} \sum_{i=1}^{p} A_{i} B_{i}, C_{s} \hat{C} C+\frac{1}{2} \sum_{i=1}^{p} C_{i} A_{i}$.

Call a square matrix positive semisimple if it has positive eigenvalues and diagonal Jordan canonical form, i.e., if it is similar to a positive-diagonal (or, equivalently, a positive-definite) matrix. The following lemma is proved in [19].

Lemma 6.1. If $n \times n \hat{\hat{O}} \hat{\hat{P}}$ are non-negative definite and rank $\theta \rho^{\circ}=n_{c}$ then there exist $n_{c} \times n$ G, $\Gamma$ and $n_{c} \times n_{c}$ positi. .-semisimple $M$ such that

$$
\begin{align*}
& \hat{Q} \hat{P}=G T_{M \Gamma}  \tag{6.6}\\
& \Gamma G T=I_{n_{c}} . \tag{6.7}
\end{align*}
$$

For convenience in stating the main result, we shall refer to $G, M$ and $\Gamma$ satisfying ( 6.6 ) and (6.7) as a projective factorization of $\hat{0}$.

For convenience in stating the optimality conditions, define the following notation for
Q. $P, \hat{Q}, \hat{P} \in R^{n \times n}$ :

$$
\begin{aligned}
& R_{2 s} R_{2}+\sum_{i=1}^{P} B_{i}^{T}\left(P+\hat{P}^{P}\right) B_{i} \\
& V_{2 s}=V_{2}+\sum_{i=1}^{P} C_{i}(Q+\hat{Q}) C_{i}^{T}
\end{aligned}
$$

$$
\begin{aligned}
& {\underset{f}{P}}_{P_{s}} B_{s}^{T}+R_{12}^{T}+\sum_{i=1}^{P} B_{1}^{T}(P+\hat{P}) A_{1}= \\
& A_{0 s} A_{s}-g_{s} V_{2 s}^{-1} C_{s} \text {. } \\
& A_{p s} E A_{s}-8_{s} R_{2 s}^{-1} p_{s} \text {. } \\
& T E G^{\top} \Gamma, T_{1}=I_{n}=T .
\end{aligned}
$$

Theorem 6.1. Suppose $\left(A_{C}, B_{C}, C_{C}\right) \in S$ solves The optimal dynamic-compensation problem. Then there exist $n \times n$ non-negative-definite matrices $Q$, $P$. 6 and $\%$ such that, for some projective factorization, $G, M, \Gamma$ of $Q P, A_{C}, B_{C}$, and $C_{C}$ are given by

$$
\begin{align*}
& A_{c}=\Gamma\left(A_{s}-B_{s} R_{2 s-s}^{-1} P_{-}-Q_{s} V_{2 s}^{-1} C_{s}\right) G^{\top},  \tag{6.8}\\
& B_{c}=\Gamma Q_{s} V_{2 s}^{-1},  \tag{6.9}\\
& C_{c}=-R_{2 s}^{-1} P_{s} G^{\top}, \tag{6.10}
\end{align*}
$$

and such sat the following conditions are satisfied:

$$
\begin{align*}
0 & =A_{s} Q+Q A_{s}^{T}+V_{i}+\sum_{i=1}^{P}\left[A_{i} Q A_{i}^{T}+\left(A_{i}-B_{i} R_{2 s}^{-1} P_{s}\right) \hat{Q}\left(A_{i}-B_{i} R_{2 s}^{-1} P_{s}\right)^{T}\right] \\
& -Q_{s} V_{2 s}^{-1} Q_{s}^{T}+T_{1} Q_{s} V_{2 s}^{-1} Q_{s}^{T} T_{1}^{T} \tag{6.11}
\end{align*}
$$

$0=A_{s}^{T} P+P A_{s}+R_{1}+\sum_{i=1}^{p}\left[A_{i} P A_{i}+\left(A_{i}-Q_{s} V_{2 s}^{-1} C_{i}\right)^{T} \hat{P}^{T}\left(A_{i}-Q_{s} V_{2 s}^{-1} C_{i}\right)\right]$

$$
\begin{equation*}
-P_{s}^{T} R_{2 s}^{-1} P_{s}+T_{\perp}^{T} P_{s}^{T} R_{2 s}^{-1} P_{s} T_{1} \tag{6.12}
\end{equation*}
$$

$0=A_{P_{s}} \tilde{Q}^{Q}+\hat{Q}_{A_{P s}}^{T}+Q_{s} V_{2 s}^{-1} Q_{s}^{T}-r_{1} Q_{S} V_{2 s}^{-1} Q_{s}^{T} T_{1}^{T}$,


$$
\begin{equation*}
\operatorname{rank} \hat{Q}=\operatorname{rank} \hat{F}=\operatorname{rank} \hat{Q P}=n_{c} \tag{6.14}
\end{equation*}
$$

Remark 6.1. Since $R_{2 s} \geq R_{2}$, so that $R_{2 s}^{-1} \leq R_{2}^{-T}$. it is clear that the
control dependent noise effectively suppresses
the regulator gain $C_{C}$. Similarly, since $V_{2 s} \geq V_{2}$, the measurement-dependent noise suppresses the observer gain $B_{c}$. The terms $A_{1} Q A_{1}^{\top}$ are responsible for the decorrelation effects discussed in Section 5.

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# NUMERICAL SOLUTION OF THE OPTIMAL PROJECTION/ MAXIMUM ENTROPY DESIGN EQUATIONS FOR LOW-ORDER, ROBUST CONTROLLER DESIGN 

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## ABSTMCT

This peper summarizes some recent results obtained using the optimal projection/maximum entropy control-design equations. The main reaults include: low-order controllers for CSDL Model ${ }^{2}$; robust controllere for the 8COLE and VCOSS A models with modal-frequency uncertainties; and Doyle's example.

## 1. Introduction

The optimal projection/maximm entropy design equations are discussed in [1]-[5] and a complete, self-contained derivation appears in [5]. In brief, these equations generalize classical LQG theory in two distinct ways. Firat, the controller is constrained to have a fixed, reduced order and the reaulting necessary conditions involve an oblique projection [1]. And, second, multiplicative white noise is introduced into the plant to capture the statistical effecta of parameter uncertainty. The resulting dynamical equation is interpreted according to symetric Stratonovich etochastic integration and, using the theory of stochastic approximation, has been cotivated by the maximum entropy principle of Jaynes.

## 2. CSDL $\$ 2$

The optimal projection (OP) reduced-order design equations vere solved for the 20 -atate veraion of CSDL Model 12 treated in $(6,7]$. For various controlauthority levels, $O P$ designs vere obtained for orders $n_{c}=10,6$ and 4 . Figure 1 anmarizes the results obtaiaed in [6,7]. Note that for compensator order $\mathbf{n}_{c}$ - 4, the allowable control bandwidth is severely restricted. The $O P$ designs, however, all lie vichin the shaded band close to the LQG performance over a considerably expanded range of control bandwidthe. Delative to LQG: the performance of the OP designs is given in Figure 2. Details of the numerical algorithm used to obrain these resulta are given in [8].

## 3. 8COL

The sCOLE configuration is discussed at length in [9.10]. The model utilized in [9] involvee 16 states. 12 actuators and 17 sensors. The LQG deaign reveals inetability reaulting from $5 \%$ modal-frequency perturbatione. Deing the maximumentropy (ME) design equations, a pair of controllers vere obtained in the presence of stochatically modelled modal-frequency ugertaineiee. The first design exhibite near-LQG pryformance vith $60 \%$ increase in robustrese, vhile the eecoed design is comaiderably more robuet (behaving eore like the open-loop atructure) with nominal performance within $6 \%$ of LQG. These designe (Fig. 3) illustrate the performance/robustness tradeoff capabilitien of the Me method. It should be noted
that, for lightly damped structures, significant modal-frequency uncertainty corresponds to pronounced spectral-resonance shifting. Yrequency-domin bounds for such perturbations are consequently large and hence may result in conservative performance estimates.

## 4. $y \operatorname{coss} A$

The vCOSS A model [11] is a version of CSDL model *2 involving 9 colocated sensor/actuator pairs plun 2 line-of-sight sensors. Por the 28 -state ( 14 -mode) model and corresponding 28-state LQG design obtained in [11], the sensitivity to modal-frequency perturbstion is shown in Figure 4 . Note that instability results from $3 \mathbf{Z}$ modal-frequency perturbations of one of the modelled poles. For the maximum entropy design (Figure 5) the robustness is considerably improved with approximetely 20\% performance trade. Of course, there are a continuum of intermediate designs that could be obtained for desired performance/robustness tradeoffs. The closedloop ctability margins for the full 142 -mode evaluation model are shown in Figure 6.

## 5. Doyle's Example

As final application of the ME design equations, we consider the problem used in [12] to demonstrate the lack of robustaess of LQG designs. As shown in [12] (see [4] for notation), LQG regulator: for the example

$$
\begin{aligned}
& A=\left[\begin{array}{ll}
1 & 1 \\
0 & 0
\end{array}\right], \quad B=\left[\begin{array}{l}
0 \\
b
\end{array}\right], \quad \begin{array}{ll}
\mathrm{C}=\left[\begin{array}{ll}
1 & 0
\end{array}\right], \\
v_{1}=\sigma\left[\begin{array}{ll}
1 & 1 \\
1 & 1
\end{array}\right], & v_{12}=0,
\end{array} \quad v_{2}=1 \\
& R_{1}=\rho\left[\begin{array}{ll}
1 & 1 \\
1 & 1
\end{array}\right],
\end{aligned}
$$

have arbitrarily small stability margin vith regard to variations $b+\Delta b$ when $\sigma$ and $\rho$ are sufficiently large and $b=1$. Setting $\sigma=\rho=60$, it follow that the LQG regulator is only stable for . $93 \leq b+\Delta b \leq 1.01$. Uncertainty in $b$ can be modelled by eetting $p=1$, $A_{1}=0, B_{1}=\left[\begin{array}{ll}0 & b_{1}\end{array}\right]^{T}$ and $C_{1}=0$. Solving the $\boldsymbol{R}$ deaign equations with $b_{1}=.05, .10, .15$ and .20 yields a series of incressingly robust controller designs with respect to both positive and megative variations ab (see Figures 7 and 8). Tor more detaile, see [13].


Figure 1
Optimal Projection Designs (Shaded Band) versus Suboptimal Methods


Figure 2
Optimal Projection Performance
Relative to LQG


Figure 3
LQG and Maximum Entropy Designs for SCOLE: Clased-Loop Robustness to Modal-Frequency Perturbations


Figure 4
LQG Deaign for VCOSS A: Closed-Loop Robustrean to Modal-Frequency Perturbationa


Figure 5
Maximum Entropy Design for VCOS8 A: Closed-Loop Robustnese to Modal-Frequency Perturbationa

Cl mirie


Figure 6
vCoss a re Deaiga: Stability Margin for 142-State Model


Figure 7
Robustness for Doyle's Example:
LQG plus ME Designs


Figure 8
Stability Bounde for Doyle' Example

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# The Optimal Projection Equations <br> For Reduced-Order Modelling, Estimation and Control <br> Of Linear Systems With Multiplicative White Noise ${ }^{1.2}$ 

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Communicated by J. L. Speyer

[^15]Abstract. The optimal projection equations for quadratically optimal reducedorder modelling, estimation and control are generalized to include the effects of state-, control- and measurement-dependent noise.

Key Words. Feedback, control, robust, fixed-order, optimal.

## 1. Introduction

As is well know, LQR and LQG controllers lack guaranteed robustness with respect to arbitrary parameter variations (Refs. 1 and 2). A widely studied approach to correcting this defect involves introducing noise into the plant via the imperfectly known parameters (Refs. 3-10). Intuitively speaking, the quadratically optimal feedback controller designed in the presence of such disturbances is automatically desensitized to actual parameter variations. This was demonstrated in Ref. 11 for the example given in Ref. 1.

The contribution of the present paper is a generalization of classical steady-state LQG theory to include the effects of state-, control- and measurement-dependent noise. In contrast to the classical solution involving a pair of separated Riccati equations, the necessary conditions for quadratic optimality in the presence of multiplicative white noise consist of a system of two modified Riccati equations and two modified Lyapunov equations coupled by stochastic effects. The coupling serves as a graphic portrayal of the breakdown of the separation principle in the multiplicative noise case. When the multiplicative noise terms are set to zero, the modified lyapunov equations drop out and the modified Riccati equations immediately reduce to the standard pair of separated LQG Riccati equations. Related results were obtained for the discrete-time, finite-interval problem in Ref. 10.

To attain further generality, a constraint is imposed on controller order as in Ref. 12. Hence, the results of the present paper also constitute a direct generalization of the coupled system of modified Riccati and Lyapunov equations which arise in characterizing reduced-order controllers.

For the special case of full-order compensation in the presence of state-dependent noise only, versions of these equations were discovered
 interesting difference between Refs. 13 and 14 and Ref. 15 is that Mil'stein interpreted the plant model as an Ito stochastic differential equation whereas Hyland utilized the Fisk-Stratonovich definition (Refs, 16-18). In earlier work on modelling flexible mechanical structures (Refs. 19 and 20), justification for this interpretation as an appropriate model for parameter uncertainty was based
upon the Maximum Entropy Principle of Jaynes (Ref. 21) and the theory of stochastic approximation (Ref. 22). A summary of this approach and its relationship to Refs. 23 and 24 can be found in Ref. 25 , Rigorous guarantees of robustness over a prescribed range of parameter variations have been obtained using Lyapunov functions (Refs. 26-29). Although the present paper utilizes an Ito model for simplicity, results based on Stratonovich models are readily obtained by means of standard transformations.

An immediate practical benefit of the structured form of the necessary conditions is the means for constructing numerical algorithms which differ fundamentally from gradient search techniques. One such iterative algorithm, proposed in Refs. 30-32, exploits the characterization of the oblique projection as the sum of rank-1 eigenprojections of the product of the rank-deficient "pseudogramians" satisfying the modified Lyapunov equations. As discussed in Ref. 32, the necessary conditions fail to specify which eigenprojections comprise the oblique projection; indeed, each choice may correspond to a local extremal. In practice, judicious choice of the eigenprojections can eliminate extremals with high cost and hence efficiently identify the global minimum. These issues are a reault of the reduced-order constraint only; the stochastic effects alone do not appear to introduce extremal multiplicity.

The scope of the present paper involves deriving the optimal projection equations for reduced-order modelling, estimation and control obtained in Refs. 32, 33 and 12 to include state- control- and measurementdependent noise. The main results. Theorems 2.1-2.3. present the necessary conditions for optimality as systems of 2, 3, and 4 matrix equations (modified Riccati and Lyapunov equations) coupled by both the optimal projection and stochastic effects. The necessary conditions in this generality are presented here for the first time. The dynamic compensation result supports the numerical results obtained in Refs. 11 and 34. Appendix $D$ contains the proof of Theorem 2.3; the proofs of Theorems 2.1 and 2.2 are similar and hence are omitted. Although the derivations in Refg. 32, 33 and 12 utilizing Lagrange multipliers could have been adapted to the present case, we have devised a new proof based upon Rronecker products which is thought to be more direct.

## 2. Problem Statement and Main Results

The following notation and definitions will be used throughout the paper.
$\underline{8}$
R
${ }_{R_{2}}^{\alpha \times \beta}$
$R^{\alpha}$
$I_{\alpha}$
${ }^{2}(i)$
$Z_{(i, j)}$
$z^{T}$
$z^{-T}$
$\rho(2)$
tr 2
$\|2\|$
$\operatorname{diag}\left(a_{1}, \ldots, a_{\alpha}\right)$
$\mathbf{E}_{i}$
$\Pi_{i}(\Psi)$
$\mathbf{X} \otimes \mathbf{Y}$
expected value
real numbers
$\alpha \times \beta$ real matrices
$R^{\alpha \times 1}$
$\alpha \times \alpha$ identity matrix
ith element of vector $Z$
( $i, j$ ) element of matrix $Z$
transpose of vector or matrix 2
$\left(z^{T}\right)^{-1}$ or $\left(z^{-1}\right)^{T}$
rank of matrix $Z$
trace of square matrix 2
$\left(\operatorname{tr} \mathrm{Zz}^{\mathrm{T}}\right)^{\frac{1}{2}}$ (Frobenius norm)
$\alpha \times \alpha$ diagonal matrix with listed diagonal elements
matrix with unity in the (i,i) position and zeros elsewhere
$\Psi_{E} \mathbf{i}^{-1}$

$$
\left[\begin{array}{ccc}
X_{(1,1)}^{Y} & \cdots & X_{(1, \beta)^{Y}} \\
\vdots & & \vdots \\
X_{(\alpha, 1)^{Y}} & \cdots & X^{Y} \\
(\alpha, \beta)^{Y}
\end{array}\right]
$$

$X \in \mathbb{R}^{\alpha \times \beta}, Y \in R^{\gamma \times \delta}$ (Rronecker product. Refe. 35 and 36)

| $X \oplus Y$ | $\begin{aligned} & X \otimes I_{\beta}+I_{\alpha} \otimes Y \\ & X \in R^{\alpha \times a}, Y \in R^{\beta \times \beta} \quad \text { (Rronecker sum) } \end{aligned}$ |
| :---: | :---: |
| $2^{\frac{3}{3}}$ | group generalized inverse (Ref. 37) |
| row $_{i}(z)$ | ith row of matrix $Z$ |
| $\operatorname{col}_{i}(2)$ | ith column of matrix 2 |
| vec (z) | $\left[\begin{array}{c} \operatorname{col}_{1}(z) \\ \cdot \\ \cdot \\ \operatorname{col}_{\beta}(z) \end{array}\right] \in \mathbb{R}^{\alpha \beta \beta}, z \in \mathbb{R}^{\alpha \times \beta}$ |
| $\operatorname{vec}_{(\alpha, \beta)}^{-1}(z)$ | $\left[\begin{array}{llll}\mathrm{z}_{(1)} & \cdots & \mathrm{z}_{(\alpha \beta-\alpha+1)} \\ \bullet & & \bullet \\ \bullet & & \cdot \\ \mathrm{z}_{(\alpha)} & \cdots & \cdot & \mathrm{z}_{(\alpha \beta)}\end{array}\right] \in \mathrm{R}^{\alpha \times \beta}, \mathrm{z} \in \mathrm{R}^{\alpha \beta}$ |
| stable matrix | matrix with eigenvalues in open left half plane |
| nonnegative-definite matrix | symmetric matrix with nonnegative eigenvalues ( $z \geq 0$ ) |
| positive-definite matrix | symmetric matrix with positive eigenvalues ( $z>0$ ) |
| semisimple eigenvalue | eigenvalue with equal algebraic and geometric multiplicity |
| simple eigenvalue | eigenvalue with unity algebraic multiplicity |
| group-invertible matrix | matrix $Z$ satisfying $\rho(z)=\rho\left(z^{2}\right)$, i.e., matrix which is either invertible or whose zero eigenvalue is semisimple (Ref. 37) |
| semisimple matrix | matrix with semisimple eigenvalues (i.e., nondefective matrix) |
| real-semisimple matrix | semisimple matrix with real eigenvalues |
| nonuegative-semisimple matrix | semisimple matrix with nonnegative eigenvalues |

> positive-semisịmple matrix simple matrix
> $r$
> $n, m, l, n_{m}, n_{e}, n_{c}, k, p, q$
> $n_{r}$
> $x, x_{m}, x_{e}, x_{c}$
> u, y, ye
> $A, A_{1}, \ldots, A_{p}$
> $B, B_{1}, \ldots, B_{p}$
> $C, C_{1}, \ldots, c_{p}$
> $A_{m}, B_{m}, C_{m}$
> $A_{e}, B_{e}, C_{e}$
> $A_{c}, B_{c}, C_{c}$
> $\mathrm{R}, \mathrm{H}, \mathrm{R}_{2}$
> ${ }^{R} 1$
> ${ }^{\mathrm{R}_{12}}$
> $\tilde{\mathbf{R}}$
> $L$
> $\nabla_{1 t}, \cdots, v_{p t}, w_{1 t}, \cdots, w_{q t}$
> $w_{t}$
> $G, G_{1}, G_{2}$
semisimple matrix with positive eigenvalues
matrix with distinct (i.e,, simple) eigenvalues
generic subscript denoting $m$, e or $c$ positive integers; $1 \leq n_{m}, n_{e}, n_{c} \leq n ; \ell, m \leq q$ ${ }^{n+n_{r}}$
$n, n_{m}, n_{e}, n_{c}$-dimensional vectors
$m, l, k$-dimensional vectors
nxin matrices
nxm matrices
lxn matrices
$n_{m} \times n_{m}, n_{m} \times m, l \times n_{m}$ matrices
$n_{e} \times n_{e}, n_{e} \times l, k \times n_{e}$ matrices
$n_{c} \times n_{c}, n_{c} \times \ell, m \times n_{c}$ matrices

nxn nonnegative-definite matrix
uxm matrix; $R_{1}-R_{12} R_{2}^{-1} R_{12}^{T} \geq 0$
$\left[\begin{array}{ll}R_{1} & R_{12} \\ R_{12}^{T} & R_{2}\end{array}\right] \geq 0$
kxin matrix
standard, independent Wiener processes, $t \geq 0$
$\left(w_{1 t}, \ldots, w_{q}\right)^{T}$
$\mathbf{m \times q}, \mathrm{n} \times \mathrm{q}, \ell \times \mathrm{q}$ matrices; $\rho(G)-m, \rho\left(G_{2}\right)=\ell$

$$
\begin{aligned}
& \nabla_{1} \triangleq G_{1} G_{1}^{T} \geq 0, \quad \nabla_{12} \triangleq G_{1} G_{2}^{T}, \quad \nabla_{2} \triangleq G_{2} G_{2}^{T}>0, \quad \nabla \geqslant G G^{T}>0 . \\
& \tilde{G} \triangleq\left[\begin{array}{l}
G_{1} \\
G_{2}
\end{array}\right], \quad \tilde{\nabla} \triangleq \tilde{G G}^{T}=\left[\begin{array}{ll}
\nabla_{1} & \nabla_{12} \\
\nabla_{12}^{T} & \nabla_{2}
\end{array}\right] \geq 0, \\
& \tilde{B}_{\mathrm{m}} \triangleq\left[\begin{array}{l}
B_{1} \\
B_{m}
\end{array}\right], \quad \bar{B}_{e} \triangleq\left[\begin{array}{ll}
I_{n} & 0 \\
0 & B_{e}
\end{array}\right], \quad \tilde{B}_{c} \triangleq\left[\begin{array}{ll}
I_{n} & 0 \\
0 & B_{c}
\end{array}\right],
\end{aligned}
$$

$$
\begin{aligned}
& \bar{c}_{m} \triangleq\left[\begin{array}{ll}
c & -c_{m}
\end{array}\right], \quad \tilde{c}_{e} \triangleq\left[\begin{array}{ll}
L & -c_{e}
\end{array}\right], \quad \bar{c}_{c} \triangleq\left[\begin{array}{ll}
I_{n} & 0 \\
0 & c_{c}
\end{array}\right],
\end{aligned}
$$

$$
\begin{aligned}
& A \triangleq A \oplus A+\sum_{i=1}^{p} A_{i} \otimes A_{i}, \quad \tilde{x}_{r}=\left[\begin{array}{l}
x \\
x_{r}
\end{array}\right],
\end{aligned}
$$

$$
\begin{aligned}
& \tilde{A}_{m i} \triangleq\left[\begin{array}{ll}
A_{i} & 0 \\
0 & 0
\end{array}\right] \quad \tilde{A}_{e i} \triangleq\left[\begin{array}{ll}
A_{i} & 0 \\
B_{e} C_{i} & 0
\end{array}\right], \quad \tilde{A}_{c i} \triangleq\left[\begin{array}{ll}
A_{i} & B_{i} C_{c} \\
B_{c} C_{i} & 0
\end{array}\right],
\end{aligned}
$$

For the following definitions, let $Q, P, \hat{Q}, \hat{P} \in \mathbb{R}^{n \times n}$ :

$$
\begin{array}{ll}
\hat{R}_{2} \triangleq R_{2}+\sum_{i=1}^{P} B_{i}^{T}(P+\hat{P}) B_{i}, & \hat{\nabla}_{2} \triangleq \nabla_{2}+\sum_{i=1}^{P} C_{i}(Q+\hat{Q}) C_{i}^{T}, \\
Q \triangleq Q C^{T}+\nabla_{12}+\sum_{i=1}^{P} A_{i}(Q+\hat{Q}) C_{i}^{T}, & \underline{P} \triangleq B^{T} P+R_{12}^{T}+\sum_{i=1}^{P} B_{i}^{T}(P+\hat{P}) A_{i}, \\
A_{Q} \triangleq A-Q \hat{V}_{2}^{-1} C_{0} & A_{P} \triangleq A-B \hat{R}_{2}^{-1} \underline{P} .
\end{array}
$$

Using the above notation we can state the reduced-order modelling, estimation and control problems.

Reduced-Order Modelling Problem. Given the nth-order model

$$
\begin{align*}
& d x_{t}=A x_{t} d t+\sum_{i=1}^{p} A_{i} x_{t} d v_{i t}+G_{1} d w_{t^{\prime}}  \tag{1}\\
& y_{t}=C x_{t} \tag{2}
\end{align*}
$$

where $t \in[0, \infty)$, determine an $n_{m}$ th-order model

$$
\begin{align*}
& d x_{m t}=A_{m} x_{m t} d t+B_{m} G d w_{t},  \tag{3}\\
& y_{m t}=C_{m} x_{m t}, \tag{4}
\end{align*}
$$

which minimizes the model-reduction criterion

$$
\begin{equation*}
J_{m}\left(A_{m}, B_{m}, C_{m}\right) \triangleq \lim _{t \rightarrow \infty} \sup \underline{E}\left[\left(y_{t}-y_{m t}\right) T_{\left.R\left(y_{t}-y_{m t}\right)\right] .}\right. \tag{5}
\end{equation*}
$$

Reduced-Order State-Estimation Problem. Given the nth-order observed system

$$
\begin{align*}
& d x_{t}=A x_{t} d t+\sum_{i=1}^{p} A_{i} x_{t} d v_{i t}+B G d w_{t}  \tag{6}\\
& d y_{t}=C x_{t} d t+\sum_{i=1}^{p} C_{i} x_{t} d v_{i t}+G w_{2} d w_{t} \tag{7}
\end{align*}
$$

where $t \in[0, \infty)$, design an $n_{e}$ th-order state estimator

$$
\begin{align*}
& d x_{e t}=A_{e} x_{e t} d t+B_{e} d y_{t},  \tag{8}\\
& y_{e t}=C_{e} x_{e t}, \tag{9}
\end{align*}
$$

which minimizes the state-estimation criterion

Reduced-Order Dynamic-Compensation Problem. Given the nth-order observed and controlled system

$$
\begin{align*}
& d x_{t}=A x_{t} d t+\sum_{i=1}^{p} A_{i} x_{t} d v_{i t}+B u_{t} d t+\sum_{i=1}^{p} B_{i} u_{t} d v_{i t}+G_{1} d w_{t}  \tag{11}\\
& d y_{t}=C x_{t} d t+\sum_{i=1}^{p} C_{i} x_{t} d v_{i t}+G_{2} d w_{t} \tag{12}
\end{align*}
$$

where $t \in[0, \infty)$, design an $n_{c} t h$-order dynamic compensator

$$
\begin{align*}
& d x_{c t}=A_{c} x_{c t} d t+B_{c} d y_{t},  \tag{13}\\
& u_{t}=C_{c} x_{c t}, \tag{14}
\end{align*}
$$

which minimizes the dynamic-compensation criterion

$$
\begin{equation*}
J_{c}\left(A_{c}, B_{c}, c_{c}\right) \triangleq \underset{t \rightarrow \infty}{\lim \sup } \underset{=}{E}\left[x_{t}^{T_{R_{1}} x_{t}}+2 x_{t}^{T_{R}} 12 u_{t}+u_{t}^{T} R_{2} u_{t}\right] . \tag{15}
\end{equation*}
$$

Clearly, $J_{m}, J_{e}$ and $J_{c}$ are nonnegative, extended-real-valued functionals defined on appropriate Ruclidean spaces. Explicit expressions for these functionals are now given. Henceforth we assume that $E\left\|\bar{x}_{r 0}\right\|^{2}<\infty$ and that $\bar{x}_{r 0}$ and $v_{1 t}, \ldots, v_{p t}, w_{t}$ are uncorrelated, $t \geq 0$.

Proposition 2.1. The nonnegative-definite covariance
$\bar{Q}_{r}(t) \triangleq E\left[\bar{x}_{r t} \bar{T}_{r t}^{T}\right], t \geq 0$, is given by

$$
\begin{equation*}
\dot{\bar{Q}}_{r}(r)=\tilde{A}_{r} \bar{Q}_{r}(t)+\bar{Q}_{r}(t) \tilde{A}_{r}^{T}+\sum_{i=1}^{p} \tilde{A}_{r i} \tilde{Q}_{r}(t) \tilde{A}_{r i}^{T}+\tilde{\nabla}_{r^{0}} t \geq 0 . \tag{16}
\end{equation*}
$$

or, explicitly, by

$$
\begin{equation*}
\tilde{Q}_{r}(t)=\operatorname{vec} \tilde{\mid \tilde{n}}_{r}, \tilde{n}_{r} \mid<\left(e^{\tilde{A}_{r}} \operatorname{vec} \tilde{Q}_{r}(0)+\int_{0}^{t} e^{\tilde{A}_{r} \sigma} d \sigma \text { vec } \tilde{\nabla}_{r}\right) . \tag{17}
\end{equation*}
$$

The cost criteria $J_{m}, J_{e}$ and $J_{c}$ are given by

$$
\begin{equation*}
J_{r}\left(A_{r}, B_{r}, C_{r}\right)=\underset{t \rightarrow \infty}{\lim \sup } t r \tilde{Q}_{r}(t) \tilde{R}_{r} \tag{18}
\end{equation*}
$$

or, equivalently, by

$$
\begin{equation*}
J_{\mathbf{r}}\left(A_{r}, B_{r}, C_{r}\right)=1 \lim _{t \rightarrow \infty}\left(\operatorname{vec} \tilde{R}_{r}\right)^{T}\left(e^{\tilde{A}_{r}} \operatorname{vec} \tilde{Q}_{r}(0)+\int_{0}^{t} e^{\tilde{A}_{r} \sigma} d \sigma \operatorname{vec} \tilde{\nabla}_{r}\right) . \tag{19}
\end{equation*}
$$

Proof. See Apperdix A.

The finiteness and smoothness of $J_{m^{\prime}} J_{e}$ and $J_{c}$ clearly depend upon the interrelationships among $\bar{Q}_{r}(0), \tilde{A}_{r}, \bar{R}_{r}$ and $\bar{\nabla}_{r}$. To avoid a derailed analysis and to guarantee that $J_{m^{0}} J_{e^{0}}$ and $J_{c}$ are finite and independent of initial date, we restrict our consideration to second-moment stable or second-moment stabilizing design triples. Furthermore, to avoid degeneracy in later developments (and without loss of generality) only minimal (i.e., controllable and observable) realizations are admitted. Hence, for the modelling, estimation and control prob1ems define the open sets

$$
\underline{S}_{r}=\left\{\left(A_{r}, B_{r}, C_{r}\right): \tilde{A}_{r} \text { is stable and }\left(A_{r}, B_{r}, C_{r}\right) \text { is minimal }\right\}
$$

In the following result we abuse notation slightly and let $\bar{Q}_{r}$ denote $\lim _{t \rightarrow \infty} \bar{Q}_{r}(t)$.
Proposition 2.2. Suppose $\underline{S}_{r}$ is nonempty. If $\left(A_{r}, B_{r}, C_{r}\right) \in S_{r}$ then $\tilde{Q}_{r} \Delta$ $\lim _{t \rightarrow \infty} \tilde{Q}_{r}(t)$ exists and is given by the unique, nonnegative-definite solution to

$$
\begin{equation*}
0=\bar{A}_{r} \tilde{Q}_{r}+\tilde{Q}_{r} \tilde{A}_{r}^{T}+\sum_{i=1}^{p} \tilde{A}_{r i} \tilde{Q}_{r} \tilde{A}_{r i}^{T}+\tilde{V}_{r} \tag{20}
\end{equation*}
$$

or, explicitly, by

$$
\begin{equation*}
\bar{Q}_{r}=\operatorname{vec}{\left.\tilde{\left(n_{r}\right.}, \tilde{n}_{r}\right)}_{\left(-\tilde{A}_{r}^{-1} \operatorname{vec} \tilde{\nabla}_{r}\right)} \tag{21}
\end{equation*}
$$

Hence,

$$
\begin{equation*}
J_{r}\left(A_{r}, B_{r}, C_{r}\right)=\operatorname{tr} \tilde{Q}_{r} \tilde{R}_{r} . \tag{22}
\end{equation*}
$$

or, equivalently,

$$
\begin{equation*}
J_{\mathbf{r}}\left(A_{\mathbf{r}} \cdot B_{\mathbf{r}}, C_{\mathbf{r}}\right)=-\left(\operatorname{vec} \tilde{R}_{\mathbf{r}}\right)^{T_{A}}{ }_{-\mathbf{r}}^{-1} \text { vec } \tilde{\nabla}_{\mathbf{r}} . \tag{23}
\end{equation*}
$$

Proof. See Appendix A.
As a side note we examine the evolution of the mean value of $\tilde{x}_{r}$.

Proposition 2.3. The mean

$$
\begin{equation*}
\tilde{m}(t) \triangleq \tilde{E x}_{r t}, t \geq 0, \tag{24}
\end{equation*}
$$

satisfies

$$
\begin{equation*}
\dot{m}(t)=\tilde{A}_{r} \tilde{m}(t), t \geq 0 . \tag{25}
\end{equation*}
$$

Furthermore, if $\left(A_{r}, B_{r}, C_{r}\right) \in \underline{S}_{\mathbf{r}}$ then $\tilde{A}_{\mathbf{r}}$ is stable and thus

$$
\begin{equation*}
\lim _{t \rightarrow \infty} \tilde{m}(t)=0 . \tag{26}
\end{equation*}
$$

Proof. See Appendix A.

Of course; it is useful to know when the sets $S_{-m}, \underline{S}_{e}$ and $\underline{S}_{c}$ are nonempty. Although for the closed-loop control problem the question is complex because of stabilizability concerns, the modelling and estimation problems permit considerable simplification.

Proposition 2.4. $\underline{S}_{m}$ (alternatively, $\underline{S}_{e}$ ) is nonempty if and only if A is stable. In this case $S_{m}$ and $\underline{S}_{e}$ are given by
$\underline{S}_{m}=\left\{\left(A_{m}, B_{m}, C_{m}\right): A_{m}\right.$ is stable and $\left(A_{m}, B_{m}, C_{m}\right)$ is minimal $\}$,
$S_{e}=\left\{\left(A_{e}, B_{e}, C_{e}\right): A_{e}\right.$ is stable and $\left(A_{e}, B_{e}, C_{e}\right)$ is minimal $\}$.

## Proof. See Appendix B.

The following observation concerns the smoothness of the cost functionals.

Proposition 2.5. The functionals $J_{r}$ are infinitely Frechet differentiable on $\underline{S}_{r}$.

Proof. From Lemma 3.7.2. p. 203 of Ref. 38, it follows that the map $W-W^{-1}$ defined on the set of invertible matrices is $C^{\infty}$. The result follows from the chain rule and (23).

It is now possible to proceed with the principal aim of the paper which is to characterize solutions of the reduced-order modelling, estimation and control problems by means of a first-order variational analysis. To this end, one additional assumption is required. In order to obtain closed-form expressions for extremal values of the closed-loop control gains, the dynamiccompensation problem requires the technical assumption

$$
\begin{equation*}
\left[B_{i} \neq 0 \Rightarrow c_{i}=0\right], i=1, \ldots, p \tag{27}
\end{equation*}
$$

or, equivalently,

$$
\begin{equation*}
\left[C_{i} \notin 0 \Longrightarrow B_{i}=0\right], i=1, \ldots, p \tag{28}
\end{equation*}
$$

i.e., for each $i \in\{1, \ldots, p\}, B_{i}$ and $C_{i}$ are not both nonzero. Essentially, (27). expresses the condition that the control-dependent and measurementdependent noises are independent. There are no constraints, however, on correlations with the state-dependent noise.

In order to state the main results we require sone additional notation and a lemma concerning a pair of nonnegative-definite matrices. For a realsemisimple matrix $X \in \mathbb{R}^{\mathbb{n} \times \mathbf{n}}$ define the set of diagonalizing matrices

$$
\underline{D}(\mathbb{X}) \hat{\sum}\left\{\Psi \in \underline{R}^{\mathrm{n} \times \mathrm{n}}: \Psi^{-1} \mathbb{X} \Psi \text { is diagonal }\right\}
$$

and, for a pair of nonnegative-definite matrices $X, Y \in \mathbb{R}^{\mathbb{R} \times \mathbb{n}}$ define the set of contragrediently diagonalizing matrices

$$
\underline{C}(X, Y) \triangleq\left\{\Psi \in \underline{E}^{n \times n}: \Psi^{-1} X \Psi^{-T} \text { and } \Psi^{T} Y \Psi \text { are diagonal }\right\}
$$

and the subset of balancing transformations

$$
\underline{B}(X, Y) \triangleq\left\{\Psi \in \underline{C}(X, Y): \Psi^{-1} X \Psi^{-T}=\Psi^{T} Y \Psi\right\} .
$$

The following result unifies and extends similar results found in Refs. 32, 33 and 12.

Lemme 2.1. Suppose $\hat{Q}, \hat{P} \in \mathbb{R}^{\mathrm{R} \times \mathbb{R}}$ are nonnegative definite and
$\rho(\hat{Q} \hat{P})=n_{r}$. Then the following statements hold:
(i) $\emptyset \neq \underline{\mathrm{C}}(\hat{Q}, \hat{\mathrm{P}}) \subset \underline{\mathrm{D}}(\hat{\mathrm{Q}} \hat{P})$.
(ii) $\widehat{Q P}$ is nonnegative semisimple.
(iii) The $n \times n$ matrix

$$
\begin{equation*}
\tau \triangleq \hat{Q} \hat{P}(\hat{Q} \hat{P})^{*} \tag{29}
\end{equation*}
$$

is idempotent, ie., $T$ is an oblique projection.
(iv) There exists $\Psi \in \underline{C}(\hat{Q}, \hat{P})$ with $\left(\Psi^{-1} \hat{Q} \hat{P} \Psi\right)_{(i, i)} \neq 0, i=1, \ldots, n_{r}$, such that $T$ is given by

$$
\begin{equation*}
r=\sum_{i=1}^{n_{r}} n_{i}(\Psi) . \tag{30}
\end{equation*}
$$

(v) If $\rho(\hat{Q})=\rho(\hat{P})=n_{r}$ then $\underline{B}(\hat{Q}, \hat{P}) \neq \emptyset$.

$$
\begin{align*}
& \text { (vi) If } \rho(\hat{Q})=\rho(\hat{p})=n_{I} \text { then there exists } \Psi \in \underline{B}(\hat{Q}, \hat{P}) \text { with } \\
& \left(\Psi^{-1} \hat{Q} \hat{P}^{( } \Psi\right)_{(i, i)} \neq 0, i=1, \ldots, n_{r}, \text { such that } T \text { is given by (30). } \\
& \text { (vii) If } \rho(\hat{Q})=\rho(\hat{P})=n_{r} \text { then } \\
& \hat{Q}=\tau \hat{Q}=\hat{Q} \tau^{T}=\tau \hat{Q} T^{T},  \tag{31}\\
& \hat{\mathbf{P}}=\tau^{T} \boldsymbol{P}=\hat{\mathrm{P}} \boldsymbol{\tau}=\tau^{\boldsymbol{T}}{ }^{\hat{P}} \boldsymbol{\tau} . \tag{32}
\end{align*}
$$



$$
\begin{align*}
& \hat{Q P}=G^{T} M,  \tag{33}\\
& \Gamma G^{T}=I_{\mathbf{m}_{r}} \tag{34}
\end{align*}
$$

(ix) If $\bar{G}, \bar{J} \in R^{n} r^{\times n}$ and $\bar{M} \in R^{n_{i} \times n_{r}}$ satisfy

$$
\begin{align*}
& \hat{Q P}=\bar{G}^{T} \bar{M} \bar{I},  \tag{35}\\
& \bar{\Gamma}^{T}=I_{\mathbf{n}_{\mathbf{I}}}, \tag{36}
\end{align*}
$$


 invertible, $(\hat{Q P})^{\frac{1}{f}}=G^{T} M^{-1} T$ and

$$
\begin{equation*}
T=G^{T} \Gamma \tag{37}
\end{equation*}
$$

Proof. See Appendix C.

For convenience we shall call $G, \Gamma \in \mathbb{R}^{n^{r} \times n}$ and $M \in \mathbb{R}^{n^{n} \times n_{r}}$ satisfying (33) and (34) a projective factorization of $\hat{Q P}$. Furthermore, define the complementary oblique projection

$$
\begin{equation*}
\tau_{\perp} \triangleq I_{n}-\tau \tag{38}
\end{equation*}
$$

and let $J_{r}\left(A_{r}, B_{r}, C_{r}\right)$ denote the Frechet derivative of $J_{r}$ evaluated at ( $A_{r}, B_{r}, C_{r}$ ).

It is now possible to atate the main results, which provide a parameterization of triples $\left(A_{r}, B_{r}, C_{r}\right) \in \underline{S}_{\mathbf{r}}$ for which the first Frechet derivative of $\mathrm{J}_{\mathbf{r}}$ vanishes.

Theorem 2.1. Assume $A$ is stable. Then, for $\left(A_{m}, B_{m}, C_{m}\right) \in S_{m}$.

$$
\begin{equation*}
J_{m}^{-}\left(A_{m}, B_{m}, C_{m}\right)=0 \tag{39}
\end{equation*}
$$

if and only if there exist $n \times n$ nonnegative-definite matrices $\hat{Q}$ and $\hat{P}$ such that, for some projective factorization $G, M, \Gamma$ of $\hat{Q} \hat{P}, A_{m}, B_{m}$ and $C_{m}$ are given by

$$
\begin{align*}
& A_{m}=\Gamma A G^{T}  \tag{40}\\
& B_{m}=\Gamma B  \tag{41}\\
& C_{m}=C G^{T} \tag{42}
\end{align*}
$$

and such that, with $\tau \hat{\sum} \hat{Q P}(\hat{Q} \hat{Q})^{*}=G^{T} \Gamma$ and $\tau_{\perp} \hat{\wedge} I_{n}-\tau$, the following conditions are satisfied:

$$
\begin{align*}
& 0=A \hat{Q}+\hat{Q} A^{T}+B V B^{T}-T_{\perp} B V B^{T} T_{\perp}^{T}  \tag{43}\\
& 0=A^{T} \hat{P}+\hat{P} A+C^{T} R C-T_{1}^{T} C^{T} R C T_{\perp}  \tag{44}\\
& \rho(\hat{Q})=\rho(\hat{P})=\rho(\hat{Q P})=n_{m} . \tag{45}
\end{align*}
$$

Furthermore, if ( $A_{m}, B_{m}, C_{m}$ ) $\underline{S}_{m}$ satisfies (39) then the extremal cost is given by

$$
\begin{align*}
J_{m}\left(A_{m}, B_{m}, C_{m}\right) & =\operatorname{tr}\left[\left(W_{c}-\hat{Q}\right) C^{T} R C\right] \\
& =\operatorname{tr}\left[\left(W_{0}-\hat{P}\right) B \forall B^{T}\right]  \tag{46}\\
& =2 \operatorname{tr}\left[\left(\hat{Q} \hat{P}-W_{c} W_{0}\right) A\right]-2 \sum_{i=1}^{p} t r W_{c} A_{i}^{T} W_{0} A_{i} .
\end{align*}
$$

where $W_{c}, W_{0} \in \mathbb{R}^{n \times n}$ are the unique, nonnegative-definite solutions to

$$
\begin{align*}
& 0=A W_{c}+W_{c} A^{T}+\sum_{i=1}^{p} A_{i} W_{c} A_{i}^{T}+B V B^{T}  \tag{47}\\
& 0=A^{T} W_{0}+W_{0} A+\sum_{i=1}^{p} A_{i}^{T} W_{0} A_{i}+c^{T} R c . \tag{48}
\end{align*}
$$

Theorem 2.2. Assume $A$ is stable. Then, for ( $A_{e}{ }^{B_{e}} e^{C_{e}}$ ) $\in \underline{S}_{e}{ }^{\text {- }}$

$$
\begin{equation*}
J_{e}^{-}\left(\mathbb{A}_{e}, B_{e}, C_{e}\right)=0 \tag{49}
\end{equation*}
$$

if andonly if there exist $n \times n$ nonnegative-definite matrices $Q, \hat{Q}$ and $\hat{P}$ such that, for some projective factorization $G, M, \Gamma$ of $\hat{Q P}, A_{e}, B_{e}$ and $C_{e}$ are given by

$$
\begin{align*}
& A_{e}=\Gamma\left(A-Q \hat{V}_{2}^{-1} C\right) G^{T}  \tag{50}\\
& B_{e}=\Gamma \hat{V}_{2}^{-1}  \tag{51}\\
& C_{e}=L G^{T} \tag{52}
\end{align*}
$$

and such that, with $\tau \triangleq \hat{\underline{Q}} \hat{P}(\hat{Q} \hat{P})^{\neq}=G^{T} \Gamma$ and $\tau_{\perp} \triangleq I_{n}-\tau$ the following conditions are satisfied:

$$
\begin{align*}
& 0=A Q+Q A^{T}+V_{1}+\sum_{i=1}^{P} A_{i}(Q+\hat{Q}) A_{i}^{T}-Q \hat{V}_{2}^{-1} Q^{T}+\tau_{1} \hat{V}_{2}^{-1} Q^{T} \tau_{1}^{T}  \tag{53}\\
& 0=A \hat{Q}+\hat{Q A}^{T}+Q \hat{V}_{2}^{-1} Q^{T}-\tau_{1} \hat{Q V}_{2}^{-1} Q^{T} \tau_{1}^{T}  \tag{54}\\
& 0=A_{Q}^{T} \hat{P}+\hat{P} A_{Q}+L{ }^{T} N L-\tau_{1}^{T} T_{N L} \tau_{1} .  \tag{55}\\
& \rho(\hat{Q})=\rho(\hat{Q})=\rho(\hat{Q} \hat{P})=n_{e} . \tag{56}
\end{align*}
$$

Furthermore, if ( $A_{e}, B_{e}, C_{e}$ ) $\in \underline{S}_{e}$ satisfies (49) then the extremal cost is given by

$$
\begin{equation*}
J_{e}\left(\mathbb{A}_{e}, B_{e}, C_{e}\right)=\operatorname{tr} Q^{T}{ }^{T} N L . \tag{57}
\end{equation*}
$$

Theorem 2.3. Assume (27) holds and $\mathrm{S}_{\mathrm{c}}$ is nonempty. Then for $\left(A_{c}, B_{c}, C\right) \in S_{c}$.

$$
\begin{equation*}
J_{c}^{-}\left(A_{c}, B_{c}, C_{c}\right)=0 \tag{58}
\end{equation*}
$$

if and only if there exist $n \times n$ nonnegative-definite matrices $Q, P, \hat{Q}$ and $\hat{P}$ such that, for some projective factorization $G, M, \Gamma$ of $\widehat{Q P}, A_{c}, B_{c}$ and $C_{c}$ are given by

$$
\begin{align*}
& A_{c}=\Gamma\left(A-B R_{2}^{-1} \underline{P}-Q \hat{V}_{2}^{-1} c\right) G^{T}  \tag{59}\\
& B_{c}=\Gamma \hat{Q}_{2}^{-1}  \tag{60}\\
& C_{c}=-\hat{R}_{2}^{-1} \underline{P G}^{T} \tag{61}
\end{align*}
$$

and such that, with $r \triangleq \hat{Q P} \hat{P}(\hat{Q P})^{\#}=G^{T} \Gamma$ and $T_{\perp} \triangleq I_{n}-\tau_{0}$ the following conditions ase satisfied:

$$
\begin{align*}
& 0=A Q+Q A^{T}+V_{1}+\sum_{i=1}^{P}\left[A_{i} Q A_{i}^{T}+\left(A_{i}-B_{i} \hat{R}_{2}^{-1} \underline{P}\right) \hat{Q}\left(A_{i}-B_{i} \hat{R}_{2}^{-1} \underline{Q}\right)^{T}\right]-Q \hat{V}_{2}^{-1} 9^{T}+\tau_{1} \hat{V}_{2}^{-1} g^{T} T V_{1} .  \tag{62}\\
& 0=A^{T} P+P A+R_{1}+\sum_{i=1}^{P}\left[A_{i}^{T} P A_{i}+\left(A_{i}-\hat{Q}_{2}^{-1} C_{i}\right)^{T} \hat{P}\left(A_{i}-\hat{Q}_{2}^{-1} C_{i}\right)\right]-\underline{P}^{T} R_{2}^{-1} \underline{P}+\tau_{1}^{T} \underline{P}^{T} R_{2}^{-1} \underline{P} \tau_{1} .  \tag{63}\\
& 0=A_{P} \hat{Q}+\hat{Q A}_{P}^{T}+Q \hat{V}_{2}^{-1} Q^{T}-\tau_{1} Q \hat{V}_{2}^{-1} Q^{T} \tau_{1}^{T} .  \tag{64}\\
& 0=A_{Q}^{T} \hat{P}+\hat{P} A_{Q}+\underline{P}^{T} \hat{R}_{2}^{-1} \underline{P}-\tau_{1}^{T} \underline{T}^{T} \mathbb{R}_{2}^{-1} \underline{P} \tau_{1} .  \tag{65}\\
& \rho(\hat{Q})=\rho(\hat{P})=\rho(\hat{Q} \hat{P})=n_{c} . \tag{66}
\end{align*}
$$

Furthermore, if $\left(A_{c}, B_{c}, C_{c}\right) \in \underline{S}_{c}$ satisfies (58) then the extremal cost is given by

$$
\begin{align*}
& \hat{J}_{c}\left(A_{c}, B_{c}, C_{c}\right)=\operatorname{tr}\left[(Q+\hat{Q}) R_{1}-2 R_{12} \hat{R}_{2}^{-1} \underline{P Q}\right. \\
&\left.+\underline{P}^{T} R_{2}^{-1} R_{2} \hat{R}_{2}^{-1} \underline{Q} \hat{Q}\right] \tag{67}
\end{align*}
$$

## 3. Appendix A: Proof of Propositions 2.1, 2.2 and 2.3

First note that (1)-(4). (6)-(9) and (11)-(14) can be written as

$$
\begin{equation*}
\tilde{d x}_{r t}=\tilde{A}_{r} \tilde{x}_{r t} d t+\sum_{i=1}^{p} \tilde{A}_{r i} \tilde{X}_{r t} d v_{i t}+\tilde{G}_{r} d w_{t} \tag{68}
\end{equation*}
$$

From Theorem 8.5.5. p. 142 of Ref. 17 (or from the Ito differential rule) it follows that the nonnegative-definite covariance $\tilde{Q}_{r}(t)$ is given by (16). Eurthermore, (5). (10) and (15) are equivalent to (18). Rewriting $\bar{Q}_{r}(t)$ in the form (see Refs. 35 and 36)
$\operatorname{vec} \dot{\tilde{Q}}_{\mathbf{r}}(t)=\tilde{\underline{A}}_{\mathbf{r}} \operatorname{vec} \tilde{Q}_{\mathbf{Q}}(t)+\operatorname{vec} \tilde{\nabla}_{\mathbf{r}}$
leads to (19).

To prove Proposition 2.2, note that the stability of $\tilde{\sim}_{\mathbf{r}}$ implies, by (69). that $\tilde{Q}_{r} \triangleq \lim _{t \rightarrow \infty} \tilde{Q}_{r}(t)$ exists and is given by (21). which satisfies (20). Clearly, $\tilde{Q}_{I} \geq 0$ since $\bar{Q}_{r}(t) \geq 0, t \geq 0$. Now (22) and (23) follow from (18) and (19).

To prove Proposition 2.3, note that the differential equation for $\bar{m}(t)$ is an immediate consequence of (68). To show that $\tilde{A}_{r}$ is stable, we proceed as in Lemma 2.2 of Ref. 4. Repeating the steps leading to (22) with $\tilde{\nabla}_{r}$ replaced by $I_{n_{r}}^{\sim}$. it follows (see Ref. 4) that $\left(\bar{A}_{\mathbf{r}^{\prime}}, I_{\mathbf{n}_{\mathbf{r}}}^{\sim}\right)$ is stabilizable. Hence,

$$
\int_{0}^{t} e^{\tilde{A}_{r} \sigma} e^{\tilde{A}_{r}^{T} \sigma} d \sigma
$$

is bounded as $t \rightarrow \infty$, which implies $\bar{A}_{r}$ is stable.

## 4. Appendix B: Proof of Proposition 2.4

We require some elementary properties of the Kronecker product (Refs. 35 and 36) applied to partitioned matrices. For $X \in \mathbb{R}^{8 \times t}$ and $Y \in \underline{R}^{p \times q}$ partitioned by

$$
X=s_{1}\left[\begin{array}{ll}
s_{1} & t_{2} \\
X_{1} & X_{12} \\
X_{21} & X_{2}
\end{array}\right], \quad Y=\begin{aligned}
& p_{1}
\end{aligned}\left[\begin{array}{ll}
q_{1} & q_{2} \\
Y_{1} & Y_{12} \\
Y_{21} & Y_{2}
\end{array}\right]
$$

it follows that

$$
\begin{aligned}
& X \otimes Y=\left[\begin{array}{ll}
X_{1} \otimes Y & X_{12} \otimes Y \\
X_{21} \otimes Y & X_{2} \otimes Y
\end{array}\right] \\
& =\left[\begin{array}{ll}
U_{\mathbf{B}_{1} \times p} & 0 \\
0 & U_{\mathbf{B}_{2} \times p}
\end{array}\right]\left[\begin{array}{ll}
Y \otimes X_{1} & Y \otimes X_{12} \\
Y \otimes X_{21} & Y \otimes X_{2}
\end{array}\right]\left[\begin{array}{ll}
U_{q \times t_{1}} & 0 \\
0 & U_{q \times t_{2}}
\end{array}\right] \\
& =\left[\begin{array}{ll}
U_{B_{1} \times p} & 0 \\
0 & U_{s_{2 \times p}}
\end{array}\right]\left[\begin{array}{llll}
Y_{1} \otimes X_{1} & Y_{12} \otimes X_{1} & Y_{1} \otimes X_{12} & Y_{12} \otimes X_{12} \\
Y_{21} \otimes X_{1} & Y_{2} \otimes X_{1} & Y_{21} \otimes X_{12} & Y_{2} \otimes X_{12} \\
Y_{1} \otimes X_{21} & Y_{12} \otimes X_{21} & Y_{1} \otimes X_{2} & Y_{12} \otimes X_{2} \\
Y_{21} \otimes X_{21} & Y_{2} \otimes X_{21} & Y_{21} \otimes X_{2} & Y_{2} \otimes X_{2}
\end{array}\right]\left[\begin{array}{ll}
U_{q \times t_{1}} & 0 \\
0 & U_{q \times t_{2}}
\end{array}\right],
\end{aligned}
$$

Where $U_{i x j}$ is the permutation matrix defined in Refs. 35 and 36 . Since $U_{i x j}=$ $\mathrm{V}_{\mathrm{ixj}}^{-1}$ (ide.. $\mathrm{V}_{\mathrm{ixj}}$ is involutory), the stability of (square) $X \otimes Y$ is equivalent to the stability of the above block $4 \times 4$ matrix. Hence note that

$$
\overline{A_{m}}=\text { block-diagonal }\left(A_{2}, A \oplus A_{m}, A_{m} \oplus A_{p} A_{m} \oplus A_{m}\right) .
$$

If ( $A_{m}, B_{m}, C_{m}$ ) is such that $\tilde{A}_{m}$ is stable then clearly $A$ is also stable and, by elementary properties of the Kronecker sum, $A_{m}$ is stable. Conversely, if A and $A_{m}$ are stable then $\tilde{A}_{m}$ is stable. The result for $\tilde{A}_{e}$ is obtained analogously noting only that ${\underset{-}{e}}^{\text {is }}$ lower block triangular.

## 5. Appendix C: Proof of Lemma 2.1

(i) From Theorem 6.2.5, p. 123 of Ref. 39. it follows that there exists an n xn invertible matrix $\Psi$ such that the nonnegative-definite matrices $D_{\hat{Q}}^{\hat{2}} \Psi^{-1} \hat{Q} \Psi^{-T}$ and $D_{\hat{P}}^{\hat{\omega}} \Psi^{T} \hat{P} \Psi$ are both diagonal. Hence $\bar{C}(\hat{Q}, \hat{P}) \neq \phi$. Since $\mathrm{D}_{\hat{Q}}^{\sim} \mathrm{D}_{\mathrm{P}}=\Psi^{-1} \hat{Q P} \Psi$ is also diagonal, $\underline{C}(\hat{Q}, \hat{P}) \subset \underline{D}(\hat{Q P})$.
(ii) Since $\hat{Q P}=\Psi \Lambda \Psi^{-1}$, where $\Lambda \hat{\otimes} D_{\hat{Q}}^{\hat{Q}} \hat{P}$ is nonnegative diagonal, $\hat{Q} \hat{P}$ is nonnegative semisimple.
(iii) Since $\hat{Q P}$ is semisimple, it is group invertible. By properties of the group inverse (Ref. 37. P. 124). $T^{2}=T$.
(iv) Note that by means of a basis rearrangement, it can be assumed that $\Psi$ in (i) is such that $\Lambda=\operatorname{diag}\left(\lambda_{1}, \ldots, \lambda_{n_{r}}, 0, \ldots, 0\right)$, where


$$
\tau=\hat{Q} P(\hat{Q P})^{*}=\Psi \Lambda \Lambda^{*} \Psi^{-1}=\sum_{i=1}^{n} \Psi E_{i} \Psi^{-1}
$$

(v) Since $\rho(\hat{Q})=\rho(\hat{P})=\rho(\hat{Q P})$, it follows that $\left(D_{\hat{Q}}^{\hat{Q}}(i, i) \neq 0\right.$, ${\left(D_{P}^{e}\right)}_{(i, i)} \neq 0, i=1, \ldots, n_{r}$. Hence, let $\Lambda_{\hat{Q}}^{\hat{a}}$ and $\Lambda_{\hat{P}}^{\hat{p}}$ denote the upper left positivediagonal blocks of $D_{\hat{Q}}$ and $D_{\hat{P}}$, respectively, and define

$$
\hat{\psi} \triangleq \psi\left[\begin{array}{cc}
\left(\Lambda_{\hat{Q}} \Lambda_{P}^{z 1}\right)^{\frac{1}{4}} & 0 \\
0 & I_{n-n_{r}}
\end{array}\right]
$$

It now follows that

$$
\hat{\psi}^{-1} \hat{Q}^{-T}=\hat{\psi}^{T} \hat{P}^{-1}=\left[\begin{array}{ll}
\left(\Lambda_{\hat{Q}}^{\hat{Q}} \hat{P}^{-1}\right)^{\frac{1}{2}} & 0  \tag{70}\\
0 & 0
\end{array}\right]
$$

as desired.
(vi) This is an immediate consequence of (70).
(vii) This is an immediate consequence of (70) and (30)
(viii) With $\psi$ as in (iv) and $\Lambda_{0}{ }_{0}$ diag $\left(\tau_{1}, \ldots, \tau_{n_{r}}\right)$, it follows that for arbitrary invertible $S \in \mathbb{R}^{\mathrm{n} \mathrm{\times n}}$,

$$
\hat{Q P}=\psi\left[\begin{array}{l}
s \\
0
\end{array}\right]\left(s^{-1} \Lambda_{0} s\right)\left[s^{-1} \quad 01 \psi^{-1}\right.
$$

and thus (33) and (34) hold with $G=\left[S^{T} 0\right], M=S^{-1} \Lambda_{0} S$ and $\Gamma=\left[S^{-1} 0\right] \psi^{-1}$.
(ix) The result follows from $S=\bar{M}^{-1} \bar{\Gamma} G^{T} M^{-1}$ with $S^{-1}=M \bar{\Gamma} \bar{G}^{T} \bar{M}^{-1}$.
( $x$ ) The result is a consequence of (viii) and (ix).
6. Appendix D: Proof of Theorem 2.3.

First note that by arguments similar to those used in Appendix A the dual of (22) given by

$$
\begin{equation*}
0=\tilde{A}_{r} P_{r}+\tilde{P}_{r} \tilde{A}_{r}+\sum_{i=1}^{P} \tilde{A}_{r i}^{T} \tilde{P}_{r} \tilde{A}_{r i}+\tilde{R}_{r} \tag{71}
\end{equation*}
$$

has a unique, nonnegative-definite solution given explicitly by

$$
\begin{equation*}
\left.\tilde{P}_{r}=\operatorname{vec}_{\left(\tilde{n}_{r}, \tilde{n}_{r}\right.}^{-1}\right)\left(-\tilde{A}_{r}^{-T} \operatorname{vec} \tilde{R}_{r}\right) \tag{72}
\end{equation*}
$$

Define the partitionings

$$
\left.\dot{Q}_{r}=\begin{array}{c}
n \\
n_{r}
\end{array} \begin{array}{cc}
n_{r} \\
Q_{1} & Q_{12} \\
Q_{12}^{T} & Q_{2}
\end{array}\right], \quad \tilde{P}_{r}=\begin{gathered}
n_{r} \\
n_{r}\left[\begin{array}{cc}
P_{1} & P_{12} \\
P_{12}^{T} & P_{i}
\end{array}\right], ~ . ~
\end{gathered}
$$

where, for notational convenience, we suppress the subscript $r$. Also define the notation

$$
\tilde{P Q}=\left[\begin{array}{ll}
z_{1} & z_{12} \\
z_{21} & z_{2}
\end{array}\right]
$$

where

$$
\begin{array}{ll}
Z_{1} \triangleq P_{1} Q_{1}+P_{12} Q_{12}^{T}, & Z_{12} \triangleq P_{1} Q_{12}+P_{12} Q_{2}, \\
Z_{21} \triangleq P_{12}^{T} Q_{1}+P_{2} Q_{12}^{T}, & z_{2} \triangleq P_{12}^{T} Q_{12}+P_{2} Q_{2},
\end{array}
$$

 control problem.

Lemma 6.1. Under the as sumptions of Theorem 2.3,

$$
\begin{align*}
& J_{c}^{-}\left(A_{c}, B_{c}, C_{c}\right)\left(\delta_{A_{c}}, \delta_{B_{c}}, \delta_{C}\right)=2 \operatorname{tr}\left[Z_{2}^{T} \delta_{A_{c}}\right]  \tag{73}\\
& \quad+2 \operatorname{tr}\left[\left(\hat{\nabla}_{2} B_{c}^{T} P_{2}+C z_{21}^{T}+\left[v_{12}^{T}+\sum_{i=1}^{P} C_{i} Q_{1} A_{i}^{T}\right] P_{12}\right) \delta_{B_{c}}\right] \\
& \quad+2 \operatorname{tr}\left[\left(Q_{2} C_{c}^{T} R_{2}+z_{12}^{T} B+Q_{12}^{T}\left[R_{12}+\sum_{i=1}^{P} A_{i}^{T} P_{1} B_{i}\right]\right) \delta_{C_{c}}\right]
\end{align*}
$$

Proof. From Lemma 3.7.2, p. 203 of Ref. 38, it follows that the Frechet derivative of the map $W \rightarrow W^{-1}$ is given by $\delta_{W} \rightarrow-W^{-1} \delta_{W} W^{-1}$. Also, recall from Refs. 35 and 36 the identities

$$
(\operatorname{vec} X)^{T} \operatorname{vec} Y=\operatorname{tr} X^{T} Y
$$

$(X \otimes Y)$ ven $Z=\operatorname{vec} Y Z X^{T}$.
Hence, using (23) and noting that $\tilde{\nabla}_{c}$ and $\tilde{R}_{c}$ are independent of $A_{c}$, we compute

$$
\begin{aligned}
&\left.\frac{\partial J_{c}\left(A_{c}, B_{c}, C\right.}{}\right) \\
& \partial A_{c} \delta_{A_{c}}
\end{aligned}=\left(\operatorname{vec} \tilde{R}_{c}\right)^{T} \underline{A}_{c}^{-1}\left(\frac{\partial \tilde{A}_{c}}{\partial A_{c}} \delta_{A_{c}}\right) \tilde{A}_{c}^{-1} \operatorname{vec} \tilde{\nabla}_{c} .
$$

Purthermore, noting that.

$$
\operatorname{vec} \tilde{\nabla}_{c}=\operatorname{vec} \tilde{B}_{c} \tilde{V B}_{c}^{T}=\left(\tilde{B}_{c} \otimes \tilde{B}_{c}\right) \operatorname{vec} \tilde{\nabla}
$$

we obtain

$$
\begin{aligned}
& \frac{\partial J_{c}\left(A_{c}, B_{c}, C_{c}\right)}{\partial B_{c}} \delta_{B}=\left(\operatorname{vec} \tilde{P}_{c}\right)^{T}\left(\frac{\partial \tilde{\tilde{A}}_{c}}{\partial \tilde{B}_{c}} \delta_{B_{c}}\right) \text { vec } \tilde{Q}_{c} \\
& -\left(\text { vec } \tilde{P}_{c}\right)^{T}\left(\frac{\partial\left(\tilde{B}_{c} \otimes \tilde{B}_{c}\right)}{\partial B_{c}} \delta_{B_{c}}\right) \text { vec } \tilde{\nabla} \\
& =\left(\operatorname{vec} \tilde{P}_{C}\right)^{T}\left(\left[\begin{array}{ll}
0 & 0 \\
\delta_{B_{C}}{ }^{C} & 0
\end{array}\right] \oplus\left[\begin{array}{ll}
0 & 0 \\
\delta_{B_{C}} C & 0
\end{array}\right]\right. \\
& \left.+\sum_{i=1}^{p} \tilde{A}_{c i} \otimes\left[\begin{array}{lll}
0 & & 0 \\
\delta_{B_{c}} c_{i} & 0
\end{array}\right]+\left[\begin{array}{lll}
0 & & 0 \\
\delta_{B} c_{i} & 0
\end{array}\right] \otimes \tilde{A}_{c i}\right) \operatorname{vec} \tilde{Q}_{c} \\
& -\left(\operatorname{vec} \tilde{P}_{c}\right)^{T}\left(\tilde{B}_{c} \otimes\left[\begin{array}{ll}
0 & 0 \\
0 & \delta_{B_{c}}
\end{array}\right]+\left[\begin{array}{ll}
0 & 0 \\
0 & \delta_{B_{c}}
\end{array}\right] \otimes \tilde{B}_{c}\right) \text { vec } \dot{\nabla} \\
& =2 t r \tilde{Q}_{c} \tilde{P}_{c}\left[\begin{array}{ll}
0 & 0 \\
\delta_{B_{c}} c & 0
\end{array}\right]+2 t r \sum_{i=1}^{p} \tilde{Q}_{c} \tilde{A}_{c i}^{T} \tilde{P}_{c}\left[\begin{array}{ll}
0 & 0 \\
\delta_{B} c_{i} & 0
\end{array}\right] \\
& +2 \operatorname{tr}{\underset{\nabla B}{c}}_{\sim}^{T_{c}} \tilde{P}_{c}\left[\begin{array}{ll}
0 & 0 \\
0 & \delta_{B}
\end{array}\right] \\
& =2 \operatorname{tr} C z_{21}^{T} \delta_{B_{c}}+2 \operatorname{tr} \sum_{i=1}^{P}\left(C_{i} Q_{1} A_{i}^{T} P_{12}+C_{i} Q_{1} C_{i}^{T}\right) \\
& +2 \operatorname{tr}\left(\nabla_{12}^{T} P_{12}+\nabla_{2} B_{c}^{T} P_{2}\right) \delta_{B_{c}} \\
& =2 \operatorname{tr}\left(\hat{\nabla}_{2} \mathrm{~B}_{\mathrm{C}} \mathrm{P}_{2}+\mathrm{Cz}_{21}^{\mathrm{T}}+\left[\nabla_{12}^{\mathrm{T}}+\sum_{i=1}^{\mathrm{P}} \mathrm{C}_{\mathrm{i}} \mathrm{Q}_{1} \mathrm{~A}_{\mathrm{i}}^{\mathrm{T}}\right] \mathrm{P}_{12}\right) \delta_{\mathrm{B}_{\mathrm{c}}} .
\end{aligned}
$$

A similar computation for ( $\left.\partial J_{c}\left(A_{c}, B_{c}, C_{c}\right) / \partial C_{c}\right) \delta_{C_{c}}$ yields (73).

We can now proceed with the proof of Theorem 2.3. Obviously, (58) is equivalent to

$$
\begin{align*}
& 0=Z_{2}  \tag{74}\\
& 0=P_{2} B c \hat{\nabla}_{2}+Z_{21} C^{T}+P_{12}^{T}\left(V_{12}+\sum_{i=1}^{P} A_{i} Q_{1} C_{i}^{T}\right)  \tag{75}\\
& 0=\hat{R}_{2} C_{c} Q_{2}+B^{T} z_{12}+\left(R_{12}^{T}+\sum_{i=1}^{P} B_{i}^{T} P_{1} A_{i}\right) Q_{12} \tag{76}
\end{align*}
$$

Expanding the $n \times n, n \times n_{c}$ and $n_{c} \times n_{c}$ blocks of (20) and (71) yields

$$
\begin{align*}
& 0=A Q_{1}+Q_{1} A^{T}+\nabla_{1}+B C_{c} Q_{12}^{T}+Q_{12}\left(B C_{c}\right)^{T}  \tag{77}\\
& +\sum_{i=1}^{p}\left[A_{i} Q_{1} A_{i}^{T}+B_{i} C_{c} Q_{12}^{T} A_{i}^{T}+A_{i} Q_{12}\left(B_{i} C_{c}\right)^{T}+B_{i} C_{c} Q_{2}\left(B_{i} C_{c}\right)^{T}\right] . \\
& 0=A Q_{12}+Q_{12} A_{c}^{T}+B C_{c} Q_{2}+Q_{1}\left(B_{c} C\right)^{T}  \tag{78}\\
& +\sum_{i=1}^{p} A_{i} Q_{1}\left(B_{c} C_{i}\right)^{T}+\nabla_{12} B_{c}^{T} . \\
& 0=A_{c} Q_{2}+Q_{2} A_{c}^{T}+B_{c} C_{12}+Q_{12}^{T}\left(B_{c} C\right)^{T}+B_{c} \hat{\nabla}_{2} B_{c}{ }^{T} \text {. }  \tag{79}\\
& 0=A^{T} P_{1}+P_{1} A+R_{1}+\left(B_{c} C\right)^{T} P_{12}^{T}+P_{12} B_{c} C  \tag{80}\\
& +\sum_{i=1}^{p}\left[A_{i}^{T} P_{1} A_{i}+\left(B_{c} C_{i}\right)^{T} P_{12}^{T} A_{i}+A_{i}^{T} P_{12} B_{c} C_{i}+\left(B_{c} C_{i}\right)^{T} P_{2} B_{c} C_{i} .\right.
\end{align*}
$$

$$
\begin{align*}
& 0=A^{T} P_{12}+P_{12} A_{c}+\left(B_{c} C\right)^{T} P_{2}+P_{1} B C_{c} \\
& +\sum_{i=1}^{p} A_{i}^{T} P_{1} B_{i} C_{c}+R_{12} C_{c} .  \tag{81}\\
& 0=A_{c}^{T} P_{2}+P_{2} A_{c}+\left(B C_{c}\right)^{T} P_{12}+P_{12}^{T} B C_{c}+C_{c}^{T} \hat{R}_{2} C_{c} . \tag{82}
\end{align*}
$$

Obviously, $\nabla_{2}>0$ and $R_{2}>0$ imply $\vec{\nabla}_{2}>0$ and $\hat{R}_{2}>0$. Next note that since ( $A_{c}, B_{c}$ ) is controllable and $V_{2}>0$ it follows that ( $A_{c}+B_{c} C Q_{12} Q_{2}^{\#}, B_{c}{ }_{c}^{\frac{1}{2}}$ ) is controllable. Using $Q_{12}=Q_{12} Q_{2}^{\#} Q_{2}$ (Refs. 39 and 40). (79) can be rewritten as

$$
\begin{equation*}
0=\left(A_{c}+B_{c} C Q_{12} Q_{2}^{*}\right) Q_{2}+Q_{2}\left(A_{c}+B_{c} C_{12} Q_{2}^{*}\right)^{T}+B_{c} \hat{V}_{2} B_{c}^{T} \tag{83}
\end{equation*}
$$

Now, using Lemma 12.2 of Ref. 41 it follows from (83) that $Q_{2}$ is positive definite. Similarly, $P_{2}$ is positive definite.

Since $\hat{R}_{2}, \hat{\nabla}_{2}, Q_{2}$ and $P_{2}$ are invertible, (74)-(76) can be written as

$$
\begin{equation*}
-P_{2}^{-1} P_{12}^{T} Q_{12} Q_{2}^{-1}=I_{n_{c}} \tag{84}
\end{equation*}
$$

$$
\begin{align*}
& B_{c}=-P_{2}^{-1}\left[z_{21} C^{T}+P_{12}^{T}\left(V_{12}+\sum_{i=1}^{p} A_{i} Q_{1} C_{i}^{T}\right)\right] \hat{V}_{2}^{-1}  \tag{85}\\
& C_{c}=-\hat{R}_{2}^{-1}\left[B^{T} z_{12}+\left(R_{12}^{T}+\sum_{i=1}^{P} B_{i}^{T} P_{1} A_{i}\right) Q_{12}\right] Q_{2}^{-1} \tag{86}
\end{align*}
$$

Now define new variables

$$
\begin{array}{ll}
Q \triangleq Q_{1}-Q_{12} Q_{2}^{-1} Q_{12^{\circ}}^{T} & P \triangleq P_{1}-P_{12} P_{2}^{-1} P_{12}^{T} \\
\hat{Q} \triangleq Q_{12} Q_{2}^{-1} Q_{12^{*}}^{T} & \hat{P} \triangleq P_{12^{\circ}} P_{2}^{-1} P_{12^{\circ}}^{T} \tag{88}
\end{array}
$$

which are $n \times n$ nonnegative-definite matrices. Note that because of (84). (33) and (34) bold with

$$
\begin{equation*}
G \triangleq Q_{2}^{-1} Q_{12}^{T}, \quad M \triangleq Q_{2} P_{2}, \quad \Gamma \triangleq-P_{2}^{-1} P_{12}^{T} \tag{89}
\end{equation*}
$$

where $M$ is positive semisimple since

$$
Q_{2} P_{2}=Q_{2}^{\frac{1}{2}}\left(Q_{2}^{\frac{1}{2}} P_{2} Q_{2}^{\frac{1}{2}}\right) Q_{2}^{\frac{-1}{2}}
$$

It is helpful to note the identities

$$
\begin{align*}
& \hat{Q}=Q_{12} \mathrm{G}=G^{T} Q_{12}^{T}=G^{T} Q_{2} G, \quad \hat{P}=-P_{12} \Gamma=-\Gamma^{T_{P}} P_{12}^{T}=\Gamma^{T_{P}}{ }_{2} \Gamma,  \tag{90}\\
& \tau G^{T}=G^{T}, \Gamma \tau=\Gamma,  \tag{91}\\
& \hat{Q}=\tau \hat{Q}, \hat{P}=\hat{P} \tau,  \tag{92}\\
& \hat{Q} \hat{P}=-Q_{12} P_{12}^{T} . \tag{93}
\end{align*}
$$

Using (34) and Sylvester's inequality, it follows that $\rho(G)=\rho(\Gamma)=\rho\left(Q_{12}\right)=\rho\left(P_{12}\right)=n_{c}$ which in turn imply (66).

The components of $\bar{Q}_{c}$ and $\bar{P}_{c}$ can be written in terms of $Q, P, \vec{Q}, \vec{P}, G$ and $\Gamma$ as

$$
\begin{array}{ll}
Q_{1}=Q+\hat{Q}, & P_{1}=P+\hat{P} \\
Q_{12}=\hat{Q} \Gamma^{T}, & P_{12}=-\hat{P} G^{T} \\
Q_{2}=\hat{R} \tilde{C}^{T}, & P_{2}=\hat{G P G} \tag{96}
\end{array}
$$

The gain expressions (60) and (61) can now be seen to be equivalent to (85) and (86). Substituting (94)-(96) into (77)-(82) yields
$0=A Q+Q A^{T}+\nabla_{1}+\sum_{i=1}^{P}\left[A_{i} Q A_{i}^{T}+\left(A_{i}-B_{i} \hat{R}_{2}^{-1} \underline{P}\right) \hat{Q}\left(A_{i}-B_{i} \hat{R}_{2}^{-1} \underline{P}\right)^{T}\right]+A_{P} \hat{Q}+\hat{Q} A_{P}^{T}$
$0=\left[A_{P} \hat{Q}+\hat{Q}\left(\Gamma^{T} A_{C} G+C^{T} V_{2}^{-1} Q^{T}\right)+\hat{Q}_{2}^{-1} Q^{T}\right] \Gamma^{T}$.
$0=\Gamma\left[\left(G^{T} A_{c} \Gamma+Q \hat{\nabla}_{2}^{-1} C\right) \hat{Q}+\hat{Q}\left(G^{T} A_{c} \Gamma+Q \hat{V}_{2}^{-1} C\right)^{T}+Q \hat{V}_{2}^{-1} Q^{T}\right] \Gamma^{T}$,
$0=A^{T} P+P A+R_{1}+\sum_{i=1}^{P}\left[A_{i}^{T} P A_{i}+\left(A_{i}-Q V_{2}^{-1} C_{i}\right)^{T} P\left(A_{i}-Q V_{2}^{-1} C_{i}\right)\right]+A_{Q}^{T a}+\hat{P} A_{Q}$.
$0=\left[A_{Q}^{T} \hat{P}+\hat{P}\left(G^{T} A_{C} \Gamma+B \hat{R}_{2}^{-1} \underline{P}\right)+\underline{P}^{T} R_{2}^{-1} \underline{P}\right] G^{T}$.
$0=G\left[\left(G^{T} A_{C} \Gamma+B \hat{R}_{2}^{-1} \underline{P}\right)^{T} \hat{P}+\hat{P}\left(G^{T} A_{c} \Gamma+B \hat{R}_{2}^{-1} \underline{P}\right)+\underline{P}^{T} \hat{R}_{2}^{-1} \underline{P}\right] G^{T}$.

The remaining calculations proceed as follows. Computing either (99) $-\boldsymbol{T}(98)$ or (102)-G(101) yields (59). Inserting this expression for $A$ into (98). (99). (101) and 102 and computing the (equivalent) equations $G^{T}(98)^{T}$. $G^{T}(99) G$. (101) $\Gamma$ and $\Gamma^{T}(102) \Gamma$ it follows that $G^{T}(99) G=G^{T}(98) T^{T}$ and $\Gamma^{T}(102) \Gamma=$ $T(101) \Gamma$. Hence, (99) and (102) are superfluous. Furthermore, $G^{T}(98){ }^{T}$ and (101) $\Gamma$ are equivalent. respectively, to

$$
\begin{align*}
& 0=T\left[A_{P} \hat{Q}+\hat{Q} A_{P}^{T}+Q \hat{V}_{2}^{-1} Q^{T}\right]  \tag{103}\\
& 0=\left[A_{Q}^{T} \hat{P}+\hat{P} A_{Q}+\underline{P}^{T} \hat{R}_{2}^{-1} \underline{P}\right] \tau \tag{104}
\end{align*}
$$

Finally. to obtain (62)-(65) note that $(64)=(103)+(103)^{T}-(103) T^{T},(62)=$ $(97)-(64) .(65)=(104)+(104)^{T}-\tau^{T}(104)$ and $(63)=(100)-(65)$.

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## Robust Controller Synthesis Using the Maximum Entropy Design Equations

## DENNIS S. BERNSTEIN AND SCOTT W. GREELEY


#### Abstract

This note presents an application of the optimality coudrtious obtained in [1] for dynamic compensation in the presence of atate-, coatrol-, and acasurement-dependent aolse. By solving these equations, which represent a fundamental geacralization of standard asendy-state LQG theory, a series of increasiagly robust control designs is obtained for the example considered in [2].


## 1. INTRODUCTION

Perhaps the most significant aspect of LQG theory is the explicit synthesis of dynamic feedback compensators. In practice, however, LQG suffers from serious defects concerning closed-loop robustness with respect to plant deviations. In particular, LQG controllers may possess arbitrarily small stability margin with respect to parameter variations [2].
One approach to correcting this defect is to rederive the optimality conditions for dynamic compensation in the presence of state-, control-, and measurement-dependent noise [1]. Intuitively speaking, the quadratically optimal feedback controller designed in the presence of such multiplicative disturbances is automatically desensitized to actual paramerer variations. The optimality conditions now comprise a system of four matrix equations, specifically, two modified Riccati equations and two modified Lyapunov equations, coupled by stochastic effects. This coupling is a graphic reminder of the breakdown of the separation principle in the uncertain plant case. When the uncertainty terms are absent, the equations immediately reduce to the standard pair of separated Riccati equations.
For the special case of full-order compensation in the presence of statedependent noise only, versions of these equations were discovered independenuly by Hyland [3]-[5] and Mil'stein [6]. A crucial feature of [1], [3]-[5] is the interprectation of the closed-loop stochastic differential equation according to the Fisk-Stratonovich definition of stochastic integration. For modeling flexible mechanical structures, justification of this interpretation as an appropriate model for a priori parameter uncertainty was based upon the maximum entropy principle of Jaynes [1].

A time-varying version of these design equations involving uncorrelated state- and control-dependent noise has been given in [7]. The stochastic interpretation is in the sense of Ito as in [6].

The purpose of the present note is to summarize the maximum entropy equations for full-order dynamic feedback compensation. These equations are then applied to Doyle's example [2] to produce a series of quadratically optimal robust controllers. The full optimal projection/ maximum entropy design equations, which also account for a constraint on controlier order [1]. [8], are applied to a more realistic design problem in [9].

## II. Problem Statement and Maximum Entropy Design Equations

To state the optimal dynamic-compensation problem, we require the following notation. Let $x \in R^{n}, y \in R^{\prime}, u \in R^{m}, A, A_{1}, \cdots, A_{p} \in$ $B^{n \times n}, B, B_{1}, \cdots, B_{p} \in A^{n \times m}, C, C_{1}, \cdots, C_{p} \in B^{1 \times n}, R_{1} \in B^{n \times n}$, $R_{1} \geq 0, R_{2} \in \Omega^{m \times m}, R_{2}>0$. Furthermore, let $v_{1} \cdots, v_{p}$ be unit-intensity, zero-mean, and mutually uncorrelated white noise processes and let $w_{1} \in B^{n}$ and $w_{2} \in \Omega^{\prime}$ be zero-mean white noise processes with

[^16]intensities $V_{1} \geq 0$ and $V_{2}>0$, respectively, and cross intensity $V_{12} \in$ $\boldsymbol{A}^{\times \times 1}$. It is further assumed that $v_{i}, w_{i}$ and $x(0)$ are mutually uncorrelated. We require the technical assumption that, for each $i, B_{i} \neq 0$ implies $C_{i}=$ 0 , i.e., the control- and measurement-dependent noises are uncorrelated.

## Optimal Dynamic-Compensation Problem

Given the controlled system

$$
\begin{gather*}
x^{2}=\left(A+\sum_{i=1}^{f} v_{i} A_{i}\right) x+\left(B+\sum_{i=1}^{P} v_{i} B_{i}\right) u+w_{1}  \tag{2,1}\\
y=\left(C+\sum_{i=1}^{p} v_{i} C_{i}\right) x+w_{2} \tag{2.2}
\end{gather*}
$$

design an nth-order dynamic compensator

$$
\begin{gather*}
\dot{x}_{c}=A_{e} x_{c}+B_{c} y  \tag{2.3}\\
u=C_{c} x_{c} \tag{2.4}
\end{gather*}
$$

which minimizes the performance criterion

$$
\begin{equation*}
J\left(A_{c}, B_{c}, C_{c}\right)=J_{x}\left(A_{c}, B_{c}, C_{c}\right)+J_{s}\left(A_{c}, B_{c}, C_{c}\right)+J_{n}\left(A_{c}, B_{c}, C_{c}\right) \tag{2.5}
\end{equation*}
$$

where

$$
\begin{aligned}
J_{x}\left(A_{c}, B_{c}, C_{c}\right) & \triangleq \lim _{t \rightarrow \infty} I\left[x^{\top} R_{1} x\right] \\
J_{x}\left(A_{c}, B_{c}, C_{c}\right) & \triangleq \lim _{t \rightarrow \infty} E\left[2 x^{\top} R_{12} u\right] \\
J_{n}\left(A_{c}, B_{c}, C_{s}\right) & \triangleq \lim _{r \rightarrow \infty} E\left[u^{\top} R_{2} u\right]
\end{aligned}
$$

To guarantee that $J$ is finite and independent of initial conditions, we restrict ( $A_{c}, B_{c}, C_{c}$ ) to the (open) set of second-moment-stabilizing triples

$$
\delta \widehat{\cong}\left\{\left(A_{c}, B_{c}, C_{c}\right): \overline{A_{s}} \oplus \bar{A}_{2}+\sum_{i=1}^{n} \mathcal{A}_{i} \otimes A_{i} \text { is stable }\right\}
$$

where $f$ and 8 denote Kronecker sum and product and

$$
\begin{gathered}
A_{s} \triangleq\left[\begin{array}{cc}
A_{s} & B_{s} C_{c} \\
B_{c} C_{i} & A_{c}
\end{array}\right], \bar{A}_{i} \dot{\underline{2}}\left[\begin{array}{cc}
A_{i} & B_{i} C_{c} \\
B_{c} C_{i} & 0
\end{array}\right], \\
A_{3} \hat{\geqslant} A_{i}+\frac{1}{2} \sum_{i=1}^{p} A_{i}^{2}, B_{3} \hat{2} B+\frac{1}{2} \sum_{i=1}^{p} A_{i} B_{i}, C_{3} \triangleq C+\frac{1}{2} \sum_{i=1}^{p} C_{i} A_{i} .
\end{gathered}
$$

For convenience in stating the optimality conditions, define the following notation for $Q, P, Q, P \in \beta^{n \times n}$ :

$$
R_{2 s} \triangleq R_{2}+\sum_{i=1}^{n} B_{1}^{T}\left(P+P_{i} B_{i}, V_{2} \triangleq V_{2}+\sum_{i=1}^{p} C_{i}(Q+Q) C_{i}^{T}\right.
$$

Q. $\triangleq Q C_{i}^{r}+V_{12}+\sum_{i=1}^{H} A_{i}(Q+Q) C_{i}^{T}, Q_{i} \triangleq B_{i}^{r} P+R_{12}^{T}+\sum_{i=1}^{P} B_{i}^{T}(P+P) A_{i}$.

$$
A_{Q} \triangleq A_{3}-Q_{3} V_{2 s}^{-1} C_{3}, A_{P_{3}} \triangleq A_{5}-B_{1} R_{25}^{-1} Q_{3}
$$

Theorem 2.I: Suppose $\left(A_{c}, B_{c}, C_{c}\right) \in S$ solves the optimal dynamiccompensation problem. Then there exist $n \times n$ nonnegative-definite matrices $Q, P, Q$. and $P$ such that $A_{c}, B_{c}, C_{c}$ are given by

$$
\begin{equation*}
A_{5}=A_{3}-B_{1} R_{23}^{-1} \Phi_{3}-Q_{2} V_{23}^{-1} C_{3} \tag{2.6}
\end{equation*}
$$

$$
\begin{align*}
& B_{c}=Q_{3} V_{2 s}^{-1}  \tag{2.7}\\
& C_{r}=-R_{2 s}^{-1} Q_{s} \tag{2.8}
\end{align*}
$$

and such that the following conditions are satisfied:

$$
\begin{align*}
& 0=A_{s} Q+Q A_{s}^{r}+V_{1}+\sum_{i=1}^{P}\left(A_{i} Q A_{i}^{r}\right. \\
& \left.+\left(A_{i}-B_{i} R_{2 i}^{-1} Q_{s}\right) Q\left(A_{i}-B_{i} R_{2}^{-1} Q_{s}\right)^{\top}\right]-Q_{i} V_{2 k}^{-1} Q_{s}^{\top},  \tag{2.9}\\
& 0=A_{j}^{T} P+P A_{s}+R_{i}+\sum_{i=1}^{R}\left[A_{i}^{T} P A_{i}\right. \\
& \left.+\left(A_{i}-Q_{s} V_{2 i}^{-1} C_{i}\right)^{T} P_{( }\left(A_{i}-Q_{1} V_{2 s}^{-1} C_{i}\right)\right]-Q_{s}^{r} R_{2 t}^{-1} P_{s},  \tag{2.10}\\
& 0=A_{P_{3}} \phi+Q A_{P_{3}}^{T}+Q_{r} V_{2 S}^{-1} Q_{s}^{T},  \tag{2.11}\\
& 0=A_{Q_{3}}^{T} P+P_{A_{5}}+P_{5}^{r} R_{2}^{-1} P_{5} . \tag{2.12}
\end{align*}
$$

Remark 2.1: Letting $A_{i}=0, B_{i}=0$ and $C_{i}=0, i=1, \cdots, p$, it can readily be seen that (2.11) and (2.12) are superfluous and that (2.9) and (2.10) yield the standard separated LQG Riccati equations.

Remark 2.2. Since $R_{2} \geq R_{2}$, so that $R_{2 s}^{-1} \leq R_{2}^{-1}$, it is clear that the control-dependent noise effectively suppresses the regulator gain $C_{c}$. Similarly, since $V_{2 s} \geq V_{2}$, the measurement-dependent noise suppresses the observer gain $B_{c}$. The effect of the terms $A_{i} Q A_{i}^{T}$ is discussed in [1] for modal systems.

## III. The Maximum Entropy Design Equations Applied to DOYLE'S EXAMPLE

As shown in [2], LQG regulators for the example

$$
\begin{gathered}
A=\left[\begin{array}{ll}
1 & 1 \\
0 & 1
\end{array}\right], B=\left[\begin{array}{l}
0 \\
b
\end{array}\right], C=\left[\begin{array}{ll}
1 & 0
\end{array}\right], \\
V_{1}=\sigma\left[\begin{array}{ll}
1 & 1 \\
1 & 1
\end{array}\right], V_{12}=0, V_{2}=1, \\
R_{1}=\rho\left[\begin{array}{ll}
1 & 1 \\
1 & 1
\end{array}\right], R_{12}=0, R_{2}=1,
\end{gathered}
$$

have arbitrarily small stability margin with regard to variations $b+\Delta b$ when $\sigma$ and $\rho$ are sufficiently large and $b=1$.

Setting $\sigma=\rho=60$, it follows that the LQG regulator is only stable for $0.93 \leq b+\Delta b \leq 1.01$. Uncertainty in $b$ can be modeled by setting $p=$ $1, A_{1}=0, B_{1}=\left[\begin{array}{ll}0 & b_{1}\end{array}\right]^{\top}$, and $C_{1}=0$. Solving the optimality conditions (2.9)-(2.12) with $b_{1}=0.05,0.10,0.15$, and 0.20 yields a series of increasingly robust controller designs with respect to both positive and negative variations $\Delta b$ (see Table I and Figs. 1 and 2).

## Conclusion

As demonstrated on the example of [2], the maximum entropy design equations provide a novel method for synthesizing robust feedback controllers. Since the design equations represen: a fundamental generalization of standard LQG theory, the approach represents an alternative to LQG-modification techniques. Indeed, these equations are not intended as a device for recovering the gain and phase margins of LQ state-feedbeck regulators, but rather as a method for designing output-feedback dynamic compensators which are robust with respect to parmetric devistions in

TABLE 1
DYNAMIC COMPENSATOR OANS FOR LOG AND MAXIMUM ENTROPY DESIGNS $(b=1, \sigma=p=60)$

| 1 | 4 | ${ }^{3}$ | $c_{e}$ | ceability Rame |
| :---: | :---: | :---: | :---: | :---: |
| (10) | $\left[\begin{array}{cc}-9 & 1 \\ -20 & -9\end{array}\right]$ | $\left[\begin{array}{l}10 \\ 10\end{array}\right]$ | $(-10)$ | (.83, 1.01) |
| .08 | $\left[\begin{array}{cc}-9.251 & 1.0 \\ -20.69 & -7.382\end{array}\right]$ | $\left[\begin{array}{l}10.25 \\ 12.31\end{array}\right]$ | [-8.362 -6.382] | (.68, 1.03) |
| . 10 | $\left[\begin{array}{cc}-9.638 & 1.0 \\ -23.87 & -6.310\end{array}\right]$ | $\left[\begin{array}{l}10.64 \\ 15.95\end{array}\right]$ | [-7.318 -7.318] | (.42, 1.48) |
| . 15 | $\left[\begin{array}{cc}-20.10 & 1.0 \\ -27.24 & -5.710\end{array}\right]$ | $\left[\begin{array}{l}21.10 \\ 20.35\end{array}\right]$ | $(-6.710-6.710)$ | (.77. 1.13) |
| . 20 | $\left[\begin{array}{cc}-30.69 & 1.0 \\ -32.87 & -3.295\end{array}\right]$ | $\left[\begin{array}{l}21.69 \\ 28.67\end{array}\right]$ | [-6.395 -6.2931 | (.72, 1.21) |



Fig. 1. Robusthess of LQG versus maximum entropy designs $\left(b_{1}=0.05,0.1,0.15\right.$, 0.2).
the plant model. As discussed in [10], these are significantly different objectives.

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Fig. 2. Scability bounds for LQG and maximum entropy designs.
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# Robust Static and Dynamic Output-Feedback Stabilization: Deterministic and Stochastic Perspectives 

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#### Abstract

Three paraliel gaps in robust feedback control theory are examined: sufficiency versus necessity, deterministic versus stochastic vacertainty modeling, and stability versus performance. Deterministic and stochastic outpur-feedback control problems are considered with both static and dynamic controllers. The static and dyuamic robust stabilization probiems involve deterministically modeled bounded but maknown measurable time-varying parameter variations, while the static and dynamic stochastic optimal control problems feature state-, conirol-, and messurement-dependent white noise. General sufficiency conditions for the deterministic problems are obtained using Lyapunov's direct method, while mecessary conditions for the stochastic problems are derived as a consequence of minimiziag a quadratic performance criterion. The suificieacy tests are then applied to the necessary conditions to determine when solutions of the stochastic optimization problems also solve the deterministic robust stability problems. As an additional application of the deterministic result, the modified Riccati equation approach of Petersen and Hollot is generalized in the static case and exteseded to dynamic compensation.


## I. Introduction

THE gain and phase margins of full-state-feedback LQ regulators are well known [1], [2]. Although dynamic output-feedback LQG controllers lack such margins [3], considerable effort has been devoted to recovering the full-state-feedback properties [4]-[6]. A crucial point discussed in [7]-[9] is that such margins may be meaningless for guaranteeing robustness with respect to arbitrary plant parameter variations. This was demonstrated by means of a simple example in [7]. In addition, as is well known, the use of singular value bounds to characterize plant uncertainty contributes directly to conservatism with respect to real-valued structured parameter variations.
For the parametric-uncertainty stability robustness problem. there exists a considerable body of literature (see, e.g., [10][251). These results often rely upon Lyapunov's direct method and thus are usually in the form of sufficient conditions. Two factors are often lacking, however: a measure of performance beyond stability and design considerations involving controller effon.'

Performance and controller effort are, of course, the natural domain of stochastic optimal control via the cost criterion. In addition, parameter uncertainties can be directly incorporated into the stochastic model by means of multiplicative white noise [26][40]. Heuristically speaking, the performance of a quadratically optimal feedback controller designed in the presence of such

[^17]multiplicative disturbances is desensitized to actual constant or time-varying parameter variations. It should be emphasized that the white noise parameter-uncertainty model is not interpreted literally as a physical model. Rather, the multiplicative noise model serves as a device which captures the effect of parameter uncertainty on the second-moment matrix, and hence on the closed-loop performance. From a practical point of vieu. the multiplicative white noise model is extremely tractable since the second-moment equation is closed and the optimal feedback gains can be given explicitly by closed-form expressions involving solutions of algebraic equations. For example, the necessary conditions derived in [40] for quadratically optimal steady-state fixed-order (i.e., reduced-order) dynamic compensation in the presence of state-, control-, and measurement-dependent white noise involve a coupled system of two modified algebraic Riccati equations and two modified algebraic Lyapunov equations. The coupling is due to both the optimal projection, which enforces the fixed-order constraint [41], and the multiplicative white noise terms. Unfortunately, however, stochastic optimal control is predicated upon stability of the second moment of the state [42][47], which may be weaker than deterministic robust asymptotic stability. As a matter of fact, it has been shown, rather surprisingly, that a nominally unstable system can be rendered stochastically stable by multiplicative white noise interpreted in the sense of Stratonovich without actually applying feedback control [48]. ${ }^{2}$ Hence, there is no prior guarantee that a secondmoment stable optimal controller predicated upon a multiplicative white noise model will provide deterministic robust or even nominal asymptotic stability.

Three parallel gaps can thus be perceived between the above approaches: sufficiency versus necessity, deterministic versus stochastic uncertainty modeling, and stability versus performance. In attempting to bridge these gaps we ask the following question: When is the solution to a stochastic optimal control problem also the solution to a deterministic robust stabilization problem? In the present paper we show that our necessary conditions for stochastic optimality become sufficient conditions for deterministic robustness when we include an exponential weighting factor in the quadratic cost criterion. As shown in [49]. the weighting factor $e^{2 n}$ leads to replacement of the closed-loop dynamics matrices

$$
A+B K C,\left[\begin{array}{cc}
A & B C_{c} \\
B_{c} C & A_{c}
\end{array}\right]
$$

for static and dynamic controllers, respectively, by the shifted dynamics matrices

$$
A+\alpha I_{n}+B K C,\left[\begin{array}{cc}
A+\alpha I_{n} & B C_{c} \\
B c & A_{r}+\alpha I_{n_{c}}
\end{array}\right]
$$

When there are no parametric plant uncertainties, a right shift ( $\alpha$

* This phenomenon does nol occur, however, with the Ito interpretation.
$>0)$ yields a prescribed stability margin, i.e., all closed-loop poles having real part less than $-\alpha$. Unfortunately, a right shift alone does not appear to provide guaranteed stability robustness levels with respect to arbitrary plant variations. Since multiplicative white noise also does not ensure robustness, the present paper goes beyond previous work by employing the right shift in conjunction with multiplicative white noise to guarantee robust stability over a specified range of deterministic parameter variations. In particular, we consider perturbations of $A, B$, and $C$ of the form

$$
A+\sum_{i=1}^{p} \sigma_{i}(t) A_{i}, B+\sum_{i=1}^{p} \sigma_{i}(t) B_{i}, C+\sum_{i=1}^{p} \sigma_{i}(t) C_{i}
$$

where $A_{i}, B_{i}$, and $C_{i}$ denote the pattern of parametric uncertainty in $A, B$, and $C$ and $\sigma_{i}(\cdot)$ are real-valued Lebesgue measurable functions satisfying

$$
\begin{equation*}
\left|\sigma_{i}(t)\right| \leq \delta_{i}, \quad i=1, \cdots, p \tag{1.1}
\end{equation*}
$$

In this formulation the patterns $\mathcal{A}_{i}, D_{i}$, and $C_{i}$ are assumed to be known while the deterministically modeled uncertain parameters $\sigma,(\cdot)$ are known except for the bounds (1.1). Our principal result for both static and dynamic controllers states that a solution of the necessary conditions for stochastic optimal control with exponential weighting and multiplicative white noise provides guaranteed robust asymptotic stability for parametric variations satisfying (1.1) as long as

$$
\begin{equation*}
\alpha>\frac{1}{2} \sum_{i=1}^{p} \delta_{i}^{2} / \gamma_{i} \tag{1.2}
\end{equation*}
$$

where $\alpha$ is the right shift and $\gamma_{i}$ is the intensity of the noise $v_{i}(t)$ multiplying $A_{i}, B_{i}$, and $C_{i}$ in the multiplicative white noise formulation of the stochastic optimal control problem. Clearly, the rectangular robust stability region $\left[-\delta_{1}, \delta_{1}\right] \times \cdots \times\left[-\delta_{p}\right.$, $\delta_{p}$ ] in $p$-dimensional parameter space can be enlarged by increasing either $\alpha$ or $\gamma_{1}, \cdots, \gamma_{p}$. Note that for given values of $\alpha$, $\gamma_{1}, \cdots, \gamma_{p}(1.2)$ does not define a unique robustness region when $p>1$. The robust stability guarantee holds, however, for simultaneous parameter variations $\sigma_{1}(\cdot), \cdots, \sigma_{p}(\cdot)$ within each region satisfying (1.2).

The above result is based upon the observation that secondmoment stability of a stochastic system with multiplicative disturbances and right-shifted mean dynamics induces a Lyapunov function which guarantees robust stability of a deterministic system subject to time-varying parameter variations. This observation, which appears to have been previously overlooked in the literature, may be utilized in the context of robustness analysis for linear uncertain systems. In the present paper this result is developed within the context of robust controller synthesis to achieve a unified approach to robust, fixed-order controller design consistent with [41].

The derivation of our results is quite simple and is based upon the standard quadratic Lyapunov function

$$
\begin{equation*}
V(x)=x^{\top} P x \tag{1.3}
\end{equation*}
$$

where $P$ is given by the modified Lyapunov equation

$$
\begin{equation*}
0=A^{T} P+P A+2 \alpha P+\sum_{i=1}^{D} \gamma_{i} A_{i}^{T} P A_{i}+R \tag{1.4}
\end{equation*}
$$

where $\bar{A}$ is the closed-loop dynamics matrix and $R$ is a closedloop weighting matrix. Note that the third and fourth terms in (1.4) correspond to exponential weighting and multiplicative white noise, respectively. The result that $\dot{V}(x) \leq-\gamma\|x\|^{2}$, with
$\gamma>0$, follows directly from the inequality

$$
\begin{equation*}
\sum_{i=1}^{p} \sigma_{i}(t)\left(A_{i}^{T} P+P A_{i}\right) \leq \sum_{i=1}^{p}\left(\delta_{i}^{2} / \gamma_{i}\right) P+\sum_{i=1}^{p} \gamma_{i} A_{i}^{T} P A_{i} \tag{1.5}
\end{equation*}
$$

along with (1.2). Inequality (1.5) follows immediately from
$\sum_{i=1}^{p^{\prime}}\left[\gamma_{i}^{1 / 2} A_{i}-\left(\sigma_{i}(t) / \gamma_{i}^{1 / 2}\right) I_{n}\right]^{\tau} P\left[\gamma_{i}^{1 / 2} A_{i}-\left(\sigma_{i}(t) / \gamma_{i}^{1 / 2}\right) I_{n}\right] \geq 0$
and $\sigma_{i}^{2}(t) \leq \delta_{i}^{2}$.
An alternative approach to guaranteeing robustness of designs predicated upon a multiplicative white noise model is to interpret the stochastic differential equation according to Stratonovich stochastic integration instead of Ito integration [36]. [37]. Now the closed-loop dynamics matrices become

$$
A_{s}+B_{s} K C_{v}\left(\begin{array}{cc}
A_{s} & B_{s} C_{c} \\
B_{c} C_{s} & A_{s}
\end{array}\right)
$$

where
$A_{s}=A+\frac{1}{2} \sum_{i=1}^{p} \gamma_{i} A_{i}^{2}, B_{s}=B+\frac{1}{2} \sum_{i=1}^{p} \gamma_{i} A_{i} B_{i}, C_{s}=C+\frac{1}{2} \sum_{i=1}^{p} \gamma_{i} C_{i} A_{i}$
The closed-loop Stratonovich correction evident in $A_{s}, B_{s}$, and $C_{s}$, which can negate the so-called "uncertainty threshold principle" [50], [51], appears to be crucial for designing vibration-suppression controllers for flexible structures [37], [38], [52]. Because of inherent damping, such systems are usually nominally open-loop stable with nondestabilizing uncertainties so that robust stability is less of a challenge than robust performance. Although conditions under which the Stratonovich model yields robust controllers are beyond the scope of the present paper. it should be noted that the differences between the two models are far from trivial. For example, for frequency uncertainties the Stratonovich correction, which corresponds to a variable left shift rather than a uniform right shift, automatically induces a positivereal controller for the high-frequency, poorly modeled flexible modes. Since quadratic Lyapunov functions do not appear to be adequate for guaranteeing the robustness of such designs, the majorant Lyapunov function has been developed [53].
Inasmuch as deterministic robust stability of stochastically optimal controllers is guaranteed by right shift/multiplicative white noise modifications to the closed-loop Lyapunov equation, a natural question which arises is the following: Do there exist alternative modifications to the closed-loop Lyapunov equation which guarantee robust stability? One possibility which immediately suggests itself is to replace the bound (1.5) by

$$
\begin{equation*}
\sum_{i=1}^{p} \sigma_{i}(\prime)\left(A_{i}^{T} P+P A_{i}\right) \leq \sum_{i=1}^{p} \delta_{i}\left(P D_{i} D_{i}^{T} P+E_{i}^{T} E_{i}\right) \tag{1.7}
\end{equation*}
$$

where $A_{i}=D_{i} E_{i}$. Carrying out full-state-feedback control design with (1.7) leads to an alternative generalization of the standard algebraic regulator Riccati equation. Indeed, a version of this modified Riccati equation has already been developed by Petersen and Hollot [23], [25] as a means for designing robust static and dynamic controllers. A fourth-order aircraft example considered in [25] shows the practical potential of their approach. In the present paper we extend the results of [23], [25] to more general problem formulations encompassing a wider class of parametric uncertainty structures within the context of static output feedback and reduced-order dynamic compensation. Most interestingly, the results we obtain are completely analogous to the static and
dynamic compensation results obtained using a multiplicative white noise model with quadratic cost. This raises the following interesting question: Does there exist an optimization problem whose necessary conditions coincide with the Petersen-Hollottype equations? Indeed, our results were obtained by optimizing over the class of closed-loop Lyapunov equations modified (i.e., robustified) in the sense of Petersen and Hollot. Full justification for this technique is developed in [54], [55] where robust performance bounds are obtained.

## II. Notation and Definitions

Note: all matrices have real entries.


$$
\begin{aligned}
& K \\
& A_{i}, A_{i} \\
& A_{c}, B_{c}, C_{c}
\end{aligned}
$$

$$
\hat{\bar{A}}, \overline{\bar{A}_{i}}
$$

$$
\sigma_{i}(\cdot)
$$

$$
\begin{aligned}
& \delta_{i} \\
& \alpha \\
& A_{a}, A_{a}, \tilde{A_{a}} \\
& R_{1} \\
& R_{2} \\
& R_{12}
\end{aligned}
$$

$$
\hat{R}
$$

$\stackrel{\rightharpoonup}{\boldsymbol{R}}$
real numbers, $r \times s$ real matrices, $\mathbb{B}^{r \times 1}$
$r \times r$ identity matrix, transpose
Kronecker sum, Kronecker product [56]
$r \times r$ symmetric matrices
$r \times r$ symmetric nonnegative-definite matrices
$r \times r$ symmetric positive-definite matrices
$Z_{1}-Z_{2} \in \mathbb{N}^{\prime}, Z_{1}, Z_{2} \in \boldsymbol{S}^{\prime}$
$Z_{1}-Z_{2} \in \operatorname{Pr}^{r}, Z_{1}, Z_{2} \in \boldsymbol{S}^{r}$
matrix with eigenvalues in open lett-half plane
positive integers, $i=1, \cdots, p$
$n+n_{c}, n+n_{i}, i=1, \cdots, p$
$n, m, l, n_{c}$-dimensional vectors
$n \times n$ matrices; $n \times m$ matrices, $l \times n$
matrices; $i=1, \cdots, p$
$m \times 1$ matrix
$A+B K C, A_{i}+B_{j} K C, i=1, \cdots, p$
$n_{c} \times n_{c}, n_{c} \times l, m \times n_{c}$ matrices

$$
\left[\begin{array}{cc}
A & B C_{c} \\
B_{c} C & A_{c} \\
i=1, \cdots, D
\end{array}\right],\left[\begin{array}{cc}
A_{i} & B C_{c} \\
B_{c} C_{i} & 0
\end{array}\right]
$$

Lebesgue measurable function on
$[0, \infty), i=1, \cdots, p$
positive number, $i=1, \cdots, p$
real number
$A+\alpha I_{n}, A+\alpha I_{n}, A+\alpha I_{A}$
state weighting matrix in $\mathbb{N}^{n}$ control weighting matrix in $P^{m}$
$n \times m$ cross weighting matrix such that $R_{1}$ $-R_{12} R_{2}^{-1} R_{12}^{T} \geq 0$
$\underset{(K C)^{T_{R}}{ }_{R} K C}{R_{12} K C}+(K C)^{T_{1}} R_{12}+$ $(K C){ }^{T} R_{2} K C$

$$
\left[\begin{array}{cc}
R_{1} & R_{12} C_{c} \\
C_{c}^{T} R_{12}^{T} & C_{c}^{T} R_{2} C_{c}
\end{array}\right]
$$

$n$-dimensional Wiener process incremental covariance of $w_{1}$ in $\mathbb{N}^{n}$ $n$-dimensional Wiener process $l$-dimensional Wiener process incremental covariance of $w_{11}$ in $\mathbb{N}^{n}$ incremental covariance of $w_{2}$ in $P^{\prime}$
$n \times l$ incremental cross-covariance of $w_{11}, w_{21}$
$\left[\begin{array}{cc}V_{1} & V_{12} B_{c}^{T} \\ B_{c} V_{12}^{T} & B_{c} V_{2} B_{c}^{T}\end{array}\right]$
murually uncorrelated scalar Wiener processes, $i=1, \cdots, p$ incremental covariance of $v_{l \prime}, \gamma_{i}>0$, $i=1, \cdots, p$ expected value

For the following definitions, let $Q, P, Q, P \in \mathbb{Q}^{n \times n}$ :

$$
\begin{gathered}
R_{2 z} \triangleq R_{2}+\sum_{i=1}^{p} \gamma_{i} B_{i}^{T} P B_{i}, \\
P_{s} \triangleq B^{T} P+R_{i 2}^{T}+\sum_{i=1}^{p} \gamma_{i} B_{i}^{T} P A_{i}, \\
R_{2 d} \triangleq R_{2}+\sum_{i=1}^{p} \gamma_{i} B_{i}^{T}(P+P) B_{i}, \\
P_{d} \triangleq B^{T} P+R_{12}^{T}+\sum_{i=1}^{p} \gamma_{i} B_{i}^{T}(P+P) A_{i}, \\
V_{2 d} \triangleq V_{2}+\sum_{i=1}^{p} \gamma_{i} C_{i}(Q+Q) C_{i}^{T}, \\
Q_{d} \triangleq Q C^{T}+V_{i 2}+\sum_{i=1}^{p} \gamma_{i} A_{i}(Q+\emptyset) C_{i}^{T}, \\
A_{P_{\alpha}} \triangleq A_{\alpha}-B R_{2 d}^{-1} P_{d}, \\
A_{Q_{\alpha}} \triangleq A_{\alpha}-Q_{d} V_{2 d}^{-1} C . \\
\text { II. STATIC OUTPUT FEEDBACK }
\end{gathered}
$$

## A. Static Robust Stabilization Problem: Determinisric Sufficiency Theory

Consider the following problem.
Static Robust Stabilization Problem: Determine $K \in \mathbb{R}^{m \times 1}$ such that the closed-loop system consisting of the controlled plant

$$
\begin{align*}
\dot{x}(t)=(A+ & \left.\sum_{i=1}^{\rho} \sigma_{i}(t) A_{i}\right) x(t) \\
& +\left(B+\sum_{i=1}^{p} \sigma_{i}(t) B_{i}\right) u(t), \quad t \in[0, \infty), \tag{3.1.1}
\end{align*}
$$

measurements

$$
\begin{equation*}
y(t)=C x(t) \tag{3.1.2}
\end{equation*}
$$

and static output-feedback law

$$
\begin{equation*}
u(t)=K y(t) \tag{3.1.3}
\end{equation*}
$$

is asymptotically stable for all measurable $\left(\sigma_{1}, \cdots, \sigma_{p}\right):[0, \infty) \rightarrow$ $\mathbb{R}^{P}$ satisfying

$$
\begin{equation*}
\left|\sigma_{i}(t)\right| \leq \delta_{i}, t \in[0, \infty), i=1, \cdots, p \tag{3.1.4}
\end{equation*}
$$

Remark 3.1.1: The nominal stabilization problem, i.e., the case in which parameter uncertainties are absent, can be recovered by setting $A_{i}=0$ and $B_{i}=0$. All of the results in this paper can readily be specialized to this case. For brevity, however, the details are omitted.

Remark 3.1.2: The symmetric bounds (3.1.4) are for convenience only. The constraints

$$
\begin{equation*}
p_{i} \leq \theta_{i}(t) \leq \delta_{i}, t \in[0, \infty), i=1, \cdots, p \tag{3.1.4}
\end{equation*}
$$

can be recast in the form (3.1.4) by redefining $A$ and $B$. Further notational simplification is possible by scaling $A_{i}$ and $B_{i}$ so that $\delta_{i}$
$=1, i=1, \cdots, p$. For clarity, however, we choose not to do this.

Remark 3.1.3: Standard existence theory guarantees that the absolutely continuous solution $x(\cdot)$ of (3.1.1) exists on $[0, \infty)$ for all $K \in \mathbb{R}^{m \times 1}$ and for all measurable ( $\sigma_{1}, \cdots, \sigma_{p}$ ) satisfying (3.1.4).

The robustness results are based upon the following easily proved theorem which concerns the construction of a Lyapunov function for robust stability.

Theorem 3.1.1: Let $K \in \mathbb{R}^{m \times t}$ and assume there exist

$$
\begin{gather*}
\Phi: \mathbb{P}^{n} \rightarrow S^{n},  \tag{3.1.5}\\
P \in \mathbb{P}^{n},  \tag{3.1.6}\\
M_{i}, N_{i} \in \mathbb{S}_{\mathbb{R}_{i} \times n}, \quad i=1, \cdots, p \tag{3.1.7}
\end{gather*}
$$

such that

$$
\begin{gather*}
M_{i}^{T} N_{i}+N_{i}^{T} M_{i}=\bar{A}_{i}^{T} P+P \bar{A}_{i}, \quad i=1, \cdots, p  \tag{3.1.8}\\
0=\bar{A}^{T} P+P \bar{A}+\Phi(P)  \tag{3.1.9}\\
\sum_{i=1}^{p} \delta_{i}\left(M_{i}^{T} M_{i}+N_{i}^{T} N_{i}\right)<\Phi(P) \tag{3.1.10}
\end{gather*}
$$

Then $K$ solves the static robust stabilization problem.
Proof: Using (3.1.8), compute for $t \in[0, \infty)$ and $i=1$, $\cdots, p$,

$$
\begin{aligned}
0 & \leq\left[\delta_{i}^{1 / 2} M_{i}-\delta_{i}^{-1 / 2} \sigma_{i}(t) N_{i}\right]^{T}\left[\delta_{i}^{1 / 2} M_{i}-\delta_{i}^{-1 / 2} \sigma_{i}(t) N_{i}\right] \\
& =\delta_{i} M_{i}^{T} M_{i}+\delta_{i}^{-1} \sigma_{i}^{2}(t) N_{i}^{T} N_{i}-\sigma_{i}(t)\left[M_{i}^{T} N_{i}+N_{i}^{T} M_{i}\right] \\
& \leq \delta_{i}\left[M_{i}^{T} M_{i}+N_{i}^{T} N_{i}\right]-\sigma_{i}(t)\left[\tilde{A}_{i}^{T} P+P \bar{A}_{i}\right] .
\end{aligned}
$$

Thus,
$\sigma_{i}(\ell)\left[\bar{A}_{i}^{\top} P+P \bar{A}_{i}\right] \leq \delta_{i}\left(M_{i}^{\top} M_{i}+N_{i}^{\top} N_{i}\right)$,

$$
\begin{equation*}
t \in \mid 0, \infty), i=1, \cdots, p \tag{3.1.11}
\end{equation*}
$$

Defining the Lyapunov function

$$
V(x) \triangleq x^{T} P x
$$

its derivative along solutions $x(t)$ of (3.1.1)-(3.1.3) is given by

$$
\begin{aligned}
\dot{V}(x(t))= & \dot{x}(t)^{\top} P x(t)+x(t)^{T} P \dot{x}(t) \\
= & x(t)^{\top}\left[\bar{A}^{T} P+F \bar{A}\right] x(t) \\
& +x(t)^{T}\left[\sum_{i=1}^{p} \sigma_{i}(t)\left(\tilde{A}_{t}^{\top} P+P \bar{A}_{i}\right)\right] x(t)
\end{aligned}
$$

Using (3.1.9) and (3.1.11) y elds

$$
\dot{V}(x(t)) \leq-x(t)^{i}\left[\Phi(P)-\sum_{i=1}^{D} \delta_{i}\left(M_{i}^{T} M_{i}+N_{i}^{T} N_{i}\right)\right] x(t)
$$

By (3.1.10), there exists $\gamma>0$ such that $\dot{V}(x(t)) \leq-\gamma\|x(t)\|^{2}$, $t \in(0, \infty)$, as required.

Remork 3.1.4: As will be seen in later sections. this result is applied by choosing $M_{1}$ and $N$, to satisfy

$$
\begin{equation*}
N_{1}^{T} M_{1}=P \bar{A}_{1} \tag{3.1.12}
\end{equation*}
$$

so that (3.1.8) holds. The bound $\Phi$ is then constructed to satisfy (3.1.10).

## B. Static Optimal Control Problem: Stochastic Necessity Theory

We now turn to the static optimal control problem with statedependent and control-dependent white noise and exponentially weighted quadratic cost.

Static Optimal Control Problem: Determine $K \in \mathbb{R}^{m \times 1}$ such that, for the closed-loop system consisting of the controlled plant

$$
\begin{align*}
& d x_{t}=A x_{1} d t+\sum_{i=1}^{p} A_{i} x_{1} d v_{i t}+B u_{1} d t \\
&+\sum_{i=1}^{p} B_{i} u_{t} d v_{i l}+e^{-\alpha t} d w_{t}, t \in[0, \infty), \tag{3.2.1}
\end{align*}
$$

measurements

$$
\begin{equation*}
y_{t}=C x_{1} \tag{3.2.2}
\end{equation*}
$$

and output-feedback law

$$
\begin{equation*}
u_{1}=K y_{r}, \tag{3.2.3}
\end{equation*}
$$

the performance criterion

$$
\begin{equation*}
J_{a}(K) \triangleq \lim _{t \rightarrow \infty} E e^{2 a t}\left[x_{1}^{T} R_{1} x_{t}+2 x_{1}^{T} R_{12} u_{t}+u_{1}^{T} R_{2} u_{t}\right] \tag{3.2.4}
\end{equation*}
$$

is minimized.
Remark 3.2.1: The exponential time weighting of the disturbance noise in (3.2.1) is required to balance the exponential weighting in the cost (3.2.4). It has no physical significance as such.
To develop necessary conditions for this problem, $K$ must be restricted to the set of second-moment-stabilizing gains

$$
\begin{aligned}
\mathrm{S}_{\alpha} \triangleq\left\{K \in \Omega^{m \times 1}: \bar{A}_{\alpha}\right. & \oplus \bar{A}_{\alpha} \\
& \left.+\sum_{i=1}^{p} \gamma_{i} \bar{A}_{i} \otimes \bar{A}_{i} \text { is asymptotically stable }\right\}
\end{aligned}
$$

For the shifted plant dynamics. The requirement $K \in \delta_{a}$ implies the existence of the steady-state nonnegative-definite covariance $Q \triangleq \lim _{r \rightarrow \infty}\left[e^{2 \alpha 1} x_{t} x_{l}^{T}\right]$. Furthermore, $Q$ is the unique solution to the modified Lyapunov equation

$$
\begin{equation*}
0=\tilde{A}_{a} Q+Q \bar{A}_{\alpha}^{T}+\sum_{i=1}^{D} \gamma_{i} \tilde{A}_{i} Q \tilde{A}_{i}^{T}+V . \tag{3.2.5}
\end{equation*}
$$

An additional technical assumption is that $K$ be confined to the set

$$
S_{\alpha}^{+} \triangleq\left\{K \in \delta_{a}: C Q C^{\top}>0, \text { where } Q \text { satisfies (3.2.5) }\right\}
$$

The positive definiteness condition holds, for example, when $Q$ is positive definite and $C$ has full row rank.

Theorem 3.2.I: Suppose $K \in \delta^{+}$solves the stas . optimal control problem. Then these exist $P, Q \in \oiint^{n}$ such that $K$ is given by

$$
\begin{equation*}
K=-R_{2 s}^{-1} P_{s} Q C^{T}\left(C Q C^{T}\right)^{-1} \tag{3.2.6}
\end{equation*}
$$

and such that $P$ and $Q$ satisfy

$$
\begin{align*}
0=A_{s}^{T} P+P A_{a}+ & \sum_{i=1}^{p} \gamma_{i} A_{i}^{T} P A_{1}+R_{1} \\
& -P_{s}^{T} R_{2}^{-1} P_{s}+\dot{f}_{\lambda}^{T} P_{s}^{T} R_{2 j}^{-1} P_{s} \hat{\gamma}_{1}, \tag{3.2.7}
\end{align*}
$$

$$
\begin{align*}
0=\left(A_{a}\right. & \left.-B R_{2 r}^{-1} P_{s} \hat{r}\right) Q+Q\left(A_{\alpha}-B R_{2 s}^{-1} P_{s} \hat{r}\right)^{T} \\
& +\sum_{i=1}^{p} \gamma_{1}\left(A_{i}-B_{i} R_{2 s}^{-1} P_{s} \hat{r}\right) Q\left(A_{i}-B_{i} R_{2 i}^{-1} P_{s} \hat{r}\right)^{T}+V \tag{3.2.8}
\end{align*}
$$

where

$$
\begin{equation*}
\hat{\tau} \triangleq Q C^{T}\left(C Q C^{T}\right)^{-1} C, \hat{\tau}_{\perp} \triangleq I_{n}-\hat{\tau} . \tag{3.2.9}
\end{equation*}
$$

Proof: First note that from [49] it follows that the exponential factors in (3.2.1) and (3.2.4) are equivalent to replacing $A$ by $A_{\alpha}$. From [57, Theorem 8.5.5, p. 142]. it follows that the state covariance $Q(t) \triangleq E\left[e^{2 a t} x_{t} x_{t}^{T}\right]$ satisfies

$$
Q(t)=A_{\alpha} Q(t)+Q(t) \bar{A}_{\alpha}^{T}+\sum_{i=1}^{p} \gamma_{i} A_{i} Q(t) A_{i}^{T}+V
$$

Since $K \in \delta_{a}, Q \triangleq \lim _{t \rightarrow-} Q(t)$ exists, is nonnegative definite, and is the unique solution to (3.2.5). Note that

$$
J_{\sigma}(K)=\operatorname{tr} Q \hat{R}
$$

Now define the Lagrangian

$$
L(Q, K) \triangleq \mathrm{tr}\left[\lambda Q \bar{R}+\left(\bar{A}_{a} Q+Q \bar{A}_{a}^{T}+\sum_{i=1}^{p} \gamma_{i} \bar{A}_{i} Q \bar{A}_{i}^{T}+V\right) P\right]
$$

with multipliers $\lambda \geq 0$ and $P \in \mathbb{R}^{n \times n}$, and compute

$$
\partial L / \partial Q=A_{\alpha}^{T} P+P \bar{A}_{\alpha}+\sum_{i=1}^{D} \gamma_{i} A_{i}^{T} P \bar{A}_{i}+\lambda R .
$$

Setting $\partial L / \partial Q=0, \lambda=0$ implies $P=0$ since $K \in \delta_{\alpha}$. Hence, without loss of generality set $\lambda=1$ so that

$$
\begin{equation*}
0=\bar{A}_{\alpha}^{T} P+P \bar{A}_{\alpha}+\sum_{i=1}^{p} \gamma_{i} \bar{A}_{i}^{T} P \bar{A}_{i}+\hat{R} \tag{3.2.10}
\end{equation*}
$$

Since $P$ is the (unique) steady-state covariance of the dual system, it is nonnegative definite. Also, since $\delta_{\alpha}^{+}$is open, evaluating $\partial L /$ $\partial K=0$ yields

$$
0=R_{2 s} K C Q C^{T}+P_{s} Q C^{T}
$$

Since $K \in S_{\alpha}^{+}, C Q C^{r}$ is invertible, and hence (3.2.6) holds. Finally. ( 3.2 .8 ) is equi alent to (3.2.5) and. performing some algebraic manipulation, (3.2.7) is equivalent to (3.2.10).

Remark 3.2.2: Because $R_{2 s} \geq R_{2}$, and thus $R_{25}^{-1} \leq R_{2}^{-1}$, the gain $K$ in (3.2.6) leads to a net decrease in controller authority due to the control-dependent noise. For problems which are open-loop stable and which remain stable under plant variations. this is an intuitively expected consequence of parameter uncertainty.

Remark 3.2.3: Theorem 3.2.1 generalizes and unifies several previous results. In particular, the noise correlation and output feedback constraint constitute generalizations of [26]-[33]. Furthermore, because of the presence of multiplicative noise, the results of [58], [59] are extended. The role of the oblique projection $\boldsymbol{\tau}$ has been discussed in [58], [59]. Connections with the oblique projection $r$ arising in the dynamic-compensation problem [14] are discussed in [60].

## C. Sufficiency Meets Necessity: A Marriage of the Deterministic and Stochastic

We now answer our main question: Can a feedback law predicated on a stochastic multiplicative noise model provide guarameed deterministic robust asymptotic stability? The answer is "yes" provided the exponential is of sufficient magnitude.

Theorem 3.3.1: Suppose there exists $P \in P^{n}$ and $Q \in \mathbb{N}^{n}$ satisfying $C Q C{ }^{\top}>0,(3.2 .7)-(3.2 .9)$ and

$$
\begin{equation*}
0<\left(2 \alpha-\sum_{i=1}^{p} \delta_{i}^{2} / \gamma_{i}\right) P+\vec{R} \tag{3.3.1}
\end{equation*}
$$

where $K$ in $R$ is given by (3.2.6). Then $K$ solves the static robust stabilization problem.

Proof: In Theorem 3.1.1 define

$$
\begin{gathered}
\Phi(P)=2 \alpha P+\sum_{i=1}^{p} \gamma_{i} A_{i}^{T} P A_{i}+\bar{R}, \\
M_{i}=\left(\gamma_{i} / \delta_{i}\right)^{1 / 2} P^{1 / 2} A_{i}, N_{i}=\left(\delta_{i} / \gamma_{i}\right)^{1 / 2} P^{1 / 2} .
\end{gathered}
$$

Note that (3.1.5)-(3.1.8) hold. Furthermore, because of the equivalence of (3.2.10) and (3.2.7), it follows that (3.1.9) is equivalent to (3.2.7). Finally, (3.1.10) is a consequence of (3.3.1).

Remark 3.3.1: Note that (3.2.7)-(3.2.9) serve to construct a Lyapunov function guaranteeing robust stability. Hence, it is not necessary to actually verify that $K \in \mathcal{S}^{+}$.
By strengthening (3.3.1) the following simplification is immediate.

Corollary 3.3. I: Suppose there exists $P \in P^{n}$ and $Q \in \mathrm{Jgn}^{n}$ satisfying $C Q C^{r}>0$, (32.7)-(3.2.9) and

$$
\begin{equation*}
\alpha>\frac{1}{2} \sum_{i=1}^{p} \sigma_{i}^{2} / \gamma_{i} \tag{3.3.2}
\end{equation*}
$$

Then $K$ given by (3.2.6) solves the static robust stabilization problem.

It is interesting to note that the feedback gain given by Corollary 3.3.1 may be an extremal, i.e., local minimum, local maximum, etc., and not necessarily a solution of the static optimal control problem. The result is valid, however, for all extremals of the optimization problem. By specializing Corollary 3.3.1 to a solution, i.e., global minimum, of the optimal control problem, we can bridge the gap between sufficiency and necessity.

Corollary 3.3.2: Suppose $K \in \delta_{a}^{+}$solves the static optimal control problem where $\alpha$ satisfies (3.3.1), and suppose that the corresponding solution $P$ of (3.2.7) is positive definite. Then $K$ also solves the static robust stabilization problem.

## IV. Dynamic Output Feedback

## A. Dynamic Robust Stabilization Problem: Deterministic Sufficiency Theory

Consider the following problem.
Dynamic Robust Stabilization Problem: Determine ( $A_{c}, B_{c}$, $\left.C_{c}\right)$ such that the closed-loop system consisting of the controlled plant (3.1.1), measurements

$$
\begin{equation*}
y(t)=\left(C+\sum_{i=1}^{p} \sigma_{i}(t) C_{i}\right) x(t) \tag{4.1.1}
\end{equation*}
$$

and dynamic output-feedback law

$$
\begin{gather*}
\dot{x}_{c}(t)=A_{c} x_{c}(t)+B_{c} y(t)  \tag{4.1.2}\\
u(t)=C_{c} x_{c}(t) \tag{4.1.3}
\end{gather*}
$$

is asymptotically stable for all measurable $\left(\sigma_{1}, \cdots, \sigma_{p}\right):[0, \infty) \rightarrow$ $\mathbb{R}^{\rho}$ satisfying (3.1.4).

Remark 4.1.1: Note that the problem statement places no restriction on the order $n_{r}$ of the dynamic compensator. Also, we now permit uncerainties in the observation matrix $C$ by including perturbations $\sigma_{i}(f) C_{i}$ in (4.1.1).

The following result is completely analogous to Theorem 3.1.1. Theorem 4.1.1: Given $\left(A_{c}, B_{c}, C_{c}\right)$ assume there exist

$$
\begin{gather*}
\Phi: \mathbb{P}^{\boldsymbol{n}} \rightarrow \Phi^{\boldsymbol{n}},  \tag{4.1.4}\\
\boldsymbol{P}^{\boldsymbol{A}} \in \mathbb{P}^{\boldsymbol{n}},  \tag{4.1.5}\\
\bar{M}_{i}, \bar{N}_{i} \in \mathbb{R}^{n_{i} \times \boldsymbol{n}}, \quad i=1, \cdots, p \tag{4.1.6}
\end{gather*}
$$

such that

$$
\begin{align*}
& \bar{M}_{i}^{T} \bar{N}_{i}+\bar{N}_{i}^{T} \bar{M}_{i}=\overline{\tilde{A}}_{i}^{T} \bar{P}^{+}+\overline{\bar{A}_{i}}, \quad i=1, \cdots, p,  \tag{4.1.7}\\
& 0=\overline{\bar{A}}^{\top} \boldsymbol{P}+\overline{\mathcal{P}} \overline{\bar{A}}+\bar{\Phi}(\bar{P}),  \tag{4.1.8}\\
& \sum_{i=1}^{p} \delta_{i}\left(\tilde{M}_{i}^{T} \tilde{M_{i}}+\tilde{N}_{i}^{\tau} \tilde{N_{i}}\right)<\boldsymbol{\Phi}(\tilde{P}) . \tag{4.1.9}
\end{align*}
$$

Then $\left(A_{c}, B_{c}, C_{c}\right)$ solves the dynamic robust-stabilization problem.

## B. Dynamic Optimal Control Problem: Stochastic Necessity Theory

We now consider the dynamic optimal control problem with state-, control-, and measurement-dependent white noise and exponentially weighted quadratic cost. The optimization is performed over the class of dynamic compensators of fixed order $n_{c}$ $\leq n$.

Dynamic Optimal Control Problem: Determine $\left(A_{c}, B_{c}, C_{c}\right)$ such that, for the closed-loop system consisting of the controlled plant

$$
\begin{align*}
d x_{t}=A x_{i} d t & +\sum_{i=1}^{p} A_{i} x_{t} d v_{i t}+B u_{t} d t \\
& +\sum_{i=1}^{p} B_{i} u_{t} d v_{i l}+e^{-a t} d w_{1 t}, t \in[0, \infty) \tag{4.2.1}
\end{align*}
$$

measurements

$$
\begin{equation*}
d y_{t}=C x_{1} d t+\sum_{i=1}^{p} C_{i} x_{1} d v_{a l}+e^{-a t} d w_{2 t} \tag{4.2.2}
\end{equation*}
$$

and dynamic output-feedback law

$$
\begin{gather*}
d x_{c l}=A_{c} x_{c t} d t+B_{c} d y_{l}  \tag{4.2.3}\\
u_{t}=C_{r} x_{c t} \tag{4.2.4}
\end{gather*}
$$

the performance criterion

$$
\begin{align*}
J_{a}\left(A_{c}, B_{\mathrm{c}}, C_{c}\right) \triangleq \lim _{t \rightarrow \infty} E e^{2 a t}\left[x_{1}^{T} R_{1} x_{t}\right. & \\
& \left.+2 x_{t}^{T} R_{12} u_{t}+\mu_{t}^{T} R_{2} \mu_{t}\right] \tag{4.2.5}
\end{align*}
$$

is minimized.
To develop necessary conditions we restrict $\left(A_{c}, B_{c}, C_{c}\right)$ to the set

$$
D_{a} \triangleq\left\{\left(\overline{A_{c}}, B_{c}, C_{c}\right): \hat{A}_{n} \oplus \bar{A}_{a}+\sum_{i=1}^{p} \gamma_{i} \vec{A}_{i} \otimes \bar{A}_{i}\right. \text { is }
$$

asymptotically stable and $\left(A_{c}, B_{c}, C_{c}\right)$ is minimal $\}$ and invoke the technical assumption

$$
\begin{equation*}
\left[B_{i} \neq 0 \Rightarrow C_{i}=0\right], \quad i=1, \cdots, p . \tag{4.2.6}
\end{equation*}
$$

The following lemma will be needed.

Lemma 4.2.1: If $Q, P \in \mathbb{M}^{n}$ and $\operatorname{rank} \phi P=n_{c}$, then there exist $G, \Gamma \in \mathbb{R}^{n_{c} \times n}$ and invertible $M \in \mathbb{R}^{n_{c} \times n_{c}}$ such that

$$
\begin{gather*}
Q P=G^{T} M \Gamma  \tag{4.2.7}\\
\Gamma G^{T}=I_{r c} . \tag{4.2.8}
\end{gather*}
$$

Furthermore, $G, M$, and $\Gamma$ are unique modulo a change of basis in $p^{n_{c} \times{ }^{n} c \text {. }}$

Proof: The result is an immediate consequence of $\mathbf{6 1}$, Theorem 6.2.5, p. 123].

Note that because of (4.2.8), the $n \times n$ matrix $r=G^{T} \Gamma$ is idempotent, i.e., $\boldsymbol{T}^{2}=\tau$. Since $\tau$ is not necessarily symmetric, it is an oblique projection.

Theorem 4.2.1: Suppose (4.2.6) holds and $\left(A_{c}, B_{c}, C_{c}\right) \in D_{a}$ solves the dynamic optimal control problem. Then there exjst $Q$, $P, \hat{Q}, \hat{P} \in \mathrm{~N}^{n}$ such that $A_{c}, B_{c}, C_{c}$ are given by

$$
\begin{gather*}
A_{c}=\Gamma\left(A-B R_{2 d}^{-1} P_{d}-Q_{d} V_{2 d}^{-1} C\right) G^{T},  \tag{4.2.9}\\
B_{c}=\Gamma Q_{d} V_{2 d}^{-1},  \tag{4.2.10}\\
C_{c}=-R_{2 d}^{-1} P_{d} G^{T} \tag{4.2.11}
\end{gather*}
$$

and such that $Q, P, \boldsymbol{Q}, P_{\text {satisfy }}$

$$
\begin{align*}
& 0=A_{a} Q+Q A_{a}^{T}+V_{1}+\sum_{i=1}^{p} \gamma_{i}\left[A_{i} Q A_{i}^{T}+\left(A_{i}-B_{i} R_{2 d}^{-\frac{1}{d}} P_{d}\right) Q\left(A_{i}\right.\right. \\
& \left.\left.-B_{i} R_{2 d}^{-1} P_{d}\right)^{T}\right]-Q_{d} V_{2 d}^{-1} Q_{d}^{T}+\tau_{\perp} Q_{d} V_{2 d}^{-1} Q_{d}^{T} \tau_{1}^{T},  \tag{4.2.12}\\
& 0=A_{a}^{T} P+P A_{\alpha}+R_{1}+\sum_{i=1}^{p} \gamma_{i}\left[A_{i}^{T} P A_{i}+\left(A_{i}-Q_{d} V_{2 d}^{-1} C_{i}\right)^{T} P\left(A_{i}\right.\right. \\
& \left.\left.-Q_{d} V_{2 d}^{-1} C_{i}\right)\right]-P_{d}^{T} R_{2 d}^{-1} P_{d}+\tau_{\perp}^{T} P_{d}^{T} R_{2 d}^{-1} P_{d} \tau_{\perp},  \tag{4.2.13}\\
& 0=A_{P_{G}} Q+Q A_{P_{Q}}^{T}+Q_{d} V_{2 d}^{-1} Q_{d}^{T}-T_{\perp} Q_{d} V_{2 d}^{-1} Q_{d}^{T} T_{\perp}^{T},  \tag{4.2.14}\\
& 0=A_{Q_{a}}^{T} P_{+}+P A_{Q_{a}}+P_{d}^{T} R_{2 d}^{-1} P_{d}-\tau_{\perp}^{T} P_{d}^{T} R_{2 d}^{-1} P_{d} \tau_{\perp}  \tag{4.2.15}\\
& \operatorname{rank} \phi=\operatorname{rank} P=\operatorname{rank} \phi \boldsymbol{P}=\boldsymbol{n}_{\boldsymbol{c}} \tag{4.2.16}
\end{align*}
$$

(4.2.7) and (4.2.8), where

$$
\begin{equation*}
\tau \triangleq G^{T} \Gamma, \tau_{\perp} \triangleq I_{n}-\tau \tag{4.2.17}
\end{equation*}
$$

Proof: As in the proof of Theorem 3.2.1 we note that the exponential factors in (4.2.1), (4.2.2), and (4.2.5) are equivalent to replacing $A$ and $A_{c}$ by $A+\alpha I_{n}$ and $A_{c}+\alpha I_{n c}$, respectively. Theorem 4.2.1. now follows immediately from [40. Theorem 2.3]. It need only be noted that (4.2.9) follows from

$$
A_{c}+\alpha I_{n_{c}}=\Gamma\left(A_{\alpha}-B R_{2 d}^{-1} P_{d}-Q_{d} V_{2 d}^{-1} C\right) G^{T}
$$

and the fact that $\Gamma A_{\alpha} G^{T}=\Gamma A G^{T}+\alpha I_{n_{c}}$ because of (4.2.8).

## C. Sufficiency Meets Necessity: The Dynamic Case

We now bridge the gap between sufficiency and necessity for dynamic controllers.

Theorem 4.3.1: Assume (4.2.6) holds and suppose there exists $Q, P, Q, P \in \mathbb{N}^{\prime}$ satisfying (4.2.12)-(4.2.17), (4.2.7), (4.2.8),

$$
P \Leftrightarrow\left[\begin{array}{ll}
P+\beta & -P G^{r}  \tag{4.3.1}\\
-G P & G P G^{r}
\end{array}\right]>0
$$

and

$$
\begin{equation*}
0<\left(2 \alpha-\sum_{i=1}^{p} \delta_{i}^{2} / \gamma_{i}\right) \hat{F}+\overline{\bar{R}} \tag{4.3.2}
\end{equation*}
$$

Then $\left(A_{c}, B_{c}, C_{c}\right)$ given by (4.2.9)-(4.2.11) solves the dynamic robust stabilization problem.

Remark 4.3.1: Note that $\overline{\mathcal{P}}$ is always at least nonnegative definite since

$$
\bar{P}=\left[\begin{array}{ll}
P & 0  \tag{4.3.3}\\
0 & 0
\end{array}\right]+\left[\begin{array}{c}
P_{1}^{1 / 2} \\
-G P^{1 / 2}
\end{array}\right]\left[\begin{array}{c}
P_{1 / 2} \\
-G P^{1 / 2}
\end{array}\right]^{r} \geq 0
$$

Proof of Theorem 4.3.1: As shown in [40], (4.2.6)(4.2.17) are equivalent to

$$
\begin{align*}
& 0=\bar{A}_{\alpha} Q+Q \bar{A}_{\alpha}^{T}+\sum_{i=1}^{p} \gamma_{i} \bar{A}_{i} \bar{Q}_{\bar{A}}^{T}+\bar{V} \tag{4.3.5}
\end{align*}
$$

where

$$
\phi=\left[\begin{array}{cc}
Q+Q & \phi \Gamma^{r} \\
\Gamma Q & \Gamma \phi \Gamma^{r}
\end{array}\right]
$$

The result now follows from Theorem 4.1.1 as in the proof of Theorem 3.3.1.

Remark 4.3.2: As in Corollary 3.3.1 the inequality (4.3.2) can be replaced by the stronger condition (3.3.2).
V. The Petersen-Hollot Approach to Robust Stabilization

## A. Static Output Feedback

The deterministic Riccati equation approach of Petersen and Hollot is based upon factoring $A_{i}$ and $B_{i}$ as

$$
\begin{equation*}
A_{i}=D_{i} E_{i}, B_{i}=D_{i} F_{i}, \quad i=1, \cdots, p \tag{5.1.1}
\end{equation*}
$$

where $D_{i} \in \Omega^{n \times n_{i}}, E_{i} \in \mathbb{R}^{n_{i} \times n}$, and $F_{i} \in \mathbb{B}^{n_{i} \times m}$. Obviously. such a factorization may not be unique, and the nonuniqueness is an element of the sufficiency test. To state the sufficiency condition we shall require the notation

$$
\begin{aligned}
& R_{20} \triangleq R_{2}+\sum_{i=1}^{p} \delta_{i} F_{i}^{T} F_{i}, P_{0} \leqq B^{T} P+R_{12}^{T}+\sum_{i=1}^{p} \delta_{i} F_{i}^{T} E_{i} \\
& Q_{e}=Q C^{T}+V_{12}, D \triangleq \sum_{i=1}^{p} \delta_{i} D_{i} D_{i}^{T}, E \triangleq \sum_{i=1}^{p} \delta_{i} E_{i}^{T} E_{i}
\end{aligned}
$$

Theorem 5.1.1: Assume there exist $P \in P^{n}$ and $Q \in \mathbb{N}^{n}$ satisfying $C Q C^{r}>0$,

$$
\begin{equation*}
0=A^{T} P+P A+P D P+R_{1}+E-P_{a}^{T} R_{2 a}^{-1} P_{a}+\hat{i}_{1}^{T} P_{a}^{T} R_{2 a}^{-1} P_{a} \tilde{i}_{1} \tag{5.1.2}
\end{equation*}
$$

$$
\begin{equation*}
0=\left(A+D P-B R R_{2 a}^{-1} P_{a} \hat{\tau}\right) Q+Q\left(A+D P-B R_{2 a}^{-1} P_{a} \hat{\tau}\right)^{T}+V \tag{5.1.3}
\end{equation*}
$$

and

$$
\begin{equation*}
R>0 \tag{5.1.4}
\end{equation*}
$$

where $\hat{r}$ is given by (3.2.9) and $K$ in $\hat{R}$ is given by

$$
\begin{equation*}
K=-R_{2 a}^{-1} P_{a} Q C^{T}\left(C Q C^{T}\right)^{-1} . \tag{5.1.5}
\end{equation*}
$$

Proof: In Theorem 3.1.1 define

$$
\begin{gathered}
\Phi(P)=R+\sum_{i=1}^{p} \delta_{i}\left[\left(E_{i}+F_{i} K C\right)^{T}\left(E_{i}+F_{i} K C\right)+P D_{i} D_{i}^{\tau} P\right], \\
M_{i}=E_{i}+F_{i} K C, N_{i}=D_{i}^{T} P
\end{gathered}
$$

so that (3.1.5)-(3.1.8) hold. Note that (5.1.2) and (5.1.3) are equivalent to (3.1.9). Finally, because of (5.1.4), (3.1.10) is satisfied.
Remark 5.1.1: Theorem 5.1.1 generalizes [25, Theorem 3.3] in four distinct ways. First, in [25], the uncertainties $A$ and $B$ are assumed to be independent. In our setup this corresponds to the additional condition that $A_{i} \neq 0$ implies $B_{i}=0$. In this case (5.1.1) is satisfied with either $E_{i}=0$ or $F_{i}=0$ for each $i$. Second, in [25], $A_{i}$ and $B_{i}$ are confined to have unit rank. This would be the case if we required $n_{i}=1$ in (5.1.1). Third, in [25], $R_{12}=0$. And, fourth, the results of [25] apply only to the statefeedback case, i.e., $C=I_{n}$. In this case (5.1.3) is superfluous and $\hat{\boldsymbol{\tau}}=I_{n}$.

## B. Dynamic Output Feedback

We now extend the Petersen-Hollot approach to reduced-order dynamic compensation. Our only constraint is that we do not permit uncertainty in the observation matrix. Define

$$
A_{P} \triangleq A-B R_{20}^{-1} P_{a}, A_{Q} \triangleq A-Q_{a} V_{2}^{-1} C
$$

Theorem 5.2.1: Assume $C_{1}=\cdots=C_{p}=0$ and suppose there exist $Q, P, Q, P \in \mathbb{N}^{n}$ satisfying

$$
\begin{align*}
& 0=[A+D(P+P)] Q+Q[A+D(P+F)]^{T} \\
&+V_{1}-Q_{a} V_{2}^{-1} Q_{a}^{T}+r_{1} Q_{a} V_{2}^{-1} Q_{0}^{T} T_{1}^{T}, \quad(5.2 .1  \tag{5.2.1}\\
& 0=A^{T} P+P A+R_{1}+E+P D P-P_{a}^{r} R_{2}^{-1} P_{a}+\tau_{1}^{T} P_{a}^{T} R_{2 a}^{-1} P_{a} T_{l}, \tag{5.2.2}
\end{align*}
$$

$0=\left(A_{P}+D P\right) Q+Q\left(A_{P}+D P\right)^{T}+Q_{0} V_{2}^{-1} Q_{a}^{T}-\tau_{\perp} Q_{0} V_{2}^{-1} Q_{a}^{T_{r}}{ }_{2}$,

$$
\begin{equation*}
0=\left(A_{Q}+D P\right)^{\top} P+P\left(A_{Q}+D P\right)+P D P \tag{5.2.3}
\end{equation*}
$$

$$
\begin{equation*}
+P_{a}^{T} R_{2 a}^{-1} P_{a}-\tau{ }_{\perp}^{T} P_{a}^{T} R_{2 a}^{-1} P_{a} \tau_{\perp} \tag{5.2.4}
\end{equation*}
$$

(4.2.7), (4.2.8). (4.2.16), (4.3.1), and

$$
\left[\begin{array}{cc}
R_{1} & -R_{12} R_{2 a}^{-1} P_{a} G^{T} \\
-G P_{a}^{T} R_{2 a}^{-1} R_{12}^{T} & G P_{a}^{r} R_{2 a}^{-1} R_{2} R_{2 a}^{-1} P_{a} G^{T}
\end{array}\right]>0 .
$$

Then $\left(A_{c}, B_{c}, C_{\dot{c}}\right)$ given by

$$
\begin{gather*}
A_{c}=\Gamma\left(A-Q_{a} V_{2}^{-1} C-B R_{2 a}^{-1} P_{a}+D P\right) G^{T}  \tag{5.2.5}\\
B_{c}=\Gamma Q_{a} V_{2}^{-1} \tag{5.2.6}
\end{gather*}
$$

solves the dynamic robust stabilization problem.
Proof: In Theorem 4.1.1 define

$$
\begin{gathered}
\Phi(F)=R+\sum_{i=1}^{p} \delta_{i}\left(E_{i}^{T} E_{i}^{*}+P D_{i} D_{i}^{T} P\right), M_{i}=E_{i}, N_{i}=D_{i}^{T} P, \\
B_{i}=\left[\begin{array}{cc}
D_{i} & 0 \\
0 & 0
\end{array}\right], E_{i}=\left[\begin{array}{cc}
E_{1} & F_{i} C_{i} \\
0 & 0
\end{array}\right]
\end{gathered}
$$

and note that (5.2.2) and (5.2.4) correspond to (4.1.8). Also, note that (5.2.1) and (5.2.3) correspond to

$$
\begin{equation*}
0=\left(A+\sum_{i=1}^{p} \delta_{i} D_{i} D_{i} T \bar{P}\right) Q+Q\left(\tilde{A}+\sum_{i=1}^{p} \delta_{i} D_{i} \tilde{D}_{i} T \bar{P}\right)^{T}+\tilde{V}^{E} \tag{5.2.8}
\end{equation*}
$$

Remark 5.2.1: Theorem 5.2.1 was discovered by optimizing $J$ $=\mathbf{t} \boldsymbol{P} \boldsymbol{V}$ subject to

$$
\begin{equation*}
0=\tilde{A}^{T} \tilde{P}+\hat{P A}+\hat{R}+\sum_{i=1}^{n} \delta_{i}\left[E_{i}^{T} E_{i}+\dot{P} \dot{D}_{i} \hat{D}_{i}^{\tau} \tilde{P}\right] \tag{5.2.9}
\end{equation*}
$$

and using the techniques of [40], [41]. As shown in [54], this approach also yields robust performance bounds.

## Acknowledgment

The author wishes to thank Prof. C. V. Hollot for numerous helpful discussions and for carefully reviewing the original manuscript of the paper.

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Dennis S. Bernstein (M•82), for a photograph and biography, see p. 1013 of the November 1987 issue of this TrANSACTIONS.

## APPENDIX G: Robust Analysis

115. D. S. Bernstein and W. M. Haddad, "Robust Stability and Performance Analysis for Linear Dynamic Systems," IEEEE Trans, Autam. Contr., Vol. 34, 1989, to appear.
116. D. S. Bernstein and W. M. Haddad, "Robust Stability and Performance Analysis for State Space Systems via quadratic Iyapunov Bounds," submitted to SIAM J. Matrix Anal. Appl.
117. D. C. Hyland and D. S. Bernstein, "The Majorant Iyapunov Equation: A Nomegative Matrix Equation for Guaranteed Robust Stability and Performance of Large Scale Systems," IFEE Trans, Autom. Contr., Vol. AC-32, pp. 1005-1013, 1987.

# Robust Stability and Performance Analysis for 

Linear Dynamic Systems

## by

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#### Abstract

In a recent paper Zhou and Khargonekar obtained sufficient conditions for robust stability over specified sets of matrix perturbations. In the present paper these results are extended to include, in addition, performance bounds. Here performance is defined as the worst-case expected value of a quadratic functional involving the state variables when the system is subjected to white noise disturbances. The results are illustrated by considering the gain margin of both an LQG controller and a robustified design obtained by Bernstein and Greeley for Doyle's example.


Suppurted in part by the Air Force Office of Scientific Research under contract F49620-86-C0002.

## 1. Introduction

It is well known that unavoidable discrepancies between mathematical models and real-world systems can result in the degradation of control-system performance. Ideally, feedback control systems should be designed to be robust with respect to uncertainties in the plant characteristics. Thus robustness analysis must play a key role in control-system design. That is, given an existing or proposed control system, determine the performance degradation due to variations in the plant. The most fundamental concern in this regard is clearly that of stability. For linear state space systems with which the present paper is concerned, this problem has received increasing attention over the past several years (see, e.g., [1-12]).

One of the principal techniques used to assess robust stability is based upon quadratic Lyapunov functions (see [1-4,10]). Quadratic Lyapunov functions have also been used extensively for robust control-system synthesis; see [13] for relevant references. The problem of robust synthesis is, however, beyond the scope of the present paper.

In addition to assessing robust stability, it is often desirable to quantify performance by considering the degradation of a cost functional as the plant parameters deviate from their nominal values. Although any robustly stable system over a compact set of parameters possesses a worst-case performance, it is desirable in practice to actually determine a bound for the worst-case performance. The concern for both robust stability and performance goes back to the early work of Michael and Merriam [14], while more recent references include the work of Chang and Peng [15], Noldus [16], and Petersen [17]. The results of [15-17] can be shown to depend upon a modified Lyapunov equation of the form

$$
\begin{equation*}
0=A Q+Q A^{T}+\hat{\Omega}(Q)+V \tag{1.1}
\end{equation*}
$$

where the operator $\hat{\Omega}(Q)$ is chosen to bound terms of the form $\Delta A Q+Q \Delta A^{T}$, where $\Delta A$ is an uncertain perturbation of the dynamics matrix $A$. Since robust performance per se was not discussed in $[16,17]$, the work most closely related to the present paper is that of Chang and Peng [15]. They essentially show that consideration of (1.1) leads to a bound on worst-case performance. Although the development in [15] was carried out for full-state feedback, specialization of their approach to robust performance analysis is straightforward. A systematic, in-depth treatment of robust performance analysis involving the approach of [15] as well as other bounds is given in [18].

The starting point for the present paper is the recent paper by Zhou and Khargonekar [10]. By analyzing the Lyapunov equation they obtain a series of stability robustness tests which improve
significantly upon earlier work [2-4]. In the present paper we extend the results of [10] to obtain, in addition, a bound on worst-case performance. As in (1.1) we consider a Lyapunov equation of the form

$$
\begin{equation*}
0=A Q+Q A^{T}+\Omega+V \tag{1.2}
\end{equation*}
$$

where $\Omega$ bounds uncertainty terms of the form $\triangle A Q+Q \Delta A^{T}$. The principal difference between (1.1) and (1.2) is that $\Omega$ in (1.2) is a constant matrix independent of the solution $Q$. The case considered in [15] in which $\Omega$ is a function of $Q$ is discussed in [18].

The cost functional used in the present paper to quantify robust performance is the trace of the output covariance of a system subjected to white noise disturbances. This measure of performance is identical in form to the standard performance criterion of LQG theory. Since we also obtain a bound for the state covariance matrix, our results yield bounds on the variances (mean square response levels) of system states. Although the results of [15] were obtained within a deterministic setting, it is easy to see that the performance criterion of [15] is also of this form.

The contents of the paper are as follows. After introducing notation at the end of this section we consider the robust stability and performance problems in Section 2. In Section 3 we present the main result (Theorem 3.1) which provides sufficient conditions for robust stability over a set of parameter variations along with a performance bound. In Section 4 we present a dual result (Theorem 4.1) in terms of the dual matrix $P$. This result serves two purposes. First, it clarifies connections with the previous literature where results are presented in terms of the quadratic Lyapunov function $V(x)=x^{T} P x$. And, second, we show that the dual performance bound may be much better than the primal bnund (and vice versa) for particular problems. The results of Theorems 3.1 and 4.1 are given in terms of a robustness set $U$ which is a subset of a maximal set $\tilde{U}$. Since $\tilde{u}$ is defined implicitly, we provide explicit characterizations of subsets $U$ in Section 5. Here we restate the principal results of $[2-4,10]$ which, in our context, correspond to particular characterizations of subsets of $\tilde{u}$. We also introduce an additional subset of $\tilde{u}$ which provides a new robust stability result. Finally, in Section 6 we consider a pair of illustrative examples. The first example, which was previously considered in [10], involves two uncertain parameters. It is shown that the new guaranteed stability region is considerably larger for certain parameter values than the regions given in [10] (see Figure 1). Furthermore, we obtain a robust performance bound, a result which has no counterpart in [10]. The second example involves controllers for a second-order openloop unstable plant originally considered in [19] to demonstrate the lack of a guaranteed stability margin for LQG controllers. We apply Theorems 3.1 and 4.1 to analyze both the LQG design and
a robustified design obtained in [20]. We show that the new robust stability test is effective in the sense that the guaranteed gain margin for the robustified controller is a factor of 5 larger than the actual gain margin of the LQG design.

## Notation

Note: All matrices have real entries

| $\mathbb{R}, \mathbb{R}^{r \times}, \mathbb{R}^{r}, \mathbb{E}$ | real numbers, $r \times s$ real matrices, $\mathbb{R}^{r \times 1}$, expectation |
| :--- | :--- |
| $I_{r}$ | $r \times r$ identity matrix |
| asymptotically stable matrix | matrix with eigenvalues in open left half plane |
| $\mathbb{S}^{r}, \mathbb{N}^{r}, \mathbb{P}^{r}$ | $r \times r$ symmetric, nonnegative-definite, positive-definite matrices |
| $Z_{1} \geq Z_{2}, Z_{1}>Z_{2}$ | $Z_{1}-Z_{2} \in \mathbb{N}^{r}, Z_{1}-Z_{2} \in \mathbb{P}^{r}, \quad Z_{1}, Z_{2} \in \mathbb{S}^{r}$ |
| $\operatorname{tr} Z, Z^{r}$, co | trace of $Z$, tranpose of $Z$, convex hull |
| $\lambda_{\text {min }}(Z), \lambda_{\text {max }}(Z)$ | smallest and largest eigenvalues of $Z \in \mathbb{S}^{r}$ |
| $\\|Z\\|_{0}$ | spectral norm |
| $Z_{(i, j)}$ | $(i, j)$ element of matrix $Z \quad$ |
| $Z \geq \geq 0$ | $Z_{(i, j)} \geq 0, i, j=1, \ldots, r, \quad Z \in \mathbb{R}^{r \times r}$ |
| $Z \gg 0$ | $Z_{(i, j)}>0, i, j=1, \ldots, r, \quad Z \in \mathbb{R}^{r \times r}$ |
| $\|Z\|_{m}$ | $\left\{\left\|Z_{(i, j)}\right\|\right\}_{i, j=1}^{r}, \quad Z \in \mathbb{R}^{r \times r}($ matrix modulus) |

## 2. Robust Stability and Performance Problems

Let $U \subset \mathbb{R}^{n \times n}$ denote a set of perturbations $\Delta A$ of the nominal dynamics matrix $A$. Throughout the paper it is assumed that $A$ is asymptotically stable. We begin by considering the question of whether or not $A+\Delta A$ is asymptotically stable for all $\Delta A \in U$.

Robust Stability Problem. Determine whether the linear system

$$
\begin{equation*}
\dot{x}(t)=(A+\Delta A) x(t), \quad t \in[0, \infty) \tag{2.1}
\end{equation*}
$$

is asymptotically stable for all $\Delta A \in U$.
The problem of robust performance involves a quadratic form $x^{T}(t) R x(t)$, where $R \in \mathbb{N}^{n}$, when the system is subjected to a white noise disturbance $w(t)$ with nonnegative-definite intensity $V$. The matrix $R$ can be viewed as a means for selecting output variables of interest while the matrix $V$ can be used to specify disturbance levels.

Robust Performance Problem. For the disturbed linear system

$$
\begin{equation*}
\dot{x}(t)=(A+\Delta A) x(t)+w(t), \quad t \in[0, \infty) \tag{2.2}
\end{equation*}
$$

determine a performance bound $\boldsymbol{\beta}$ satisfying

$$
\begin{equation*}
J(U) \triangleq \sup _{\Delta A \in U} \limsup _{t \rightarrow \infty} \mathbb{E}\left[x^{T}(t) R x(t)\right] \leq \beta \tag{2.3}
\end{equation*}
$$

The system (2.2) may, for example, denote a control system in closed-loop configuration subjected to external white noise disturbances (see Section 6). Such specializations are not required for this development, however. Note that $J(U)$ represents the worst case (over $U$ ) of the average (over the white noise statistics) of quadratically weighted steady-state deviations of the state from the origin. Thus $\beta$ represents an upper bound on selected output variances.

Of course, since $R$ and $V$ are only assumed to be nonnegative definite, there may be cases in which a finite performance bound $\beta$ satisfying (2.3) exists while (2.1) is not asymptotically stable over U. In practice, however, robust performance is mainly of interest when (2.1) is robustly stable. In this case the performance $J(u)$ is given in terms of the steady-state second moment of the state. The following result from linear system theory will be useful.

Lemma 2.1. Suppose (2.1) is asymptotically stable for all $\Delta A \in U$. Then

$$
\begin{equation*}
J(u)=\sup _{\Delta A \in u} \operatorname{tr} Q_{\Delta A} R \tag{2.4}
\end{equation*}
$$

where $n \times n Q_{\Delta A} \triangleq \lim _{t \rightarrow \infty} \mathbb{E}\left[x(t) x^{T}(t)\right]$ is the unique, nonnegative-definite solution to

$$
\begin{equation*}
0=(A+\Delta A) Q_{\Delta A}+Q_{\Delta A}(A+\Delta A)^{T}+V \tag{2.5}
\end{equation*}
$$

In the present paper our approach is to obtain sufficient conditions for robust stability as a consequence of sufficient conditions for robust performance. Such conditions are developed in the following sections.

## 3. Sufficient Conditions for Robust Stability and Performance

The key step in obtaining robust stability and performance is to replace the uncertain terms in the Lyapunov equation (2.5) by a bounding matrix $\Omega$. The nonnegative-definite solution $Q$ of this bounding Lyapunov equation is then guaranteed to be an upper bound for $Q_{\Delta A}$. The uncertainty
set $U$ over which robustness is guaranteed then depends upon $Q$. The following easily proved result is fundamental and forms the basis for all later developments. The hypotheses of this result are of a general nature and are not intended to be directly verifiable. Suitably verifiable specializations of the hypotheses are discussed in Section 5.

Theorem 3.1. Let $\Omega \in \mathbb{N}^{n}$, let $Q \in \mathbb{N}^{n}$ be the unique solution to

$$
\begin{equation*}
0=A Q+Q A^{T}+\Omega+V \tag{3.1}
\end{equation*}
$$

and let $U$ be a subset of

$$
\begin{equation*}
\tilde{u} \triangleq\left\{\Delta A \in \mathbb{R}^{n \times n}: \Delta A Q+Q \Delta A^{T} \leq \Omega\right\} . \tag{3.2}
\end{equation*}
$$

Then

$$
\begin{equation*}
\left(A+\Delta A,\left[V+\Omega-\left(\Delta A Q+Q \Delta A^{T}\right)\right]^{\frac{1}{2}}\right) \text { is stabilizable, } \quad \Delta A \in U \tag{3.3}
\end{equation*}
$$

if and only if

$$
\begin{equation*}
A+\Delta A \text { is asymptotically stable, } \quad \Delta A \in U . \tag{3.4}
\end{equation*}
$$

In this case,

$$
\begin{equation*}
Q_{\Delta A} \leq Q, \quad \Delta A \in U \tag{3.5}
\end{equation*}
$$

where $Q_{\Delta A} \in \mathbb{N}^{n}$ is given by (2.5), and

$$
\begin{equation*}
J(u) \leq \operatorname{tr} Q R . \tag{3.6}
\end{equation*}
$$

If, in addition, there exists $\Delta A \in \tilde{u}$ such that $\left(A+\Delta A,\left[V+\Omega-\left(\Delta A Q+Q \Delta A^{T}\right)\right]^{\frac{1}{2}}\right)$ is controllable, then $Q$ is positive definite.

Proof. This result is a minor variation of Theorem 3.1 of [21] and hence the proof is omitted.

To apply Theorem 3.1, one first chooses a nonnegative-definite matrix $\Omega$ and then solves (3.1) for $Q$. Next, as shown in Section 4, one examines $\tilde{U}$ to determine subsets $U$ of perturbations $\Delta A$ over which robustness is guaranteed. Note that if $U_{1}$ and $U_{2}$ are subsets of $\tilde{U}$ then so is the convex hull of their union. (To see this note that $\tilde{U}$ is convex.) The set $\tilde{U}$ is the largest set over which robustness can be guaranteed by Theorem 3.1 for the particular choice of $\Omega$. One may also select several matrices $\Omega$ and determine subsets of each resulting $\tilde{u}$ as a constructive approach to
determining larger robustness sets. In the next section we examine subsets $U$ of $\tilde{u}$ of specified structure. Before doing so, we have the following observations.

In applying Theorem 3.1 it may be convenient to replace condition (3.3) with stronger conditions which are easier to verify in practice. The following result is immediate.

Proposition 3.1. Consider the conditions

$$
\begin{gather*}
V>0  \tag{3.7}\\
\left(A+\Delta A, V^{\frac{1}{2}}\right) \text { is stabilizable, } \quad \Delta A \in U  \tag{3.8}\\
\Delta A Q+Q \Delta A^{T}<\Omega, \quad \Delta A \in U  \tag{3.9}\\
\Delta A Q+Q \Delta A^{T}<\Omega+V, \quad \Delta A \in U \tag{3.10}
\end{gather*}
$$

Then $(3.7) \Rightarrow(3.8) \Rightarrow(3.3),(3.7) \Rightarrow(3.10) \Rightarrow(3.3)$, and $(3.9) \Rightarrow(3.10) \Rightarrow(3.3)$.
If only robust stability is of interest, then the noise intensity $V$ need not have physical significance. In this cse one may either set $V=\varepsilon I_{n}$, where $\varepsilon>0$ is small to satisfy (3.7), or set $V=0$ and confine $U$ to perturbations $\Delta A$ for which (3.9) holds. This is the case in $[3,4,10]$ where $V=0, \Omega=2 I_{n}$, and the parametric robustness sets are characterized by strict inequality.

Remark 3.1. Since $A$ is asymptotically stable, $Q$ is given by

$$
\begin{equation*}
Q=\int_{0}^{\infty} e^{A t}(\Omega+V) e^{A^{T} t} d t=\int_{0}^{\infty} e^{A t} \Omega e^{A^{T} t} d t+Q_{0} \tag{3.11}
\end{equation*}
$$

where $Q_{0} \in \mathbb{N}^{n}$ is given by

$$
\begin{equation*}
0=A Q_{0}+Q_{0} A^{T}+V \tag{3.12}
\end{equation*}
$$

Note that $Q_{0} \leq Q$ and that the nominal performance is given by $\operatorname{tr} Q_{0} R$.
Remark 3.2. Using (3.11) it is also useful to note that the bound for $J(u)$ given by (3.6) can be written as

$$
\begin{equation*}
\operatorname{tr} Q R=\operatorname{tr} \int_{0}^{\infty} e^{A t}(\Omega+V) e^{A^{T_{t}}} d t R=\operatorname{tr} P_{0}(\Omega+V) \tag{3.13}
\end{equation*}
$$

where $P_{0} \in \mathbb{N}^{n}$ is given by

$$
\begin{equation*}
0=A^{T} P_{0}+P_{0} A+R . \tag{3.14}
\end{equation*}
$$

The bound $\operatorname{tr} P_{0}(\Omega+V)$ can be viewed as a dual formulation of the bound $\operatorname{tr} Q R$ since the roles of $A$ and $A^{T}$ are reversed. Dual bounds are developed in the following section. Note that $\operatorname{tr} Q_{0} R=\operatorname{tr} P_{0} V$.

## 4. Dual Sufficient Conditions for Robust Stability and Performance

As noted in Remark 3.2, the performance bound tr $Q R$ given by (3.6) can be expressed equivalently in terms of a dual variable $P_{0}$ for which the roles of $A$ and $A^{T}$ are reversed. Using a similar technique, additional conditions for robust stability and performance can be obtained by developing a dual version of Theorem 3.1. A prime motivation for developing such dual bounds is to draw direct connections with previous results in the literature relating to robust stability. Traditionally, the use of the quadratic Lyapunov function $V(x)=x^{T} P x$ for robust stability leads naturally to the dual formulation. In addition, the dual bounds may, for certain problems, be much sharper than the bounds introduced in the previous section. This point is illustrated at the end of this section by examining an extreme case and in Section 6 by means of numerical examples. We note, in addition, that robust performance bounds are more difficult to motivate within the dual formulation without first developing the primal results. The following result is immediate.

Lemma 4.1. Suppose (2.1) is asymptotically stable for all $\Delta A \in U$. Then

$$
\begin{equation*}
J(u)=\sup _{\Delta A \in U} \operatorname{tr} P_{\Delta A} V, \tag{4.1}
\end{equation*}
$$

where $n \times n P_{\Delta A}$ is the unique, nonnegative-definite solution to

$$
\begin{equation*}
0=(A+\Delta A)^{T} P_{\Delta A}+P_{\Delta A}(A+\Delta A)+R . \tag{4.2}
\end{equation*}
$$

The dual of Theorem 3.1 can now be stated.
Theorem 4.1. Let $\Lambda \in \mathbb{N}^{n}$, let $P \in \mathbb{N}^{n}$ be the unique solution to

$$
\begin{equation*}
0=A^{T} P+P A+\Lambda+R, \tag{4.3}
\end{equation*}
$$

and let $U$ be a subset of

$$
\begin{equation*}
\tilde{u}^{\prime} \triangleq\left\{\Delta A \in \mathbb{R}^{n \times n}: \Delta A^{T} P+P \Delta A \leq \Lambda\right\} \tag{4.4}
\end{equation*}
$$

Then

$$
\begin{equation*}
\left(\left[R+\Lambda-\left(\Delta A^{T} P+P \Delta A\right)\right]^{\frac{1}{2}}, A+\Delta A\right) \text { is detectable, } \quad \Delta A \in U \tag{4.5}
\end{equation*}
$$

if and only if

$$
\begin{equation*}
A+\Delta A \text { is asymptotically stable, } \quad \Delta A \in U \text {. } \tag{4.6}
\end{equation*}
$$

In this case,

$$
\begin{equation*}
P_{\Delta A} \leq P, \quad \Delta A \in U \tag{4.7}
\end{equation*}
$$

where $P_{\Delta A} \in \mathbb{N}^{n}$ is given by (4.2), and

$$
\begin{equation*}
J(u) \leq \operatorname{tr} P V . \tag{4.8}
\end{equation*}
$$

If, in addition, there exists $\Delta A \in U$ such that $\left(\left[R+\Lambda-\left(\Delta A^{T} P+P \Delta A\right)\right]^{\frac{1}{2}}, A+\Delta A\right)$ is observable, then $P$ is positive definite.

The usefulness of Theorem 4.1 resides in the fact that it provides stability and performance bounds which are generally different from those given by Theorem 3.1. Hence, depending upon $\Omega$ and $\Lambda$ either bound (3.6) or bound (4.8) may be better for a particular problem. To illustrate how dual bounds can improve estimates of robust performance, consider the case in which $V=0$, i.e., plant disturbances are absent. In this case $Q_{\Delta A}=0$ satisfies (2.5) and thus $J(u)=0$ so long as $A+\Delta A$ is stable for all $\Delta A \in U$. The performance bound $\operatorname{tr} Q R$ given by (3.6) may, however, be arbitrarily large depending upon $R$ since $Q$ may be nonzero due to $\Omega$. Hence this performance bound may be arbitrarily conservative. The dual bound (4.8), on the other hand, is zero in this case, which completely eliminates the conservatism.

## 5. Characterization of Subsets of $\tilde{u}$ and $\tilde{u}^{\prime}$

To apply Theorems 3.1 and 4.1 it is necessary to explicitly characterize subsets $U$ of $\tilde{U}$ and $\tilde{U}^{\prime}$ over which robustness is guaranteed. In this section we provide several such characterizations by collecting together and extending known results from the literature.

For the following result let $\Omega=\omega I_{n}$, where $\omega>0$, let $W \in \mathbb{R}^{n \times n}, W \gg 0$, and let $A_{1}, \ldots, A_{p} \in \mathbb{R}^{n \times n}$ be arbitrary. Furthermore, for $Q \in \mathbb{P}^{n}$ satisfying (3.1) define for $i=1, \ldots, p$ :

$$
\begin{array}{rlrl}
\alpha_{i} \triangleq \lambda_{\min } & \left(A_{i} Q+Q A_{i}^{T}\right), & \beta_{i} \triangleq \lambda_{\max }\left(A_{i} Q+Q A_{i}^{T}\right), \\
I_{i} \triangleq(-\infty, \infty), & & \alpha_{i}=\beta_{i}=0, \\
\triangleq\left(-\infty, \omega / \beta_{i}\right), & & \alpha_{i} \geq 0, \beta_{i}>0, \\
\triangleq\left(\omega / \alpha_{i}, \infty\right), & \alpha_{i}<0, \beta_{i} \leq 0 \\
\triangleq\left(\omega / \alpha_{i}, \omega / \beta_{i}\right), & & \alpha_{i}<0<\beta_{i}
\end{array}
$$

Finally, let $e_{i}^{(p)}$ denote the $i$ th column of the $p \times p$ identity matrix.

Proposition 5.1. Let $Q \in \mathbb{P}^{n}$ satisfy (3.1) with $\Omega=\omega I_{n}$, where $\omega>0$. Then the following sets are subsets of $\tilde{u}$ which also satisfy (3.9):

$$
\begin{aligned}
& u_{1} \triangleq\left\{\Delta A \in \mathbb{R}^{n \times n}:\|\Delta A\|_{0}<\frac{\omega}{2}\|Q\|_{0}^{-1}\right\} \\
& u_{2} \triangleq\left\{\Delta A \in \mathbb{R}^{n \times n}:|\Delta A|_{m} \ll \omega\left\|\left.\left.W\right|_{Q}\right|_{m}+|Q|_{m} W^{T}\right\|_{0}^{-1} W\right\}, \\
& u_{3} \triangleq\left\{\Delta A \in \mathbb{R}^{n \times n}: \Delta A=\sum_{i=1}^{p} \sigma_{i} A_{i}, \quad\left(\sigma_{1}, \ldots, \sigma_{p}\right)^{T} \in R\right\}
\end{aligned}
$$

where $R$ is one of the following regions in $\mathbb{R}^{p}$ :

$$
\begin{aligned}
& R_{1} \triangleq\left\{\left(\sigma_{1}, \ldots, \sigma_{p}\right): \sum_{i=1}^{p}\left|\sigma_{i}\right|\left\|A_{i} Q+Q A_{i}^{T}\right\|_{2}<\omega\right\}, \\
& R_{2} \triangleq\left\{\left(\sigma_{1}, \ldots, \sigma_{p}\right): \sum_{i=1}^{p} \sigma_{i}^{2}<\omega^{2}\left\|\sum_{i=1}^{p}\left(A_{i} Q+Q A_{i}^{T}\right)^{2}\right\|_{-}^{-1}\right\}, \\
& R_{3} \triangleq\left\{\left(\sigma_{1}, \ldots, \sigma_{p}\right):\left|\sigma_{i}\right|<\omega\left\|\sum_{i=1}^{p}\left|A_{i} Q+Q A_{i}^{T}\right|_{m}\right\|_{0}^{-1}, \quad i=1, \ldots, p\right\}, \\
& R_{4} \triangleq \operatorname{co}\left\{\sigma_{i} e_{i}^{(p)}: \sigma_{i} \in I_{i}, \quad i=1, \ldots, p\right\} .
\end{aligned}
$$

For the dual case we set $\Lambda=\lambda I_{n}$, where $\lambda>0$, and define the dual sets $\tilde{u}_{1}^{\prime}, \tilde{u}_{2}^{\prime}, \tilde{u}_{3}^{\prime}, I_{i}^{\prime}, R_{1}^{\prime}$, $R_{2}^{\prime}, R_{3}^{\prime}$, and $R_{4}^{\prime}$ in an analogous fashion.

Remark 5.1. The proof of Proposition 5.1 is omitted since the results are either known or are immediate. Specifically, $U_{1}^{\prime}$ can be found in [2] while $U_{2}^{\prime}$ appears in [3,4]. The sets $R_{1}^{\prime}, R_{2}^{\prime}$, and $R_{3}^{\prime}$ are given in [10]. The set $R_{4}^{\prime}$ has not appeared previously in the literature although the result is immediate. It is only necessary to diagonalize $A_{i}^{T} P+P A_{i}$ by means of an orthogonal transformation and compare diagonal elements to obtain $I_{i}^{\prime}$. Taking the convex hull over the intervals $I_{i}^{\prime}$ thus yields $R_{4}^{\prime}$. Of course, the required eigenproblem entails additional computation.

Remark 5.2. Although most of the dual of Proposition 5.1 has appeared previously, the primal result Proposition 5.1 has not been discussed in the literature. For robust stability this result can be obtained by considering the stability of $A^{T}$ in place of $A$. As will be shown in Section 6, the primal and dual results lead in general to different robust stability regions and performance bounds. It should also be stressed that although most of the dual of Proposition 5.1 has appeared previously, the present paper extends its applicability to the problem of robust performance in addition to robust stability.

Remark 5.3. As mentioned previously, the convex hull of the union of any collection of subsets of $\tilde{u}$ is also a subset of $\tilde{u}$ since $\tilde{u}$ is convex. This observation applies to $u_{3}$ in the sense that if $u_{3}$ is
a subset of $\tilde{U}$ with regions $R=\hat{R}$ and $R=\hat{R}$ separately, then $U_{3}$ is also a subset with $R$ equal to the convex hull of the union of $\hat{R}$ and $\hat{\hat{R}}$. Note that these observations follow from the convexity of $\tilde{u}$ and do not contradict the fact that the set of asymptotically stable matrices is not convex.

Remark 5.4. The requirement that $\Omega$ be of the form $\omega I_{n}$ is not a constraint in applying Proposition 5.1. Indeed, it is only required that $\Omega$ be positive definite. To see this let invertible $\phi \in$ $\mathbb{R}^{n \times n}$ be such that $\phi \Omega \phi^{T}=I_{n}$. Then Proposition 5.1 can be applied with suitable transformations of $\Delta A, Q, W$, and $A_{i}$.

Remark 5.5. As in [2-4,10], the sets $U_{1}, U_{2}, R_{1}, R_{2}$, and $R_{3}$ are defined in terms of strict inequalities. In this case $\tilde{U}_{1}, U_{2}$, and $U_{3}$ consist of elements of $\tilde{u}$ satisfying $\Delta A Q+Q \Delta A^{T}<\Omega$ so that (3.9) is satisfied. Thus, by Proposition 3.1, the stabilizability condition (3.3) is automatically satisfied without reference to $V$.

Remark 5.6. In the special case $p=1$ it is clear that $R_{1}=R_{2}$. Furthermore, in this case $R_{3}$ is always a subset of $R_{1}$ and $R_{2}$ and hence leads to a more conservative stability region. The largest possible set of perturbations $\Delta A$ of the form $\sigma_{i} A_{i}$ contained in $\tilde{U}$ is given by $R_{4}$.

Remark 5.7. It is shown in Remark 2.12 of [10] that $U_{2}$ can be obtained as a consequence of $U_{3}$ with $R=R_{3}$ and a suitable choic\% of $A_{i}$. Hence $U_{2}$ need not actually be considered separately. Our assumption that $W \gg 0$ (and not $W \geq \geq 0$ ) is for convenience only.

Remark 5.8. Note that all of the subsets of $\tilde{U}$ given by Proposition 5.1 are symmetric except for $U_{3}$ with $R=R_{4}$. When the actual stability region is highly asymmetric, it follows that a symmetric robust stability region is necessarily highly conservative. This observation is illustrated by an example in Section 6.

Remark 5.9. The regions given by $R_{1}, R_{2}$, and $R_{3}$ correspond, respectively, to 1 -norm, 2 norm, and $\infty$-norm neighborhoods. These results can easily be extended to include more general regions. For example, in the definition of $R_{2}$ replace $\sigma_{i}$ by $\sigma_{i} / a_{i}$ and $A_{i} Q+Q A_{i}^{T}$ by $a_{i}\left(A_{i} Q+Q A_{i}^{T}\right)$, where $a_{i}$ is an arbitrary positive constant, $i=1, \ldots, p$. With this modification $R_{2}$ corresponds to an elliptical robust stability region. Detailed investigation of such regions is beyond the scope of this paper.

Remark 5.10. When each interval $I_{i}$ is finite, or when only a finite interval is of interest, $R_{4}$ can be expressed as the convex hull of a finite number of points. Specifically, letting $I_{i}=$
$\left[a_{i}, b_{i}\right], i=1, \ldots, p$, it follows that

$$
R_{4}=\operatorname{co}\left\{a_{1} e_{2}^{(p)}, b_{1} e_{1}^{(p)}, \ldots, a_{p} e_{p}^{(p)}, b_{p} e_{p}^{(p)}\right\}
$$

This set is illustrated by means of an example in the next section.

## 6. Examples

As a first example we adopt Example 2 of [10]. This example, which involves two uncertain parameters, was used in [10] to illustrate the robust stability regions $R_{1}^{\prime}, R_{2}^{\prime}$, and $R_{3}^{\prime}$. The problem was originally cast in the form of a static output feedback controller with uncertain gains. Here for convenience in discussing robust performance we reformulate the example to involve uncertainty in the control input matrix. Hence consider the control system

$$
\begin{align*}
& \dot{x}(t)=A_{0} x(t)+B_{0} u(t)  \tag{6.1}\\
& y(t)=C_{0} x(t)  \tag{6.2}\\
& u(t)=K y(t) \tag{6.3}
\end{align*}
$$

where

$$
A_{0}=\left[\begin{array}{rrr}
-1 & 0 & 0 \\
0 & -2 & 0 \\
0 & 0 & -3
\end{array}\right], \quad B_{0}=\left[\begin{array}{ll}
1 & 0 \\
0 & 1 \\
1 & 1
\end{array}\right], \quad C_{0}=\left[\begin{array}{lll}
1 & 0 & 1 \\
0 & 1 & 0
\end{array}\right], \quad K=\left[\begin{array}{cc}
-1 & 0 \\
0 & -1
\end{array}\right],
$$

and the uncertainty $\Delta B_{0}$ in $B_{0}$ is given by

$$
\Delta B_{0}=\left[\begin{array}{cc}
-\sigma_{1} & 0 \\
0 & -\sigma_{2} \\
-\sigma_{1} & -\sigma_{2}
\end{array}\right] .
$$

The closed-loop dynamics matrix is then given by

$$
A+\Delta A=\left[\begin{array}{ccc}
-2+\sigma_{1} & 0 & -1+\sigma_{1} \\
0 & -3+\sigma_{2} & 0 \\
-1+\sigma_{1} & -1+\sigma_{2} & -4+\sigma_{1}
\end{array}\right]
$$

where $\Delta A=\sigma_{1} A_{1}+\sigma_{2} A_{2}$ and $A_{1}, A_{2}$ have the evident definitions. It can easily be shown that the exact stability region is given by $\sigma_{1} \in(-\infty, 1.75)$ and $\sigma_{2} \in(-\infty, 3)$. Thus the nominal dynamics matrix corresponding to $\sigma_{1}=\sigma_{2}=0$ lies in the upper right-hand corner of the exact stability region so that, as noted in Remark 5.8, a high degree of conservatism can be expected using symmetric robustness regions. To consider robust stability alone, set $V=R=0$ and $\omega=\lambda=2$. In this case
regions $R_{1}^{\prime}, R_{2}^{\prime}$, and $R_{3}^{\prime}$, as computed in [10], are shown in Figure 1. Region $R_{4}^{\prime}$ for this problem is given (see Remark 5.10) by

$$
R_{4}^{\prime}=c o\left\{\binom{-29.6}{0},\binom{1.65}{0},\binom{0}{-20.5},\binom{0}{2.85}\right\},
$$

which accounts somewhat better for the asymmetry of the stability region. The regions $R_{1}, R_{2}$, and $R_{3}$ were found to be smaller than the corresponding dual regions, while $R_{4}$ is given by

$$
R_{4}=\operatorname{co}\left\{\binom{-31.1}{0},\binom{1.64}{0},\binom{0}{-10.4},\binom{0}{2.63}\right\},
$$

which yields slight improvement in $\sigma_{1}$.
To evaluate robust performance replace (6.1) by

$$
\begin{equation*}
\dot{x}(t)=A_{0} x(t)+B_{0} u(t)+w(t), \tag{6.4}
\end{equation*}
$$

and define

$$
J=\lim _{t \rightarrow \infty} \mathbb{E}\left[x^{T}(t) R_{1} x(t)+u^{T}(t) R_{2} u(t)\right]
$$

which corresponds to (2.3) with $R=R_{1}+C_{0}^{T} K^{T} R_{2} K C_{0}$. Hence setting $R_{1}=I_{3}$ and $R_{2}=I_{2}$ yields

$$
R=\left[\begin{array}{lll}
2 & 0 & 1 \\
0 & 2 & 0 \\
1 & 0 & 2
\end{array}\right]
$$

We also set $V=I_{3}$ and $\omega=\lambda=2$. The resulting stability region for these values of $V$ and $R$ is given by

$$
\begin{aligned}
& R_{1}^{\prime}=\left\{\left(\sigma_{1}, \sigma_{2}\right):\left|\sigma_{1}\right| / .70+\left|\sigma_{2}\right| / 1.46<1\right\} \\
& R_{2}^{\prime}=\left\{\left(\sigma_{1}, \sigma_{2}\right): \sigma_{1}^{2}+\sigma_{2}^{2}<(.70)^{2}\right\} \\
& R_{3}^{\prime}=\left\{\left(\sigma_{1}, \sigma_{2}\right):\left|\sigma_{1}\right|<.68, i=1,2\right\}, \\
& R_{4}^{\prime}=c u\left\{\binom{-20.5}{0},\binom{.70}{0},\binom{0}{-13.7},\binom{0}{1.46}\right\} .
\end{aligned}
$$

Over these combined regions the performance bound was computed to be $\operatorname{tr} P V=2.26$. The primal result produced the regions

$$
\begin{aligned}
& R_{1}=\left\{\left(\sigma_{1}, \sigma_{2}\right):\left|\sigma_{1}\right| / 1.09+\left|\sigma_{2}\right| / 1.75<1\right\} \\
& R_{2}=\left\{\left(\sigma_{1}, \sigma_{2}\right): \sigma_{1}^{2}+\sigma_{2}^{2}<(1.08)^{2}\right\} \\
& R_{3}=\left\{\left(\sigma_{1}, \sigma_{2}\right):\left|\sigma_{1}\right|<1.0, \quad i=1,2\right\} \\
& R_{4}=c o\left\{\binom{-20.8}{0},\binom{1.09}{0},\binom{0}{-6.93},\binom{0}{1.75}\right\} .
\end{aligned}
$$

Over these regions the performance bound was computed to be $\operatorname{tr} Q R=3.18$. Contour plots of actual performance for perturbed values of $\sigma_{1}$ and $\sigma_{2}$ are shown in Figure 2. Note that when determining robust performance Theorems 3.1 and 4.1 yield performance bounds over robust stability regions which are generally smaller than the robust stability regions determined with $R=0$ and $V=0$. This mechanism represents the natural tradeoff between stability and performance. In general, to determine the largest stability regions, $V$ and $R$ should be set to zero initially.

As a second example we consider the control system given in [19] to demonstrate the lack of guaranteed gain margin for LQG controllers. Hence consider

$$
\begin{align*}
\dot{x}_{0}(t) & =A_{0} x_{0}(t)+B_{0} u(t)+w_{1}(t)  \tag{6.5}\\
y(t) & =C_{0} x_{0}(t)+w_{2}(t) \tag{6.6}
\end{align*}
$$

with controller

$$
\begin{align*}
\dot{x}_{c}(t) & =A_{o} x_{c}(t)+B_{c} y(t)  \tag{6.7}\\
u(t) & =C_{o} x_{c}(t) \tag{6.8}
\end{align*}
$$

and performance

$$
\begin{equation*}
J=\lim _{t \rightarrow \infty} \mathbb{E}\left[x_{0}^{T}(t) R_{1} x_{0}(t)+u^{T}(t) R_{2} u(t)\right] \tag{6.9}
\end{equation*}
$$

The data are

$$
\begin{gathered}
A_{0}=\left[\begin{array}{ll}
1 & 1 \\
0 & 1
\end{array}\right], \quad B_{0}=\left[\begin{array}{l}
0 \\
1
\end{array}\right], \quad C_{0}=\left[\begin{array}{ll}
1 & 0
\end{array}\right], \\
V_{1}=R_{1}=\rho\left[\begin{array}{ll}
1 & 1 \\
1 & 1
\end{array}\right], \quad V_{2}=R_{2}=1
\end{gathered}
$$

where $V_{1}$ and $V_{2}$ are the intensities of $w_{1}(t)$ and $w_{2}(t)$, respectively. Uncertainty $\Delta B_{0}$ in $B_{0}$ is thus represented by $\sigma_{1} B_{1}$, where $B_{1}=\left[\begin{array}{lll}0 & 1\end{array}\right]^{T}$. Thus, the closed-loop system corresponds to

$$
\begin{gathered}
A=\left[\begin{array}{cc}
A_{0} & B_{0} C_{c} \\
B_{c} C_{0} & A_{c}
\end{array}\right], \quad A_{1}=\left[\begin{array}{cc}
0 & B_{1} C_{c} \\
0 & 0
\end{array}\right], \\
R=\left[\begin{array}{cc}
R_{1} & 0 \\
0 & 0
\end{array}\right], \quad V=\left[\begin{array}{cc}
V_{1} & 0 \\
0 & B_{c} V_{2} B_{c}^{T}
\end{array}\right],
\end{gathered}
$$

where the zero in the $(2,2)$ block of $R$ denotes the fact that we are considering the robust performance bound for the state regulation cost only. Choosing $\rho=60$, it follows that the LQG gains are given by

$$
A_{\mathrm{c}}=\left[\begin{array}{cc}
-9 & 1 \\
-20 & -9
\end{array}\right], \quad B_{\mathrm{c}}=\left[\begin{array}{l}
10 \\
10
\end{array}\right], \quad C_{c}=\left[\begin{array}{ll}
-10 & -10
\end{array}\right] .
$$

For this controller the actual stability region corresponds to $\sigma_{1} \in(-.07, .01)$ (see Figure 3). Applying the results of Section 5 with $V=R=0$ (for robust stability only) and $\omega=\lambda=2$, we obtain

$$
\begin{array}{ll}
R_{1}=R_{2}=R_{3}=(-.000242, .000242), & R_{4}=(-.000242, .000728) \\
R_{1}^{\prime}=R_{2}^{\prime}=(-.0000247, .0000247), & R_{3}^{\prime}=(-.0000219, .0000219) \\
R_{4}^{\prime}=(-.0000247, .0000265)
\end{array}
$$

Note that although the primal results are better than the dual results by an order of magnitude, they are conservative by two orders of magnitude with respect to the actual gain margin. For robust performance we again set $\omega=\lambda=2$ and, using $R$ and $V$ given above, we obtained the bound $\operatorname{tr} Q R=7633$ over the stability region $R_{4}=(-.000192, .000613)$. The nominal performance was given by $\operatorname{tr} Q_{0} R=\operatorname{tr} P_{0} V=4875$, while the dual performance bound was $\operatorname{tr} P V=10510$ over $R_{4}^{\prime}=(-.0000222, .0000238)$.

Robustified controllers for the example of [19] were obtained in [20] using the approach discussed in [13]. As shown in Figure 3 (see also [20]), the closed-loop system with the controller

$$
A_{\mathrm{c}}=\left[\begin{array}{cc}
-10.69 & 1 \\
-32.97 & -5.295
\end{array}\right], \quad B_{\mathrm{c}}=\left[\begin{array}{l}
11.69 \\
26.67
\end{array}\right], \quad C_{c}=\left[\begin{array}{ll}
-6.245 & -6.245
\end{array}\right],
$$

is stable over the range $\sigma_{1} \in(-.28, .21)$. Hence we wish to determine whether the robust stability tests are capable of detecting this increase in gain margin. Applying Theorems 3.1 and 4.1 with $\omega=\lambda=2$ and $V=R=0$ yields stability for $\sigma_{1}$ in the regions $R_{1}=R_{2}=R_{3}=(-.0115, .0115)$ and $R_{4}=(-.0115, .057)$. This guarantee of stability is two orders of magnitude greater than the guarantee for the LQG design but is still an order of magnitude conservative with respect to the actual stability region for this controller. Note, however, that for $\sigma_{1}>0$ the guaranteed gain margin for the robustified design given by $R_{4}$ (i.e., .057 ) is greater than the actual gain margin of the LQG design (.01). Hence the robustness test given by the $R_{4}$ was able to detect a factor of 5 stability augmentation provided by the robustified design compared to the LQG controller. Finally, the robust performance bound for this controller was computed to be $\operatorname{tr} Q R=11185$ over the region $R_{4}=(-.00165, .00493)$, while the dual bound was found to be $\operatorname{tr} P V=11223$ over $R_{4}^{\prime}=(-.000724, .00123)$. For this problem the nominal performance is $\operatorname{tr} Q_{0} R=\operatorname{tr} P_{0} V=9997$.

Acknowledgment. We wish to thank J. Straehla for wordprocessing support, T. Rhodes for preparing Figure 1, A. Tellez for performing the numerical calculations, and A. Daubendiek for producing Figures 2 and 3.

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## $\sigma_{1}\left(\times 10^{-3}\right)$



乙 əxn6Ts



Figure 1


# Robust Stability and Performance Analysis <br> for <br> State Space Systems <br> via <br> Quadratic Lyapunov Bounds 

by

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#### Abstract

For a given asymptotically stable linear dynamic system it is often of interest to determine whether stability is preserved as the system varies within a specified class of uncertainties. If, in addition, there also exist associated performance measures (such as the steady-state variances of selected state variables), it is desirable to assess the worst-case performance over a class of plant variations. These are problems of robust stability and performance analysis. In the present paper we consider quadratic Lyapunov bounds to obtain a simultaneous treatment of both robust stability and performance. The approach is based upon the construction of modified Lyapunov equations which provide sufficient conditions for robust stability along with robust performance bounds. One f the principal features of the paper is the unified treatment and extension of several quadratic Lyapunov bounds developed previously for feedback control design.


Key Words: robust analysis, stability, performance, Lyapunov equations, structured uncertainty

Supported in part by the Air Force Office of Scientific Research under contract F49620-86-C0002.

## 1. Introduction

Unavoidable discrepancies between mathematical models and real-world systems can result in degradation of control-system performance including instability ([1,2]). Ideally, feedback control systems should be designed to be robust with respect to uncertainties, or perturbations, in the plant characteristics. Such uncertainties may arise either due to limitations in performing system identification prior to control-system implementation or because of unpredictable plant changes which occur during operation. Thus robustness analysis must play a key role in control-system design. That is, given an existing or proposed control system, determine the performance degradation due to variations in the plant.

In performing robustness analysis there are two principal concerns, namely, stability robustness and performance robustness. Stability robustness addresses the qualitative question as to whether or not the system remains stable for all plant perturbations within a specified class of uncertainties. A related problem involves determining the largest class of plant perturbations under which stability is preserved. Once robust stability has been ascertained, it is of interest to determine quantitatively the degradation of performance within a given robust stability range. In practice it is often desirable to determine the worst case performance as a measure of degradation.

The concern for both robust stability and performance can be traced back to the earliest developments in control theory. Design specifications such as gain and phase margin have traditionally been used to gauge system reliability in the face of uncertainty. In the modern control literature considerable effort has focused on rigorous robustness analysis and design techniques in a variety of settings. Analysis and synthesis results have been developed for both state-space and frequencydomain plant models to address structured parameter variations as well as normed-neighborhood uncertainty ([3-7]).

The present paper is concerned solely with the analysis of structured real-valued parameter uncertainty within the context of state space models. Motivation for such problems is clearly illustrated by simple examples given in [1,2]. These examples show that standard linear-quadratic methods used to design either full-state feedback controllers or dynamic compensators may result in closed-loop systems which are arbitrarily sensitive to structured real-valued plant parameter variations. A particularly effective technique for analyzing robust stability is to construct a Lyapunor function which guarantees stability of the system as the uncertain parameters vary over a specified range. Using the quadratic Lyapunov function $V(x)=x^{T} P x$ this technique has been extensively
developed for both analysis and synthesis (see, e.g., [8-37]).
Although both robust stability and performance are of interest in practice, most of the literature involving quadratic Lyapunov functions is confined to the problem of robust stability. A notable exception is the early work of Chang and Peng ([9]) which also provides bounds on worst-case quadratic performance within full-state feedback control design. In the present paper we further extend the approach of [9] to obtain a series of results for analyzing both robust stability and performance. As will be seen, these results also provide substantial unification of more recent results pertaining only to robust stability.

To illustrate the basis for our approach, consider the system

$$
\begin{equation*}
\dot{x}(t)=(A+\Delta A) x(t)+D_{0} w(t), \tag{1.1}
\end{equation*}
$$

where $x(t)$ is an $n$-vector, $A$ is an $n \times n$ matrix denoting the nominal dynamics matrix, $\Delta A$ denotes an uncertain perturbation of $A$ belonging to a specified set $U$, and $D_{0} w(t)$ is (for now) a white noise signal of specified intensity $V \triangleq D_{0} D_{0}^{T}$. System (1.1) may, for example, denote a control system in closed-loop configuration. For the system (1.1) the performance invoives the steady-state covariance of specified outputs $E_{0} x(t)$. In practice the diagonal elements of the output covariance are measures of the ability of the external disturbances $D_{0} w(t)$ to affect specified states. In the presence of uncertainties $\Delta A$, it is of interest to determine the worst case steady-state values of selected state variances. Thus, we define the scalar performance criterion

$$
\begin{equation*}
J_{S}(u) \triangleq \sup _{\Delta A \in u} \limsup _{t \rightarrow \infty} \mathbb{E}\left\{\left[E_{0} x(t)\right]^{\mathrm{T}}\left[E_{0} x(t)\right]\right\} \tag{1.2}
\end{equation*}
$$

where " $E$ " denotes expectation. To evaluate (1.2) define the state covariance

$$
\begin{equation*}
Q(t) \triangleq \mathbb{E}\left[x(t) x^{\mathbb{T}}(t)\right] \tag{1.3}
\end{equation*}
$$

which satisfies the Lyapunov differential equation

$$
\begin{equation*}
\dot{Q}_{\Delta A}(t)=(A+\Delta A) Q_{\Delta A}(t)+Q_{\Delta A}(t)(A+\Delta A)^{T}+V, \tag{1.4}
\end{equation*}
$$

so that (1.2) becomes

$$
\begin{equation*}
J_{\mathbf{S}}(u)=\sup _{\Delta A \in u} \limsup _{t \rightarrow \infty} \operatorname{tr} Q_{\Delta A}(t) R, \tag{1.5}
\end{equation*}
$$

where $R \triangleq E_{0} E_{0}^{T}$. To guarantee both robust stability and performance we consider modified algebraic Lyapunov equations of the form

$$
\begin{equation*}
0=A Q+Q A^{\mathbf{T}}+\Omega(Q)+V \tag{1.6}
\end{equation*}
$$

where $\Omega(\cdot)$ is a matrix operator satisfying

$$
\begin{equation*}
\Delta A Q+Q \Delta A^{T} \leq \Omega(Q) \tag{1.7}
\end{equation*}
$$

for all $\Delta A \in U$ and all nonnegative-definite matrices $Q$. The ordering in (1.7) is defined with respect to the cone of nonnegative-definite matrices. Our results are based on the following robust stability and performance result. If there exists a nonnegative-definite solution $Q$ to (1.6) where $\Omega(\cdot)$ satisfies (1.7), then $A+\Delta A$ is asymptotically stable for all $\Delta A \in U$ and, furthermore,

$$
\begin{equation*}
J_{\mathbf{s}}(U) \leq \operatorname{tr} Q R . \tag{1.8}
\end{equation*}
$$

The performance bound (1.8) follows from the fact that since $A+\Delta A$ is asymptotically stable, $Q_{\Delta A} \triangleq \lim _{t \rightarrow \infty} Q_{\Delta A}(t)$ exists and satisfies

$$
\begin{equation*}
0=(A+\Delta A) Q_{\Delta A}+Q_{\Delta A}(A+\Delta A)^{T}+V \tag{1.9}
\end{equation*}
$$

Now subtracting (1.9) from (1.6) implies

$$
\begin{equation*}
Q_{\Delta A} \leq Q^{\prime} \tag{1.10}
\end{equation*}
$$

which with (1.5) yields (1.8).
Since the ordering induced by the cone of nonnegative-definite matrices is only a partial ordering, it should not be expected that there exists an operator $\Omega(\cdot)$ satisfying (1.7) which is a least upper bound. Indeed, connections between the result outlined above and the approach of [9] as well as more recent work arise from alternative definitions of the operator $\Omega(\cdot)$. To illustrate these connections assume for convenience that $\Delta A$ is of the form

$$
\begin{equation*}
\Delta A=\sigma_{1} A_{1}, \quad\left|\sigma_{1}\right| \leq \delta_{1} \tag{1.11}
\end{equation*}
$$

where $\sigma_{1}$ is an uncertain real scalar parameter assumed only to satisfy the stated bounds, and $A_{1}$ is a known matrix denoting the structure of the parametric uncertainty. The original definition of $\Omega(\cdot)$ in [9] was given by

$$
\begin{equation*}
\Omega(Q)=\delta_{1}\left|A_{1} Q+Q A_{1}^{\mathrm{T}}\right| \tag{1.12}
\end{equation*}
$$

where $|\cdot|$ denotes the nonnegative-definite matrix obtained by replacing each eigenvalue by its absolute value. This bound was studied in $[9,12]$ for full-state feedback design. More recently, the quadratic bound

$$
\begin{equation*}
\Omega(Q)=\delta_{1}[D+Q E Q] \tag{1.13}
\end{equation*}
$$

has been considered, where $D \triangleq D_{1} D_{1}^{\mathrm{T}}, E \triangleq E_{1}^{\mathrm{T}} E_{1}$, and $D_{1}, E_{1}$ are a factorization of $A_{1}$ of the form $A_{1}=D_{1} E_{1}$. Bound (1.13) was studied in [29] for robustness analysis and in [17,25,28,30,33,36] for robust controller synthesis. A third bound which has also been considered is the linear bound

$$
\begin{equation*}
\Omega(Q)=\delta_{1}\left[\alpha Q+\alpha^{-1} A_{1} Q A_{1}^{\mathrm{T}}\right], \tag{1.14}
\end{equation*}
$$

where $\alpha$ is an arbitrary positive scalar. As shown in [33], bound (1.14) arises from a multiplicative white noise model with exponential disturbance weighting. Control-design applications of bound (1.14) are given in $[23,27,33,34,35]$.

The principal contribution of the present paper is thus a unified development of bounds (1.12)(1.14) for both robust stability and performance analysis. In addition, we present a systematic approach which pays careful attention to the structure of the uncertainty set $\mathcal{U}$. For example, we show that bound (1.12) guarantees stability over a rectangular uncertainty set while (1.14) is most naturally associated with an ellipsoidal region. Furthermore, to provide a methodical development, we consider three classes of bounds (Type I, II and III) which operate by exploiting, respectively, the symmetry of $\Delta A Q+Q \Delta A^{T}$, the structure of $Q$, and the structure of $\Delta A$. This approach clarifies the relationships among different bounds and suggests several new bounds.

Finally, the present paper also considers an alternative functional for robust performance analysis. Specifically, in place of white noise disturbances, we reinterpret $w(t)$ in (1.1) as a deterministic $L_{2}$ signal as in $H_{\infty}$ theory ([6]). By imposing an $L_{\infty}$ norm on the output $E_{0} x(t)$, the corresponding performance measure is given by (see [38])

$$
\begin{equation*}
J_{\mathrm{D}}(U)=\sup _{\Delta A \in U} \limsup _{t \rightarrow \infty} \lambda_{\max }\left(Q_{\Delta A}(t) R\right) \tag{1.15}
\end{equation*}
$$

in contrast to (1.5). Both performance measures $J_{S}(u)$ and $J_{D}(u)$ are considered in the paper.
The contents of the paper are as follows. After summarizing notation later in this section, the Robust Stability Problem, Stochastic Robust Performance Problem, and Deterministic Robust Performance Problem are introduced in Section 2. In Section 3 the basic result guaranteeing robust stability and performance (Theorem 3.1) is stated. This result is easily stated and forms the basis for all later developments. A dual version of Theorem 3.1 (Theorem 4.1) provides additional sufficient conditions and clarifies connections to traditional robust stability results. The bound $\Omega(\cdot)$ and its dual $\Lambda(\cdot)$ are given concrete forms in Section 5. In Section 6, the bounds of Section 5 are merged with Theorem 3.1 to yield the main results guaranteeing robust stability and performance (Theorems
6.1-6.5) via modified Lyapunov equations. In Section 7 we analyze the modified Lyapunov equations with regard to existence, uniqueness, and monotonicity of solutions. Additional bounds are derived in Section 8 by utilizing a recursive substitution technique, while both upper and lower bounds are obtained in Section 9. Finally, illustrative numerical examples are considered in Sections 10 and 11.

## Notation

Note: All matrices have real entries

| $\mathbb{R}, \mathbb{R}^{r \times s}, \mathbb{R}^{r}, \mathbb{E}$ | real numbers, $r \times s$ real matrices, $\mathbb{R}^{r \times 1}$, expectation |
| :--- | :--- |
| $I_{r}$ | $r \times r$ identity matrix |
| asymptotically <br> stable matrix | matrix with eigenvalues in open left half plane |
| $\mathbb{S}^{r}$ | $r \times r$ symmetric matrices |
| $\mathbb{N}^{r}$ | $r \times r$ symmetric nonnegative-definite matrices |
| $\mathbb{P}^{r}$ | $r \times r$ symmetric positive-definite matrices |
| $Z_{1} \geq Z_{2}$ | $Z_{1}-Z_{2} \in \mathbb{N}^{r}, Z_{1}, Z_{2} \in \mathbb{S}^{r}$ |
| $Z_{1}>Z_{2}$ | $Z_{1}-Z_{2} \in \mathbb{P}^{r}, Z_{1}, Z_{2} \in \mathbb{S}^{r}$ |
| $\operatorname{tr} Z, Z^{T}$ | trace of $Z$, transpose of $Z$ |
| $\lambda_{i}(Z)$ | eigenvalue of matrix $Z$ |
| $\lambda_{m a x}(Z)$ | maximum eigenvalue of matrix $Z$ having real spectrum |

## 2. Robust Stability and Performance Problems

Let $U \subset \mathbb{R}^{n \times n}$ denote a set of perturbations $\Delta A$ of a given nominal dynamics matrix $A \in$ $\mathbb{R}^{n \times n}$. Throughout the paper it is assumed that $A$ is asymptotically stable and that $0 \in U$. We begin by considering the question of whether or not $A+\Delta A$ is asymptotically stable for all $\Delta A \in U$.

Robust Stability Problem. Determine whether the linear system

$$
\begin{equation*}
\dot{x}(t)=(A+\Delta A) x(t), \quad t \in[0, \infty), \tag{2.1}
\end{equation*}
$$

is asymptotically stable for all $\Delta A \in U$.
To consider the problem of robust performance it is necessary to introduce external disturbances. In this paper we consider both stochastic and deterministic disturbance models. The
stochastic disturbance model involves white noise signals as in standard LQG theory while the deterministic disturbance model involves $L_{2}$ signals as in $H_{\infty}$ theory ([6]). By defining an appropriate performance measure for each disturbance class it turns out that we can provide a simultaneous treatment of both cases.

We first consider the case of stochastic disturbances. In this case the robust performance problem concerns the worst-case magnitude of the expected value of a quadratic form $\left[E_{0} x(t)\right]^{\mathrm{T}}\left[E_{0} x(t)\right]$, where the matrix $E_{0} \in \mathbb{R}^{q \times n}$ defines the output states, when the system is subjected to a standard white noise disturbance $w(t) \in \mathbb{R}^{d}$ with weighting $D_{0} \in \mathbb{R}^{n \times d}$.

Stochastic Robust Performance Problem. For the disturbed linear system

$$
\begin{equation*}
\dot{x}(t)=(A+\Delta A) x(t)+D_{0} w(t), \quad t \in[0, \infty) \tag{2.2}
\end{equation*}
$$

where $w(\cdot)$ is a d-dimensional white noise signal with intensity $I_{d}$, determine a performance bound $\beta_{S}$ satisfying

$$
\begin{equation*}
J_{R S}(U) \triangleq \sup _{\Delta A \in U} \limsup _{t \rightarrow \infty} \mathbb{E}\left\{\left[E_{0} x(t)\right]^{\mathrm{T}}\left[E_{0} x(t)\right]\right\} \leq \beta_{\mathrm{S}} \tag{2.3}
\end{equation*}
$$

The system (2.2) may denote, for example, a control system in closed-loop configuration subjected to external white noise disturbances for which $z(t) \triangleq E_{0} x(t)$ may be the state regulation error. Such specializations are not required for this development, however.

Of course, since $E_{0}$ and $D_{0}$ may be rank deficient, there may be cases in which a finite performance bound $\beta_{\mathrm{S}}$ satisfying (2.3) exists while (2.1) is not asymptotically stable over $U$. In practice, however, robust performance is mainly of interest when (2.1) is robustly stable. In this case the performance $J_{\mathbf{s}}(U)$ is given in terms of the steady-state second moment of the state. The following result from linear system theory will be useful. For convenience define the $n \times n$ nonnegative-definite matrices

$$
\begin{equation*}
R \triangleq E_{0}^{\mathbf{T}} E_{0}, \quad V \triangleq D_{0} D_{0}^{\mathbf{T}} \tag{2.4}
\end{equation*}
$$

Lemma 2.1. Suppose (2.1) is asymptotically stable for all $\Delta A \in U$. Then

$$
\begin{equation*}
J_{S}(u)=\sup _{\Delta A \in U} \operatorname{tr} Q_{\Delta A} R \tag{2.5}
\end{equation*}
$$

where $n \times n Q_{\Delta A} \triangleq \lim _{t \rightarrow \infty} \mathbb{E}\left[x(t) x^{T}(t)\right]$ is given by

$$
\begin{equation*}
Q_{\Delta A}=\int_{0}^{\infty} e^{(A+\Delta A) t} V e^{(A+\Delta A)^{T} t} d t \tag{2.6}
\end{equation*}
$$

which is the unique, nonnegative-definite solution to

$$
\begin{equation*}
0=(A+\Delta A) Q_{\Delta A}+Q_{\Delta A}(A+\Delta A)^{\mathbf{T}}+V \tag{2.7}
\end{equation*}
$$

To state the Deterministic Robust Performance Problem some additional notation is required. For a measurable function $w:[0, \infty) \rightarrow \mathbb{R}^{d}$ define

$$
\begin{equation*}
\| w\left(\cdot \|_{2} \triangleq\left\{\int_{0}^{\infty} w^{\mathbf{T}}(t) w(t) d t\right\}^{\frac{1}{2}}\right. \tag{2.8}
\end{equation*}
$$

Note that definition (2.8) is an $L_{2}$ function norm with a Euclidean spatial norm. We now reconsider (2.2) with w(•) now interpreted as a square-integrable function. In this case the robust performance problem concerns the worst-case $L_{\infty}$ norm of a quadratic form $\left[E_{0} x(t)\right]^{\mathrm{T}}\left[E_{0} x(t)\right]$.

Deterministic Robust Performance Problem. For the disturbed linear system (2.2), where $\|w(\cdot)\|_{2} \leq 1$, determine a performance bound $\beta_{\mathrm{D}}$ satisfying

$$
\begin{equation*}
J_{\mathrm{D}}(u) \triangleq \sup _{\Delta \mathbf{A} \in u} \sup _{\|凶(\cdot)\|_{2} \leq 1} \sup _{t \in[0, \infty)}\left\{\left[E_{0} x(t)\right]^{\mathbf{T}}\left[E_{0} x(t)\right]\right\} \leq \beta_{\mathrm{D}} \tag{2.9}
\end{equation*}
$$

The following result is proved in [38].
Lemma 2.2. Suppose (2.1) is asymptotically stable for all $\Delta A \in U$. Then

$$
\begin{equation*}
J_{D}(u)=\sup _{\Delta A \in U} \lambda_{\max }\left(Q_{\Delta A} R\right) \tag{2.10}
\end{equation*}
$$

where $Q_{\Delta A}$ is the unique, nonnegative-definite solution to (2.7).
Remark 2.1. Although $J_{\mathbf{S}}(U)$ and $J_{\mathrm{D}}(U)$ arise from different mathematical settings they are quite similar in form. Note that in general $J_{\mathrm{D}}(U) \leq J_{\mathbf{s}}(U)$, and $J_{\mathrm{D}}(U)=J_{\mathrm{s}}(U)$ if rank $R=1$.

Remark 2.2. In Lemma $2.2 Q_{\Delta A}$ can be viewed as the controllability Gramian for the pair ( $A+\Delta A, D_{0}$ ) rather than the state covariance.

In the present paper our approach is to obtain robust stability as a consequence of sufficient conditions for robust performance. Such conditions are developed in the following sections.

## 3. Sufficient Conditions for Robust Stability and Performance

The key step in obtaining robust stability and performance is to bound the uncertain terms $\Delta A Q+Q \Delta A^{T}$ in the Lyapunov equation (2.7) by means of a function $\Omega(Q)$. The nonnegativedefinite solution $Q$ of this modified Lyapunov equation is then guaranteed to be an upper bound
for $Q_{\Delta A}$. The following easily proved result is fundamental and forms the basis for all later developments. The hypothesis of this result are of a general nature and are not intended to be directly verifiable. Suitably verifiable hypotheses are discussed later.

Theorem 3.1. Let $\Omega: \mathbb{N}^{n} \rightarrow \mathbb{N}^{n}$ be such that

$$
\begin{equation*}
\Delta A Q+Q \Delta A^{\mathrm{T}} \leq \Omega(Q), \quad \Delta A \in U, \quad Q \in \mathbb{N}^{n} \tag{3.1}
\end{equation*}
$$

and suppose there exists $Q \in \mathbb{N}^{n}$ satisfying

$$
\begin{equation*}
0=A Q+Q A^{\mathbf{T}}+\Omega(Q)+V . \tag{3.2}
\end{equation*}
$$

Then

$$
\begin{equation*}
\left(A+\Delta A,\left[V+\Omega(Q)-\left(\Delta A Q+Q \Delta A^{T}\right)\right]^{\frac{1}{2}}\right) \text { is stabilizable, } \quad \Delta A \in U \tag{3.3}
\end{equation*}
$$

if and only if

$$
\begin{equation*}
A+\Delta A \text { is asymptotically stable, } \quad \Delta A \in U . \tag{3.4}
\end{equation*}
$$

In this case,

$$
\begin{equation*}
Q_{\Delta A} \leq Q, \quad \Delta A \in U \tag{3.5}
\end{equation*}
$$

where $Q_{\Delta A} \in \mathbb{N}^{n}$ is given by (2.7), and

$$
\begin{gather*}
J_{\mathrm{S}}(U) \leq \operatorname{tr} Q R,  \tag{3.6}\\
J_{\mathrm{D}}(U) \leq \lambda_{\max }(Q R) . \tag{3.7}
\end{gather*}
$$

In addition, there exists $\Delta A \in U$ such that $\left(A+\Delta A,\left[V+\Omega(Q)-\left(\Delta A Q+Q \Delta A^{T}\right)\right]^{\frac{1}{2}}\right)$ is controllable if and only if $Q$ is positive definite. In this case $\left(A+\Delta A,\left[V+\Omega(Q)-\left(\Delta A Q+Q \Delta A^{T}\right)\right]^{\frac{1}{2}}\right)$ is controllable for all $\Delta A \in U$.

Proof. We stress that in (3.1) " $Q$ " denotes an arbitrary element of $\mathbb{N}^{n}$ while $Q$ in (3.2) denotes a specific solution of the modified Lyapunov equation. This minor abuse of notation considerably simplifies the presentation. Now note that for all $\Delta A \in \mathbb{R}^{n \times n}$ (3.2) is equivalent to

$$
\begin{equation*}
0=(A+\Delta A) Q+Q(A+\Delta A)^{\mathrm{T}}+\Omega(Q)-\left(\Delta A Q+Q \Delta A^{\mathrm{T}}\right)+V \tag{3.8}
\end{equation*}
$$

Hence, by assumption (3.8) has a solution $Q \in \mathbb{N}^{n}$ for all $\Delta A \in \mathbb{R}^{n \times n}$. If $\Delta A$ is restricted to the set $U$ then by (3.1) $\Omega(Q)-\left(\Delta A Q+Q \Delta A^{T}\right)$ is nonnegative definite. Now if the stabilizability condition (3.3) holds for all $\Delta A \in U$, it follows from Lemma 12.2 of [39] that $A+\Delta A$ is asymptotically stable
for all $\Delta A \in U$. Conversely, if $A+\Delta A$ is asymptotically stable for all $\Delta A \in U$, then (3.3) holds. Next, subtracting (2.7) from (3.8) yields

$$
0=(A+\Delta A)\left(Q-Q_{\Delta A}\right)+\left(Q-Q_{\Delta A}\right)(A+\Delta A)^{\mathrm{T}}+\Omega(Q)-\left(\Delta A Q+Q \Delta A^{\mathrm{T}}\right), \quad \Delta A \in U
$$

or, equivalently, (since $A+\Delta A$ is asymptotically stable for all $\Delta A \in U$ )

$$
Q-Q_{\Delta A}=\int_{0}^{\infty} e^{(A+\Delta A) t}\left[\Omega(Q)-\left(\Delta A Q+Q \Delta A^{T}\right)\right] e^{(A+\Delta A)^{T} t} d t \geq 0, \quad \Delta A \in U
$$

which implies (3.5). The performance bound (3.6) is now an immediate consequence of (3.5). To prove (3.7) we note that it follows from Corollary 7.7 .4 of [40] that if $0 \leq M_{1} \leq M_{2}$ then $\lambda_{\max }\left(M_{1}\right) \leq \lambda_{\max }\left(M_{2}\right)$. Thus

$$
\begin{equation*}
J_{\mathrm{D}}(U)=\lambda_{\max }\left(Q_{\Delta A} R\right)=\lambda_{\max }\left(E_{0} Q_{\Delta A} E_{0}^{\mathrm{T}}\right) \leq \lambda_{\max }\left(E_{0} Q E_{0}^{\mathrm{T}}\right)=\lambda_{\max }(Q R) \tag{3.9}
\end{equation*}
$$

Finally, it follows from (3.8) that the controllability condition holds for some $\Delta A \in U$ if and only if the Gramian $Q$ is positive definite. Since $Q$ is also the Gramian corresponding to $A+\Delta A$ for all $\Delta A \in U$, then controllability holds for all $\Delta A \in U$.

For convenience we shall say that $\Omega(\cdot)$ bounds $U$ if (3.1) is satisfied. To apply Theorem 3.1, one first specifies a function $\Omega(\cdot)$ and an uncertainty set $U$ such that $\Omega(\cdot)$ bounds $U$. If the existence of a nonnegative-definite solution $Q$ to (3.2) satisfying (3.3) or, equivalently, (3.4) can be determined analytically or numerically, then robust stability is guaranteed. One can then enlarge $U$, modify $\Omega(\cdot)$, and again attempt to solve (3.2). If, however, a nonnegative-definite solution to (3.2) cannot be determined, then $U$ must be decreased in size until (3.2) is solvable. For example, $\Omega(\cdot)$ can be replaced by $\varepsilon \Omega(\cdot)$ to bound $\varepsilon U$, where $\varepsilon>1$ enlarges $U$ and $\varepsilon<1$ shrinks $U$. Of course, the actual range of uncertainty which can be bounded depends upon the nominal matrix $A$, the function $\Omega(\cdot)$, and the structure of $U$. In Section 5 the uncertainty set $U$ and bound $\Omega(\cdot)$ satisfying (3.1) are given concrete forms. We complete this section with several observations.

Remark 3.1. In applying Theorem 3.1 it may be convenient to replace condition (3.3) with a stronger condition which is easier to verify in practice. Clearly, (3.3) is satisfied if $V+\Omega(Q)-$ $\left(\Delta A Q+Q \Delta A^{T}\right)$ is positive definite for all $\Delta A \in U$. This will be the case, for example, if either $V$ is positive definite or strict inequality holds in (3.1). Also, it follows from Theorem 3.6 of [39] that (3.3) is implied by the stronger condition

$$
\begin{equation*}
\left(A+\Delta A, V^{\frac{1}{2}}\right) \text { is stabilizable, } \quad \Delta A \in U \tag{3.10}
\end{equation*}
$$

Similar remarks apply to the controllability condition.
Remark 3.2. If only robust stability is of interest, then the noise intensity $V$ need not have physical significance. In this case one may either set $V=\varepsilon I_{n}$ where $\varepsilon>0$ is small to satisfy $V>0$, or set $V=0$ and require that strict inequality hold in (3.1).

Remark 3.3. Since $A$ is stable $Q$ satisfying (3.2) is given by

$$
\begin{equation*}
Q=\int_{0}^{\infty} e^{\boldsymbol{A} t}[\Omega(Q)+V] e^{\boldsymbol{A}^{\mathrm{T}} \mathrm{t}} d t \tag{3.11}
\end{equation*}
$$

or, equivalently,

$$
\begin{equation*}
Q=\int_{0}^{\infty} e^{A t} \Omega(Q) e^{A^{T}} d t+Q_{0} \tag{3.12}
\end{equation*}
$$

where $Q_{0} \in \mathbb{N}^{n}$ is defined by

$$
\begin{equation*}
Q_{0} \triangleq \int_{0}^{\infty} e^{A t} V e^{A^{\mathrm{T}} \mathrm{t}} d t \tag{3.13}
\end{equation*}
$$

and satisfies

$$
\begin{equation*}
0=A Q_{0}+Q_{0} A^{T}+V . \tag{3.14}
\end{equation*}
$$

Note that $Q_{0} \leq Q$ and that the nominal performances $J_{\mathrm{S}}(\{0\})$ and $J_{\mathrm{D}}(\{0\})$ are given by tr $Q_{0} R$ and $\lambda_{\text {max }}\left(Q_{0} R\right)$, respectively.

Remark 3.4. Using (3.11) it is also useful to note that the bound for $J_{\mathrm{s}}(u)$ given by (3.6) can be written as

$$
\begin{equation*}
\operatorname{tr} Q R=\operatorname{tr} \int_{0}^{\infty} e^{\boldsymbol{A t}}[\Omega(Q)+V] e^{\boldsymbol{A}^{\tau} t} d t R=\operatorname{tr} P_{0}[\Omega(Q)+V] \tag{3.15}
\end{equation*}
$$

where $P_{0} \in \mathbb{N}^{n}$ is defined by

$$
\begin{equation*}
P_{0} \triangleq \int_{0}^{\infty} e^{A^{\boldsymbol{T}}} R e^{A t} d t \tag{3.16}
\end{equation*}
$$

and satisfies

$$
\begin{equation*}
0=A^{\mathbf{T}} P_{0}+P_{0} A+R . \tag{3.17}
\end{equation*}
$$

The bound $\operatorname{tr} P_{0}[\Omega(Q)+V]$ can be viewed as a dual formulation of the bound $\operatorname{tr} Q R$ since the roles of $A$ and $A^{T}$ are reversed. Dual bounds are developed in the following section. Note that $\operatorname{tr} Q_{0} R=\operatorname{tr} P_{0} V$.

Remark 3.5. If $\Omega(\cdot)$ bounds $U$ then clearly $\Omega(\cdot)$ bounds the convex hull of $U$. Hence, only convex uncertainty sets $U$ need be considered. Next, we shall later use the obvious fact that if $\Omega^{\prime}(\cdot)$ bounds $U^{\prime}$ and $\Omega^{\prime \prime}(\cdot)$ bounds $U^{\prime \prime}$, then $\Omega^{\prime}(\cdot)+\Omega^{\prime \prime}(\cdot)$ bounds $U^{\prime}+U^{\prime \prime}$. Hence if $U$ can be decomposed
additively then it suffices to bound each component separately. Finally, if $\Omega(\cdot)$ bounds $U$ and there exists $\Omega^{\prime}: \mathbb{N}^{n} \rightarrow \mathbb{N}^{n}$ such that $\Omega(Q) \leq \Omega^{\prime}(Q)$ for all $Q \in \mathbb{N}^{n}$, then $\Omega^{\prime}(\cdot)$ also bounds $U$. That is, any overbound $\Omega^{\prime}(\cdot)$ for $\Omega(\cdot)$ also bounds $U$. Of course, as we shall see, it is quite possible that an overbound $\Omega^{\prime}(\cdot)$ for $\Omega(\cdot)$ may bound a set $U^{\prime}$ which is larger than the "original" uncertainty set $U$.

## 4. Dual Sufficient Conditions for Robust Stability and Performance

As noted in Remark 3.4, the performance bound $\operatorname{tr} Q R$ given by (3.6) can be expressed equivalently in terms of a dual variable $P_{0}$ for which the roles of $A$ and $A^{T}$ are reversed. Using a similar technique, additional conditions for robust stability and performance can be obtained by developing a dual version of Theorem 3.1. A prime motivation for developing such dual bounds is to draw connections with previous results in the literature relating to robust stability. In particular, note that traditional robust stability techniques based upon the Lyapunov function $V(x)=x^{T} P x$ lead to dual conditions. Robust performance bounds within the dual formulation are difficult to motivate without first developing the primal performance bounds. In addition, the dual bounds may, for certain problems, yield larger stability regions and sharper performance bounds than the "primal" bounds introduced in the previous section.

Lemma 4.1. Suppose (2.1) is asymptotically stable for all $\Delta A \in U$. Then

$$
\begin{equation*}
J_{\mathbf{s}}(u)=\sup _{\Delta A \in u} \operatorname{tr} P_{\Delta A} V \tag{4.1}
\end{equation*}
$$

where $n \times n P_{\Delta A}$ is the unique, nonnegative- definite solution to

$$
\begin{equation*}
0=(A+\Delta A)^{\mathrm{T}} P_{\Delta A}+P_{\Delta A}(A+\Delta A)+R . \tag{4.2}
\end{equation*}
$$

Proof. It need only be noted that

$$
\operatorname{tr} Q_{\Delta A} R=\operatorname{tr} \int_{0}^{\infty} e^{(A+\Delta A) t} V e^{(A+\Delta A)^{T} t} d t R=\operatorname{tr} \dot{P_{\Delta A}} V
$$

where

$$
P_{\Delta A} \triangleq \int_{0}^{\infty} e^{(A+\Delta A)^{T} t} R e^{(A+\Delta A) t} d t
$$

satisfies (4.2).
The proof of Lemma 4.1 relied on the fact that $\operatorname{tr} Q_{\Delta A} R=\operatorname{tr} P_{\Delta A} V$. However, it is not necessarily true that $\lambda_{\max }\left(Q_{\Delta A} R\right)=\lambda_{\max }\left(P_{\Delta A} V\right)$ even when $\Delta A=0$. For example, if $A=$
$\left[\begin{array}{cc}-1 & 0 \\ 0 & -2\end{array}\right], R=I_{2}$, and $V=\left[\begin{array}{ll}2 & 1 \\ 1 & 1\end{array}\right]$ then $Q_{0} R=\left[\begin{array}{cl}1 & 1 / 3 \\ 1 / 3 & 1 / 4\end{array}\right]$ and $P_{0} V=\left[\begin{array}{cc}1 & 1 / 2 \\ 1 / 4 & 1 / 4\end{array}\right]$ and thus $\lambda_{\text {max }}\left(Q_{0} R\right)=$ $(15+\sqrt{145}) / 24$ and $\lambda_{\text {max }}\left(P_{0} V\right)=(5+\sqrt{17}) / 8$. Thus to obtain a suitable dual version of $J_{\mathrm{D}}(U)$ we need to define a dual deterministic cost $\hat{J}_{\mathrm{D}}(U)$. This can be done if the disturbance signals are taken to be integrable rather than square integrable. Thus, for measurable $w:[0, \infty) \rightarrow \mathbb{R}^{d}$ define

$$
\begin{equation*}
\|w(\cdot)\|_{1} \triangleq \int_{0}^{\infty}\left[w^{\mathbf{T}}(t) w(t)\right]^{\frac{1}{2}} d t \tag{4.3}
\end{equation*}
$$

which is an $L_{1}$ function norm with a Euclidean spatial norm. The dual deterministic cost $\hat{J}_{\mathrm{D}}(u)$ is thus defined by

$$
\begin{equation*}
\hat{J}_{\mathrm{D}}(u) \triangleq \sup _{\Delta A \in u} \sup _{\|ш(\cdot)\|_{1} \leq 1}\|E x(\cdot)\|_{2}^{2} \tag{4.4}
\end{equation*}
$$

where $E x(\cdot)$ is measured according to the energy norm (2.8). The following dual result can also be found in [38].

Lemma 4.2. Suppose (2.1) is asymptotically stable for all $\triangle A \in U$. Then

$$
\begin{equation*}
\hat{J}_{\mathrm{D}}(U)=\lambda_{\max }\left(P_{\Delta A} V\right) \tag{4.5}
\end{equation*}
$$

where $P_{\Delta A}$ is the unique, nonnegative-definite solution to (4.2).
The dual version of Theorem 3.1 can now be stated.
Theorem 4.1. Let $\Lambda: \mathbb{N}^{n} \rightarrow \mathbb{N}^{n}$ be such that

$$
\begin{equation*}
\Delta A^{\mathbf{T}} P+P \Delta A \leq \Lambda(P), \quad \Delta A \in U, \quad P \in \mathbb{N}^{n} \tag{4.6}
\end{equation*}
$$

and suppose there exists $P \in \mathbb{N}^{n}$ satisfying

$$
\begin{equation*}
0=A^{\mathbf{T}} P+P A+\Lambda(P)+R \tag{4.7}
\end{equation*}
$$

Then

$$
\begin{equation*}
\left(\left[R+\Lambda(P)-\left(\Delta A^{\mathbf{T}} P+P \Delta A\right)\right]^{\frac{1}{2}}, A+\Delta A\right) \text { is detectable, } \quad \Delta A \in U \tag{4.8}
\end{equation*}
$$

if and only if

$$
\begin{equation*}
A+\Delta A \text { is asymptotically stable, } \quad \Delta A \in U . \tag{4.9}
\end{equation*}
$$

In this case,

$$
\begin{equation*}
P_{\Delta A} \leq P, \quad \Delta A \in U \tag{4.10}
\end{equation*}
$$

where $P_{\Delta A}$ is given by (4.2), and

$$
\begin{equation*}
J_{\mathbf{s}}(u) \leq \operatorname{tr} P V, \tag{4.11}
\end{equation*}
$$

$$
\begin{equation*}
\hat{J}_{\mathrm{D}}(U) \leq \lambda_{\max }(P V) \tag{4.12}
\end{equation*}
$$

In addition, there exists $\Delta A \in U$ such that $\left(\left[R+\Lambda(P)-\left(\Delta A^{T} P+P \Delta A\right)\right]^{\frac{1}{2}}, A+\Delta A\right)$ is observable if and only if $P$ is positive definite. In this case $\left(\left[R+\Lambda(P)-\left(\Delta A^{T} P+P \Delta A\right)\right]^{\frac{1}{2}}, A+\Delta A\right)$ is observable for all $\Delta A \in U$.

Proof. The proof is completely analogous to the proof of Theorem 3.1.
Remark 4.1. Note that $\hat{J}_{\mathrm{D}}(u) \leq J_{\mathrm{S}}(u)$ and $\hat{J}_{\mathrm{D}}(u)=J_{\mathrm{S}}(u)$ if rank $V=1$. Combining this fact with Remark 2.1, it follows that $J_{\mathrm{D}}(u)=\hat{J}_{\mathrm{D}}(U)$ if both rank $R=1$ and rank $V=1$.

It is quite possible that the bounds $\operatorname{tr} Q R$ and $\operatorname{tr} P V$ for $J_{S}(u)$ given by (3.6) and (4.11) may be different in spite of the fact as shown in the proof of Lemma 4.1 that $\operatorname{tr} Q_{\Delta A} R=\operatorname{tr} P_{\Delta A} V$. That is, depending upon $\Omega(\cdot)$ and $\Lambda(\cdot)$ either bound (3.6) or bound (4.11) may be better for a particular problem. In general, we have the following result.

Proposition 4.1. Let $\Omega(\cdot), \Lambda(\cdot), Q$, and $P$ be as in Theorems 3.1 and 4.1, and let $Q_{0}$ and $P_{0}$ be given by (3.13) and (3.16), respectively. Then

$$
\begin{align*}
& \operatorname{tr} Q_{0} \Lambda(P)<\operatorname{tr} P_{0} \Omega(Q) \Longleftrightarrow \operatorname{tr} Q R>\operatorname{tr} P V  \tag{4.13}\\
& \operatorname{tr} Q_{0} \Lambda(P)=\operatorname{tr} P_{0} \Omega(Q) \Longleftrightarrow \operatorname{tr} Q R=\operatorname{tr} P V  \tag{4.14}\\
& \operatorname{tr} Q_{0} \Lambda(P)>\operatorname{tr} P_{0} \Omega(Q) \Longleftrightarrow \operatorname{tr} Q R<\operatorname{tr} P V \tag{4.15}
\end{align*}
$$

Proof. Note that

$$
\operatorname{tr} Q R=\int_{0}^{\infty} e^{A t}[\Omega(Q)+V] e^{A^{T} t} d t R=\operatorname{tr} P_{0} \Omega(Q)+\operatorname{tr} \int_{0}^{\infty} e^{A t} V e^{A^{T} t} d t R
$$

and

$$
\operatorname{tr} P V=\operatorname{tr} \int_{0}^{\infty} e^{A^{\mathrm{T}} \mathrm{t}}[\Lambda(P)+R] e^{A t} d t V=\operatorname{tr} Q_{0} \Lambda(P)+\operatorname{tr} \int_{0}^{\infty} e^{A^{\mathrm{T}} t} R e^{A t} d t V
$$

so that

$$
\operatorname{tr} Q R-\operatorname{tr} P V=\operatorname{tr} P_{0} \Omega(Q)-\operatorname{tr} Q_{0} A(P),
$$

which yields (4.13)-(4.15).
Remark 4.2. Finally, as shown by example above, it is not generally true that $J_{\mathrm{D}}(U)=\hat{J}_{\mathrm{D}}(U)$. Similarly, we should not expect that the bounds $\lambda_{\text {max }}(Q R)$ and $\lambda_{m a x}(P V)$ for $J_{\mathrm{D}}(U)$ and $\hat{J}_{\mathrm{D}}(U)$ given by (3.7) and (4.12) are equal.

## 5. Construction of the Bounds $\Omega(\cdot)$ and $\Lambda(\cdot)$

As discussed in Section 1, we consider three distinct classes of bounds $\Omega(\cdot)$ denoted by Type I, Type III, and Type III. Roughly speaking, these bounds exploit, respectively, the symmetry of the Lyapunov terms $\Delta A Q+Q \Delta A^{T}$, the structure of $Q$, and the structure of $\Delta A$. The dual bounds $A(\cdot)$ can be constructed similarly by replacing $Q$ and $\Delta A$ by $P$ and $\Delta A^{T}$. Hence these bounds will not be discussed separately. For convenience in discussing the set $U$, we shall use the terms rectangle and ellipse to refer to closed regions bounded by such figures in multiple dimensions. As usual, a polytope is the convex hull of a finite number of points.

### 5.1 Type I Bounds

We begin by constructing bounds $\Omega(\cdot)$ which exploit only the symmetry of the Lyapunov terms $\Delta A Q+Q \Delta A^{T}$. First we require the following definition of a function of a symmetric matrix as an extension of a real-valued function ([40], p. 300). Specifically, if $f: \mathbb{R} \rightarrow \mathbb{R}$ then (with a minor abuse of notation) $f: \mathbb{S}^{\boldsymbol{n}} \rightarrow \mathbb{S}^{\boldsymbol{n}}$ can be defined by setting

$$
f(S) \triangleq U f(D) U^{\mathbf{T}}
$$

where $S=U D U^{\mathbf{T}}, U$ is orthogonal, $D$ is real diagonal, and $f(D)$ is the diagonal matrix obtained by applying $f$ to each diagonal element of $D$. Note that if $f$ is the polynomial $f(x)=\sum_{i=0}^{r} a_{i} x^{i}$ then $f(S)=\sum_{i=0}^{r} a_{i} S^{i}$. Note also that if $f(x)=|x|$ then $f(S)=\left(S^{2}\right)^{\frac{1}{2}}$, where $(\cdot)^{\frac{1}{2}}$ denotes the (unique) nonnegative-definite square root. As in [41], p. 262, we use the notation $|S|$ to denote $\left(S^{2}\right)^{\frac{1}{2}}$. Finally, note that if $f: \mathbb{R} \rightarrow \mathbb{R}$ and $g: \mathbb{R} \rightarrow \mathbb{R}$ are such that $f(x) \leq g(x), x \in \mathbb{R}$, then $f(S) \leq g(S), S \in \mathbb{S}^{n}$.

As a concretization of the uncertainty set $U$, consider the set

$$
\begin{equation*}
u_{1} \triangleq\left\{\Delta A \in \mathbb{R}^{n \times n}: \quad \Delta A=\sum_{i=1}^{p} \sigma_{i} A_{i}, \quad\left|\sigma_{i}\right| \leq \delta_{i}, \quad i=1, \ldots, p\right\} \tag{5.1}
\end{equation*}
$$

:Vhere, for $i=1, \ldots, p: A_{i} \in \mathbb{R}^{n \times n}$ is a given matrix denoting the structure of the parametric uncertainty; $\sigma_{i}$ is a real uncertain parameter; and $\delta_{i}$ denotes the range of parameter uncertainty. Clearly, the multidimensional set of uncertain parameters ( $\sigma_{1}, \ldots, \sigma_{p}$ ) is the rectangle $\left[-\delta_{1}, \delta_{1}\right] \times$ $\ldots \times\left[-\delta_{p}, \delta_{p}\right]$ and $U_{1}$ is a symmetric polytope of matrices in $\mathbb{R}^{n \times n}$. Note that the symmetry of the uncertainty interval $\left[-\delta_{i}, \delta_{i}\right]$ entails no loss of generality since the nominal value of $A$ can be redefined if necessary. Furthermore, it is also possible, without loss of generality, to define $\delta_{i}=1$ by replacing $A_{i}$ by $\delta_{i} A_{i}$. For clarity, however, we choose not to employ this scaling.

We begin by considering the absolutevalue bound utilized by Chang and Peng in [9].
Proposition 5.1. The function

$$
\begin{equation*}
\Omega_{1}(Q) \triangleq \sum_{i=1}^{p} \delta_{i}\left|A_{i} Q+Q A_{i}^{T}\right| \tag{5.2}
\end{equation*}
$$

bounds $u_{1}$.
Proof. For $i=1, \ldots, p$ and $\left|\sigma_{i}\right| \leq \delta_{i}$,

$$
\sigma_{i}\left(A_{i} Q+Q A_{i}^{\mathrm{T}}\right) \leq\left|\sigma_{i}\left(A_{i} Q+Q A_{i}^{\mathrm{T}}\right)\right|=\left|\sigma_{i}\right|\left|A_{i} Q+Q A_{i}^{\mathrm{T}}\right| \leq \delta_{i}\left|A_{i} Q+Q A_{i}^{\mathrm{T}}\right|
$$

Thus, summing over $i$ yields

$$
\Delta A Q+Q \Delta A^{\mathrm{T}}=\sum_{i=1}^{p} \sigma_{i}\left(A_{i} Q+Q A_{i}^{\mathrm{T}}\right) \leq \sum_{i=1}^{p} \delta_{i}\left|A_{i} Q+Q A_{i}^{\mathrm{T}}\right|
$$

which yields (3.1) with $\Omega(\cdot)=\Omega_{1}(\cdot)$ and $u=U_{1}$.
Remark 5.1. It is tempting to prove Proposition 5.1 by writing

$$
\sum_{i=1}^{p} \sigma_{i}\left(A_{i} Q+Q A_{i}^{\mathrm{T}}\right) \leq\left|\sum_{i=1}^{p} \sigma_{i}\left(A_{i} Q+Q A_{i}^{\mathrm{T}}\right)\right| \leq \sum_{i=1}^{p}\left|\sigma_{i}\left(A_{i} Q+Q A_{i}^{\mathrm{T}}\right)\right| .
$$

However, counterexamples show that the inequality $\left|M_{1}+M_{2}\right| \leq\left|M_{1}\right|+\left|M_{2}\right|$ is not generally true for arbitrary symmetric matrices $M_{1}, M_{2}$.

Remark 5.2. Because of its simplicity it is tempting to conjecture that $\Omega_{1}(\cdot)$ is the best bound for $\Delta A Q+Q \Delta A^{T}$ over the set $U_{1}$. To show that this is not the case, let $Q=\frac{1}{2} I_{2}, p=1, A_{1}=$ $\left[\begin{array}{cc}1 & 0 \\ 0 & -1\end{array}\right]$, and $\delta_{1}=1$. Then $\sigma_{1}\left(A_{1} Q+Q A_{1}^{\mathrm{T}}\right) \leq \delta_{1}\left|A_{1} Q+Q A_{1}^{\mathrm{T}}\right|=I_{2},\left|\sigma_{1}\right| \leq 1$. However, it is also true that $\sigma_{1}\left(A_{1} Q+Q A_{1}^{\mathrm{T}}\right) \leq\left[\begin{array}{cc}2 & 3 / 2 \\ 3 / 2 & 2\end{array}\right],\left|\sigma_{1}\right| \leq 1$. Neither bound, however, is an overbound for the other. This is a consequence of the fact that the nonnegative-definite matrix ordering is only a partial order.

As mentioned above, an overbound for $\Omega_{1}(\cdot)$ will also bound $U_{1}$.
Lemma 5.1. For $i=1, \ldots, p$, let $f_{i}: \mathbb{R} \rightarrow \mathbb{R}$ satisfy

$$
\begin{equation*}
f_{i}(x) \geq|x|, \quad x \in \mathbb{R} \tag{5.3}
\end{equation*}
$$

Then the function

$$
\begin{equation*}
\Omega_{2}(Q) \triangleq \sum_{i=1}^{p} \delta_{i} f_{i}\left(A_{i} Q+Q A_{i}^{\mathrm{T}}\right) \tag{5.4}
\end{equation*}
$$

is an overbound for $\Omega_{1}(\cdot)$ and hence also bounds $U_{1}$.
One particular choice of $f_{i}$ satisfying (5.3) will be considered here, namely, the polynomial

$$
\begin{equation*}
f_{i}(x)=\frac{1}{4} \beta_{i}+\beta_{i}^{-1} x^{2} \tag{5.5}
\end{equation*}
$$

where $\beta_{i}$ is an arbitrary positive constant. Thus $\Omega_{2}(\cdot)$ has the following specialization.
Corollary 5.1. Let $\beta_{1}, \ldots, \beta_{p}$ be arbitrary positive constants. Then the function

$$
\begin{equation*}
\Omega_{3}(Q) \triangleq \frac{1}{4} \sum_{i=1}^{p} \delta_{i} \beta_{i} I_{n}+\sum_{i=1}^{p}\left(\delta_{i} / \beta_{i}\right)\left(A_{i} Q+Q A_{i}^{T}\right)^{2} \tag{5.6}
\end{equation*}
$$

is an overbound for $\Omega_{1}(\cdot)$ and hence also bounds $U_{1}$.
Although overbounding $\Omega_{1}(\cdot)$ by $\Omega_{3}(\cdot)$ results in a looser bound for $u_{1}$, it turns out that $\Omega_{3}(\cdot)$ actually bounds a set which is larger than $U_{1}$. Specifically, in place of $U_{1}$ consider

$$
\begin{equation*}
u_{2} \triangleq\left\{\Delta A \in \mathbb{R}^{n \times n}: \quad \Delta A=\sum_{i=1}^{p} \sigma_{i} A_{i}, \quad \sum_{i=1}^{p} \sigma_{i}^{2} / \alpha_{i}^{2} \leq 1\right\} \tag{5.7}
\end{equation*}
$$

where $\alpha_{1}, \ldots, \alpha_{p}$ are given positive constants. Note that (5.7) replaces the rectangle of uncertain parameters $\left(\sigma_{1}, \ldots, \sigma_{p}\right)$ by an ellipse. Thus the set $U_{2}$ of matrix perturbations is an ellipse of matrices in $\mathbb{R}^{n \times n}$ in contrast to the polytope $u_{1}$. Of course, $u_{i}=u_{2}$ if $p=1$ and $\alpha_{1}=\delta_{1}$. Again it is possible to take $\alpha_{i}=1$ without loss of generality by replacing $A_{i}$ by $\alpha_{i} A_{i}$. We again choose not to do this, however. The following lemma provides a convenient characterization of the relationship between the rectangle $U_{1}$ and the ellipse $U_{2}$.

Proposition 5.2. Suppose $U_{1}$ is defined by the positive constants $\delta_{1}, \ldots, \delta_{p}$ and let $U_{2}$ be characterized by

$$
\begin{equation*}
\alpha_{i}=\left(\alpha \delta_{i} / \beta_{i}\right)^{\frac{1}{2}}, \quad i=1, \ldots, p, \tag{5.8}
\end{equation*}
$$

where $\alpha$ is defined by

$$
\begin{equation*}
\alpha=\sum_{i=1}^{p} \delta_{i} \beta_{i} \tag{5.9}
\end{equation*}
$$

and $\beta_{1}, \ldots \beta_{p}$ are arbitrary positive constants. Then $U_{2}$ contains $\mathcal{U}_{1}$. Specifically, the ellipse $\mathcal{U}_{2}$ circumscribes the polytope $U_{1}$. Furthermore, $\Omega_{3}(\cdot)$ actually bounds $U_{2}$.

Proof. If $\left|\sigma_{i}\right| \leq \delta_{i}, i=1, \ldots, p$, then it follows from (5.8) that

$$
\sum_{i=1}^{p} \sigma_{i}^{2} / \alpha_{i}^{2}=\alpha^{-1} \sum_{i=1}^{p} \beta_{i} \sigma_{i}^{2} / \delta_{i} \leq \alpha^{-1} \sum_{i=1}^{p} \beta_{i} \delta_{i}=1
$$

If, in addition, $\left|\sigma_{i}\right|=\delta_{i}, i=1, \ldots, p$, then $\sum_{i=1}^{p} \delta_{i}^{2} / \alpha_{i}^{2}=1$, which corresponds to a point on the boundary of the ellipse. To show that $\Omega_{3}(\cdot)$ actually bounds $U_{2}$ note that

$$
\begin{aligned}
0 & \leq \sum_{i=1}^{p}\left[\frac{1}{2}\left(\alpha^{\frac{1}{2}} \sigma_{i} / \alpha_{i}\right) I_{n}-\left(\alpha_{i} / \alpha^{\frac{1}{2}}\right)\left(A_{i} Q+Q A_{i}^{\mathrm{T}}\right)\right]^{2} \\
& =\frac{\&}{4} \sum_{i=1}^{p}\left(\sigma_{i}^{2} / \alpha_{i}^{2}\right) I_{n}+\alpha^{-1} \sum_{i=1}^{p} \alpha_{i}^{2}\left(A_{i} Q+Q A_{i}^{\mathrm{T}}\right)^{2}-\left(\Delta A Q+Q \Delta A^{\mathrm{T}}\right) .
\end{aligned}
$$

Since $\sum_{i=1}^{p} \sigma_{i}^{2} / \alpha_{i}^{2} \leq 1$ in $\mathcal{U}_{2}$, it follows that

$$
\Delta A Q+Q \Delta A^{T} \leq \frac{\alpha}{4} I_{n}+\alpha^{-1} \sum_{i=1}^{p} \alpha_{i}^{2}\left(A_{i} Q+Q A_{i}^{T}\right)^{2}
$$

Utilizing (5.8) and (5.9) to substitute for $\alpha$ and $\alpha_{i}$ yields (3.1) with $\Omega(\cdot)=\Omega_{s}(\cdot)$ and $U=U_{2}$.
Proposition 5.2 shows that each choice of constants $\beta_{1}, \ldots, \beta_{p}>0$ leads to a particular ellipse $U_{2}$ which contains the polytope $U_{1}$. Furthermore, $\Omega_{3}(\cdot)$, which by Corollary 5.1 bounds $U_{1}$, actually bounds the larger set $\mathcal{U}_{2}$. For convenience, we now dispense with the constants $\beta_{1}, \ldots, \beta_{p}$ which relate the rectangle $U_{1}$ to the ellipse $U_{2}$ and we characterize $\Omega_{3}(\cdot)$ entirely in terms of $\alpha_{1}, \ldots, \alpha_{p}$.

Corollary 5.2. Let $\alpha$ be an arbitrary positive constant. Then the function

$$
\begin{equation*}
\Omega_{4}(Q) \triangleq \frac{\propto}{4} I_{n}+\alpha^{-1} \sum_{i=1}^{p} \alpha_{i}^{2}\left(A_{i} Q+Q A_{i}^{T}\right)^{2} \tag{5.10}
\end{equation*}
$$

bounds $U_{2}$.
Remark 5.3. Within the context of Corollary 5.2, the positive constant $\alpha$ plays no role in defining the set $U_{2}$, although $\Omega_{4}(\cdot)$ is guaranteed to bound $U_{2}$ for all choices of $\alpha$. It can be expected, however, that certain choices of $\alpha$ provide better bounds than other choices. This will be seen by example in Section 10.

### 5.2 Type II Bounds

We now consider additional bounds for $U$ which exploit the structure of $Q$. For these bounds the natural uncertainty set is given by $U_{2}$.

Proposition 5.3. Let $\alpha$ be an arbitrary positive number and, for each $Q \in \mathbb{N}^{n}$, let $Q_{1} \in$ $\mathbb{R}^{n \times m}$ and $Q_{\mathbf{2}} \in \mathbb{R}^{m \times n}$ satisfy

$$
\begin{equation*}
Q=Q_{1} Q_{2} . \tag{5.11}
\end{equation*}
$$

Then the function

$$
\begin{equation*}
\Omega_{\mathrm{s}}(Q) \triangleq \alpha Q_{2}^{\mathrm{T}} Q_{2}+\alpha^{-1} \sum_{i=1}^{p} \alpha_{i}^{2} A_{i} Q_{1} Q_{1}^{\mathrm{T}} A_{i}^{\mathrm{T}} \tag{5.12}
\end{equation*}
$$

bounds $\boldsymbol{U}_{2}$.
Proof. Note that

$$
\begin{aligned}
0 & \leq \sum_{i=1}^{p}\left[\left(\alpha^{\frac{1}{2}} \sigma_{i} / \alpha_{i}\right) Q_{2}^{\mathrm{T}}-\left(\alpha_{i} / \alpha^{\frac{1}{2}}\right) A_{i} Q_{1}\right]\left[\left(\alpha^{\frac{1}{2}} \sigma_{i} / \alpha_{i}\right) Q_{2}^{\mathrm{T}}-\left(\alpha_{i} / \alpha^{\frac{1}{2}}\right) A_{i} Q_{1}\right]^{\mathrm{T}} \\
& =\alpha \sum_{i=1}^{p}\left(\sigma_{i}^{2} / \alpha_{i}^{2}\right) Q_{2}^{\mathrm{T}} Q_{2}+\alpha^{-1} \sum_{i=1}^{p} \alpha_{i}^{2} A_{i} Q_{1} Q_{1}^{\mathrm{T}} A_{i}^{\mathrm{T}}-\sum_{i=1}^{p} \sigma_{i}\left(A_{i} Q+Q A_{i}^{\mathrm{T}}\right)
\end{aligned}
$$

which, since $\sum_{i=1}^{p} \sigma_{i}^{2} / \alpha_{i}^{2} \leq 1$, yields (3.1) with $\Omega(\cdot)=\Omega_{5}(\cdot)$ and $U=U_{2}$.
We consider three specializations of $\Omega_{5}(\cdot)$. Specifically, we set $m=n$ and define

$$
\begin{gather*}
Q_{1}=Q, \quad Q_{2}=I_{n}  \tag{5.13}\\
Q_{1}=Q_{2}=Q^{\frac{1}{2}}  \tag{5.14}\\
Q_{1}=I_{n}, \quad Q_{2}=Q \tag{5.15}
\end{gather*}
$$

Corollary 5.3. Let $\alpha$ be an arbitrary positive number. Then the functions

$$
\begin{align*}
& \Omega_{6}(Q) \triangleq \alpha I_{n}+\alpha^{-1} \sum_{i=1}^{p} \alpha_{i}^{2} A_{i} Q^{2} A_{i}^{\mathrm{T}}  \tag{5.16}\\
& \Omega_{7}(Q) \triangleq \alpha Q+\alpha^{-1} \sum_{i=1}^{p} \alpha_{i}^{2} A_{i} Q A_{i}^{\mathrm{T}}  \tag{5.17}\\
& \Omega_{8}(Q) \triangleq \alpha Q^{2}+\alpha^{-1} \sum_{i=1}^{p} \alpha_{i}^{2} A_{i} A_{i}^{\mathrm{T}} \tag{5.18}
\end{align*}
$$

bound $U_{2}$.
Remark 5.3. Note that the term $A_{i} Q^{2} A_{i}^{T}$ appearing in $\Omega_{6}(\cdot)$ also appears in $\Omega_{4}(\cdot)$. Furthermore, both $\Omega_{4}(\cdot)$ and $\Omega_{0}(\cdot)$ involve a term proportional to $I_{n}$. Despite these similarities, neither bound $\Omega_{4}(\cdot)$ nor $\Omega_{0}(\cdot)$ is an overbound for the other.

Remark 5.4. The bound $\Omega_{7}(\cdot)$ given by (5.17) has the distinction that it is linear in $Q$. This bound was originally studied in $[23,27]$ for systems with multiplicative white noise and was shown to yield robust stability and performance in [33,35]. A similar bound was studied in [34].

Remark 5.5. Using (5.11) additional bounds can be developed. For example, by setting

$$
\begin{equation*}
Q_{1}=Q^{\frac{1}{2}}, \quad Q_{2}=Q^{\frac{1}{2}}, \tag{5.19}
\end{equation*}
$$

$\Omega_{5}(Q)$ becomes

$$
\begin{equation*}
\Omega_{9}(Q)=\alpha Q^{\frac{1}{2}}+\alpha^{-1} \sum_{i=1}^{p} \alpha_{i}^{2} A_{i} Q^{\frac{1}{2}} A_{i}^{T} . \tag{5.20}
\end{equation*}
$$

Remark 5.6. When $p=1$ and $\alpha$ is replaced by $\alpha \alpha_{i}, \Omega_{7}(\cdot)$ becomes

$$
\Omega_{7}^{\prime}(Q)=\alpha_{1}\left[\alpha Q+\alpha^{-1} A_{1} Q A_{1}^{T}\right] .
$$

Utilizing a sum of such terms with $\alpha_{i}=\delta_{i}$ can be used to bound the smaller rectangulius set $\mathcal{U}_{1}$. Similar remarks apply to $\Omega_{8}(\cdot), \Omega_{8}(\cdot)$, and $\Omega_{9}(\cdot)$.

### 5.3 Type III Bounds

We now consider bounds which exploit the structure of $\Delta A$ itself. It turns out that these bounds permit consideration of an uncertainty set $U$ which is larger than $U_{2}$. Specifically, define

$$
\begin{equation*}
u_{\mathrm{s}} \triangleq\left\{\Delta A \in \mathbb{R}^{\mathrm{n} \times n}: \quad \Delta A \doteq M N, \quad M M^{\mathrm{T}} \leq D, \quad N^{\mathrm{T}} N \leq E\right\} \tag{5.21}
\end{equation*}
$$

where $M \in \mathbb{R}^{n \times r}$ and $N \in \mathbb{R}^{r \times n}$ are uncertain matrices, $r$ is an arbitrary positive integer, and $D, E \in \mathbb{N}^{n}$ are given uncertainty bounds. The bound $\Omega_{10}(\cdot)$ for $U_{3}$ is given by the following result.

Proposition 5.4. Let $\alpha$ be an arbitrary positive constant. Then the function

$$
\begin{equation*}
\Omega_{10}(Q) \triangleq \alpha^{-1} D+\alpha Q E Q \tag{5.22}
\end{equation*}
$$

bounds $u_{s}$.
Proof. Note that

$$
\begin{aligned}
0 & \leq\left[\alpha^{-\frac{1}{2}} M-\alpha^{\frac{1}{2}} Q N^{\mathrm{T}}\right]\left[\alpha^{-\frac{1}{2}} M-\alpha^{\frac{1}{2}} Q N^{\mathrm{T}}\right]^{\mathrm{T}} \\
& =\alpha^{-1} M M^{\mathrm{T}}+\alpha Q N^{\mathrm{T}} N Q-\left[M N Q+Q(M N)^{\mathrm{T}}\right] \\
& \leq \alpha^{-1} D+\alpha Q E Q-\left(\Delta A Q+Q \Delta A^{\mathrm{T}}\right),
\end{aligned}
$$

which yields (3.1) with $\Omega(\cdot)=\Omega_{10}(\cdot)$ and $u=u_{3}$.
Remark 5.7. The bound $\Omega_{10}(\cdot)$ was developed independently in [29] for robust analysis and in $[25,28]$ for robust full-state feedback. Applications to fixed-order dynamic compensation are given in [36].

Remark 5.8. Without loss of generality we can set $\alpha=1$ in (5.22) by replacing $D$ and $E$ by $\alpha^{-1} D$ and $\alpha E$, respectively. Again for clarity we have chosen not to employ this normalization.

Note that $\Omega_{10}(\cdot)$ is an overbound for $\Omega_{\mathrm{B}}(\cdot)$ when $D$ and $E$ are chosen to satisfy

$$
\begin{equation*}
\sum_{i=1}^{p} \alpha_{i}^{2} A_{i} A_{i}^{\mathrm{T}} \leq D, \quad I_{n} \leq E . \tag{5.23}
\end{equation*}
$$

Hence, in this case $\Omega_{10}(\cdot)$ necessarily bounds $U_{2}$. As in the case of $\Omega_{3}(\cdot)$ overbounding $\Omega_{1}(\cdot)$, we should not be surprised to find that $\Omega_{10}(\cdot)$ actually bounds a set larger than $U_{2}$. Indeed, we now show that $U_{2}$ is actually a very special subset of $U_{3}$ when $D$ and $E$ defining $U_{2}$ satisfy (5.23).

Proposition 5.5. Suppose $D$ and $E$ satisfy (5.23). Then $U_{2}$ is a subset of $u_{3}$ and thus the bound $\Omega_{10}(\cdot)$ for $u_{3}$ also bounds $u_{2}$.

Proof. If $\Delta A \in U_{2}$ then $\Delta A=\sum_{i=1}^{p} \sigma_{i} A_{i}$, where $\sum_{i=1}^{p} \sigma_{i}^{2} / \alpha_{i}^{2} \leq 1$. Alternatively, we can write $\Delta A=M N$, where $r=p n$ and

$$
M=\left[\alpha_{1} A_{1} \ldots \alpha_{p} A_{p}\right], \quad N=\left[\begin{array}{c}
\left(\sigma_{1} / \alpha_{1}\right) I_{n}  \tag{5.24}\\
\vdots \\
\left(\sigma_{p} / \alpha_{p}\right) I_{n}
\end{array}\right]
$$

Note that with $D$ and $E$ satisfying (5.23) and $M$ and $N$ defined by (5.24) it follows that $M M^{T} \leq D$ and $N^{\mathrm{T}} N \leq E$. Thus $\Delta A \in U_{3}$.

The following result provides general conditions under which $\Omega_{10}(\cdot)$ bounds $\mathcal{U}_{2}$.
Proposition 5.6. Suppose $A_{i}=D_{i} E_{i}, i=1, \ldots, p$, where $D_{i} \in \mathbb{R}^{n \times n_{i}}$ and $E_{i} \in \mathbb{R}^{n_{i} \times n}$, and suppose that

$$
\begin{equation*}
\sum_{i=1}^{p} \alpha_{i}^{2} D_{i} D_{i}^{\mathrm{T}} \leq D, \quad \sum_{i=1}^{p} E_{i}^{\mathrm{T}} E_{i} \leq E . \tag{5.25}
\end{equation*}
$$

Then $U_{2}$ is a subset of $U_{3}$ and thus the bound $\Omega_{10}(\cdot)$ for $U_{3}$ also bounds $U_{2}$.
Proof. The result follows as in the proof Proposition 5.5.
Remark 5.9. When $p=1, D=\alpha_{1}^{2} D_{1} D_{1}^{T}$ and $E=E_{1}^{\mathrm{T}} E_{1}$, it is convenient to replace $\alpha$ by $\alpha \alpha_{1}$ so that $\Omega_{10}(\cdot)$ becomes

$$
\begin{equation*}
\Omega_{10}^{\prime}(Q)=\alpha_{1}\left[\alpha^{-1} D_{1} D_{1}^{\mathrm{T}}+\alpha Q E_{1}^{\mathrm{T}} E_{1} Q\right] \tag{5.26}
\end{equation*}
$$

In certain situations it is desirable to consider subsets of $u_{3}$ of special structure. For example, define

$$
U_{4} \triangleq\left\{\Delta A \in \mathbb{R}^{n \times n}: \quad \Delta A=D_{0} M N E_{0}, \quad D_{0} M M^{\mathrm{T}} D_{0}^{\mathrm{T}} \leq D, \quad E_{0}^{\mathrm{T}} N^{\mathrm{T}} N E_{0} \leq E\right\}
$$

where $D_{0} \in \mathbb{R}^{\boldsymbol{n} \times \boldsymbol{n}_{2}}$ and $E_{0} \in \mathbb{R}^{\boldsymbol{n}_{2} \times \boldsymbol{n}}$ are known matrices denoting the structure of the uncertainty, and $M \in \mathbb{R}^{n_{1} \times r}$ and $N \in \mathbb{R}^{r \times n_{2}}$ are uncertain matrices. Finer structure can be included within $U_{4}$ by replacing $D_{0} M N E_{0}$ by a sum of terms $D_{i} M_{i} N_{i} E_{i}$, where $D_{i}, E_{i}$ are known and $M_{i}, N_{i}$ are uncertain ([29]). Note, however, that even though $U_{4}$ is a proper subset of $U_{3}$, the form of the bound $\Omega_{10}(\cdot)$ does not change. Thus such refinements render the bound $\Omega_{10}(\cdot)$ conservative with respect to $U_{4}$ since the larger uncertainty set $U_{3}$ is actually being bounded.

## 6. Robust Stability and Performance via Modified Lyapunov Equations

We now combine the principal results of Sections 3, 4 and 5 to obtain a series of conditions guaranteeing robust stability and performance. In particular, we focus on bounds $\Omega_{1}, \Omega_{4}, \Omega_{8}, \Omega_{7}$, and $\Omega_{10}$. For simplicity we shall frequently assume that $V$ is positive definite so that (3.3) is satisfied. In this case it follows that the solution $Q$ of (3.2) is positive definite. Our first result is a corollary of Theorem 3.1 with $\Omega(\cdot)=\Omega_{1}(\cdot)$ and $U=U_{1}$.

Theorem 6.1. Let $V \in \mathbb{P}^{n}, \delta_{1}, \ldots, \delta_{p}>0$, and suppose there exists $Q \in \mathbb{P}^{n}$ satisfying

$$
\begin{equation*}
0=A Q+Q A^{\mathrm{T}}+\sum_{i=1}^{p} \delta_{i}\left|A_{i} Q+Q A_{i}^{\mathrm{T}}\right|+V . \tag{MLE1}
\end{equation*}
$$

Then $A+\Delta A$ is asymptotically stable for all $\Delta A \in U_{1}$, and

$$
\begin{gather*}
J_{\mathrm{S}}\left(U_{1}\right) \leq \operatorname{tr} Q R,  \tag{6.1}\\
J_{\mathrm{D}}\left(U_{1}\right) \leq \lambda_{\max }(Q R) . \tag{6.2}
\end{gather*}
$$

For the next result define

$$
\begin{equation*}
A_{\alpha} \triangleq A+\frac{\alpha}{2} I_{n} \tag{6.3}
\end{equation*}
$$

and

$$
\begin{equation*}
\gamma_{i} \triangleq \alpha_{i}^{2} / \alpha, \quad i=1, \ldots, p \tag{6.4}
\end{equation*}
$$

Setting $\Omega(\cdot)=\Omega_{4}(\cdot), \Omega_{6}(\cdot), \Omega_{7}(\cdot)$ and $U=U_{2}$ yields the following corollary of Theorem 3.1.
Theorem 6.2. Let $V \in \mathbb{P}^{n}, \alpha, \alpha_{1}, \ldots, \alpha_{p}>0$, and suppose there exists $Q \in \mathbb{P}^{n}$ satisfying either

$$
\begin{equation*}
0=A Q+Q A^{T}+\sum_{i=1}^{p} \gamma_{i}\left(A_{i} Q+Q A_{i}^{T}\right)^{2}+\frac{a}{4} I_{n}+V \tag{MLE2}
\end{equation*}
$$

$$
\begin{equation*}
0=A Q+Q A^{\mathrm{T}}+\sum_{i=1}^{p} \gamma_{i} A_{i} Q^{2} A_{i}^{\mathrm{T}}+\alpha I_{n}+V \tag{MLE3}
\end{equation*}
$$

or

$$
\begin{equation*}
0=A_{\alpha} Q+Q A_{\alpha}^{\mathrm{T}}+\sum_{i=1}^{p} \gamma_{i} A_{i} Q A_{i}^{\mathrm{T}}+V . \tag{MLE4}
\end{equation*}
$$

Then $A+\Delta A$ is asymptotically stable for all $\Delta A \in \mathcal{U}_{2}$, and

$$
\begin{gather*}
J_{\mathrm{S}}\left(U_{2}\right) \leq \operatorname{tr} Q R,  \tag{6.5}\\
J_{\mathrm{D}}\left(U_{2}\right) \leq \lambda_{\max }(Q R) . \tag{6.6}
\end{gather*}
$$

Next we set $\Omega(\cdot)=\Omega_{10}(\cdot)$ and $U=U_{3}$.
Theorem 6.3. Let $V \in \mathbb{P}^{n}, D \in \mathbb{N}^{n}$, and $E \in \mathbb{N}^{n}$, and suppose there exists $Q \in \mathbb{P}^{n}$ satisfying

$$
\begin{equation*}
0=A Q+Q A^{T}+\alpha Q E Q+\alpha^{-1} D+V . \tag{MLE5}
\end{equation*}
$$

Then $A+\Delta A$ is asymptotically stable for all $\Delta A \in U_{3}$, and

$$
\begin{gather*}
J_{\mathrm{S}}\left(U_{3}\right) \leq \operatorname{tr} Q R,  \tag{6.7}\\
J_{\mathrm{D}}\left(U_{3}\right) \leq \lambda_{\max }(Q R) . \tag{6.8}
\end{gather*}
$$

Additional sufficient conditions can be obtained by considering "mixed" bounds. That is, one can construct modified Lyapunov equations by combining two or more different bounds. Although mixed bounds will not be considered further in this paper, we present one ouch result for illustrative purposes.

Theorem 6.4. Let $V \in \mathbb{P}^{n}, \delta_{1}, \ldots, \delta_{p}>0, D \in \mathbb{N}^{n}$, and $E \in \mathbb{N}^{n}$, and suppose there exists $Q \in \mathbb{P}^{n}$ satisfying

$$
\begin{equation*}
0=A Q+Q A^{\mathbf{T}}+\sum_{i=1}^{p} \delta_{i}\left|A_{i} Q+Q A_{i}^{\mathrm{T}}\right|+\alpha Q E Q+\alpha^{-1} D+V \tag{MLE1,5}
\end{equation*}
$$

Then $A+\Delta A$ is asymptotically stable for all $\Delta A \in U_{1}+U_{3}$, and

$$
\begin{gather*}
J_{\mathrm{s}}\left(U_{1}+U_{3}\right) \leq \operatorname{tr} Q R  \tag{6.9}\\
J_{\mathrm{D}}\left(U_{1}+U_{3}\right) \leq \lambda_{\max }(Q R) . \tag{6.10}
\end{gather*}
$$

As noted previously, the bound $\Lambda(\cdot)$ can readily be constructed by replacing $\Delta A$ by $\Delta A^{T}$ in the definitions of $\Omega_{1}(\cdot)$ through $\Omega_{10}(\cdot)$. Denote these bounds by $\Lambda_{1}(\cdot)$ through $\Lambda_{10}(\cdot)$, respectively. For illustration we state the dual of Theorem 6.1 involving $\Lambda_{1}(\cdot)$. The dual versions of MLE1 through MLE5 will be denoted by MLED1 through MLED5.

Theorem 6.5. Let $R \in \mathbb{P}^{n}, \delta_{1}, \ldots, \delta_{p}>0$ and suppose there exists $P \in \mathbb{P}^{n}$ satisfying

$$
\begin{equation*}
0=A^{\mathrm{T}} P+P A+\sum_{i=1}^{p} \delta_{i}\left|A_{i}^{\mathrm{T}} P+P A_{i}\right|+R \tag{MLEDI}
\end{equation*}
$$

Then $A+\Delta A$ is asymptotically stable for all $\Delta A \in U_{1}$, and

$$
\begin{gather*}
J_{\mathrm{S}}\left(U_{1}\right) \leq \operatorname{tr} P V  \tag{6.11}\\
\hat{J}_{\mathrm{D}}\left(U_{1}\right) \leq \lambda_{\max }(P V) \tag{6.12}
\end{gather*}
$$

It is reasonable to expect that the sufficient conditions given by Theorems 3.1 and 4.1 are generally different. For example, the modified Lyapunov equations and their duals need not both possess a solution, while the bounds $\operatorname{tr} Q R$ and $\operatorname{tr} P V$ need not be equal. An exception is the case in which $\Omega(\cdot)=\Omega_{7}(\cdot)$ and $\Lambda(\cdot)=\Lambda_{7}(\cdot)$. Note that the dual of MLE4 is given by

$$
\begin{equation*}
0=A_{\alpha}^{\mathrm{T}} P+P A_{\alpha}+\sum_{i=1}^{p} \gamma_{i} A_{i}^{\mathrm{T}} P A_{i}+V . \tag{MLED4}
\end{equation*}
$$

Proposition 6.1. Let $\alpha_{1}, \ldots, \alpha_{p}>0$ and assume there exist $Q, P \in \mathbb{N}^{n}$ satisfying MLE4 and MLED4, respectively. Then

$$
\begin{equation*}
\operatorname{tr} Q R=\operatorname{tr} P V \tag{6.13}
\end{equation*}
$$

Proof. Note that

$$
\begin{aligned}
\operatorname{tr} Q R & =-\operatorname{tr} Q\left(A_{\alpha}^{\mathrm{T}} P+P A_{\alpha}+\sum_{i=1}^{p} \gamma_{i} A_{i}^{\mathrm{T}} P A_{i}\right) \\
& =-\operatorname{tr} P\left(A_{\alpha} Q+Q A_{\alpha}^{\mathrm{T}}+\sum_{i=1}^{p} \gamma_{i} A_{i} Q A_{i}^{\mathrm{T}}\right) \\
& =\operatorname{tr} P V . \square
\end{aligned}
$$

Remark 6.1. By setting $\Omega(\cdot)=\Omega_{7}(\cdot)$ and $\Lambda(\cdot)=\Lambda_{7}(\cdot)$ it follows from (4.14) that

$$
\begin{equation*}
\operatorname{tr} Q_{0}\left(\alpha P+\sum_{i=1}^{p} \gamma_{i} A_{i}^{\mathrm{T}} P A_{i}\right)=\operatorname{tr} P_{0}\left(\alpha Q+\sum_{i=1}^{p} \gamma_{i} A_{i} Q A_{i}^{\mathrm{T}}\right) \tag{6.14}
\end{equation*}
$$

## 7. Existence, Uniqueness, and Monotonicity of Solutions to the Modified Lyapunov Equations

It is important to stress that the sufficient conditions for robustness given by Theorems 6.1-6.5 assume only that there exist nonnegative-definite solutions $Q, P$ satisfying the modified Lyapunov equations. Indeed, no explicit assumptions on the problem data $A, V, R$, and $U$ were utilized for assuring robust stability and performance. In applying Theorems 6.1-6.5 to specific problems it thus suffices to show that a nonnegative-definite solution $Q$ exists in order to obtain robust stability, while, for robust performance, the bounds (6.1), (6.2), (6.5)-(6.8) require explicit knowledge of Q. Thus, any computational method which yields a nonnegative-definite solution will suffice to guarantee both robust stability and performance.

Before considering the numerical solution of the modified Lyapunov equations, several relevant issues require discussion. For example, before seeking to compute solutions to MLE1-MLE5 it would be desirable to determine a priori whether these equations actually possess nonnegativedefinite solutions. For example, it may be useful to obtain sufficient and/or necessary conditions for the existence of nonnegative-definite solutions. Thus, if the sufficient conditions are satisfied then existence (and hence robustness) is assured, while if the necessary conditions are not satisfied then existence is ruled out. If, on the other hand, either the sufficient conditions are not satisfied or the necessary conditions are satisfied, then nothing can be surmised. Finally, such conditions need to be easily verifiable and reasonably nonconservative since otherwise it would be more prudent to attempt to numerically solve the modified Lyapunov equation itself.

It is quite possible that at least some of the modified Lyapunov equations possess multiple nonnegative-definite solutions. In this case one may seek the minimal solution (i.e., the smallest with respect to the nonnegative-definite matrix ordering) in order to minimize the performance bounds. If multiple solutions exist, none of which is minimal, then the best bound would depend upon the matrix $R$.

Since the matrix $Q$ determines the performance bound, it is reasonable to expect that $Q$ is monotonic in $U$. That is, if $U$ decreases in size, then the solution $Q$ is more likely to exist while decreasing in the nonnegative-definite matrix ordering. For example, consider $U_{1}^{\prime}$ characterized by $\delta_{i}^{\prime}$, where $\delta_{i}^{\prime} \leq \delta_{i}, i=1, \ldots, p$. Then one might expect $Q^{\prime} \leq Q$ where $Q^{\prime}$ is the solution to MLE1 with $\delta_{i}$ replaced by $\delta_{i}^{\prime}$. Finally, monotonicity with respect to $V$ should also be expected.

Because of linearity, the analysis of bound $\Omega_{7}(\cdot)$ is simplest and it is possible to obtain necessary and sufficient conditions for the existence of solutions to MLE4. The basic tool required is the Kronecker matrix algebra ([42]). For convenience, define

$$
\begin{equation*}
A \triangleq A_{\alpha} \oplus A_{\alpha}+\sum_{i=1}^{p} \gamma_{i} A_{i} \otimes A_{i} \tag{7.1}
\end{equation*}
$$

Proposition 7.1. If $V \in \mathbb{N}^{n}$ and $\mathbb{A}$ is asymptotically stable, then there exists a unique $n \times n Q$ satisfying MLE4, and $Q \geq 0$. Conversely, if for all $V \in \mathbb{N}^{n}$ there exists $Q \geq 0$ satisfying MLE4, then $\mathcal{A}$ is asymptotically stable.

Proof. Since MLE4 is equivalent to

$$
\begin{equation*}
Q=-\operatorname{vec}^{-1}\left[A^{-1} \operatorname{vec} V\right] \tag{7.2}
\end{equation*}
$$

existence and uniqueness hold. Here, vec and vec ${ }^{-1}$ denote the column-stacking operation ([42]) and its inverse. To prove that $Q$ is nonnegative definite, we rewrite (7.2) as

$$
\begin{equation*}
Q=\int_{0}^{\infty} \operatorname{vec}^{-1}\left[e^{\wedge t} \mathrm{vec} V\right] d t \tag{7.3}
\end{equation*}
$$

and show that the integrand is nonnegative-definite for all $t \in[0, \infty)$. [Note that the following argument does not require that $\mathcal{A}$ be stable.] Using the exponential product formula, the exponential in (7.3) can be written as

$$
\begin{equation*}
e^{A t}=\lim _{k \rightarrow \infty}\left\{\exp \left[\frac{1}{k}\left(A_{\alpha} \oplus A_{\alpha}\right) t\right] \exp \left[\frac{1}{k} \sum_{i=1}^{p} \gamma_{i}\left(A_{i} \otimes A_{i}\right) t\right]\right\}^{k} \tag{7.4}
\end{equation*}
$$

For convenience, let $S$ and $N$ be $r \times r$ matrices with $N \geq 0$. Since (see [42])

$$
\begin{equation*}
\operatorname{vec}^{-1}[(S \otimes S) \operatorname{vec} N]=S N S^{\mathrm{T}} \geq 0 \tag{7.5}
\end{equation*}
$$

and

$$
\begin{equation*}
\left(S^{k} \otimes S^{k}\right)(S \otimes S)=S^{k+1} \otimes S^{k+1} \tag{7.6}
\end{equation*}
$$

it follows that

$$
\begin{equation*}
\operatorname{vec}^{-1}\left[e^{s \otimes S} \operatorname{vec} N\right]=\sum_{k=0}^{\infty}(k!)^{-1} S^{k} N S^{k T} \geq 0 \tag{7.7}
\end{equation*}
$$

Furthermore,

$$
\begin{equation*}
\operatorname{vec}^{-1}\left[e^{s \oplus S} \operatorname{vec} N\right]=\operatorname{vec}^{-1}\left[\left(e^{s} \otimes e^{s}\right) \operatorname{vec} N\right]=e^{s} N e^{s^{T}} \geq 0 \tag{7.8}
\end{equation*}
$$

Applying (7.7) and (7.8) alternately with (7.4) and using induction on $k$ it follows that the integrand of (7.3) is nonnegative definite. To prove the converse, note that it follows from MLE4 that $Q$ satisfies

$$
\begin{equation*}
Q=\operatorname{vec}^{-1}\left[e^{A t} \operatorname{vec} Q\right]+\int_{0}^{t} \operatorname{vec}^{-1}\left[e^{A} \cdot \operatorname{vec} V\right] d s, \quad t \in[0, \infty) \tag{7.9}
\end{equation*}
$$

Since the integral term on the right hand side of (7.9) is nonnegative definite, is bounded from above by $Q$, and $V \in \mathbb{N}^{n}$ is arbitrary, it follows that $A$ is asymptotically stable.

We now show that if $\mathcal{A}$ is asymptotically stable then actually $A_{\alpha}$ is asymptotically stable. This shows that the assumption that $A$ is asymptotically stable is consistent with the original hypothesis that $A$ is asymptotically stable.

Proposition 7.2. Assume $A$ is asymptotically stable, let $\alpha_{i}^{\prime} \in\left[0, \alpha_{i}\right], i=1, \ldots, p$, and define

$$
R^{\prime} \triangleq A_{\alpha} \oplus A_{\alpha}+\sum_{i=1}^{p}\left(\alpha_{i}^{\prime 2} / \alpha\right) A_{i} \otimes A_{i} .
$$

Then $A^{\prime}$ is also asymptotically stable. In particular, $A_{\alpha}$ and $A$ are asymptotically stable.
Proof. Let $V \in \mathbb{N}^{n}$ and let $Q$ be the unique, nonnegative-definite solution of MLE4. Equivalently, $Q$ satisfies

$$
0=A_{\alpha} Q+Q A_{\alpha}^{\mathrm{T}}+\sum_{i=1}^{p}\left(\alpha_{i}^{\prime 2} / \alpha\right) A_{i} Q A_{i}^{\mathrm{T}}+V^{\prime}
$$

where

$$
V^{\prime} \triangleq \sum_{i=1}^{p} \alpha^{-1}\left(\alpha_{i}^{2}-\alpha_{i}^{\prime 2}\right) A_{i} Q A_{i}^{\mathrm{T}}+V
$$

Since $V^{\prime} \in \mathbb{N}^{\boldsymbol{n}}$, the stability of $\boldsymbol{R}^{\prime}$ now follows from the converse of Proposition 7.1. Finally, if $V \in \mathbb{P}^{n}$ then $\sum_{i=1}^{p}\left(\alpha_{i}^{\prime 2} / \alpha\right) A_{i} Q A_{i}^{T}+V^{\prime}$ is positive definite and it follows from Lemma 12.2 of [39] that $A_{\alpha}$, and hence $A$, is asymptotically stable.

Hence it follows from Proposition 7.2 that a necessary condition for $A$ to be asymptotically stable is that

$$
\begin{equation*}
\alpha<2 \max _{i=1, \ldots, n} \operatorname{Re} \lambda_{i}(A) . \tag{7.10}
\end{equation*}
$$

We now have the following monotonicity result.
Proposition 7.3. Let $u_{2}^{\prime} \subset U_{2}$, where $u_{2}^{\prime}$ is defined as in (5.7) with $\alpha_{i}$ replaced by $\alpha_{i}^{\prime} \in$ $\left[0, \alpha_{i}\right], i=1, \ldots, p$. Furthermore, let $V \in \mathbb{P}^{n}$, assume $A$ is asymptotically stable, and let $Q \in \mathbb{N}^{n}$
satisfy MLE4. Then there exists $Q^{\prime} \in \mathbb{N}^{n}$ satisfying

$$
\begin{equation*}
0=A_{\alpha} Q^{\prime}+Q^{\prime} A_{\alpha}^{\mathbf{T}}+\sum_{i=1}^{p}\left(\alpha_{i}^{\prime 2} / \alpha\right) A_{i} Q^{\prime} A_{i}^{\mathrm{T}}+V \tag{7.11}
\end{equation*}
$$

and, furthermore,

$$
\begin{equation*}
Q^{\prime} \leq \boldsymbol{Q} \tag{7.12}
\end{equation*}
$$

Consequently,

$$
\begin{gather*}
\operatorname{tr} Q^{\prime} R \leq \operatorname{tr} Q R,  \tag{7.13}\\
\lambda_{\max }\left(Q^{\prime} R\right) \leq \lambda_{\max }(Q R) . \tag{7.14}
\end{gather*}
$$

Proof. Subtracting (7.11) from MLE4 yields

$$
0=A_{\alpha}\left(Q-Q^{\prime}\right)+\left(Q-Q^{\prime}\right) A_{\alpha}^{\mathbf{T}}+\sum_{i=1}^{p}\left(\alpha_{i}^{\prime 2} / \alpha\right) A_{i}\left(Q-Q^{\prime}\right) A_{i}^{\mathbf{T}}+V^{\prime}
$$

where $V^{\prime}$ is defined in the proof of Proposition 7.2. Since, by the converse of Proposition 7.1, $R^{\prime}$ is asymptotically stable, $Q-Q^{\prime} \geq 0$, which yields (7.12) and thus (7.13) and (7.14).

Returning now to the existence question, Proposition 7.1 shows that a solution to MLE4 exists so long as $\alpha_{1}, \ldots \alpha_{p}$ are sufficiently small that $A$ remains stable for some $\alpha>0$. To this end one can treat this as a stability perturbation problem and apply results from [3]. Within our modified Lyapunov equation approach we have the following related result. For this and the following result let $\|\cdot\|$ denote an arbitrary vector norm and the corresponding induced matrix norm.

Proposition 7.4. If

$$
\begin{equation*}
\left\|(A \oplus A)^{-1}\left(\alpha I_{n^{2}}+\alpha^{-1} \sum_{i=1}^{p} \alpha_{i}^{2} A_{i} \otimes A_{i}\right)\right\|<1 \tag{7.15}
\end{equation*}
$$

then for all $V \in \mathbb{N}^{n}$ there exists $Q \in \mathbb{N}^{n}$ satisfying MLE4 and hence $A$ is asymptotically stable.
Proof. Define $\left\{Q_{k}\right\}_{k=0}^{\infty}$ where $Q_{0}$ satisfies (3.14) and $Q_{k+1}$ satisfies

$$
0=A Q_{k+1}+Q_{k+1} A^{\mathrm{T}}+\Omega_{7}\left(Q_{k}\right)+V
$$

Note that $Q_{k} \geq 0, k=1,2, \ldots$ Hence it follows that

$$
\operatorname{vec} Q_{k+1}-\operatorname{vec} Q_{k}=-(A \oplus A)^{-1}\left[\operatorname{vec} \Omega_{7}\left(Q_{k}\right)-\operatorname{vec} \Omega_{7}\left(Q_{k-1}\right)\right]
$$

and thus
$\left\|\operatorname{vec} Q_{k+1}-\operatorname{vec} Q_{k}\right\| \leq\left\|(A \oplus A)^{-1}\left(\alpha I_{n^{2}}+\alpha^{-1} \sum_{i=1}^{p} \alpha_{i}^{2} A_{i} \otimes A_{i}\right)\right\|\left\|\operatorname{vec} Q_{k}-\operatorname{vec} Q_{k-1}\right\|$.
Using (7.15) it follows that $Q \triangleq \lim _{k \rightarrow \infty} Q_{k}$ exists. Thus $Q \geq 0$ and satisfies MLE4. Finally, by the converse of Proposition 7.1, $\AA$ is asymptotically stable.

Since MLE5 is nonlinear a slightly different approach is required for existence. For the following result let $\kappa, \beta>0$ satisfy

$$
\begin{equation*}
\left\|e^{A t}\right\| \leq k e^{-\beta t}, \quad t \geq 0 \tag{7.16}
\end{equation*}
$$

where $\|\cdot\|$ denotes an arbitrary submultiplicative matrix norm, and define $\rho \triangleq 2 \beta / \kappa^{2}$.
Proposition 7.5. Suppose $V \in \mathbb{N}^{n}$ and

$$
\begin{equation*}
4 \alpha\|E\|\left\|\alpha^{-1} D+V\right\|<\rho^{2} . \tag{7.17}
\end{equation*}
$$

Then there exists $Q \in \mathbb{N}^{n}$ satisfying MLE5.
Proof. Consider the sequence $\left\{Q_{k}\right\}_{k=0}^{\infty}$ where $Q_{0}$ satisfies (3.14) and $Q_{k+1}$ is given by

$$
0=A Q_{k+1}+Q_{k+1} A^{T}+\alpha Q_{k} E Q_{k}+\alpha^{-1} D+V
$$

Clearly, $Q_{k} \geq 0, k=0,1, \ldots$. Next we have

$$
\begin{equation*}
Q_{k+1}=\int_{0}^{\infty} e^{A t}\left[\alpha Q_{k} E Q_{k}+\alpha^{-1} D+V\right] e^{A^{T}} d t \tag{7.18}
\end{equation*}
$$

which yields

$$
\begin{equation*}
1 \cdot{ }_{+1}\left\|\leq \alpha \rho^{-1}\right\| E\| \| Q_{k}\left\|^{2}+\rho^{-1}\right\| \alpha^{-1} D+V \| . \tag{7.19}
\end{equation*}
$$

Similarly, we obtain

$$
\left\|Q_{0}\right\| \leq \rho^{-1}\left\|\alpha^{-1} D+V\right\|
$$

Now suppose that

$$
\left\|Q_{k}\right\| \leq 2 \rho^{-1}\left\|\alpha^{-1} D+V\right\| .
$$

Then (7.17) and (7.19) imply

$$
\begin{aligned}
\left\|Q_{k+1}\right\| & \leq \alpha \rho^{-1}\|E\|\left[2 \rho^{-1}\left\|\alpha^{-1} D+V\right\|\right]^{2}+\rho^{-1}\left\|\alpha^{-1} D+V\right\| \\
& <2 \rho^{-1}\left\|\alpha^{-1} D+V\right\| .
\end{aligned}
$$

Thus $\left\|Q_{k}\right\| \leq 2 \rho^{-1}\left\|\alpha^{-1} D+V\right\|, k=0,1, \ldots$. Next, (7.18) yields

$$
\begin{aligned}
Q_{k+1}-Q_{k} & =\alpha \int_{0}^{\infty} e^{A t}\left[Q_{k} E Q_{k}-Q_{k-1} E Q_{k-1}\right] e^{A^{T} t} d t \\
& =\alpha \int_{0}^{\infty} e^{A t}\left[Q_{k} E\left(Q_{k}-Q_{k-1}\right)+\left(Q_{k}-Q_{k-1}\right) E Q_{k-1}\right] e^{A^{T} t} d t
\end{aligned}
$$

and thus

$$
\begin{aligned}
\left\|Q_{k+1}-Q_{k}\right\| & \leq \alpha \rho^{-1}\|E\|\left(\left\|Q_{k}\right\|+\left\|Q_{k-1}\right\|\right)\left\|Q_{k}-Q_{k-1}\right\| \\
& \leq 4 \alpha \rho^{-2}\|E\|\left\|\alpha^{-1} D+V\right\|\left\|Q_{k}-Q_{k-1}\right\| \\
& \leq \varepsilon\left\|Q_{k}-Q_{k+1}\right\|
\end{aligned}
$$

where $\varepsilon \triangleq 4 \alpha \rho^{-2}\|E\|\left\|\alpha^{-1} D+V\right\|$. Since by (7.16) $\varepsilon<1, \lim _{k \rightarrow \infty} Q_{k}$ exists, is nonnegative definite, and satisfies MLE5.

## 8. Additional Upper Bounds via Recursive Substitution

In this section we obtain additional upper bounds for $J_{S}(u)$ and $J_{D}(u)$ by utilizing a recursive substitution technique. The main idea involves rewriting (2.7) as

$$
\begin{equation*}
Q_{\Delta A}=-\operatorname{vec}^{-1}\left\{(A \oplus A)^{-1}(\Delta A \oplus \Delta A) \operatorname{vec} Q_{\Delta A}\right\}+Q_{0} \tag{8.1}
\end{equation*}
$$

and substituting this expression into the terms $\triangle A Q_{\Delta A}+Q_{\Delta A} \Delta A^{T}$ appearing in (2.7). This technique yields an equation which is, as expected, equivalent to (2.7) but which permits the development of additional bounds. As will be seen, the ability to develop new bounds exploits the fact that the substitution technique leads to terms which are quadratic in $\Delta A$. We begin the development with the following technical result which does not require the assumption that $A$ is asymptotically stable.

Proposition 8.1. Suppose $A \oplus A$ is invertible and let $\Delta A \in \mathbb{R}^{n \times n}$. If $Q_{\Delta A}$ satisfies (2.7) then $Q_{\Delta A}$ also satisfies

$$
\begin{gather*}
0=A Q_{\Delta A}+Q_{\Delta A} A^{\mathrm{T}}-\operatorname{vec}^{-1}\left[(\Delta A \oplus \Delta A)(A \oplus A)^{-1}(\Delta A \oplus \Delta A) \operatorname{vec} Q_{\Delta A}\right. \\
\left.+(\Delta A \oplus \Delta A)(A \oplus A)^{-1} \operatorname{vec} V\right]+V \tag{8.2}
\end{gather*}
$$

Conversely, if $Q_{\Delta A}$ satisfies (8.2) and $(A-\Delta A) \oplus(A-\Delta A)$ is invertible, then $Q_{\Delta A}$ also satisfies (2.7).

Proof. To obtain (8.2) substitute (8.1) into (2.7) as noted above. Conversely, adding the zero term $(\Delta A \oplus \Delta A)(A \oplus A)^{-1}(A \oplus A)$ vec $Q_{\Delta A}-(\Delta A \oplus \Delta A) \operatorname{vec} Q_{\Delta A}$ to (8.2), it follows that (8.2) can be written as

$$
0=[(A-\Delta A) \oplus(A-\Delta A)](A \oplus A)^{-1}\left[(A+\Delta A) \oplus(A+\Delta A) \operatorname{vec} Q_{\Delta A}+\operatorname{vec} V\right]
$$

which, under the invertibility assumption, implies that $Q_{\Delta A}$ satisfies (2.7).
The following result is analogous to Theorem 3.1.
Theorem 8.1. Suppose $U$ is symmetric, let $\Omega_{0} \in \mathbb{N}^{n}$ satisfy

$$
\begin{equation*}
\Delta A Q_{0}+Q_{0} \Delta A^{T} \leq \Omega_{0}, \quad \Delta A \in U \tag{8.3}
\end{equation*}
$$

where $Q_{0}$ satisfies (3.14), let $\hat{\Omega}: \mathbb{N}^{\boldsymbol{n}} \rightarrow \mathbb{N}^{\boldsymbol{n}}$ satisfy

$$
\begin{equation*}
-\operatorname{vec}^{-1}\left[(\Delta A \oplus \Delta A)(A \oplus A)^{-1}(\Delta A \oplus \Delta A) \text { vec } Q\right] \leq \hat{\Omega}(Q), \quad \Delta A \in U, \quad Q \in \mathbb{N}^{n} \tag{8.4}
\end{equation*}
$$

and suppose there exists $Q \in \mathbb{N}^{n}$ satisfying

$$
\begin{equation*}
0=A Q+Q A^{\mathbf{T}}+\hat{\Omega}(Q)+\Omega_{0}+V \tag{8.5}
\end{equation*}
$$

Then

$$
\begin{align*}
(A+\Delta A,\{\hat{\Omega}(Q)+ & \left.\left.\operatorname{vec}^{-1}\left[(\Delta A \oplus \Delta A)(A \oplus A)^{-1}(\Delta A \oplus \Delta A) \operatorname{vec} Q\right]+\Omega_{0}+V\right\}^{\frac{1}{2}}\right)  \tag{8.6}\\
& \text { is stabilizable, } \quad \Delta A \in U,
\end{align*}
$$

if and only if

$$
\begin{equation*}
A+\Delta A \text { is asymptotically stable, } \quad \Delta A \in U \tag{8.7}
\end{equation*}
$$

In this case,

$$
\begin{equation*}
Q_{\Delta A} \leq Q, \quad \Delta A \in U \tag{8.8}
\end{equation*}
$$

where $Q_{\Delta A}$ satisfies (2.7), and

$$
\begin{gather*}
J_{\mathrm{S}}(U) \leq \operatorname{tr} Q R,  \tag{8.9}\\
J_{\mathrm{D}}(U) \leq \lambda_{\max }(Q R) . \tag{8.10}
\end{gather*}
$$

Proof. The equivalence of (8.6) and (8.7) follows from (8.5) as in the proof of Theorem 3.1. Next (8.8) follows by comparing (8.5) and (8.2) while using (8.3) and (8.4). Since $U$ is assumed to be symmetric, it follows that $A-\Delta A$ is asymptotically stable, $\Delta A \in U$, and hence $(A-\Delta A) \oplus(A-\Delta A)$ is invertible, $\Delta A \in U$. Thus, the converse of Proposition 8.1 implies that $Q_{\Delta A}$ satisfying (8.2) also satisfies (2.7). Thus, the bound (8.8) can be used to obtain (8.9) and (8.10).

The principal difference between (8.4) and (3.1) is that $\Delta A$ appears linearly in (3.1) while it appears quadratically in (8.4). By exploiting this structure we can obtain new bounds for $Q_{\Delta A}$.

To simplify matters, we now consider the bound in (8.4) in two special cases. In the first case we set $U=U_{1}$ and $p=1$ so that $\Delta A=\sigma_{1} A_{1},\left|\sigma_{1}\right| \leq \delta_{1}$. In this case (8.4) becomes

$$
\begin{equation*}
-\sigma_{1}^{2} \mathrm{vec}^{-1}\left[\left(A_{1} \oplus A_{1}\right)(A \oplus A)^{-1}\left(A_{1} \oplus A_{1}\right) \operatorname{vec} Q\right] \leq \hat{\Omega}(Q), \quad\left|\sigma_{1}\right| \leq \delta_{1}, \quad Q \in \mathbb{N}^{n} \tag{8.11}
\end{equation*}
$$

One choice of $\hat{\Omega}(\cdot)$ which inmediately suggests itself can be obtained by defining the matrix function $|\cdot|_{+}$on the set of aymmetric matrices by

$$
\begin{equation*}
|S|_{+} \triangleq \frac{1}{2}(S+|S|), \tag{8.12}
\end{equation*}
$$

which effectively replaces the negative eigenvalues of $S$ by zeros. We shall thus utilize the fact that

$$
\begin{equation*}
\sigma_{1}^{2} S \leq \delta_{1}^{2}|S|_{+}, \quad\left|\sigma_{1}\right| \leq \delta_{1}, \tag{8.13}
\end{equation*}
$$

for all symmetric $S$.
Corollary 8.1. Let $V \in \mathbb{P}^{n}, u=U_{1}, p=1$, let $\Omega_{0} \in \mathbb{N}^{n}$ satisfy (8.3), and suppose there exists $Q \in \mathbb{N}^{n}$ aatisfying

$$
\begin{equation*}
0=A Q+Q A^{\mathrm{T}}+\delta_{1}^{2}\left|-\mathrm{vec}^{-1}\left[\left(A_{1} \oplus A_{1}\right)(A \oplus A)^{-1}\left(A_{1} \oplus A_{1}\right) \mathrm{vec} Q\right]\right|_{+}+\Omega_{0}+V \tag{8.14}
\end{equation*}
$$

Then (8.7)-(8.10) are satisfied.
For the next specialization we shall assume that

$$
\begin{equation*}
(\Delta A) A=A(\Delta A), \quad \Delta A \in U \tag{8.15}
\end{equation*}
$$

which holds, for example, for modal systems with frequency uncertainty (see Section 10). It thus follows that $(A \oplus A)^{-1}(\Delta A \oplus \Delta A)=(\Delta A \oplus \Delta A)(A \oplus A)^{-1}$ and thus (8.4) can be rewritten as

$$
\begin{equation*}
\Delta A^{2} \hat{Q}+2 \Delta A \hat{Q} \Delta A^{T}+\hat{Q} \Delta A^{2 T} \leq \hat{\Omega}(Q), \quad \Delta A \in \mathbb{N}, \quad Q \in \mathbb{N}^{n} \tag{8.16}
\end{equation*}
$$

where $\hat{\boldsymbol{Q}} \in \mathbb{N}^{\boldsymbol{n}}$ satisfies

$$
\begin{equation*}
0=A \hat{Q}+\hat{Q} A^{\mathrm{T}}+Q . \tag{8.17}
\end{equation*}
$$

Assuming in addition to (8.15) that $\Delta A=\sigma_{1} A_{1},\left|\sigma_{1}\right| \leq \delta_{1}$, (8.14) becomes

$$
\begin{equation*}
0=A Q+Q A^{\mathrm{T}}+\delta_{1}^{2}\left|A_{1}^{2} \hat{Q}+2 A_{1} \hat{Q} A_{1}^{\mathrm{T}}+\hat{Q} A_{1}^{2 \mathrm{~T}}\right|_{+}+\Omega_{0}+V . \tag{8.18}
\end{equation*}
$$

Remark 8.1. It is interesting to note that the left hand side of (8.16) is essentially of the same form as $\Omega_{7}(\cdot)$. Specifically, the term $\Delta A^{2} \hat{Q}+\hat{Q} \Delta A^{2 T}$ is analogous to $\alpha Q$ while $2 \Delta A \hat{Q} \Delta A^{T}$ is similar to $\alpha^{-1} \sum_{i=1}^{p} \alpha_{i}^{2} A_{i} Q A_{i}^{T}$ appearing in $\Omega_{7}(\cdot)$. The form of the left hand side of (8.16) is also of interest since it is similar to terms which arise from a multiplicative white noise model with a Stratonovich correction. Specifically, while the term $\triangle A Q \Delta A^{T}$ arises from an Ito model ([33]), the new terms $\frac{1}{2} \Delta A^{2}$ can be viewed as a correction to the nominal $A$ due to the Stratonovich interpretation of stochastic integration ([43]). These terms have interesting ramifications in designing robust controllers for fiexible structures ([23]).

## 9. An Alternative Approach Yielding Upper and Lower Bounds

In this section we develop a variation on the results of Section 3 which has the additional benefit of yielding both upper and lower performance bounds. The basic approach was suggested by results obtained in [44]. To simplify the presentation we assume that if $\Delta A \in U$ then $-\Delta A \in U$. This symmetry assumption of course holds for all of the uncertainty sets considered in previous sections. The underlying idea involves bounding the deviation of $Q_{\Delta A}$ from $Q_{0}$ rather than bounding $Q_{\Delta A}$ directly.

Theorem 9.1. Let $\Omega_{0} \in \mathbb{N}^{n}$ satisfy

$$
\begin{equation*}
\Delta A Q_{0}+Q_{0} \Delta A^{T} \leq \Omega_{0}, \quad \Delta A \in U \tag{9.1}
\end{equation*}
$$

let $\Omega: \mathbb{N}^{n} \rightarrow \mathbb{N}^{n}$ be such that (3.1) is satisfied, and suppose there exists $\Delta Q \in \mathbb{N}^{n}$ satisfying

$$
\begin{equation*}
0=A \Delta Q+\Delta Q A^{\mathbf{T}}+\Omega(\Delta Q)+\Omega_{0} . \tag{9.2}
\end{equation*}
$$

Then

$$
\begin{equation*}
\left(A+\Delta A,\left[\Omega_{0}+\Omega(\Delta \mathcal{Q})-\left(\Delta A \Delta \mathcal{Q}+\Delta \mathcal{Q} \Delta A^{\mathrm{T}}\right)\right]^{\frac{1}{2}}\right) \text { is stabilizable, } \quad \Delta A \in U \tag{9.3}
\end{equation*}
$$

if and only if

$$
\begin{equation*}
A+\Delta A \text { is asymptotically stable, } \quad \Delta A \in U \tag{9.4}
\end{equation*}
$$

In this case,

$$
\begin{equation*}
Q_{0}-\Delta Q \leq Q_{\Delta A} \leq Q_{0}+\Delta Q, \quad \Delta A \in U \tag{9.5}
\end{equation*}
$$

where $Q_{\Delta A}$ is given by (2.7), and

$$
\begin{equation*}
\operatorname{tr}\left(Q_{0}-\Delta Q\right) R \leq J_{\mathrm{S}}(U) \leq \operatorname{tr}\left(Q_{0}+\Delta Q\right) R, \tag{9.6}
\end{equation*}
$$

$$
\begin{equation*}
\lambda_{\max }\left[\left(Q_{0}-\Delta Q\right) R\right] \leq J_{\mathrm{D}}(\mathcal{U}) \leq \lambda_{\max }\left[\left(Q_{0}+\Delta Q\right) R\right] . \tag{9.7}
\end{equation*}
$$

Proof. Define

$$
\begin{equation*}
\Delta Q \triangleq Q_{\Delta A}-Q_{0} \tag{9.8}
\end{equation*}
$$

and subtract (3.14) from (2.7) to obtain

$$
\begin{equation*}
0=(A+\Delta A) \Delta Q+\Delta Q(A+\Delta A)^{\mathrm{T}}+\Delta A Q_{0}+Q_{0} \Delta A^{\mathrm{T}} . \tag{9.9}
\end{equation*}
$$

Now rewrite (9.2) as

$$
\begin{equation*}
0=(A+\Delta A) \Delta \mathcal{Q}+\Delta Q(A+\Delta A)^{T}+\Omega(\Delta Q)-\left(\Delta A \Delta Q+\Delta Q \Delta A^{T}\right)+\Omega_{0} \tag{9.10}
\end{equation*}
$$

Using (9.10), the equivalence of (9.3) and (9.4) is immediate as in the proof of Theorem 3.1. Next, subtracting (9.9) from (9.10) yields

$$
\begin{align*}
0=(A+\Delta A)(\Delta Q-\Delta Q) & +(\Delta Q-\Delta Q)(A+\Delta A)^{\mathrm{T}}+\Omega(\Delta Q)-\left(\Delta A \Delta Q+\Delta Q \Delta A^{\mathrm{T}}\right) \\
& +\Omega_{0}-\left(\Delta A Q_{0}+Q_{0} \Delta A^{\mathrm{T}}\right) . \tag{9.11}
\end{align*}
$$

Using (3.1) and (9.1) it follows from (9.11) that

$$
\Delta Q-\Delta Q \geq 0
$$

or, equivalently,

$$
\begin{equation*}
Q_{\Delta::} \leq Q_{0}+\Delta \mathcal{Q} . \tag{9.12}
\end{equation*}
$$

To obtain the lower bound rewrite (9.9) as

$$
\begin{equation*}
0=(A+\Delta A)(-\Delta Q)+(-\Delta Q)(A+\Delta A)^{\mathbf{T}}-\left(\Delta A Q_{0}+Q_{0} \Delta A^{\mathrm{T}}\right) \tag{9.13}
\end{equation*}
$$

Note that because of the assumed symmetry of $U$, (9.1) holds with $\Delta A$ appearing in the inequality replaced by $-\Delta A$. Hence it can be shown similarly that

$$
\Delta Q+\Delta Q \geq 0,
$$

or, equivalently,

$$
\begin{equation*}
Q_{0}-\Delta Q \leq Q_{\Delta A} . \tag{9.14}
\end{equation*}
$$

Finally, (9.6) follows immediately from (9.5) while (9.7) is a consequence of Theorem 4.3.1 of [40].

Remark 9.1. To compare the upper bound in (9.5) with (3.5), rewrite (9.2) as

$$
\begin{equation*}
0=A\left(Q_{0}+\Delta Q\right)+\left(Q_{0}+\Delta Q\right) A^{\mathbf{T}}+\Omega(\Delta Q)+\Omega_{0}+V \tag{9.15}
\end{equation*}
$$

If $\Omega(\Delta Q)+\Omega_{0}=\Omega\left(Q_{0}+\Delta Q\right)$ then (9.15) has the same form as (3.2) and thus the two upper bounds are identical. This will be the case, for example, if $\Omega(\cdot)=\Omega_{7}(\cdot)$ and $\Omega_{0}$ is chosen to be $\Omega_{7}\left(Q_{0}\right)$ since $\Omega_{7}(\cdot)$ is linear. If, for example, $\Omega(\Delta Q)+\Omega_{0}<\Omega\left(Q_{0}+\Delta Q\right)$ then the upper bound in (9.5) will be sharper. In any case it is clear that the individual treatment of $\Delta Q$ and $Q_{0}$ yields potentially new upper bounds.

Remark 9.2. Theorem 9.1 does not guarantee that the lower bound $Q_{0}-\Delta \mathcal{Q}$ for $Q_{\Delta A}$ is nonnegative definite. However, $Q_{\Delta A}$ is always nonnegative definite and thus the lower bound in (9.5) may be of limited usefulness. Nevertheless, if $Q_{0}-\Delta Q$ is indefinite then, depending on $R$, the lower bounds in (9.6) and (9.7) may still be positive and thus be meaningful lower bounds.

## 10. Analytical Examples

In this section we consider simple analytical examples which illustrate the principal results of the paper. These examples also provide insight into the individual characteristics of different bounds as a prelude to numerical examples considered in the following section.

To begin we consider the simplest possible example. Set $n=1, A<0, R>0, V>0, A_{1}=1$, and $u=\left\{\Delta A:|\Delta A| \leq \delta_{1}\right\}$. For $\delta_{1}<-A, Q_{\Delta A}=V / 2(|A|-\Delta A)$ and $J_{S}(u)=J_{D}(u)=$ $R V / 2\left(|A|-\delta_{1}\right)$, where this worst-case performance is achieved for $\Delta A=\delta_{1}$. Solving MLE1 yields $Q=V / 2\left(|A|-\delta_{1}\right)$ which is a nonconservative result for both robust stability and performance. The same result is obtained from MLE4 by setting $\alpha=\alpha_{1}=\delta_{1}$. To apply MLE5, set $\delta_{1}=\sqrt{D E}$. Choosing $\alpha=2 \delta_{1}\left(|A|-\delta_{1}\right) E V$ again yields the nonconservative result. Finally, the same result follows from Theorem 8.1.

For the second example we consider nondestabilizing uncertainty in the imaginary component of an uncertain eigenvalue, i.e., frequency uncertainty, in contrast to uncertainty in the real part considered in the previous example. Let $n=2, A=\left[\begin{array}{cc}-\nu & \omega \\ -\omega & -\nu\end{array}\right], \nu>0, \omega \geq 0, V=R=I_{2}$, and $u=\left\{\Delta A: \Delta A=\sigma_{1} A_{1},\left|\sigma_{1}\right| \leq \delta_{1}\right\}$, where $A_{1}=\left[\begin{array}{cc}0 & 1 \\ -1 & 0\end{array}\right]$. Obviously, $A+\Delta A$ remains asymptotically stable for all values of $\sigma_{1}$ since $\Delta A$ affects only the imaginary part of the poles of $\Delta A$. The question then is whether the robustness tests are able to guarantee this robustness. Note also that because of the choice of $V, Q_{\Delta A}=Q_{0}=(2 \nu)^{-1} I_{2}$ for all $\Delta A \in U$.

For this example we note that MLE1 is satisfied by $Q=(2 \nu)^{-1} I_{2}$ which is independent of $\delta_{1}$. Thus MLE1 possesses a nonnegative-definite solution for all $\delta_{1}>0$ which shows that MLE1 is nonconservative with respect to robust stability and performance. Since $A(\Delta A)=(\Delta A) A$, it can also be seen that the same result holds for (8.18). The situation is considerably different for MLE4 and MLE5. To analyze MLE4 note that $\&$ has an eigenvalue $-2 \nu+\alpha+\delta_{1}$. (This can be shown by diagonalizing $A$ and $A_{1}$ and thus $A$.) Since, by Proposition 7.1, $A$ must be asymptotically stable, we require $\delta_{1}<2 \nu$. This is, of course, an extremely conservative result, especially when the damping $\nu$ is small. For MLE5 we can factor $A_{1}=D_{1} E_{1}$. Thus, let $D_{1}=I_{2}$ and $E_{1}=A_{1}$ and define $D=\delta_{1}^{2} I_{2}$ and $E=I_{2}$. Assuming that $Q$ is a multiple of $I_{2}$, it follows that $Q$ is nonnegative definite only if $\delta_{1} \leq \nu$, which is again an extremely conservative result. The reason for this conservatism becomes clear by noting that $D$ and $E$ as given above will also serve as bounds for perturbations of the form $\sigma_{1} I_{2}$ for which the range of nondestabilizing $\sigma_{1}$ is $\left|\sigma_{1}\right|<\delta_{1}$. This will also be the case for all factorizations $D_{1} E_{1}$ of $A_{1}$ since $D_{1} D_{1}^{\mathrm{T}}$ and $E_{1}^{\mathrm{T}} E_{1}$ must be positive definite and thus will also serve as bounds for destabilizing perturbations such as $\sigma_{1} I_{2}$.

Finally, we consider a nondestabilizing uncertainty affecting the interaction of a pair of real poles. Let $n=2, A=-I_{2}, V=R=I_{2}$, and $U=\left\{\Delta A: \Delta A=\sigma_{1} A_{1},\left|\sigma_{1}\right| \leq \delta_{1}\right\}$, where $A_{1}=\left[\begin{array}{ll}0 & 1 \\ 0 & 0\end{array}\right]$. Obviously, $A+\Delta A$ remains stable for all values of $\sigma_{1}$ since $\Delta A$ does not affect the nominal poles. Note that $Q_{\Delta A}=\left[\begin{array}{cc}\sigma_{2}^{2} / 4+1 / 2 & \sigma_{1} / 4 \\ \sigma_{1} / 4 & 1 / 2\end{array}\right]$ and $J_{S}(U)=\frac{1}{4} \delta_{1}^{2}+1$, where this worst-case performance is achieved for $\sigma_{1}=\delta_{1}$. In this case MLE1 has the solution $Q=\left(2-\delta_{1}\right)^{-1} I_{2}$ which is valid only for $\delta_{1}<2$, an extremely conservative robust stability result. Furthermore, the corresponding performance bound $\operatorname{tr} Q R=2\left(2-\delta_{1}\right)^{-1}$ is conservative with respect to the actual worst-case performance $\frac{1}{4} \delta_{1}^{2}+\frac{1}{2}$. In constrast MLE4 has the solution $Q=\left[\begin{array}{cc}\left(2-\alpha \delta_{1}\right)^{-1}+\alpha^{-1} \delta_{1}\left(2-\alpha \delta_{1}\right)^{-2} & 0 \\ 0 & \left(2-\alpha \delta_{1}\right)^{-1}\end{array}\right]$ which is nonnegative definite for all $\delta_{1}$ so long as $\alpha<2 / \delta_{1}$. Hence MLE4 is nonconservative with respect to robust stability. For robust performance, $\operatorname{tr} Q R=2\left(2-\alpha \delta_{1}\right)^{-1}+\alpha^{-1} \delta_{1}\left(2-\alpha \delta_{1}\right)^{-2}$ which is clearly an upper bound for $\frac{1}{4} \delta_{1}^{2}+1$. Choosing, for example, $\alpha=\delta_{1}^{-1}$ yields $\operatorname{tr} Q R=\delta_{1}^{2}+2$. The parameter $\alpha$ can also be chosen to minimize $\operatorname{tr} Q R$, although this is somewhat tedious to carry out analytically. Finally, MLE5 has the solution $Q=\left[\begin{array}{cc}\frac{1}{2}\left(1+a^{-1} \delta_{1}\right) & 0 \\ 0 & {\left[1-\left(1-a \delta_{1}\right)^{\frac{1}{2}}\right] / a \delta_{1}}\end{array}\right]$ which exists so long as $\alpha \leq 1 / \delta_{1}$. Hence MLE5 is also nonconservative with respect to robust stability. Choosing $\alpha=1 / \delta_{1}$ yields $\operatorname{tr} Q R=\frac{1}{2} \delta_{1}^{2}+\frac{s}{2}$ which lies above the minimal bound $\frac{1}{4} \delta_{1}^{2}+1$. Again, $\alpha$ can be chosen to minimize tr $Q R$.

## 11. Numerical Example

In this section we consider additional examples illustrating the results developed in earlier sections. In contrast to the analytical examples considered in Section 10, however, we consider more complex examples by numerically solving the modified Lyapunov equations. Here we focus on MLE4 and MLE5 which are the easiest to solve numerically. Specifically, we solved MLE4 by using the representation (7.2) (although this may not be practical when $n$ is large), and we solved MLE5 by means of a standard Riccati package. To simplify matters we consider only uncertainties $\Delta A$ of the form $\sigma_{1} A_{1}$. Presentation and evaluation of robust stability and performance results for multiparameter uncertainty can be fairly complex and thus is deferred to a future numerical study.

Since both robustness tests MLE4 and MLE5 depend upon an arbitrary positive constant $\alpha$, it is desirable to determine the value of $\alpha$ which yields the tightest (i.e., lowest) performance bound for each robust stability range. To this end we performed a simple one-dimensional search to determine the best such $\alpha$. Although analytical techniques may assist in determining optimal values of $\alpha$ more efficiently, the search technique proved to be adequate for the examples considered here.

As a first example we consider the control system given in [1] to demonstrate the lack of a guaranteed gain margin for LQG controllers. Hence consider

$$
\begin{align*}
\dot{x}_{0}(t) & =A_{0} x_{0}(t)+B_{0} u(t)+w_{1}(t)  \tag{11.1}\\
y(t) & =C_{0} x_{0}(t)+w_{2}(t) \tag{11.2}
\end{align*}
$$

with controller

$$
\begin{align*}
\dot{x}_{\mathrm{c}}(t) & =A_{c} x_{c}(t)+B_{c} y(t),  \tag{11.3}\\
u(t) & =C_{c} x_{c}(t), \tag{11.4}
\end{align*}
$$

and performance

$$
\begin{equation*}
J=\lim _{t \rightarrow \infty} \mathbb{E}\left[x_{0}^{\mathrm{T}}(t) R_{1} x_{0}(t)+u^{\mathrm{T}}(t) R_{2} u(t)\right] . \tag{11.5}
\end{equation*}
$$

The data are

$$
\begin{gathered}
A_{0}=\left[\begin{array}{ll}
1 & 1 \\
0 & 1
\end{array}\right], \quad B_{0}=\left[\begin{array}{l}
0 \\
1
\end{array}\right], \quad C_{0}=\left[\begin{array}{ll}
1 & 0
\end{array}\right], \\
V_{1}=R_{1}=\rho\left[\begin{array}{ll}
1 & 1 \\
1 & 1
\end{array}\right], \quad V_{2}=R_{2}=1,
\end{gathered}
$$

where $V_{1}$ and $V_{2}$ are the intensities of $w_{1}(t)$ and $w_{2}(t)$, respectively. Uncertainty $\Delta B_{0}$ in $B_{0}$ is thus represented by $\sigma_{1} B_{1}$, where $B_{1}=\left[\begin{array}{ll}0 & 1\end{array}\right]^{\mathrm{T}}$. Thus the closed-loop system corresponds to

$$
\begin{gathered}
A=\left[\begin{array}{cc}
A_{0} & B_{0} C_{c} \\
B_{c} C_{0} & A_{c}
\end{array}\right], \quad A_{1}=\left[\begin{array}{cc}
0 & B_{1} C_{c} \\
0 & 0
\end{array}\right], \\
R=\left[\begin{array}{cc}
R_{1} & 0 \\
0 & 0
\end{array}\right], \quad V=\left[\begin{array}{cc}
V_{1} & 0 \\
0 & B_{c} V_{2} B_{c}^{\mathrm{T}}
\end{array}\right],
\end{gathered}
$$

where the zero in the $(2,2)$ block of $R$ denotes the fact that we are considering the robust performance bound for the state regulation cost only. Choosing $\rho=60$, it follows that the LQG gains are given by

$$
A_{a}=\left[\begin{array}{cc}
-9 & 1 \\
-20 & -9
\end{array}\right], \quad B_{c}=\left[\begin{array}{l}
10 \\
10
\end{array}\right], \quad C_{c}=\left[\begin{array}{ll}
-10 & -10
\end{array}\right] .
$$

For this controller the actual stability region corresponds to $\sigma_{1} \in(-.07, .01)$ so that the largest symmetric region about $\sigma_{1}=0$ is $\left|\sigma_{1}\right| \leq .01$. The worst- case performance over each stability region $\left|\sigma_{1}\right| \leq \delta_{1}$ is denoted by the solid line in Figure 1, while the performance bounds obtained from MLE4 and MLE5 are shown for several values of $\delta_{1}$. For MLE5 we set $D_{1}=\left[\begin{array}{llll}0 & 1 & 0 & 0\end{array}\right]^{\mathrm{T}}$ and $E_{1}=\left[\begin{array}{lll}0 & 0 & C_{c}\end{array}\right]$. Note that MLE5 yields considerably tighter estimates of worst-case performance, particularly as $\delta_{1}$ approaches .01. For MLE4, optimal values of $\alpha$ were in the range .0012 to .0058 , while for MLE5 (with $\Omega_{10}^{\prime}(\cdot)$, see (5.26)), $\alpha$ was in the range .0143 to .0020 .

As a second example we consider a pair of nominally uncoupled oscillators with uncertain coupling. This example was considered in [45] using the majorant Lyapunov technique. Let

$$
\begin{gathered}
A=\left[\begin{array}{cccc}
-\nu & \omega_{1} & 0 & 0 \\
-\omega_{1} & -\nu & 0 & 0 \\
0 & 0 & -\nu & \omega_{2} \\
0 & 0 & -\omega_{2} & -\nu
\end{array}\right], \quad A_{1}=\left[\begin{array}{llll}
0 & 0 & 1 & 0 \\
0 & 0 & 0 & 1 \\
1 & 0 & 0 & 0 \\
0 & 1 & 0 & 0
\end{array}\right], \\
\nu=.2, \quad \omega_{1}=.2, \quad \omega_{2}=1.8, \quad R=V=I_{4},
\end{gathered}
$$

and, for MLE5, define $D=A_{1}$ and $E_{1}=I_{4}$. We consider bounds on $J_{S}(U)$ only.
Figure 2 illustrates the exact worst case performance along with performance bounds obtained from MLE4 and MLE5. For MLE4 optimal values of $\alpha$ ranged from .036 to .141 while for MLE5 optimal $\alpha$ was between .361 and .096 . Although MLE4 was slightly less conservative than MLE5, both bounds were able to guarantee robust stability only for $\delta_{1}=.15$ while the largest stability region is actually $\delta_{1}=.54$. It is interesting to contrast this result with [45] where the majorant Lyapunov technique yielded a robust stability range of $\delta_{1}=.4$ for a richer class of off-diagonal blocks having maximum singular value less than $\delta_{1}$.

Conclusion. It seems clear that no single quadratic Lyapunov bound is superior to the others. Although the conservatism of each bound is problem dependent, it is desirable to better understand the nature of the conservatism in order to utilize the bounds in an effective manner. Finally, the example illustrated in Figure 2 may indicate fundamental limitations of the quadratic Lyapunov function approach to robustness as compared to the majcrant technique of [45]. These remain questions for future research.

Acknowledgment. We wish to thank A. W. Daubendiek for producing the numerical results in Section 11.


$1 0 \longdiv { - 2 }$
ITY REGION $\boldsymbol{\delta}_{\mathbf{1}}$
Figure 2
aNnOg ヨDNEWAOJyga

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# The Majorant Lyapunov Equation: A Nonnegative Matrix Equation for Robust Stability and Performance of Large Scale Systems 

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#### Abstract

A new robust stability and performance analysis technique Is developed. The approach lavolves replacing the state covariance by its block-norm matrix, i.e., the nonaegative matrix whose elements are the morms of subblocks of the covariance matrix partitioned according to subsystem dyoamics. A bound (i.e., majorant) for the block-norm matrix is given by the majorant Lyapunov equation, a Lyapunov-type nonnegative matrix equation. Existence, uniqueness, and computational tractabilits of solutions to the majorant Lyapunov equation are shown to be completely characterized in terms of $M$ matrices. Two examples are considered. For a damped simple harmonic oscillator with uncertain but constant natural frequency, the majorant lyapunov equation predicts unconditional stability. And, for a pair of nominally uncoupled oscillazors with uncertain coupliag, the majorant Lyapunov equation shows that the range of noondestabilizing coupliags is proporional to the frequency separation between the oscillators, a result not predictable from quadratic or vector Lyapunov functions.


## I. Introduction

THE importance of robustness in control-system analysis and design cannot be overemphasized. The past ten years' literature reflects considerable frequency-domain development [1]-[11], while recent publications indicate increasing timedomain activity [12]-[19]. Wide variations in underlying assumptions. mathematical settings, and problem data render it difficult, if not impossible, to clearly delineate the relative effectiveness of different methods. Our own philosopiucal outlook has thus been guided by two general criteria:

1) effectiveness for simple examples:
2) efficiency when applied to large scale problems.

The first criterion involves applying robustness techniques to simple. perhaps trivially obvious, examples to serve as "acid tests." A given method's effectiveness on a collection of such examples can possibly reveal inherent shoncomings. As an illustration of this criterion, consider a damped harmonic oscillator with constant but uncertain natural frequency. Using the nocation of [6], stability is guaranteed so long as

$$
\begin{equation*}
\sigma_{\text {mas }}\left[R(j \omega)(I+G(j \omega) K(j \omega))^{-1} G(j \omega) L^{-1}(j \omega)\right]<1, \quad \omega \geq 0 \tag{1.1}
\end{equation*}
$$

where, for $\boldsymbol{y}>\mathbf{0}$,

$$
G(s)=\left(s^{2}+2 \nu s+\omega_{n}^{2}\right)^{-1}
$$

[^18]and uncertainty in the nominal natural frequency $\omega_{n}$ is modeled by
\[

$$
\begin{gathered}
\Delta(s)=L^{-1}(s) \theta(s) R(s)=\delta \omega_{n}^{2}, \\
L(s)=1 / \alpha, \theta(s)=\delta / \alpha, R(s)=\omega_{n}^{2}, K(s)=0, \\
\delta \in[-\min (1, \alpha), \alpha] . \quad \alpha>0 .
\end{gathered}
$$
\]

Note that

$$
\sigma_{\max }[\theta(j \omega)] \leq 1, \quad \omega \geq 0
$$

as required in [6]. The perturbation $\Delta(s)$ (modeled as a feedback gain) effectively replaces $\omega_{n}^{2}$ in $G(s)$ by $(1+\delta) \omega_{n}^{2}$. Hence. for a given $\alpha>0$ this uncertainty model permits perturbed natural frequencies in the range $\left[0 .(1+\alpha)^{12} \omega_{n}\right]$. Evaluating (1.1) yields the upper bound

$$
\begin{equation*}
\alpha<\left[\left(\omega_{n}^{2}-\omega^{2}\right)^{2}+4 \nu^{2} \omega^{2}\right]^{12} / \omega_{n}^{2}, \quad \omega \geq 0 \tag{1.2}
\end{equation*}
$$

or, equivalently,

$$
\begin{equation*}
\alpha<2 \xi\left(1-\zeta^{2}\right)^{1 / 2} \tag{1.3}
\end{equation*}
$$

where $\zeta \doteq \nu / \omega_{n}$. The conservatism of (1.3) is obviously most pronounced when the damping ratio $\zeta$ is small. In all cases. however, the conservatism is infinite.

The second criterion is obviously subjective and depends upon a variety of factors such as problem structure. designer experience. and computational resources. This criterion is, in our opinion, most important since the need for robustness techniques becomes increasingly critical as system complexity grows. Indeed, the ultimate test of a given approach is to scale it up to larger and larger problems to reveal inherent limitations. Obviously, such tests are not only difficult, but may entail a significant commitment of human and financial resources. Nevertheless, crude predictions are sometimes available, and a case in point is the "curse of dimensionality" encountered in the approach of [9]. Another example involves computational difficulties in obtaining bounds for the $\mu$-function with more than three blocks [10].

The contribution of the present paper is a new robustness analysis method developed specifically for large scale systems. The basic idea, motivated by the work of Siljak [30] on connective stability, is as follows. The system is assumed to be in the form of a collection of subsystems with uncerain local dynamics and uncertain interactions. ${ }^{1}$ Parameter uncertainties are modeled as either structured or unstructured constant variations contained in prescribed sets. The state covariance, partitioned conformably with the subsystem dynamics, is replaced by its block-norm matrix, i.e., the nonnegative matrix each of whose elements is the norm of the corresponding subblock of the original matrix. This nonnegative matrix satisfies a novel inequality designated the

[^19]covariance block-norm inequality. The existence of a solution to the majorant Lyapunov equation, i.e., the covariance blocknorm inequality interpreted as an equation, yields an element-byelement bound (i.e., majorant) for the covariance block-norm, hence, assuring robust stability and performance. The relevance of this technique to large scale systems stems from the fact that replacing each subblock of the covariance by its norm can significantly reduce the dimension of the problem. Indeed, the dimension of the majorant Lyapunov equation is equal to the number of subsystems which may be significantly less than the dimension of the original system.

To illustrate the above ideas in more detail, consider the covariance equation

$$
\begin{equation*}
0=(A+G) Q+Q(A+G)^{T}+V \tag{1.4}
\end{equation*}
$$

where $A$ denotes the nominal dynamics, $G$ denores uncertainty in $\boldsymbol{A}, V$ is the disturbance intensity, and $Q$ is the state covariance. Assuming that $\bar{A}$ is block diagonal with $r$ diagonal blocks leads to the covariance block-norm inequality (see Proposition 4.2)

$$
\begin{equation*}
Q-Q \leq \leq S Q+Q S^{r}+V \tag{1.5}
\end{equation*}
$$

In (1.5), $Q, Q, G$, and $v$ are $r \times r$ nonnegative matrices, i.e., each element is a nonnegative number. The matrices $Q$ and $V$ are formed by taking the Frobenius norm of each subblock of $Q$ and $V$, while each component of $\mathcal{S}$ is a given constant which bounds the spectral norm (largest singular value) of the corresponding subblock of the uncertain perturbation $G$. Hence, $\mathcal{G}$ is a majorant for $G$ in the sense of [21]-[23]. Each element of the matrix $Q$ is bounded above by the smallest singular value of the Kronecker sum [24]-[26] of pairs of diagonal blocks of $\bar{A}$. The operation " "" is the Hadamard product [27], [28], and the ordering " $\leq \leq$ " denotes element-by-element comparison, i.e., the ordering induced by the cone of nonnegative matrices [29], [30].
The majorant Lyapunov equation is obtained by replacing the inequality (1.5) by the $r \times r$ nonnegative matrix equation

$$
\begin{equation*}
Q * \mathbb{Q}=S \mathbb{Q}+\mathbb{Q} S^{r}+V \tag{1.6}
\end{equation*}
$$

A key result (Corollary 5.1) states that

$$
\begin{equation*}
Q \leq \leq \mathbb{Q} \tag{1.7}
\end{equation*}
$$

for all stable $\boldsymbol{A}+\boldsymbol{G}$. Consequently (see Theorem 5.1), the existence of a unique solution to (1.6) leads directly to a guarantee of robust stability over the range specified by $S$ and to a performance bound involving $\mathbb{Q}$. Moreover, solutions of (1.6) exist if and only if the $r^{2} \times r^{2}$ matrix

$$
\begin{equation*}
A \equiv \operatorname{diag}(\operatorname{vec} Q)-S \oplus G \tag{1.8}
\end{equation*}
$$

is an $M$ matrix [29], [30].
Even when the number of subsystems is large, the majorant Lyapunov equation is generally computationally tractable. Specifically, although $A$ is an $r^{2} \times r^{2}$ matrix, no computations whotsoever need to be corried out with matrices of this dimension. Rather, it suffices to solve only the majorant Lyapunov equation (1.6). In this regard we show that $\mathbb{Q}$ is given by

$$
\begin{equation*}
\mathbb{Q}=\lim _{i \rightarrow \infty} \mathbb{Q}_{i} \tag{1.9}
\end{equation*}
$$

where if $\boldsymbol{G}$ has only off-diagonal nonzero blocks the sequence $\left\{Q_{i}\right\}$ is generated by

$$
\begin{equation*}
Q=Q_{i+1}=S Q_{i}+Q_{i} S^{T}+V, \quad Q_{0}=0 \tag{1.10}
\end{equation*}
$$

and is monotonically increasing. Furthermore, the convergence of this sequence is equivalent to $\alpha$ being an $M$ matrix so that it is not even necessary to form $A$. Note that (1.6) does not require the development of new solution techniques. Indeed, since (1.10) is a
straightforward iteration, (1.6) is even easier to solve than the original Lyapunov equation (1.4).

To illustrate these results we consider two examples. The first example is the damped oscillator already considered in this section. With little effort the majorant Lyapunov equation yields the (obvious) result that the oscillator is stable for all constant natural frequencies. The second example involves a pair of oscillators with known parameters but with uncertain coupling. The majorant Lyapunov equation yields bounds over which stability is guaranteed, and these bounds are compared to the actual stability region as a function of frequency separation. The main result shows that the robustness to uncertain coupling is proportional to the frequency separation. This weak subsystem interaction robustification mechanism is the principal contribution of the majorant theory. This example has immediate application to the problem of vibration control in flexible structures. For this class of problems the open-loop dynamics can be viewed as a collection of uncoupled oscillators which become coupled via feedback and structural uncertainties.

The majorant bound developed in the present paper is quite different from the widely used quadratic Lyapunov function (see, e.g., [12]-[20]). As can readily be shown using the methods of [12], [17]-[20], the quadratic Lyapunov function yields robust stability and performance by replacing (1.4) by

$$
\begin{equation*}
0=A Q+Q A^{T}+\Omega(Q)+V \tag{1.11}
\end{equation*}
$$

where $\Omega(\cdot)$ satisfies

$$
\begin{equation*}
G Q+\phi G^{T} \leq \Omega(\phi) \tag{1.12}
\end{equation*}
$$

for all variations $G$. It can then be shown that

$$
\begin{equation*}
Q \leq \phi \tag{1.13}
\end{equation*}
$$

where now, in contrast to (1.7), the ordering in (1.13) is defined with respect to the cone of nonnegative-definite matrices. Indeed, the majorant bound may be more closely related to vector Lyapunov functions [30], [31] and the Lyapunov matrix function [32], [33]. It does not appear possible, however, to use these techniques to obtain the majorant results on robustness due to subsystem frequency separation.

The reader will observe that this paper exploits a wide variety of techniques including nonnegative matrices, block norms, matrix majorants, the Hadamard product, the Kronecker sum, and $M$ matrices. Each of these techniques, except majorants, has, however, been previously applied to control problems in numerous instances. In the special case of scalar subblocks, the blocknorm matrix has, moreover, been utilized by Yedavalli [13]-[15] and others for robustness analysis and design. In this case the block norm is known as the matrix modulus. The variety of algebraic structures employed in the present paper should not be surprising since the quest for increasingly refined robustness techniques can be expected to invoke correspondingly refined uncertainty bounds. Related techniques are employed in [11]. Furthermore, nonnegative matrix equations involving $M$ matrices arise naturally in a variety of settings (see, e.g., [38], [39]).

The contents of the paper are as follows. Section II presents notation, definitions, and lemmas for use throughout the paper. In Section III robust stability and performance are defined for the homogeneous and nonhomogeneous systems. Detailed system structure and uncertainty characterization are given in Section IV and the covariance block-norm inequality is derived. Section V analyzes the majorant Lyapunov equation to obtain a majorant for the steady-state covariance. The main result, Theorem S.1, guarantees robust stability and provides a performance bound. Finally, the examples appear in Section VI.

## II. Preliminaries

The following notation will be used throughout. All matrices are assumed to have real entries.

| T | expected value |
| :---: | :---: |
| Q, Es ${ }^{p \times q}, \mathbb{H}^{p}$ | real numbers, $p \times q$ real matrices, $0^{p \times 1}$ |
| $I_{p}, 0_{p \times q}, 0_{p}$ | $p \times p$ identity matrix, $p \times q$ zero matrix, $0_{p \times p}$ |
| - , 8 | Kronecker sum, Kronecker product [24]-[27] |
| * | Hadamard product [27], [28] |
| $\operatorname{col}_{i}(\mathrm{Z})$ | $i$ th column of matrix $\bar{Z}$ |
| vec (Z) | $\left[\begin{array}{c} \operatorname{col}_{1}(Z) \\ \vdots \\ \operatorname{col}_{q}(Z) \end{array}\right] \in B^{p q}, Z \in \mathbb{B}^{p \times q}$ |
| $Z_{(i, j)}$ | (i, $j$ ) element of matrix $Z$ |
| $Z^{T}$ | transpose of vector or matrix $Z$ |
| $2^{-7}$ | $\left(Z^{T}\right)^{-1}$ or $\left(Z^{-1}\right)^{T}$ |
| 42 | trace of matrix $\boldsymbol{Z}$ |
| diag ( $\boldsymbol{Z}_{1}, \cdots, Z_{p}$ ) | diagonal matrix with listed diagonal elements |
| block-diag ( $Z_{1}, \cdots, Z_{p}$ ) | block-diagonal matrix with listed diagonal blocks |
| $\rho(Z)$ | spectral radius of $Z$ |
| asymptotically stable matrix | matrix with eigenvalues in open lefthalf plane |
| nonnegative-definite matrix | symmetric matrix with nonnegative eigenvalues ( $Z \geq 0$ ) |
| positive-definite matrix | symmetric matrix with positive eigen: values ( $Z>0$ ) |
| $\boldsymbol{Z}_{1} \geq \boldsymbol{Z}_{2}$ | $Z_{1}-Z_{2} \geq 0, Z_{1}, Z_{2}$ symmetric |
| $Z_{1}>Z_{2}$ | $Z_{1}-Z_{2}>0, Z_{1}, Z_{2}$ symmetric |
| nonnegative matrix | matrix with nonnegative elements $(Z \geq \geq 0)[29] .[30]$ |
| positive matrix | matrix with positive elements ( $Z>0$ ) |
| $Z_{1} \geq \geq Z_{2}$ | $Z_{1}-Z_{2} \geq \geq 0$ |
| $Z_{1}>Z_{2}$ | $\boldsymbol{Z}_{1}-\boldsymbol{Z}_{2}>0$ |
| $Z^{H 1}$ | Hadamard inverse, $\left(Z^{H I}\right)_{(i, j)} \triangleq$ $\left[Z_{(i, j)}\right]^{-1}, Z>0$ |
| block-norm matrix | nonnegative matrix each of whose elements is the norm of a corresponding subblock of a given partitioned matrix |
| majorant | nonnegative matrix each of whose elements bounds the corresponding element of a block-norm matrix |
| $\\|Z\\|_{\sigma_{i}(Z)}$ | Euclidean norm of vector $Z$ singular value of matrix $Z$ |
| $\sigma_{\text {min }}(Z), \sigma_{\text {max }}(Z)$ | smallest and largest singular values of matrix 2 |
| $\lambda_{\text {max }}(Z)$ | largest eigenvalue of symmetric matrix $Z$ |
| $\\|Z\\|_{s}$ | $\sigma_{\max }(Z)$ (spectral norm induced by $\left.\\|\cdot\\|_{2}\right)$ |
| $\\|\boldsymbol{Z}\\|_{F}$ | $\left(\operatorname{tr} 2 Z^{T}\right)^{1 / 2}=\left[\sum_{i, j=1}^{0.4} Z_{(i, j)}^{2}\right]^{1 / 2}$ |
|  | $=\left[\sum_{i=1}^{p} \sigma_{i}^{2}\right]^{1 / 2}$ |

(Frobenius norm [34]).
In subsequent sections we shall exploit the fact that the norms $\|\cdot\|_{2},\left\|_{1} \cdot\right\|_{s}$, and $\|\cdot\|_{F}$ coincide for vectors. Hence, if $Z \in \mathbb{R}^{p}$, then by interpreting $\mathbb{R}^{p}=\mathbb{R}^{p \times 1}$ it follows that

$$
\begin{equation*}
\|Z\|_{2}=\|Z\|_{s}=\|Z\|_{F} . \tag{2.1}
\end{equation*}
$$

Furthermore, if $Z \in \mathbb{R}^{p \times q}$, then

$$
\begin{equation*}
\|Z\|_{s} \leq\|Z\|_{F}=\|\operatorname{vec} Z\|_{F}=\|\operatorname{vec} Z\|_{2}=\|\operatorname{vec} Z\|_{s} . \tag{2.2}
\end{equation*}
$$

Lemma 2.1: If $Z \in H^{p \times q}$ and $Z \in J^{\rho \times \gamma}$ then

$$
\begin{equation*}
\sigma_{\min }(Z)\|Z\|_{F} \leq\|Z Z\|_{F} \leq\|Z\|_{s}\|Z\|_{F} \tag{2.3,2.4}
\end{equation*}
$$

If, furthermore, $p=q=r Z \geq 0$, and $Z$ is symmetric, then

$$
\begin{equation*}
\operatorname{tr} Z Z \leq(\operatorname{tr} Z) \lambda_{\max }(Z) \leq(\operatorname{tr} Z)\|Z\|_{s} . \tag{2.5}
\end{equation*}
$$

Proof: Inequality (2.4) can be found in [35, p. 263]. To prove (2.3), note that when $Z$ is singular the result is immediate. Otherwise, if $p=q$ replace $Z$ and $\mathcal{Z}$ in (2.4) by $Z^{-1}$ and $Z \mathcal{Z}$, respectively. The result now follows from $\left[\sigma_{\text {max }}(Z)\right]^{-1}=$ $\sigma_{\min }\left(Z^{-1}\right)$. If $p \neq q$, then related arguments apply. Finally, (2.5) is given in [36].
Recall [30] that a matrix $S \in B^{r \times r}$ is an $N$ matrix if $S_{(i, j)} \leq 0$. $i, j=1, \cdots, r, i \neq j$. If, in addition, all principal minors of $S$ are positive, then $S$ is an $M$ matrix.
Lemma 2.2: Suppose $S \in \beta^{r \times r}$ is an $N$ matrix. Then the following are equivalent:
i) $S$ is an $M$ matrix;
ii) det $S \neq 0$ and $S^{-1} \geq \geq 0$;
iii) for each $y \in B^{\prime}, y \geq \geq 0$, there exists a unique $x \in \beta^{r}, x$ $\geq \geq 0$. such that $S x=y$;
iv) there exists $x \in \beta^{\prime}, x \geq \geq 0$. such that $S x>0$;
v) $I_{r} * S \geqslant 0$ and each diagonal matrix $D>I_{r} * S$ satisfies $\rho\left[D^{-1}\left(I_{r} * S-S\right)\right]<1$.

Proof: The equivalence of statements i), ii), iv), and v) follow's from [30, p. 396]. The implication ii) $=\mathrm{iii}$ ) is immediate, and iii) $=$ iv) follows by setting $y=\left[\begin{array}{lll}11\end{array}\right]^{T}$.
Lemma 2.3: Suppose $S \in{\underset{R}{r x r}}^{R^{\prime}}$ is an $M$ matrix and let $S \in$ $\beta^{r \times r}$ be an $N$ matrix such that $S \geq \geq S$. Then $S$ is an $M$ matrix. Proof: See [30, p. 400].

## III. Robust Stability and Performance Bounds

Consider the $n$ th-order homogeneous system ${ }^{2}$

$$
\begin{gather*}
\dot{x}(t)=(A(\theta)+G) x(t), \quad t \in[0, \infty),  \tag{3.1}\\
G \in \Theta \subset B^{n \times n},  \tag{3.2}\\
\theta \in \theta \subset \Omega^{m}, \tag{3.3}
\end{gather*}
$$

where $A: \Theta \rightarrow \bar{B}^{n \times n}$ is continuous, $\bar{A} \equiv A(\bar{\theta})$ denotes the known nominal dynamics for $\hat{\theta} \in \Theta, \theta$ denotes the unstructured parametric uncertainty in $\bar{A}, G$ denotes the structured parametric uncertainty in $A$, and $0 \in G$ is the nominal value of $G$. We first consider the stability of (3.1) over 3 and $\theta$.

Definition 3.1: If $A(\theta)+G$ is asymptotically stable for all $G$ $\epsilon$ and $\theta \in \theta$, then the homogeneous system (3.1) is robustly stable over 3 and $\theta$.

Now consider the $n$ th-order nonhomogeneous system

$$
\begin{equation*}
\dot{x}(t)=(A(\theta)+G) x(t)+w(t), \quad t \in[0, \infty) \tag{3.4}
\end{equation*}
$$

where $G \in G, \theta \in \theta$, and $w(\cdot)$ is white noise with intensity $V \geq$ 0 . For given $G \in G$ and $\theta \in \Theta$, the steady-state average quadratic performance is defined by

$$
\begin{equation*}
J(G, \theta) \triangleq \lim \sup E\left[x^{T}(t) R x(t)\right] \tag{3.5}
\end{equation*}
$$

where $R=R^{T} \geq 0$. The system (3.4) may, for example, denote a control system in closed-loop configuration. There is no need in our development, however, to make such distinctions.

In practice, steady-state performance is only of interest when the system is robustly stable. The following result is immediate.

Proposition 3.1: Suppose the system (3.1) is robustly stable

[^20]over $b$ and $\theta$. Then for each $G \in \operatorname{tand} \theta \in \theta$,
\[

$$
\begin{equation*}
J(G, \theta)=\operatorname{tr} Q R \tag{3.6}
\end{equation*}
$$

\]

where $n \times n$ nonnegative-definite $Q$ is the unique solution to

$$
\begin{equation*}
0=(A(\theta)+G) Q+Q(A(\theta)+G)^{T}+V . \tag{3.7}
\end{equation*}
$$

We shall only be concerned with the case in which and $\theta$ are compact. Since $Q$ is a continuous function of $G$ and $\theta$, we can define the worst-case average steady-state quadratic performance

$$
\begin{equation*}
J_{\max } \triangleq \max _{\sigma \in,, \delta \in \theta} J(G, \theta) \tag{3.8}
\end{equation*}
$$

Since it is difficult to determine $J_{\text {mex }}$ explicitly, we shall seek upper bounds.

Definition 3.2: If $J_{\max } \leq \hat{\alpha}$, then $\hat{\alpha}$ is a performance bound for the nonhomogeneous system (3.4) over and $\theta$.

## IV. System Structure, Uncertanty Characterrzation, and the Covariance Block-Norm Inequality

A discussed in Section I, (3.1) and (3.4) are assumed to be in the form of a large scale system with uncoupled local dynamics and uncertain interactions. Hence, with the subsystem partitioning

$$
\begin{equation*}
n=\sum_{i=1}^{\infty} n_{i} \tag{4.1}
\end{equation*}
$$

the local system dynamics $A(\theta)$ can be decomposed into subsystem dynamics according to

$$
\begin{equation*}
A(\theta)=\underset{i=1, \cdots, r}{\text { block-diag }}\left\{A_{i}(\theta)\right\} \tag{4.2}
\end{equation*}
$$

where $A_{l}(\theta) \in B^{n_{i} \times n_{i}}, \theta \in \theta$. For convenience, denote

$$
A \xlongequal[\substack{\text { block-diag } \\ i=1, \cdots, r}]{ }\left\{A_{i}\right\}
$$

Accordingly, $R$ is assumed to be of the form

$$
\begin{equation*}
R=\underset{\substack{\text { block-diag } \\ i=1, \cdots, r}}{ }\left\{R_{i}\right\} \tag{4.3}
\end{equation*}
$$

where $R_{i} \in \beta^{n_{1} \times n_{i}}, R_{i} \geq 0, i=1, \cdots, r$. The intensity $V$ and steady-state covariance $Q$ satisfying (3.7) are assumed to be conformably partitioned, i.e.,

$$
\begin{array}{ll}
V=\left\{V_{i j}\right\}_{i, j=1}^{r}, & V_{i j} \in B^{n_{i} \times n_{j}} \\
Q=\left\{Q_{i j}\right\}_{i, j=1}^{r}, & Q_{i j} \in B^{n_{i} \times n_{j}} \tag{4.5}
\end{array}
$$

For notational simplicity define

$$
\begin{equation*}
V_{1} \triangleq V_{i i}, \quad Q_{i} \equiv Q_{i i}, \quad i=1, \cdots, r \tag{4.6}
\end{equation*}
$$

Taking the Frobenius norm of each subblock of $V$ and $Q$ leads to the $r \times r$ symmetric nonnegative matrices $V$ and $Q$ defined by

$$
\begin{equation*}
V \equiv\left\{\left\|V_{i j}\right\|_{F}\right\}_{i, j=1}^{\prime}, \quad Q \triangleq\left\{\left\|Q_{i j}\right\|_{r}\right\}_{i, j=1}^{r} \tag{4.7}
\end{equation*}
$$

Note that

$$
\begin{equation*}
\|Q\|_{F}=\|Q\|_{F},\|v\|_{F}=\|V\|_{F} . \tag{4.8}
\end{equation*}
$$

A few observations concerning the nominal system, i.e., with $G=0$ and $\theta=\theta$, are worth noting. If $\bar{A}$ is stable then so is $\boldsymbol{A}_{i}, i$ $=1, \cdots, r$, and there exist unique, nonnegative-definite $\phi_{i}, \hat{P}_{i}$ $\in \mathbb{R}^{m_{i} \times n_{i}}, i=1, \cdots, r$, satisfying

$$
\begin{align*}
& 0=A_{i} Q_{i}+Q_{i} A_{i} T+V_{i}  \tag{4.9}\\
& 0=A_{i} P_{i}+\beta_{i} A_{i}+R_{i} . \tag{4.10}
\end{align*}
$$

Proposition 4.I: Suppose $A$ is asymptotically stable. Then the nominal performance $J_{\text {nom }}$ is given by

$$
\begin{equation*}
J_{\text {nom }} \triangleq J(0, \bar{\theta})=\sum_{i=1}^{r} \operatorname{tr} Q_{i} R_{1}=\sum_{i=1}^{r} \operatorname{tr} \beta_{i} V_{i} \tag{4.11}
\end{equation*}
$$

Proof: First note that with $G=0$ and $\theta=\theta$ the diagonal blocks of $Q$ satisfying (3.7) coincide with $\boldsymbol{Q}_{1}, \cdots, Q_{\text {r }}$. Thus

$$
\begin{aligned}
J(0, \delta) & =\sum_{i=1}^{r} \operatorname{tr} Q_{i} R_{i} \\
& =\sum_{i=1}^{r}\left(\operatorname{vec} Q_{i}\right)^{T} \text { vec } R_{i} \\
& =\sum_{i=1}^{r}\left[\left(\bar{A}_{i} \oplus A_{i}\right)^{-1} \text { vec } V_{i}\right]^{T} \text { vec } R_{i} \\
& =\sum_{i=1}^{r}\left(\operatorname{vec} V_{i}\right)^{T}\left(A_{i}^{T} \oplus A_{i}^{T}\right)^{-1} \text { vec } R_{i} \\
& =\sum_{i=1}^{r}\left(\operatorname{vec} V_{i}\right)^{T} \text { vec } P_{i} \\
& =\sum_{i=1}^{r} \operatorname{tr} P_{i} V_{i}
\end{aligned}
$$

The matrices $G \in \Phi$ are also conformably partitioned so that

$$
\begin{equation*}
G=\left\{G_{i j}\right\}_{i, j=1}^{\prime}, G_{i j} \in \AA^{n_{i} \times n_{j}} \tag{4.12}
\end{equation*}
$$

and is characterized by

$$
\begin{equation*}
\triangleq\left\{G \in \mathbb{R}^{n \times n}: a_{\max }\left(G_{i j}\right) \leq \gamma_{i j}, \quad i, j=1, \cdots, r\right\} \tag{4.13}
\end{equation*}
$$

where $\gamma_{i j} \geq 0, i, j=1, \cdots, r$, are given constants. For convenience, define the $r \times r$ nonnegative matrix

$$
\begin{equation*}
S \equiv\left\{\gamma_{i j}\right\}_{i, j=1}^{\prime} \tag{4.14}
\end{equation*}
$$

The bound $\mathcal{G}$ is a matrix majorant for $G \in G$ in the sense of [21]-[23].

Remark 4.1: is compact and convex.
Finally, let symmetric, positive $Q \in \mathcal{F}^{r \times r}$ satisfy

$$
\begin{equation*}
a_{(i, j)} \leq \min _{\theta \in \theta}\left\{\sigma_{\min }\left(A_{j}(\theta) \oplus A_{i}(\theta)\right)\right\}, \quad i, j=1, \cdots, r \tag{4.15}
\end{equation*}
$$

Proposition 4.2: Let $G \in$ and $\theta \in \theta$ be such that $A(\theta)+$ $G$ is asymptotically stable and let $n \times n Q \geq 0$ satisfy (3.7). Then $Q$ defined by (4.7) satisfies

$$
\begin{equation*}
Q * Q \leq \leq Q Q+Q S^{T}+v \tag{4.16}
\end{equation*}
$$

or, equivalently,

$$
\begin{equation*}
A \operatorname{vec} Q \leq \leq \text { vec } V \tag{4.17}
\end{equation*}
$$

where

$$
\begin{equation*}
A \triangleq[\operatorname{diag}(\operatorname{vec} Q)]-G \oplus G . \tag{4.18}
\end{equation*}
$$

Proof: Expanding (3.7) yields

$$
\begin{array}{r}
-\left[A_{i}(\theta) Q_{i j}+Q_{v} A_{j}^{T}(\theta)\right]=\sum_{k=1}^{r}\left[G_{i k} Q_{k j}+Q_{i k} G_{j k}^{T}\right]+V_{v} \\
i, j=1, \cdots, r . \tag{4.19}
\end{array}
$$

Bounding the right-hand side of (4.19) from above using (2.4) yields for all $\boldsymbol{G} \in$

$$
\begin{aligned}
\| \sum_{k=1}^{r}\left[G_{\mu} Q_{k j}+Q_{j k} G_{j k}^{T}\right] & +V_{v} \|_{F} \\
& \leq \sum_{k=1}^{r}\left[S_{(1, k)} Q_{(k, j)}+Q_{(i, k)} \mathcal{S}_{(j, k)}\right]+V_{(1, j)}
\end{aligned}
$$

while bounding the left-hand side of (4.19) from below using (2.3) implies for all $\theta \in \boldsymbol{\theta}$

$$
\begin{aligned}
\left\|-\left[A_{i}(\theta) Q_{i j}+Q_{i} A_{j}^{T}(\theta)\right]\right\|_{F} & =\left\|\operatorname{vec}\left(A_{i}(\theta) Q_{i j}+Q_{i j} A_{j}^{T}(\theta)\right)\right\|_{F} \\
& =\|\left(A_{i}(\theta) \oplus A_{i}(\theta)\right) \text { vec } Q_{i j} \|_{F} \\
& \geq \sigma_{\min }\left(A_{j}(\theta) \oplus A_{i}(\theta)\right)\left\|\operatorname{vec} Q_{i j}\right\|_{F} \\
& =\sigma_{\min }\left(A_{j}(\theta) \oplus A_{i}(\theta)\right) Q_{(i, j)} \\
& \geq Q_{i, j)} Q_{(i, j)} .
\end{aligned}
$$

Combining the above inequalities yields (4.16).
Remark 4.2: Since $\mathcal{G} \geq 0$, the $r^{2} \times r^{2}$ matrix $A$ is an $N$ matrix [30].

## V. The Majorant Lyapunov Equation

In this section we interpret (4.16) as an equality rather than an inequality and consider the Lyapunov-type nonnegative matrix equation

$$
\begin{equation*}
\mathbb{Q} * \mathbb{Q}=S \mathbb{Q}+\mathbb{Q} \mathbb{S}^{T}+\boldsymbol{V} \tag{5.1}
\end{equation*}
$$

or, equivalently,

$$
\begin{equation*}
A \operatorname{vec} Q=\operatorname{vec} V . \tag{5.2}
\end{equation*}
$$

Note that since $Q$ and $v$ are symmetric a unique solution of (5.1) is necessarily symmetric.

Proposition 5.1: The following are equivalent:
i) $\lambda$ is an $M$ matrix;
ii) $\operatorname{det} \lambda \neq 0$ and $A^{-1} \geq \geq 0$;
iii) for each $r \times r$ symmetric $v \geq \geq 0$ there exists a unique $r$ $\times r \mathbb{Q} \geq \geq 0$ satisfying (5.1);
iv) there exist $r \times r$ symmetric $v>0$ and $r \times r$ symmetric $\mathbb{Q}$ $\geq \geq 0$ satisfying (5.1);
v) diag (vec $Q)-\left(I_{r} * \mathcal{G}\right) \oplus\left(I_{r} * S\right)>0$ and each diagonal matrix $\mathcal{D} \geq \geq$ diag (vec $Q)-\left(I_{r} * S\right) \oplus\left(I_{r} * G\right)$ satisfies

$$
\begin{equation*}
\rho\left(D^{-1}\left[S \oplus S-\left(I_{r} * S\right) \oplus\left(I_{r} * \mathcal{S}\right)\right]\right)<1 \tag{5.3}
\end{equation*}
$$

vi) for each $r \times r$ symmetric $Q_{0} \geq \geq 0$ and $r \times r$ symmetric $v$ $\geq \geq 0$, the sequence $\left\{\hat{Q}_{i}\right\}_{i=1}^{\infty}$ generated by

$$
\begin{align*}
Q & : \mathbb{Q}_{i+1}-\left(I_{r} * \mathcal{S}\right) \mathbb{Q}_{i+1}-\mathbb{Q}_{i+1}\left(I_{r} * \mathcal{S}\right) \\
& =\left(\mathcal{S}-I_{r} * S\right)_{i}+\mathbb{Q}_{1}\left(S-I_{r} * \mathcal{E}\right)^{T}+V, \quad i=0,1, \cdots, \tag{5.4}
\end{align*}
$$

converges;
vii) for each $r \times r$ symmetric $Q_{0} \geq \geq 0$ there exists $r \times r$ symmetric $\nabla>0$ such that the sequence $\left\{\mathcal{Q}_{i}\right\}_{i=1}^{\infty}$ generated by (5.4) converges.

Proof: Statements $i$ )-v) are equivalent to $i$ )-v) of Lemma 2.2. Clearly, vi) implies iii), and vii) implies iv). To show v) implies vi) and vii) note that $I_{\mu} *(\mathcal{G} \oplus \mathcal{S})=\left(I_{r} * S\right) \bullet\left(I_{r} * \mathcal{S}\right)$ and

Thus, (5.4) is equivalent to
$\operatorname{vec} \mathcal{Q}_{i+1}=\left[\operatorname{diag}(\operatorname{vec} Q)-\left(I_{r} * \mathcal{G}\right) \oplus\left(I_{r} * \mathcal{S}\right)\right]^{-1}$
$-\left[\mathcal{E} \oplus \mathcal{S}-I_{r^{2}} *(S \oplus S)\right]$ vec $\mathcal{Q}_{i}+[$ diag (vec $\left.Q)\right]^{-1}$ vec $\mathcal{V}$.
Thus, vi) and vii) follow from $v$ ) with $D=\operatorname{diag}(v e c Q)-\left(I_{r} *\right.$ 8) $\oplus\left(I_{r} * \mathcal{S}\right)$.

Since statements i)-vii) depend only upon $Q$ and $\mathcal{G}$ we have the following definition inspired by v -vii).

Definition 5.1: $(\mathbb{Q}, \mathcal{G})$ is stable if $\mathcal{A}$ is an $M$ matrix.
Remark 5.1: When $I_{r} * S=0$, i.e., when the local dynamics have no structured uncertainty, (5.4) simplifies to

$$
\begin{equation*}
Q=Q_{i+1}=\mathcal{Q}_{i}+Q_{i} S^{T}+V, \quad i=0,1, \cdots \tag{5.5}
\end{equation*}
$$

or, equivalently,

$$
\begin{equation*}
Q_{i+1}=Q^{H I} *\left(S Q_{i}+\mathcal{Q}_{i} G^{T}+V\right), \quad i=0,1, \cdots \tag{5.5a}
\end{equation*}
$$

The following result shows that for zero initial condition, the iterative sequence is monotonic.

Proposition 5.2: Suppose diag (vec $Q$ ) $-I_{2} *(S \oplus S)>0$. Then the sequence $\left\{\hat{Q}_{i}\right\}_{i=1}^{\infty}$ generated by (5.4) with $\hat{Q}_{0}=0$ and $\vartheta$ $\geq \geq 0$ is monotonically increasing.

Proof: To simplify notation we consider the case mentioned in Remark 5.1. Hence, assume $Q>0$. Clearly, if $Q_{0}=0$, then (5.5a) implies that $\hat{Q}_{1}=Q^{H /} * V \geq \geq 0$. Hence, $\hat{Q}_{1} \geq \geq \hat{Q}_{0}$. Defining $\Delta \mathcal{Q}_{i+1} \triangleq \hat{Q}_{i+1}-\mathcal{Q}_{i}$ (5.5a) yields

$$
\Delta \mathbb{Q}_{i+1}=\mathbb{Q}^{H l} *\left(S \Delta \mathbb{Q}_{i}+\Delta \mathbb{Q}_{i} S^{T}\right) .
$$

Since $\Delta Q_{1} \geq \geq 0$, the result follows from induction.
Remark 5.2: Proposition 5.2 is a particularly useful result in applications and can be utilized as follows. Setting $Q_{0}=0$, the sequence $\left\{\tilde{Q}_{i}\right\}$ can be evaluated by a simple numerical procedure. As will be shown in Theorem 5.1 below, each $\hat{Q}_{i}$ corresponds to a robust performance measure $\hat{\alpha}_{i}$. For practical purposes the increasing sequence $\left\{\hat{\alpha}_{i}\right\}$ can be generated until either convergence is attained (in which case $\hat{\alpha}=\lim _{i \rightarrow \infty} \hat{\alpha}_{i}$ is a robust performance bound) or a maximum permissible performance level is exceeded. In the latter case the question of convergence is irrelevant since the closed-loop system is known to either be unstable for some $G \in$ (i.e., $\dot{\alpha}=\infty$ ) or exceed acceptable performance specifications, thereby necessitating system redesign.
We now prove a comparison result for solutions of (5.1).
Lemma 5.1: Assume ( $Q, \mathcal{S}$ ) is stable, let $\hat{a}, \hat{S}$ be $r \times r$ nonnegative matrices where $\hat{Q}$ is symmetric, and assume that

$$
\begin{equation*}
Q \leq \leq \bar{Q}, \quad \widehat{S} \leq \leq S . \tag{5.6}
\end{equation*}
$$

Then ( $\hat{Q}, \hat{\mathcal{S}}$ ) is stable. Furthermore, let $r \times r$ symmetric 0 satisfy

$$
\begin{equation*}
\theta \leq \leq V_{1} \tag{5.7}
\end{equation*}
$$

let $\mathbb{Q}$ be the unique, nonnegative solution to (5.1), and let © be the unique solution to

$$
\begin{equation*}
\hat{Q} * \hat{Q}=S \hat{Q}+\hat{Q} S^{T}+\hat{O} \tag{5.8}
\end{equation*}
$$

Then if $\mathbb{Q} \geq \geq 0$, it follows that

$$
\begin{equation*}
\mathbb{Q} \leq \leq \mathbb{Q} . \tag{5.9}
\end{equation*}
$$

Proof: Since

$$
A \triangleq \text { diag }(\operatorname{vec} \mathbb{C})-\delta \cdot \xi
$$

is an $N$ matrix, $\AA$ is an $M$ matrix, and

$$
A-A=\operatorname{diag}(\operatorname{vec}(A-Q))+(8-\delta) \cdot(8-\delta) \geq \geq 0
$$

it follows from Lemma 2.3 that $\hat{A}$ is an $M$ matrix, and thus ( $\hat{A}, \mathcal{S}$ ) is stable. Next note that (5.1) and (5.8) imply

$$
\operatorname{vec}(Q-Q)=A^{-1}(A-A) \operatorname{vec} A+A^{-1} \operatorname{vec}(v-\theta) .
$$

Since $A-A \geq \geq 0, A^{-1} \geq \geq 0$ (see Lemma 2.2), $V-V \geq \geq$ 0 , and $Q \geq \geq 0$, it follows that (5.9) is satisfied.

Corollary 5.1: Suppose ( $\mathbb{Q}, \mathbf{8}$ ) is stable and let $Q$ be the unique, nonnegative solution to (5.1). Furthermore, let $G \in \in$ and $\theta \in \Theta$ be such that $A(\theta)+G$ is asymptotically stable and define $Q$ by (4.7) for $n \times n Q \geq 0$ satisfying (3.7). Then

$$
\begin{equation*}
Q \leq \leq \mathbb{Q} . \tag{5.10}
\end{equation*}
$$

Proof: By Proposition 4.2, Q satisfies the covariance blocknorm inequality (4.16). In the notation of Lemma 5.1 define

$$
\begin{equation*}
Q=Q, S=S, V=Q * Q-\left(S Q+Q S^{T}\right) \tag{5.11}
\end{equation*}
$$

so that (5.6) is satisfied and (4.16) implies (5.7). Note that with the notation (5.11), equation (5.8) has the unique solution $\hat{Q}=Q$ $\geq \geq 0$. Hence ( 5.9 ) implies ( 5.10 ).

Theorem 5.J: Assume $A$ is asymptotically stable, $\theta$ is continuously arcwise connected, and ( $Q, \mathcal{E}$ ) is stable. Then the homogeneous system (3.1) is robustly stable over and $\theta$, and the nonhomogeneous system (3.4) has the performance bound

$$
\begin{equation*}
\hat{\alpha}=\max _{\theta \in \Theta}\left\{\sum_{i=1}^{r}\left[\operatorname{tr}\left(Q_{i}(\theta) R_{i}\right)+2\left(\operatorname{tr} P_{i}(\theta)\right)(S \mathbb{Q})_{(i, i)}\right]\right\} \tag{5.12}
\end{equation*}
$$

where $n_{i} \times n_{i}$ nonnegative-definite $\boldsymbol{Q}_{i}(\theta)$ and $\boldsymbol{P}_{i}(\theta)$ satisfy

$$
\begin{gather*}
0=A_{i}(\theta) Q_{i}(\theta)+Q_{i}(\theta) A_{i}^{T}(\theta)+V_{i}  \tag{5.13}\\
0=A_{i}^{T}(\theta) P_{i}(\theta)+P_{i}(\theta) A_{i}(\theta)+R_{i} \tag{5.14}
\end{gather*}
$$

and $r \times r \bar{Q}$ is the unique, nonnegative solution to (5.1).
Proof: First note that since robust stability is independent of the disturbances, we can set $V=I_{n}$ for convenience in proving the first result. Hence, suppose (3.1) is not robustly stable. Since B is convex (see Remark 4.1), $A$ is asymptotically stable, and $\theta$ is continuously arcwise connected, there exist $G_{0} \in \theta$ and $\hat{\theta}:[0,1]$ $\rightarrow \theta$ such that $\hat{A}(\mu) \underset{\underline{\underline{1}}}{ } A(\hat{\theta}(\mu))+\mu G_{0}$ is asymptotically stable for all $\mu \in[0,1)$, and $\hat{A}(1)$ is not asymptotically stable. Define

$$
Q(\mu, t) \triangleq \int_{0}^{t} e^{\lambda(\mu) s} e^{\lambda^{T}(\mu) s} d s, \quad t \geq 0, \mu \in[0,1]
$$

which is monotonically increasing in the nonnegative-definite cone with respect to $t$. Clearly, the limit

$$
Q(\mu) \equiv \lim _{t \rightarrow \infty} Q(\mu, t), \quad \mu \in[0,1)
$$

exists and satisfies

$$
0=\hat{A}(\mu) Q(\mu)+Q(\mu) \hat{A}^{T}(\mu)+I_{n}, \quad \mu \in[0,1)
$$

Now define $r \times r$ nonnegative symmetric $Q(\mu)$ by

$$
Q(\mu)=\left\{\left\|Q_{U}(\mu)\right\|_{F}\right\}_{i, j=1}^{\prime}
$$

where $Q_{v}(\mu) \in \mathbb{R}^{n_{i} \times n_{j}}$ and $Q(\mu)$ is partitioned as in (4.5). By Corollary 5.1 with $\theta=\theta(\mu), G=\mu G_{0}, Q \triangleq Q(\mu), \mu \in[0,1)$. and $V=I_{n}$, it follows from (5.10) that

$$
\begin{equation*}
Q(\mu) \leq \leq \mathbb{Q}, \mu \in[0,1) \text {. } \tag{5.15}
\end{equation*}
$$

Hence, by (4.8), (5.15) implies

$$
\begin{equation*}
\|Q(\mu)\|_{F}=\|Q(\mu)\|_{F} \leq \| \mathbb{Q}_{F, \mu} \in[0,1) . \tag{5.16}
\end{equation*}
$$

On the other hand, for $\mu \in[0,1)$ it follows that

$$
\begin{aligned}
Q(\mu) & =Q(\mu)-Q(\mu, t)+Q(\mu, t)-Q(1, t)+Q(1, t) \\
& \geq Q(\mu, t)-Q(1, t)+Q(1, t)
\end{aligned}
$$

which implies, for arbitrary $x \in \mathbb{R}^{n}$,

$$
x^{T} Q(\mu) x \geq x^{T}[Q(\mu, t)-Q(1, t)] x+x^{T} Q(1, t) x
$$

Thus, by continuity of $Q(\mu, t)$ in $\mu$,

$$
\begin{equation*}
\lim _{\mu \rightarrow 1} x^{T} Q(\mu) x \geq x^{T} Q(1, t) x, x \in \beta^{n} \tag{5.17}
\end{equation*}
$$

Now, since $\hat{A}(1)$ is not asymptotically stable and $\left(\hat{A}(1), I_{n}\right)$ is stabilizable, it follows from [37, Proposition 3.2, p. 67] that for some $\bar{x} \in$ Bn $^{n}$,

$$
\lim _{t \rightarrow \infty} \bar{x}^{T} Q(1, t) \bar{x}=\infty
$$

Thus, by (5.17)

$$
\lim _{\mu \rightarrow 1} \bar{x}^{T} Q(\mu) x^{\tilde{x}}=\infty
$$

and thus

$$
\begin{equation*}
\lim _{\mu \rightarrow 1}\|Q(\mu)\|_{F}=\infty . \tag{5.18}
\end{equation*}
$$

However, (5.18) contradicts (5.16). Hence, (3.1) is robustly stable over and $\theta$.

To derive (5.12) note that since $R$ is block diagonal.

$$
J(G, \theta)=\sum_{i=1}^{r} \operatorname{tr} Q_{i} R_{i}=\sum_{i=1}^{r}\left(\operatorname{vec} Q_{i}\right)^{T} \operatorname{vec} R_{i}
$$

where $Q$ satisfies (3.7). Furthermore, (4.19) implies
vec $Q_{i}=-\left[A_{i}(\theta) \oplus A_{i}(\theta)\right]^{-1}$

$$
\left[\operatorname{vec} V_{i}+\sum_{k=1}^{\prime} \operatorname{vec}\left(G_{i k} Q_{k i} r Q_{i k} G_{i k}^{T}\right)\right]
$$

Hence, using Lemma 2.1,

$$
\begin{aligned}
J(G, \theta)= & \sum_{i=1}^{r}\left[\operatorname{tr}\left(Q_{i}(\theta) R_{i}\right)\right. \\
& \left.+\sum_{k=1}^{r}\left(\operatorname{vec}\left[G_{i k} Q_{k i}+Q_{i k} G_{i k}^{T}\right]\right)^{T} \operatorname{vec} P_{i}(\theta)\right] \\
= & \sum_{i=1}^{r}\left[\operatorname{tr}\left(Q_{i}(\theta) R_{i}\right)+\sum_{k=1}^{\prime} \operatorname{tr} P_{i}(\theta)\left(G_{i k} Q_{k i}+Q_{i k} G_{i k}^{T}\right)\right] \\
\leq & \sum_{i=1}^{r}\left[\operatorname{tr}\left(Q_{i}(\theta) R_{i}\right)\right. \\
& \left.+\sum_{k=1}^{r}\left(\operatorname{tr} P_{i}(\theta)\right) \sigma_{\max }\left(G_{i k} Q_{k i}+Q_{i k} G_{i k}^{T}\right)\right] \\
\leq & \sum_{i=1}^{r}\left[\operatorname{rr}\left(Q_{i}(\theta) R_{i}\right)+2\left(\operatorname{tr} P_{i}(\theta)\right) \sum_{k=1}^{\prime} \sigma_{\max }\left(G_{i k}\right) \sigma_{\max }\left(Q_{k i}\right)\right] \\
& \leq \sum_{i=1}^{r}\left[\operatorname{rr}\left(Q_{i}(\theta) R_{i}\right)+2\left(\operatorname{tr} P_{i}(\theta)\right) \sum_{k=1}^{r} \sigma_{\max }\left(G_{i k}\right)\left\|Q_{k i}\right\| f\right.
\end{aligned}
$$

$$
\begin{aligned}
& \leq \sum_{i=1}^{r}\left[\operatorname{tr}\left(Q_{i}(\theta) R_{i}\right)+2\left(\operatorname{tr} P_{i}(\theta)\right) \sum_{k=1}^{r} S_{(i, k)} \mathbb{Q}_{(k, i)}\right] \\
& =\sum_{i=1}^{r}\left[\operatorname{tr}\left(Q_{i}(\theta) R_{i}\right)+2\left(\operatorname{tr} P_{i}(\theta)\right)\left(S Q_{(i, i)}\right]\right.
\end{aligned}
$$

which yields (5.12).

## VI. Examples

We first confirm that the damped harmonic oscillator is asymptotically stable for all constant frequency perturbations. Hence, let

$$
r=1, n=n_{1}=2
$$

and

$$
A=\bar{A}_{1}=\left[\begin{array}{cc}
-\nu & \omega \\
-\omega & -\nu
\end{array}\right]
$$

where $\nu>0$ and $\omega \in$. To represent frequency uncertainty let 3 $=\{0\}, \theta=B, \theta=0$, and

$$
A(\theta)=A+\theta\left[\begin{array}{rr}
0 & 1 \\
-1 & 0
\end{array}\right]
$$

Note that $A(\theta)$ is stable for all $\theta \in B$ with poles $-\nu \pm j(\omega+\theta)$. Nore that $A(\theta)$ can be diagonalized by means of the unitary transformation

$$
\phi=\frac{1}{\sqrt{2}}\left[\begin{array}{cc}
1 & 1 \\
j & -j
\end{array}\right], \phi^{-1}=\frac{1}{\sqrt{2}}\left[\begin{array}{cc}
1 & -j \\
1 & j
\end{array}\right]
$$

so that

$$
\hat{A}(\theta) \equiv \phi^{-1} A(\theta) \phi=\left[\begin{array}{cc}
-\nu+j(\omega+\theta) & 0 \\
0 & -\nu-j(\omega+\theta)
\end{array}\right]
$$

Hence, using

$$
A(\theta) \oplus A(\theta)=\left(\phi^{-1} \otimes \phi^{-1}\right)(\tilde{A}(\theta) \oplus \hat{A}(\theta))(\phi \otimes \phi)
$$

it follows that

$$
\sigma_{\min }(A(\theta) \oplus A(\theta))=2 v, \quad \theta \in \mathcal{B}
$$

Defining [see (4.15)]

$$
Q=Q_{(1, n)}=2 v
$$

and $\mathcal{S}=0$, the scalar majorant Lyapunov equation (5.1) has the solution

$$
Q=v / 2 v
$$

where $V=\|V\|_{F}$. Choosing $V=I_{2}$ and noting that $Q=Q=2 \nu$ $>0$ is an $M$ matrix, Theorem 5.1 guarantees robust stability for all frequency variations $\theta \in \mathbb{R}$.

The next example has been chosen to demonstrate the robustness of a pair of nominally uncoupied oscillators with respect to uncertain coupling. Hence, let

$$
n=4, r=2, n_{1}=n_{2}=2
$$

and

$$
A_{i}=\left[\begin{array}{cc}
-v & \omega_{i} \\
-\omega_{i} & -v
\end{array}\right], \quad i=1,2
$$

where $\nu, \omega_{1}, \omega_{2} \geq 0$. Furthermore, let $\theta=\{\theta\}$ and

$$
\mathrm{E}=\left[\begin{array}{cc}
0 & \gamma_{12} \\
\gamma_{21} & 0
\end{array}\right]
$$

which denotes the fact that the local subsystem (oscillator) dynamics are assumed to be known. Since

$$
\sigma_{\min }\left(\bar{A}_{j} \oplus A_{i}\right)=\left[4 \nu^{2}+\left(\omega_{j}-\omega_{i}\right)^{2}\right]^{1 / 2}
$$

define

$$
Q=\left[\begin{array}{cc}
2 \nu & {\left[4 \nu^{2}+\left(\omega_{1}-\omega_{2}\right)^{2}\right]^{1 / 2}} \\
{\left[4 \nu^{2}+\left(\omega_{1}-\omega_{2}\right)^{2}\right]^{1 / 2}} & 2 \nu
\end{array}\right]
$$

Letting $V=I_{4}$ yields $V=2 I_{2}$. Solving (5.1) yields

$$
\begin{gathered}
\mathcal{Q}_{(1,1)}=\left(2 \nu^{2} \delta-\gamma_{12} \gamma_{21}+\gamma_{12}^{2}\right) / 2 \sqrt{2 \nu}\left(\nu^{2} \delta-\gamma_{12} \gamma_{21}\right), \\
\mathcal{Q}_{(1,2)}=\left(\gamma_{12}+\gamma_{21}\right) / 2 \sqrt{2}\left(\nu^{2} \delta-\gamma_{12} \gamma_{21}\right), \\
\mathcal{Q}_{(2,2)}=\left(2 \nu^{2} \delta-\gamma_{12} \gamma_{21}+\gamma_{21}^{2}\right) / 2 \sqrt{2} \nu\left(\nu^{2} \delta-\gamma_{12} \gamma_{21}\right)
\end{gathered}
$$

where

$$
\delta \equiv\left[1+\delta^{2}\right]^{12}, \delta \doteq\left(\omega_{1}-\omega_{2}\right) / 2 v
$$

Clearly, $\mathcal{Q}$ is nonnegative if and only if

$$
\begin{equation*}
\gamma_{12} \gamma_{21}<\nu^{2} \delta \tag{6.1}
\end{equation*}
$$

The bound (6.1) characterizes the magnitude of coupling uncertainty for which stability is guaranteed. Note that the parameter $\delta$ is a measure of the frequency separation of the oscillators relative to the damping. When $\delta>1$. (6.1) becomes asymptotically

$$
\begin{equation*}
\gamma_{12} \gamma_{21}<\frac{\nu}{2}\left|\omega_{1}-\omega_{2}\right| \tag{6.2}
\end{equation*}
$$

which confirms the infuitive expectation that robust stability is proportional to damping and subsystem frequency separation. This result does not appear to be predictable from quadratic or vector Lyapunov functions.

To evaluate the conservatism inherent in the bound (6.1) we solve for the actual stability region. To render the calculation tractable we assume that $G_{12}$ and $G_{21}$ have the structured form

$$
G_{i j}=\left[\begin{array}{cc}
\alpha_{i j} & \beta_{i j}  \tag{6.3}\\
-\beta_{i j} & \alpha_{i j}
\end{array}\right]
$$

By constraining (6.3) the set of coupling variations is reduced, which may or may not lead to a larger stability region. Thus, our estimate of conservatism may itself be conservative, i.e., the actual conservatism may indeed be less than the following analysis indicates. However, without (6.3) the development becomes intractable. This calculation will thus be called semiexact.

By considering the characteristic equation for $A+G$, lengthy manipulation shows that $\bar{A}+G$ is stable if and only if

$$
\begin{equation*}
\gamma_{12} \gamma_{21}<2 \nu^{2}\left(-\epsilon+\left[1+\delta^{2}\left(1-\epsilon^{2}\right)\right]^{1 / 2}\right) /\left(1-\epsilon^{2}\right) \tag{6.4}
\end{equation*}
$$

where $\epsilon \in(0,1]$ is the smallest positive real root of

$$
\begin{equation*}
\epsilon=\left(1+\epsilon^{2}\right)\left[1+\delta^{2}\left(1-\epsilon^{2}\right)\right]^{1 / 2} /\left[2+\delta^{2}\left(1-\epsilon^{2}\right)\right] . \tag{6.5}
\end{equation*}
$$

The majorant bound (6.1) and semiexact bound (6.4) are illustrated in unified form in Fig. 1. For $\delta>1$ note that $\epsilon=$ $O\left(\delta^{-1}\right)$ and thus (6.4) becomes asymptotically

$$
\begin{equation*}
\gamma_{12} \gamma_{21}<\nu\left|\omega_{1}-\omega_{2}\right| . \tag{6.0}
\end{equation*}
$$



Fig. 1.


Hence, for large $\delta$ the majorant bound (6.2) is, at worst, and the system has the performance bound conservative by a factor of 2 compared to the semiexact bound.

To determine the performance bound ( 5.12 ) set $R=I_{4}$. Hence.

$$
\begin{equation*}
\dot{\alpha}=J_{\operatorname{cosin}}+\sqrt{2}\left(\rho_{12}+\rho_{21}\right)^{2 / h}\left(1-2 \rho_{12} \rho_{21}\right) \tag{6.7}
\end{equation*}
$$ it can be shown that

where

$$
J_{\text {nom }}=2 / \nu
$$

$$
\rho_{12}=\gamma_{12} / \sqrt{2} \kappa^{1 / 2}, \rho_{21}=\gamma_{21} / \sqrt{2} \kappa^{1 / 2} .
$$

On the other hand, the semiexact calculation yields

$$
J_{\max }=\max _{\lambda \in \rho 0,1)}\left\{\left[\rho_{12}^{2}+\rho_{21}^{2}+2 \rho_{12} \rho_{21} \lambda+2 \delta\left(\rho_{12} \rho_{21}\right)^{2}\left(1-\lambda^{2}\right)\right] /\right.
$$

$$
\begin{equation*}
\left.\left[2 \delta-4 \rho_{12} \rho_{21} \lambda-2 \delta\left(\rho_{12} \rho_{21}\right)^{2}\left(1-\lambda^{2}\right)\right]\right\} \tag{6.8}
\end{equation*}
$$

Fig. 2 compares the semiexact worst-case performance (6.8) to the majorant Lyapunov equation bound (6.7). To efficiently illustrate the results the data are specialized to the case $\rho_{12}=\rho_{21}$. Note that the semiexact performance is plotted for several values of $\delta$ because of the explicit dependence of (6.8) on $\delta$ via $\delta$.

## Acknowledgment

The authors wish to thank J. Straehla for excellent typing of the original manuscript and B. Bland for executing the numerical calculations.

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APPENDIX H: Robust Synthesis: Linear Bound
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# Robust, Reduced-Order, Nonstrictly Proper State Estimation Via the Optimal Projection Equations with Guaranteed Cost Bounds 

WASSIM M. HADDAD and DENNIS S. BERNSTEIN

Abstract-A state-estimation design problem involving parametric plant uncertainties is considered. An estimation error bound suggested by multiplicative white noise modeling is utilized for guaranteeing robust estimation over a specified range of parameter uncertainties. Necessary conditions which generalize the optimal projection equations for reducedorder state estimation are used to characterize the estimator which minimizes the error bound. The design equations thus effectively serve as sufficient conditions for synthesizing robust estimators. Additional features include the presence of a static estimation gain in conjunction with the dynamic (Kalman) estimator to obtain a nonstrictly proper estimator.

## 1. INTRODUCTION

As is well known [1]-[12]. the performance of optimal filters based upon nominal parameter values may be severely degraded in the presence of parameter deviations. Thus. it is desirable to obtain robust state estimators which provide acceptable performance over the range of parametric uncertainty. The approach of the present paper is related to the guaranteed cost approach developed for control in [13]. [14] and applied to estimation in [3]. Specifically, the main idea is to bound the effect of the uncertain parameters on the estimation error over the uncentainty range and then choose estimator gains to minimize the estimation bound. Thus. the actual estimation error is guaranteed to lie belou the prescribed upper bound.

The technique used to determine minimizing estimator gains is a generalization of the optimal projection equations for reduced-order state estimation [15]. Thus. the results of the present paper effectively extend the results of [15] to the case of parameter uncertainties. It should be noted that the optimal projection equations, which are necessary conditions for optimality, now serve as sufficient conditions for robust estimation by virtue of the fact that a bound on the estimation error is being minimized rather than the estimation error itself. The bound utilized in the present paper was originally suggested by multiplicative white noise

[^21]modeling and was used in [16]-[18] for constructing Lyapunov functions for robust fixed-order dynamic compensation. A similar bound was used for full-state feedback in [19].

An additional feature of the present paper is the inclusion of a static feedback gain in conjunction with the dynamic estimator. Thus, the results of the present paper represent a generalization of standard results to the case of nonstriclly proper estimation. Similar treatments in the context of multiplicative noise models were given in [10] and [11] for discrete-time and continuous-time systems, respectively.

## II. Notation and Defintions

Note: All matrices have real entries.
the state-estimation error criterion
$J\left(A_{c}, B_{c}, C_{0}, D_{f}\right) \triangleq \sup$

$$
\mathfrak{B A} . \Delta C \in u
$$

$$
\begin{equation*}
\lim _{t \rightarrow \infty} \mathbb{s u p}\left[L x(t)-y_{e}(t)\right]^{\top} R\left[L x(t)-y_{e}(t)\right] \tag{3.6}
\end{equation*}
$$

is minimized.
For each estimator $\left(A_{e}, B_{e}, C_{e}, D_{c}\right)$ and system variation $(\Delta A, \Delta C) \in$ U, the disturbed augmented system (3.1)-(3.5) is given by

$$
\begin{equation*}
\dot{x}(t)=(A+\Delta \tilde{A}) \tilde{x}(t)+\tilde{w}(t), \quad t \in[0, \infty) \tag{3.7}
\end{equation*}
$$

where $\tilde{x}(t) \triangleq\left[x^{\top}(t), x_{e}^{T}(t)\right]^{T}$ and $\tilde{w}(t)$ has intensity $\bar{\nabla} \in \mathbb{N}^{\boldsymbol{N}}$. The cost can be expressed in terms of the second-moment matrix.

|  |
| :---: |
|  |
| $\oplus$, $\otimes$ |
| Sr, Nr, Pr |
| $Z_{1} \leq Z_{21} Z_{1}<Z_{2}$ |
| $n, l, f, n_{f}, p, q ;{ }_{\text {n }}$ |
| $x, y, y, y_{e}, x_{e}, \vec{x}$ |
| $A, \Delta A ; C, \Delta C$ |
|  |
| $A_{e}, B_{e}, C_{e}, D_{e}$ |
| $A, \Delta A$ |
| L, R |
| R |
| $w_{1}(\cdot), w_{2}(\cdot)$ |
| $V_{1}, V_{2}$ |
| $V_{12}$ |
| $w(\cdot), D$ |
| $\boldsymbol{\alpha}$ |
| $\alpha_{i}$ |
| $\mathrm{g}_{\boldsymbol{i}}$ |
| $\gamma_{i}$ |
| $A_{a}, A_{\text {ea }}$ |
| $A_{\text {a }}$ |

Real numbers, $r \times s$ real matrices, $\mathbb{R}^{r \times 1}$.
$r \times r$ identity matrix, transpose, expected value.
Kronecker sum, Kronecker product [20].
$r \times r$ symmetric, nonnegative-definite, positive-definite matrices.
$Z_{2}-Z_{1} \in \mathbb{N}^{\prime}, Z_{2}-Z_{1} \in \mathcal{P}^{\prime}, Z_{1}, Z_{2} \in \Phi^{\prime}$.
Positive integers; $n+n_{\mathrm{e}}$.
$n, l, l, q, n_{e}, n$-dimensional vectors.
$n \times n$ matrices; $I \times n$ matrices.
$\Gamma \times n$ matrix.
$n_{e} \times n_{e}, n_{e} \times 1, q \times n_{e}, q \times r$ matrices.
$\left[\begin{array}{cc}A & 0 \\ B_{e} C & A_{\varepsilon}\end{array}\right],\left[\begin{array}{cc}\Delta A & 0 \\ B_{r} \Delta C & 0\end{array}\right]$.
$q \times n$ matrix, estimation-error weighting in $P Q$.
$\left[\begin{array}{cc}L^{\top} R L-L^{\top} R D_{e} C-C^{T} D_{e}^{r} R L+C^{\top} D_{e}^{r} R D_{e} C & -L^{\top} R C_{e}+C^{\top} D_{e}^{r} R C_{e} \\ -C_{e}^{\top} R L+C_{e}^{\top} R D_{e} C & C_{e}^{\top} R C_{e}\end{array}\right]$.
$n, 1$-dimensional white noise.
Intensity of $w_{1}(\cdot), w_{2}(\cdot) ; V_{1} \in \mathbb{M}^{n}, V_{2} \in \mathbb{P}^{\prime}$.
Cross intensity of $w_{t}(\cdot), w_{z}(\cdot)$.
$\left[\begin{array}{c}w_{1}(\cdot) \\ B_{\mathrm{e}} w_{2}(\cdot)\end{array}\right],\left[\begin{array}{cc}V_{1} & V_{12} B_{!}^{T} \\ B_{\mathrm{e}} V_{12}^{T} & B_{\mathrm{e}} V_{2} B_{?}^{T}\end{array}\right]$.
Positive number.
Positive number, $i=1, \cdots, p$.
Real number, $i=1, \cdots, p$.
$\alpha_{1}^{2 / \alpha}, i=1, \cdots, p$.
$A+(\alpha / 2) I_{n}, A_{e}+(\alpha / 2) I_{n_{e}}$.
$\left[\begin{array}{cc}A_{a} & 0 \\ B_{r} C & A_{80}\end{array}\right]$.

## iII. Robust Estimation Problem

Let $u \subset \mathbb{S}^{n \times n} \times \mathbb{N}^{1 \times n}$ denote the set of uncertain perturbations ( $\Delta A$, $\Delta C$ ) of the nominal plant matrices $A$ and $C$.
Robust Estimation Problem: For fixed $n_{e} \leq n$, determine ( $A_{1,}, B_{e}$, $C_{e}, D_{e}$ ) such that, for the system consisting of the $n$ th-order disturbed plant

$$
\begin{equation*}
\dot{x}(l)=(A+\Delta A) x(t)+w_{1}(l), \quad l \in[0, \infty) \tag{3.1}
\end{equation*}
$$

noisy and nonnoisy measurements

$$
\begin{gather*}
y(t)=(C+\Delta C) x(t)+w_{2}(t),  \tag{3.2}\\
y(t)=C x(t) \tag{3.3}
\end{gather*}
$$

and $n_{\text {a }}$ th-order nonstrictly proper state estimator

$$
\begin{align*}
& x_{0}(t)=A_{e} x_{0}(f)+B_{e} y(t)  \tag{3.4}\\
& y_{4}(f)=C_{e} x_{t}(f)+D_{e} f(t) \tag{3.5}
\end{align*}
$$

Proposition 3.1: For given $\left(A_{e}, B_{e}, C_{e}, D_{e}\right)$ and $(\Delta A, \Delta C) \in \mathcal{U}$, the second-moment matrix
satisfies

$$
\begin{equation*}
\dot{Q}_{\Delta A}(l)=(\tilde{A}+\Delta \bar{A}) \varrho_{\Delta \lambda}(t)+Q_{\Delta X}(t)(\tilde{A}+\Delta A)^{T}+D, \quad l \in[0, \infty) \tag{3.9}
\end{equation*}
$$

Furthermore,

$$
\begin{equation*}
J\left(A_{e}, B_{e}, C_{e}, D_{e}\right) \in \sup _{i \Delta A, \Delta C\rangle \in q} \lim _{\substack{ \\\rightarrow \infty}} \text { tr } \mathcal{Q u p}_{\Delta \Lambda}(f) R \tag{3.10}
\end{equation*}
$$

## IV. Sufficient Conditions for Robust performance

The following result is immediate.
Lemma 4.I: Suppose $X+\Delta X$ is stable for all $(\Delta A, \Delta C) \in \mathcal{U}$. Then

$$
\begin{equation*}
J\left(A_{1}, B_{r}, C_{r}, D_{d}\right)=\operatorname{mup}_{1 \Delta A, \Delta \cap \in v} \operatorname{tr} Q_{\Delta A} R \tag{4.1}
\end{equation*}
$$

where $\phi_{\Delta \lambda} \in \mathbb{N}^{A}$ is the unique solution to

$$
\begin{equation*}
0=(A+\Delta X) Q_{\Delta X}+Q_{\Delta \lambda}(A+\Delta X)^{T}+D . \tag{4.2}
\end{equation*}
$$

We seek upper bounds for $J\left(A_{c,}, B_{e}, C_{e}, D_{c}\right)$.
Theorem 4.1: Let $\Omega: \mathbb{N}^{n} \times \mathbb{R}^{n e} \times t^{\prime} \rightarrow \mathbb{S}^{n}$ be such that
$(\Delta A, \Delta C) \in$ U.

$$
\begin{equation*}
\Delta \bar{A} Q+Q \Delta A^{T} \leq Q\left(Q, B_{e}\right),\left(Q, B_{e}\right) \in \mathbb{N}^{n} \times \mathbb{B}^{n_{e} \times 1} \tag{4.3}
\end{equation*}
$$

and, for given ( $A_{e}, B_{e}, C_{e}, D_{e}$ ), suppose there exists $Q \in \mathbb{N}^{1}$ satisfying

$$
\begin{equation*}
0=\tilde{A} Q+Q \bar{A}^{r}+\Omega\left(Q . B_{e}\right)+\hat{V}, \tag{4.4}
\end{equation*}
$$

and suppose the pair $\left(\tilde{V}^{1 / 2}, \bar{A}+\Delta \bar{A}\right)$ is stabilizable for all $(\Delta A, \Delta C) \in$ $\mathcal{U}$. Then $A_{\mathrm{e}}$ is asymptotically stable, $A+\Delta \mathcal{A}$ is asymptotically stable for all $(\Delta A, \Delta C) \in \mathcal{U}$,

$$
\begin{equation*}
\oint_{\Delta \pi} \leq Q . \quad(\Delta A, \Delta C) \in u \tag{4.5}
\end{equation*}
$$

where $Q_{A X}$ satisfies (4.2) and

$$
\begin{equation*}
J\left(A_{e}, B_{e}, C_{e}, D_{e}\right) \leq \operatorname{tr} Q R . \tag{4.6}
\end{equation*}
$$

Proof: For all $(\Delta A, \Delta C) \in \mathcal{U},(4.4)$ is equivalent to

$$
\begin{equation*}
0=(\bar{A}+\Delta \bar{A}) Q+Q(\bar{A}+\Delta \bar{A}))^{r}+\boldsymbol{\Psi}\left(\mathbb{Q}, B_{c}, \Delta \bar{A}\right)+\bar{V} \tag{4.7}
\end{equation*}
$$

where

$$
\Psi\left(Q, B_{e}, \Delta \tilde{A}\right) \triangleq \Omega\left(Q, B_{\mathrm{c}}\right)-\left(\Delta \tilde{A} Q+Q \Delta A^{r}\right) .
$$

Note that by (4.3), $\Psi\left(Q, B_{e}, \Delta \bar{A}\right) \geq 0$ for all $(\Delta A, \Delta C) \in \mathcal{U}$. Since ( $\nabla^{1 / 2}, \tilde{A}+\Delta A$ ) is stabilizable for all $(\Delta A, \Delta C) \in \cup$, it follows from [21, Theorem 3.6] that $\left(\left(\bar{V}+\Psi\left(Q, B_{e}, \Delta \bar{A}\right)\right)^{1 / 2}, A+\Delta \bar{A}\right)$ is stabilizable for all $(\Delta A, \Delta C) \in \mathfrak{U}$. Hence, [21, Lemma 12.2] implies $A+\Delta A$ is ssymptotically stable for all $(\Delta A, \Delta C) \in U$. Since $\bar{A}+\Delta \tilde{A}$ is lower block triangular, $A_{e}$ is asymptotically stable and $A+\Delta A$ is asymptotically stable for all $(\Delta A, \Delta C) \in \mathcal{U}$. Next, (4.7) minus (4.2) yields

$$
0=(\bar{A}+\Delta \bar{A})\left(Q-Q_{\Delta A}\right)+\left(Q-Q_{\Delta A}\right)(\bar{A}+\Delta \bar{A})^{T}+\Psi\left(Q, B_{e}, \Delta \bar{A}\right)
$$

or, equivalently (since $\bar{A}+\Delta \tilde{A}$ is asymptotically stable).

$$
Q-Q_{د A}=\int_{0}^{\infty} e^{(\lambda+\Delta A) r} \Psi\left(Q, B_{r}, \Delta \tilde{A}\right) e^{(\lambda+\Delta \dot{A}) T_{t}} d t \geq 0
$$

which implies (4.5). Finally, (4.5) and (4.1) yield (4.6).

## V. Uncertainty Structure and Guaranteed Cost bound

The uncertainty set $\mathcal{U}$ is assumed to be of the form

$$
\begin{align*}
u & \triangleq\left\{(\Delta A, \Delta C) \in \beta^{n \times n} \times n^{i \times n}: \Delta A=\sum_{i=1}^{p} \sigma_{i} A_{1}, \Delta C\right. \\
& \left.=\sum_{i=1}^{p} \sigma_{1} C_{i}, \sum_{i=1}^{p} \sigma_{i}^{2} / \alpha_{i}^{2} \leq 1\right\} \tag{5.1}
\end{align*}
$$

where, for $i=1, \cdots, p: A_{i} \in \mathbb{R}^{n \times n}$ and $C_{i} \in \mathbb{R}^{1 \times n}$ are fixed matrices denoting the structure of the parametric uncertainty in the dynamics and measutement matrices; $\alpha_{i}$ is a given positive number; and $\sigma_{i}$ is an uncertain real parameter. In practice, the form of $\Delta A$ and $\Delta C$ permits the modeling of linear parameter uncertainties of arbitrary structure. Note that the uncertain parameters $\sigma_{i}$ are assumed to lie in a specified ellipsoidal region in $1^{p}$. The augmented system thus has structured uncertainty of the form

$$
\begin{equation*}
\Delta X=\sum_{1=1}^{\infty} \sigma_{1} A_{1} \tag{5.2}
\end{equation*}
$$

where

$$
A_{i} \triangleq\left[\begin{array}{cc}
A_{i} & 0 \\
B_{e} C_{i} & 0
\end{array}\right], \quad i=1, \cdots, p
$$

Remark 5.1: Note that (5.1) allows a particular parameter $\sigma_{;}$to appear in both $\Delta A$ and $\Delta C$. Thus, it is possible to consider the case in which the uncertainties $\Delta A$ and $\Delta C$ are known to be correlated. Of course, for a given $i, A_{i}$ or $C_{i}$ can be set to zero so that the similar form of $\Delta A$ and $\Delta C$ represents no restriction.

We now specify the bounding function $\Omega$ satisfying (4.3).
Proposition 5.1: Let $\alpha$ be an arbitrary positive scalar. Then the function

$$
\begin{equation*}
Q\left(Q, B_{e}\right) \triangleq \alpha Q+\alpha^{-1} \sum_{i=1}^{n} \alpha_{i}^{2} A_{i} Q A_{1}^{T} \tag{5.3}
\end{equation*}
$$

satisfies (4.3) with $\mathcal{U}$ given by (5.1).
Proof: Note that

$$
\begin{aligned}
0 & \leq \sum_{i=1}^{D}\left[\left(\alpha^{1 / 2} \sigma_{i} / \alpha_{i}\right) I_{A}-\left(\alpha_{i} / \alpha^{1 / 2}\right) A_{i}\right] \cdot Q\left[\left(\alpha^{1 / 2} \sigma_{i} / \alpha_{i}\right) I_{A}-\left(\alpha_{i} / \alpha^{1 / 2}\right) A_{i}\right]^{T} \\
& =\alpha \sum_{i=1}^{p}\left(\sigma_{i}^{2} / \alpha_{i}^{2}\right) Q+\alpha^{-1} \sum_{i=1}^{p} \alpha_{i}^{2} A_{i} Q A_{i}^{T}-\sum_{i=1}^{p} \sigma_{i}\left(A_{i} Q+Q \bar{A}_{i}^{T}\right)
\end{aligned}
$$

which, since $\sum_{i=1}^{p} \sigma_{i}^{2 /} \alpha_{i}^{2} \leq 1$, implies (4.3).
Remark 5.2: Note that with (5.3), the modified Lyapunov equation (4.4) becomes

$$
\begin{equation*}
0=A_{a} Q+Q A_{\alpha}^{T}+\sum_{i=1}^{p} \gamma_{i} A_{i} Q A_{i}^{T}+\nabla . \tag{5.4}
\end{equation*}
$$

## VI. The auxiliary Minimization Problem

Our goal is to minimize the error bound (4.6).
Auxiliary Minimization Problem: Determine ( $\left.Q, A_{f}, B_{e}, C_{e}, D_{e}\right)$ with $\mathbb{Q} \in \mathbb{N}^{A}$ which minimizes

$$
\begin{equation*}
J\left(Q, A_{e}, B_{c}, C_{e}, D_{e}\right) \triangleq \operatorname{tr} Q R \tag{6.1}
\end{equation*}
$$

subject to (5.4) and

$$
\begin{equation*}
\left(\bar{V}^{1 / 2}, \tilde{A}+\Delta \bar{A}\right) \text { is stabilizable, } \quad(\Delta A, \Delta C) \in U \tag{6.2}
\end{equation*}
$$

Proposition 6.1: If (Q, $\left.A_{e}, B_{e} . C_{e}, D_{e}\right)$ satisfies (5.4) and (6.2) with $Q$ $\geq 0$, then $\bar{A}+\Delta \tilde{A}$ is asymptotically stable for all $(\Delta A, \Delta C) \in \cup$ and

$$
\begin{equation*}
, \quad J\left(A_{e}, B_{e}, C_{e}, D_{e}\right) \leq g\left(Q, A_{f}, B_{e}, C_{e}, D_{e}\right) \tag{6.3}
\end{equation*}
$$

Proof: With $\Omega$ given by (5.3), Proposition 5.1 implies that (4.3) is satisfied. Hence, with (6.2), the hypotheses of Theorem 4.1 are satisfied so that the system (3.7) is stable over $u$ with estimation bound (4.6). Note that ( 6.3 ) is merely a restatement of (4.6).
Remark 6.I: The conservatism of the bound (6.3) is difficult to predict for two reasons. First, the overbounding (4.3) holds with respect to the partial ordering of the nonnegative-definite matrices for which no scalar measure of conservatism is available. And second, the bound (4.3) is required to hold for all nonnegative-definite matrices $Q$ and estimator gains $B_{e}$. The conservatism will thus depend upon the actual values of $Q$ and $B_{c}$ determined by solving (5.4).

## VII. Necessary Conditions for the Auxillary Minlmbzation PROBLEM

Rigorous application of the Lagrange multiplier technique requires additional technical assumptions. Specifically, we further restrict ( $Q A_{1}$,
$B_{e}, C_{e}, D_{e}$ ) to the open set
$8 \triangleq\left\{\left(Q_{,}, A_{f}, B_{r}, C_{f}, D_{f}\right): Q \in \mathbb{P}^{\prime}\right.$.
$\overline{\mathbb{a}}$ is asymprotically stable,
( $A_{e}, B_{e}, C_{e}$ ) is controliable and observable, and $C\left(Q_{1}-Q_{12} Q_{2}^{-1} Q_{12}^{T} C^{\top}>0\right\}$
where

$$
\mathbb{Q} \triangleq A_{\mathrm{a}} \oplus \mathcal{A}_{\mathrm{a}}+\sum_{i=1}^{\infty} \gamma_{1} \bar{A}_{i} \otimes A_{i}
$$

and $Q$ is partitioned as

$$
\mathfrak{Q} \triangleq\left[\begin{array}{ll}
Q_{1} & \boldsymbol{Q}_{12} \\
\boldsymbol{Q}_{12}^{r} & \boldsymbol{Q}_{2}
\end{array}\right]
$$

where $Q_{1}, Q_{12}$, and $Q_{2}$ are $n \times n_{1} n \times n_{e}$, and $n_{e} \times n_{e}$, respectively. As shown in [II], $Q_{2}$ is invertible since ( $A_{e}, B_{e}$ ) is controllable. The definiteness condition holds when $\mathcal{C}$ has full row rank and $Q$ is positive definite. As shown in [11], this condition implies the existence of the projection $f$ defined below.
Remark 7.1: Proposition 6.1 shows that the constraint ( $Q, A_{e}, B_{e}, C_{e}$, $\left.D_{e}\right) \in \delta$ is not required for robust estimation. As can be seen from the proof given in [11], the set $\delta$ constimutes sufficient conditions under which the Lagrange multiplier technique is applicable to the auxiliary minimization problem. Specifically, $Q \in \mathbb{P}^{*}$ replaces $Q \in \mathbb{N}^{*}$ by an open set constraint, while asymptotic stability of $\mathbb{C}$ serves as a normality condition which further implies that the dual $\mathcal{Q}$ of $Q$ is nonnegative definite. Thus, it is not necessary for guaranteed robust estimation that an admissible quadruple obtained by solving the necessary conditions actually be shown to be an element of $\delta$.
The following factorization lemma is needed for the statement of the main result. For details, see [15].
Lemma 7.1: If $Q, P \in \mathbb{N}{ }^{n}$ and rank $\phi \mathcal{P}=n_{e}$, then there exist $n_{e} \times$ $n G . \Gamma$, and $n_{e} \times n_{e}$ invertible $M$ such that

$$
\begin{gather*}
Q P=G^{T} M \Gamma  \tag{7.1}\\
\Gamma G^{r}=I_{\text {me }} \tag{7.2}
\end{gather*}
$$

Recall from [15] that

$$
\begin{equation*}
T \triangleq Q P(Q P)^{4}=G^{\top} \Gamma \tag{7.3}
\end{equation*}
$$

is an oblique projection. Define the complementary projection $\tau_{1} \triangleq I_{n}-$ $\tau$ and call ( $G, M, \Gamma$ ) satisfying (7.1), (7.2) a projective factorization of $Q \boldsymbol{Q}$. Furthermore, for arbitrary $Q, Q \in \beta^{n \times n}$, define the notation

$$
\begin{gathered}
V_{2 s} \triangleq V_{2}+\sum_{i=1}^{\infty} \gamma_{i} C_{i}(Q+Q) C_{i}^{T} \\
Q \triangleq Q C^{T}+V_{12}+\sum_{i=1}^{f} \gamma_{i} A_{i}(Q+Q) C_{i}^{T}
\end{gathered}
$$

Theorem 7.1: If ( $\left.Q, A_{e}, B_{e}, C_{e}, D_{e}\right) \in \&$ solves the auxiliary minimization problem with $\mathcal{U}$ given by (5.1) and $\Omega$ given by (5.3), then there exist $Q, Q, P \in \mathbb{8 9}$ " such that. Ior some projecuve factorizalion ( $G$, $M, \Gamma)$ of $Q \boldsymbol{P},\left(Q, A_{e}, B_{e}, C_{f}, D_{e}\right)$ are given by

$$
\begin{align*}
& Q=\left[\begin{array}{cc}
Q+Q & Q \Gamma^{r} \\
\Gamma Q & \Gamma Q \Gamma^{r}
\end{array}\right] .  \tag{7.4}\\
& A_{0}=\Gamma\left(A-Q V_{2 S}^{-1} C\right) G^{r}, \tag{7.5}
\end{align*}
$$

$$
\begin{gather*}
B_{\mathrm{f}}=\Gamma Q_{d} V_{2 s}^{-1},  \tag{7.6}\\
C_{f}=L \hat{q}_{L} G^{\top},  \tag{7.7}\\
D_{f}=L Q C^{\top}\left(C Q C^{\top}\right)^{-1}, \tag{7.8}
\end{gather*}
$$

and such that $\boldsymbol{Q}, \boldsymbol{Q}$, and $\boldsymbol{P}$ satisfy

$$
\begin{equation*}
\operatorname{rank} \phi=\operatorname{rank} \phi=\operatorname{rank} Q P=n_{e} \tag{7.12}
\end{equation*}
$$

where

$$
\begin{equation*}
\mathfrak{f} \triangleq Q C^{T}\left(C Q C^{T}\right)^{-1} C, \hat{\tau}_{+} \triangleq I_{n}-\hat{\tau} \tag{7.13}
\end{equation*}
$$

Furthermore, the auxiliary cost is given by

$$
\begin{equation*}
g\left(Q_{,} A_{e}, B_{e}, C_{e}, D_{e}\right)=\operatorname{tr} Q \hat{\tau}_{1}^{\top} L^{\top} R L \hat{\tau}_{1} . \tag{7.14}
\end{equation*}
$$

Conversely, if there exist $Q, Q, P \in \mathbb{N}^{n}$ satisfying (7.9)-(7.12), then ( $Q$, $\left.A_{e}, B_{e}, C_{e}, D_{e}\right)$ given by (7.4)-(7.8) satisfy (5.4) with $Q \in \mathbb{N}^{\wedge}$ and with $g\left(Q_{0} A_{f}, B_{e}, C_{e}, D_{e}\right)$ given by (7.14).

Proof: The derivation requires only a minor modification of the derivation given in [11]. The only change involves treatment of $\boldsymbol{A}_{\alpha}$ in place of $\boldsymbol{A}$.

Remark 7.1: The necessary conditions given in Theorem 7.1 directly generalize the result given in [15]. To recover the result of [15], set $A_{i}=$ $0, C_{i}=0, i=1, \cdots, p$ (to delete the plant uncertainties), and set $\mathcal{C}=0$ (to eliminate the static estimation term $D_{e}$ ). It follows from the proof given in [11] that $\hat{C}=0$ yields $\hat{\boldsymbol{r}}=0$, and thus $\dot{\tau}_{\perp}=I_{n}$.
Remark 7.2: Note that $Q$ given by (7.4) is nonnegative definite.

## Vili. Sufficient Conditions for Robust, Reduced-Order estimation

The main result guaranteeing robust estimation can now be stated.
Theorem 8.1: Suppose there exist $Q, \phi, \beta \in M^{n}$ satisfying (7.9)(7.12), let $A_{e}, B_{r}, C_{r}, D_{e}$ be given by (7.5)-(7.8), and suppose that ( $\bar{V}^{1 / 2}$, $A+\Delta A)$ is stabilizable for all $(\Delta A, \Delta C) \in \mathcal{U}$ with $\mathfrak{U}$ given by ( 5.1 ). Then $A_{\&}$ is asymptotically stable, $A+\Delta A$ is asymptotically stable for all $(\Delta A, \Delta C) \in \mathcal{U}$, and the estimation error satisfies the performance bound

$$
\begin{equation*}
J\left(A_{\mathrm{f}}, B_{\mathrm{e}}, C_{e}, D_{e}\right) \leq \operatorname{tr} Q \hat{\mathrm{~T}}_{\mathrm{L}}^{\top} L^{\top} R L \hat{\tau}_{4} . \tag{8.1}
\end{equation*}
$$

Proof: Theorem 7.1 and Remark 7.2 imply that $Q$ given by (7.4) is nonnegative definite and satisfies (5.4). With the stabilizability assumption, the result follows from Proposition 6.1.
Remark 8.1: Suppose $f=n, \dot{C}=I_{n}$ (so that perfect measurements of the entire state are available), and $Q$ satisfying (7.9) is positive definite. Then it follows from Theorem 7.1 that $f=I_{n}, f_{A}=0, C_{e}=0$ (i.e., the dynamic filter is disabled), $D_{e}=L$, and by ( 8.1 ),$J=0$. This is, of course, the expected result since perfect estimation is achievable in this case.

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$$
\begin{align*}
& 0=A_{0} Q+Q A_{i}^{T}+V_{1}+\sum_{i=1}^{\rho} \gamma_{i} A_{i}(Q+Q) A_{i} \\
& -Q V_{21}^{-1} Q_{1}^{T}+T_{1} Q, V_{2}^{-1} Q_{i}^{T} T_{1}^{T},  \tag{7.9}\\
& 0=A_{a} Q+Q_{a}^{T}+Q_{s} V_{2}^{-1} Q_{s}^{T}-T_{1} Q_{s} V_{2}^{-1} Q_{s}^{T} T_{d}^{T},  \tag{7.10}\\
& 0=\left(A_{a}-Q_{s} V_{2 s}^{-1} C\right)^{T} P+P\left(A_{\alpha}-Q_{s} V_{2 s}^{-1} C\right) \tag{7.11}
\end{align*}
$$

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## Approximate and Limit Results for Nonlinear Filters with Small Observation Noise: The Linear Sensor and Constant Diffusion Coefficient Case

## OFFR ZEITOUNI


#### Abstract

Recursive approximations for a class of filtering problems are presented. This class is characterized by linear observation sensor, constant diffusion terms, and for the multidimensional problem, poten-tial-like conditions on the drift. For the case of small observation noise, these approximations are used to demonstrate the Gaussian limiting structure of the optimal noninear filter.


## I. Introduction

The classical nonlinear filtering problem is of the form

$$
\begin{gather*}
d x_{1}=f\left(x_{1}\right) d t+o\left(x_{1}\right) d w_{1}, x_{1} \in \mathbb{N}^{*}, \quad p\left(x_{0}\right)=p_{0}\left(x_{0}\right)  \tag{1.1}\\
d y_{1}=g\left(x_{1}\right) d t+N_{0}^{1 / 2} d \theta_{1}, y_{1} \in \mathbb{R}^{m}, \quad y_{0}=0 \tag{1.2}
\end{gather*}
$$

Manuscripe received February 5. 1987: revised Seprember 2. 1987. Paper recornmended by Associate Editor, T. J. Tam. This work was mpported in part by a Weizmann Postdoctoral Fellowship.

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IEEE Log Number 8718566.
where $\boldsymbol{w}_{n}, \boldsymbol{\theta}$, are independent Brownian motions and the filtering problem consists of compuing statistics of $x_{1}$ when the observation $o$ algebra $y_{0}^{\prime} \Delta$ $\left\{y_{s}, 0 \leqslant s \leqslant t\right\}$ is given. By now, it is clear that for all but a few problems, an explicit final dimensional solution does not exist [4]. Therefore, one is led to consider approximations and to consider simplified limiting cases. Especially, the low observation noise case ( $N_{0}$ $\rightarrow 0$ ) has been considered in the literature [3], [6], [11].

In this paper, we restrict out attention to the special case of linear observations and constant diffusion coefficients, i.e.,

$$
\begin{gather*}
d x_{1}=f\left(x_{1}\right) d t+\sigma d w_{1}, x_{1} \in \mathbb{R}^{n}, \quad p\left(x_{0}\right)=p_{0}\left(x_{0}\right)  \tag{1.1'}\\
d y_{i}=g x_{1} d t+N_{0}^{1 / 2} d \theta_{1}, y_{1} \in \mathbb{R}^{m}, \quad y_{0}=0
\end{gather*}
$$

where $\sigma, g$ are matrices of appropriate dimensions. In the multidimensional case, we impose some additional potential-like stuctural conditions on $f(\cdot)$. For this restricted class of filtering problems, we derive recursive approximations to the conditional density $p_{x_{f}}\left(z \mid y_{q}\right)$ and to its unnormalized version $\rho_{x}\left(z \mid y g\right.$. For the limiting case $N_{0} \rightarrow 0$, those approximations are used to show that the conditional density, rescaled in a suitable manner. converges to a Gaussian density, with tight estimates on the "tails" of the density. This fact demonstrates the usefulness of Gaussian-based approximations (like the extended Kalman filter or the second-order Gaussian filter).

Related results were obtained by Mayer-Wolf [9] in his dissertation. There, bounds on the filtering error and the Cramer-Rao inequality are used to prove a basic Gaussian limit result, although under different assumptions.

The paper is organized as follows. The one-dimensional problem ( $n=$ $m=1$ ) is treated in Sections II and III. In Section II, we present our basic approximation theorem, which holds whether $N_{0}$ is small or not. We further demonstrate that, if $N_{0} \rightarrow 0$ the approximations exhibit certain nice limiting behavior, then the rescaled conditional density converges (weakly and pointwise) to a Gaussian density. In Section III, we check out explicitly the limiting behavior of the approximations and derive explicit conditions on $f(\cdot)$ under which the density indeed converges to a Gaussian one. Finally, in Section $I V$, we extend our results to a class of multidimensional problems.
We make throughout, the following assumption.
(A1) $f(\cdot)$ is continuously differentiable with bounded first partial derivatives.

## II. an approximation Theorem-The One-Dimensional Case

In this section, an approximation theorem for the unnormalized conditional density $\rho(z \mid y \partial$ is presented. Throughout, the one-dimensional case is treated. Multidimensional extensions are postponed to Section IV.

Without loss of generality, we assume $\sigma=1$ in (1.1'). Recall that under (A-1), a solution to (1.1') exists and is unique. Moreover, the measure $P_{1}$ defined by the pair (1.1'), (1.2') is absolutely continuous w.r.t. the reference measure $P_{0}$ defined by

$$
\begin{array}{cc}
d x_{1}=\alpha x_{1} d t+d w_{1}, & p\left(x_{0}\right)=p_{0}\left(x_{0}\right) \\
d y_{1}=N_{0}^{1 / 2} d \theta_{1}, & y_{0}=0 \tag{2.2}
\end{array}
$$

where $\alpha$ is some constant to be defined. The Radon-Nikodym derivalive $d P_{1} / d P_{0}$ is [8].

$$
\begin{align*}
& \frac{d P_{1}}{d P_{0}}=\exp \left(\int_{0}^{1}\left(f\left(x_{s}\right)-\alpha x_{s}\right) d x_{3}-1 / 2 \int_{0}^{1}\left(f^{2}\left(x_{s}\right)-\alpha^{2} x_{1}^{2}\right) d s\right. \\
&\left.+\int_{0}^{1} N_{0}^{-1} g x_{1} d y_{s}-1 / 2 \int_{0}^{1} N_{0}^{-1} g^{2} x_{1}^{2} d s\right) \tag{2.3}
\end{align*}
$$

As is well known [2], [12], [13], the unnormalized conditional density $\rho(z \mid y)$ satisfies

$$
\begin{equation*}
\rho\left(z \mid y_{0}^{\prime}\right) \propto E_{0}\left(\left.\frac{d P_{1}}{d P_{0}} \right\rvert\, y_{0,}^{\prime} x_{1}=z\right) p_{x_{1}}(z) \tag{2.4}
\end{equation*}
$$

# Robust Stability and Performance via Fixed-Order Dynamic Compensation <br> with <br> Guaranteed Cost Bounds 

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Key words: robust, guaranteed bounds, LQG control, reduced-order, optimal, structured uncertainty, real parameters


#### Abstract

A feedback control-design problem involving structured real-valued plant parameter uncertainties is considered. Two robust control-design issues are addressed. The Robust Stability Problem involves deterministic bounded structured parameter variations, while the Robust Performance Problem includes, in addition, a quadratic performance criterion averaged over stochastic disturbances and maximized over the admissible parameter variations. The optimal projection approach to fixed-order dynamic compensation is merged with the guaranteed cost control approach to robust stability and performance to obtain a theory of full- and reduced-order robust control design. The principal result is a sufficient condition for characterizing dynamic controllers of fixed dimension which are guaranteed to provide both robust stability and performance. The sufficient conditions involve a system of modified Riccati and Lyapunov equations coupled by an oblique projection, i.e., idempotent matrix, as well as the uncertainty bounds. Finally, in contrast to the usual separated Riccati equations, the full-order result involves a coupled system consisting of two modified Riccati equations and two modified Lyapunov equations coupled by the uncertainty bounds. The coupling illustrates the breakdown of the separation principle for LQG control with real-valued structured plant parameter variations.


Supported in part by the Air Force Office of Scientific Research under contract F49620-86-C0002.

## 1. Introduction

The direct method of Lyapunov has proven to be an effective approach to robust analysis and design of feedback control laws. References [B1], [B2], [BCL], [BG2], [CL], [CP], [ER], [GB], $[\mathrm{H}],[\mathrm{KB}],[\mathrm{KBH}],[\mathrm{L}],[\mathrm{PH}],[\mathrm{TB}],[\mathrm{VW}]$ comprise only a representative collection of the increasing literature in this area. In performing robust synthesis there are two principal issues, namely, stability robustness and performance robustness. Stability robustness addresses the problem of guaranteeing stability of the closed-loop system for plant perturbations within a specified class of uncertainties. In addition to guaranteeing robust stability, it is often desirable to minimize the worst-case performance degradation within a given robust stability range. Although both robust stability and performance are of interest in practice, most of the literature involving quadratic Lyapunov functions is confined to the problem of robust stability. A notable exception is the early work of Chang and Peng ([CP]) which also provides bounds on worst-case quadratic performance within full-state feedback control design.

The contribution of the present paper is a methodology for designing controllers which provide both robust stability and robust performance over a prescribed range of real-valued structured plant parameter variations. The feedback law is in the form of a fixed-order (i.e., full- or reducedorder) strictly proper dynamic compensator. The overall approach is based upon the merging of two distinct control-design techniques, namely, the guaranteed cost control approach to robust performance ([CP]) and the optimal projection approach to fixed-order dynamic compensation ([BH1],[HB]). The principal motivation for our approach is to permit greater flexibility in the design of robust feedback laws by providing an alternative to full-state feedback and full-order dynamic compensation.

The guaranteed cost control approach ([CP]) adopted in the present paper utilizes a performance bound to provide robust performance in addition to robust stability. Here, robust performance refers to a guaranteed bound on the worst-case value of the expectation of a quadratic cost criterion over a prescribed uncertainty set. This quadratic criterion is precisely the standard cost functional of linear-quadratic-Gaussian control theory. By bounding the worst-case value of this criterion over a specified range of plant uncertainties, we effectively bound the variances of specified states and control signals.

To bound the worst-case closed-loop performance, we require a bound on the effect of plant uncertainties on the steady-state closed-loop covariance matrix. The form of the guaranteed cost
control bound utilized herein was originally motivated by the effect of multiplicative white noise on the state covariance ([B2]),[BG2]). Since this bound is differentiable with respect to the covariance matrix and compensator gains, it permits optimal design via first-order necessary conditions. This approach is not possible using the nondifferentiable bound originally proposed in [CP]. An alternative differentiable bound proposed in $[\mathrm{PH}]$ for full-state feedback has been extended to fixed-order dynamic compensation in [BH2].

In the present paper, the guaranteed cost technique is used to bound the closed-loop performance and characterize robustly stabilizing controllers. This performance bound is then interpreted as an auxiliary cost which is to be minimized by the choice of compensator gains. The actual performarse for a given realization of the parameter uncertainty is thus guaranteed to lie below this bound. In the presence of a stabilizability (disturbability) assumption, the robust performance bound automatically implies robust stability. The auxiliary cost and the Lyapunov equation constraint together form the Auxiliary Minimization Problem. Since the Auxiliary Minimization Problem is a nonconvex mathematical programming problem with differentiable data, it is amenable to first-order necessary conditions.

One feature of this approach is that since the necessary conditions are obtained for the Auxiliary Minimization Problem rather than the original problem, extremals are guaranteed to provide both robust stability and performance. Note that this is true for every extremal of the Auxiliary Minimization Problem whether it corresponds to a local minimum, local maximum, or otherwise. Of course, the global minimum is most likely to provide the best worst-case performance over the robust stability range. In any case, necessary conditions for the Auxiliary Minimization Problem effectively serve as sufficient conditions for robust stability with a guaranteed performance bound.

The present paper encompasses a rigorous development of sufficient conditions for robust stability and performance via fixed-order dynamic compensation. These sufficient conditions are in the form of a coupled system of algebraic matrix equations consisting of two modified Riccati equations and two modified Lyapunov equations. The coupling is due to the optimal projection, which characterizes reduced-order controllers, and the uncertainty bounds, which account for the effect of parameter uncertainties on the performance functional. When the compensator order is constrained to be equal to the dimension of the plant and the uncertainty bounds are absent, the equations specialize to the usual pair of separated Riccati equations of steady-state LQG theory.

We are quick to point out, however, that our approach is constructive in nature rather than
existential. That is, our sufficient conditions provide explicit formulae for robust, fixed-order feedback gains when the Auxiliary Minimization Problem has a solution. In this sense our constructive conditions can be viewed as complementary to existential results on robust stabilizability. Specifically, the existence of a solution to the Auxiliary Minimization Problem and associated design equations depends upon stabilizability via fixed-order controllers as well as the sharpness of the quadratic Lyapunov bounds. The stabilizability problem has been studied using independent methods (see, e.g., [BHK]), while the conservatism of the bounds is considered in [BH3]. In addition, we state a local existence result for solvability of the design equations which assumes only nominal stabilizability.

The contents and scope of the paper are as follows. In Section 2 we state the robust stability and performance problems for fixed-order dynamic compensation with plant parameter uncertainty. In Section 3 a modified Lyapunov equation is introduced whose solution, when it exists, is guaranteed to bound the steady-state closed-loop covariance over the specified range of plant uncertainty. A performance bound is then given in terms of the covariance bound. In Section 4 we view the performance bound as an auxiliary cost and consider the problem of minimizing the auxiliary cost subject to the modified Lyapunov equation and a definiteness condition as side constraints. These side constraints have the property that all admissible elements provide robust stability and performance (Proposition 4.1). In Section 5 the uncertainty set and bound for constructing the modified Lyapunov equation are given concrete forms. Specifically, the uncertainty set has the form of an ellipsoidal region in parameter space while the modified Lyapunov equation includes additional linear terms to bound the uncertainty. A sufficient condition involving Kronecker sums and products implies the existence of a unique, nonnegative-definite solution to the modified Lyapunov equation. Section 6 presents the first-order necessary conditions (Theorem 6.1) for the Auxiliary Minimization Problem under minor additional technical conditions to ensure the applicability of the Lagrange multiplier technique. As discussed above, these necessary conditions are in the form of extended optimal projection equations. A partial converse of the necessary conditions shows that solutions of these algebraic equations provide, by construction, a solution of the original modified Lyapunov equation. This result is combined in Section 7 (Theorem 7.1) with a stabilizability assumption to guarantee robust stability with a robust performance bound. In addition, we state an existence result for local solvability of the design equations by applying a result from [R1], [R2] (Theorem 7.2). To draw connections with standard LQG theory, in Section 8 we specialize Theorem 7.1 to the full-order case. In contrast to the pair of separated Riccati equations of standard LQG theory,
the full-order result in the presence of plant parameter variations'is given by a coupled system of four modified Riccati/Lyapunov equations. In Section 9 the theory is illustrated by means of an example due to Doyle ([D]). This problem was also considered in [BG1] before the robustness theory developed herein was available. Hence the present paper can be viewed as the rigorous mathematical foundation which legitimizes the heretofore ad hoc robustness approach of [BG1].

Notation. Note: All matrices have real entries

| $\boldsymbol{R}, \mathbb{R}^{r \times 1}, \mathbb{R}^{r}, \mathbb{E}$ | real numbers, $r \times s$ real matrices, $\mathbb{R}^{r \times 1}$, expected value |
| :---: | :---: |
| $I_{r}()^{T}, 0_{r \times a}, 0_{r}$ | $r \times r$ identity matrix, transpose, $r \times s$ zero matrix, $0_{r \times r}$ |
| $Z_{(i, j)}$, tr $Z$ | ( $i, j$ )-element of matrix $Z$, trace of square matrix $Z$ |
| $\oplus$ ¢ $\otimes$ | Kronecker sum, Kronecker product ([B3]) |
| $\mathbb{S}^{\mathbf{r}}, \mathbb{N}^{\mathbf{r}}, \mathbb{P}^{\mathbf{r}}$ | $r \times r$ symmetric, nonnegative-definite, positive-definite matrices |
| $Z_{1} \leq Z_{2}, Z_{1}<Z_{2}$ | $Z_{2}-Z_{1} \in \mathbb{N}^{r}, \quad Z_{2}-Z_{1} \in \mathbb{P}^{r}, \quad Z_{1}, Z_{2} \in \mathbb{S}^{r}$ |
| $n, m, \ell, n_{c}, p ; \tilde{n}$ | positive integers; $n+n_{c}$ |
| $x, u, y, x_{c}, \tilde{x}$ | $n, m, \ell, n_{c}, \tilde{n}$-dimensional vectors |
| $A, \Delta A ; B, \Delta B$ | $n \times n$ matrices; $n \times m$ matrices |
| $C, \Delta C ; D, \Delta D$ | $\ell \times n$ matrices; $\ell \times m$ matrices |
| $A_{e}, B_{c}, C_{c}$ | $n_{c} \times n_{c}, n_{c} \times \ell, m \times n_{c}$ matrices |
| $\boldsymbol{\alpha}$ | positive number |
| $A_{\alpha}, A_{c \alpha}$ | $A+\frac{\infty}{2} I_{n}, \quad A_{c}+\frac{\infty}{2} I_{n}$ |
| $\alpha_{i}$ | positive number, $i=1, \ldots, p$ |
| $\boldsymbol{\gamma}$ | $\alpha_{i}^{2} / \alpha, i=1, \ldots, p$ |
| $\sigma_{i}$ | real number, $i=1, \ldots, p$ |
| $R_{1}, R_{2}$ | state, control weighting matrices; $R_{1} \in \mathbb{N}^{\mathbf{n}}, R_{2} \in \mathbb{P}^{\mathbf{m}}$ |
| $R_{12}$ | $n \times m$ cross-weighting matrix; $\quad R_{1}-R_{12} R_{2}^{-1} R_{12}^{\mathrm{T}} \geq 0$ |
| $\tilde{\boldsymbol{R}}$ | $\left[\begin{array}{cc} R_{1} & R_{12} C_{c} \\ C_{c}^{\mathrm{T}} R_{12}^{\mathrm{T}} & C_{c}^{\mathrm{T}} R_{2} C_{c} \end{array}\right]$ |
| $w_{1}(\cdot), w_{2}(\cdot)$ | $n, \ell$-dimensional white noise |
| $V_{1}, V_{2}$ | intensity of $w_{1}(\cdot), w_{2}(\cdot) ; \quad V_{1} \in \mathbb{N}^{n}, V_{2} \in \mathbb{P}^{\ell}$ |
| $V_{12}$ | $n \times \ell$ cross intensity of $w_{1}(\cdot), w_{2}(\cdot)$ |
| $\tilde{w}(\cdot), \tilde{V}$ | $\left[\begin{array}{c} w_{1}(\cdot) \\ B_{c} w_{2}(\cdot) \end{array}\right], \quad\left[\begin{array}{cc} V_{1} & V_{12} B_{c}^{\mathrm{T}} \\ B_{c} V_{12}^{\mathrm{T}} & B_{c} V_{2} B_{c}^{\mathrm{T}} \end{array}\right]$ |

## 2. Robust Stability and Robust Performance Problems

In this section we state the Robust Stability Problem and Robust Performance Problem. Both problems involve a set $U \subset \mathbb{R}^{n \times n} \times \mathbb{R}^{n \times m} \times \mathbb{R}^{\ell \times n} \times \mathbb{R}^{\ell \times m}$ of uncertain perturbations ( $\Delta A, \Delta B, \Delta C, \Delta D$ ) of the nominal system matrices ( $A, B, C, D$ ). The goal of the Robust Stability Problem is to determine a fixed-order, strictly proper dynamic compensator ( $A_{c}, B_{c}, C_{c}$ ) which stabilizes the plant for all variations in $U$. In this section and the following section no explicit assumptions are required for the set $U$. In Section 5 the structure of variations in $U$ will be specified.

Robust Stability Problem. For fixed $n_{c} \leq n$ determine ( $A_{c}, B_{c}, C_{c}$ ) such that the closedloop system consisting of the $n$ th-order controlled plant

$$
\begin{equation*}
\dot{x}(t)=(A+\Delta A) x(t)+(B+\Delta B) u(t), \quad t \in[0, \infty) \tag{2.1}
\end{equation*}
$$

measurements

$$
\begin{equation*}
y(t)=(C+\Delta C) x(t)+(D+\Delta D) u(t) \tag{2.2}
\end{equation*}
$$

and $n_{\mathrm{c}}$ th-order dynamic compensator

$$
\begin{align*}
\dot{x}_{c}(t) & =A_{c} x_{c}(t)+B_{c} y(t)  \tag{2.3}\\
u(t) & =C_{c} x_{c}(t) \tag{2.4}
\end{align*}
$$

is asymptotically stable for all $(\Delta A, \Delta B, \Delta C, \Delta D) \in U$.
The Robust Performance Problem involves, in addition, white plant disturbances and measurement noise. The goal of this problem is to determine a fixed-order, strictly proper compensator ( $A_{c}, B_{c}, C_{c}$ ) which minimizes the worst-case value over the uncertainty set $U$ of a steady-state average quadratic performance criterion.

Robust Performance Problem. For fixed $n_{c} \leq n$, determine ( $A_{c}, B_{c}, C_{c}$ ) such that, for the closed-loop system consisting of the $n$ th-order controlled and disturbed plant

$$
\begin{equation*}
\dot{x}(t)=(A+\Delta A) x(t)+(B+\Delta B) u(t)+w_{1}(t), \quad t \in[0, \infty) \tag{2.5}
\end{equation*}
$$

noisy measurements

$$
\begin{equation*}
y(t)=(C+\Delta C) x(t)+(D+\Delta D) u(t)+w_{2}(t) \tag{2.6}
\end{equation*}
$$

and $n_{c}$ th-order dynamic compensator (2.3), (2.4), the performance criterion

$$
\begin{equation*}
J\left(A_{c}, B_{c}, C_{c}\right) \triangleq \sup _{(\Delta A, \Delta B, \Delta C, \Delta D) \in u} \limsup _{t \rightarrow \infty} \mathbb{E}\left[x^{\mathrm{T}}(t) R_{1} x(t)+2 x^{\mathrm{T}}(t) R_{12} u(t)+u^{\mathrm{T}}(t) R_{2} u(t)\right] \tag{2.7}
\end{equation*}
$$

is minimized.

Remark 2.1. The cost functional (2.7) is identical to the standard LQG criterion with the exception of the supremum for evaluating worst-case quadratic performance over $U$. Note that (2.7) can also be written in terms of an averaged integral; i.e.,

$$
\begin{aligned}
& J\left(A_{e}, B_{c}, C_{c}\right)= \\
& \quad \sup _{(\Delta A, \Delta B, \Delta C, \Delta D) \in u} \limsup _{t \rightarrow \infty} \frac{1}{t} \mathbb{E}\left\{\int_{0}^{t}\left[x^{\mathrm{T}}(s) R_{1} x(s)+2 x^{\mathrm{T}}(s) R_{12} u(s)+u^{\mathrm{T}}(s) R_{2} u(s)\right] d s\right\} .
\end{aligned}
$$

For practical application, the cost (2.7) provides the means for minimizing the variances of selected state variables and control signals. This can be achieved by appropriate selection of the matrices $R_{1}$ and $R_{2}$ which serve as design weights. For robust performance the goal is to minimize the worst-case variances of selected variables over the plant uncertainty.

For each uncertain variation $(\Delta A, \Delta B, \Delta C, \Delta D) \in U$, the undisturbed closed-loop system (2.1)-(2.4) can be written as

$$
\begin{equation*}
\dot{\tilde{x}}(t)=(\tilde{A}+\Delta \tilde{A}) \tilde{x}(t), \quad t \in[0, \infty), \tag{2.8}
\end{equation*}
$$

where

$$
\tilde{x}(t) \triangleq\left[\begin{array}{c}
x(t) \\
x_{c}(t)
\end{array}\right] \quad \tilde{A} \triangleq\left[\begin{array}{cc}
A & B C_{c} \\
B_{c} C & A_{c}+B_{c} D C_{c}
\end{array}\right], \quad \Delta \tilde{A} \triangleq\left[\begin{array}{cc}
\Delta A & \Delta B C_{c} \\
B_{c} \Delta C & B_{c} \Delta D C_{c}
\end{array}\right]
$$

Similarly, the disturbed closed-loop system (2.3)-(2.6) can be written as

$$
\begin{equation*}
\dot{\tilde{x}}(t)=(\tilde{A}+\Delta \tilde{A}) \tilde{x}(t)+\tilde{w}(t), \quad t \in[0, \infty) \tag{2.9}
\end{equation*}
$$

where the closed-loop disturbance $\tilde{w}(t)$ has intensity $\tilde{V} \in \mathbb{N}^{\tilde{n}}$.

## 3. Sufficient Conditions for Robust Stability and Performance

In practice, steady-state performance is only of interest when the undisturbed closed-loop system (2.8) is robustly stable over $U$. The following result, which expresses the performance in terms of the steady-state closed-loop second-moment matrix; is immediate.

Lemma 3.1. Let ( $A_{c}, B_{c}, C_{c}$ ) be given and suppose the system (2.8) is stable for all ( $\Delta A, \Delta B$, $\Delta C, \Delta D) \in U$. Then

$$
\begin{equation*}
J\left(A_{c}, B_{c}, C_{c}\right)=\sup _{(\Delta A, \Delta B, \Delta C, \Delta D) \in U} \operatorname{tr} \tilde{Q}_{\Delta \tilde{A}} \tilde{R} \tag{3.1}
\end{equation*}
$$

where $\tilde{\boldsymbol{Q}}_{\Delta \tilde{A}} \triangleq \lim _{t \rightarrow \infty} \mathbb{E}\left[\tilde{x}(t) \tilde{x}^{\mathrm{T}}(t)\right] \in \mathbb{N}^{\tilde{n}}$ is the unique solution to

$$
\begin{equation*}
0=(\tilde{A}+\Delta \tilde{A}) \tilde{Q}_{\Delta \tilde{A}}+\tilde{Q}_{\Delta \tilde{A}}(\tilde{A}+\Delta \tilde{A})^{\mathbf{T}}+\tilde{V} \tag{3.2}
\end{equation*}
$$

Remark 3.1. When $U$ is compact, "sup" in (3.1) can be replaced by "max".
The key step in guaranteeing robust stability and performance is to replace the uncertain terms in the covariance Lyapunov equation (3.2) by a bounding function $\Omega$. Note that since $\Delta \tilde{A}$ is independent of $A_{c}$, the bounding function need only depend upon $B_{c}$ and $C_{c}$.

Theorem 3.1. Let $\Omega: \mathbb{N}^{n} \times \mathbb{R}^{n_{0} \times \ell} \times \mathbb{R}^{m \times n_{0}} \rightarrow \mathbb{S}^{n}$ be such that

$$
\Delta \tilde{A} Q+Q \Delta \tilde{A}^{\mathrm{T}} \leq \Omega\left(\mathcal{Q}, B_{c}, C_{c}\right)
$$

$$
\begin{equation*}
(\Delta A, \Delta B, \Delta C, \Delta D) \in U, \quad\left(\Omega, B_{c}, C_{c}\right) \in \mathbb{N}^{\tilde{n}} \times \mathbb{R}^{n_{c} \times \ell} \times \mathbb{R}^{m \times n_{e}} \tag{3.3}
\end{equation*}
$$

and, for given $\left(A_{c}, B_{c}, C_{c}\right)$, suppose there exists $Q \in \mathbb{N}^{\bar{n}}$ satisfying

$$
\begin{equation*}
0=\tilde{A} \underline{Q}+Q \tilde{A}^{\mathrm{T}}+\Omega\left(Q, B_{c}, C_{c}\right)+\tilde{V} . \tag{3.4}
\end{equation*}
$$

Then

$$
\begin{gather*}
\left(\tilde{A}+\Delta \tilde{A},\left[\tilde{V}+\Omega\left(Q, B_{c}, C_{c}\right)-\left(\Delta \tilde{A} Q+Q \Delta \tilde{A}^{\mathrm{T}}\right)\right]^{\frac{1}{2}}\right) \text { is stabilizable, }  \tag{3.5}\\
(\Delta A, \Delta B, \Delta C, \Delta D) \in U
\end{gather*}
$$

if and only if

$$
\begin{equation*}
\tilde{A}+\Delta \tilde{A} \text { is asymptotically stable, } \quad(\Delta A, \Delta B, \Delta C, \Delta D) \in U \tag{3.6}
\end{equation*}
$$

In this case,

$$
\begin{equation*}
\tilde{Q}_{\Delta \bar{A}} \leq \mathcal{Q}, \quad(\Delta A, \Delta B, \Delta C, \Delta D) \in U \tag{3.7}
\end{equation*}
$$

where $\tilde{Q}_{\Delta \tilde{A}}$ is given by (3.2), and

$$
\begin{equation*}
J\left(A_{c}, B_{c}, C_{c}\right) \leq \operatorname{tr} Q \tilde{R} \tag{3.8}
\end{equation*}
$$

Proof. First note for clarity that in (3.3) $Q$ denotes an arbitrary element of $\mathbb{N}^{n}$ since (3.3) holds for all $Q \in \mathbb{N}^{\tilde{n}}$, while in (3.4) $\mathcal{Q}$ denotes a specific solution to (3.4). Now for $(\Delta A, \Delta B, \Delta C, \Delta D) \in$ $U$, (3.4) is equivalent to

$$
\begin{equation*}
0=(\tilde{A}+\Delta \tilde{A}) Q+Q(\tilde{A}+\Delta \tilde{A})^{\mathbf{T}}+\Omega\left(Q, B_{c}, C_{c}\right)-\left(\Delta \tilde{A} Q+Q \Delta \tilde{A}^{\mathrm{T}}\right)+\tilde{V} \tag{3.9}
\end{equation*}
$$

Hence, by assumption, (3.9) has a solution $Q \in \mathbb{N}^{n}$ for all $(\Delta A, \Delta B, \Delta C, \Delta D) \in U$ and, by (3.3), $\Omega\left(\mathcal{Q}, B_{c}, C_{c}\right)-\left(\Delta \tilde{A} Q+\mathcal{Q} \Delta \tilde{A}^{T}\right)$ is nonnegative definite. Now if the stabilizability condition
(3.5) holds for all $(\Delta A, \Delta B, \Delta C, \Delta D) \in U$, it follows from Lemma 12.2 of [W] that $\tilde{A}+\Delta \tilde{A}$ is asymptotically stable for all $(\Delta A, \Delta B, \Delta C, \Delta D) \in U$. Conversely, if $\tilde{A}+\Delta \tilde{A}$ is asymptotically stable for all ( $\Delta A, \Delta B, \Delta C, \Delta D) \in U$, then (3.5) holds. Next, subtracting (3.2) from (3.9) yields

$$
0=(\tilde{A}+\Delta \tilde{A})\left(Q-\tilde{Q}_{\Delta \tilde{A}}\right)+\left(\mathcal{Q}-\tilde{Q}_{\Delta \tilde{A}}\right)(\tilde{A}+\Delta \tilde{A})^{\mathrm{T}}+\Omega\left(\mathcal{Q}, B_{c}, C_{c}\right)-\left(\Delta \tilde{A} Q+Q \Delta \tilde{A}^{\mathrm{T}}\right),
$$

or, equivalently, since $\tilde{A}+\Delta \tilde{A}$ is asymptotically stable for all $(\Delta A, \Delta B, \Delta C, \Delta D) \in U$,

$$
Q-\tilde{Q}_{\Delta \bar{A}}=\int_{0}^{\infty} e^{(\bar{A}+\Delta \tilde{A}) t}\left[\Omega\left(\mathcal{Q}, B_{c}, C_{c}\right)-\left(\Delta \tilde{A} Q+Q \Delta \tilde{A}^{T}\right)\right] e^{(\bar{A}+\Delta \tilde{A})^{T} t} d t \geq 0
$$

which implies (3.7). The performance bound (3.8) is now an immediate consequence of (3.7).
Remark 3.2. In applying Theorem 3.1 it may be convenient to replace condition (3.5) with a stronger condition which is easier to verify in practice. Clearly, (3.5) is satisfied if $\tilde{V}+\Omega\left(\mathcal{Q}, B_{c}, C_{c}\right)-$ $\left(\Delta \tilde{A} Q+Q \Delta \tilde{A}^{\mathrm{T}}\right)$ is positive definite for all $(\Delta A, \Delta B, \Delta C, \Delta D) \in U$. This will be the case, for example, if either $\tilde{V}$ is positive definite or strict inequality holds in (3.3). Also, it follows from Theorem 3.6 of [ $W$ ] that (3.5) is implied by the stronger condition

$$
\begin{equation*}
\left(\tilde{A}+\Delta \tilde{A}, \tilde{V}^{\frac{1}{2}}\right) \text { is stabilizable, } \quad(\Delta A, \Delta B, \Delta C, \Delta D) \in U \tag{3.10}
\end{equation*}
$$

Remark 3.3. The covariance bound (3.7) can also be used to analyze the effect of disturbances on specified state variables. For example, if $\boldsymbol{E}_{1} \in \mathbb{R}^{9 \times n}$ then (3.7) implies

$$
\left[\begin{array}{ll}
E_{1} & 0_{q \times n_{0}}
\end{array}\right] \tilde{Q}_{\Delta \tilde{A}}\left[\begin{array}{c}
E_{1}^{\mathrm{T}}  \tag{3.11}\\
0_{n_{*} \times q}
\end{array}\right] \leq\left[\begin{array}{ll}
E_{1} & 0_{q \times n_{d}}
\end{array}\right] \odot\left[\begin{array}{c}
E_{1}^{\mathrm{T}} \\
0_{n_{0} \times q}
\end{array}\right]
$$

so that the right hand side of (3.11) serves as a bound on selected state variances. For controldesign purposes we effectively set $R_{1}=E_{1}^{T} E_{1}$. Similar remarks apply to obtaining bounds on the variances of control signals.

## 4. The Auxiliary Minimization Problem

The key step in our development involves consideration of the performance bound (3.8) in place of the actual worst-case performance $J\left(A_{c}, B_{c}, C_{c}\right)$. This leads to the following problem.

Auxiliary Minimization Problem. Determine ( $\mathcal{Q}, A_{c}, B_{c}, C_{c}$ ) which minimizes

$$
\begin{equation*}
J\left(Q, A_{c}, B_{c}, C_{c}\right) \triangleq \operatorname{tr} Q \tilde{R} \tag{4.1}
\end{equation*}
$$

subject to (3.4) and

$$
\begin{equation*}
Q \in \mathbb{N}^{n} \tag{4.2}
\end{equation*}
$$

The relationship between the Auxiliary Minimization Problem and the Robust Stabiiity and Performance Problems is straightforward as shown by the following observation.

Proposition 4.1. If ( $\mathcal{Q}, A_{c}, B_{c}, C_{c}$ ) satisfies (3.4), (4.2) and the stabilizability condition (3.5) holds, then $\tilde{A}+\Delta \tilde{A}$ is asymptotically stable for all $(\Delta A, \Delta B, \Delta C, \Delta D) \in U$, and

$$
\begin{equation*}
J\left(A_{c}, B_{c}, C_{c}\right) \leq J\left(\mathcal{Q}, A_{c}, B_{c}, C_{c}\right) \tag{4.3}
\end{equation*}
$$

Proof. Since (3.4) has a solution $\mathcal{Q} \in \mathbb{N}^{\bar{n}}$ and the stabilizability condition (3.5) holds, the hypotheses of Theorem 3.1 are satisfied so that robust stability with robust performance bound (3.8) is guaranteed. Note that (4.3) is merely a restatement of (3.8).

Several comments are in order. Since the auxiliary cost (4.1) is an upper bound for the actual cost (2.7), it is clearly desirable to minimize (4.1) over $Q$ and the controller gains. Note, howevc=, that the Auxiliary Minimization Problem is a nonconvex mathematical programming problem on a noncompact set. Hence guarantees of existence of solutions and sufficient conditions for global optimality cannot be obtained without imposing additional confining assumptions. To develop nonrestrictive res lts, we proceed in Sectiva 6 by deriving necessary conditions for optimality which require no further assumptions except that $\Omega$ be differentiable and that the minimization be performed over an open set. In the next section we cusstruct a bound $\Omega$ which possesses the required smoothness.

## 5. Uncertainty Structure and the Guaranteed Cost Bound

Having established the theoretical basis for our approach, we now assign explicit structure to the set $U$ and bounding function $\Omega$. Specifically, the uncertainty set $U$ is assumed to be of the form

$$
\begin{gather*}
u=\left\{(\Delta A, \Delta B, \Delta C, \Delta D): \quad \Delta A=\sum_{i=1}^{p} \sigma_{i} A_{i}, \quad \Delta B=\sum_{i=1}^{p} \sigma_{i} B_{i}, \quad \Delta C=\sum_{i=1}^{p} \sigma_{i} C_{i}\right.  \tag{5.1}\\
\left.\Delta D=\sum_{i=1}^{p} \sigma_{i} D_{i}, \quad \sum_{i=1}^{p} \sigma_{i}^{2} / \alpha_{i}^{2} \leq 1\right\}
\end{gather*}
$$

where, for $i=1, \ldots, p: A_{i} \in \mathbb{R}^{n \times n}, B_{i} \in \mathbb{R}^{n \times m}, C_{i} \in \mathbb{R}^{\ell \times n}$, and $D_{i} \in \mathbb{R}^{\prime}$ ( $m$ are fixed matrices denoting the structure of the parametric uncertainty; $\alpha_{i}$ is a given positive number; and $\sigma_{i}$ is an uncertain real parameter. Note that the uncertain parameters $\sigma_{i}$ are assumed to lie in a specified ellipsoidal region in $\mathbb{R}^{\boldsymbol{p}}$. The closed-loop system (2.8) thus has structured uncertainty of the form

$$
\begin{equation*}
\Delta \tilde{A}=\sum_{i=1}^{p} \sigma_{i} \tilde{A}_{i}, \tag{5.2}
\end{equation*}
$$

where

$$
\tilde{A}_{i} \triangleq\left[\begin{array}{cc}
A_{i} & B_{i} C_{c} \\
B_{c} C_{i} & B_{c} D_{i} C_{c}
\end{array}\right], \quad i=1, \ldots, p
$$

The uncertainty set $U$ is general in the sense that no explicit assumptions such as the matching conditions used in [BCL] will be made with regard to the structure of $A_{i}, B_{i}, C_{i}$, and $D_{i}$. We do, however, require (as is evident from (5.1)) that the uncertain parameters $\sigma_{i}$ appear linearly in the off-nominal perturbations which is more confining than matching assumptions. Note that the symmetry of the uncertainty set entails no loss of generality by requiring only a redefinition of the nominal plant matrices.

In order to obtain explicit gain expressions for ( $A_{c}, B_{c}, C_{c}$ ) in Section 6, we shall require one additional technical assumption. Specifically, we assume that for each $i \in\{1, \ldots, p\}$, at most one of the matrices $B_{i}, C_{i}$, and $D_{i}$ is nonzero. This condition thus implies that a given uncertain parameter $\sigma_{i}$ may appear explicitly in both $\Delta A$ and $\Delta B$, or both $\Delta A$ and $\Delta C$, or both $\Delta A$ and $\Delta D$, or only $\Delta A$, but not (say) in both $\Delta B$ and $\Delta D$. Thus we can account partially (but not totally) for correlated parameter uncertainties in different plant matrices. If a given uncertain parameter does arise in both (say) $\Delta B$ and $\Delta D$, then it must be represented by two distinct uncertain parameters. If this assumption is not.imposed, then optimality conditions can still be derived, but at the expense of closed-form gain expressions.

For the structure of $U$ as specified by (5.1), the bound $\Omega$ satisfying (3.3) can now be given a concrete for n .

Proposition 5.1. Let $\alpha$ be an arbitrary positive scalar. Then the function

$$
\begin{equation*}
\Omega\left(Q, B_{c}, C_{c}\right)=\alpha Q+\alpha^{-1} \sum_{i=1}^{p} \alpha_{i}^{2} \tilde{A}_{i} Q \tilde{A}_{i}^{\mathrm{T}} \tag{5.3}
\end{equation*}
$$

satisfies (3.3) with $U$ given by (5.1).
Proof. Note that

$$
\begin{aligned}
0 & \leq \sum_{i=1}^{p}\left[\left(\alpha^{\frac{1}{2}} \sigma_{i} / \alpha_{i}\right) I_{\tilde{n}}-\left(\alpha_{i} / \alpha^{\frac{1}{2}}\right) \tilde{A}_{i}\right] Q\left[\left(\alpha^{\frac{1}{2}} \sigma_{i} / \alpha_{i}\right) I_{\tilde{n}}-\left(\alpha_{i} / \alpha^{\frac{1}{2}}\right) \tilde{A}_{i}\right]^{\mathrm{T}} \\
& =\alpha \sum_{i=1}^{p}\left(\sigma_{i}^{2} / \alpha_{i}^{2}\right) Q+\alpha^{-1} \sum_{i=1}^{p} \alpha_{i}^{2} \tilde{A}_{i} Q \tilde{A}_{i}^{\mathrm{T}}-\sum_{i=1}^{p} \sigma_{i}\left(\tilde{A}_{i} Q+Q \tilde{A}_{i}^{\mathrm{T}}\right),
\end{aligned}
$$

which, since $\sum_{i=1}^{p} \sigma_{i}^{2} / \alpha_{i}^{2} \leq 1$, implies (3.3).

Remark 5.1. Note that the bound $\Omega$ given by (5.3) consists of two distinct terms. The first term $\alpha \mathcal{Q}$ can be thought of as arising from an exponential time weighting of the cost, or, equivalently, from a uniform right shift of the open-lonp dynamics ([ AM ]). The second term $\alpha^{-1} \sum_{i=1}^{p} \alpha_{i}^{2} \tilde{A}_{i} Q \tilde{A}_{i}^{\mathrm{T}}$ arises naturally from a multiplicative white noise model ([BG1],[BG2],[B2]). Such interpretations have no bearing on the results obtained here since only the bound $\Omega$ defined by (5.3) is required. Note that the bound (5.3) is valid for all positive $\alpha$. A similar bound was also considered in [KB].

With $\Omega$ defined by (5.3), the modified Lyapunov equation (3.4) becomes

$$
\begin{equation*}
0=\tilde{A} Q+Q \tilde{A}^{\mathrm{T}}+\alpha Q+\alpha^{-1} \sum_{i=1}^{p} \alpha_{i}^{2} \tilde{A}_{i} Q \tilde{A}_{i}^{\mathrm{T}}+\tilde{V} \tag{5.4}
\end{equation*}
$$

or, equivalently,

$$
\begin{equation*}
0=\tilde{A}_{\alpha} Q+Q \tilde{A}_{\alpha}^{\mathrm{T}}+\sum_{i=1}^{p} \gamma_{i} \tilde{A}_{i} Q \tilde{A}_{i}^{\mathrm{T}}+\tilde{V}, \tag{5.5}
\end{equation*}
$$

where

$$
\tilde{A}_{\alpha} \triangleq \tilde{A}+\frac{\alpha}{2} I_{\tilde{n}}=\left[\begin{array}{cc}
A_{\alpha} & B C_{c}  \tag{5.6}\\
B_{c} C & A_{c \alpha}+B_{c} D C_{c}
\end{array}\right]
$$

and $\gamma_{i} \triangleq \alpha_{i}^{2} / \alpha$. Note that (5.5) is equivalent to

$$
\begin{equation*}
0=A \operatorname{vec} Q+\operatorname{vec} \tilde{V}, \tag{5.7}
\end{equation*}
$$

where "vec" is the column-stain'ing operation defined in [B3] and $\AA$ is defined by

$$
A \triangleq \tilde{A}_{\alpha} \oplus \tilde{A}_{\alpha}+\sum_{i=1}^{p} \gamma_{i} \tilde{A}_{i} \otimes \tilde{A}_{i} .
$$

Next, using the bound $\Omega$ given by (5.3) and $U$ given by (5.1) we present a result which guarantees the existence of a nonnegative-definite solution to (3.4) or, equivalently, (5.5) for a given controller ( $A_{c}, B_{c}, C_{c}$ ). For the converse we view $\tilde{V}$ as an arbitrary element of $\mathbb{N}^{n}$.

Proposition 5.2. Let $\left(A_{c}, B_{c}, C_{c}\right)$ be given and let $\alpha>0$. If $A$ is asymptotically stable, then there exists a unique $\bar{n} \times \tilde{n} Q$ satisfying (5.5) and, furthermore, $\mathcal{Q} \geq 0$. Conversely, if for all $\tilde{V} \subseteq \mathbb{N}^{\tilde{n}}$ there exists $\mathcal{Q} \geq 0$ satisfying (5.5), then $A$ is asymptotically stable.

Proof. Since (5.5) is equivalent to

$$
\begin{equation*}
Q=-\operatorname{vec}^{-1}\left[A^{-1} \operatorname{vec} \tilde{V}\right] \tag{5.8}
\end{equation*}
$$

existence and uniqueness hold. To prove that $Q$ is nonnegative definite, we rewrite (5.8) as

$$
\begin{equation*}
\mathcal{Q}=\int_{0}^{\infty} \operatorname{vec}^{-1}\left[e^{A t} \operatorname{vec} \tilde{V}\right] d t \tag{5.9}
\end{equation*}
$$

and show that the integrand is nonnegative definite for all $t \in[0, \infty)$. [Note that the following argument does not require that $A$ be stable.] Using the Lie exponential product formula, the exponential in (5.9) can be written as

$$
\begin{equation*}
e^{A t}=\lim _{k \rightarrow \infty}\left\{\exp \left[\frac{1}{k}\left(\tilde{A}_{\alpha} \oplus \tilde{A}_{\alpha}\right) t\right] \prod_{i=1}^{p} \exp \left[\frac{1}{k} \gamma_{i}\left(\tilde{A}_{i} \otimes \tilde{A}_{i}\right) t\right]\right\}^{k} \tag{5.10}
\end{equation*}
$$

For convenience, let $S$ and $N$ be $r \times r$ matrices with $N \geq 0$. Since (see [B3])

$$
\begin{equation*}
\operatorname{vec}^{-1}[(S \otimes S) \operatorname{vec} N]=S N S^{\mathrm{T}} \geq 0 \tag{5.11}
\end{equation*}
$$

and

$$
\begin{equation*}
\left(S^{k} \otimes S^{k}\right)(S \otimes S)=S^{k+1} \otimes S^{k+1} \tag{5.12}
\end{equation*}
$$

it follows that

$$
\begin{equation*}
\operatorname{vec}^{-1}\left[e^{S \otimes s} \operatorname{vec} N\right]=\sum_{k=0}^{\infty}(k!)^{-1} S^{k} N S^{k T} \geq 0 \tag{5.13}
\end{equation*}
$$

Furthermore,

$$
\begin{equation*}
\operatorname{vec}^{-1}\left[e^{S \oplus S} \operatorname{vec} \mathrm{~N}\right]=\operatorname{vec}^{-1}\left[\left(e^{S} \otimes e^{S}\right) \operatorname{vec} N\right]=e^{S} N e^{S^{\mathrm{T}}} \geq 0 \tag{5.14}
\end{equation*}
$$

Applying (5.13) and (5.14) alternately with (5.10) and using induction on $k$ it follows that the integrand of (5.9) is nonnegative definite. To prove the converse, note that it follows from (5.5) that $\mathcal{Q}$ satisfies

$$
\begin{equation*}
\mathcal{Q}=\operatorname{vec}^{-1}\left[e^{A t} \operatorname{vec} \mathcal{Q}\right]+\int_{0}^{t} \operatorname{vec}^{-1}\left[e^{A} \operatorname{vec} \tilde{V}\right] d s, \quad t \in[0, \infty) \tag{5.15}
\end{equation*}
$$

Since the integral term on the right hand side of (5.15) is nonnegative definite, is bounded from above by $\mathcal{Q}$, and $\tilde{V} \in \mathbb{N}^{\bar{n}}$ is arbitrary, it follows that $R$ is asymptotically stable.

Proposition 5.2 shows that a solution to (5.5) exists so long as $\alpha_{1}, \ldots, \alpha_{p}$ are sufficiently small that $A$ remains stable for some $\alpha>0$. The following result characterizes values of $\alpha_{1}, \ldots, \alpha_{p}$ for which $\mathcal{A}$ is asymptotically stable. Let $\|\cdot\|$ denote an arbitrary vector norm and its induced matrix norm.

Proposition 5.3. Let $\left(A_{c}, B_{c}, C_{c}\right)$ be given, assume $\tilde{A}$ is asymptotically stable, and let $\alpha, \alpha_{1}, \ldots, \alpha_{p}>0$. If

$$
\begin{equation*}
\left\|(\bar{A} \oplus \tilde{A})^{-1}\left(\alpha I_{\tilde{n}^{2}}+\sum_{i=1}^{p} \gamma_{i} \bar{A}_{i} \otimes \bar{A}_{i}\right)\right\|<1 \tag{5.16}
\end{equation*}
$$

then there exists $Q \in \mathbb{N}^{\tilde{n}}$ satisfying (5.5) and $A$ is asymptotically stable.
Proof. Define $\left\{\mathcal{Q}_{k}\right\}_{j=0}^{\infty}$ where $\mathcal{Q}_{0}$ satisfies

$$
0=\tilde{A} Q_{0}+\Omega_{0} \tilde{A}^{\mathrm{T}}+\tilde{V},
$$

and $\boldsymbol{Q}_{k+1}$ satisfies

$$
0=\tilde{A} Q_{k+1}+\mathcal{Q}_{k+1} \tilde{A}^{\mathrm{T}}+\Omega\left(\mathcal{Q}_{k}, B_{c}, C_{c}\right)+\tilde{V}
$$

Note that $Q_{k} \geq 0, k=1,2, \ldots$. Hence it follows that

$$
\text { vec } Q_{k+1}-\operatorname{vec} Q_{k}=-(\tilde{A} \oplus \tilde{A})^{-1}\left[\operatorname{vec} \Omega\left(Q_{k}, B_{c}, C_{c}\right)-\operatorname{vec} \Omega\left(Q_{k-1}, B_{c}, C_{c}\right)\right]
$$

and thus

$$
\left\|\operatorname{vec} \mathcal{Q}_{k+1}-\operatorname{vec} \mathcal{Q}_{k}\right\| \leq \|(\tilde{A} \oplus \tilde{A})^{-1}\left(\alpha I_{\tilde{n}^{2}}+\sum_{i=1}^{p} \gamma_{i} \tilde{A}_{i} \otimes \tilde{A}_{i}\| \| \operatorname{vec} \mathcal{Q}_{k}-\operatorname{vec} \mathcal{Q}_{k-1} \|\right.
$$

Using (5.16) it follows that $\mathcal{Q} \triangleq \lim _{k \rightarrow \infty} \mathcal{Q}_{k}$ exists. Thus $\mathcal{Q} \geq 0$ and satisfies (5.5). Furthermore, since $\tilde{V} \in \mathbb{N}^{\tilde{n}}$ can be considered arbitrary, Proposition 5.2 implies that $A$ is asymptotically stable.

## 6. Necessary Conditions for the Auxiliary Minimization Problem

The derivation of the necessary conditions for the Auxiliary Minimization Problem is based upon the Fritz John form of the Lagrange multiplier theorem. Rigorous application of this theorem requires that we further restrict $\left(\mathcal{Q}, A_{c}, B_{c}, C_{c}\right)$ to the open set

$$
\begin{gathered}
S \triangleq\left\{\left(\mathcal{Q}, A_{c}, B_{c}, C_{c}\right): \quad Q \in \mathbb{P}^{\bar{n}}, A \text { is asymptotically stable },\right. \\
\text { and } \left.\left(A_{c}, B_{c}, C_{c}\right) \text { is controllable and observable }\right\} .
\end{gathered}
$$

As will be seen, the constraint $\left(\Omega, A_{c}, B_{c}, C_{c}\right) \in S$ is not required for either robust stability or robust performance since Proposition 4.1 shows that only (3.4), (3.5) and (4.2) are needed. Rather, the set $S$ constitutes sufficient conditions under which the Lagrange multiplier technique is applicable to the Auxiliary Minimization Problem. Specifically, the condition $\mathcal{Q}>0$ replaces (4.2) by an open set constraint, the asymptotic stability of $\AA$ serves as a normality condition which further implies that the dual $\mathcal{P}$ of $\mathcal{Q}$ is nonnegative definite, and ( $A_{c}, B_{c}, C_{c}$ ) minimal is a nondegeneracy condition which implies that the lower right $n_{c} \times n_{c}$ subblocks of $Q$ and $P$ are positive definite thus yielding explicit expressions for $B_{c}$ and $C_{c}$. Note that by Proposition 5.2 the condition that $A$ be
asymptotically stable also implies that (5.5) has a unique, nonnegative solution. Finally, we point out that the stabilizability condition (3.5) and stability condition (3.6) play no role in determining solutions of the suxiliary Minimization Problem.

In order to state the main results we require some additional notation and a lemma concerning pairs of nonnegative-definite matrices. For a real $n \times n$ matrix $Z$ define the set of real diagonalizing matrices

$$
D(Z) \triangleq\left\{\Psi \in \mathbb{R}^{n \times n}: \quad \Psi^{-1} Z \Psi \text { is diagonal }\right\}
$$

and, for a pair of $n \times n$ symmetric matrices $X, Y$ define the set of real contragrediently diagonalizing matrices

$$
C(X, Y) \triangleq\left\{\Psi \in D(X Y): \quad \Psi^{-1} X \Psi^{-T} \text { and } \Psi^{T} Y \Psi \text { are diagonal }\right\}
$$

and the subset of real balancing transformations

$$
B(X, Y) \triangleq\left\{\Psi \in C(X, Y): \quad \Psi^{-1} X \Psi^{-T}=\Psi^{T} Y \Psi\right\}
$$

Of course, a necessary condition for $B(X, Y)$ to be nonempty is that $X, Y$, and $X Y$ all have the same rank. Note that in general

$$
\begin{equation*}
B(X, Y) \subset C(X, Y) \subset D(X Y) \tag{6.1}
\end{equation*}
$$

Obviously, a diagonalizable matrix is either invertible (has no zero eigenvalues) or has semisimple zero eigenvalues. Hence if $D(Z) \neq \varnothing$ then the group generalized inverse $Z^{\#}$ exists as a special case of the Drazin generalized inverse ([CM]). Note that we limit our consideration to diagonalizable matrices with real eigenvalues. Also, note that there is no assumption here that $Z$ is symmetric. Of course, when $Z$ is symmetric the group, Drazin, and Moore-Penrose generalized inverses coincide.

Lemma 6.1. Let $\hat{Q}, \hat{P} \in \mathbb{N}^{n}$ and let $r=\operatorname{rank} \hat{\mathcal{Q}} \hat{P}$. Then the following statements hold:
(i) $\hat{\boldsymbol{Q}} \hat{P}$ has nonnegative eigenvalues.
(ii) $C(\hat{Q}, \hat{P}) \neq \emptyset$.
(iii) $\hat{Q} \hat{P}$ is diagonalizable.
(iv) The $n \times n$ matrix

$$
\begin{equation*}
\tau \triangleq \hat{Q} \hat{P}(\hat{Q} \hat{P})^{\#}=(\hat{Q} \hat{P})^{\#} \hat{Q} \hat{P} \tag{6.2}
\end{equation*}
$$

is idempotent, i.e., $r$ is an oblique projection, and

$$
\begin{equation*}
\operatorname{rank} r=r . \tag{6.3}
\end{equation*}
$$

(v) There exists $G, \Gamma \in \mathbb{R}^{r \times n}$ and invertible $M \in \mathbb{R}^{r \times r}$ such that

$$
\begin{align*}
\hat{Q} \hat{P} & =G^{\mathbf{T}} M \Gamma,  \tag{6.4}\\
\Gamma G^{\mathbf{T}} & =I_{r} . \tag{6.5}
\end{align*}
$$

(vi) If $G, \Gamma \in \mathbb{R}^{r \times n}$ and $M \in \mathbb{R}^{r \times r}$ satisfy (6.4) and (6.5), then

$$
\begin{gather*}
\operatorname{rank} G=\operatorname{rank} \Gamma=\operatorname{rank} M=r,  \tag{6.6}\\
(\hat{Q} \hat{P})^{*}=G^{\mathbf{T}} M^{-1} \Gamma  \tag{6.7}\\
r=G^{\mathbf{T}} \Gamma  \tag{6.8}\\
r G^{\mathbf{T}}=G^{\mathbf{T}}, \quad \Gamma \tau=\Gamma \tag{6.9}
\end{gather*}
$$

(vii) The matrices $G, \Gamma \in \mathbb{R}^{r \times n}$ and $M \in \mathbb{R}^{r \times r}$ satisfying (6.4) and (6.5) are unique except for a change of basis in $\mathbb{R}^{r}$. Furthermore, all such $M$ are diagonalizable with positive eigenvalues.
(viii) If rank $\hat{Q}=\operatorname{rank} \hat{P}=r$ then $B(\hat{Q}, \hat{P}) \neq \emptyset$ and

$$
\begin{align*}
& \hat{Q}=\tau \hat{Q}=\hat{Q} \tau^{T}=\tau \hat{Q} \tau^{T},  \tag{6.10}\\
& \hat{P}=r^{T} \hat{P}=\hat{P} r=\tau^{T} \hat{P} \tau . \tag{6.11}
\end{align*}
$$

Proof. See Appendix A.
A tripie ( $G, M, \Gamma$ ) satisfying (6.4) and (6.5) with $G, \Gamma \in \mathbb{R}^{r \times n}, M \in \mathbb{R}^{r \times r}$, and $r=\operatorname{rank} \hat{Q} \hat{P}$ will be called a projertive factorization of $\hat{Q} \hat{P}$. In particular, we shall set $r=n_{c}$. Furthermore, define the complementary projection

$$
\begin{equation*}
\tau_{\perp} \triangleq I_{n}-\tau \tag{6.12}
\end{equation*}
$$

and, for arbitrary $Q, P, \hat{Q}, \hat{P} \in \mathbb{R}^{n \times n}, G, \Gamma \in \mathbb{R}^{n_{0} \times n}, B_{c} \in \mathbb{R}^{n_{c} \times \ell}, C_{c} \in \mathbb{R}^{m \times n_{c}}$, and $\alpha>0$, define the following notation:

$$
\begin{aligned}
& V_{2 s} \triangleq V_{2}+\sum_{i=1}^{p} \gamma_{i}\left[C_{i}(Q+\hat{Q}) C_{i}^{\mathrm{T}}+D_{i} C_{c} \Gamma \hat{Q} \Gamma^{\mathrm{T}} C_{c}^{\mathrm{T}} D_{i}^{\mathrm{T}}\right], \\
& R_{2 g} \triangleq R_{2}+\sum_{i=1}^{p} \gamma_{i}\left[B_{i}^{\mathrm{T}}(P+\hat{P}) B_{i}+D_{i}^{\mathrm{T}} B_{c}^{\mathrm{T}} G \hat{P} G^{\mathrm{T}} B_{c} D_{i}\right], \\
& Q_{0} \triangleq Q C^{\mathrm{T}}+V_{12}+\sum_{i=1}^{p} \gamma_{i}\left[A_{i}(Q+\hat{Q}) C_{i}^{\mathrm{T}}+A_{i} \hat{Q} \Gamma^{\mathrm{T}} C_{c}^{\mathrm{T}} D_{i}^{\mathrm{T}}\right], \\
& P_{c} \triangleq B^{\mathrm{T}} P+R_{12}^{\mathrm{T}}+\sum_{i=1}^{\mathrm{p}} \gamma_{i}\left[B_{i}^{\mathrm{T}}(P+\hat{P}) A_{i}-D_{i}^{\mathrm{T}} B_{c}^{\mathrm{T}} G \hat{P} A_{i}\right] \text {, } \\
& A_{Q} \triangleq A_{\alpha}-Q_{Q} V_{2 \rho}^{-1} C, \quad A_{P} \triangleq A_{\alpha}-B R_{2,}^{-1} P_{\rho} .
\end{aligned}
$$

The above definitions are for convenience in stating the necessary conditions for the Auxiliary Minimization Problem. This result provides explicit formulae for extremals ( $\mathcal{Q}, A_{c}, B_{c}, C_{c}$ ) of the Auxiliary Minimization Problem. A partial converse shows that this form of the necessary conditions represents no loss of generality with regard to the constraint equation (5.5).

Theorem 6.1. (I) Suppose ( $\left.\mathcal{Q}, A_{c}, B_{c}, C_{c}\right) \in S$ solves the Auxiliary Minimization Problem with $U$ given by (5.1) and $\Omega$ given by (5.3). Then there exist $Q, P, \hat{Q}, \hat{P} \in \mathbb{N}^{n}$ such that, for some projective factorization ( $G, M, \Gamma$ ) of $\hat{Q} \hat{P},\left(\Omega, A_{c}, B_{c}, C_{c}\right)$ are given by

$$
\begin{align*}
\mathcal{Q} & =\left[\begin{array}{cc}
Q+\hat{Q} & \hat{Q} \Gamma^{T} \\
\Gamma \hat{Q} & \Gamma \hat{Q} \Gamma^{T}
\end{array}\right],  \tag{6.13}\\
A_{0} & =\Gamma\left(A-B R_{2}^{-1} P_{s}-Q_{s} V_{2 b}^{-1} C+Q_{s} V_{2 b}^{-1} D R_{2 b}^{-1} P_{s}\right) G^{T},  \tag{6.14}\\
B_{c} & =\Gamma Q_{0} V_{2 \theta}^{-1},  \tag{6.15}\\
C_{c} & =-R_{2 \theta}^{-1} P_{s} G^{T}, \tag{6.16}
\end{align*}
$$

and such that $Q, P, \hat{Q}$, and $\hat{P}$ satisfy

$$
\begin{align*}
& 0=A_{\alpha} Q+Q A_{\alpha}^{\mathrm{T}}+V_{1}+\sum_{i=1}^{p} \gamma_{i}\left[A_{i} Q A_{i}^{\mathrm{T}}+\left(A_{i}-B_{i} R_{2 \varepsilon}^{-1} P_{s}\right) \hat{Q}\left(A_{i}-B_{i} R_{2 s}^{-1} P_{e}\right)^{\mathrm{T}}\right] \\
& -Q_{0} V_{2 e}^{-1} Q_{e}^{T}+r_{\perp} Q_{e} V_{2 e}^{-1} Q_{e}^{T} r_{\perp}^{T},  \tag{6.17}\\
& 0=A_{\alpha}^{\mathrm{T}} P+P A_{\alpha}+R_{1}+\sum_{i=1}^{p} \gamma_{i}\left[A_{i}^{\mathrm{T}} P A_{i}+\left(A_{i}-Q_{s} V_{2 e}^{-1} C_{i}\right)^{\mathrm{T}} \hat{P}\left(A_{i}-Q_{e} V_{2 e}^{-1} C_{i}\right)\right] \\
& -P_{9}^{\mathrm{T}} R_{2 e}^{-1} P_{s}+\tau_{1}^{\mathrm{T}} P_{s}^{\mathrm{T}} R_{2 \rho}^{-1} P_{s} \tau_{\perp},  \tag{6.18}\\
& 0=A_{P} \hat{Q}+\hat{Q} A_{P}^{T}+Q_{0} V_{2 e}^{-1} Q_{i}^{T}-\tau_{\perp} Q_{0} V_{2 e}^{-1} Q_{i}^{T} \tau_{\perp}^{T},  \tag{6.19}\\
& 0=A_{Q}^{\mathrm{T}} \hat{P}+\hat{P} A_{Q}+P_{\rho}^{\mathrm{T}} R_{2}^{-1} P_{\rho}-\tau_{\perp}^{\mathrm{T}} P_{\rho}^{\mathrm{T}} R_{2 \rho}^{-1} P_{\rho} \tau_{\perp}, \tag{6.20}
\end{align*}
$$

$$
\begin{equation*}
\operatorname{rank} \hat{Q}=\operatorname{rank} \hat{P}=\operatorname{rank} \hat{Q} \hat{P}=n_{c} \tag{6.21}
\end{equation*}
$$

Furthermore, the auxiliary cost is given by

$$
\begin{equation*}
J\left(\mathcal{Q}, A_{c}, B_{c}, C_{c}\right)=\operatorname{tr}\left[(Q+\hat{Q}) R_{1}-2 R_{22} R_{2 a}^{-1} P_{s} \hat{Q}+P_{s}^{\mathrm{T}} R_{2 \varepsilon}^{-1} R_{2} R_{2 e}^{-1} P_{c} \hat{Q}\right] . \tag{6.22}
\end{equation*}
$$

(II) Conversely, if there exist $Q, P, \hat{Q}, \hat{P} \in \mathbb{N}^{n}$ satisfying (6.17)-(6.21) with $B_{c}$ and $C_{c}$ given by (6.15) and (6.16), then ( $\mathcal{O}, A_{c}, B_{c}, C_{c}$ ) given by (6.13)-(6.16) satisfy (4.2) and (5.5) with $J\left(\mathcal{Q}, A_{c}, B_{c}, C_{c}\right)$ given by (6.22).

Proof. See Appendix B.
Remark 6.1. Theorem 6.1 presents necessary conditions for the Auxiliary Minimization Problem which explicitly characterize extremal quadruples ( $\mathcal{Q}, A_{c}, B_{c}, C_{c}$ ). These necessary conditions consist of a system of two modified Riccati equations and two modified Lyapunov equations coupled by both the optimal projection $\tau$ and uncertainty bounds. If the uncertainty bounds are deleted, then the results of $[\mathrm{HB}]$ are recovered.

Remark 6.2. When solving (6.17)-(6.21) numerically, the uncertainty terme can be adjusted to examine tradeoffs between performance and robustness. Specifically, the bounds $\alpha_{i}$ and the structure matrices $A_{i}, B_{i}, C_{i}$, and $D_{i}$ appearing in $V_{2 \rho}, R_{2 \Omega}, Q_{0}$, and $P_{s}$ can be varied systematically to determine the region of solvability of (6.17)-(6.21).

Remark 6.3. Although equations (6.17)-(6.21) appear formidable, they are, in fact, quite numerically tractable. For related problems involving coupled Riccati equations, homotopic cortinuation methods have been shown to be effective ([KLJ],[MB]). Similar algorithms for solving (6.17)(6.21) have been developed in [GH], [R], while iterative algorithms are discussed in [G2],[GV], [CY].

Remark 6.4. Because of the presence of $B_{c}$ and $C_{c}$ in the definitions of $V_{2 f}, R_{2 \rho}, Q_{a}$, and $P_{a}$, the optimality conditions (6.17)-(6.20) are coupled with the gain expressions (6.15) and (6.16) for $B_{e}$ and $C_{c}$. When the problem is specialized to the case $D_{i}=0, i=1, \ldots, p$, this coupling disappears and equations (6.17)-(6.20) can be solved without reference to the gain expressions (6.15) and (6.16).

## 7. Sufficient Conditions for Robust Stability and Performance

In this section we combine Theorem 3.1 with Theorem 6.1 (II) to obtain our main result guaranteeing robust stability and performance.

Theorem 7.1. Suppose there exist $Q, P, \hat{Q}, \hat{P} \in \mathbb{N}^{\bar{n}}$ satisfying (6.17)-(6.21) with $B_{c}$ and $C_{c}$ given by (6.15) and (6.16). Then, with ( $\mathcal{Q}, A_{c}, B_{c}, C_{c}$ ) given by (6.13)-(6.16), ( $\tilde{A}+\Delta \tilde{A},[\tilde{V}+$ $\left.\left.\Omega\left(\mathcal{Q}, B_{c}, C_{c}\right)-\left(\Delta \bar{A} Q+Q \Delta \tilde{A}^{T}\right)\right]^{\frac{1}{2}}\right)$ is stabilizable for all $(\Delta A, \Delta B, \Delta C, \Delta D) \in U$ if and only if $\tilde{A}+\Delta \tilde{A}$ is asymptotically stable for all $(\Delta A, \Delta B, \Delta C, \Delta D) \in U$. In this case, the performance (2.7) of the closed-loop system (2.9) satisfies the bound

$$
\begin{equation*}
J\left(A_{c}, B_{c}, C_{c}\right) \leq \operatorname{tr}\left[(Q+\hat{Q}) R_{1}-2 R_{12} R_{2 \rho}^{-1} P_{s} \hat{Q}+P_{s}^{\mathrm{T}} R_{2 \rho}^{-1} R_{2} R_{2}^{-1} P_{\rho} \hat{Q}\right] . \tag{7.1}
\end{equation*}
$$

Proof. The converse portion of Theorem 6.1 implies that $\mathcal{Q}$ given by (6.13) is nonnegative definite and satisfies (5.5) or, equivalently, (3.4). It now follows from Theorem 3.1 that the stabilizability condition (3.5) is equivalent to the asymptotic stability of $\tilde{A}+\Delta \tilde{A}$ for all $(\Delta A, \Delta B, \Delta C, \Delta D) \in U$. In this case Proposition 4.1 yields robust stability and performance. The robust performance bound (7.1) is a restatement of (4.3) utilizing (6.22).

Note that Theorem 7.1 is constructive in nature rather than existential. Specifically, Theorem 7.1 involves a coupled system of modified Riccati/Lyapunov equations (6.17)-(6.21) whose solutions, when they exist, are used to explicitly construct the dynamic feedback gains (6.14)-(6.16) which are guaranteed to provide both robust stability and performance. The following existence result concerns the solvability of (6.17)-(6.21). Let $n_{u}$ denote the dimension of the unstable subspace of the plant dynamics matrix $A$.

Theorem 7.2. Assume $n_{c} \geq n_{u}, R_{1}>0, V_{1}>0$, suppose the nominal plant, i.e., (2.1), (2.2) with $\alpha_{i}=0, i=1, \ldots, p$, is stabilizable and detectable and, in addition, is stabilizable by means of an $n_{c}$ th-order strictly proper dynamic compensator (2.3),(2.4). Then there exist $\bar{\alpha}_{1}, \ldots, \bar{\alpha}_{p}>0$ such that if $\alpha_{i} \in\left[0, \bar{\alpha}_{i}\right), i=1, \ldots, p$, then (6.17)-(6.21) have a solution $Q, P, \hat{Q}, \hat{P} \in \mathbb{N}^{n}$ for which ( $A_{c}, B_{c}, C_{c}$ ) given by (6.14)-(6.16) solves the Robust Stability Problem with robust performance bound (6.22)

Proof. From Theorem 3.1 of [ $R 1$ ],[R2] it follows that there exists a solution to (6.17)-(6.21) which stabilizes the nominal plant. By continuity there exists a neighborhood over which robust stability with performance bound (6.22) holds.

Theorem 7.2 is an existence result which guarantees solvability of the sufficiency conditions over a range of parameter uncertainties. The actual range of uncertainty which can be bounded and the conservatism of the performance bound are, of course, problem dependent.

## 8. Specialization to Full-Order Dynamic Compensation

To draw connections with standard full-order LQG theory, we specialize the results of Sections 6 and 7 to the full-order case, i.e., $n_{c}=n$. As discussed in [HB], in the full-order case $G=\Gamma^{-1}$ and thus $G=\Gamma=\tau=I_{n}$ and $\tau_{\perp}=0$ without loss of generality. To develop further connections with standard LQG theory assume

$$
\begin{equation*}
R_{12}=0, \quad V_{12}=0, \quad D=\Delta D=0 \tag{8.1}
\end{equation*}
$$

Since $\Delta D=0$ we shall write $(\Delta A, \Delta B, \Delta C)$ in place of $(\Delta A, \Delta B, \Delta C, \Delta D)$. Also, for arbitrary $Q, P, \hat{Q}, \hat{P} \in \mathbb{R}^{n \times n}$ and $\alpha>0$ define the following notation:

$$
\begin{array}{ll}
\hat{V}_{2 s} \triangleq V_{2}+\sum_{i=1}^{p} \gamma_{i} C_{i}(Q+\hat{Q}) C_{i}^{\mathrm{T}}, & \hat{R}_{2 \varepsilon} \triangleq R_{2}+\sum_{i=1}^{p} \gamma_{i} B_{i}^{\mathrm{T}}(P+\hat{P}) B_{i} \\
\hat{Q}_{\varepsilon} \triangleq Q C^{\mathrm{T}}+\sum_{i=1}^{p} \gamma_{i} A_{i}(Q+\hat{Q}) C_{i}^{\mathrm{T}}, & \hat{P}_{s} \triangleq B^{\mathrm{T}} P+\sum_{i=1}^{p} \gamma_{i} B_{i}^{\mathrm{T}}(P+\hat{P}) A_{i} \\
\hat{A}_{Q} \triangleq A_{\alpha}-\hat{Q} \hat{Q}_{2 e}^{-1} C, & \hat{A}_{P} \triangleq A_{\alpha}-B \hat{R}_{2 e}^{-1} \hat{P}_{s} .
\end{array}
$$

Theorem 8.1. Let $n_{c}=n$, assume (8.1) is satisfied, and suppose there exist $Q, P, \hat{Q}, \hat{P} \in \mathbb{N}^{n}$ satisfying

$$
\begin{align*}
& 0=A_{\alpha} Q+Q A_{\alpha}^{\mathrm{T}}+V_{1}+\sum_{i=1}^{p} \gamma_{i}\left[A_{i} Q A_{i}^{\mathrm{T}}+\left(A_{i}-B_{i} \hat{R}_{2 \theta}^{-1} \hat{P}_{s}\right) \hat{Q}\left(A_{i}-B_{i} \hat{R}_{2 d}^{-1} \hat{P}_{s}\right)^{\mathrm{T}}\right]-\hat{Z}_{\Delta} \hat{V}_{2 \theta}^{-1} \hat{Q}_{s}^{\mathrm{T}},  \tag{8.2}\\
& 0=A_{\alpha}^{\mathrm{T}} P+P A_{\alpha}+R_{1}+\sum_{i=1}^{p} \gamma_{i}\left[A_{i}^{\mathrm{T}} P A_{i}+\left(A_{i}-\hat{Q}_{e} \hat{V}_{2 \theta}^{-1} C_{i}\right)^{\mathrm{T}} \hat{P}\left(A_{i}-\hat{Q} . \hat{V}_{2 \rho}^{-1} C_{i}\right)\right]-\hat{P}_{e}^{\mathrm{T}} \hat{R}_{2}^{-1} \hat{P}_{\varepsilon},  \tag{8.3}\\
& 0=\hat{A}_{P} \hat{Q}+\hat{Q} \hat{A}_{P}^{\mathrm{T}}+\hat{Q}_{8} \hat{V}_{2 e}^{-1} \hat{Q}_{s}^{\mathrm{T}},  \tag{8.4}\\
& 0=\hat{A}_{Q}^{\mathrm{T}} \hat{P}+\hat{P} \hat{A}_{Q}+\hat{P}_{s}^{\mathrm{T}} \hat{R}_{{ }_{2}}{ }^{\mathbf{1}} \hat{P}_{s}, \tag{8.5}
\end{align*}
$$

and let ( $\mathcal{Q}, A_{c}, B_{c}, C_{c}$ ) be given by

$$
\begin{align*}
Q & =\left[\begin{array}{cc}
Q+\hat{Q} & \hat{Q} \\
\hat{Q} & \hat{Q}
\end{array}\right],  \tag{8.6}\\
A_{c} & =A-B \hat{R}_{2 a}^{-1} \hat{P}_{c}-\hat{Q} \hat{V}_{2 e}^{-1} C,  \tag{8.7}\\
B_{c} & =\hat{Q} \hat{Q}_{2 a}^{-1}  \tag{8.8}\\
C_{c} & =-\hat{R}_{2 \cdot}^{-1} \hat{P}_{e} \tag{8.9}
\end{align*}
$$

Then, $\left(\tilde{A}+\Delta \tilde{A},\left[\tilde{V}+\Omega\left(Q, B_{c}, C_{c}\right)-\left(\Delta \tilde{A} Q+Q \Delta \tilde{A}^{T}\right)\right]^{\frac{1}{2}}\right)$ is stabilizable for all $(\Delta A, \Delta B, \Delta C) \in U$ if and only if $\tilde{A}+\Delta \tilde{A}$ is asymptotically stable for all $(\Delta A, \Delta B, \Delta C) \in U$. In this case the performance
of the closed-loop system (2.9) satisfies the bound

$$
\begin{equation*}
J\left(A_{c}, B_{c}, C_{c}\right) \leq \operatorname{tr}\left[(Q+\hat{Q}) R_{1}+\hat{P}_{s}^{\mathrm{T}} \hat{R}_{2 \rho}^{-1} R_{2} \hat{R}_{2 \rho}^{-1} \hat{P}_{s} \hat{Q}\right] . \tag{8.10}
\end{equation*}
$$

Proof. The proof follows from the reduced-order case given in Appendix B.
Remark 8.1. Theorem 8.1 presents sufficient conditions for robust stability and performance for full-order dynamic compensation. These sufficient conditions comprise a system of two modified Riccati equations and two Lyapunov equations coupled by the uncertainty bounds. This coupling illustrates the breakdown of regulator/estimator separation and shows that the certainty equivalence principle is no longer valid for the LQG result with real-valued structured plant parameter variations. If, however, the uncertainty terms $A_{i}, B_{i}, C_{i}$ are set to zero, it can be seen that (8.4) and (8.5) drop out, while (8.2) and (8.3) reduce to the standard separated Riccati equations of LQG theory.

## 9. Ilustrative Numerical Example

To demonstrate the above results we present an illustrative numerical example. The example chosen was originally used in [D] to illustrate the lack of a guaranteed gain margin for LQG controllers. This example was also considered in [BG1] for a preliminary robustness study. Define

$$
\begin{gathered}
n=2, \quad m=\ell=p=1, \\
A=\left[\begin{array}{ll}
1 & 1 \\
0 & 1
\end{array}\right], \quad B=\left[\begin{array}{l}
0 \\
1
\end{array}\right], \quad C=\left[\begin{array}{ll}
1 & 0
\end{array}\right], \quad D=0, \\
A_{1}=\left[\begin{array}{ll}
0 & 0 \\
0 & 0
\end{array}\right], \quad B_{1}=\left[\begin{array}{l}
0 \\
1
\end{array}\right], \quad C_{1}=\left[\begin{array}{ll}
0 & 0
\end{array}\right], \quad D_{1}=0, \\
R_{1}=V_{1}=\left[\begin{array}{ll}
60 & 60 \\
60 & 60
\end{array}\right], \quad R_{12}=V_{12}=\left[\begin{array}{l}
0 \\
0
\end{array}\right], \quad R_{2}=V_{2}=1 .
\end{gathered}
$$

Note that the system is open-loop unstable and becomes unstabilizable at $\sigma_{1}=-1$. As can easily be seen using root locus, a strictly proper stabilizing controller must be of at least second order. Hence we consider (6.17)-(6.21) with $n_{c}=n$ and thus $\tau_{\perp}=0$. Using algorithms described in $[\mathrm{GH}],[\mathrm{R}]$, controllers were obtained for $\left(\alpha, \alpha_{1}\right)=(.1,1),(.4,2)$ and (1.6,.4). Figure 1 compares the guaranteed robust stability region to the actual robust stability region. Note that the design approach yields greater stability than is guaranteed a priori. This phenomenon is not surprising since even the LQG result may provide arbitrarily high levels of robustness for particular problems while failing to guarantee even minimal robustness for all problems. These results thus demonstrate
the ability of the theory to robustify the LQG result. Interestingly, the form of the actual stability region mimics the classical 6 dB downward/infinite dB upward gain margin of full-state-feedback LQR controllers. Finally, Figure 2 compares guaranteed closed-loop performance to actual closedloop performance over the guaranteed closed-loop robust stability region. Controller gains are given in Table 1.

Acknowledgment. We wish to thank Jill M. Straehla for preparing the manuscript version of this paper and Scott $W$. Greeley for carrying out the numerical computations.


Eigure 1


Figure 2


Table 1

## APPENDDK A: Proof of Lemma 6.1

(i) Clearly $\hat{Q} \hat{P}$ and $\hat{P}\} \hat{Q} \hat{P} \hat{\}}$ have the aame nonzero eigenvalues. Since $\hat{P} \ \dot{Q} \hat{P}\}$ is nonnegative definite, $\boldsymbol{Q} \hat{P}$ has nonnegative eigenvalues.
(ii) The result follows from Theorem 6.2.5 of [RM], p. 123. See also Theorem 4.3 of [G1].
(iii) This result follows from (ii) and (6.1).
(iv) This result follows from the definition of the group generalised inverse (see [CM]). Alternatively, $\operatorname{lot} \hat{Q} \hat{P}=\boldsymbol{\Phi} \boldsymbol{\Psi}^{-1}$, where $\boldsymbol{\Psi} \in D(\hat{\phi} \hat{P}), D=\operatorname{diag}\left(d_{1}, \ldots, d_{n}\right)$. Then $(\dot{Q} \dot{P})^{*}=D^{-1}$, where $D_{(i, i)}^{*}=1 / d_{i}$ if $d_{i} \neq 0$, and $D_{(i, i)}^{*}=0$, if $d_{i}=0, i=1, \ldots, n$. Hence $\dot{Q} \dot{P}(\dot{Q} \dot{P})^{*}=E \varphi^{-2}$ is idempotent, where $E$ is a diagonal matrix with $r$ ones and $n-r$ zeroe on the diagonal. Clearly, (6.3) is valid.
(v) Without loes of generality choose $\Psi$ in the preceding argument so that $D=\operatorname{block-ding}\left(\dot{D}, 0_{n}-r\right)$, where $\hat{D}=\operatorname{diag}\left(d_{1}, \ldots, d_{r}\right), d_{i}>0, i=1, \ldots, r$. Hence

$$
\hat{Q} \hat{P}=\left[\begin{array}{cc}
\dot{D} & 0_{r \times(n-r)} \\
0_{(n-r)} \times r & 0_{n-r}
\end{array}\right]{ }^{-1},
$$

and thus (6.5) holds with

$$
G=\left[\begin{array}{ll}
I_{r} & 0_{r \times(n-r)}
\end{array}\right]^{\mathrm{T}}, \quad M=B, \quad \Gamma=\left[\begin{array}{ll}
I_{r} & 0_{r \times(n-r)}
\end{array} \|^{-1}\right.
$$

(vi) Sylvester's inequality and (6.4) imply that

$$
r=\operatorname{rank} \hat{Q} \hat{P} \leq\{\operatorname{rank} G, \operatorname{rank} M, \operatorname{rank} \Gamma\} \leq r,
$$

which yields (6.6). The expression (6.7) for $(\hat{Q} \hat{P})^{*}$ follows directly form the definition of the group generalized inverse. Furthermore, (6.2), (6.5) and (6.7) imply (6.8), while (6.5) and (6.8) imply (6.9).
(vii) Let both ( $G, M, \Gamma$ ) and an identically dimensioned triple ( $\hat{G}, \hat{M}, \hat{\Gamma}$ ) satisfy (6.4). Then it is easy to verify that $\hat{G}=S^{-1} G, \hat{M}=S M S^{-1}$ and $\hat{\Gamma}=S \Gamma$, where $S=\hat{\Gamma} G^{T}$ and $S^{-1}=\Gamma \hat{G}^{T}$.
(viii) It follows from (ii) that there exists $\Psi \in \mathcal{C}(\hat{Q}, \hat{P})$ such that

$$
\hat{Q}=\Psi\left[\begin{array}{cc}
D_{Q} & 0_{r \times(n-r)} \\
0_{(n-r) \times r} & 0_{n-r}
\end{array}\right] \Psi^{T}, \quad \hat{P}=\Psi^{-T}\left[\begin{array}{cc}
D_{\beta} & 0_{r \times(n-r)} \\
O_{(n-r) \times r} & O_{n-r}
\end{array}\right] \Psi^{-1},
$$

where $D_{Q}$ and $D_{P}$ are positive diagonal. Define

$$
\psi=\varphi\left[\begin{array}{cc}
\left(D_{\Delta} D_{\beta}^{-1}\right)^{t} & 0_{r \times(n-r)} \\
0_{(n-r) \times r} & I_{n-r}
\end{array}\right]
$$

so that

$$
\hat{\psi}^{-1} \hat{Q}^{-T}=\hat{\Psi}^{T} \hat{P}_{\hat{Y}}^{\hat{Y}}=\left[\begin{array}{cc}
\left(D_{\phi} D_{\rho}\right) & 0_{r \times(n-r)} \\
0_{(n-r) \times r} & 0_{n-r}
\end{array}\right]
$$

and thus $\hat{\Psi} \in B(\hat{Q}, \hat{P})$. Finally, (6.10) and (6.11) are immediate. $\square$

## APPENDIX B: Proof of Theorem 6.1

To optimise (4.10) over the open set $S$ subject to the contraint (5.5), form the Lagrangian

$$
\begin{equation*}
\mathcal{L}\left(Q, A_{e}, B_{a}, C_{a}, P, \lambda\right) \triangleq \operatorname{tr}\left\{\lambda Q \tilde{R}+\left[\tilde{\Lambda}_{a} Q+Q \tilde{A}_{a}^{T}+\sum_{i=1}^{\rho} \gamma_{i} \bar{\Lambda}_{i} Q \tilde{A}_{i}^{T}+\bar{V}\right] P\right\}, \tag{B.1}
\end{equation*}
$$

where the Lagrange multipliers $\lambda \geq 0$ and $P \in \mathbb{R}^{\boldsymbol{\pi} \times \boldsymbol{\pi}}$ are not both zero. We thus obtain

$$
\begin{equation*}
\frac{\partial \mathcal{L}}{\partial Q}=\bar{A}_{a}^{\mathrm{T}} P+P \bar{A}_{a}+\sum_{i=1}^{P} \gamma_{i} \bar{A}_{i}^{\mathrm{T}} P \bar{A}_{i}+\lambda \tilde{R} . \tag{B.2}
\end{equation*}
$$

Setting $\partial \mathcal{L} / \partial Q=0$ yields

$$
\begin{equation*}
0=\tilde{A}_{a}^{\mathrm{T} P}+P \tilde{A}_{\alpha}+\sum_{i=1}^{P} \gamma_{i} \tilde{A}_{i}^{\mathrm{T} P} \tilde{A}_{i}+\lambda \tilde{R} \tag{B.3}
\end{equation*}
$$

or, equivalently,

$$
A^{T} \text { vec } P=-\lambda \text { vec } \tilde{R} .
$$

Since $A$ is sasumed to be stable, $A^{T}$ is invertible, and thus $\lambda=0$ implies $P=0$. Hence, it can be assumed without lose of generality that $\lambda=1$. Purthermore, it followe from Proposition 5.2 with $A, \tilde{V}$ replaced by $A^{T}, \tilde{R}$ that $\rho$ is nonnegative definite.

Now partition $\tilde{n} \times \tilde{n} \mathcal{Q}, \mathcal{P}$ into $n \times n, n \times n_{c}$, and $n_{c} \times n_{c}$ subblocks as

$$
Q=\left[\begin{array}{ll}
Q_{1} & Q_{12} \\
Q_{12}^{T} & Q_{2}
\end{array}\right], \quad P=\left[\begin{array}{cc}
P_{1} & P_{12} \\
P_{12}^{T} & P_{2}
\end{array}\right],
$$

and define the positive-definite matrices

$$
V_{2 s} \triangleq V_{2}+\sum_{i=1}^{p} \gamma_{i}\left[C_{i} Q_{1} C_{i}^{\mathrm{T}}+D_{i} C_{c} Q_{2} C_{c}^{\mathrm{T}} D_{i}^{\mathrm{T}}\right], \quad R_{2 s} \triangleq R_{2}+\sum_{i=1}^{p} \gamma_{i}\left[B_{i}^{\mathrm{T}} P_{1} B_{i}+D_{i}^{\mathrm{T}} B_{c}^{\mathrm{T}} P_{2} B_{c} D_{i}\right] .
$$

Thus, the stationarity conditions for $A_{e}, B_{c}, C_{c}$ are given by

$$
\begin{align*}
\frac{\partial \mathcal{L}}{\partial A_{c}}= & P_{12}^{\mathrm{T}} Q_{12}+P_{2} Q_{2}=0,  \tag{B.4}\\
\frac{\partial \mathcal{L}}{\partial B_{c}}= & P_{2} B_{c} V_{2 \varepsilon}+\left(P_{12}^{\mathrm{T}} Q_{1}+P_{2} Q_{12}^{\mathrm{T}}\right) C^{\mathrm{T}} \\
& +P_{12}^{\mathrm{T}}\left[V_{12}+\sum_{i=1}^{p} \gamma_{i}\left(A_{i} Q_{1} C_{i}^{\mathrm{T}}+A_{i} Q_{12} C_{c}^{\mathrm{T}} D_{i}^{\mathrm{T}}\right)\right]=0,  \tag{B.5}\\
\frac{\partial \mathcal{L}}{\partial C_{c}}= & R_{2} C_{c} Q_{2}+B^{\mathrm{T}}\left(P_{1} Q_{12}+P_{12} Q_{2}\right) \\
& +\left[R_{12}^{\mathrm{T}}+\sum_{i=1}^{p} \gamma_{i}\left(B_{i}^{\mathrm{T}} P_{1} A_{i}+D_{i}^{\mathrm{T}} B_{c}^{\mathrm{T}} P_{12}^{\mathrm{T}} A_{i}\right)\right] Q_{12}=0 . \tag{B.6}
\end{align*}
$$

Expanding (5.5) and (B.3) yields

$$
\begin{align*}
& 0=A_{\alpha} Q_{1}+Q_{1} A_{a}^{\mathrm{T}}+B C_{c} Q_{12}^{\mathrm{T}}+Q_{12} C_{c}^{\mathrm{T}} B^{\mathrm{T}} \\
& +\sum_{i=1}^{P} \gamma_{i}\left[A_{i} Q_{1} A_{i}^{\mathrm{T}}+B_{i} C_{0} Q_{12}^{\mathrm{T}} A_{i}^{\mathrm{T}}+A_{i} Q_{13} C_{e}^{\mathrm{T}} B_{i}^{\mathrm{T}}+B_{i} C_{0} Q_{2} C_{e}^{\mathrm{T}} B_{i}^{\mathrm{T}}\right]+V_{1},  \tag{B.7}\\
& 0=A_{a} Q_{12}+Q_{13} A_{6 \alpha}^{\mathrm{T}}+Q_{12} C_{8}^{\mathrm{T}} D^{\mathrm{T}} B_{6}^{\mathrm{T}}+Q_{1} C^{\mathrm{T}} B_{6}^{\mathrm{T}}+B C_{6} Q_{2} \\
& +\sum_{i=1}^{p} \gamma_{i}\left[A_{i} Q_{1} C_{i}^{\mathrm{T}} B_{a}^{\mathrm{T}}+A_{i} Q_{12} C_{d}^{\mathrm{T}} D_{i}^{\mathrm{T}} B_{a}^{\mathrm{T}}\right]+V_{12} B_{a}^{\mathrm{T}},  \tag{B.8}\\
& 0=A_{\sigma \alpha} Q_{2}+Q_{2} A_{c \alpha}^{\mathbf{T}}+B_{6} C Q_{12}+Q_{12}^{\mathrm{T}} C^{\mathrm{T}} B_{s}^{\mathrm{T}}+B_{s} D C_{6} Q_{2} \\
& +Q_{2} C_{e}^{\mathrm{T}} D^{\mathrm{T}} B_{e}^{\mathrm{T}}+B_{e} V_{2} B_{e}^{\mathrm{T}},  \tag{B.9}\\
& 0=A_{\alpha}^{\mathbf{T}} P_{1}+P_{1} A_{a}+C^{T} B_{d}^{T} P_{12}^{T}+P_{12} B_{c} C \\
& +\sum_{i=1}^{p} \gamma_{i}\left[A_{i}^{\mathrm{T}} P_{1} A_{i}+C_{i}^{\mathrm{T}} B_{d}^{\mathrm{T}} P_{12}^{\mathrm{T}} A_{i}+A_{i}^{\mathrm{T}} P_{12} B_{d} C_{i}+C_{i}^{\mathrm{T}} B_{d}^{\mathrm{T}} P_{2} B_{c} C_{i}\right]+R_{1},  \tag{B.10}\\
& 0=A_{a}^{T} P_{12}+P_{12} A_{a}+P_{12} B_{c} D C_{c}+P_{1} B C_{c}+C^{T} B_{a}^{T} P_{2} \\
& +\sum_{i=1}^{P} \gamma_{i}\left[A_{i}^{\top} P_{1} B_{i} C_{c}+A_{i}^{\top} P_{12} B_{6} D_{i} C_{6}\right]+R_{12} C_{6},  \tag{B.11}\\
& 0=A_{\epsilon \alpha}^{\mathrm{T}} P_{2}+P_{2} A_{\epsilon \sigma}+C_{e}^{\mathrm{T}} B^{\mathrm{T}} P_{12}+P_{12}^{\mathrm{T}} B C_{c}+C_{d}^{\mathrm{T}} R_{2,} C_{a} . \tag{B.12}
\end{align*}
$$

Lemma B.1. $Q_{2}$ and $P_{2}$ are positive definite.
Proof. By a minor extension of resulta from [A], (B.9) can be rewritten as

$$
0=\left(A_{c \alpha}+B_{c} D C_{c}+B_{c} C Q_{12} Q_{2}^{+}\right) Q_{2}+Q_{2}\left(A_{c a}+B_{c} D C_{c}+B_{c} C Q_{12} Q_{2}^{+}\right)^{\mathrm{T}}+B_{c} V_{2 \varepsilon} B_{c}^{\mathrm{T}}
$$

where $Q_{2}^{+}$is the Moore-Penrose or Drazin generalized inverse of $Q_{2}$. Next note that since ( $A_{c a}, B_{c}$ ) is controllable then, by Theorem 3.6 of $[W],\left(A_{c a}+B_{c} D C_{c}+B_{c} C Q_{12} Q_{2}^{+}, B_{c} V_{2 a}^{\frac{1}{2}}\right)$ is also controllable. Now, since $Q_{2}$ and $B_{c} V_{2 c} B_{c}^{\mathrm{T}}$ are nonnegative definite, it follows from Lemma 12.2 of [ W ], that $Q_{2}$ is positive definite. Using (B.12), similar arguments show that $P_{2}$ is positive definite.

Since $R_{24}, V_{2 e}, Q_{2}$, and $P_{2}$ are invertible, (B.4)-(B.6) can be written as

$$
\begin{gather*}
-P_{2}^{-1} P_{12}^{\mathrm{T}} Q_{12} Q_{2}^{-1}=I_{n_{0}}  \tag{B.13}\\
B_{c}=-P_{2}^{-1}\left\{\left(P_{12}^{\mathrm{T}} Q_{1}+P_{2} Q_{12}^{\mathrm{T}}\right) C^{\mathrm{T}}+P_{12}^{\mathrm{T}}\left[V_{12}+\sum_{i=1}^{p} \gamma_{i}\left(A_{i} Q_{1} C_{i}^{\mathrm{T}}+A_{i} Q_{12} C_{c}^{\mathrm{T}} D_{i}^{\mathrm{T}}\right)\right]\right\} V_{2 c}^{-1},  \tag{B.14}\\
C_{c}=-R_{2 \cdot}^{-1}\left\{B^{\mathrm{T}}\left(P_{1} Q_{12}+P_{12} Q_{2}\right)+\left[R_{12}^{\mathrm{T}}+\sum_{i=1}^{p} \gamma_{i}\left(B_{i}^{\mathrm{T}} P_{1} A_{i}+D_{i}^{\mathrm{T}} B_{\mathrm{c}}^{\mathrm{T}} P_{12}^{\mathrm{T}} A_{i}\right)\right] Q_{12}\right\} Q_{2}^{-1} . \tag{B.15}
\end{gather*}
$$

Now define the $n \times n$ matrices

$$
\begin{array}{cc}
Q \triangleq Q_{1}-Q_{12} Q_{2}^{-1} Q_{12}^{\mathrm{T}}, & P \triangleq P_{1}-P_{12} P_{2}^{-1} P_{12}^{\mathrm{T}}, \\
\hat{Q} \triangleq Q_{12} Q_{2}^{-1} Q_{12}^{\mathrm{T}}, & P \triangleq P_{12} P_{2}^{-1} P_{12}^{\mathrm{T}}, \\
T \triangleq-Q_{12} Q_{2}^{-1} P_{2}^{-1} P_{12}^{\mathrm{T}},
\end{array}
$$

and the $n_{c} \times n, n_{c} \times n_{c}$, and $n_{c} \times n$ matrices

$$
G \triangleq Q_{2}^{-1} Q_{12}^{\mathrm{T}}, \quad M \triangleq Q_{2} P_{2}, \quad \Gamma \triangleq-P_{2}^{-1} P_{12}^{\mathrm{T}} .
$$

Note that $\boldsymbol{r}=\boldsymbol{G}^{\mathbf{T}} \Gamma$.
Clearly, $Q, P, \hat{Q}$, and $\hat{P}$ are symmetric and $\hat{Q}$ and $\hat{P}$ are nonnegative definite. To show that $Q$ and $P$ are also nonnegative definite, note that $Q$ is the upper left- hand block of the nonnegativedefinite matrix $\tilde{\mathbf{Q}} \mathbf{Q} \tilde{\mathbf{Q}}^{\mathbf{T}}$, where

$$
\tilde{\mathscr{Q}}=\left[\begin{array}{cc}
I_{n} & -Q_{12} Q_{2}^{-1} \\
O_{n, \times n} & I_{n_{0}}
\end{array}\right] .
$$

Similarly, $P$ is nonnegative definite.
Next note that with the above definitions, (B.13) is equivalent to (6.5) and that (6.4) holds. Hence $\tau=G^{\mathbf{T}} \Gamma$ is idempotent, i.e., $r^{\mathbf{2}}=r$. Furthermore, it is helpful to note the identities

$$
\begin{array}{cl}
\hat{Q}=Q_{12} G=G^{\mathrm{T}} Q_{12}^{\mathrm{T}}=G^{\mathrm{T}} Q_{2} G, & \hat{P}=-P_{12} \Gamma=-\Gamma^{\mathrm{T}} P_{12}^{\mathrm{T}}=\Gamma^{\mathrm{T}} P_{2} \Gamma, \\
r G^{\mathrm{T}}=G^{\mathrm{T}}, & \Gamma \tau=\Gamma, \\
\hat{Q}=\tau \hat{Q}, & \hat{P}=\hat{P} r, \\
\hat{Q} \hat{P}=-Q_{12} P_{12}^{\mathrm{T}} . \tag{B.19}
\end{array}
$$

Using (B.13) and Sylvester's inequality, it follows that

$$
\operatorname{rank} G=\operatorname{rank} \Gamma=\operatorname{rank} Q_{12}=\operatorname{rank} P_{12}=n_{\mathrm{c}}
$$

Now using (B.16) and Sylvester's inequality yields

$$
n_{c}=\operatorname{rank} Q_{12}+\operatorname{rank} G-n_{c} \leq \operatorname{rank} \hat{Q} \leq \operatorname{rank} Q_{12}=n_{c},
$$

which implies that rank $\hat{Q}=n_{c}$. Similarly, rank $\hat{P}=n_{c}$, and rank $\hat{Q} \hat{P}=n_{c}$ follows from (B.19).

The components of $Q$ and $P$ can be written in terms of $Q, P, \hat{Q}, \hat{P}, G$, and $\Gamma$ as

$$
\begin{array}{ll}
Q_{1}=Q+\hat{Q}, & P_{1}=P+\hat{P} \\
Q_{12}=\hat{Q} \Gamma^{\mathrm{T}}, & P_{12}=-\hat{P} G^{\mathrm{T}} \\
Q_{2}=\Gamma \hat{Q} \Gamma^{\mathrm{T}}, & P_{2}=G \hat{P} G^{\mathrm{T}} \tag{B.22}
\end{array}
$$

The expressions (6.13), (6.15), and (6.16) follow from the definition of $\mathcal{Q}$, (B.14) and (B.15). Substituting (B.20)-(B.22) into (B.7)-(B.12) yields

$$
\begin{align*}
& 0=A_{\alpha} Q+Q A_{\alpha}^{\mathbf{T}}+V_{1}+\sum_{i=1}^{p} \gamma_{i}\left[A_{i} Q A_{i}^{\mathrm{T}}+\left(A_{i}-B_{i} R_{2 d}^{-1} P_{s}\right) \hat{Q}\left(A_{i}-B_{i} R_{2}^{-1} P_{d}\right)^{\mathrm{T}}\right] \\
& +A_{P} \hat{Q}+\hat{\boldsymbol{Q}} A_{P}^{\mathbf{T}},  \tag{B.23}\\
& 0=\left[A_{P} \hat{Q}+\hat{Q}\left(G^{\mathrm{T}} A_{c \alpha} \Gamma+Q_{\theta} V_{2 \theta}^{-1} C\right)^{\mathrm{T}}+Q_{\theta} V_{2 \theta}^{-1} Q_{\theta}^{\mathrm{T}}\right] \Gamma^{\mathrm{T}}, \tag{B.24}
\end{align*}
$$

$$
\begin{align*}
& 0=A_{\alpha}^{\mathrm{T}} P+P A_{\alpha}+R_{1}+\sum_{i=1}^{p} \gamma_{i}\left[A_{i}^{\mathrm{T}} P A_{i}+\left(A_{i}-Q_{i} V_{2 j}^{-1} C_{i}\right)^{\mathrm{T}} \hat{P}\left(A_{i}-Q_{0} V_{2 e}^{-1} C_{i}\right)\right]  \tag{B.25}\\
& +A_{Q}^{\mathrm{T}} \hat{P}+\hat{P} A_{Q}, \tag{B.26}
\end{align*}
$$

$$
\begin{align*}
& 0=G\left[\left(G^{\mathrm{T}} A_{c \alpha} \Gamma+B R_{2 \varepsilon}^{-1} P_{f}\right)^{\mathrm{T}} \hat{P}+\hat{P}\left(G^{\mathrm{T}} A_{c a} \Gamma+B R_{2 \varepsilon}^{-1} P_{s}\right)+P_{a}^{\mathrm{T}} R_{2 \varepsilon}^{-1} P_{s}\right] G^{\mathrm{T}} . \tag{B.27}
\end{align*}
$$

Next, computing either $\Gamma$ (B.24) - (B.25) or $G(B .27)-(B .28)$ yields (6.14). Substituting this expression for $A_{c}$ into (B.23), (B.24), (B.27) and (B.28) it follows that (B.25) $=\Gamma$ (B.24) and (B.28) $=G($ B.27 ). Thus, (B.25) and (B.28) are superfluous and can be omitted. Next, using $(\mathrm{B} .23)+G^{\mathrm{T}} \Gamma(\mathrm{B} .24) G-(\mathrm{B} .24) G-[(\mathrm{B} .24) G]^{\mathrm{T}}$ and $G^{\mathrm{T}} \Gamma(\mathrm{B} .24) G-(\mathrm{B} .24) G-[(\mathrm{B} .24) G]^{\mathrm{T}}$ yields (6.17) and (6.19). Using (B.26) $+\Gamma^{\mathrm{T}} G(\mathrm{~B} .27) \Gamma-(\mathrm{B} .27) \Gamma-[(\mathrm{B} .27) \Gamma]^{\mathrm{T}}$ and $\Gamma^{\mathrm{T}} G(\mathrm{~B} .27) \Gamma-(\mathrm{B} .27) \Gamma$ [(B.27) $\Gamma]^{\mathrm{T}}$ yields (6.18) and (6.20).

Finally, to prove the converse we use (6.13)-(6.21) to obtain (5.5) and (B.3)-(B.6). Let $A_{c}, B_{s}, C_{c}, G, \Gamma, \tau, Q, P, \hat{Q}, \hat{P}, Q$ be as in the statement of Theorem 6.1 and define $Q_{1}, Q_{12}, Q_{2}, P_{1}$, $P_{12}, P_{2}$ by (B.20)-(B.22). Using (6.5), (6.15) and (6.16), it is easy to verify (B.5), (B.6). Finally, substitute the definitions of $Q, P, \hat{Q}, \hat{P}, G$, and $\tau$ into (6.17)-(6.20), reverse the steps taken earlier in the proof, and use (6.13)-(6.16) along with (6.5) and (6.8)- (6.11) to obtain (5.5) and (B.3). Finally, note that

$$
Q=\left[\begin{array}{cc}
Q & 0_{n_{\times n_{0}}} \\
0_{n_{0} \times n} & 0_{n_{a_{0}}}
\end{array}\right]+\left[\begin{array}{c}
I_{n} \\
\Gamma
\end{array}\right] \hat{Q}\left[\begin{array}{ll}
I_{n} & \Gamma^{\mathrm{T}}
\end{array}\right]
$$

which shows that $Q \geq 0$ thus verifying (4.2).

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# Unified optimal projection equations for simultaneous reduced-order, robust modelling, estimation and control 

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An optimal design problem which unifies reduced-order modelling, estimation and control problems is stated. Necessary conditions for optimality are obtained in the form of a coupled system of modified Riccati and Lyapunov equations. The results permit treatment of several new problems, such as reduced-order dynamic compensation with partially known disturbances and unified reduced-order control and estimation. Upon appropriate specialization, results obtained previously for the individual problems of reduced-order modelling, estimation and control are recovered. An additional feature is the inclusion of parameter uncertainty bounds so that the necessary conditions for an auxiliary minimization problem serve as sufficient conditions for simultaneous robust, reduced-order modelling, estimation and conirol.

## Notation and definitions

Note. All matrices have real entries.

$$
\mathbf{R} ; \mathbf{R}^{r \times s} ; \mathbf{R}^{r} \quad \text { real numbers; } r \times s \text { real matrices; } \mathbf{R}^{r \times 1}
$$

$I_{r},()^{\mathrm{T}} \quad r \times r$ identity matrix, transpose
$\oplus ; \otimes$ Kronecker sum; Kronecker product (Brewer 1978)
$\mathbf{S}^{r} \quad r \times r$ symmetric matrices
$\mathbf{N}^{r} \quad r \times r$ symmetric non-negative-definite matrices
Pr $r \times r$ symmetric positive-dêinite matrices

$$
Z_{1} \leqslant Z_{2} \quad Z_{2}-\mathbf{Z}_{1} \in \mathbf{N}, \quad Z_{1}, Z_{2} \in \mathbf{S}^{r}
$$

$$
Z_{1}<Z_{2} \quad Z_{2}-Z_{1} \in \boldsymbol{P}^{r}, \quad Z_{1}, Z_{2} \in \mathbf{S}^{\prime}
$$

asymptotically stable matrix
$n, m, \hat{m}, l, l, n_{c}, q, p$ positive integers
$n+n_{c}$
$x, u, y, \hat{y}, x_{c}, y_{e}, y_{m}, \tilde{x} \quad n, m, l, l, n_{c}, q, l, n$-dimensional vectors
$A, \Delta A ; B, \Delta B ; C, \Delta C \quad n \times n$ matrices; $n \times m$ matrices; $l \times n$ matrices
$A_{i}, B_{i}, C_{i} \quad n \times n, n \times m, l \times n$ matrices, $i=1, \ldots, p$
$\delta_{i}, \alpha_{i}$ positive numbers, $i=1, \ldots, p$
$A_{a} \quad A+\frac{1}{2} \sum_{i=1}^{R} \delta_{i} \alpha_{i} I_{n}$
$A_{c}, B_{c}, B_{m}, C_{c}, C_{e}, C_{m} \quad n_{c} \times n_{c} n_{c} \times l, n_{c} \times \hat{m}, m \times n_{c}, q \times n_{c}, l \times n_{c}$ matrices $\hat{B}, \hat{C} \quad n \times \hat{m}, \hat{l} \times n$ matrices
$A_{c \alpha} \quad A_{c}+\frac{1}{2} \sum_{i=1}^{l} \delta_{i} \alpha_{i} I_{n_{c}}$
$L q \times n$ matrix

[^22]$R$ estimation-error weryhting is $p$
$R$ model-error werghting in $F$
$R_{1}$ state weightiag matris in $\boldsymbol{N}^{*}$
$R_{2}$ control weighting matin in
$R_{12} n \times m$ cross weighlias matrix wach that $R_{1}-R_{1}: R_{:}{ }^{1} R_{1}^{1}:$
$\geqslant 0$
(.) indimeasional white notir andependent of $m_{1} 1 \cdot$ and
$w_{2}(\cdot)$
$w_{1}(\cdot)$ n-dimensional white ooste
$w_{2}(\cdot) 1$-dimensional whice notie
$V_{1}$ incensty of $w_{1}(\cdot)$ in $\boldsymbol{N}^{+}$
$V_{1}$ intensity of $w_{2}(\cdot)$ in $p$
$V_{12} n \times 1$ croms intenuity of $w_{1} 1 \cdot 1 . m_{3} 1 \cdot 1$
$\checkmark$ intensity of $m \cdot$ in $\mathrm{p}^{-}$
$E$ expected value

$\lambda_{1} \Delta A\left[\begin{array}{cc}A & B C_{1} \\ B, C & A_{1}\end{array}\right] \cdot\left[\begin{array}{cc}\Delta A & \Delta B C_{1} \\ B, \Delta C & 0\end{array}\right]$

$$
\begin{aligned}
& \dot{w}(t)\left[\begin{array}{c}
w_{1}(t)+B_{m}(t) \\
B_{2} w_{2}(t)+B_{m} m(1)
\end{array}\right]
\end{aligned}
$$

1. Introdection

The problems of quadratically optimal reduced-order modelling extimation and control have been treated in a common (ramework by Hyland and Berniten (1985). Bernstein and Hyiand (1985), and Hyland and Bernstetn (1984). reapectively Specifically, by carrying out a judicious iransformation of varrebles. If was shown that the necessary conditions for optimality could be cast as coupled systems of 2.3 and 4 modified Lyapunov and Riccati equations. respectively. The coupling is via an oblique projection (i.e. idempotent matrix) which anses as a firect consequence of opimality and which determines the geometric structure of the reduced-order model. extimator. or compensator. When the order of the estimator or compensator is set equal to the order of 1 plant, the additional modified Lyapunov equations drop out and the remaining modified Riccati equations reduce to the standard steady-state Riccati equations of Kalman filter and LQG theory.

An immediate benefit of this formulation of the necessary conditions is clanfication of the relationship between the operations of model reduction and estimator or controller design. Specifically, although the additional pair of modified Lyapunov equations arising in the reduced-order estimation and control problems are analogous to the pair of modified Lyapunov equations characterizng the optimal reducedorder model, these equations are now inextricably coupled with the modified Riocati equations characterizing the estimator and controller design. Hence. because of the coupling, this approach is distinct from LQG controller-reduction techniques (see, for
example, Lill and Anderwn 1986, and Jonethere and Siverman 19831 A commpanson between the LQC reduction methods reviwed b) LuU and Anderion (1986) and the optimal projection approach has been given by Greetey and Hyland (1988)

The goal of the present paper ti to unif) the results obtainad pretiously for reduced-order modelling, autumaison and conitrol by denviag a seater result which. upon appropriate specialvation, yreds the reduced-order modelling estimation and controt results an specal caser Thus as accomplathed by definung a pencrahered performance functional which secorporites kealurer of all throe ontern Thus the optimization problem involver determining a swate reduced-ordet oysiem whech simultaneounly server as a reduced-order model, extimator and controike (or an) two of these as desured) The necemary conduriss now take the form of a coupled syttem of two modified Lyapunor equasions and rwo moduked liccull equations whech can be speciatized to the necessary condtions obtaned prevoudy for the reduand-order modelling estumation and coatrol problema

There are several motuvations for developing a unitied formulation encompatsias all three results. For canmple, to the full-order case the certanity equralescr prinople impties that the states of the opinal dyname compenstior are ahto optumal extumates of the plant states. This is definutely sor the case for as optumal reducod-order controlke in which the states may bear no resemblance to the plast states The untiod formulation of the present paper, bowever. expremes the deure that compensitor states also provide estimates of selected piant stater. Of courve. except th the full. order case, such a compensator will eenerally be suboptimal from a stnctly conirol point of view sunce the design abso serves as an extmator $A$ anmiar formulation has been considered by Wibon and Kumar (1983)

Additional problems which can be handled to the unutied wettiag tavoive reducedorder estimation and control in the presence of partially known plant draturbences When measurements of disturbance componeats are avalable duriag real-ame operation. such measurements can be used as inputs to the estumstor or controlier to improve performance. Note that this problem incorporates aspects of the modelreduction formulation in which the same white nose signal as injected into both the plant and the design system.

A practical motivation for the unified problem setting is convencence in deveioping numerical algorithms for treating different problems In particular, a zngle algonthm for solving the unified optimal projection equations can readily be used for all special cases without reprogramming (For discussions of numencal alponthms for the optimal projection equations, see Greekey and Hyland 1988. and Richter 1987.)

An additional feature of the results given herein is the treatment of parametric uncertainty in the plant matrices. By bounding the effects of parameter uncertainty on worst-case system performance. the necessery conditions for optimality effectively serve as sufficient conditions for robust stability and performance. A similar approach has been carried out by Bernstein and Haddad (1988). using structured stability radius bounds. In the present paper we use an alternative bound which corresponds to a right shift of the dynamics matrix (or equivalently, an exponential cost weighting) in conjunction with multiplicative white-noise type terms. The effect of multiplicative noise on the optimal projection equations has been developed by Bernstein and Hyland (1988). In the present paper such underlying interpretations will be suppressed since only the bound per se will be needed. Hence. although we use the phrase 'multiplicative white noise' for convenience in referring to the type of bound used. it should be stressed that our treatment of parameter uncertainty is wholly deterministic. (See Bernstein 1987 a, and Haddad 1987, for further background and discussion.)

## 

In this section we state the 'robust performance probiem' for sumulaneous reduced-order modelling. estimation and control along with retated notation for later use. Let $U \subset \subset^{+* *} \times R^{+\infty} \times R^{1: 0}$ denote the sel of uncertain perturbations ( $\Delta A, \Delta B, \Delta C$ ) of the nominal system matrices $A, B$ and $C$.

## Robuss performance oroblem

For fixed $n_{c} \leqslant m$ determine ( $A_{8}, B_{c}, B_{n}, C_{d}, C_{0}, C_{a}$ ) such that, for the augmented system consisting of the nth-ordier controlled and disfurbed plant

$$
\begin{equation*}
\dot{x}(t)=(A+\Delta A) x(t)+(B+\Delta B) M(t)+B w(1)+w_{1}(t) . \quad \mid \in[0, x \mid \tag{2.1}
\end{equation*}
$$

with noisy and non-noisy measurements

$$
\begin{gather*}
M(1)=(C+\Delta C) \times(t)+w_{2}(t)  \tag{2.2}\\
\dot{Y}(1)=C(1) \tag{2.3}
\end{gather*}
$$

and $n_{f}$ th-order design system

$$
\begin{gather*}
\dot{x}_{d}(t)=A_{s} x_{d}(t)+B_{c} x(t)+B_{m} w(t)  \tag{2.4}\\
u(t)=C_{c} x_{s}(t)  \tag{2.5}\\
y_{d}(t)=C_{s} x_{d}(t)  \tag{2.6}\\
y_{m}(t)=C_{\infty} x_{s}(t) \tag{2.7}
\end{gather*}
$$

the performance criterion

$$
\begin{equation*}
J\left(A_{c}, B_{c}, B_{m}, C_{c}, C_{c}, C_{m}\right) \triangleq J_{c}+J_{e}+J_{m} \tag{28}
\end{equation*}
$$

is minimized, where

$$
\begin{align*}
& J_{c} \triangleq \sup _{\text {isa. An. scicu }} \lim \sup _{t \rightarrow \infty} E\left[x^{\top}(t) R_{1} x(t)+2 x^{\top}(t) R_{12} M(t)+u^{\top}(t) R_{2} u(t)\right] \tag{29}
\end{align*}
$$

## Remark 2.1

Suppose there are no uncertainties present, i.e. $\Delta A, \Delta B, \Delta C=0$. By setting $L=0$ and $C=0$, it follows that $J_{e}$ and $J_{m}$ play no role in the optimization problem when $C_{e}$ and $C_{m}$ are both taken to be zero. As will be seen in Theorem 6.1. this is indeed the optimal solution in this case. If, furthermore, $B=0$, then the reduced-order dynamiccompensation problem of Hyland and Bernstein (1984), is recovered. If, alternatively, $R_{1}=0, R_{12}=0, B=0, B=0$ and $C=0$ then the reduced-order state-estimation problem of Bernstein and Hyland (1985) is obtained. Finally, setting $R_{1}=0, R_{12}=0$, $L=0, V_{1}=0, B=0$ and $C=0$ yields the model-reduction problem considered by Hyland and Bernstein (1985).

## Remark 2.2

Suppose $L=0$ and $\mathcal{C}=0$ (so that with $C_{c}$ and $C_{e}$ both zero $J_{E}$ and $J_{m}$ are ineffective) but that $B \neq 0$. In this case, a portion of the plant disturbance, which is
assumed to be measurable during on-line operation, is being fed directly into the compensator. Hence this problem, which generalizes that of Hyland and Bernstein (1984), can be thought of as reduced-order dynamic compensation with partially known disturbances. Similarly, the case $R_{1}=0, R_{12}=0, B=0$ and $\mathcal{C}=0$ but $B \neq 0$ provides a generalization of Bernstein and Hyland (1985), which can be thought of as reduced-order state estimation with partially known disturbances.

For each variation $(\Delta A, \Delta B, \Delta C) \in U$, the augmented system (2.1)-(2.5) can be written as

$$
\begin{equation*}
\dot{x}(t)=(\lambda+\Delta \bar{A}) \dot{x}(1)+\bar{w}(t), \quad t \in[0, \infty) \tag{2.12}
\end{equation*}
$$

where

$$
\begin{equation*}
x(t) \triangleq\left[x^{\top}(t), x_{c}^{\top}(t)\right]^{\top} \tag{2.13}
\end{equation*}
$$

and $\tilde{w}(t)$ is white noise with intensity $\nabla \in \boldsymbol{N}^{4}$.
For the 'robust performance problem' the cost can be expressed in terms of the second-moment matrix of $\bar{x}(\mathrm{f})$. The following result is immediate.

## Proposition 2.1

For given $\left(A_{c}, B_{c}, B_{m}, C_{c}, C_{n}, C_{m}\right)$ and $(\Delta A, \Delta B, \Delta C) \in U$ the second-moment matrix

$$
\begin{equation*}
\boldsymbol{Q}_{\Delta x}(t) \triangleq E\left[\tilde{x}(t) x^{\top}(t)\right], \quad t \in[0, \infty) \tag{2.14}
\end{equation*}
$$

satisfies

$$
\begin{equation*}
\dot{Q}_{\Delta x}(t)=(\Lambda+\Delta \Lambda) \underline{Q}_{\Delta x}(t)+Q_{\Delta x}(t)(\lambda+\Delta \Lambda)^{r}+D_{0}, \quad t \in[0, \infty) \tag{2.15}
\end{equation*}
$$

Furthermore,

$$
\begin{equation*}
J\left(A_{c}, B_{c}, B_{m}, C_{c}, C_{0}, C_{m}\right)=\sup _{1 \Delta \Lambda, \Delta n, \Delta O, U} \lim _{c} \sup _{t \rightarrow \infty} \operatorname{tr} \delta_{a i}(t) \hat{R} \tag{2.16}
\end{equation*}
$$

## 3. Suricient conditions for robmet stability and performance

In practice, steady-state performance is only of interest when the augmented system is stable over $\mathbf{U}$. The following result is immediate.

## Lemma 3.1

Suppose the system (2.12) is stable for all $(\Delta A, \Delta B, \Delta C) \in U$. Then

$$
\begin{equation*}
J\left(A_{c}, B_{c}, B_{m}, C_{r}, C_{e}, C_{m}\right)=\sup _{1 A A, \Delta R . \Delta \cap) U} \operatorname{tr} X_{\Delta i} R \tag{3.1}
\end{equation*}
$$

where $\tilde{Q}_{\Delta x} \in \mathbb{R}^{n}$ is the unique solution to

$$
\begin{equation*}
0=(\bar{A}+\Delta \bar{A}) \bar{Q}_{\Delta x}+\bar{X}_{\Delta x}(\bar{\lambda}+\Delta \bar{X})^{T}+\nabla \tag{3.2}
\end{equation*}
$$

Remark 3.1
When $\mathbf{U}$ is compact, 'sup' in (3.1) can be replaced by 'max'.
Since it is difficult to determine $J\left(A_{c}, B_{c}, B_{m}, C_{c}, C_{e}, C_{m}\right)$ explicitly, we shall seek upper bounds.

Theorem 3.1
Let $\Omega: \boldsymbol{N}^{N} \times \boldsymbol{R}^{n \times 1} \times \boldsymbol{R}^{m \times m_{4}} \rightarrow \mathbf{S}^{4}$ be such that

$$
\begin{gather*}
\Delta \bar{A} Q+Q \Delta \lambda^{\top} \leqslant \Omega\left(Q, B_{c}, C_{c}\right) \\
(\Delta A, \Delta B, \Delta C) \in U, \quad\left(Q, B_{c}, C_{c}\right) \in R^{4} \times R^{a+1} \times R^{m \times a} \tag{3.3}
\end{gather*}
$$

and, for given $\left(A_{c}, B_{c}, B_{m}, C_{c}, C_{c}, C_{m}\right)$, suppose there exists $Q \in \mathcal{N}^{\wedge}$ satisfying

$$
\begin{equation*}
0=\lambda Q+Q A^{\top}+\Omega\left(Q_{c} B_{c}, C_{s}\right)+D \tag{3.4}
\end{equation*}
$$

Furthermore, suppose the pair $\left(\bar{\nabla}^{1 / 2}, A+\Delta \pi\right)$ is detectable for all $(\Delta A, \Delta B, \Delta C) \in U$. Then $X+\Delta \bar{A}$ is asymptotically stable for all $(\Delta A, \Delta B, \Delta C) \in U$.

$$
\begin{equation*}
\boldsymbol{Q}_{\Delta X} \leqslant Q,(\Delta A, \Delta B, \Delta C) \in U \tag{3.5}
\end{equation*}
$$

and

$$
\begin{equation*}
J\left(A_{c}, B_{c}, B_{m}, C_{c}, C_{c}, C_{m}\right) \leqslant \operatorname{tr} Q \hat{R} \tag{3.6}
\end{equation*}
$$

Proof
For all $(\Delta A, \Delta B, \Delta C) \in U$, (3.4) is equivalent to

$$
\begin{equation*}
0=(\bar{A}+\Delta \bar{A}) \mathbf{Q}+\mathbf{Q}(\bar{A}+\Delta \tilde{A})^{\mathrm{T}}+\Psi\left(Q_{,}, B_{c}, C_{c}, \Delta \bar{A}\right)+\bar{V} \tag{3.7}
\end{equation*}
$$

where

$$
\left.\Psi\left(Q, B_{c}, C_{c}, \Delta \lambda\right) \triangleq \Omega Q, B_{c}, C_{c}\right)-\left(\Delta \tilde{A} \mathbf{Q}+\mathbf{Q} \Delta \tilde{A}^{\mathrm{T}}\right)
$$

Note that by (3.3), $\Psi\left(O, B_{c}, C_{c}, \Delta X\right) \geqslant 0$ for all $(\Delta A, \Delta B, \Delta C) \in U$. Since ( $D^{1 / 2}, \tilde{A}$ $+\Delta \tilde{A})$ is detectable for all $(\Delta A, \Delta B, \Delta C) \in U$, it follows from Theorem 3.6 of Wonham (1979), that $\left(\left(\tilde{V}+\Psi\left(Q, B_{c}, C_{c}, \Delta \bar{A}\right)\right)^{1 / 2}, \bar{A}+\Delta \bar{X}\right)$ is detectable for all $(\Delta A, \Delta B, \Delta C) \in U$. Hence Lemma 12.2 of Wonham (1979), implies that $\bar{A}+\Delta \bar{A}$ is asymptotically stable for all $(\Delta A, \Delta B, \Delta C) \in U$.

Next, subtracting (3.2) from (3.7) yields

$$
0=(\tilde{A}+\Delta \tilde{A})\left(\mathbf{Q}-\tilde{Q}_{\Delta \bar{J}}\right)+\left(\mathbf{Q}-\tilde{Q}_{\Delta \lambda}\right)(\tilde{A}+\Delta \tilde{A})^{\mathrm{T}}+\Psi\left(\mathbf{Q}, B_{c}, C_{c}, \Delta \tilde{A}\right)
$$

or, equivalently (since $\bar{A}+\Delta \tilde{A}$ is asymptotically stable),

$$
\mathbf{Q}-\bar{Q}_{\Delta X}=\int_{0}^{\infty} \exp (\tilde{A}+\Delta \tilde{A}) t \Psi\left(Q, B_{c}, C_{c}, \Delta \tilde{A}\right) \exp (\tilde{A}+\Delta \tilde{A})^{T} t d t \geqslant 0
$$

which implies (3.5). Finally, (3.5) and (3.1) yield (3.6).
Remark 3.2
For the dynamic-compensation problem the result that $\tilde{A}+\Delta \tilde{A}$ is asymptotically stable for all $(\Delta A, \Delta B, \Delta C) \in U$ is equivalent to robust stability of the closed-loop system. For the state-estimation and model-reduction problems, however, $\tilde{A}+\Delta \tilde{A}$ is lower block triangular (since $B=0$ ) and block diagonal (since $C=0$ ), respectively. Thus robust stability is equivalent to $A_{c}$ stable and $A+\Delta A$ stable for all $(\Delta A, \Delta B, \Delta C) \in U$.

We also note a sufficient condition for the solution $\mathbf{Q}$ of (3.4) to be positive definite.

Proposition 3.1
Let $\Omega$ be as in Theorem 3.1, let $\left(A_{c}, B_{c}, B_{m}, C_{c} C_{c} C_{m}\right)$ be given, and suppose there exists $Q \in \mathbb{N}^{\mathbf{d}}$ satisfying (3.4). If $\left(V^{1 / 2}, \lambda+\Delta X\right)$ is observable for some $(\Delta A, \Delta B, \Delta C) \in U$, then $Q$ is positive definite.

Proof
If $\left(\bar{V}^{1 / 2}, \bar{A}+\Delta \bar{A}\right)$ is observable for some $(\Delta A, \Delta B, \Delta C) \in U$, then, by Theorem 3.6 of Wonham (1979), $\left(\left(\tilde{V}+\Psi\left(Q, B_{c}, C_{c}, \Delta \bar{A}\right)\right)^{1 / 2}, \tilde{A}+\Delta \bar{A}\right)$ is also observable for the same $(\Delta A, \Delta B, \Delta C) \in U$. It thus follows from (3.7) and Lemma 12.2 of Wonham (1979), that $\mathbf{Q}$ is positive definite.

Remark 3.3
If $\tilde{\nabla}$ is positive definite then the detectability and observability hypotheses of Theorem 3.1 and Proposition 3.1 are automatically satisfied.

Remark 3.4
Theorem 3.1 can be strengthened by noting that the detectability assumption is, in a sense, superfluous. To see this, first note that robust stability concerns only the undisturbed system while $\tilde{\nabla}$ involves the disturbance noise. Hence robust stability is guaranteed by the existence of a solution $Q \in \mathbb{N}$ satisfying (3.4) with $\hat{V}_{\text {replaced by } \alpha I_{n}}$ for some $\alpha>0$. For this replacement detectability is automatic (see previous remark). For robust performance, however, $\mathbf{Q}$ in (3.5) must be obtained from (3.4).

## 4. Uncertainty structure asd right shift maltiplicative white noise bound

The uncertainty set $U$ is assumed to be of the form

$$
\begin{align*}
& U=\left\{(\Delta A, \Delta B, \Delta C) \in R^{n \times n} \times R^{n \times m} \times R^{1 \times n}: \Delta A=\sum_{i=1}^{N} \sigma_{i} A_{i}\right. \\
&\left.\Delta B=\sum_{i=1}^{R} \sigma_{i} B_{i}, \quad \Delta C=\sum_{i=1}^{N} \sigma_{i} C_{i}, \quad\left|\sigma_{i}\right| \leqslant \delta_{i}, \quad i=1, \ldots, p\right\} \tag{4.1}
\end{align*}
$$

where, for $i=1, \ldots, p: A_{i} \in R^{n \times n}, B_{i} \in R^{n \times m}$ and $C_{i} \in R^{1 \times n}$ are fixed matrices denoting the structure of the parametric uncertainty; $\delta_{i}$ is a given uncertainty bound; and $\sigma_{i}$ is an uncertain real parameter. The closed-loop system thus has structured uncertainty of the form

$$
\Delta \tilde{A}=\sum_{i=1}^{N} \sigma_{i} \tilde{A}_{i}
$$

where

$$
A_{i} \triangleq\left[\begin{array}{cc}
A_{i} & B_{i} C_{c}  \tag{4.2}\\
B_{c} C_{i} & 0
\end{array}\right], \quad i=1, \ldots, p
$$

To obtain an explicit gain expression for $\left(A_{c}, B_{c}, B_{m}, C_{c}, C_{e}, C_{m}\right)$ we require that

$$
\begin{equation*}
\left[B_{1} \neq 0 \Rightarrow C_{1}=0\right], \quad i=1, \ldots, p \tag{4.3}
\end{equation*}
$$

That is, for each $i \in\{1, \ldots, p\}$ either $B_{i}$ or $C_{i}$ is zero. Of course, both $B_{i}=0$ and $C_{1}=0$ are possible for a given $i$, and there are no restrictions on $A_{i}$.

Given the structure of $U$ defined by (4.1) we can define the bound satisfying (3.3).

## Proposition 4.1

Let $\alpha_{1}, \ldots, \alpha_{p}$ be arbitrary positive scalars. Then the function

$$
\begin{equation*}
\Omega\left(\mathbf{Q}, B_{c}, C_{c}\right)=\sum_{i=1}^{2} \delta_{i}\left(\alpha_{i} \mathbf{Q}+\alpha_{i}^{-1} A_{i} \mathbf{Q} A_{i}^{\top}\right) \tag{4.4}
\end{equation*}
$$

satisfies (3.3) with $U$ given by (4.1).

## Proof

Note that

$$
\begin{aligned}
0 & \leqslant\left[\sigma_{i}\left(\alpha_{i} / \delta_{i}\right)^{1 / 2} l_{n}-\left(\delta_{i} / \alpha_{i}\right)^{1 / 2} \tilde{A}_{i}\right] Q\left[\sigma_{i}\left(\alpha_{i} / \delta_{i}\right)^{1 / 2} I_{i}-\left(\delta_{i} / \alpha_{i}\right)^{1 / 2} \tilde{A}_{i}\right]^{\mathrm{T}} \\
& =\sigma_{i}^{2}\left(\alpha_{i} / \delta_{i}\right) \mathbf{Q}+\left(\delta_{i} / \alpha_{i}\right) \tilde{A}_{i} Q \tilde{A}_{i}^{\mathrm{T}}-\sigma_{i}\left(\tilde{A}_{i} \mathbf{Q}+\mathbf{Q} \tilde{A}_{i}^{\mathrm{T}}\right)
\end{aligned}
$$

which, since $\sigma_{i}^{2} \leqslant \delta_{i}^{2}$, implies (3.3).

## 5. Auxiliary minimization problem

Rather than minimize the actual cost (2.8), we shall consider the upper bound (3.6). This leads to the following problem.

## Auxiliary minimization problem

Determine ( $A_{c}, B_{c}, B_{m}, C_{c}, C_{e}, C_{m}$ ) and $\mathbf{Q} \in \mathbf{N}^{N}$ which minimize

$$
\begin{equation*}
J\left(A_{c}, B_{c}, B_{m}, C_{c}, C_{e}, C_{m}, \mathbf{Q}\right) \triangleq \operatorname{Tr} \mathbf{Q} \tag{5.1}
\end{equation*}
$$

subject to

$$
\begin{equation*}
0=\tilde{A} \mathbf{Q}+\mathbf{Q} \tilde{A}^{T}+\sum_{i=1}^{R} \delta_{i}\left(\alpha_{i} \mathbf{Q}+\alpha_{i}^{-1} \tilde{A}_{i} \mathbf{Q} \tilde{A}_{i}^{\mathrm{T}}\right)+\tilde{V} \tag{5.2}
\end{equation*}
$$

and

$$
\begin{equation*}
\left(\tilde{V}^{1 / 2}, \tilde{A}+\Delta \tilde{A}\right) \text { is detectable, }(\Delta A, \Delta B, \Delta C) \in U \tag{5.3}
\end{equation*}
$$

## Proposition 5.1

If $\left(A_{c}, B_{c}, B_{m}, C_{c}, C_{e}, C_{m}, \mathbf{Q}\right)$ is admissible, i.e. ( $\left.A_{c}, B_{c}, B_{m}, C_{c}, C_{e}, C_{m}, \mathbf{Q}\right)$ satisfies (5.2) and (5.3), then $\tilde{A}+\Delta \bar{A}$ is asymptotically stable for all $(\Delta A, \Delta B, \Delta C) \in U$ and

$$
\begin{equation*}
J\left(A_{c}, B_{c}, B_{m}, C_{c}, C_{e}, C_{m}\right) \leqslant J\left(A_{c}, B_{c}, B_{m}, C_{c}, C_{e}, C_{m}, \mathbf{Q}\right) \tag{5.4}
\end{equation*}
$$

Proof
With $\Omega$ given by (4.4), Proposition 4.1 implies that (3.3) is satisfied. Furthermore, admissibility implies that (3.4) has a solution $Q \in \hat{N}^{n}$. Hence, with (5.3), the hypotheses of Theorem 3.1 are satisfied so that robust stability with the performance bound (3.6) is guaranteed. Note that with the definition (5.1), (5.4) is merely a restatement of (3.6).

## 6. Necessary conditions for the auxiliary minimization problem

The derivation of the necessary conditions for the 'auxiliary minimization problem' is based upon the Fritz John form of the Lagrange multiplier theorem.

Rigorous application of this technique requires additional technical assumptions. Specifically, we further restrict ( $A_{c}, B_{c}, B_{m}, C_{c}, C_{e}, C_{m}$ ) to the set

$$
\begin{aligned}
& \mathbf{S} \triangleq\left\{\left(A_{c}, B_{c}, B_{m}, C_{c}, C_{e}, C_{m}, \mathbf{Q}\right): \mathbf{Q} \in \mathbf{P}^{N}, \mathbb{A}\right. \text { is asymptotically stable, and } \\
& \\
& \left.\quad\left(A_{c}, B_{c}, C_{c}\right) \text { is minimal }\right\}
\end{aligned}
$$

where

$$
\mathbf{A} \Delta\left(\tilde{A}+\frac{1}{2} \sum_{i=1}^{p} \delta_{i} \alpha_{i} I_{n}\right) \oplus\left(\tilde{A}+\frac{1}{2} \sum_{i=1}^{N} \delta_{i} \alpha_{i} I_{n}\right)+\sum_{i=1}^{R} \gamma_{i} \tilde{A}_{i} \otimes \tilde{A}_{i}
$$

with, for convenience,

$$
\gamma_{i} \triangleq \delta_{i} / \alpha_{i}
$$

The following observation assures us that we can apply Lagrange multipliers over an open constraint set.

## Proposition 6.1

The set $\mathbf{S}$ is open.

Proof
It need only be noted that $\mathbf{S}$ is the intersection of three open sets.

## Remark 6.1

The constraint $\left(A_{c}, B_{c}, B_{m}, C_{c}, C_{e}, C_{m}, \mathbf{Q}\right) \in \mathbf{S}$ is not required for either robust stability or robust performance. As can be seen from the proof of Theorem 6.1 in the Appendix, the set $\mathbf{S}$ constitutes sufficient conditions under which the Lagrange multiplier technique is applicable to the 'auxiliary minimization problem'. Specifically, asymptotic stability of $\bar{A}$ serves as a normality condition which further implies that the dual $\mathbf{P}$ of $\mathbf{Q}$ is non-negative definite. Furthermore, $\left(A_{c}, B_{c}, C_{c}\right)$ minimal is a nondegeneracy condition which implies that the lower right $n_{c} \times n_{c}$ subblocks of $Q$ and $P$ are positive definite. It is extremely important to emphasize that Proposition 5.1 implies that it is not necessary for guaranteed robust stability and performance that an admissible quadruple, obtained by solving the necessary conditions, actually be shown to be an element of $\mathbf{S}$.

For arbitrary $Q, P, Q, P \in R^{n \times n}$ define the following notation:

$$
\begin{array}{cl}
R_{2 s} \triangleq R_{2}+\sum_{i=1}^{p} \gamma_{i} B_{i}^{\mathrm{T}}(P+\hat{P}) B_{i}, & V_{2 s} \triangleq V_{2}+\sum_{i=1}^{R} \gamma_{i} C_{i}(Q+\hat{Q}) C_{i}^{\mathrm{T}} \\
Q_{s} \triangleq Q C^{\mathrm{T}}+V_{12}+\sum_{i=1}^{R} \gamma_{i} A_{i}(Q+\hat{Q}) C_{i}^{\mathrm{T}}, & P_{s} \triangleq B^{\mathrm{T}} P+R_{12}^{\mathrm{T}}+\sum_{i=1}^{p} \gamma_{i} B_{i}^{\mathrm{T}}(P+\hat{P}) A_{i} \\
A_{Q} \triangleq A_{a}-Q V_{2 s}^{-1} C, & A_{P} \triangleq A_{a}-B R_{2 s}^{-1} P_{s}
\end{array}
$$

The following factorization lemma is needed for the statement of the main result.

## Lemma 6.1

If $\hat{Q}, \hat{P} \in \mathcal{N}^{n}$ then $\hat{Q} \hat{P}$ is diagonalizable with non-negative eigenvalues. If, in addition, rank $Q \hat{P}=n_{c}$, then there exist $n_{c} \times n G, \Gamma$ and $n_{c} \times n_{c}$ invertible $M$ such that

$$
\begin{equation*}
\oint \mathscr{P}=G^{\mathrm{T}} M \Gamma \tag{6.1}
\end{equation*}
$$

$$
\begin{equation*}
\Gamma G^{\mathrm{T}}=I_{n_{c}} \tag{6.2}
\end{equation*}
$$

Furthermore, $G, M$ and $\Gamma$ are unique except for a change of basis in $R^{n_{c}}$.

## Proof

The result is an immediate consequence of Rao and Mitra (1971), Theorem 6.2.5, p. 123.

A triple ( $G, M, \Gamma$ ) satisfying (6.1) and (6.2) will be called a projective factorization of $\hat{Q} \hat{P}$. Since $\hat{Q} \hat{P}$ is diagonalizable it has a group generalized inverse $(\hat{Q} \hat{P})^{*}=G^{\mathrm{T}} M^{-1} \Gamma$ and

$$
\tau \triangleq \hat{Q} \hat{P}(\hat{Q} \hat{P})^{*}=G^{\mathrm{T}} \Gamma
$$

is an oblique projection. Furthermore, define the complementary projection

$$
\tau_{\perp} \triangleq I_{n}-\tau
$$

Theorem 6.1
Suppose $\left(A_{c}, B_{c}, B_{m}, C_{c}, C_{e}, C_{m}, \mathbf{Q}\right) \in \mathbf{S}$ solves the 'auxiliary minimization problem'. Then there exist $Q, P, \hat{Q}, \hat{P} \in \mathbb{N}^{n}$ such that, for some projective factorization $(G, M, \Gamma)$ of $\hat{Q} \hat{P}, A_{c}, B_{c}, B_{m}, C_{c}, C_{e}, C_{m}$ and $Q$ are given by

$$
\begin{gather*}
A_{\mathrm{c}}=\Gamma\left(A-B R_{2 s}^{-1} P_{s}-Q_{s} V_{2 s}^{-1} C\right) G^{\mathrm{T}}  \tag{6.3}\\
B_{c}=\Gamma Q_{s} V_{2 s}^{-1}  \tag{6.4}\\
B_{m}=\Gamma \hat{B}  \tag{6.5}\\
C_{c}=-R_{2 s}^{-1} P_{s} G^{\mathrm{T}}  \tag{6.6}\\
C_{e}=L G^{\mathrm{T}}  \tag{6.7}\\
C_{m}=\hat{C} G^{\mathrm{T}}  \tag{6.8}\\
\mathbf{Q}=\left[\begin{array}{cc}
Q+\hat{Q} & \hat{Q} \Gamma^{\mathrm{T}} \\
\Gamma \hat{Q} & \Gamma \hat{Q} \Gamma^{\mathrm{T}}
\end{array}\right] \tag{6.9}
\end{gather*}
$$

and such that $Q, P, \hat{Q}$ and $\hat{P}$ satisfy

$$
\begin{align*}
& 0=A_{a} Q+Q A_{a}^{\mathrm{T}}+V_{1}+\sum_{i=1}^{p} \gamma_{i}\left[A_{i} Q A_{i}^{\mathrm{T}}+\left(A_{i}-B_{i} R_{2 s}^{-1} P_{s}\right) \hat{Q}\left(A_{i}-B_{i} R_{2 s}^{-1} P_{s}\right)^{\mathrm{T}}\right] \\
&-Q_{s} V_{2 s}^{-1} Q_{s}^{\mathrm{T}}+\tau_{1}\left[Q_{s} V_{2 s}^{-1} Q_{s}^{\mathrm{T}}+\hat{B} V \hat{B}^{\mathrm{T}}\right] \tau_{\perp}^{\mathrm{T}}  \tag{6.10}\\
& 0=A_{a}^{\mathrm{T}} P+P A_{a}+R_{1}+\sum_{i=1}^{p} \gamma_{i}\left[A_{i}^{\mathrm{T}} P A_{i}+\left(A_{i}-Q_{s} V_{2 s}^{-1} C_{i}\right)^{\mathrm{T}} \hat{P}\left(A_{i}-Q_{s} V_{2 s}^{-1} C_{i}\right)\right. \\
&-P_{s}^{\mathrm{T}} R_{2 s}^{-1} P_{s}+\tau_{1}^{\mathrm{T}}\left[P_{s}^{\mathrm{T}} R_{2 s}^{-1} P_{s}+L^{\mathrm{T}} R L+\hat{C}^{\mathrm{T}} \hat{R} \hat{C}\right] \tau_{\perp}  \tag{6.11}\\
& \quad 0=A_{P} \hat{Q}+\hat{Q} A_{P}^{\mathrm{T}}+Q_{s} V_{2 s}^{-1} Q_{s}^{\mathrm{T}}+\hat{B} V \hat{B}^{\mathrm{T}}-\tau_{\perp}\left[Q_{s} V_{2 s}^{-1} Q_{s}^{\mathrm{T}}+\hat{B} V \hat{B}^{\mathrm{T}}\right] \tau_{1}^{\mathrm{T}}  \tag{6.12}\\
& 0= A_{Q}^{\mathrm{T}} \hat{P}+\hat{P} A_{Q}+P_{s}^{\mathrm{T}} \mathrm{P}_{2 s}^{-1} P_{s}+L^{\mathrm{T}} R L+\hat{C}^{\mathrm{T}} \hat{R} \hat{C}-\tau_{1}^{\mathrm{T}}\left[P_{s}^{\mathrm{T}} R_{2 s}^{-1} P_{s}+L^{\mathrm{T}} R L+\hat{C}^{\mathrm{T}} \hat{R} \hat{C}\right] \tau_{\perp} \tag{6.13}
\end{align*}
$$

$$
\begin{equation*}
\operatorname{rank} \hat{Q}=\operatorname{rank} \hat{P}=\operatorname{rank} \phi \hat{P}=n_{c} \tag{6.14}
\end{equation*}
$$

## Proof

For the proof see the Appendix.

## Remark 6.2

As in Remark 2.1 suppose $\Delta A, \Delta B, \Delta C=0$. By setting $L=0, \hat{B}=0$ and $C=0$, (6.10)-(6.13) specialize to the optimal projection equations (2.18)-(2.21) derived by Hyland and Bernstein (1984), with the added features of correlated plant/measurement noise ( $V_{12}$ ) and cross weighting ( $R_{12}$ ). If $R_{1}=0, R_{12}=0, B=0, \hat{B}=0$ and $\bar{C}=0$ then, since $P_{s}=0$, (6.11) drops out and the remaining equations (6.10), (6.12) and (6.13) specialize to (2.10)-(2.12) of Bernstein and Hyland (1985). Finally, if $R_{1}=0$, $R_{12}=0, L=0, V_{1}=0, B=0$ and $C=0$, then, since $Q_{s}=P_{s}=0,(6.10)$ and (6.11) drop out and the remaining equations (6.12) and (6.13) specialize to (2.21) and (2.22) of Hyland and Bernstein (1985).

## Remark 6.3

A more restrictive formulation for unified modelling, estimation and control is to require $C_{c}=C_{e}=C_{m}$ so that $u=y_{e}=y_{m}$. In this case the three outputs of the design system (2.4)-(2.7) are replaced by a single output. Again, the necessary conditions involve a system of four coupled matrix equations similar to (6.10)-(6.13) which specialize to previously known results. Since this formulation requires $m=q=\{$, it appears to be less useful than the three-output formulation.

## 7. Sufficient conditions for robust stability and performance

The main result guaranteeing robust stability and performance for the unified problem can now be stated.

## Theorem 7.1

Suppose there exist $Q, P, \hat{Q}, \hat{P} \in \mathbb{N}^{n}$ satisfying (6.10)-(6.14) and assume ( $\tilde{V}^{1 / 2}, \tilde{A}$ $+\Delta \tilde{A})$ is detectable for all $(\Delta A, \Delta B, \Delta C) \in U$ with $A_{c}, B_{c}, B_{m}, C_{c}, C_{e}, C_{m}$ given by (6.3)-(6.18). Then $\tilde{A}+\Delta \tilde{A}$ is asymptotically stable for all $(\Delta A, \Delta B, \Delta C) \in U$ and the closed-loop system satisfies the performance bound

$$
\begin{align*}
J\left(A_{c}, B_{c}, B_{m}, C_{c}, C_{e}, C_{m}\right) \leqslant & \operatorname{tr}\left[(Q+\hat{Q}) R_{1}+P_{s}^{\mathrm{T}} R_{2 s}^{-1} R_{2} R_{s}^{-1} P_{s} \hat{Q}-2 R_{12} R_{2 s}^{-1} P_{s} \hat{Q}\right. \\
& \left.+Q L^{\mathrm{T}} R L+\hat{C}^{\mathrm{T}} \hat{R} \hat{C}\left(W_{c}-\hat{Q}\right)\right] \tag{7.1}
\end{align*}
$$

where the controllability gramian $W_{c}$ satisfies

$$
\begin{equation*}
0=A W_{\mathrm{c}}+W_{\mathrm{c}} A^{\mathrm{T}}+\hat{B} V \hat{B}^{\mathrm{T}} \tag{7.2}
\end{equation*}
$$

Proof
Theorem 7.1 implies $\mathbf{Q}$ given by (6.9) satisfies (5.2). With the detectability assumption the result follows from Proposition 5.1.
8. Directions for further research Several generalizations remain to be explored. These include:
(a) permit $w(\cdot)$ to be correlated with $w_{1}(\cdot)$ and $w_{2}(\cdot)$;
(b) replace (2.2) with

$$
\begin{equation*}
y(t)=(C+\Delta C) x(t)+(D+\Delta D) u(t)+w_{2}(t) \tag{8.1}
\end{equation*}
$$

(c) replace (2.4) with

$$
\begin{equation*}
\dot{x}_{c}(t)=A_{c} x_{c}(t)+B_{c} y(t)+B_{m} w(t)+w_{3}(t) \tag{8.2}
\end{equation*}
$$

(d) replace (2.5) with

$$
\begin{equation*}
u(t)=C_{c} x_{c}(t)+D_{c} y(t) \tag{8.3}
\end{equation*}
$$

The extension (8.3) has been studied by Bernstein (1987b), for control and by Haddad and Bernstein (1987), for estimation.

## Acknowledgment

The authors are indebted to Scott W. Greeley for bringing to their attention the reduced-order state estimation problem with partially known disturbances which motivated this paper. We also wish to thank Jill Straehla for excellent typing of the original manuscript.

Dr Bernstein was supported in part by the Air Force Office of Scientific Research under Contract AFOSR F49620-86-C-0002.

## Appendix

Proof of Theorem 6.1
Partition $\tilde{n} \times \tilde{n} \mathbf{Q}, \mathbf{P}$ into $n \times n, n \times n_{c}$ and $n_{c} \times n_{c}$ sub-blocks as

$$
\mathbf{Q}=\left[\begin{array}{cc}
Q_{1} & Q_{12} \\
Q_{12}^{\mathrm{T}} & Q_{2}
\end{array}\right], \quad \mathbf{P}=\left[\begin{array}{cc}
P_{1} & P_{12} \\
P_{12}^{\mathrm{T}} & P_{2}
\end{array}\right]
$$

and define the $n \times n$ matrices

$$
\begin{aligned}
Q \triangleq Q_{2}-Q_{12} Q_{2}^{-1} Q_{12}^{\mathrm{T}}, & P \triangleq P_{1}-P_{12} P_{2}^{-1} P_{12}^{\mathrm{T}} \\
Q \triangleq Q_{12} Q_{2}^{-1} Q_{12}^{\mathrm{T}}, & \hat{P} \triangleq P_{12} P_{2}^{-1} P_{12}^{\mathrm{T}}
\end{aligned}
$$

and the $n_{c} \times n, n_{c} \times n_{c}, n_{c} \times n$ matrices

$$
G \triangleq Q_{2}^{-1} Q_{12}^{\mathrm{T}}, \quad M \triangleq Q_{2} P_{2}, \quad \Gamma \triangleq-P_{2}^{-1} P_{12}^{\mathrm{T}}
$$

The existence of $Q_{2}^{-1}$ and $P_{2}^{-1}$ is shown below.
Clearly, $Q, P, \hat{Q}$ and $\hat{P}$ are symmetric and $\hat{Q}$ and $\hat{P}$ are non-negative definite. To show that $Q$ and $P$ are non-negative definite, note that $Q$ is the upper left-hand block of the non-negative-definite matrix $Q^{*} \mathbf{Q} Q^{* T}$, where

$$
Q^{*} \triangleq\left[\begin{array}{cc}
I_{n} & -Q_{12} Q_{2}^{-1} \\
0 & -I_{n_{c}}
\end{array}\right]
$$

Similarly, $P$ is non-negative definite.
To optimize (5.1) over the open set $\mathbf{S}^{\prime}$, where

$$
\mathbf{S}^{\prime} \triangleq\left\{\left(A_{c}, B_{c}, B_{m}, C_{c}, C_{e}, C_{m}, \mathbf{Q}\right) \in \mathbf{S}:(5.3) \text { is satisfied }\right\}
$$

and subject to the constraint (5.2), form the lagrangian

$$
\begin{aligned}
& L\left(A_{c}, B_{c}, B_{m}, C_{c}, C_{e}, C_{m}, \mathbf{Q}, \mathbf{P}, \lambda\right) \\
& \\
& \qquad \operatorname{tr}\left[\lambda \mathbf{Q} \tilde{R}+\left(\tilde{A} \mathbf{Q}+\mathbf{Q} \tilde{A}^{\mathrm{T}}+\sum_{i=1}^{p} \delta_{i}\left(\alpha_{i} \mathbf{Q}+\alpha_{i}^{-1} \tilde{A}_{i} \mathbf{Q} \tilde{A}_{i}^{\mathrm{T}}\right)+\tilde{V}\right) \mathbf{P}\right]
\end{aligned}
$$

where the Lagrange multipliers $\lambda \geqslant 0$ and $\mathbf{P} \in \boldsymbol{P}^{i \times i}$ are not both zero. We thus obtain

$$
\frac{\partial L}{\partial \mathbf{Q}}=\tilde{A}^{\mathbf{T} \mathbf{P}}+\mathbf{P} \tilde{A}+\sum_{i=1}^{R} \delta_{i}\left(\alpha_{i} \mathbf{P}+\alpha_{i}^{-1} \tilde{A}_{i}^{\mathrm{T}} \mathbf{P} \bar{A}_{i}\right)+\lambda \tilde{R}
$$

Setting $\partial \mathrm{L} / \partial \mathbf{Q}=0$ yields

$$
\mathbf{A}^{\top} \operatorname{vec} \mathbf{P}=-\lambda \operatorname{vec} \tilde{R}
$$

where 'vec' is the column-stacking operation (see Brewer 1978). Since $\AA$ is assumed to be stable and thus invertible, $\lambda=0$ implies $P=0$. Hence, it can be assumed without loss of generality that $\lambda=1$.

Furthermore, the stability of $\bar{A}$ implies that $P$ is non-negative definite. The stationarity conditions are given by

$$
\begin{gather*}
\frac{\partial L}{\partial \mathbf{P}}=\tilde{A} \mathbf{Q}+\mathbf{Q} \tilde{A}^{\mathrm{T}}+\sum_{i=1}^{R}\left[\delta_{i} \alpha_{i} Q+\gamma_{i} \tilde{A}_{i} Q \tilde{A}_{i}^{\mathrm{T}}\right]+\tilde{V}=0  \tag{A1}\\
\frac{\partial L}{\partial \mathbf{Q}}=\tilde{A}^{\mathrm{T}} \mathbf{P}+\mathbf{P} \tilde{A}^{\mathrm{T}}+\sum_{i=1}^{\mathrm{R}}\left[\delta_{i} \alpha_{i} \mathbf{P}+\gamma_{i} \tilde{A}_{i}^{\mathrm{T}} \mathbf{P} \tilde{A}_{i}\right]+\tilde{R}=0  \tag{A2}\\
\frac{\partial L}{\partial A_{c}}=P_{12}^{\mathrm{T}} Q_{12}+P_{2} Q_{2}=0  \tag{A3}\\
\frac{\partial L}{\partial B_{c}}=P_{2} B_{c} V_{2 s}+\left(P_{12}^{\mathrm{T}} Q_{1}+P_{2} Q_{12}^{\mathrm{T}}\right) C+P_{12}^{\mathrm{T}}\left(V_{12}+\sum_{i=1}^{R} \gamma_{i} A_{i} Q_{1} C_{i}^{\mathrm{T}}\right)=0  \tag{A4}\\
\frac{\partial L}{\partial C_{c}}=P_{12}^{\mathrm{T}} \hat{B} V+P_{2} B_{m} V=0  \tag{A5}\\
\frac{\partial L}{\partial C_{m}} C_{c} Q_{2}+B^{\mathrm{T}}\left(P_{1} Q_{12}+P_{12} Q_{2}\right)+\left(R_{12}^{\mathrm{T}}+\sum_{i=1}^{R} \gamma_{i} B_{i}^{\mathrm{T}} P_{1} A_{i}\right) Q_{12}=0  \tag{A6}\\
\frac{\partial L}{\partial C_{e}}=-R L Q_{12}+R C_{e} Q_{12}=0  \tag{A7}\\
\frac{\partial L}{\partial C_{m}}=-R \hat{C} Q_{12}+R C_{m} Q_{2}=0 \tag{A8}
\end{gather*}
$$

Expanding (A 1) and (A 2) vields

$$
\begin{align*}
& 0=A_{a} Q_{1}+B C_{c} Q_{12}^{\top}+Q_{1} A_{a}^{\top}+Q_{12} C_{c}^{\mathrm{T}} B^{\mathrm{T}}+\sum_{i=1}^{p} \gamma_{i} \\
& \times\left[A_{i} Q_{1} A_{i}^{\top}\right.\left.+B_{i} C_{c} Q_{12}^{\top} A_{i}^{\top}+A_{i} Q_{12} C_{c}^{\top} B_{i}^{\top}+B_{i} C_{c} Q_{2} C_{c}^{\top} B_{i}^{\top}\right]+V_{1}+\hat{B} V \hat{B}^{\top}  \tag{A9}\\
& 0= A_{a} Q_{12}+Q_{12} A_{c a}^{\top}+B C_{c} Q_{2}+Q_{1} C^{\top} B_{c}^{\top}+\sum_{i=1}^{R} \gamma_{i} A_{i} Q_{1} C_{i}^{\top} B_{c}^{\top} \\
& \quad+V_{12} B_{c}^{\top}+\hat{B} V B_{m}^{\top} \\
& 0= B_{c} C Q_{12}+A_{c a} Q_{2}+Q_{12}^{\top} C^{\top} B_{c}^{\top}+Q_{2} A_{c a}^{\top}+B_{c} V_{2 s} B_{c}^{\top}+B_{m} V B_{m}^{\top} \tag{A11}
\end{align*}
$$

$$
\begin{align*}
0= & P_{1} A_{a}+P_{12} B_{c} C+A_{a}^{\mathrm{T}} P_{1}+C^{\mathrm{T}} B_{c}^{\mathrm{T}} P_{12}^{\mathrm{T}}+\sum_{i=1}^{R} \gamma_{i} \\
& \times\left[A_{i}^{\mathrm{T}} P_{1} A_{i}+C_{i}^{\mathrm{\top}} B_{c}^{\mathrm{T}} P_{i 2}^{\mathrm{T}} A_{i}+A_{i}^{\mathrm{T}} P_{12} B_{c} C_{i}+C_{i}^{\mathrm{T}} B_{c}^{\mathrm{T}} P_{2} B_{c} C_{i}\right]+R_{1}+L^{\mathrm{\top}} R L+\hat{C}^{\mathrm{\top}} R \hat{C} \hat{l} \tag{A12}
\end{align*}
$$

$$
\begin{align*}
& 0= P_{1} B C_{c}+P_{12} A_{c a}+A^{\mathrm{T}} P_{12}+C^{\mathrm{T}} B_{c}^{\mathrm{T}} P_{2}+\sum_{i=1}^{p} \gamma_{i} A_{i}^{\mathrm{T}} P_{1} B_{i} C_{c}+R_{12} C_{c} \\
& \quad-L^{\mathrm{T}} R C_{m}-\hat{C}^{\mathrm{T}} R \hat{C}_{m}  \tag{A13}\\
& 0=A_{c a}^{\mathrm{T}} P_{2}+P_{2} A_{c a}+C_{c}^{\mathrm{T}} B^{\mathrm{T}} P_{12}+P_{12}^{\mathrm{T}} B C_{c}+C_{c}^{\mathrm{T}} R_{2 s} C_{c}+C_{e}^{\mathrm{T}} R C_{e}+C_{m}^{\mathrm{T}} \hat{R} C_{m} \tag{A14}
\end{align*}
$$

## Lemma A 1

$Q_{2}$ and $P_{2}$ are positive definite.

Proof
By a minor extension of the results from Albert (1969), (A 11) can be rewritten as

$$
0=\left(A_{c a}+B_{c} C Q_{12} Q_{2}^{+}\right) Q_{2}+Q_{2}\left(A_{c z}+B_{c} C Q_{12} Q_{2}^{+}\right)^{\mathrm{T}}+\Lambda
$$

where $\Lambda=B_{c} V_{2 s} B_{c}^{\top}+B_{m} V B_{m}^{\top}$ and $Q_{2}^{+}$is the Moore-Penrose or Drazin generalized inverse of $Q_{2}$. Next note that since $\left(A_{c}, B_{c}\right)$ is controllable then, by Theorem 3.6 of Wonham (1979), $\left(A_{c a}+B_{c} C Q_{12} Q_{2}^{+}, \Lambda^{1 / 2}\right)$ is also controllable. Now, since $Q_{2}$ and $\Lambda$ are non-negative definite, it follows from Lemma 12.2 of Wonham (1979), that $Q_{2}$ is positive definite. Using (A 14), similar arguments show $P_{2}$ is positive definite.

Since $\hat{R}, R, R_{2 s}, V, V_{2 s}, Q_{2}$ and $P_{2}$ are invertible, (A 3)-(A 8) can be written as

$$
\begin{gather*}
-P_{2}^{-1} P_{12}^{\mathrm{T}} Q_{12} Q_{2}^{-1}=I_{n c}  \tag{A15}\\
B_{c}=-P_{2}^{-1}\left[\left(P_{12}^{\mathrm{T}} Q_{1}+P_{2} Q_{12}^{\mathrm{T}}\right) C+P_{12}^{\mathrm{T}}\left(V_{12}+\sum_{i=1}^{R} \gamma_{i} A_{i} Q_{1} C_{i}^{\mathrm{T}}\right)\right] V_{2 s}^{-1}  \tag{A16}\\
B_{m}=-P_{2}^{-1} P_{12} \hat{B}  \tag{A17}\\
C_{c}=-R_{2 s}^{-1}\left[B^{\mathrm{T}}\left(P_{1} Q_{12}+P_{12} Q_{2}\right)+\left(R_{12}^{\mathrm{T}}+\sum_{i=1}^{p} \gamma_{i} B_{i}^{\mathrm{T}} P_{1} A_{i}\right) Q_{12}\right] Q_{2}^{-1}  \tag{A18}\\
C_{e}=L Q_{12} Q_{2}^{-1}  \tag{A19}\\
C_{m}=\hat{C} Q_{12} Q_{2}^{-1} \tag{A20}
\end{gather*}
$$

Note that because of (A15), (6.1) and (6.2) hold. Since $Q_{2}$ and $P_{2}$ are positive definite and

$$
Q_{2} P_{2}=P_{2}^{-1 / 2}\left(P_{2}^{1 / 2} Q_{2} P_{2}^{1 / 2}\right) P_{2}^{1 / 2}
$$

$M$ is diagonalizable with positive eigenvalues. It is helpful to note the identities

$$
\begin{gather*}
\hat{Q}=Q_{12} G=G^{\mathrm{T}} Q_{12}^{\mathrm{T}}=G^{\mathrm{T}} Q_{2} G, \quad \hat{P}=-P_{12} \Gamma=-\Gamma^{\mathrm{T}} P_{12}^{\mathrm{T}}=\Gamma^{\mathrm{T}} P_{2} \Gamma  \tag{A21}\\
\tau G^{\mathrm{T}}=G^{\mathrm{T}}, \quad \Gamma \tau=\Gamma  \tag{A22}\\
\hat{Q}=\tau \hat{Q}, \quad \hat{P}=\hat{P}_{\tau}  \tag{A23}\\
\hat{Q} \hat{P}=-Q_{12} P_{12}^{\mathrm{T}} \tag{A24}
\end{gather*}
$$

Using (6.2) and Silvester's inequality, it follows that $\operatorname{rank} G=\operatorname{rank} \Gamma=\operatorname{rank} Q_{12}$ $=$ rank $P_{12}=n_{c}$ which in turn imply (6.14).

The components of $Q$ and $P$ can be written in terms of $Q, P, \hat{Q}, \hat{P}, G$ and $\Gamma$ as

$$
\begin{array}{rlrl}
Q_{1} & =Q+\hat{Q}, \quad P_{1} & =P+\hat{P} \\
Q_{12} & =\hat{Q} \Gamma^{\mathrm{T}}, & P_{12} & =-\hat{P} G^{\mathrm{T}} \\
Q_{2} & =\Gamma \hat{Q} \Gamma^{\mathrm{T}}, \quad P_{2} & =G \hat{P} G^{\mathrm{T}} \tag{A27}
\end{array}
$$

The gain expressions (6.3)-(6.8) and (6.9) follow from (A 16)-(A 20) and the definition of $\mathbf{Q}$. Substituting (A 25)-(A 27) into (A 9)-(A 14) yields

$$
\begin{align*}
0= & A_{a} Q+Q A_{a}^{\mathrm{T}}+V_{1}+\hat{B} V \hat{B}^{\mathrm{T}}+\sum_{i=1}^{p} \gamma_{i} \\
& \times\left[A_{i} Q A_{i}^{\mathrm{T}}+\left(A_{i}-B_{i} R_{2 s}^{-1} P_{s} \hat{Q}\left(A_{i}-B_{i} R_{2 s}^{-1} P_{s}\right)^{\mathrm{T}}\right]+A_{P} \hat{Q}+\hat{Q} A_{P}^{\mathrm{T}}\right.  \tag{A28}\\
& 0=\left[A_{P} \hat{Q}+\hat{Q}\left(\Gamma^{\mathrm{T}} A_{c a}^{\mathrm{T}} G+C^{\mathrm{T}} V_{2 s}^{-1} Q_{s}^{\mathrm{T}}\right)+Q_{s} V_{2 s}^{-1} Q_{s}^{\mathrm{T}}+\hat{B} V \hat{B}^{\mathrm{T}}\right] \Gamma^{\mathrm{T}} \quad(\mathrm{~A}  \tag{A29}\\
0= & {\left[\left(G^{\mathrm{T}} A_{c x} \Gamma+Q_{s} V_{2 s}^{-1} C\right) \hat{Q}+\hat{Q}\left(G^{\mathrm{T}} A_{c s} \Gamma+Q_{s} V_{2 s}^{-1} C\right)^{\mathrm{T}}+Q_{s} V_{2 s}^{-1} Q_{s}^{\mathrm{T}}+\hat{B} V \hat{B}^{\mathrm{T}}\right] \Gamma^{\mathrm{T}} }
\end{align*}
$$

$$
\begin{align*}
0= & A_{a}^{\mathrm{T}} P+P A_{a}+R_{1}+L^{\mathrm{T}} R L+\hat{C}^{\mathrm{T}} \hat{R} \hat{C}+\sum_{i=1}^{R} \gamma_{i}  \tag{A30}\\
& \times\left[A_{i}^{\mathrm{T}} P A_{i}+\left(A_{i}-Q_{s} V_{2 s}^{-1} C_{i}\right)^{\mathrm{T}} \hat{P}\left(A_{i}-Q_{s} V_{2 s}^{-1} C_{i}\right)\right]+A_{Q}^{\mathrm{T}} \hat{P}+\hat{P} A_{Q}  \tag{A31}\\
& \quad 0=\left[A_{Q}^{\mathrm{T}} \hat{P}+\hat{P}\left(G^{\mathrm{T}} A_{c a} \Gamma+B R_{2 s}^{-1} P_{s}\right)+P_{s}^{\mathrm{T}} R_{2 s}^{-1} P_{s}+L^{\mathrm{T}} R L+\hat{C}^{\mathrm{T}} \hat{R} \hat{C}\right] G^{\mathrm{T}}  \tag{A32}\\
0= & G\left[\left(G^{\mathrm{T}} A_{c a} \Gamma+B R_{2 s}^{-1} P_{s}\right)^{\mathrm{T}} \hat{P}+\hat{P}\left(G^{\mathrm{T}} A_{c a} \Gamma+B R_{2 s}^{-1} P_{s}\right)+P_{s}^{\mathrm{T}} R_{2 s}^{-1} P_{s}+L^{\mathrm{T}} R L\right. \\
& \left.+\hat{C}^{\mathrm{T}} \hat{R} \hat{C}\right] G^{\mathrm{T}} \tag{A33}
\end{align*}
$$

Next, computing either $\Gamma$ (A 29) - (A 30) or $G(A 32)$ - (A 33) yields (6.3). Substituting this expression for $A_{c}$ into (A 28), (A 29), (A 31) and (A 32) and using

$$
(\mathrm{A} 28)+G^{\mathrm{T}} \Gamma(\mathrm{~A} 29) G-(\mathrm{A} 29) G-(\mathrm{A} 29 G)^{\mathrm{T}}
$$

and

$$
G^{\mathrm{T}} \Gamma(\mathrm{~A} 29) G-(\mathrm{A} 29) G-(\mathrm{A} 29 G)^{\mathrm{T}}
$$

yields (6.10) and (6.12). Using

$$
(\mathrm{A} 31)+\Gamma^{\mathrm{T}} G(\mathrm{~A} 32) \Gamma-(\mathrm{A} 32) \Gamma-(\mathrm{A} 32 \Gamma)^{\mathrm{T}}
$$

and

$$
\Gamma^{\top} G(A 32) \Gamma-(A 32) \Gamma-(A 32 \Gamma)^{\top}
$$

yields (6.11) and (6.13).
Finally, to show that the preceding development entails no loss of generality in the optimality condition we now use (6.3)-(6.14) to obtain (A 1)-(A 8). Let $A_{c}, B_{c}, B_{m}, C_{c}$, $C_{e}, C_{m}, G, \Gamma, \tau, Q, P, \hat{Q}, \hat{P}, \mathbf{Q}$ be as in the statement of Theorem 6.1 and define $Q_{1}, Q_{12}$, $Q_{2}, P_{1}, P_{12}, P_{2}$ by (A 25)-(A 27). Using (6.2) and (6.4)-(6.8) it is easy to verify (A 4)-(A 8). Finally, substitute the definitions for $Q, P, \hat{Q}, \hat{P}, G$ and $\Gamma$ into (6.10)-(6.13), reverse the steps taken earlier in the proof and use (6.3) (6.8) to obtain (A 1) and (A 2), which completes the proof.

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APPENOXX I: Pobunt Synthaia: Quadratic Bound
 the optimal Projection Equetions with Petermen-fiollot Bourcis," pise Drentenento oritu, Vol. 33, pp. 692-695, 1988.
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D. S. Bernetein and W. K. Hedded, Whe Opeimal Profection Equations with Petersentiollot Bounds: Potunt stability and Pertormance via Fixed-Orter Dymenic Compermation for Syuters with Structured Real-
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D. S. Burnetain, "Robust Stability and Purfocmance via Fixed-ordar Dynemic Comperation," STM J, contre ootin, to appear.
of robuss reduced-order estimation and conurol in the presence of real. valued. structured parancter uncervinity. This was accomplished by incorporating the quadrack uncertaimy biund proposed in [7] within tie opeimal proyection framework

The purpose of de presem note st wo complete thes cycke of resuits by similarly exteadias the resules of [1]. Cur goed is thus to cteran mobust reduced-order anodels over apecified range of paranneric ghan uncerainty. As in $[4]-(6]$, the main idea is to bound the effect of the uncertain parameters on the model-reduction error over the uncertainy rage and then determunc a roduced-order model which minimizes the modet-reduction bound The resulung generalization of the opumal projection equations now serves as a su/ficrent condurion for robust model reduction by virtue of the fact the a bound on the model-roduction error is being minumuzed rather chan the model-reduction error uself These optimalily condicions now comprise a coupled symem of three algebras matnx equasions which reduce to the result of 111 when the uncercaines) bounds are absenx

## il Notation and Definitions

| ม. $\boldsymbol{n}^{\prime \cdots} \cdot \boldsymbol{H}^{\prime}$ | real numbers. x s rea |
| :---: | :---: |
| l. ( ) ${ }^{\text {r }}$ : | OX indemuty matrix. trampose. experted value |
|  | Kronecker sum |
| E', 时, P' | r $\times$ r symmetric. nonnegpive definutc. positivedefinuc mutrices |
| $Z_{1} \leq Z_{2}, Z_{1}<Z_{2}$ |  |
| n, l, no. m; ${ }_{\text {d }}$ | postive inmegers: $n+n_{\text {m }}$ |
| $x_{1}, y_{1}, y_{0}, x_{m}, z$ | n. I, I, na. 1 -dimersional vectors |
| A. $\triangle A: B, C$ | $n \times n$ macrices, $n \times m, l \times n$ matrices |
| $A_{\text {a }} \mathrm{Ba}_{\text {a }} C_{\text {c }}$ | $n_{m} \times n_{m} . n_{m} \times m_{0} / \times n_{m}$ matices |
| $\boldsymbol{\lambda}, \Delta \boldsymbol{\lambda}, \boldsymbol{B}$ | $\left[\begin{array}{cc}A & 0 \\ 0 & A_{-}\end{array}\right] \cdot\left[\begin{array}{cc}\Delta A^{\prime} & 0 \\ 0 & 0\end{array}\right] \cdot\left[\begin{array}{l}B \\ B_{-}\end{array}\right]$ |

$\boldsymbol{R}$
model-reduction errmer-weighting marrix in $\mathbf{p}^{\prime}$
$w(\cdot)$
$V$ m-dimensional white $\mathrm{r} . \mathrm{os}$ e unemsity of $w(\cdot)$ in $8^{-}$

# Robust Reduced-Order Modeling Vie tbe Optimal Projection Equations with Petersem-Hollot Bomnds 

## WASSIM M HADDAD and DENNIS S. BERNSTEIN

Abstract-An optimat reduced-orter modelian problem whb parnmetnic slawl uncertaialy is considered. A modet-redwetion bomsed magested by recest work of Petersen and Hoilot is atilized for guaraateedeg robust reduced-order modeling over a specified range of uncertad paremeters. Necescary comditiogs which gemeralize the optinal projection equations for model reduction are ened to charactertae the redwend-order modet wifich nimimizes the modet-reduction tomad. The optimality equations thus effectively serve as sufficient conditions for characteriaing pobust redoced-order models.

## I. Introdiction

It has been shown in [1]-13] that the first-order necessary conditions for quadraticalls optimal reduced-orter modeling, estimution, and control can be transformed into coupled systems of two. three, and four matrix equations, respectively. This coupling is due to an oblique projection which arises as a direct consequence of optimality. In a series of pupers (4)-|6] the optimal projection approech was generalized to the problems

[^23]For each reduced-order model $\left(A_{m}, B_{m}, C_{m}\right)$ and system variation $\Delta A$ $\in \mathcal{U}$, the augmented system (3.1)-(3.4) is given by

$$
\begin{equation*}
\dot{x}(1)=(\mathbb{A}+\Delta X) \dot{x}(1)+B w(1), \quad t \in[0, \infty) \tag{3.6}
\end{equation*}
$$

where $\boldsymbol{x}(t) \triangleq\left\{x^{T}(t), x_{m}^{T}(t)\right\}^{T}$.

## IV. Sufficient Conditions for Robust Performance

The following result is immediate.
Lemma 4.I: Suppose $A_{m}$ is asymprorically stable and $A+\Delta A$ is asymprotically stable for all $\Delta \mathbb{A} \in \mathcal{U}$. Then

$$
\begin{equation*}
J\left(A_{m}, B_{m}, C_{m}\right)=\sup _{\Delta A \in v} \operatorname{tr} Q_{\Delta \lambda} R \tag{4.1}
\end{equation*}
$$

where $Q_{\Delta X} \hat{E} \lim _{1 \rightarrow \infty} E\left[x^{2}(t) X(t)^{r}\right] \in Y^{A}$ is the unique solution to

$$
\begin{equation*}
0=(\bar{A}+\Delta A) Q_{د X}+Q_{\Delta A}(A+\Delta A)^{\Gamma}+D \tag{4.2}
\end{equation*}
$$

We now determine an upper bound for $J$ given by (4.1).
Theorem 4.1: Let R: y $^{n} \rightarrow \mathbb{E}^{i}$ be such that

$$
\begin{equation*}
\Delta A Q+Q \Delta A^{\top} \leq \Omega(Q), \Delta A \in U, Q \in H^{\prime} \tag{4.3}
\end{equation*}
$$

and. for given ( $A_{m}, B_{m}, C_{m}$ ), suppose there exists $Q \in \mathcal{H}^{\wedge}$ satisfying

$$
\begin{equation*}
0=\bar{A} Q+Q \bar{A}^{r}+Q(Q)+\bar{D} \tag{4.4}
\end{equation*}
$$

and suppose the pair ( ${ }^{10}, \vec{A}+\Delta \mathcal{A}$ ) is stabilizable for all $\mathcal{A} \in \mathcal{U}$. Then $A_{m}$ is asymptotically stable, $A+\Delta A$ is asymptotically stable for all $\Delta \mathcal{A} \in \mathcal{U}$.

$$
\begin{equation*}
\boldsymbol{Q}_{\Delta X} \leq Q, \Delta A \in u \tag{4.5}
\end{equation*}
$$

where $\bar{Q}_{\lambda A}$ satisfies (4.2), and

$$
\begin{equation*}
J\left(A_{m}, B_{m}, C_{m}\right) \leq \operatorname{tr} Q R . \tag{4.6}
\end{equation*}
$$

Proof: See [5].
Remark 4.1: Theorem 4.1 provides sufficient conditions for reducedorder modeling with an upper bound on modeling error. The result yields. in addition. the result that $A_{m}$ and $A+\Delta \mathcal{A}$ are asymptotically stable. Thus, it is important to emphasize that our results are effectively limited to systems which remain asymprotically stable over the class of uncertainties. Relevant applications include, for example, damped flexible structures with uncertain modal data.

## V. Uncertainty Structure and the Petersen-Hollot bound

The uncertainty set $\mathcal{U}$ is assumed to be of the form
$\mathcal{U}=\left\{\Delta A \in \beta^{n \cdot n}: \Delta A=\sum_{1}^{n} D_{1} M, N, E_{i}, M, M_{1}^{?} \leq \bar{M}_{1}, N_{1}^{\top} N, \leq \bar{N}\right.$,

$$
\begin{equation*}
i=1, \cdots, p\} \tag{5.1}
\end{equation*}
$$

where. for $i=1, \cdots, p: D_{1} \in \|^{n \times n}$ and $E_{1} \in\left\{\Omega^{1, \times n}\right.$ are fixed matrices denoting the structure of the uncertainty: $\overline{\mathcal{M}}, \in \boldsymbol{j}^{\prime \prime}$ and $\hat{N}, \in \boldsymbol{N}^{\prime \prime}$ are
 matrices. The augmented system thus has structured uncertainty of the form

$$
\begin{equation*}
\Delta A=\sum_{i=1}^{p} D_{i} M_{i} N_{i} E_{i} \tag{5.2}
\end{equation*}
$$

where

$$
D_{1} \cong\left[\begin{array}{c}
D_{i}  \tag{5.3}\\
0
\end{array}\right] . E_{1} \triangleq\left[\begin{array}{ll}
E_{1} & 0
\end{array}\right]
$$

Proposition 5.1: The function

$$
\begin{equation*}
Q(Q) \triangleq \sum_{i=1}^{n} D_{i} \bar{M}_{1} D_{1}^{r}+Q E_{1}^{r} \bar{N}_{1} E, Q \tag{5.4}
\end{equation*}
$$

satisfies (4.3) with $\mathcal{U}$ given by (5.1).
Proof: For $i=1, \cdots, p$.

$$
\begin{aligned}
0 & \left.\leq\left(D_{1} M_{1}-Q E_{1}^{r} N_{!}^{r}\right] I D, M_{1}-Q E_{1}^{Y} N_{1}^{r}\right]^{r} \\
& =D_{1} M_{1} M_{1}^{r} D_{1}^{r}+Q E_{1}^{r} N_{1}^{T} N_{1} E_{1} Q-\left(D_{1} M_{1} N_{1} E_{1} Q+Q E_{1}^{r} N_{1}^{T} M_{1}^{T} D_{1}^{r}\right) \\
& \leq D_{1} M_{1} D_{1}^{r}+Q E_{1}^{r} N_{1} E_{1} Q-\left(D_{1}, N_{1} E_{1} Q+Q E_{1}^{T} N_{1}^{T} M_{1}^{r} D_{1}^{r}\right)
\end{aligned}
$$

Summing over $i$ yields (4.3).
Remark 5.1: The bound (5.4) was originally used in [7] for unit-ranh perturbations with scalar uncertain parameters. For further details, see [4]-\{6]

## Vi. The Auxiliary Minimization Problem

Our gual is to minimize the error bound (4 6).
Auxiliary Alinimization Problem: Determine (Q, $A_{m}, B_{m}, C_{m}$ ) with $Q \in \mathrm{P}^{\prime \prime}$ which minimize

$$
\begin{equation*}
J\left(Q, A_{m}, B_{m}, C_{m}\right) \in \operatorname{tu} \tag{6.1}
\end{equation*}
$$

subject to

$$
\begin{equation*}
0=A Q+Q \Lambda^{\top}+\sum_{i=1}^{n}\left[D, M, D_{1}^{r}+Q E_{1}^{r} N_{1} E_{1} Q\right]+D \tag{6.2}
\end{equation*}
$$

and

$$
\begin{equation*}
\left(D^{1 / 2}, \Lambda+\Delta \Lambda\right) \text { is stabilizable. } \Delta A \in \text { थ. } \tag{6.3}
\end{equation*}
$$

Proposition 6.1: If $\left(Q, A_{m}, B_{m}, C_{m}\right)$ satisfy $(6.2)$ and (6.3) with $Q \geq$ 0 , then $A_{m}$ is asymptotically stable. $A+\Delta A$ is asymptotically stable for all $\triangle A \in$ U. and

$$
\begin{equation*}
J\left(A_{m}, B_{m}, C_{m}\right) \leq g\left(Q_{0}, A_{m} B_{m}, C_{m}\right) \tag{6.4}
\end{equation*}
$$

Proof: With $\Omega$ given by (5.4). Proposition 5.1 implies that (4.3) is satisfied. Hence, with (6.3), the hypotheses of Theorem 4.I are satisfied so that the system (3.6) is stable over $\mathscr{U}$ with model-reduction bound (4.6) Note that with (6.1), (6.4) is merely a restatement of (4.6)

VII Necessary Conditions for the Auxiliary Minimization Problem

Rigorous application of the Lagrange multiplier technique requires additional technical assumptions. Specifically, we furiter restrict ( $\mathcal{Q} . \mathcal{A}_{m}$. $B_{m}, C_{m}$ ) to the open set
$\delta \vdots\left\{\left(Q . A_{m}, B_{m}, C_{m}\right): Q \in p^{N}, X^{\alpha}\right.$ is asymptotically stable. and ( $A_{m}, B_{m}, C_{m}$ ) is controllable and observable)
where

$$
\mathbb{Q} \vdots\left(\bar{A}+\sum_{i=1}^{n} E_{1}^{r} \bar{N}, E, Q\right) \oplus\left(\bar{A}+\sum_{i=1}^{n} E_{1}^{r} \bar{N}, E_{1} Q\right)
$$

Remark 7.1: The constraint ( $Q, A_{m}, B_{m}, C_{m}$ ) $\in \delta$ is not required for robust reduced-order modeling since, as shown by Proposition 6.1, only (6.2) and (6.3) are required. As will be seen from the proof of Theorem 7.1, the set $\delta$ constitutes sufficient conditions under which the Lagrange multiplier technique is applicable to the auxiliary minimizesion problem. Specifically, $Q>0$ replaces $Q \geq 0$ by an open set constraim, asymptotic stubility of $\dot{A}$ serves as a normality condition which further implies that the dual $P$ of $Q$ is nonnegative definite, and $\left(A_{m}, B_{m}, C_{m}\right)$ minimal is a nondegeneracy condition.

The following factorization lemma is needed for the atacment of the main result. For details, see [1].

Lemma 7.1: If $\boldsymbol{Q}, \boldsymbol{P} \in \mathbb{M}^{\wedge}$ and rank $\boldsymbol{Q} \boldsymbol{P}=n_{m}$, then there exist $n_{m} \times$ $n G, \Gamma$, and $n_{m} \times n_{m}$ invertible $M$ such that

$$
\begin{gather*}
Q \beta=G^{T} M \Gamma,  \tag{7.1}\\
\Gamma G^{r}=I_{m m} . \tag{7.2}
\end{gather*}
$$

Furthermore, $G, M$, and $\Gamma$ are unique except for a change of basis in $\beta^{n} m$. Recall from [1] that

$$
\begin{equation*}
r \triangleq Q P(Q P)^{r}=G^{r} \Gamma \tag{7.3}
\end{equation*}
$$

is an oblique projection, where ( $)^{\prime}$ denotes group generalized inverse. Define the complementary projection $\tau_{\perp} \dot{\otimes} I_{n}-\tau$ and call $(G, M, \Gamma)$ satisfying (7.1). (7.2) a projective factorization of $Q P$. Furthermore, define the notation

$$
D \doteq \sum_{i=1}^{n} D_{1} \bar{M}, D_{1}^{\prime}, E \doteq \sum_{i=1}^{n} E_{1}^{T} \bar{N}, E, .
$$

Theorem 7.I: Suppose (Q. $A_{m}, B_{m}, C_{m}$ ) $\in \delta$ solves the auxiliary minimization problem with $\mathcal{U}$ given by (5.1). Then there exist $Q, Q, P \in$ $\mathrm{N}^{n}$ such that $Q, A_{m}, B_{m}, C_{m}$ are given by

$$
\begin{gather*}
Q=\left[\begin{array}{cc}
Q+Q & Q \Gamma^{T} \\
\Gamma Q & \Gamma Q \Gamma^{r}
\end{array}\right],  \tag{7.4}\\
A_{m}=\Gamma(A+Q E) G^{T},  \tag{7.5}\\
B_{m}=\Gamma B  \tag{7.6}\\
C_{m}=\dot{C} G^{r} \tag{7.7}
\end{gather*}
$$

for sorne projective factorization ( $G, M, \Gamma$ ) of $Q P$, and such that $Q, Q, \beta$ satisfy

$$
\begin{gather*}
0=A Q+Q A^{r}+D+Q E Q+r_{\perp} B V B^{T_{T} T},  \tag{7.8}\\
0=(A+Q E) Q+Q(A+Q E)^{T}+Q E Q+B V B^{T}-T_{1} B V B^{T} T_{A}^{T},  \tag{7.9}\\
0=(A+Q E)^{T} P+P(A+Q E)+C^{r} R C-r_{1}^{T} C^{T} R C \tau_{\perp},  \tag{7.10}\\
\text { rank } Q=\operatorname{rank} P=\operatorname{rank} Q^{\beta}=n_{m} . \tag{7.11}
\end{gather*}
$$

Furhermore. the auxiliary cost is given by

$$
\begin{equation*}
\mathfrak{J}\left(Q, A_{m}, B_{m}, C_{m}\right)=\operatorname{tr} Q C^{\top} R C . \tag{7.12}
\end{equation*}
$$

Conversely, if $Q, Q, P \in$ in satisfy ( 7.8 )-( 7.11 ) then ( $Q, A_{m}, B_{m}, C_{m}$ ) given by (7.4)-(7.7) satisfy $Q \geq 0$ and (6.2) with auxiliary cost given by (7.12).

Proof: See the Appendix.
Theorem 7.1 presents necessary conditions for the auxiliary minimization problem which explicilly characterize extremals ( $Q, A_{m}, B_{m}, C_{m}$ ). These necessary conditions consist of a system of iwo modified Lyapunov equations and one modified Riccati equation coupled by an oblique projection $t$ and uncertainty terms. Setting $D$ and $E$ to zero, i.e., deleting the plant uncenainties. it can be seen that (7.8) drops out while (7.9) and (7.10) reduce to the optimal projection equations for model reduction oblained in [1]. If. alternatively, $n_{m}=n$, then the full-order robust model is given by $A+Q E, B, C$ where $Q$ is given by (7.8) with $\tau_{\perp}=0$.
Remark 7.2: As in the perfect modeling case considered in [1], (7.8)(7.10) may support multiple solutions. When uncertainty is present but a full-order model is desired, then the solution is unique.
Remark 7.3: The conservatism of the bound (7.12) is difficult to predict for two reasons. First, the overbounding (4.3) holds with respect to the partial ordering of the nonnegative-definite marrices for which no scalar measure of conservatism is available. And, second, the bound (4.3) is required to hold for all nonnegative-definite matrices $Q$. The conservatism will thus depend upon the actual value of $Q$ determined by solving (6.2). Numerical experience with related bounds shows that the conservatism is highly problem dependent. See [8].

## Vili. Sufficient Conditions for Robust Reduced-Order Modeling

The main result guaranteeing robust model reduction can now be stated.

Theorem 8.I: Suppose there exist $Q, Q, \beta \in \mathrm{Mn}^{n}$ satisfying (7.8)( 7.11 ) and suppose that ( $D^{1 / 2}, \lambda+\Delta \bar{X}$ ) is stabilizable for all $\Delta A \in U$ with $A_{m}, B_{m}, C_{m}$ given by (7.5)-(7.7) and $\mathcal{U}$ defined by (5.1). Then $A_{m}$ is asymptorically stable. $A+\Delta A$ is asymptotically stable for all $\Delta A \in$ $U$, and the model-reduction criterion satisfies the bound

$$
\begin{equation*}
J\left(A_{m}, B_{m}, C_{m}\right) \leq \operatorname{tr} Q C^{\top} R C \tag{8.1}
\end{equation*}
$$

Proof: The converse of Theorem 7.1 implies that Q given by (7.4) is nonnegative definite and satisfies (6.2). With the stabilizability assumption the result follows from Proposition 6.1.
Remark 8.1: As noted in Remark 4.1. Theorem 8.1 is effeclively limited to systems which remain asymproxically stable over the class of uncerainties.

## APPENDIX

## Proof of Theorem 7.1

To optimize (6.1) over the open set

$$
\delta^{\prime} \triangleq\left\{\left(Q, A_{m}, B_{m}, C_{m}\right) \in \delta:(6.3) \text { is satisfied }\right\}
$$

subject to the constraint (6.2), form the Lagrangian

$$
\begin{aligned}
& \mathcal{L}\left(A_{m}, B_{m}, C_{m}, \mathcal{Q}, \varnothing, \lambda\right) £ \\
& \operatorname{tr}\left\{\lambda Q R+\left[\lambda Q+Q A^{\top}+\sum_{i=1} D_{1} M_{1} D^{\top}+Q E_{1}^{\top}, V_{1} E, Q+V\right] \sigma\right\} \text {. }
\end{aligned}
$$

where the Lagrange multipliers $\lambda \geq 0$ and $\mathscr{P} \in \mathbb{\theta}^{\dagger \times \AA}$ are no both zero. We thus obtain

$$
\frac{\partial \mathcal{L}}{\partial Q}=\bar{A}^{T} P+\varnothing \bar{A}+\sum_{i=1}^{f}\left[E_{1}^{T} \bar{N} E_{1}, Q+P Q E_{i}^{T} \bar{N}, E_{1}\right]+\lambda \boldsymbol{R} .
$$

Selting $\partial \mathcal{L} / \partial Q=0$ yields

$$
\begin{equation*}
\mathbf{a}^{r} \operatorname{vec} \boldsymbol{P}=-\lambda \operatorname{vec} \vec{R} \tag{A.1}
\end{equation*}
$$

where "vec" is the column-stacking operation (see [6]). Since | ब |
| :---: | is assumed to be stable. $\bar{\alpha}^{T}$ is inverible. and thus $\lambda=0$ implics $\theta=0$. Hence, it can be assumed without loss of generality that $\lambda=1$. Furthermore, the stability of $\mathbb{Q}^{T}$ implies that $\mathscr{P}^{\mathbb{C}}$ is nonnegative definite.

Now partion $n \times n \in \mathbb{Q}$. $\mathcal{Q}$ into $n \times n_{1} n \times n_{m}$, and $n_{m} \times n_{m}$ subblocks as

$$
Q=\left[\begin{array}{ll}
Q_{1} & Q_{1} \\
Q_{1:}^{\prime} & Q_{:}
\end{array}\right] \cdot Q=\left[\begin{array}{ll}
P_{1} & P_{1:} \\
P_{12}^{\prime} & P_{2}
\end{array}\right]
$$

Thus, the stationarity conditions are given by

$$
\begin{gather*}
\left.\left.\frac{\partial \mathcal{L}}{\partial Q}=\tilde{A}^{r} Q+P \tilde{i}+\sum_{i=1}^{n} \right\rvert\, E_{1}^{r} \tilde{N} E_{1} Q P+P Q E_{1}^{r} \tilde{N}, E_{1}\right]+R=0 .  \tag{A.2}\\
\frac{\partial \mathcal{L}}{\partial A_{m}}=P_{12}^{r} Q_{12}+P_{2} Q_{2}=0,  \tag{A.3}\\
\frac{\partial \mathcal{L}}{\partial B_{m}}=P_{12}^{r} B V+P_{2} B_{m} V=0,  \tag{A.4}\\
\frac{\partial \mathcal{L}}{\partial C_{m}}=-R C Q_{12}+R C_{m} Q_{2}=0 . \tag{A.S}
\end{gather*}
$$

Expanding (6.2) and (A.2) yields

$$
\begin{gather*}
0=A Q_{1}+Q_{1} A^{T}+D+Q_{1} E Q_{1}+B V B^{r}, \\
0=A Q_{12}+Q_{12} A_{m}^{r}+Q_{1} E Q_{12}+B V B_{m}^{T},  \tag{A.7}\\
0=A_{m} Q_{2}+Q_{2} A_{m}^{r}+Q_{12}^{T} E Q_{12}+B_{m} V B_{m}^{r},  \tag{A.8}\\
0=A^{T} P_{1}+P_{1} A+E\left(P_{1} Q_{1}+P_{12} Q_{12}\right)^{r}+\left(P_{1} Q_{1}+P_{12} Q_{12}\right) E+C^{r} R C, \tag{A.9}
\end{gather*}
$$

$$
\begin{gather*}
0=A^{T} P_{12}+P_{12} A_{m}+E\left(P_{12}^{T} Q_{1}+P_{2} Q_{12}^{\tau}\right)^{T}-C^{\top} R C_{m,}  \tag{A.10}\\
0=A_{m}^{\prime} P_{2}+P_{2} A_{m}+C_{m}^{T} R C_{m} \tag{A.1I}
\end{gather*}
$$

Now define the $n \times n$ nonnegative-definite matrices

$$
\begin{gathered}
Q=Q_{1}-Q_{1} Q_{2}^{\prime} Q_{1:}^{\top} P \dot{=} P_{1}-P_{1:} P_{:}^{\prime} P_{1:}^{\prime} \\
\phi=Q_{1:} Q_{2}^{-1} Q_{12}^{\prime}, P=P_{1:} P_{:}^{\prime} P_{1:}^{\top}
\end{gathered}
$$

and the $n_{m} \times n_{1} n_{m} \times n_{m}, n_{m} \times n$ matrices

$$
G=Q_{2}^{-1} Q_{12}^{r}, M=Q_{:} P_{3}, \Gamma=-P_{2}^{-1} P_{12}^{T}
$$

The existence of $Q_{2}{ }^{\prime}$ and $P_{2}{ }^{\prime}$ follows from the fact that $\left(A_{m}, B_{m}, C_{m}\right)$ is minimal. See $[1]-[4]$, \{6] for details.

Note that (A.3) implies (7.1) and (7.2). Sylvester's inequality yields (7.11). Next (7.4). (7.6). and (7.7) follow from the definition of $Q$. relations (A.4) and (A.5), and the identities

$$
Q_{1}=Q+Q, Q_{12}=\phi \Gamma^{r}, P_{12}=-P G^{r}, Q_{2}=\Gamma Q \Gamma^{r}, P_{2}=G P G^{r}
$$

Computing either $\Gamma(A .7)$. (A.8) or $G(A .10)+(A .11)$ yields (7.5). Inserting (7.5)-(7.7) into (A.6)-(A.11) it can be shown that (A.8), (A.9), and (A.11) are superfluous. Using $(A .6)+G^{T} \Gamma(A .7) C-(A .7) G-$ $\left(\left.(A 7) G\right|^{T}\right.$ and $G^{T} \Gamma(A .7) G-(A .7) G-[(A .7) G)^{T}$ yield (7.8) and (7.9). Similarls. $\Gamma^{7} G(A .10) \Gamma-(A .10) \Gamma-|(A .10) \Gamma|^{7}$ yields (7.10).

Finalis, the proxit can be reversed so that (7.5)-(7.11) yield (A.1)(A.5) and (6.2).

## Reprenenctis

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## Defect Correction Methods for the Solution of Algebraic Riccati Equations

## V. MEHRMANN and E. TAN


#### Abstract

The solution of discrete and continsous algebrale Riccati equations is considered. It is shown that if an approximate solution is obtained, then the defect for this solution again solves an algebraic Riccati equation of the same form and that the system properties of detectability and stabilizability are inherited by this defect equation. On the basis of these results, a general defect correction method is proposed and aumerical examples are given for the use of this method in combination with the SR method.


## I. Introduction

We consider the numerical solution of generalized algebraic Riccati equations

$$
\begin{equation*}
0=A^{*} X E+E^{*} X A-\left(B^{*} X E+S^{*}\right)^{*} R^{-1}\left(B^{*} X E+S^{*}\right)+C^{*} Q C=C R(X) \tag{1.1}
\end{equation*}
$$

$0=-E^{*} X E+A^{*} X A-\left(A^{*} X B+S\left(R+B^{*} X B\right)^{-1}\right.$

$$
\begin{equation*}
\cdot(A * X B+S)^{*}+C^{*} Q C=D R(X) \tag{1.2}
\end{equation*}
$$

where $X, A, E \in \mathbb{S}^{n, n}, B, S \in \mathbb{N}^{n, m}, C \in \mathbb{C}^{r \cdot n}, Q=Q^{*} \in \mathbb{S}^{r . r}, R=$ $R^{*} \in G^{m, m}$ are positive definite and $E$ is nonsingular. Both equations arise. for example. in the solutions of linear quadratic optimal control problems. Equation (1.1) stems from a continuous-time and (1.2) from a discrete-time problem (see, e.g.. [1], [4]. and [5]). The numerical solution of these two equations has been studied extensively in recent years (e.g., [7], [4]-[6], [9], [14], [16], [17], and [20]). Typically. solutions are obtained using $Q R$ - or $Q Z$-type algorithms to compute deflating subspaces of the matrix pencils

$$
\left[\begin{array}{ccc}
A & 0 & B  \tag{1.3}\\
C^{*} Q C & A^{*} & S \\
S^{*} & B^{*} & R
\end{array}\right]-\lambda\left[\begin{array}{ccc}
E & 0 & 0 \\
0 & -E^{*} & 0 \\
0 & 0 & 0
\end{array}\right]
$$

corresponding to problem (1.1) and

$$
\left[\begin{array}{ccc}
A & 0 & B  \tag{1.4}\\
C^{*} Q C & -E^{*} & S \\
S^{*} & 0 & R
\end{array}\right]-\lambda\left[\begin{array}{ccc}
E & 0 & 0 \\
0 & -A^{*} & 0 \\
0 & -B^{*} & 0
\end{array}\right]
$$

corresponding to problem (1.2). where. using the fact that $R$ is positive definite, many of the algorithms are applied to the reduced pencils

$$
\begin{align*}
& {\left[\begin{array}{cc}
F & G \\
H & -F^{*}
\end{array}\right]-\lambda\left[\begin{array}{cc}
E & 0 \\
0 & E^{*}
\end{array}\right]} \\
& \quad=\left[\begin{array}{cc}
A-B R^{\prime} S^{*} & B R^{-1} B^{*} \\
C^{*} Q C^{-}-S R^{\prime} S^{*} & -A^{*}+S R^{\prime} B^{*}
\end{array}\right]-\lambda\left[\begin{array}{cc}
E & 0 \\
0 & E^{*}
\end{array}\right] \tag{1.5}
\end{align*}
$$

corresponding to ( 1.3 ) and

$$
\begin{align*}
& {\left[\begin{array}{cc}
F & 0 \\
H & E^{*}
\end{array}\right]-\lambda\left[\begin{array}{cc}
E & -C \\
0 & F^{*}
\end{array}\right] } \\
&=\left[\begin{array}{cc}
A-B R & ' S \\
C^{*} Q C-S R^{\prime} S^{*} & E^{*}
\end{array}\right]-\lambda\left[\begin{array}{ccc}
E & -B R^{\prime} B^{*} \\
0 & A^{*}-S R^{\prime} B^{*}
\end{array}\right] \tag{1.6}
\end{align*}
$$

corresponding to (1.4).

[^24]
# Robust, reduced-order, nonstrictly proper state estimation via the optimal projection equations with Petersen-Hollot bounds 

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Received 10 March 1987
Revised 11 June 1987


#### Abstract

A state-estimation design problem involving parametric plant uncertainties is considered. An error bound suggested by recent work of Petersen and Hollot is utilized for guaranteeing robust estimation. Necessary conditions which generalize the optimal projection equations for reduced-order state estimation are used to characterize the estimator which minimizes the error bound. The design equations thus effectively serve as sufficient conditions for synthesizing robust estimators. An additional feature is the presence of a static estimation gain in conjunction with the dynamic (Kalman) estimator, i.e.. a nonstrictly proper estimator.


Keywords: Robust Kalman filter. Error bounds, Reduced-order state estimation.

## 1. Introduction

As is well known [2,5-8,11,14,15] optimal filters based upon nominal parameter values may be severely degraded in the presence of parameter deviations. Thus, it is desirable to obtain robust state estimators which provide acceptable performance over the range of parametric uncertainty. The approach of the present paper is related to the guaranteed cost approach developed for control in [4,16] and applied to estimation in [11]. Specifically, the main idea is to bound the effect of the uncertain parameters on the estimation error over the uncertainty range and then choose estimator gains to minimize the estimation bound. Thus the actual estimation error is guaranteed to lie below the prescribed upper bound.

The technique used to determine minimizing estimator gains is based upon a generalization of the optimal projection equations for reduced-order state estimation [1]. Thus the results of the present paper effectively extend the results of [1] to the case of system uncertainties. It should be noted that the optimal projection equations, which are necessary conditions for optimality, now serve as sufficient conditions for robust estimation by virtue of the fact that a bound on the estimation error is being minimized rather than the estimation error itself. The bound utilized in the present paper is an extension of the approach developed in [12,13] for constructing Lyapunov functions for full-state feedback and utilized in [10] to characterize the structured stability radius.

An additional feature of the present paper is the inclusion of a static feedback gain in conjunction with the dynamic estimator. Thus the results of the present paper represent a generalization of standard results to the case of nonstrictly proper estimation.

[^25]
## 2. Notation and definitions

Note: All matrices have real entries.
$\mathbf{R}, \mathbf{R}^{r \times s}, \mathbb{R}^{r}, \mathbf{E} \quad$ real numbers, $r \times s$ real matrices, $\mathbb{R}^{r \times 1}$, expected value.
$I_{r},()^{\top} \quad r \times r$ identity matrix, transpose
$\mathbb{S}^{r}, \mathbf{N}^{\prime}, \mathbf{P}^{r} \quad r \times r$ symmetric, nonnegative-definite, positive-definite matrices.
$Z_{1} \leq Z_{2}, Z_{1}<Z_{2} \quad Z_{2}-Z_{1} \in \mathbb{N}^{r} . Z_{2}-Z_{1} \in \mathbb{P}^{r}, Z_{1}, Z_{2} \in \mathbb{S}^{r}$.
$n, l, \hat{l}, n_{\mathrm{e}}, p, q ; \tilde{n}$
$x, y, \hat{y}, y_{c}, x_{e}, \tilde{x}$
$A, \Delta A ; C, \Delta C$
positive integers; $n+n_{\mathrm{e}}$.
$n, l, \hat{l}, q, n_{e}, \tilde{n}$-dimensional vectors.
$\hat{C}, L, R$
$A_{\mathrm{e}}, B_{\mathrm{e}}, C_{\mathrm{e}}, D_{\mathrm{e}}$
$\tilde{A} . \Delta \tilde{A}$
$n \times n$ matrices; $l \times n$ matrices.
$\hat{i} \times n$ matrix, $q \times n$ matrix, $q \times q$ matrix, $R \in \mathbb{P}^{q}$.
$n_{e} \times n_{e}, n_{e} \times l, q \times n_{e}, q \times l$ matrices.
$\left[\begin{array}{cc}A & 0 \\ B_{e} C & A_{e}\end{array}\right],\left[\begin{array}{cc}\Delta A & 0 \\ B_{e} \Delta C & 0\end{array}\right]$.
$\tilde{R}$
$w_{0}(\cdot), w_{1}(\cdot)$
$\left[\begin{array}{cc}L^{\top} R L-L^{\top} R D_{e} \hat{C}-\hat{C}^{\top} D_{e}^{\top} R L+\hat{C}^{\top} D_{\mathrm{e}}^{\top} R D_{\mathrm{e}} \hat{C} & -L^{\top} R C_{\mathrm{e}}+\hat{C}^{\top} D_{\mathrm{e}}^{\top} R C_{\mathrm{e}} \\ -C_{e}^{\top} R L+C_{e}^{\top} R D_{\mathrm{e}} \hat{C} & C_{\mathrm{e}}^{\top} R C_{e}\end{array}\right]$.
n. $l$-dimensional white noise.
$V_{0}, V_{1} \quad$ intensity of $w_{0}(\cdot), w_{1}(\cdot) ; V_{0} \in \mathbb{N}^{\prime}, V_{1} \in \mathbb{P}^{\prime}$.
$V_{01} \quad n \times l$ cross intensity of $w_{0}(\cdot), w_{1}(\cdot)$.
$\tilde{w}(\cdot), \tilde{V}$

$$
\left[\begin{array}{c}
w_{0}(\cdot) \\
B_{\mathrm{e}} w_{1}(\cdot)
\end{array}\right] \cdot\left[\begin{array}{cc}
V_{0} & V_{01} B_{\mathrm{e}}^{\mathrm{T}} \\
B_{\mathrm{e}} V_{01}^{\mathrm{T}} & B_{\mathrm{e}} V_{1} B_{\mathrm{e}}^{\mathrm{T}}
\end{array}\right] .
$$

## 3. Robust estimation problem

Let $\mathscr{l} \subset \mathbb{R}^{n \times n} \times \mathbb{R}^{1 \times n}$ denote the set of uncertain perturbations $(\Delta A, \Delta C)$ of the nominal plant matrices $A$ and $C$.

Robust Estimation Problem. For fixed $n_{\mathrm{e}} \leq n$, determine ( $A_{\mathrm{e}}, B_{\mathrm{e}}, C_{\mathrm{e}}, D_{\mathrm{e}}$ ) such that, for the system consisting of the $n$-th-order disturbed plant

$$
\begin{equation*}
\dot{x}(t)=(A+\Delta A) x(t)+w_{0}(t), \quad t \in[0, \infty) \tag{3.1}
\end{equation*}
$$

noisy and nonnoisy measurements

$$
\begin{align*}
& y(t)=(C+\Delta C) x(t)+w_{1}(t)  \tag{3.2}\\
& \hat{y}(t)=\hat{C} x(t) \tag{3.3}
\end{align*}
$$

and $n_{\mathrm{e}}$-th-order nonstrictly proper state estimator

$$
\begin{align*}
& \dot{x}_{\mathrm{e}}(t)=A_{\mathrm{e}} x_{\mathrm{e}}(t)+B_{\mathrm{e}} y(t) .  \tag{3.4}\\
& y_{\mathrm{e}}(t)=C_{\mathrm{e}} x_{\mathrm{e}}(t)+D_{\mathrm{e}} \hat{y}(t), \tag{3.5}
\end{align*}
$$

the state-estimation error criterion

$$
\begin{equation*}
J\left(A_{\mathrm{e}}, B_{\mathrm{e}}, C_{\mathrm{e}}, D_{\mathrm{e}}\right) \triangleq \sup _{(\Delta A, \Delta C) \in \mathbb{Z}_{t}} \lim \sup \mathbb{E}\left[L x(t)-y_{\mathrm{e}}(t)\right]^{\top} R\left[L x(t)-y_{\mathrm{e}}(t)\right] \tag{3.6}
\end{equation*}
$$

is minimized.
Note that the augmented systen (3.1)-(3.5) can be written as

$$
\begin{equation*}
\dot{\tilde{x}}(t)=(\tilde{A}+\Delta \tilde{A}) \tilde{x}(t)+\tilde{w}(t), \quad t \in[0, \infty) \tag{3.7}
\end{equation*}
$$

where $\left.\tilde{x}(t) \triangleq{ }^{\triangleq} x^{\top}(t), x_{\mathrm{e}}^{\top}(t)\right]^{\top}$. The cost can be expressed in terms of the augmented second-moment matrix.

Proposition 3.1. For given $\left(A_{e}, B_{e}, C_{e}, D_{e}\right)$ and $(\Delta A, \Delta C) \in \mathscr{W}$ the second-moment matrix

$$
\begin{equation*}
\tilde{Q}_{\Delta i}(t) \triangleq \mathbf{E}\left[\tilde{x}(t) \tilde{x}^{\mathrm{T}}(t)\right], \quad t \in[0, \infty) \tag{3.8}
\end{equation*}
$$

satisfies

$$
\begin{equation*}
\dot{\tilde{Q}}_{\Delta i}(t)=(\tilde{A}+\Delta \tilde{A}) \tilde{Q}_{د \dot{A}}(t)+\tilde{Q}_{د \dot{i}}(t)(\tilde{A}+\Delta \tilde{A})^{\mathrm{T}}+\tilde{V}, \quad t \in[0, \infty) \tag{3.9}
\end{equation*}
$$

Furthermore,

$$
\begin{equation*}
J\left(A_{\mathrm{e}}, B_{\mathrm{e}}, C_{\mathrm{e}}, D_{\mathrm{e}}\right)=\sup _{(\Delta A, \Delta C) \in \mathscr{Q}} \limsup _{t \rightarrow \infty} \operatorname{tr} \tilde{Q}_{د \tilde{A}}(t) \tilde{R} \tag{3.10}
\end{equation*}
$$

## 4. Sufficient conditions for robust performance

Lemma 4.1. Suppose the system (3.7) is stable for all $(\Delta A, \Delta C) \in \mathscr{U}$. Then

$$
\begin{equation*}
J\left(A_{\mathrm{e}}, B_{\mathrm{e}}, C_{\mathrm{e}}, D_{\mathrm{e}}\right)=\sup _{(\Delta A, \Delta C) \in Q^{2}} \operatorname{tr} \tilde{Q}_{\Delta \bar{A}} \tilde{R} \tag{4.1}
\end{equation*}
$$

where $\tilde{Q}_{\Delta \bar{A}} \in \mathbf{N}^{\tilde{n}}$ is the unique solution to

$$
\begin{equation*}
0=(\tilde{A}+\Delta \tilde{A}) \tilde{Q}_{د \tilde{A}}+\tilde{Q}_{د \hat{A}}(\tilde{A}+\Delta \tilde{A})^{\mathrm{T}}+\tilde{V} \tag{4.2}
\end{equation*}
$$

We now seek upper bounds for $J\left(A_{\mathrm{e}}, B_{\mathrm{e}}, C_{\mathrm{e}}, D_{\mathrm{e}}\right)$.
Theorem 4.1. Let $\Omega: \mathbf{N}^{\tilde{n}} \times \mathbf{R}^{n_{e} \times l} \rightarrow \mathbb{S}^{\boldsymbol{n}}$ be such that

$$
\begin{equation*}
\Delta \tilde{A} \mathscr{Q}+\mathscr{2} \Delta \tilde{A}^{\top} \leq \Omega\left(\mathscr{Q}, B_{\mathrm{e}}\right), \quad(\Delta A, \Delta C) \in \mathscr{G}, \quad\left(2, B_{\mathrm{e}}\right) \in \mathbb{N}^{\tilde{n}} \times \mathbb{R}^{n_{\mathrm{c}} \times I} \tag{4.3}
\end{equation*}
$$

and, for given $\left(A_{\mathrm{e}}, B_{\mathrm{e}}, C_{\mathrm{e}}, D_{\mathrm{e}}\right)$, suppose there exists $\mathscr{Q} \in \mathbf{N}^{\boldsymbol{n}}$ satisfying

$$
\begin{equation*}
0=\tilde{A} \mathscr{Q}+\mathscr{Q} \bar{A}^{\mathrm{T}}+\Omega\left(\mathscr{Q}, B_{e}\right)+\tilde{V} \tag{4.4}
\end{equation*}
$$

and suppose the pair $\left(\tilde{V}^{1 / 2}, \tilde{A}+\Delta \tilde{A}\right)$ is detectable for all $(\Delta A, \Delta C) \in \mathscr{U}$. Then $A_{\mathrm{e}}$ is asymptotically stable. $A+\Delta A$ is asymptotically stable for all $(\triangle A, \Delta C) \in \mathscr{H}$,

$$
\begin{equation*}
\bar{Q}_{\Delta i} \leq \mathscr{Q}, \quad(\Delta A . \Delta C) \in \mathscr{U} \tag{4.5}
\end{equation*}
$$

where $\tilde{Q}_{\mathcal{J}_{i}}$ satisfies (4.2), and

$$
\begin{equation*}
J\left(A_{\mathrm{e}}, B_{\mathrm{e}}, C_{e}, D_{\mathrm{e}}\right) \leq \operatorname{tr} \mathscr{Q} \tilde{R} \tag{4.6}
\end{equation*}
$$

Proof. For all $(\Delta A, \Delta C) \in \mathscr{U},(4.4)$ is equivalent to

$$
\begin{equation*}
0=(\tilde{A}+\Delta \tilde{A}) \mathscr{Q}+\mathscr{Q}(\tilde{A}+\Delta \tilde{A})^{\mathrm{T}}+\Psi\left(\mathscr{Q}, B_{c}, \Delta \tilde{A}\right)+\tilde{V} \tag{4.7}
\end{equation*}
$$

where

$$
\Psi\left(\mathscr{Q}, B_{\mathrm{e}}, \Delta \tilde{A}\right) \triangleq \Omega\left(\mathscr{Q}, B_{\mathrm{e}}\right)-\left(\tilde{\Delta} A \mathscr{Q}+\mathscr{Q} \Delta \tilde{A}^{\top}\right)
$$

Note that by (4.3), $\Psi\left(\mathscr{Q}, B_{\mathrm{e}}, \Delta \tilde{A}\right) \geq 0$ for all $(\Delta A, \Delta C) \in \mathscr{Q}$. Since $\left(\tilde{V}^{1,2}, \tilde{A}+J \tilde{A}\right)$ is detectable for all $(\Delta A, \Delta C) \in \mathscr{U}$, it follows from Theorem 3.6 of [17] that $\left(\left(\tilde{V}+\Psi\left(2, B_{c}, \Delta \bar{A}\right)\right)^{1,2}, \tilde{A}+\Delta \tilde{A}\right)$ is detectable for all $(\Delta A, \Delta C) \in \mathscr{Q}$. Hence Lemma 12.2 of [17] implies $\tilde{A}+\Delta \tilde{A}$ is asymptotically stable for all $(\Delta A, \Delta C) \in \mathscr{Q}$. Since $\tilde{A}+\Delta \tilde{A}$ is lower block triangular, $A_{c}$ is asymptotically stable and $A+J A$ is asymptotically stable for all $(\Delta A, \Delta C) \in \mathscr{Z}$. Next, substracting (4.2) from (4.7) yields

$$
0=(\tilde{A}+\Delta \tilde{A})\left(\mathscr{Q}-\bar{Q}_{د i}\right)+\left(\mathscr{Q}-\bar{Q}_{د \dot{A}}\right)(\tilde{A}+\Delta \tilde{A})^{\mathrm{T}}+\Psi\left(\mathscr{2} \cdot B_{\mathrm{c}} \cdot \Delta \tilde{A}\right)
$$

or, equivalently, (since $\tilde{A}+\Delta \tilde{A}$ is asymptotically stable)

$$
\mathscr{Q}-\tilde{Q}_{\Delta \bar{A}}=\int_{0}^{\infty} \mathrm{e}^{(\tilde{i}+\Delta \tilde{i})} \Psi\left(\mathscr{Q}, B_{\mathrm{e}}, \Delta \tilde{A}\right) \mathrm{e}^{(\tilde{A}+\Delta \tilde{A})^{\mathrm{T}} t} \mathrm{~d} t \geq 0
$$

which implies (4.5). Finally, (4.5) and (4.1) yield (4.6).

## 5. Uncertainty structure

The uncertainty set $\mathscr{l}$ is assumed to be of the form

$$
\begin{align*}
& \mathscr{Q}=\left\{(\Delta A, \Delta C) \in \mathbf{R}^{n \times n} \times \mathbf{R}^{1 \times n}: \Delta A=\sum_{i=1}^{p} D_{i} M_{i} N_{i} E_{i},\right. \\
&  \tag{5.1}\\
& \left.\quad \Delta C=\sum_{i=1}^{p} F_{i} M_{i} N_{i} E_{i}, M_{i} M_{i}^{\mathrm{\top}} \leq \bar{M}_{i}, N_{i}^{\top} N_{i} \leq \bar{N}_{i}, i=1, \ldots, p\right\},
\end{align*}
$$

where, for $i=1, \ldots, p, D_{i} \in \mathbf{R}^{n \times r_{i}}, E_{i} \in \mathbb{R}^{l_{1} \times n}$ and $F_{i} \in \mathbb{R}^{\prime \times n_{1}}$ are fixed matrices denoting the structure of the uncertainty; $\bar{M}_{i} \in \mathbf{N}^{r_{i}}$ and $\bar{N}_{i} \in \mathbf{N}^{t_{r}}$ are given uncertainty bounds; and $M_{i} \in \mathbb{R}^{r_{1} \times s_{i},} N_{i} \in \mathbb{R}^{s_{1} \times t_{i}}$ are uncertain matrices. The closed-loop system thus has structured uncertainty of the form

$$
\begin{equation*}
\Delta \tilde{A}=\sum_{i=1}^{p} \tilde{D}_{i} M_{i} N_{i} \tilde{E}_{i} \tag{5.2}
\end{equation*}
$$

where

$$
\tilde{D}_{i} \triangleq\left[\begin{array}{c}
D_{1}  \tag{5.3}\\
B_{c} F_{i}
\end{array}\right], \quad \tilde{E}_{i} \triangleq\left[\begin{array}{ll}
E_{1} & 0
\end{array}\right] .
$$

The special case $\bar{M}_{i}=\mu_{i}^{2} I_{r}, \bar{N}_{i}=\nu_{i}^{2} I_{u}$ is worth noting.
Proposition 5.1. Let $\mu_{i}, \nu_{i} \geq 0, i=1, \ldots, p$. Then $M_{i} M_{i}^{\top} \leq \mu_{i}^{2} I_{r_{i}}$ and $N_{i}^{\top} N_{t} \leq \nu_{i}^{2} I_{t}$ if and only if $\sigma_{\max }\left(M_{i}\right) \leq \mu_{i}$ and $\sigma_{\max }\left(N_{i}\right) \leq \nu_{i}$.

Remark 5.1. $\mathbb{Z}$ given by (5.1) is directly related to the structured stability radius introduced in [10]. Setting $p=1, \bar{M}_{1}=\mu_{1} I_{r}, r_{1}=s_{1}, \bar{N}_{1}=I_{1}$, and $N_{1} \equiv I_{r}$, yields the setting of [10]. For a similar formulation, see [13].

## 6. The Petersen-Hollot bound

Given $\mathscr{Z}$ as defined in (5.1), we now specify $\Omega$ satisfying (4.3).
Proposition 6.1. The bound $\Omega$ given by

$$
\begin{equation*}
\Omega\left(\mathscr{Q}, B_{e}\right) \triangleq \sum_{i=1}^{p} \bar{D}_{i} \bar{M}_{i} \bar{D}_{i}^{\top}+\mathscr{Q} \bar{E}_{i}^{\top} \bar{N}_{i} \bar{E}_{\mathrm{i}} \mathscr{Q} \tag{6.1}
\end{equation*}
$$

satisfies (4.3) with ${ }^{2} /$ given by (5.1).

Proof. For $i=1, \ldots, p$,

$$
\begin{aligned}
0 & \leq\left[\tilde{D}_{i} M_{i}-\mathscr{Q} \tilde{E}_{i}^{\mathrm{T}} N_{i}^{\mathrm{T}}\right]\left[\tilde{D}_{i} M_{i}-\mathscr{Q} \tilde{E}_{i}^{\mathrm{\top}} N_{i}^{\mathrm{\top}}\right]^{\mathrm{\top}} \\
& =\tilde{D}_{i} M_{i} M_{i}^{\mathrm{\top}} \tilde{D}_{i}^{\mathrm{T}}+\mathscr{Q} \tilde{E}_{i}^{\mathrm{T}} N_{i}^{\mathrm{T}} N_{i} \tilde{E}_{i} \mathscr{Q}-\left(\tilde{D}_{i} M_{i} N_{i} \tilde{E}_{i} \mathscr{Q}+\mathscr{Q} \tilde{E}_{i}^{\mathrm{\top}} N_{i}^{\mathrm{\top}} M_{i}^{\top} \tilde{D}_{i}^{\mathrm{T}}\right) \\
& \leq \tilde{D}_{i} \bar{M}_{i} \tilde{D}_{i}^{\mathrm{T}}+\mathscr{Q} \tilde{E}_{i}^{\mathrm{\top}} \bar{N}_{i} \tilde{E}_{i} \mathscr{Q}-\left(\tilde{D}_{i} M_{i} N_{i} \tilde{E}_{i} \mathscr{Q}+\mathscr{Q} \tilde{E}_{i}^{\mathrm{\top}} N_{i}^{\top} M_{i}^{\mathrm{T}} \tilde{D}_{i}^{\mathrm{\top}}\right) .
\end{aligned}
$$

Summing over $i$ yields (4.3).
Remark 6.1. The bound (6.1) is used in [12] for unit-rank perturbations while a more general treatment appears in [13].

## 7. The auxiliary minimization problem

Our goal is to minimize the error bound (4.6).
Auxiliary Minimization Problem. Determine (2, $A_{\mathfrak{e}}, B_{e}, C_{e}, D_{e}$ ) with $\mathscr{2} \in \mathbb{N}^{\tilde{n}}$ which minimizes

$$
\begin{equation*}
\mathscr{J}\left(\mathscr{Q}, A_{e}, B_{e}, C_{e}, D_{e}\right) \triangleq \operatorname{tr} \mathscr{Q} \tilde{R} \tag{7.1}
\end{equation*}
$$

subject to

$$
\begin{equation*}
0=\tilde{A} \mathscr{Q}+\mathscr{Q} \tilde{A}^{\top}+\sum_{i=1}^{p}\left[\tilde{D}_{i} \bar{M}_{i} \tilde{D}_{i}^{\top}+\mathscr{Q} \tilde{E}_{i}^{\top} \bar{N}_{i} \tilde{E}_{i} \mathscr{Q}\right]+\tilde{V} \tag{7.2}
\end{equation*}
$$

and

$$
\begin{equation*}
\left(\tilde{V}^{1 / 2}, \tilde{A}+\Delta \tilde{A}\right) \text { is detectable, } \quad(\Delta A, \Delta C) \in \mathscr{U} \tag{7.3}
\end{equation*}
$$

Proposition 7.1. If (2, $A_{e}, B_{e}, C_{e}, D_{e}$ ) satisfies (7.2) and (7.3) with $\mathscr{2} \geq 0$, then $\tilde{A}+\Delta \tilde{A}$ is asymptotically stable for all $(\Delta A, \Delta C) \in \mathscr{Z}$, and

$$
\begin{equation*}
J\left(A_{e}, B_{e}, C_{e}, D_{e}\right) \leq \mathscr{J}\left(2, A_{e}, B_{e}, C_{e}, D_{e}\right) \tag{7.4}
\end{equation*}
$$

Proof. With $\Omega$ given by (6.1), (7.2) is equivalent to (4.4). Hence. with (7.3), the hypotheses of Theorem 4.1 are satisfied so that the augmented system is stable over $\mathscr{G}$ with estimation bound (4.6). Note that with (7.1), (7.4) is merely a restatement of (4.6).

## 8. Necessary conditions for the auxiliary minimization problem

Rigorous application of the Lagrange multiplier technique requires additional technical assumptions. Specifically, we further restrict ( $2, A_{\mathrm{e}}, B_{\mathrm{e}}, C_{\mathrm{e}}, D_{e}$ ) to the set

$$
\begin{aligned}
& \mathscr{S} \triangleq\left\{\left(\mathscr{Q}, A_{\mathrm{e}}, B_{\mathrm{e}}, C_{\mathrm{e}}, D_{\mathrm{e}}\right): \mathscr{Q} \in \mathbb{P}^{\boldsymbol{N}}, \tilde{\mathscr{A}} \text { is asymptotically stable. }\left(A_{\mathrm{e}}, B_{\mathrm{e}}, C_{\mathrm{e}}\right)\right. \text { is controllable } \\
& \\
& \left.\quad \text { and observable, and } \hat{C}\left(Q_{1}-Q_{12} Q_{2}^{-1} Q_{12}^{\mathrm{T}}\right) \hat{C}^{\mathrm{T}}>0\right\}
\end{aligned}
$$

where ( $\oplus$ denotes Kronecker sum [3])

$$
\tilde{\mathscr{A}} \triangleq\left(\tilde{A}+\sum_{i=1}^{p} \tilde{E}_{i}^{\mathrm{T}} \bar{N}_{i} \tilde{E}_{i} \mathscr{Q}\right) \oplus\left(\tilde{A}+\sum_{i=1}^{p} \tilde{E}_{i}^{\mathrm{\top}} \bar{N}_{i} \tilde{E}_{i} \mathscr{Q}\right)
$$

and $\mathscr{2}$ is partitioned as in Appendix A. As shown in Appendix A, $Q_{2}$ is invertible since $\left(A_{e}, B_{e}\right)$ is controllable. The positive definiteness condition holds when $\hat{C}$ has full row rank and $\mathscr{Q}$ is positive definite.

As can be seen from the proof of Theorem 8.1 in Appendix $A$, this condition implies the existence of the projection $\tau_{1}$ defined below. Note that $\mathscr{P}$ is open.

Remark 8.1. The constraint ( $\mathscr{2}, A_{e}, B_{e}, C_{e}, D_{e}$ ) $\in \mathscr{S}$ is not required for robust estimation. As will be seen from the proof of Theorem 8.1, the set $\mathscr{P}$ constitutes sufficient conditions under which the Lagrange multiplier technique is applicable to the Auxiliary Minimization Problem. Specifically, asymptotic stability of $\mathscr{\mathscr { A }}$ serves as a normality condition which further implies that the dual $\mathscr{P}$ of $\mathscr{Q}$ satisfying (A.2) is nonnegative definite. Furthermore, ( $A_{e}, B_{e}, C_{e}$ ) minimal is a nondegeneracy condition which implies that the lower right $n_{e} \times n_{\mathrm{e}}$ subblocks of $\mathscr{2}$ and $\mathscr{P}$ are positive definite. It is extremely important to emphasize that Proposition 7.1 shows that it is not necessary for guaranteed robust estimation that an admissible quadruple obtained by solving the necessary conditions actually be shown to be an element of $\mathscr{P}$.

For arbitrary $Q \in \boldsymbol{R}^{n \times n}$ define the following notation:

$$
\begin{aligned}
& V_{1 \mathrm{a}} \triangleq V_{1}+\sum_{i=1}^{p} F_{1} \bar{M}_{i} F_{i}^{\mathrm{T}}, \quad Q_{\mathrm{a}} \triangleq V_{01}+Q C^{\mathrm{T}}+\sum_{i=1}^{\rho} D_{i} \bar{M}_{i} F_{1}^{\mathrm{T}} . \\
& D \triangleq \sum_{i=1}^{p} D_{i} \bar{M}_{i} D_{i}^{\mathrm{T}}, \quad E \triangleq \sum_{i=1}^{p} E_{i}^{\top} \bar{N}_{i} E_{i}, \quad A_{Q} \triangleq A-Q_{\mathrm{a}} V_{1 \mathrm{a}}^{-1} C .
\end{aligned}
$$

The following factorization lemma is needed for the statement of the main result. See [1] for details.
Lemma 8.1. If $\hat{Q}, \hat{P} \in \mathbb{N}^{n}$ and rank $\hat{Q} \hat{P}=n_{\mathrm{e}}$, then there exist $n_{\mathrm{e}} \times n G, \Gamma$ and $n_{\mathrm{e}} \times n_{\mathrm{e}}$ invertible $M$ such that

$$
\begin{equation*}
\hat{Q} \hat{P}=G^{\mathrm{T}} M \Gamma, \quad \Gamma G^{\mathrm{T}}=I_{n_{e}} . \tag{8.1}
\end{equation*}
$$

Furthermore, $G, M$ and $\Gamma$ are unique except for a change of basis in $\mathbb{R}^{n_{c}}$.
Since $\hat{Q} \hat{P}$ is diagonalizable it has a group generalized inverse $(\hat{Q} \hat{P})^{*}=G^{\top} M^{-1} \Gamma$ and

$$
\begin{equation*}
\tau \hat{Q} \hat{Q} \hat{P}(\hat{Q} \hat{P})^{\#}=G^{\mathbf{T}} \Gamma \tag{8.3}
\end{equation*}
$$

is an oblique projection. Define the complementary projection $\tau_{\downarrow} \triangleq I_{n}-\tau$ and call ( $G, M, \Gamma$ ) satisfying (8.1) and (8.2) a projective factorization of $\hat{Q} \hat{P}$.

Theorem 8.1. (2, $\left.A_{e}, B_{e}, C_{e}, D_{e}\right) \in \mathscr{S}$ is an extremal of the Auxiliary Minimization problem with $\mathscr{U}$ given by (5.1) if and only if there exist $Q, \hat{Q}, \hat{P} \in \mathbb{N}^{n}$ such that $2, A_{\mathrm{e}}, B_{\mathrm{e}}, C_{\mathrm{e}}, D_{\mathrm{e}}$ are gicen by

$$
\begin{align*}
& \mathscr{Q}=\left[\begin{array}{cc}
Q+\hat{Q} & \hat{Q} \Gamma^{\mathrm{T}} \\
\Gamma \hat{Q} & \Gamma \hat{Q} \Gamma^{\mathrm{T}}
\end{array}\right]  \tag{8.4}\\
& A_{\mathrm{e}}=\Gamma\left(A-Q_{\mathrm{a}} V_{\mathrm{la}}^{-1} C+Q E\right) G^{\mathrm{T}}  \tag{8.5}\\
& B_{\mathrm{e}}=\Gamma Q_{\mathrm{a}} V_{\mathrm{Ia}^{-1}}  \tag{8.6}\\
& C_{\mathrm{e}}=L \tau_{1 \perp} G^{\mathrm{T}}  \tag{8.7}\\
& D_{\mathrm{e}}=L Q \hat{C}^{\mathrm{T}}\left(\hat{C} Q \hat{C}^{\mathrm{T}}\right)^{-1} \tag{8.8}
\end{align*}
$$

for some projective factorization $(G, M, \Gamma)$ of $\hat{Q} \hat{P}$, and such that $Q, \hat{Q}, \hat{P}$ satisfy

$$
\begin{align*}
& 0=A Q+Q A^{\top}+V_{0}+D+Q E Q-Q_{\mathrm{a}} V_{1 \mathrm{a}}^{-1} Q_{\mathrm{a}}^{\top}+\tau_{\perp} Q_{\mathrm{a}} V_{1 \mathrm{a}}^{-1} Q_{\mathrm{a}}^{\top} \tau_{\perp}^{\top}  \tag{8.9}\\
& 0=(A+Q E) \hat{Q}+\hat{Q}(A+Q E)^{\mathrm{T}}+\hat{Q} E \hat{Q}+Q_{\mathrm{a}} V_{1 \mathrm{a}}^{-1} Q_{\mathrm{a}}^{\top}-\tau_{\perp} Q_{\mathrm{a}} V_{1 a}^{-1} Q_{\mathrm{a}}^{\top} \tau_{\perp}^{\top}  \tag{8.10}\\
& 0=\left(A_{Q}+Q E\right)^{\mathrm{T}} \hat{P}+\hat{P}\left(A_{Q}+Q E\right)+\tau_{1 \perp}^{\top} L^{\top} R L \tau_{1 \perp}-\tau_{\perp}^{\top} \tau_{1 \perp}^{\top} L^{\top} R L \tau_{1_{\perp}} \tau_{\perp}  \tag{8.11}\\
& \operatorname{rank} \hat{Q}=\operatorname{rank} \hat{P}=\operatorname{rank} \hat{Q} \hat{P}=n_{\mathrm{e}} . \tag{8.12}
\end{align*}
$$

where

$$
\begin{equation*}
\tau_{1} \triangleq Q \hat{C}^{T}\left(\hat{C} Q \hat{C}^{T}\right)^{-1} \hat{C}, \quad \tau_{1 \perp} \triangleq I_{n}-\tau_{1} \tag{8.13}
\end{equation*}
$$

Theorem 8.1 (proved in Appendix A) presents necessary conditions for the Auxiliary Minimization Problem which explicitly characterize extremals ( $2 . A_{e}, B_{e}, C_{e}, D_{e}$ ). These necessary conditions consist of a system of two modified Lyapunov equations and one modified Riccati equation coupled by two oblique projections $\tau$ and $\tau_{1}$ and uncertainty terms. The projections $\tau$ and $\tau_{1}$ correspond to reduced estimator order and singular observation noise, respectively.

Several special cases can immediately be discerned. For example, in the full-order estimator case $n_{e}=n$. set $\tau=I_{n}$ so that $\tau_{\perp}=0$. Now the last term in each of (8.9)-(8.11) can be deleted and $G$ and $\Gamma$ in (8.4)-(8.7) can be taken to be the identity. Furthermore, since $\hat{Q}$ and $\hat{P}$ now play no role in determining the optimal estimator, equations (8.10) and (8.11) are superfluous. If, furthermore. $D_{1} . E_{1}$ and $F_{1}$ are zero, then (8.9) reduces to the standard observer Riccati equation of steady-state Kalman filter theory. Alternatively, the case in which the static estimator gain $D_{e}$ is absent can be handled by ignoring (8.8) and setting $\tau_{1}=0$. If, furthermore, the uncertainty terms are deleted then the results of [1] are recovered.

## 9. Sufficient conditions for robust, reduced-order estimation

The main result guaranteeing robust estimation can now be stated.
Theorem 9.1. Suppose there exist $\mathcal{Q}, \hat{Q}, \hat{P} \in \mathbf{N}^{n}$ satisfying (8.9)-(8.12). let $A_{\mathrm{c}}, B_{\mathrm{e}}, C_{\mathrm{e}}, D_{\mathrm{e}}$ be given by (8.5)-(8.8), and suppose that $\left(\tilde{V}^{1 / 2}, \tilde{A}+\Delta \tilde{A}\right)$ is detestable for all $(\Delta A, \Delta C) \in \mathscr{Q}$ with $\mathscr{Q}$ given by (5.1). Then $A_{\mathrm{e}}$ is asymptotically stable, $A+\Delta A$ is asymptotically stable for all $(\Delta A . \Delta C) \in \mathscr{G}$. and the estimation error satisfies the bound

$$
\begin{equation*}
J\left(A_{\mathrm{e}}, B_{\mathrm{e}}, C_{\mathrm{e}}, D_{\mathrm{e}}\right) \leq \operatorname{tr} Q \tau_{1_{\perp}}^{\top} L^{\mathrm{T}} R L \tau_{1_{\perp}} \tag{9.1}
\end{equation*}
$$

Proof. Theorem 8.1 implies 2 given by (8.4) satisfies (7.2). With the detectability assumption the result follows from Proposition 7.1.

Remark 9.1. Note that if $\hat{C}=L$ then $C_{e}=0$ and the estimation bound (9.1) is zero since $\hat{C} r_{1_{i}}=0$. This is. of course, to be expected since perfect estimation is achievable in this case.

Remark 9.2. The problem of designing reduced-order, robust estimators for unstable systems remains an area for future research.

## Appendix A: Proof of Theorem 8.1

Partition $\tilde{n} \times \bar{n} 2, \mathscr{P}$ into $n \times n, n \times n_{e}$, and $n_{e} \times n_{e}$ subblocks as

$$
\mathscr{Q}=\left[\begin{array}{ll}
Q_{1} & Q_{12} \\
Q_{12}^{T} & Q_{2}
\end{array}\right] . \quad \mathscr{P}=\left[\begin{array}{cc}
P_{1} & P_{12} \\
P_{12}^{T} & P_{2}
\end{array}\right]
$$

and define the $n \times n$ nonnegative-definite matrices

$$
\begin{array}{ll}
Q \triangleq Q_{1}-Q_{12} Q_{2}^{-1} Q_{12}^{\mathrm{T}}, & P \triangleq P_{1}-P_{12} P_{2}^{-1} P_{12}^{\top} \\
\hat{Q} \triangleq Q_{12} Q_{2}^{-1} Q_{12}^{\mathrm{T}} . & \hat{P} \triangleq P_{12} P_{2}^{-1} P_{12}^{\mathrm{T}}
\end{array}
$$

and the $n_{e} \times n, n_{e} \times n_{e}, n_{e} \times n$ matrices

$$
G \triangleq Q_{2}^{-1} Q_{12}^{\top}, \quad M \triangleq Q_{2} P_{2}, \quad \Gamma \triangleq-P_{2}^{-1} P_{12}^{\top}
$$

The existence of $Q_{2}^{-1}$ and $P_{2}^{-1}$ is shown below.

To optimize (7.1) over the open set $\mathscr{S}^{\prime}$, where $\mathscr{S}^{\prime} \triangleq\left\{\left(\mathscr{Q}, A_{\mathrm{e}}, B_{\mathrm{e}}, C_{\mathrm{e}}, D_{\mathrm{e}}\right) \in \mathscr{P}:(7.3)\right.$ is satisfied $\}$. subject to the constraint (7.2), form the Lagrangian

$$
\mathscr{L}\left(\mathscr{Q}, A_{e}, B_{e}, C_{e}, D_{e}\right) \triangleq \operatorname{tr}\left[\lambda \mathscr{Q} \tilde{R}+\left(\tilde{A} \mathscr{Q}+\mathscr{Q} \tilde{A}^{\top}+\sum_{i=1}^{p} \tilde{D}_{i} \bar{M}_{i} \tilde{D}_{t}^{\top}+\mathscr{Q} \tilde{E}_{i}^{\top} \bar{N}_{i} \tilde{E}_{1} \mathscr{Q}+\tilde{V}\right) \mathscr{B}\right]
$$

where the multipliers $\lambda \geq 0$ and $\mathscr{P} \in \mathbb{R}^{\tilde{n} \times \bar{n}}$ are not both zero. We thus obtain

$$
\frac{\partial \mathscr{L}}{\partial \mathscr{Q}}=\tilde{A}^{\top} \mathscr{P}+\mathscr{P} \tilde{A}+\sum_{i=1}^{p} \tilde{E}_{i}^{\top} \bar{N}_{i} \tilde{E}_{i} \mathscr{2} \mathscr{P}+\mathscr{P} \mathscr{Q} \tilde{E}_{i}^{\top} \bar{N}_{i} \tilde{E}_{i}+\lambda \tilde{R}
$$

Setting $\partial \mathscr{L} / \partial \mathscr{Q}=0$ yields ('vec' is defined in [3])

$$
\tilde{\mathscr{A}}^{\top} \text { vec } \mathscr{P}=-\lambda \operatorname{vec} \tilde{R}
$$

Since $\tilde{\mathscr{A}}$ is assumed to be invertible, $\lambda=0$ implies $\mathscr{P}=0$. Hence, without loss of generality, set $\lambda=1$. Since, furthermore, $\check{\mathscr{A}}$ is assumed to be asymptotically stable, $\mathscr{P}$ is nonnegative definite. The stationarity conditions are given by

$$
\begin{align*}
& \frac{\partial \mathscr{L}}{\partial \mathscr{P}}=\tilde{A} \mathscr{Q}+\mathscr{2} \tilde{A}^{\mathrm{\top}}+\sum_{i=1}^{F} \tilde{D}_{i} \bar{M}_{i} \tilde{D}_{i}^{\mathrm{T}}+\mathscr{2} \tilde{E}_{i}^{\top} \bar{N}_{i} \tilde{E}_{i} \mathscr{Q}+\tilde{V}=0 .  \tag{A.1}\\
& \frac{\partial \mathscr{L}}{\partial \mathscr{Q}}=\bar{A}^{\top} \mathscr{P}+\mathscr{P} \tilde{A}+\sum_{i=1}^{\rho} \tilde{E}_{i}^{\mathrm{\top}} \bar{N}_{i} \tilde{E}_{i} \mathscr{2 P}+\mathscr{P} \mathscr{2} \tilde{E}_{i}^{\mathrm{\top}} \bar{N}_{i} \tilde{E}_{i}+\tilde{R}=0 .  \tag{A.2}\\
& \frac{\partial \mathscr{L}}{\partial A_{\mathrm{e}}}=P_{12}^{\top} Q_{12}+P_{2} Q_{2}=0,  \tag{A.3}\\
& \frac{\partial \mathscr{L}}{\partial B_{\mathrm{e}}}=P_{12}^{\mathrm{T}} V_{01}+\left(P_{12}^{\top} Q_{1}+P_{2} Q_{12}^{\top}\right) C^{\mathrm{T}}+P_{2} B_{\mathrm{e}} V_{13}=0 .  \tag{A.4}\\
& \frac{\partial \mathscr{L}}{\partial C_{\mathrm{e}}}=-R L Q_{12}+R D_{\mathrm{e}} C Q_{12}+R C_{\mathrm{e}} Q_{2}=0 .  \tag{A.5}\\
& \frac{\partial \mathscr{L}}{\partial D_{e}}=-R L Q_{1} \hat{C}^{\mathrm{T}}+R D_{\mathrm{e}} \hat{C} Q_{1} \hat{C}^{\mathrm{T}}+R C_{\mathrm{e}} Q_{12}^{\top} \hat{C}^{\mathrm{T}}=0 . \tag{A.6}
\end{align*}
$$

Expanding (A.1) and (A.2) yields

$$
\begin{align*}
& 0=A Q_{1}+Q_{1} A^{\top}+D+Q_{1} E Q_{1}+V_{0},  \tag{A.7}\\
& 0=A Q_{12}+Q_{1} C^{\top} B_{\mathrm{e}}^{\top}+Q_{12} A_{\mathrm{e}}^{\mathrm{T}}+Q_{1} E Q_{12}+V_{01} B_{\mathrm{e}}^{\mathrm{\top}},  \tag{A.8}\\
& 0=B_{\mathrm{e}} C Q_{12}+A_{\mathrm{e}} Q_{2}+Q_{12}^{\top} C^{\top} B_{\mathrm{e}}^{\mathrm{T}}+Q_{2} A_{\mathrm{e}}^{\top}+B_{\mathrm{e}} V_{1 \mathrm{a}} B_{\mathrm{e}}^{\top}+Q_{12}^{\mathrm{T}} E Q_{12},  \tag{A.9}\\
& 0=P_{12} A_{\mathrm{e}}+A^{\top} P_{12}+C^{\mathrm{T}} B_{\mathrm{e}}^{\top} P_{2}+E\left(P_{12}^{\top} Q_{1}+P_{2} Q_{12}^{\mathrm{\top}}\right)^{\top}-L^{\top} R C_{\mathrm{e}}+\hat{C}^{\top} D_{\mathrm{e}}^{\mathrm{T}} R C_{\mathrm{e}} .  \tag{A.10}\\
& 0=P_{2} A_{\mathrm{e}}+A_{\mathrm{e}}^{\mathrm{T}} P_{2}+C_{\mathrm{e}}^{\top} R C_{\mathrm{e}} . \tag{A.11}
\end{align*}
$$

Note that the (1.1) subblock of equation (A.2) characterizing $P_{1}$ has been omitted from the above equations since the estimator gains are independent of $P_{1}$. Writing (A.9) as (see [1.9])

$$
0=\left(A_{\mathrm{e}}+B_{\mathrm{e}} C Q_{12} Q_{2}^{+}\right) Q_{2}+Q_{2}\left(A_{\mathrm{e}}+B_{\mathrm{e}} C Q_{12} Q_{2}^{-}\right)^{\top}+Q_{2}\left(Q_{12} Q_{2}^{+}\right)^{\top} E Q_{12} Q_{2}^{+} Q_{2}+B_{\mathrm{e}} V_{1 \mathrm{i}} B_{e}^{\top}
$$

where $Q_{2}^{+}$is the Moore-Penrose or Drazin generalized inverse of $Q_{2}$, it follows from [17]. Lemmas 2.1 and 12.2, that $Q_{2}$ is positive definite. Similarly. (A.11) implies that $P_{2}$ is positive definite.

Next (8.4), (8.6)-(8.8) follow from the definition of 2. (A.4)-(A.6) by using the identities

$$
Q_{1}=Q+\hat{Q} . \quad P_{1}=P+\hat{P} . \quad Q_{12}=\hat{Q} \Gamma^{\top} . \quad P_{12}=-\hat{P} G^{\top} . \quad Q_{2}=\Gamma \hat{Q} \Gamma^{\top} . \quad P_{2}=G \hat{P} G^{\top}
$$


 Similarly, $\left.\Gamma^{\top} G(A .10) \Gamma-(A 10) \Gamma-(A 10)^{\prime}\right)^{\prime}$ yehle (S 11 ,


## References

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## 1 Nfeodection












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Ow prefereste for the Pcierwes-Hollon bownd over the bound orifthally praponed of Chang and neas is baved upon the fact the the former is

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Although the opimal projectiow equatioms are necessan conditime for optimality, in is imporiam to stres, that in the present paper they are otrained own for the original crow function. but rather for a bound on the cant The necesong conditims for the auxilian minimizalion protlem thus effectively serve as sufficiom conditions for the orizind problem. Hence. even if a munctical solution of the extended opimal projection equations fails in produce tic globally opximal comrolier. robust stability and performance are sill guaranteed for all kocal extremals. Our approach thus seeks to rectify one of the main drawbacks of necessity theon by cuarmerming both robust sumbility and performance. Nevertheless. a numerical algorithm for computing the global opamum is given in $[15 \mathbf{}$.

In summary. the moin contribution of the present paper is the seneralization of the optimal projection equations by means of the Petersen-Hollot quadratic Lyapunov bound to synthesize robustly stabilizing fured-order dynamic compensalors with guaranteed performence bound. It is imeresting to note that even in the full-order case. our results. which epecialize to a coupled syssem of three matrix equations. are distinct from the results of $[5]$ which involve a pair of modified Riccati equations and an auxiliary inequality. Furchermore, the
present paper provides a robust performance bound not obtained in $|\$|$ [7]. An additional, conceptual benefit of our approach is a rigorous optimization interpretation for the Petersen-Hollor Riccati equation approech. Finally, as shown in [20] for full-state feedback, the results given herein can be directly applied to the $H_{a}$ design problem. For details, see [21].

Due to space constraims, the contents of the paper will not be reviewed here. We nove only that the proof of Theorem 8.1. which has been omitred for this reason. can be found in [13], (14]. Finally, although numerical algorithms are outside the scope of this note. related results can be found in (15).

## iI. Notation and Definitions

Note: All matrices have real entries.

| H, $8^{\prime \prime} \cdot{ }^{\prime \prime}$, E | Real numbers. $r \times s$ real matrices. $f^{\prime \prime \prime}$. expected value. |
| :---: | :---: |
| J, (1) ${ }^{\text {r }}$ | $r \times r$ identity matrix, transpose. |
| E', M', | $r \times r$ symmetric, nonnegative-definite. positive-definite matrices. |
| $Z_{1} \leq Z_{2}, Z_{1}<Z_{2}$ | $\begin{aligned} & Z_{2}-Z_{1} \in M^{\prime} Z_{2}-Z_{1} \in \mu^{\prime}, Z_{1}, Z_{2} \\ & \Xi^{\prime} . \end{aligned}$ |
| n, m, l, $n_{c} ; \boldsymbol{n}$ | Positive integers: $n+n_{6}$. |
| $x, u, y, x_{c}, x^{2}$ | $n, m, t, n_{c}, i$-dimensional vectors. |
| A. $\triangle A ; B, \Delta B ; C, \Delta C$ | $n \times n$ matrices; $n \times m$ matrices: $I \times n$ matrices. |
| $A_{s}, B_{r}, C_{r}$ | $n_{r} \times n_{r} ; n_{1} \times 1 ; m \times n_{c}$ matrices. |
| $\overline{\mathcal{A}}, \Delta \boldsymbol{A}$ | $\left[\begin{array}{cc}A & B C_{r} \\ B_{r} C & A_{r}\end{array}\right] \cdot\left[\begin{array}{cc}\Delta A & \Delta B C_{c} \\ B_{r} \Delta C & 0\end{array}\right]$ |

$$
R_{1}, R_{2}
$$

$$
R_{1:} \quad n \times m \text { cross weighting matrix: } R_{1}-
$$ $R_{12} R_{:}{ }^{1} R_{12}^{T} \geq 0$. $n$, I-dimensional white noise

```
wil}(\cdot),\mp@subsup{w}{2}{\prime}(\cdot
```

$V_{1} \cdot v$
$V_{1:}$
$\dot{n} \cdot(\cdot), \vec{V}$
$\bar{R}$ Intensity of $w_{1}(\cdot), w_{2}(\cdot) ; V_{1} \geq 0, V_{2}>0$. $n \times /$ cross intensity of $w_{1}(\cdot), w_{2}(\cdot)$.
-リ

$$
\left[\begin{array}{c}
w_{1}(\cdot) \\
B_{1} w_{2}(\cdot)
\end{array}\right] \cdot\left[\begin{array}{cc}
v_{1} & V_{12} B_{!}^{\prime} \\
B_{1} v_{12}^{r} & B_{1} v_{2} B_{1}^{r}
\end{array}\right]
$$

$$
\left[\begin{array}{cc}
R_{1} & R_{12} C_{1} \\
C_{1}^{T} R_{1:}^{T} & C_{1}^{T} R_{i} C_{6}
\end{array}\right]
$$

## III. Roblist Stability aind Roblist Performince Probifms

 perturbations ( $د .4, \Delta B, \Delta C$ ) of the neminal plant matrices $A, B$, and $C$.

Rohust Stabilit! Problem: For fixed $n_{1} \leq n_{\text {, determine }}\left(A_{i}, B_{a}, C_{1}\right)$ such that the closed-lowp system consisting of the $n$ th-order controlled plant

$$
\begin{equation*}
x(1)=(A+\Delta A) x(1)+(B+\Delta B) u(1), \quad t \in[0, \infty) \tag{3.1}
\end{equation*}
$$

measurements

$$
\begin{equation*}
y(t)=(C+\Delta C) x(t) \tag{3.2}
\end{equation*}
$$

and $n_{c}$ th-order dynamic compensator

$$
\begin{gather*}
x_{\mathrm{r}}(t)=A_{\mathrm{r}} x_{\mathrm{r}}(t)+B_{\mathrm{c}} y(t)  \tag{3.3}\\
\mu(t)=C_{\mathrm{N}} x_{\mathrm{r}}(t) \tag{3.4}
\end{gather*}
$$

is asymprotically stable for all $(\Delta A, \Delta B, \Delta C) \in \mathcal{U}$.
Robust Performance Problem: For fixed $n_{c} \leq n$, determine ( $A_{c}, B_{1}$, $\left.C_{6}\right)$ such that. for the closed-loop system consisting of the mth-order disturbed plant
$x(1)=(A+\Delta A) x(t)+(B+\Delta B) u(t)+w_{1}(t), \quad t \in(0, \infty), \quad$ (3.5)
noisy measurements

$$
\begin{equation*}
y(f)=(C+\Delta C) x(1)+w_{2}(f) \tag{3.6}
\end{equation*}
$$

and neth-order dynamic compensator (3.3). (3.4). the performance criterion

$$
\begin{align*}
J\left(A_{1}, B_{r}, C_{r}\right) & \sup _{\text {und }} \lim _{1 \rightarrow \infty} \sup \left[x^{r}(1) R_{1} x(t)\right. \\
& \left.+2 x^{r}(t) R_{1} u(t)+w^{r}(t) R_{2} u(1)\right] \tag{3.7}
\end{align*}
$$

is minimized
Remark 3.1: Noxe that (3.7) is precisely the LQG criterion except for the supremum over $U$ for worst-case performance.

For each controller ( $A_{6}, B_{1}, C_{6}$ ) and plant variation ( $\triangle A, \Delta B, \Delta C$ ) $\in$ U. the undisturbed closed-loop system (3.1)-(3.4) is given b!

$$
\begin{equation*}
x^{\prime}(1)=(A+د A) x^{\prime}(1), \quad t \in(0, \infty) . \tag{3.8}
\end{equation*}
$$

while the disturbed closed-loop system (3.3)-(3.6) is

$$
\begin{equation*}
\dot{x}(t)=(\dot{A}+\Delta \dot{A}) \dot{x}(t)+\dot{w}(t) . \quad t \in(0, \infty) . \tag{3.9}
\end{equation*}
$$

where $\vec{x}(t) \dot{\operatorname{l}}\left|x^{\top}(f), x_{f}^{\top}(f)\right|^{\top}$ and $\dot{w}(\cdot)$ is white noise with intensity $\bar{D} \in$ B4.

## IV. Sufficient Condttions for Robust Stability and Performance

In practice, steady-state performance is only of interest when the closed-loop system: (3.8) is stable over $\mathcal{U}$. The following rexult expresses the performance in terms of the steady-state closed-loop second-moment matrix.

Lemma 4.I: Suppose (3.8) is suable for all $(\triangle A, \Delta B, \Delta C) \in \mathcal{U}$. Then
where $Q_{د X} \underline{i} \lim _{\rightarrow-\infty} F\left(\tilde{r}(t) \dot{x}^{\top}(t)\right\} \in j^{n}$ is the unique solution to

$$
\begin{equation*}
0=(\tilde{A}+\Delta \tilde{A}) Q_{د} A+Q_{د A}(\tilde{A}+\Delta \tilde{A})^{\top}+\tilde{V} . \tag{4.2}
\end{equation*}
$$

We now seek upper bounds for $J\left(A_{1}, B_{1}, C_{c}\right)$.
Theorem 4.1: Let $\Omega: H^{n} \times 8^{n^{*}=1} \times H^{m \times n} \rightarrow \Xi^{n}$ be such that
$\Delta \bar{A} Q+Q \Delta A^{\prime} \leq \Omega\left(Q, B_{5}, C_{8}\right)$.

and. for given ( $A_{i}, B_{i}, C_{r}$ ), suppuse there exists $Q \in$ In $^{n}$ satisfying

$$
\begin{equation*}
0=\bar{A} Q+Q \overline{4}^{T}+\Omega\left(Q_{1}, B_{1}, C_{r}\right)+\bar{V} \tag{4.4}
\end{equation*}
$$

and suppose the pair $\left(\overline{D^{1}}: \bar{A}+\Delta \bar{A}\right)$ is stabilizable for all $(\Delta A, \Delta B$. $\Delta C) \in \mathcal{U}$. Then $\bar{A}+\Delta \bar{A}$ is asymptotically stable for all $(\Delta A, \Delta B, \Delta C)$ $\in \mathcal{U}$.

$$
\begin{equation*}
\bar{Q}_{د} \leq \mathbb{Q}, \quad(\Delta A, \Delta B, \Delta C) \in \boldsymbol{U} \tag{4.5}
\end{equation*}
$$

where $\dot{Q}_{14}$ satisfies (4.2). and

$$
\begin{equation*}
J\left(A_{c}, B_{c}, C_{r}\right) \leq \operatorname{tr} Q \tilde{R} \tag{4.6}
\end{equation*}
$$

Proof: For all $(\Delta A, \Delta B, \Delta C) \in \mathcal{U}$. (4.4) is equivalent to

$$
\begin{equation*}
0=(\bar{A}+\Delta \bar{A}) Q+Q(\bar{A}+\Delta \bar{A})^{r}+\Psi\left(Q, B_{c}, C_{c}, \Delta \bar{A}\right)+\bar{V} \tag{4.7}
\end{equation*}
$$

where

$$
\Psi\left(Q, B_{c}, C_{c}, \Delta A\right) \hat{(Q}\left(Q, B_{c}, C_{c}\right)-\left(\Delta \tilde{A} Q+Q \Delta A^{T}\right)
$$

Note that by (4.3), $\Psi\left(Q, B_{r}, C_{r}, \Delta A\right) \geq 0$ for all $(\Delta A, \Delta B, \Delta C) \in \mathcal{Y}$. Since ( $D^{1,2}, \vec{A}+\Delta A$ ) is stabilizable for all $(\Delta A, \Delta B, \Delta C) \in \mathcal{H}$, it follows from (16, Theorem 3.6) that $\left(\left(D+\Psi\left(Q, B_{r}, C_{c}, \Delta \bar{A}\right)\right)^{1 / 2}, A+\right.$ $\Delta A)$ is stabilizable for all $(\Delta A, \Delta B, \Delta C) \in \mathcal{U}$. Hence, (16, Lemma 12.2] implies $\bar{A}+\Delta \bar{A}$ is asymprotically stable for all $(\Delta A, \Delta B, \Delta C) \in$

थ. Next, subtracting (4.2) from (4.7) yields

$$
0=\left(A+\Delta \Lambda \times Q-Q_{\Delta A}\right)+\left(Q-Q_{\Delta \lambda}\right)(A+\Delta \Lambda)^{r}+\Psi\left(Q, B_{r}, C_{c}, \Delta \Lambda\right)
$$

or. equivalently (since $A+\Delta A$ is asymptotically stable),

$$
Q-Q_{\Delta \lambda}=\int_{0}^{\infty} e^{(\lambda+\Delta \lambda r} \Psi\left(Q, B_{c}, C_{r 0} \Delta \Lambda\right) e^{d \lambda+\Delta \lambda)^{r_{1}}} d t \geq 0 .
$$

which implies (4.5). Finally. (4.5) and (4.1) yield (4.6).

## V. Uncertainty Structure

To obtain explicit expressions for $\left(A_{c}, B_{c}, C_{r}\right)$, we require that $\Delta B=$ 0 , $(\Delta A, \Delta B, \Delta C) \in \mathcal{U}$. Hence, for simplicity, we write $(\Delta A, \Delta C) \in$ $\mathcal{U}$. The dual case $\Delta B \neq 0$ and $\Delta C=0$ is treated in Section $X$. Thus, $\mathcal{U}$ is assumed to be of the form
$u=\left\{(\Delta A, \Delta C) \in J^{n \times n} \times \xi^{i \times n}: \Delta A=\sum_{i=1}^{n} D_{1} M, N, E_{1}\right.$,

$$
\begin{equation*}
\left.\Delta C=\sum_{i=1}^{p} F_{i} M_{i} N_{1} E_{1}, \quad M_{1} M_{1}^{T} \leq \bar{M}_{1}, N_{1}^{r} N_{1} \leq N_{1}, i=1, \cdots, p\right\} \tag{5.1}
\end{equation*}
$$

where, for $i=1, \cdots, p ; D_{i} \in\left\|^{n \times p_{i}}, E_{i} \in\right\|^{i_{i} \times n}$, and $F_{i} \in \mathcal{F}^{1 \times r_{i}}$ are fixed matrices denoting the structure of the uncertainty; $\bar{M}_{i} \in$ inf $^{\prime}$ and $\bar{N}_{1}$
 uncerain matrices. The closed-loop system thus has structured uncertainty of the form

$$
\Delta \bar{A}=\sum_{i=1}^{D} D_{i} M_{i} N_{i} E_{i}
$$

where

$$
D_{i} \triangleq\left[\begin{array}{c}
D_{1} \\
B_{c} F_{1}
\end{array}\right], E_{i} \hat{E}\left(\begin{array}{ll}
E_{1} & 0
\end{array}\right]
$$

The special case $\bar{M}_{1}=\mu_{i}^{2} I_{r_{i}}, \bar{N}_{1}=\nu_{i}^{2} I_{t_{i}}$ is worth noting.
Proposition 5.l: Let $\mu_{1}, \nu_{1} \geq 0, i=1, \cdots, p$. Then $M_{1} M_{1}^{T} \leq \mu_{1}^{2} I_{1}$ and $N_{1}^{r} N \leq \nu_{1}^{2} I_{1}$, if and only if $\sigma_{\max }\left(M_{1}\right) \leq \mu_{1}$ and $\sigma_{\max }\left(N_{1}\right) \leq \nu_{1}$.

Remiark 5.1: The form of $\mathcal{U}$ given by ( 5.1 ) is directly related to the structured stability radius introduced by Hinrichsen and Pritchard [17], [18]. Specifically, let $p=1, \bar{M}_{1}=\mu_{1} I_{r_{1}}, r_{1}=s_{1}$, and $N_{1}=\bar{N}_{1}=I_{r_{1}}$.

## VI. The Petersen-Hollot Buund

Given $\mathcal{U}$, we now specify the bound $\Omega$ satisfying (4.3). Note that becausc of $\Delta B=0, \Omega$ is independent of $C_{6}$. Hence, we write $\Omega\left(Q, B_{c}\right)$ for $\Omega\left(Q, B_{c}, C_{r}\right)$.

Proposition 6.1: The function

$$
\begin{equation*}
\Omega\left(Q, B_{r}\right) \doteq \sum_{i=1}^{\rho} \bar{D}, \bar{M}, D_{1}^{T}+Q \bar{E} r \bar{N}_{1}, E, Q \tag{6.1}
\end{equation*}
$$

satisfies (4.3) with $u$ given by (5.1).
Proof: For $i=1, \cdots, p$,

$$
\begin{aligned}
0 & \leq\left(D_{1} M_{1}-Q E_{1}^{r} N_{1}^{T}\right]\left(D_{1} M_{1}-Q E_{1}^{T} N_{1}^{T}\right]^{T} \\
& =D_{1} M_{1} M_{1}^{r} D_{1}^{T}+Q E_{i}^{T} N_{1}^{r} N_{1} E_{1} Q-\left(D_{1} M_{1} N_{1} E_{1} Q+Q E_{1}^{T} N_{i}^{T} M_{1}^{T} D_{1}^{r}\right) \\
& \leq D_{1} \bar{M}_{1} D_{1}^{T}+Q E_{1}^{r} \bar{N}_{1} E_{1} Q-\left(D_{1} M_{1} N_{1} E_{1} Q+Q E_{1}^{T} N_{1}^{T} M_{1}^{r} D_{i}^{T}\right)
\end{aligned}
$$

Summing over $i$ yields (4.3).
Remark 6.1: The bound (6.1) was originally proposed by Petersen in [5] for unit-rank perturbations with scalar uncertain parameters. A more general treatment appears in (7). Note that we absorb the epsilon used in [7] ino $D_{1}$ and $E_{1}$.

## VII. The Auxiliary Minimization problem

To optimize robust performance while guaranteeing robust stability, we consider the following problem.
Auxiliary Minimization Problem: Determine ( $Q_{1}, A_{r}, B_{c}, C_{r}$ ) which minimizes

$$
\begin{equation*}
g\left(Q, A_{c}, B_{c}, C_{c}\right) \triangleq \mathrm{Or} Q R \tag{7.1}
\end{equation*}
$$

subject to

$$
\begin{align*}
& Q \in B^{n} \text {. }  \tag{7.2}\\
& 0=A Q+Q A^{T}+\sum_{i=1}^{n}\left|D_{1} \bar{M}_{1} D_{1}^{r}+Q E_{i}^{r} \bar{N}_{1} E_{1} Q\right|+P_{\text {. }}  \tag{7.3}\\
& \left(D^{\prime}: A+\Delta A\right) \text { is stabilizable, }(\Delta A, \Delta C) \in u \text {. } \tag{7.4}
\end{align*}
$$

Proposition 7.J: If ( $2 . A_{c}, B_{1}, C_{c}$ ) satisfies (7.2)-(7.4), then $A+$ $\Delta \vec{A}$ is asympenically stable for all $(\Delta A, \Delta C) \in \mathcal{U}$ and

$$
\begin{equation*}
J\left(A_{c}, B_{c}, C_{c}\right) \leq J\left(Q, A_{r}, B_{c}, C_{r}\right) \tag{7.5}
\end{equation*}
$$

Proof: With $\Omega$ given by ( 6.1 ), the hypotheses of Theorem 4.1 are satisfied so that robust stability is guaranteed with performance bound (4.6).

## VIII. Necessary Conditions for the Auxillary Minimization PROBLEM

Rigorous derivation of the necessary conditions for the Auxiliary Minimization Problem requires additional technical assumptions. Specifically, in addition to (7.2), we restrict ( $Q, A_{c}, B_{c}, C_{c}$ ) to the open set

and $\left(A_{c}, B_{c}, C_{c}\right)$ is controllable and observable $\}$,
where (see $\mid 19\}$ for the definition of the Kronecker sum)

$$
\widehat{Q} \equiv\left(\bar{A}+\sum_{i=1}^{0} Q E_{i}^{r} \tilde{N} E_{i}\right) \oplus\left(\bar{A}+\sum_{i=1}^{n} Q E_{i}^{r} \bar{N} E_{i}\right)
$$

Furthermore, the constraint (7.4) will not be accounted for explicitly since it can be shown that the compactness of $\mathcal{U}$ implies that the set of $\left(A_{r}, B_{c}\right.$, $C_{r}$ ) satisfying ( 7.4 ) is open.

Remark 8. I: The constraint $\left(Q, A_{c}, B_{c}, C_{c}\right) \in \delta$ is no required for either robust stability or robust performance since Proposition 7.1 shows that only (7.2)-(7.4) are needed. Rather, the set $\delta$ constitutes sufficient conditions under which the Lagrange multiplier technique is applicable to the Auxiliary Minimization Problem. Specifically, the condition $Q>0$ replaces (7.2) by an open set constraint. the stability of $Q$ serves as a normality condition, and $\left(A_{c}, B_{r}, C_{c}\right)$ minimal is a nondegeneracy condition.

For arbitrary $Q, P \in n^{n \cdot n}$ define the following notation:

$$
\begin{aligned}
& D \doteq \sum_{i=1}^{n} D_{i} \tilde{M}_{1} D_{i}, \quad E \doteq \sum_{i=1}^{r} E_{i}^{r} \tilde{N}_{1} E_{i}, \\
& P_{0} \triangleq B^{T} P+R_{12}^{T} \quad Q_{0} \triangleq Q C^{T}+V_{12}+\sum_{i=1}^{f} D_{i} \tilde{M}_{1} F_{1}^{T},
\end{aligned}
$$

The following factorization lemma is needed. For details, wee [8]. Lemma 8.1: If $Q, \mathcal{F} \in$ 明 $^{n}$ and rank $Q F=n_{c}$, then there exist $n_{c} \times n$ $G, \Gamma$. and $n_{c} \times n_{c}$ invertible $M$ such that

$$
\begin{gather*}
Q P=G^{r} M \Gamma  \tag{8.1}\\
\Gamma G^{r}=I_{n_{r}} \tag{8.2}
\end{gather*}
$$

Furthermore, $G, M$, And $\Gamma$ are unique except for a change of basis in $\mathbf{J B}^{n} c$. As shown in [8], the matrix $\tau$ defined by

$$
\begin{equation*}
r \triangleq \phi P\left(Q^{\rho}\right)^{\prime}=G^{\Gamma} \Gamma \tag{8.3}
\end{equation*}
$$

is an oblique projection where ( ) denotes group generalized inverse [8]. For convenience, define the complementary projection $\tau_{\perp} \triangleq I_{n}-9$.

Theorem 8.1: If ( $Q, A_{c}, B_{c}, C_{c}$ ) $\in \delta$ solves the Auxiliary Minimization Problem with $\mathcal{U}$ given by (5.1), then there exist $Q, P, Q, P$ $\in \mathbb{N u c h}$ shat

$$
\begin{gather*}
Q=\left[\begin{array}{cc}
Q+Q & Q \Gamma^{r} \\
\Gamma Q & \Gamma Q \Gamma^{r}
\end{array}\right]  \tag{8.4}\\
A_{c}=\Gamma\left(A-B R_{2}^{-1} P_{c}-Q_{e} V_{2 e}^{-1} C+Q E\right) G^{r}  \tag{8.5}\\
B_{c}=\Gamma Q_{e} V_{2 e}^{-1}  \tag{8.6}\\
C_{c}=-R_{2}^{-1} P_{\bullet} G^{T} \tag{8.7}
\end{gather*}
$$

and such that $Q, P, Q, \mathcal{P}$ satisfy

$$
\begin{align*}
& 0=A Q+Q A^{T}+V_{1}+D+Q E Q-Q_{\bullet} V_{2 a}^{-1} Q_{\theta}^{T}+T_{+} Q_{\Delta} V_{2 e}^{-1} Q_{s}^{T} r_{+}^{T} . \\
& 0=[A+(Q+Q) E]^{T} P+P(A+(Q+Q) E] \\
& +R_{1}-P_{4}^{T} R_{2}^{-1} P_{4}+T_{1}^{T} P_{4}^{T} R_{2}^{-1} P_{\cdot} T_{+}, \tag{8.9}
\end{align*}
$$

$$
\begin{align*}
& \operatorname{rank} \phi=\operatorname{rank} P=\operatorname{rank} \phi \boldsymbol{\phi}=\boldsymbol{n}_{\mathrm{c}} . \tag{8.11}
\end{align*}
$$

Furthermore, the auxiliary cost is given by

$$
\begin{equation*}
\mathfrak{J}\left(Q, A_{c}, B_{c}, C_{c}\right)=\operatorname{tr}\left[(Q+Q) R_{1}+\phi P B R_{2}^{-1} P_{s}-R_{12} R_{2}^{-1} P_{2} \emptyset\right] \tag{8.13}
\end{equation*}
$$

Conversely, if there exist $Q, P, Q, P \in \mathcal{N}^{n}$ satisfying (8.8)-(8.12). then (Q. $A_{c}, B_{c}, C_{c}$ ) given by (8.4)-(8.7) satisfy (7.2) and (7.3) with cost (8.13).

## Proof: See [13]. [14].

Remark 8.2: Theorem 8.1 presents necessary conditions for the Auxiliary Minimization Problem which explicitly characterize exiremal quadruples ( $Q, A_{r}, B_{c}, C_{c}$ ). These necessary conditions consist of a system of two modified Riccati equations and two modified Lyapunnv equations coupled by both the optimal projection $\tau$ and uncertainty terms. Several special cases can immediately be discemed. For example, in the full-order case $n_{r}=n$, set $\tau=I_{n}$ so that $r_{1}=0$. Now the last term in each of (8.8)-(8.11) can be deleted and $G$ and $\Gamma$ in (8.5)-(8.7) can be taken to be the identity. Furthermore, $P$ plays no role so-that ( 8.11 ) is superfluous. Note that in this case. (8.8) is independent of $P$ and $\phi$. Setting further $D_{i}, E_{i}$, and $F_{1}$ to zero, it can be seen that (8.10) and (8.11) drop out, while (8.8) and (8.9) reduce to the standard separated Riccati equations of LQG theory. If, alternatively, the reduced-order constraint is retained, bui the uncertainty terms are deleted. then the results of [8] are recovered.

Remark 8.3: When solving (8.8)-(8.12) numerically, the uncertainty terms can be adjusted to examine tradeoffs between performance and robustpess. Specifically, the bounds $\bar{M}_{1}$ and $\bar{N}_{1}$ and structure matrices $D_{1}$, $E_{1}$, and $F_{1}$ appearing in $Q_{u}, D_{1}, E_{,}$and $V_{2 a}$ can be varied systematically to determine the region of solvability of (8.8)-(8.12).

## IX. SUFFICIENT CONDITIONS FOR ROBUST STABILITY AND Performance

Theorem 9.1: Suppose there exist $Q, P, Q, P \in$ ©n satisfying (8.8)(8.12), and assume that ( $\left.\nabla^{1 / 2}, A+\Delta \mathcal{A}\right)$ is stabilizable for all $(\Delta A, \Delta C)$ $\in \mathcal{U}$ with $A_{c}, B_{r}, C_{s}$ given by (8.5)-(8.7) and $\mathcal{U}$ given by (5.1). Then $A$
$+\Delta \bar{A}$ is asymptotically stable for all $(\Delta A, \Delta C) \in U$ and the closed-loop performance is bounded by (8.13).

Proof: Theorem 8.1 implies that $Q$ given by (8.4) satisfies (7.2) and (7.3). With the stabilizability essumption, the result follows from Proposition 7.1.

## X. The Dual Case

In place of (5.1), assume now that $\Delta C=0,(\Delta A, \Delta B, \Delta C) \in \mathcal{U}$, and define
$u=\left\{(\Delta A, \Delta B) \in \exists^{n \times n} \times B^{n=m}: \Delta A=\sum_{i=1}^{n} D_{i} M_{i} N_{i} E_{i}\right.$,
$\left.\Delta B=\sum_{i=1}^{p} D_{1} M_{i} N, G_{i}, \quad M_{1} M_{i}^{T} \leq \bar{M}_{i}, N_{1}^{r} N_{i} \leq \bar{N}_{1}, i=1, \cdots, p\right\}$,
 fixed matrices denoting the structure of the uncertainty: and $\bar{M}_{i}, \bar{N}_{1}, M_{1}$, and $N$, are as before. For arbitrary $Q, P \in J^{n \times n}$ define the following notation:

$$
\begin{aligned}
& \mathcal{P}_{1} \triangleq B^{r} P+R_{12}^{T}+\sum_{i=1}^{p} G_{1}^{r} \bar{N}_{1} E_{1}, \quad Q_{0} \triangleq Q C^{r}+V_{12} \\
& \hat{A}_{P} \triangleq A-B R_{20}^{-1} P_{a}, \quad \hat{A}_{Q} \triangleq A-Q_{e} V_{2}^{-1} C, \quad R_{2} \triangleq R_{2}+\sum_{i=1}^{p} G_{1}^{r} N_{i} G_{1} .
\end{aligned}
$$

The main result guaranteeing robust stability and performance for the dual problem can now be stated. For details, see [13]. [14].

Theorem 10.1: Suppose there exist $P, Q, P, Q \in \mathbb{N A}^{n}$ satisfying (8.12) and

$$
\begin{equation*}
0=A^{T} P+P A+R_{1}+E+P D P-P_{1}^{T} R_{20}^{-1} P_{4}+\tau_{1}^{r} P_{2}^{T} R_{2 n}^{-1} P_{4} T_{A} \tag{10.2}
\end{equation*}
$$

$0=[A+D(P+P)] Q+Q[A+D(P+P)]^{\gamma}$

$$
\begin{equation*}
+V_{1}-Q_{0} V_{2}^{-1} \phi_{:}^{r}+r_{1} Q_{.} V_{2}^{-1} Q_{!}^{\Gamma_{T}^{T}} \tag{10.3}
\end{equation*}
$$



$$
\begin{equation*}
0=\left(\hat{A}_{p}+D P\right) Q+Q\left(\hat{A}_{f}+D P\right)^{T}+Q_{0} V_{2}^{-1} Q_{t}^{T}-\tau_{-} Q_{a} V_{:}^{-1} Q_{i}^{r_{i}^{T}} \tag{10.4}
\end{equation*}
$$

and assume that $\left(\tilde{R}^{1 / 2}, \tilde{A}+\Delta \tilde{A}\right)$ is detectable for all $(\Delta A, \Delta B) \in \mathcal{U}$ with $A_{r}, B_{c}, C_{r}$ given by

$$
\begin{gather*}
A_{c}=\Gamma\left(A-Q_{0} V_{2}^{-1} C-B R_{2 a}^{1} F_{a}+D P\right) G^{T}  \tag{10.6}\\
B_{c}=\Gamma Q_{a} V_{2}^{-1}  \tag{10.7}\\
C_{c}=-R_{2 a}^{-1} P_{e} G^{\top} \tag{10.8}
\end{gather*}
$$

and $\mathcal{U}$ given by (10.1). Then, with (10.6)-(10.8), $\bar{A}+\Delta \bar{A}$ is asymptotically stable for all $(\Delta A, \Delta B) \in U$ and the performance of the closed-loop system satisfies

$$
\begin{equation*}
J\left(A_{f}, B_{c}, C_{c}\right) \leq t r\left[(P+P) V_{1}+Q_{0} V_{2}^{-1} C Q P-P Q_{e} V_{2}^{-1} V_{12}^{\dagger}\right] \tag{10.9}
\end{equation*}
$$

Remark 10.1: Even in the case $\Delta B=0, \Delta C=0$, the performance bounds (8.13) and (10.9) are generally different.

Remark 10.2: The case in which $\Delta B$ and $\Delta C$ are simultancously nonzero also appears to be tractable and leads to additional terms in the design equations. The bound considered in [11] also permits this case.

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## A Frequency Response-Based Model Order Selection Criterion

## DAVID J. CLOUD and BASIL KOUVARITAKIS

Abstract-The use of weighting sequence models to describe the dynamics of physical systems provides an effective means of iransiating the uncertainty associaled with the model parameter estimates derived from noisy input/output data into corresponding frequency response uncertaial) information. However, an appropriate truncation level must be extablished to accomplish this task. This paper addresses the truacation problem from a frequency response perspective and proposes a new criterion based on frequency response considerations to select the proper truncation.

Manuscript received January 6. 1987, revised May 18. 1987. This work was cupporred by the U.S. Ais Force and the U. K. Science and Engineering Research Council (SERC). The authors are with the Department of Engineering Science. Universiry of Oxford, Oxford OX1 3PJ. England.
IEEE LOf Number 8718562.

## 1. INTRODUCTION

Recent work on robustness analysis methods for multivariable systems has led to the development of structured singular value techniques (e.g., [1], [2]) which provide a means of assessing the impact of system uncertainty on closed-loop stability and performance. A key missing element in these analysis methods is the ability to describe the frequency response uncerainty associated with any given system. When system identification techniques are used to derive the system description, however, it becomes possible to quantify system uncertainty statistically, and recent efforts have demonstrated that uncertainty information on the estimated parameters of the system model can be transformed into corresponding information on the frequency response uncertainty of the system [3]. [4]. For difference equation models, the transformation from the parameter space to the frequency domain is nonlinear, a result which necessitates the use of linear approximations and produces a statistical description of uncertainty that fails to account for the interfrequency dependence of the frequency response estimates. On the other hand. Cloud and Kouvaritakis [4] have shown that these problems can be avoided by the use of weighting sequence models to describe system dynamics.

The uncertainty description developed in [4] assumes that the system can be accurately described by a finite weighting sequence model. an assumption that is valid for all stable systems. This assumption, in turn. implies that the "correct" model order (i.e., truncation level) is known. As a result, the identification process must not only be able to generate appropriate parameter estimates, it must also be able to identify the "correct" level of truncation. In effect, this second requirement is a restatement of the standard model order selection problem, a problem which has been widely investigated in the literature (e.g. [5]). But for the frequency response applications of interest here, appropriate solutions must focus on generating accurate frequency response information. When this perspective is taken, it becomes clear that the standard order selection criteria are not well suited to the task because they focus on generating accurate input/output descriptions rather than accurate frequency response descriptions.

In this paper, a new criterion is derived to identify the "correct" truncation level based on frequency response considerations. The development begins by highlighting a geometric interpretation of the standard " input/output" order selection problem. These geometric results are then transformed into the frequency domain to produce the new "frequency response-based" criterion for rruncation selection, and simulation results are presented to demonstrate its use. Armed with the "correct" truncation level generated by this criterion. it is now possible to implement the techniques described in (4) to produce a valid description of frequency response uncertainty for any given system.

## II. Model. Order Selection: A Geometric Perspective

Consider the discrete-time system with weighting sequence elements $\left\{\theta_{1}, \theta_{2}, \theta_{3}, \cdots\right\}$ whose true response at sample $k$ to the set of inputs $\{u(k$ - 1), $u(k-2), \cdots, u(0)\}$ is given by

$$
\begin{equation*}
y^{\prime}(k)=\sum_{i=1}^{n} \theta_{1} u(k-i)=d_{1}^{\prime} \theta^{\circ} \tag{1}
\end{equation*}
$$

where $\theta^{0}=\left\{\theta_{1} \cdots \theta_{k}\right\}^{\prime}$ and $d_{i}^{\prime}=\{u(k-1) \cdots u(0)]$. The measured output at sample $k$ is then given by

$$
\begin{equation*}
y^{m}(k)=y^{0}(k)+e(k) \tag{2}
\end{equation*}
$$

where $\epsilon(k)$ is assumed to be an element of a white noise sequence with variance $\sigma_{i}^{2}$. For a set of $N$ measurements, we may stack the scalars $y^{m}(k)$ and $\epsilon(k)$ as elements of the vectors $y^{m}$ and $e$, respectively. and may then rewrite (2) as

$$
\begin{equation*}
y^{m}=y^{0}+e=D^{0} y^{0}+e \tag{3}
\end{equation*}
$$

where the rows of $D^{0}$ are given by the vectors $d_{k}^{\prime}$ for $k=1, \cdots, N$.

# Robust Stability and Performance via Fixed-Order Dynamic Compensation 

by

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#### Abstract

Two robust control-design problems are considered. The Robust Stabilization Problem involves deterministically modeled bounded but unknown time-varying parameter variations, while the Robust Performance Problem includes, in addition, a quadratic performance criterion averaged over stochastic disturbances and maximized over the admissible parameter variations. For both problems the design goal is a fixed-order (i.e., reduced- or full-order) dynamic (strictly proper) feedback compensator. A sufficient condition for solving the Robust Stabilization Problem is given by means of a quadratic Lyapunov function parameterized by the compensator gains. For the Robust Performance Problem the Lyapunov function provides an upper bound for the closedloop performance. This leads to consideration of the Auxiliary Minimization Problem: Minimize the performance bound over the class of fixed-order controllers subject to the Lyapunov-function constraint. Necessary conditions for optimality in the auxiliary. problem thus serve as sufficient conditions for robust stability and performance in the original problem. Two particular bounds are considered for constructing the quadratic Lyapunov function. The first corresponds to a right shift/multiplicative white noise model, while the second was suggested by recent work of Petersen and Hollot. The main result is an extended version of the optimal projection equations for fixedorder dynamic compensation whose solutions are guaranteed to provide both robust stability and robust performance.


[^26]
## 1. Introduction

Although considerable effort has been devoted to frequency-domain robust control-design methods ( $[1-10]$ ), there remain open questions concerning stability with respect to real-valued, structured plant parameter variations ([11-13]). Specifically, it is shown in [11-13] that classical gain and phase margin specifications can be satisfied while sensitivity to structured plant parameter variations can be arbitrarily large. From a time-domain point of view, the parametric robustness problem has been widely studied using Lyapunov's second method as the principal technique ([14-28]).

In the present paper we develop an approach to control design which provides sufficient conditions for robust stability and performance over a prescribed range of time-varying structured plant parameter variations by means of a feedback law in the form of a fixed-order (i.e., reduced- or full-order) dynamic (strictly proper) compensator. The approach is based upon the merging of two techniques, namely, the guaranteed cost control approach to robust performance ( $[14,17]$ ) and the optimal projection approach to quadratically optimal fixed-order dynamic compensation ([29,30]). One of our goals is to obtain robust output-feedback compensators rather than full-state-feedback controllers. Also, since we wish to account for real-time computational burden in implementing the controller, we impose a constraint on the dimension (i.e., order) of the dynamic compensator. This approach thus generalizes standard LQG theory which yields full-order output-feedback controllers for systems without parameter uncertainty. We note that our approach is constructive in the sense that upon satisfaction of the sufficient conditions, the feedback gains required for implementing the robust feedback controller are explicitly synthesized. Existential issues are also addressed herein, although to a lesser extent. For further background, see $[29,30]$. For extensions to nonstrictly proper controllers see [31] and for extensions to $\mathrm{H}_{\infty}$ control see [32].

To explain the rationale behind the development we briefly describe the main elements of the approach. The following discussion is intended to be descriptive; precise conditions appear in the main body of the paper.

1. Robust Stability Problem. For a nominal linear time-invariant ( $A, B, C$ ) system we consider deterministically modeled bounded but otherwise unknown Lebesgue measurable time-varying parameter variations of the form

$$
\begin{equation*}
A+\sum_{i=1}^{p} \hat{\sigma}_{i}(t) A_{i}, \quad B+\sum_{i=1}^{p} \hat{\sigma}_{i}(t) B_{i}, \quad C+\sum_{i=1}^{p} \hat{\sigma}_{i}(t) C_{i} . \tag{1.1}
\end{equation*}
$$

The nominal matrices $\mathrm{A}, \mathrm{B}, \mathrm{C}$ and the perturbation matrices $A_{i}, B_{i}, C_{i}$ denoting the structure of
the parametric uncertainty are assumed known, while the time-varying uncertain parameters $\hat{\sigma}_{i}(t)$ are assumed only to satisfy the bounds

$$
\begin{equation*}
\left|\hat{\sigma}_{i}(t)\right| \leq \delta_{i}, \quad i=1, \ldots, p, \quad t \in[0, \infty) \tag{1.2}
\end{equation*}
$$

The form of (1.1) permits an arbitrary number of uncertain parameters with arbitrary linear structure. Although we do not require matching conditions as in [21], the linear structure of (1.1) is more restrictive than the functional form $A(q(t))$ used in [21]. It is this structure which we exploit to obtain sufficiency conditions. Note also that the representation (1.1) is independent of state space basis since replacing $A$ by $S A S^{-1}$ corresponds to replacing $A_{i}$ by $S A_{i} S^{-1}$. As will be seen, our robustness bounds and optimality conditions are also basis independent. Also, scaling téchniques ([6,7]) will not play a role here. Finally, we note that because of the time-varying nature of the uncertain perturbations (1.1) it is virtually impossible to determine the actual stability region of a given design by means of empirical methods.
2. Quadratic Lyapunov Function. As a sufficient condition for characterizing solutions of the Robust Stability Problem we consider a closed-loop quadratic Lyapunov function $V(\tilde{x})=\tilde{x}^{T} p \tilde{x}$, where the matrix $P$ satisfies

$$
\begin{equation*}
0=\tilde{A}^{T} P+P \tilde{A}+\Omega\left(P, B_{c}, C_{c}\right) \tag{1.3}
\end{equation*}
$$

and the function $\Omega$ is a bound satisfying

$$
\begin{equation*}
\sum_{i=1}^{p} \sigma_{i}\left(\tilde{A}_{i}^{T} p+P \tilde{A}_{i}\right)<\Omega\left(P, B_{c}, C_{c}\right) \tag{1.4}
\end{equation*}
$$

over the parameter range

$$
\begin{equation*}
\left|\sigma_{i}\right| \leq \delta_{i}, \quad i=1, \ldots, p \tag{1.5}
\end{equation*}
$$

Note that the constant $\sigma_{i}$ in (1.4) and (1.5) plays the role of $\hat{\sigma}_{i}(t)$, i.e., $t$ is "frozen" in (1.4) and (1.5). In (1.3) and (1.4) $\tilde{A}$ and $\tilde{A}_{i}$ denote the closed-loop dynamics and closed-loop parameter-uncertainty matrices given by

$$
\tilde{A}=\left[\begin{array}{cc}
A & B C_{c}  \tag{1.6}\\
B_{c} C & A_{c}
\end{array}\right], \quad \tilde{A}_{i}=\left[\begin{array}{cc}
A_{i} & B_{i} C_{c} \\
B_{c} C_{i} & 0
\end{array}\right]
$$

Since $\tilde{A}_{i}$ is independent of $A_{c}, \Omega$ depends only on $B_{c}$ and $C_{c}$. As discussed later in this section, (1.4) is automatically satisfied by construction of the function $\Omega$. Furthermore, the existence of a solution $\rho$ to (1.3) need not be verified directly but is rather a result of numerically solving the optimality conditions discussed below.
3. Robust Performance Problem. In addition to the deterministic parameter uncertainty model (1.1), (1.2), the Robust Performance Problem includes stochastic plant disturbances and measurement noise with performance measured by means of the quadratic functional

$$
\begin{equation*}
\tilde{J}(t)=x^{T}(t) R_{1} x(t)+2 x^{T}(t) R_{12} u(t)+u^{T}(t) R_{2} u(t) . \tag{1.7}
\end{equation*}
$$

To obtain a steady-state design problem we 1) average $\tilde{J}(t)$ over the disturbance and measurement noise statistics; 2) pass to the steady-state limit; and 3) maximize over the class of parameter uncertainties. Hence the performance of a given controller ( $A_{c}, B_{c}, C_{c}$ ) is given by

$$
\begin{equation*}
J\left(A_{c}, B_{c}, C_{c}\right)=\sup _{\partial(\cdot)} \limsup _{t \rightarrow \infty} \mathbb{E}[\tilde{J}(t)] . \tag{1.8}
\end{equation*}
$$

The use of "limsup" is a technicality which account for cases in which the steady-state limit may not exist. Note that although (1.8) is an averaging criterion over the disturbances as in LQG theory, it is also a worst-case measure over the uncertain parameters. Thus (1.8) is a hybrid criterion in the sense that is stochastic in the disturbance space (i.e., external uncertainties) and deterministic in the parameter space (i.e., internal uncertainties). By "internal uncertainties" we have in mind quantities such as mass, damping or stiffness, and by "external uncertainties" we are referring to phenomena such as turbulent flow for which only power spectrum statistics may be available. No claim is made, however, with regard to the universal validity of such a mathematical uncertainty model. In particular applications, uncertainty models which are either wholly deterministic or wholly stochastic may be more appropriate. In general, our setting appears to be consistent with the available literature (see [1-28]).
4. Performance Bound. To obtain a tractable design problem the matrix $P$ is used to bound the performance of each controller solving the Robust Stability Problem. Specifically, by assuming in addition to (1.4) that

$$
\begin{equation*}
\sum_{i=1}^{p} \sigma_{i}\left(\tilde{A}_{i}^{T} P+P \tilde{A}_{i}\right)+\tilde{R} \leq \Omega\left(P, B_{c}, C_{c}\right), \tag{1.9}
\end{equation*}
$$

it follows that

$$
\begin{equation*}
J\left(A_{c}, B_{c}, C_{c}\right) \leq \operatorname{tr} P \tilde{V} . \tag{1.10}
\end{equation*}
$$

In (1.9) and (1.10) $\tilde{R}$ and $\tilde{V}$ denote closed-loop weighting and disturbance intensity matrices. The idea of bounding the performance by means of a Lyapunov function is the basis for guaranteed cost control ([14,17]).*

[^27]5. Construction of the Lyapunov Function. Thus far the Lyapunov function has only been abstractly characterized by means of (1.3) and (1.4). To obtain a useful design theory $\Omega$ is now given a concrete form. Specifically, to satisfy (1.9) it is assumed that
\[

$$
\begin{equation*}
\Omega\left(P, B_{c}, C_{c}\right)=\sum_{i=1}^{p} \Lambda_{i}\left(P, B_{c}, C_{c}\right)+\tilde{R}, \tag{1.11}
\end{equation*}
$$

\]

where, for each $i$, the $\Lambda_{i}$ are chosen such that

$$
\begin{equation*}
\sigma_{i}\left(\tilde{A}_{i}^{T} P+P \tilde{A}_{i}\right) \leq \Lambda_{i}\left(P, B_{c}, C_{c}\right), \quad\left|\sigma_{i}\right| \leq \delta_{i} \tag{1.12}
\end{equation*}
$$

Note that (1.12) implies that (1.4) holds with $\Omega$ given by (1.11). Since $\tilde{A}_{i}$ depends upon $B_{c}$ and $C_{c}$, the bound $\Lambda_{i}$ will be constructed to be gain-invariant, that is, so that (1.12) holds for all $B_{c}$ and $C_{c}$. Thus, no difficulty will arise from the fact that the controller gains are yet to be determined by optimality considerations.

It should be noted that the bounding in (1.12) is defined in the sense of the cone of nonnegativedefinite matrices. Since this is only a partial ordering and not a total ordering, a least upper bound (i.e., a "sharpest" bound) does not exist in general and the conservatism of the inequality in (1.12) cannot be quantified by a scalar measure. Hence, $\Lambda_{i}$ satisfying (1.12) is not necessarily unique and two particular choices of $\Lambda_{i}$ are developed in this paper. Since we shall utilize first-order necessary conditions for optimality, we confine our consideration to bounds which are differentiable. The first choice of $\Lambda_{i}$ satisfying (1.12) is given by the linear (in $\mathcal{P}$ ) function

$$
\begin{equation*}
\Lambda_{i}\left(P, B_{c}, C_{c}\right)=\delta_{i}\left(\alpha_{i} P+\alpha_{i}^{-1} \tilde{A}_{i}^{T} P \tilde{A}_{i}\right) \tag{1.13}
\end{equation*}
$$

where $\alpha_{i}$ is an arbitrary positive number. As shown in [33], the bound (1.13) can be viewed as arising from a stochastic optimal control problem with exponentially weighted cost and state-, control- and measurement-dependent white noise. The stochastic multiplicative white noise model serves only as an interpretation, however, and need not be viewed as having physical significance. A similar bound is utilized in [28].

The second choice for $\Lambda_{i}$ satisfying (1.12) is given by the quadratic (in $P$ ) function

$$
\begin{equation*}
\Lambda_{i}\left(\rho, B_{c}, C_{c}\right)=\delta_{i}\left(\tilde{E}_{i}^{T} \tilde{E}_{i}+\mathcal{P} \tilde{D}_{i} \tilde{D}_{i}^{T} \mathcal{P}\right) \tag{1.14}
\end{equation*}
$$

where $\tilde{D}_{i}, \tilde{E}_{i}$ denote an arbitrary factorization of $\tilde{A}_{i}$ of the form

$$
\begin{equation*}
\tilde{A}_{i}=\tilde{D}_{i} \tilde{E}_{i} \tag{1.15}
\end{equation*}
$$

The bound (1.14) was utilized in [26] for full-state feedback with rank-1 uncertainties. Note that by utilizing congruence transformations it can be seen that both bounds (1.13) and (1.14) are basis independent. That is, replacing $\tilde{A}_{i}$ by $\tilde{S} \tilde{A}_{i} \tilde{S}^{-1}$ leads to replacing $\rho$ by $\tilde{S}^{-T} \rho \tilde{S}^{-1}$.
6. The Auxiliary Minimization Problem. The next step in our development for robust performance is the foilowing. Inasmuch as the performance of a robustly stabilizing controller is bounded via (1.10) over the given range of parameter variations, it is desirable to minimize the upper bound

$$
\begin{equation*}
J\left(P, A_{c}, B_{c}, C_{c}\right) \triangleq \operatorname{tr} P \tilde{V} \tag{1.16}
\end{equation*}
$$

subject to the constraint (1.3). This is referred to as the Auxiliary Minimization Problem. For a given choice (1.13) or (1.14) of $\Lambda_{i}$ for each $i$, a solution of the Auxiliary Minimization Problem provides a controller whose steady-state performance is guaranteed to remain below the bound (1.16) over the range of parameter variations, hence guaranteeing robust performance. Since the Auxiliary Minimization Problem is a smooth mathematical programming problem, a minimum always exists on compact sets. To actually characterize extremals of the Auxiliary Minimization Problem we proceed by deriving first-order necessary conditions. Because these necessary conditions are derived for the Auxiliary Minimization Problem, they effectively serve as sufficient conditions for robustness in the original problem.

It should be noted that the guaranteed cost control approach developed in [14] does not permit this line of development since $\Lambda_{i}$ is given by

$$
\begin{equation*}
\Lambda_{i}\left(P, B_{c}, C_{c}\right)=\delta_{i}\left|\tilde{A}_{i}^{T} P+P \tilde{A}_{i}\right| \tag{1.17}
\end{equation*}
$$

where $|\cdot|$ denotes the matrix obtained by replacing each eigenvalue by its absolute value. Since this bound is not differentiable with respect to the controller gains, first-order necessary conditions cannot be used.
7. The Optimality Conditions: Full-Order Case. For the full-order case, i.e., when the order of the controller is equal to the order of the plant, the first-order necessary conditions can be derived in a form which is a direct generalization of the pair of separated Riccati equations of LQG theory. Specifically, the necessary conditions comprise a coupled system of four algebraic matrix equations including a pair of modified Riccati equations and a pair of Lyapunov equations. For plant models involving multiplicative white noise these equations have been studied in [34-36]. This form of the equations thus essentially corresponds to choosing bound (1.13).
8. The Optimality Conditions: Reduced-Order Case. For design flexibility we also consider controllers of arbitrary reduced dimension. For the linear-quadratic problem without parameter uncertainty, the formulation of the necessary conditions given in [29] provides a generalization of LQG theory. Specifically, the optimal gains are characterized by a system of algebraic matrix equations consisting of a pair of modified Riccati equations and a pair of modified Lyapunov equations coupled by an oblique projection. When the order of the controller is equal to the order of the plant, the projection becomes the identity and the standard LQG result is recovered.

The outcome of the above development is a set of algebraic matrix equations which correspond to the necessary conditions for the Auxiliary Minimization Problem and hence to sufficient conditions for robust stability and performance. These necessary conditions characterize full- or reduced-order controllers with either choice of bounds (1.13) and (1.14) for each uncertain parameter. For control-system design, these equations can be used as follows. If a solution to the necessary conditions is obtained computationally and if certain definiteness conditions hold, then the explicitly synthesized controller 1) solves the Robust Stability Problem and 2) is guaranteed to provide robust performance bounded by $\operatorname{tr} P \tilde{V}$ over the stipulated uncertainty range.

The applicability of these results is, of course, limited to plants which are nominally stabilizable via controllers of the given order. Indeed, in this case it has been shown in [37] via topological degree theory that the optimality conditions for the case $\delta_{i}=0, i=1, \ldots, p$, possess at least one stabilizing solution. For the parameter uncertainty problem, i.e., $\delta_{i}>0$, it follows from continuity properties that a solution also exists for sufficiently small $\delta_{i}$. The actual range of uncertainty which can be stabilized and the tightness of the performance bound depend upon the conservatism of our bounds. As will be seen from a numerical example, our bounds are not generally sharp. This is not unexpected, however, due to both the sense of the partial ordering employed in (1.12) and the fact that our choice of gain-invariant bounds permits a one-step, non-iterative synthesis (rather than analysis) procedure. It should be noted that necessary and sufficient conditions for robust analysis of a block-structured class of uncertainties are obtainable using Doyle's $\mu$-function ([6]). This block structure, however, does not appear to include either the linear uncertainty model (1.1) or the matched uncertainty model of [21] as special cases.

In the present paper we present results of an illustrative numerical study for a well-known example of Doyle used in [2] to demonstrate the lack of gain margin for LQG controllers. This type of uncertainty is a special case of (1.1) obtained by taking $p=m$ and defining $B_{i}$ to be the
matrix whose ith column is the same as the ith column of $B$ and zero otherwise. To obtain fullorder, robustified controllers exhibiting performance/robustness tradeoffs, we utilize bound (1.13) for several values of $\boldsymbol{\delta}_{\boldsymbol{i}}$. To obtain these numerical results we utilized a straightforward iterative algorithm which requires only an LQG-type software package. The homotopy algorithm of [37] with appropriate extensions can also be used. Further descriptions of related algorithms and numerical results can be found in [38-40].

The development herein is self-contained with the exception that the detailed derivation of the optimality conditions has been omitted. In specialized cases the derivation has been given previously. For the case of bound (1.13) only, a derivation using Kronecker products appears in [36]. Also, a derivation without parameter uncertainties has been given in [29] using Lagrange multipliers. Overall, the derivation involves considerable matrix manipulation. Since the detailed derivation does not appear to warrant the required space, we give an outline of the proof to assist the sufficiently motivated reader in reconstructing the details.

## 2. Notation and Definitions

Note: All matrices have real entries
$\mathbb{R}, \mathbb{R}^{r \times e}, \mathbb{R}^{r}, \mathbb{E}$

$$
\|\cdot\|
$$

$I_{r}, 0_{r \times e}, O_{r}$
()$^{T},()^{-1},()^{-T}$
tr
$\oplus, \otimes$
$\mathbf{S}^{r}$
$\mathbb{N}^{r}$
$\mathbb{P}^{r}$
$Z_{1} \geq Z_{2}$
$Z_{1}>Z_{2}$
asymptotically stable matrix
$n, m, \ell, p, n_{c}, n_{i}, m_{i}$
$\tilde{n}, \tilde{n}_{i}$
$\boldsymbol{x}, \boldsymbol{u}, \boldsymbol{y}, \boldsymbol{x}_{\mathrm{c}}$
$A, A_{i} ; B, B_{i} ; C, C_{i}$
real numbers, $r \times s$ real matrices, $\mathbb{R}^{r \times 1}$, expectation
Euclidean vector norm
$r \times r$ identity matrix, $r \times s$ zero matrix, $0_{r \times r}$
transpose, inverse, inverse transpose
trace
Kronecker sum, Kronecker product ([41])
$r \times r$ symmetric matrices
$r \times r$ symmetric nonnegative-definite matrices
$r \times r$ symmetric positive-definite matrices
$Z_{1}-Z_{2} \in \mathbb{N}^{r}, \quad Z_{1}, Z_{2} \in \mathbb{S}^{r}$
$Z_{1}-Z_{2} \in \mathbb{P}^{r}, \quad Z_{1}, Z_{2} \in \mathbb{S}^{r}$
matrix with eigenvalues in open left half plane
positive integers, $\quad i \in\{1, \ldots, p\}$
$n+n_{c}, n_{i}+m_{i}, \quad i \in\{1, \ldots, p\}$
$n, m, \ell, n_{c}$-dimensional vectors
$n \times n$ matrices; $n \times m$ matrices; $\ell \times n$ matrices; $i \in\{1, \ldots, p\}$

| $A_{c}, B_{c}, C_{c}$ |  |
| :---: | :---: |
| $\bar{A}, \bar{A}_{i}$ | $\left[\begin{array}{cc}A & B C_{6} \\ B_{6} C & A_{6}\end{array}\right] \quad\left[\begin{array}{cc}A_{6} & B_{1} C_{6} \\ B_{6} C_{6} & 0\end{array}\right], \quad i \in\{1, \ldots, p\}$ |
| $\delta_{i}$ | positive number, $i \in\{1, \ldots, p)$ |
| $\Delta$ | $\left\|-\delta_{1}, \delta_{1}\right\| \times \cdots \times\left\|-\delta_{p}, \delta_{p}\right\|$ |
| $\sigma_{i}$ | real number, $i \in(1, \ldots, p)$ |
| 0 | ( $\sigma_{1}, \ldots, \sigma_{p}$ ) |
| $\dot{\sigma}_{i}(\cdot)$ | Lebesgue meaurable function on $(0, \infty), \quad 1 \in(1, \ldots, p)$ |
| $\hat{\sigma}(\cdot)$ | $\left(\dot{\partial}_{1}(\cdot), \ldots, \dot{\delta}_{p}(\cdot)\right)$ |
| $L_{\infty}([0, \infty), \Delta)$ | Lebesgue measurable functions on $(0, \infty)$ with valuen in $\Delta$ |
| $\alpha_{i}$ | positive number, $i \in\{1, \ldots, p\}$ |
| $D_{i}, E_{i}, H_{i}, K_{i}$ | $n \times n_{1}, n_{4} \times n, n \times m_{1}, m_{0} \times m$ matrices, $\quad i \in(1, \ldots, p)$ |
| $\tilde{D}_{\text {i }}, E_{\text {E }}$ | $\bar{n} \times \tilde{n}_{1}, \tilde{n}_{1} \times \bar{n}$ matrices, $i \in\{1, \ldots, p\}$ |
| $\Sigma^{\prime}, \Sigma^{\prime \prime}$ | see Section 6 |
| $R_{1}$ | state weighting matrix in $\mathbb{N}^{n}$ |
| $R_{2}$ | control weighting matrix in $\mathbb{P}^{m}$ |
| $R_{12}$ | $n \times m$ croes weighting matrix such that $R_{1}-R_{12} R_{2}^{-1} R_{12}^{T} \geq 0$ |
| $\bar{R}$ | $\left[\begin{array}{cc} R_{1} & R_{12} C_{e} \\ C_{6}^{T} & R_{12}^{T} \\ C_{c}^{T} & R_{2} C_{6} \end{array}\right]$ |
| $w_{1}(\cdot)$ | $n$-dimensional white noise |
| $w_{2}(\cdot)$ | $\ell$-dimensional white noise |
| $V_{1}$ | intensity of $\omega_{1}(\cdot)$ in $\mathbb{N}^{\boldsymbol{n}}$ |
| $V_{2}$ | intensity of $w_{2}(\cdot)$ in $\mathbb{P}^{\ell}$ |
| $V_{12}$ | $n \times \ell$ cross intensity of $w_{1}(\cdot), w_{2}(\cdot)$ |
| $\tilde{V}$ | $\left[\begin{array}{cc} V_{1} & V_{12} B_{c}^{T} \\ B_{c} V_{12}^{T} & B_{c} V_{2} B_{c}^{T} \end{array}\right]$ |

## 3. Robust Stability and Robust Performance Roblems

In this section we state the Robust Stability Problem and Robust Performance Problem along with related notation for later use.

Robust Stability Problem. For fixed $n_{c} \leq n$, determine $\left(A_{c}, B_{c}, C_{c}\right) \in \mathbb{R}^{n_{c} \times n_{c}} \times \mathbb{R}^{n_{c} \times \ell} \times$ $\mathbb{R}^{\boldsymbol{m} \times n_{e}}$ such that the closed-loop system consisting of the $n$ th-order controlled plant

$$
\begin{equation*}
\dot{x}(t)=\left(A+\sum_{i=1}^{p} \hat{\sigma}_{i}(t) A_{i}\right) x(t)+\left(B+\sum_{i=1}^{p} \hat{\sigma}_{i}(t) B_{i}\right) u(t), \quad \text { a.a. } t \in[0, \infty), \tag{3.1}
\end{equation*}
$$

measurements

$$
\begin{equation*}
y(t)=\left(C+\sum_{i=1}^{p} \hat{\sigma}_{i}(t) C_{i}\right) x(t) \tag{3.2}
\end{equation*}
$$

and $n_{c}$ th-order dynamic compensator

$$
\begin{gather*}
\dot{x}_{c}(t)=A_{c} x_{c}(t)+B_{c} y(t),  \tag{3.3}\\
u(t)=C_{c} x_{c}(t) \tag{3.4}
\end{gather*}
$$

is asymptotically stable* for all $\hat{\sigma}(\cdot) \in L_{\infty}((0, \infty), \Delta)$.
Robust Performance Problem. For fixed $n_{c} \leq n$, determine $\left(A_{c}, B_{c}, C_{c}\right) \in \mathbb{R}^{n_{0}} \times n_{c} \times \mathbb{R}^{n_{0} \times \ell} \times$ $\mathbb{R}^{\boldsymbol{m} \times n}$. such that, for the closec-loop system consisting of the $n$ th-order controlled and disturbed plant

$$
\begin{equation*}
\dot{x}(t)=\left(A+\sum_{i=1}^{p} \hat{\sigma}_{i}(t) A_{i}\right) x(t)+\left(B+\sum_{i=1}^{p} \hat{\sigma}_{i}(t) B_{i}\right) u(t)+w_{1}(t), \quad \text { a.a. } t \in[0, \infty) \tag{3.5}
\end{equation*}
$$

noisy measurements

$$
\begin{equation*}
y(t)=\left(C+\sum_{i=1}^{p} \hat{\sigma}_{i}(t) C_{i}\right) x(t)+w_{2}(t) \tag{3.6}
\end{equation*}
$$

and $n_{c}$ th-order dynamic compensator (3.3), (3.4), the performance criterion

$$
\begin{equation*}
J\left(A_{c}, B_{c}, C_{c}\right) \triangleq \sup _{\partial(\cdot) \in L_{\infty}([0, \infty), \Delta)} \limsup _{t \rightarrow \infty} \mathbb{E}\left[x^{T}(t) R_{1} x(t)+2 x^{T}(t) R_{12} u(t)+u^{T}(t) R_{2} u(t)\right] \tag{3.7}
\end{equation*}
$$

is minimized.
For each controller $\left(A_{c}, B_{c}, C_{c}\right)$ and parameter variation $\hat{\sigma}(\cdot) \in L_{\infty}([0, \infty), \Delta)$ the undisturbed closed-loop system (3.1)-(3.4) is given by

$$
\begin{equation*}
\dot{\tilde{x}}(t)=\left(\tilde{A}+\sum_{i=1}^{p} \hat{\sigma}_{i}(t) \tilde{A}_{i}\right) \tilde{x}(t), \quad \text { a.a. } t \in[0, \infty) \tag{3.8}
\end{equation*}
$$

while the disturbed closed-loop system (3.3)-(3.6) is

$$
\begin{equation*}
\dot{\tilde{x}}(t)=\left(\tilde{A}+\sum_{i=1}^{p} \hat{\sigma}_{i}(t) \tilde{A}_{i}\right) \tilde{x}(t)+\tilde{w}(t), \quad \text { a.a } t \in[0, \infty) \tag{3.9}
\end{equation*}
$$

[^28]Also (see, e.g., [43], p. 194), let $\tilde{\Phi}:[0, \infty) \rightarrow \mathbb{R}^{\boldsymbol{n} \times{ }^{\boldsymbol{n}}}$ be the unique absolutely continuous solution to

$$
\begin{gather*}
\dot{\Phi}(t)=\left(\tilde{A}+\sum_{i=1}^{p} \hat{\sigma}_{i}(t) \tilde{A}_{i}\right) \tilde{\Phi}(t), \quad \text { a.a. } t \in[0, \infty),  \tag{3.10}\\
\tilde{\Phi}(0)=I_{\tilde{n}}, \tag{3.11}
\end{gather*}
$$

and recall that $\tilde{\boldsymbol{\Phi}}^{-1}(\cdot)$ satisfies

$$
\begin{equation*}
\frac{d}{d t} \tilde{\Phi}^{-1}(t)=-\tilde{\Phi}^{-1}(t)\left(\tilde{A}+\sum_{i=1}^{p} \hat{\sigma}_{i}(t) \tilde{A}_{i}\right), \quad \text { a.a. } t \in[0, \infty) \tag{3.12}
\end{equation*}
$$

## 4. Sufficient Conditions for Robust Stability and Performance

For robust stability we characterize quadratic Lyapunov functions for the closed-loop system.
Theorem 4.1. Let $\Omega: \mathbb{P}^{\bar{n}} \times \mathbb{R}^{n_{0} \times \ell} \times \mathbb{R}^{m \times n_{0}} \rightarrow \mathbb{S}^{\bar{n}}$ satisfy

$$
\begin{equation*}
\sum_{i=1}^{p} \sigma_{i}\left(\tilde{A}_{i}^{T} P+P \tilde{A}_{i}\right)<\Omega\left(P, B_{c}, C_{c}\right), \quad \sigma \in \Delta, \quad\left(P, B_{c}, C_{c}\right) \in \mathbb{P}^{\bar{n}} \times \mathbb{R}^{n_{c} \times \ell} \times \mathbb{R}^{m \times n_{e}} \tag{4.1}
\end{equation*}
$$

If, for some $\left(A_{c}, B_{c}, C_{c}\right) \in \mathbb{R}^{n_{c} \times n_{\bullet}} \times \mathbb{R}^{n_{c} \times \ell} \times \mathbb{R}^{m \times n_{e}}$, there exists $P \in \mathbb{P}^{\tilde{n}}$ satisfying

$$
\begin{equation*}
0=\tilde{A}^{T} P+P \tilde{A}+\Omega\left(P, B_{c}, C_{c}\right), \tag{4.2}
\end{equation*}
$$

then ( $A_{c}, B_{c}, C_{c}$ ) solves the Robust Stability Problem.
Proof. Define the Lyapunov function

$$
V(\tilde{x}) \triangleq \tilde{x}^{T} \rho \tilde{x}, \quad \tilde{x} \in \mathbb{R}^{\dot{n}}
$$

For $t \in[0, \infty)$ and $\tilde{x}(t)$ satisfying (3.8), it follows from (4.2) that

$$
\begin{aligned}
\dot{V}(\tilde{x}(t)) & =\dot{\tilde{x}}^{T}(t) P \tilde{x}(t)+\tilde{x}^{T}(t) P \dot{\tilde{x}}(t) \\
& =\tilde{x}^{T}(t)\left[\left(\tilde{A}+\sum_{i=1}^{p} \hat{\sigma}_{i}(t) \tilde{A}_{i}\right)^{T} P+P\left(\tilde{A}+\sum_{i=1}^{p} \hat{\sigma}_{i}(t) \tilde{A}_{i}\right)\right] \tilde{x}(t) \\
& =\tilde{x}^{T}(t)\left[\sum_{i=1}^{p} \hat{\sigma}_{i}(t)\left(\tilde{A}_{i}^{T} P+P \tilde{A}_{i}\right)-\Omega\left(P, B_{c}, C_{c}\right)\right] \tilde{x}(t) .
\end{aligned}
$$

Since $\hat{\sigma}(t) \in \Delta, t \in[0, \infty)$, it follows from (4.1) that there exists $\gamma>0$ such that $\dot{V}(\tilde{x}(t)) \leq$ $-\gamma\|\tilde{x}(t)\|^{2}, t \in[0, \infty)$.

Remark 4.1. If $\left(A_{c}, B_{c}, C_{c}\right)$ solves the Robust Stability Problem then

$$
\begin{equation*}
\lim _{t \rightarrow \infty} \tilde{\Phi}(t)=0, \quad \hat{\sigma}(\cdot) \in L_{\infty}([0, \infty), \Delta) \tag{4.3}
\end{equation*}
$$

Remark 4.2. As will be seen the bound (4.1) will be guaranteed for all $P, B_{c}, C_{c}$ by suitable construction of the function $\Omega$. In addition, the existence of a solution $\mathcal{P}$ to (4.2) need not be verified in practice. Rather, (4.2) is a result of numerically solving the necessary conditions for the Auxiliary Minimization Problem given by Theorem 6.1.

For the Robust Performance Problem the cost can be expressed in terms of the closed-loop second-moment matrix.

Proposition 4.1. For $\left(A_{c}, B_{c}, C_{c}\right) \in \mathbb{R}^{n_{c} \times n_{c}} \times \mathbb{R}^{n_{c} \times \ell} \times \mathbb{R}^{m \times n_{c}}$ and $\hat{\sigma}(\cdot) \in L_{\infty}([0, \infty), \Delta)$ the second-moment matrix

$$
\begin{equation*}
\tilde{Q}(t) \triangleq \mathbb{E}\left[\tilde{x}(t) \tilde{x}^{T}(t)\right], \quad t \in[0, \infty) \tag{4.4}
\end{equation*}
$$

satisfies

$$
\begin{equation*}
\dot{\tilde{Q}}(t)=\left(\tilde{A}+\sum_{i=1}^{p} \hat{\sigma}_{i}(t) \tilde{A}_{i}\right) \tilde{Q}(t)+\tilde{Q}(t)\left(\tilde{A}+\sum_{i=1}^{p} \hat{\sigma}_{i}(t) \tilde{A}_{i}\right)^{T}+\tilde{V}, \quad \text { a.a. } t \in[0, \infty) \tag{4.5}
\end{equation*}
$$

or, equivalently,

$$
\begin{equation*}
\tilde{Q}(t)=\tilde{\Phi}(t) \tilde{Q}(0) \tilde{\Phi}^{T}(t)+\int_{0}^{t} \tilde{\Phi}(t) \tilde{\Phi}^{-1}(s) \tilde{V} \tilde{\Phi}^{-T}(s) \tilde{\Phi}^{T}(t) d s, \quad t \in[0, \infty) \tag{4.6}
\end{equation*}
$$

Furthermore,

$$
\begin{equation*}
J\left(A_{c}, B_{c}, C_{c}\right)=\sup _{\left.\partial(\cdot) \in L_{\infty}(0, \infty), \Delta\right)} \limsup _{t \rightarrow \infty} \operatorname{tr} \tilde{Q}(t) \tilde{R} \tag{4.7}
\end{equation*}
$$

or, equivalently,

$$
\begin{align*}
& J\left(A_{c}, B_{c}, C_{c}\right) \\
& \quad \triangleq \sup _{\partial(\cdot) \in \mathrm{L}_{\infty}([0, \infty), \Delta)} \limsup _{t \rightarrow \infty} \operatorname{tr}\left[\tilde{\Phi}(t) \tilde{Q}(0) \tilde{\Phi}^{T}(t) \tilde{R}+\int_{0}^{t} \tilde{\Phi}(t) \tilde{\Phi}^{-1}(s) \tilde{V} \tilde{\Phi}^{-T}(s) \tilde{\Phi}^{T}(t) d s \tilde{R}\right] . \tag{4.8}
\end{align*}
$$

Proof. The second-moment equation (4.5) is a direct consequence of the Ito differential rule (see [44], p. 142) while (4.6) follows by direct verification. Finally, (4.7) is immediate.

We now proceed to obtain an upper bound for $J$ in terms of the matrix $P$. The following lemma is required.

Lemme 4.1. Let $\Omega: \mathbb{P}^{n} \times \mathbb{R}^{n_{0} \times \ell} \times \mathbb{R}^{m \times n_{c}} \rightarrow \mathbb{S}^{n}$ and $\left(A_{c}, B_{c}, C_{c}\right) \in \mathbb{R}^{n_{0} \times n_{0}} \times \mathbb{R}^{n_{0} \times \ell} \times \mathbb{R}^{m \times n_{0}}$ be given. Then $P \in \mathbb{P}^{n}$ satisfies (4.2) if and only if $P$ satisfies

$$
\begin{align*}
& P=\tilde{\Phi}^{T}(t) P \tilde{\Phi}(t) \\
& +\int_{0}^{t} \tilde{\Phi}^{T}(t) \tilde{\Phi}^{-T}(s)\left[\Omega\left(P, B_{c}, C_{c}\right)-\sum_{i=1}^{p} \hat{\sigma}_{i}(t)\left(\tilde{A}_{i}^{T} P+P \tilde{A}_{i}\right)\right] \tilde{\Phi}^{-1}(s) \tilde{\Phi}(t) d s,  \tag{4.9}\\
& \\
& \hat{\sigma}(\cdot) \in L_{\infty}([0, \infty), \Delta), \quad t \in[0, \infty) .
\end{align*}
$$

Proof. Suppose $P$ satisfies (4.2). Then, for $t \in[0, \infty)$,

$$
\begin{aligned}
0= & \tilde{\Phi}^{-T}(t)\left(\tilde{A}+\sum_{i=1}^{p} \hat{\sigma}_{i}(t) \tilde{A}_{i}\right)^{T} P \tilde{\Phi}^{-1}(t)+\tilde{\Phi}^{-T}(t) P\left(\tilde{A}+\sum_{i=1}^{p} \hat{\sigma}_{i}(t) \tilde{A}_{i}\right) \tilde{\Phi}^{-1}(t) \\
& +\tilde{\Phi}^{-T}(t)\left[\Omega\left(P, B_{c}, C_{c}\right)-\sum_{i=1}^{p} \hat{\sigma}_{i}(t)\left(\tilde{A}_{i}^{T} P+P \tilde{A}_{i}\right)\right] \tilde{\Phi}^{-1}(t) \\
= & -\frac{d}{d t}\left[\tilde{\Phi}^{-T}(t) P \tilde{\Phi}^{-2}(t)\right]+\tilde{\Phi}^{-T}(t)\left[\Omega\left(P, B_{c}, C_{c}\right)-\sum_{i=1}^{p} \hat{\sigma}_{i}(t)\left(\tilde{A}_{i}^{T} P+P \tilde{A}_{i}\right)\right] \tilde{\Phi}^{-1}(t),
\end{aligned}
$$

which yields

$$
0=-\tilde{\Phi}^{-T}(t) P \tilde{\Phi}^{-1}(t)+P+\int_{0}^{t} \tilde{\Phi}^{-T}(s)\left[\Omega\left(P, B_{c}, C_{c}\right)-\sum_{i=1}^{p} \hat{\sigma}_{i}(s)\left(\tilde{A}_{i}^{T} P+P \tilde{A}_{i}\right)\right] \tilde{\Phi}^{-1}(s) d s
$$

Thus (4.9) is satisfied. Conversely, suppose $P$ satisfies (4.9). Differentiating with respect to $t$ using Leibniz' rule yields

$$
\begin{aligned}
0= & \left(\tilde{A}+\sum_{i=1}^{p} \hat{\sigma}_{i}(t) \tilde{A}_{i}\right)^{T} \tilde{\Phi}^{T}(t) P \tilde{\Phi}(t)+\tilde{\Phi}^{T}(t) P \tilde{\Phi}(t)\left(\tilde{A}+\sum_{i=1}^{p} \hat{\sigma}_{i}(t) \tilde{A}_{i}\right) \\
& +\left(\tilde{A}+\sum_{i=1}^{p} \hat{\sigma}_{i}(t) \tilde{A}_{i}\right)^{T} \int_{0}^{t} \tilde{\Phi}^{T}(t) \tilde{\Phi}^{-T}(s)\left[\Omega\left(P, B_{c}, C_{c}\right)-\sum_{i=1}^{p} \hat{\sigma}_{i}(s)\left(\tilde{A}_{i}^{T} P+P \tilde{A}_{i}\right)\right] \tilde{\Phi}^{-1}(s) \tilde{\Phi}(t) d s \\
& +\int_{0}^{t} \tilde{\Phi}^{T}(t) \tilde{\Phi}^{-T}(s)\left[\Omega\left(P, B_{c}, C_{c}\right)-\sum_{i=1}^{p} \hat{\sigma}_{i}(s)\left(\tilde{A}_{i}^{T} P+P \tilde{A}_{i}\right)\right] \tilde{\Phi}^{-1}(s) \tilde{\Phi}^{\prime}(t) d s\left(\tilde{A}+\sum_{i=1}^{p} \hat{\sigma}_{i}(t) \tilde{A}_{i}\right) \\
& +\Omega\left(P, B_{c}, C_{c}\right)-\sum_{i=1}^{p} \hat{\sigma}_{i}(t)\left(\tilde{A}_{i}^{T} P+P \tilde{A}_{i}\right) \\
= & \left(\tilde{A}+\sum_{i=1}^{p} \hat{\sigma}_{i}(t) \tilde{A}_{i}\right)^{T} P+P\left(\tilde{A}+\sum_{i=1}^{p} \hat{\sigma}_{i}(t) \tilde{A}_{i}\right)+\Omega\left(P, B_{c}, C_{c}\right)-\sum_{i=1}^{p} \hat{\sigma}_{i}(t)\left(\tilde{A}_{i}^{T} P+P \tilde{A}_{i}\right) \\
= & \tilde{A}^{T} P+P \tilde{A}+\Omega\left(P, B_{c}, C_{c}\right) .
\end{aligned}
$$

Hence (4.2) is satisfied.
Remark 4.3. Note the identity

$$
\begin{equation*}
\operatorname{tr} \int_{0}^{t} \tilde{\Phi}(t) \tilde{\Phi}^{-1}(s) \tilde{V} \tilde{\Phi}^{-T}(s) \tilde{\Phi}^{T}(t) d s \tilde{R}=\operatorname{tr} \int_{0}^{t} \tilde{\Phi}^{T}(t) \tilde{\Phi}^{-T}(s) \tilde{R} \tilde{\Phi}^{-1}(s) \tilde{\Phi}(t) d s \tilde{V} \tag{4.10}
\end{equation*}
$$

$$
\left(A_{c}, B_{c}, C_{c}\right) \in \mathbb{R}^{n_{0} \times n_{c}} \times \mathbb{R}^{n_{c} \times \ell} \times \mathbb{R}^{m \times n_{0}}, \quad \hat{\sigma}(\cdot) \in L_{\infty}([0, \infty), \Delta), \quad t \in[0, \infty)
$$

We are now in a position to bound the cost $J$ by means of the matrix $P$.
Theorem 4.2. Let $\Omega: \mathbb{P}^{\bar{n}} \times \mathbb{R}^{n_{0} \times \ell} \times \mathbb{R}^{m \times n_{0}} \rightarrow \mathbb{S}^{\bar{n}}$ satisfy (4.1) and

$$
\begin{equation*}
\sum_{i=1}^{p} \sigma_{i}\left(\tilde{A}_{i}^{T} P+P \tilde{A}_{i}\right)+\tilde{R} \leq \Omega\left(P, B_{c}, C_{c}\right), \quad \sigma \in \Delta, \quad\left(P, B_{c}, C_{c}\right) \in \mathbb{P}^{\bar{n}} \times \mathbb{R}^{n_{\bullet} \times \ell} \times \mathbb{R}^{m \times n_{0}} \tag{4.11}
\end{equation*}
$$

If, for some $\left(A_{c}, B_{c}, C_{c}\right) \in \mathbb{R}^{n_{c} \times n_{c}} \times \mathbb{R}^{n_{c} \times \ell} \times \mathbb{R}^{m \times n_{c}}$, there exists $\mathcal{P} \in \mathbb{P}^{\dot{n}}$ satisfying (4.2), then

$$
\begin{equation*}
J\left(A_{c}, B_{c}, C_{c}\right) \leq \operatorname{tr} P \tilde{V} \tag{4.12}
\end{equation*}
$$

Proof. From (4.8)-(4.10) and (4.3) it follows that

$$
\begin{aligned}
J\left(A_{c}, B_{c}, C_{c}\right)= & \sup _{\partial(\cdot) \in L_{\infty}([0, \infty), \Delta)} \limsup _{t \rightarrow \infty} \operatorname{tr}\left\{\tilde{\Phi}(t) \tilde{Q}(0) \tilde{\Phi}^{T}(t) \tilde{R}+P \tilde{V}-\tilde{\Phi}^{T}(t) P \tilde{\Phi}(t) \tilde{V}\right. \\
& \left.-\int_{0}^{t} \tilde{\Phi}^{T}(t) \tilde{\Phi}^{-T}(s)\left[\Omega\left(P, B_{c}, C_{c}\right)-\tilde{R}-\sum_{i=1}^{p} \hat{\sigma}_{i}(s)\left(\tilde{A}_{i}^{T} P+P \tilde{A}_{i}\right)\right] \tilde{\Phi}^{-1}(s) \tilde{\Phi}(t) d s \tilde{V}\right\} \\
\leq & \sup _{\partial(\cdot) \in L_{\infty}([0, \infty), \Delta)} \limsup _{t \rightarrow \infty} \operatorname{tr}\left[\tilde{\Phi}(t) \tilde{Q}(0) \tilde{\Phi}^{T}(t) \tilde{R}+P \tilde{V}\right] \\
= & \operatorname{tr} P \tilde{V} . \square
\end{aligned}
$$

Remark 4.4. Note that since $\tilde{R} \geq 0$, (4.11) implies

$$
\begin{equation*}
\sum_{i=1}^{p} \sigma_{i}\left(\tilde{A}_{i}^{T} P+P \tilde{A}_{i}\right) \leq \Omega\left(P, B_{c}, C_{c}\right), \quad \sigma \in \Delta \tag{4.13}
\end{equation*}
$$

which is a weak form of (4.1). If $\tilde{R}>0$ then (4.11) implies (4.1). This implication is not surprising since (4.11) implies robust performance while (4.1) implies robust stability.
5. Choice of Bounds

To satisfy (4.11) $\Omega(\cdot, \cdot, \cdot)$ is chosen to be of the form

$$
\begin{equation*}
\Omega\left(P, B_{c}, C_{c}\right)=\sum_{i=1}^{p} \Lambda_{i}\left(P, B_{c}, C_{c}\right)+\tilde{R}, \tag{5.1}
\end{equation*}
$$

where, for each $i=1, \ldots, p, \Lambda_{i}: \mathbb{P}^{\dot{n}} \times \mathbb{R}^{n_{0} \times \ell} \times \mathbb{R}^{m \times n_{0}} \rightarrow \mathbb{S}^{\dot{n}}$ satisfies

$$
\begin{equation*}
\sigma_{i}\left(\tilde{A}_{i}^{T} P+P \tilde{A}_{i}\right) \leq \Lambda_{i}\left(P, B_{c}, C_{c}\right), \quad \sigma_{i} \in\left[-\delta_{i}, \delta_{i}\right], \quad\left(P, B_{c}, C_{c}\right) \in \mathbb{P}^{\dot{n}} \times \mathbb{R}^{n_{c} \times \ell} \times \mathbb{R}^{m \times n_{e}} . \tag{5.2}
\end{equation*}
$$

Two distinct choices for the bound $\Lambda_{i}$ are considered. As pointed out in Section 1, the first choice corresponds to a right shift/multiplicative white noise model ([33]), while the second bound generalizes results found in [26].

Proposition 5.1. For all $\alpha_{i}>0$ the function

$$
\begin{equation*}
\Lambda_{i}\left(P, B_{c}, C_{c}\right)=\delta_{i}\left(\alpha_{i} P+\alpha_{i}^{-1} \tilde{A}_{i}^{T} P \tilde{A}_{i}\right) \tag{5.3}
\end{equation*}
$$

satisfies (5.2).
Proof. Note that

$$
\begin{aligned}
0 & \leq\left[\sigma_{i}\left(\alpha_{i} / \delta_{i}\right)^{\frac{1}{2}} I_{\tilde{n}}-\left(\delta_{i} / \alpha_{i}\right)^{\frac{1}{2}} \tilde{A}_{i}\right]^{T} P\left[\sigma_{i}\left(\alpha_{i} / \delta_{i}\right)^{\frac{1}{2}} I_{\tilde{n}}-\left(\delta_{i} / \alpha_{i}\right)^{\frac{1}{2}} \tilde{A}_{i}\right] \\
& =\sigma_{i}^{2}\left(\alpha_{i} / \delta_{i}\right) P+\left(\delta_{i} / \alpha_{i}\right) \tilde{A}_{i}^{T} P \tilde{A}_{i}-\sigma_{i}\left(\tilde{A}_{i}^{T} P+P \tilde{A}_{i}\right),
\end{aligned}
$$

which, since $\alpha_{i}^{2} \leq \delta_{i}^{2}$, implies (5.2).
Proposition-5.2. For all $\tilde{D}_{i} \in \mathbb{R}^{\bar{n} \times \bar{n}_{i}}$ and $\tilde{E}_{i} \in \mathbb{R}^{\bar{n}_{i} \times \bar{n}}$ satisfying

$$
\begin{equation*}
\tilde{A}_{i}=\tilde{D}_{i} \tilde{E}_{i} \tag{5.4}
\end{equation*}
$$

the function

$$
\begin{equation*}
\Lambda_{i}\left(P, B_{c}, C_{c}\right)=\delta_{i}\left(\tilde{E}_{i}^{T} \tilde{E}_{i}+P \tilde{D}_{i} \tilde{D}_{i}^{T} P\right) \tag{5.5}
\end{equation*}
$$

satisfies (5.2).
Proof. Note that

$$
\begin{aligned}
0 & \leq\left[\delta_{i}^{\frac{1}{2}} \tilde{E}_{i}-\sigma_{i} \delta_{i}^{-\frac{1}{2}} \tilde{D}_{i}^{T} \rho\right]^{T}\left[\delta_{i}^{\frac{1}{2}} \tilde{E}_{i}-\sigma_{i} \delta_{i}^{-\frac{1}{2}} \tilde{D}_{i}^{T} \rho\right] \\
& =\delta_{i} \tilde{E}_{i}^{T} \tilde{E}_{i}+\left(\sigma_{i}^{2} / \delta_{i}\right) P \tilde{D}_{i} \tilde{D}_{i}^{T} \rho-\sigma_{i}\left(\tilde{A}_{i}^{T} \rho+\rho \tilde{A}_{i}\right),
\end{aligned}
$$

which implies (5.2).

## 6. The Auxiliary Minimization Problem and Necessary Conditions for Optimality

To optimize robust performance while retaining robust stability, we consider the following problem for which the cost functional is given by the bound (4.12).

Auxiliary Minimization Problem. For $i=1, \ldots, p$ let $\Lambda_{i}$ be given by either (5.3) or (5.5). Determine $\left(P, A_{c}, B_{c}, C_{c}\right) \in \mathbb{P}^{n} \times \mathbb{R}^{n_{d} \times n_{c}} \times \mathbb{R}^{n_{d} \times \ell} \times \mathbb{R}^{m \times n_{d}}$ which minimizes

$$
\begin{equation*}
J\left(P, A_{c}, B_{c}, C_{c}\right) \triangleq \operatorname{tr} P \tilde{V} \tag{6.1}
\end{equation*}
$$

subject to

$$
\begin{equation*}
0=\tilde{A}^{T} p+\rho \tilde{A}+\sum_{i=1}^{p} \Lambda_{i}\left(P, B_{c}, C_{c}\right)+\tilde{R} \tag{6.2}
\end{equation*}
$$

and

$$
\begin{equation*}
\sum_{i=1}^{p} \sigma_{i}\left(\tilde{A}_{i}^{T} P+P \tilde{A}_{i}\right)<\sum_{i=1}^{p} \Lambda_{i}\left(P, B_{c}, C_{c}\right)+\tilde{R}, \quad \sigma \in \Delta . \tag{6.3}
\end{equation*}
$$

Remark 6.1. Note that (6.3) enforces both (4.1) and (4.11) to guarantee robust stability and performance.

To derive first-order necessary conditions for the Auxiliary Minimization Problem, note that the constraint (6.3) defines an open set.

Proposition 6.1. The set of $\left(P, B_{c}, C_{c}\right) \in \mathbb{P}^{\bar{n}} \times \mathbb{R}^{n_{c} \times \ell} \times \mathbb{R}^{m \times n_{0}}$ satisfying (6.3) is open.
Proof. Since $\Lambda_{i}(\cdot,, \cdot)$ is continuous it can be shown that the function

$$
f\left(P, B_{c}, C_{c}\right) \triangleq \min _{\sigma \in \Delta} \lambda_{\min }\left\{\sum_{i=1}^{p} \Lambda_{i}\left(P, B_{c}, C_{c}\right)+\tilde{R}-\sum_{i=1}^{p} \sigma_{i}\left(\tilde{A}_{i}^{T} P+P \tilde{A}_{i}\right)\right\}
$$

is also continuous. Since (6.3) is equivalent to $0<f\left(P, B_{c}, C_{c}\right)$, the result is immediate.
To obtain explicit feedback gain expressions we shall require two additional technical assumptions. If bound (5.3) is chosen for a given $i \in\{1, \ldots, p\}$ we require

$$
\begin{equation*}
B_{i} \neq 0 \Longrightarrow C_{i}=0 \tag{6.4}
\end{equation*}
$$

i.e., $B_{i}$ and $C_{i}$ are not simultaneously nonzero. Of course, both $B_{i}$ and $C_{i}$ may be zero. Assumption (6.4) implies that parameter uncertainties in $B$ and $C$ must be modeled as uncorrelated. Correlation between uncertainties in $A$ and $B$ or $A$ and $C$ is, of course, permitted. Furthermore, if bound (5.5) is chosen for a given $i \in\{1, \ldots, p\}$ we require

$$
\begin{equation*}
C_{i}=0 \tag{6.5}
\end{equation*}
$$

When utilizing bound (5.3) the positive constant $\alpha_{i}$ shall be considered fixed but arbitrary. Furthermore, for bound (5.5), let $D_{i} \in \mathbb{R}^{n \times n_{i}}, E_{i} \in \mathbb{R}^{n_{i} \times n}, H_{i} \in \mathbb{R}^{n \times m_{i}}$ and $K_{i} \in \mathbb{R}^{m_{i} \times m}$ satisfy

$$
\begin{equation*}
A_{i}=D_{i} E_{i}, \quad B_{i}=H_{i} K_{i}, \tag{6.6}
\end{equation*}
$$

and define $\tilde{D}_{i}, \tilde{E}_{i}$ satisfying (5.4) by

$$
\tilde{D}_{i} \triangleq\left[\begin{array}{cc}
D_{i} & H_{i}  \tag{6.7}\\
O_{n_{e} \times n_{i}} & 0_{n_{0} \times m_{i}}
\end{array}\right], \quad \tilde{E}_{i} \triangleq\left[\begin{array}{cc}
E_{i} & 0_{n_{i} \times n_{e}} \\
0_{m_{i} \times n} & K_{i} C_{c}
\end{array}\right] .
$$

In addition to the open set defined by (6.3), the derivation of the necessary conditions requires that ( $P, A_{c}, B_{c}, C_{c}$ ) be further restricted so that

$$
\begin{gather*}
\left(\tilde{A}+\frac{1}{2} \sum^{\prime} \delta_{i} \alpha_{i} I_{\tilde{n}}+\sum^{\prime \prime} \delta_{i} \tilde{D}_{i} \tilde{D}_{i}^{T} P\right) \oplus\left(\tilde{A}+\frac{1}{2} \sum^{\prime} \delta_{i} \alpha_{i} I_{\tilde{n}}+\sum^{\prime \prime} \delta_{i} \tilde{D}_{i} \tilde{D}_{i}^{T} P\right)  \tag{6.8}\\
+\sum^{\prime}\left(\delta_{i} / \alpha_{i}\right) \tilde{A}_{i} \otimes \tilde{A}_{i} \quad \text { is asymptotically stable }
\end{gather*}
$$

and

$$
\begin{equation*}
\left(A_{c}, B_{c}, C_{c}\right) \text { is controllable and observable. } \tag{6.9}
\end{equation*}
$$

In (6.8) the notation $\Sigma^{\prime}$ and $\Sigma^{\prime \prime}$ denotes summation over indices for which bounds (5.3) and (5.5), repectively, have been chosen. Note that (6.8) and (6.9) play no role in the Auxiliary Minimization Problem and thus need not be verified for robust stability or robust performance.

For arbitrary $Q, P, \hat{Q}, \hat{P} \in \mathbb{R}^{n \times n}$ define the following notation:

$$
\begin{array}{rlr}
R_{2 a} \triangleq R_{2}+\sum^{\prime}\left(\delta_{i} / \alpha_{i}\right) B_{i}^{T}(P+\hat{P}) B_{i}+\sum^{\prime \prime} \delta_{i} K_{i}^{T} K_{i}, & V_{2 a} \triangleq V_{2}+\sum^{\prime}\left(\delta_{i} / \alpha_{i}\right) C_{i}(Q+\hat{Q}) C_{i}^{T}, \\
P_{a} \triangleq B^{T} P+R_{12}^{T}+\sum^{\prime}\left(\delta_{i} / \alpha_{i}\right) B_{i}^{T}(P+\hat{P}) A_{i}, & Q_{a} \triangleq Q C^{T}+V_{12}+\sum^{\prime}\left(\delta_{i} / \alpha_{i}\right) A_{i}(Q+\hat{Q}) C_{i}^{T}, \\
D \triangleq \sum^{\prime \prime} \delta_{i}\left(D_{i} D_{i}^{T}+H_{i} H_{i}^{T}\right), & E \triangleq \sum^{\prime \prime} \delta_{i} E_{i}^{T} E_{i}, \\
\hat{A} \triangleq A+\frac{1}{2} \sum^{\prime} \delta_{i} \alpha_{i} I_{n}, \quad \hat{A}_{P} \triangleq \hat{A}-B R_{2 a}^{-1} P_{a}, & \hat{A}_{Q} \triangleq \hat{A}-Q_{a} V_{2 a}^{-1} C .
\end{array}
$$

The following lemma will be needed.
Lemma 6.1. If $\hat{Q}, \hat{P} \in \mathbb{N}^{n}$ and rank $\hat{Q} \hat{P}=n_{c}$, then there exist $G, \Gamma \in \mathbb{R}^{n_{0} \times n}$ and invertible $M \in \mathbb{R}^{n_{0} \times n_{0}}$ such that

$$
\begin{gather*}
\hat{Q} \hat{P}=G^{\boldsymbol{T}} M \Gamma  \tag{6.10}\\
\Gamma G^{\boldsymbol{T}}=I_{n_{0}} \tag{6.11}
\end{gather*}
$$

Furthermore, $G, M$ and $\Gamma$ are unique except for a change of basis in $\mathbb{R}^{n_{d}}$.
Proof. The result is an immediate consequence of [45], Theorem 6.2.5, p. 123.
Note that because of (6.11), the $n \times n$ matrix $\tau \triangleq G^{T} \Gamma$ is idempotent, i.e., $\tau^{2}=\tau$. Since $\tau$ is not necessarily symmetric, it is an oblique projection. Also, define $\tau_{\perp} \triangleq I_{n}-\tau$.

Theorem 6.1. Suppose ( $P, A_{c}, B_{c}, C_{c}$ ) solves the Auxiliary Minimization Problem subject to (6.8) and (6.9). Then there exist $P, Q, \hat{P}, \hat{Q} \in \mathbb{N}^{n}$ such that $P, A_{c}, B_{c}, C_{c}$ are given by

$$
\begin{align*}
P & =\left[\begin{array}{ll}
P+\hat{P} & -\hat{P} G^{T} \\
-G \hat{P} & G \hat{P} G^{T}
\end{array}\right]  \tag{6.12}\\
A_{c} & =\Gamma\left(A-Q_{a} V_{2 a}^{-1} C-B R_{2 a}^{-1} P_{a}+D P\right) G^{T}  \tag{6.13}\\
B_{c} & =\Gamma Q_{a} V_{2 a}^{-1}  \tag{6.14}\\
C_{c} & =-R_{2 a}^{-1} P_{a} G^{T} \tag{6.15}
\end{align*}
$$

and such that $P, Q, \hat{P}, \hat{Q}$ satisfy

$$
\begin{gather*}
\begin{aligned}
& 0= \hat{A}^{T} P+P \hat{A}+R_{1}+\sum^{\prime}\left(\delta_{i} / \alpha_{i}\right)\left[A_{i}^{T} P A_{i}+\left(A_{i}-Q_{a} V_{2 a}^{-1} C_{i}\right)^{T} \hat{P}\left(A_{i}-Q_{a} V_{2 a}^{-1} C_{i}\right)\right] \\
&+E+P D P-P_{a}^{T} R_{2 a}^{-1} P_{a}+\tau_{\perp}^{T} P_{a}^{T} R_{2 a}^{-1} P_{a} \tau_{\perp}, \\
& 0= {[\hat{A}+D(P+\hat{P})] Q+Q[\hat{A}+D(P+\hat{P})]^{T}+V_{1}+\sum^{\prime}\left(\delta_{i} / \alpha_{i}\right)\left[A_{i} Q A_{i}^{T}\right.} \\
&\left.+\left(A_{i}-B_{i} R_{2 a}^{-1} P_{a}\right) \hat{Q}\left(A_{i}-B_{i} R_{2 a}^{-1} P_{a}\right)^{T}\right]-Q_{a} V_{2 a}^{-1} Q_{a}^{T}+\tau_{\perp} Q_{a} V_{2 a}^{-1} Q_{a}^{T} \tau_{\perp}^{T}, \\
& 0=\left(\hat{A}_{Q}+D P\right)^{T} \hat{P}+\hat{P}\left(\hat{A}_{Q}+D P\right)+\hat{P} D \hat{P}+P_{a}^{T} R_{2 a}^{-1} P_{a}-\tau_{\perp}^{T} P_{a}^{T} R_{2 a}^{-1} P_{a} \tau_{\perp}, \\
& 0=\left(\hat{A}_{P}+D P\right) \hat{Q}+\hat{Q}\left(\hat{A}_{P}+D P\right)^{T}+Q_{a} V_{2 a}^{-1} Q_{a}^{T}-\tau_{\perp} Q_{a} V_{2 a}^{-1} Q_{a}^{T} \tau_{\perp}^{T}, \\
& \operatorname{rank} \hat{Q}=\operatorname{rank} \hat{P}=\operatorname{rank} \hat{Q} \hat{P}=n_{c} .
\end{aligned} .
\end{gather*}
$$

Conversely, if there exist $P, Q, \hat{P}, \hat{Q} \in \mathbb{N}^{n}$ satisfying (6.16)-(6.20), then $P$ given by (6.12) satisfies (6.2) or, equivalently, (4.2) with ( $A_{c}, B_{c}, C_{c}$ ) given by (6.13)-(6.15).

Outline of Proof. As discussed in Section 1, we limit the presentation of the proof to the salient details. First note that with the choice of bounds $\Lambda_{i}$, (6.2) becomes

$$
\begin{align*}
0= & \left(\tilde{A}+\frac{1}{2} \sum^{\prime} \delta_{i} \alpha_{i} I_{\tilde{n}}\right)^{T} P+P\left(\tilde{A}+\frac{1}{2} \sum^{\prime} \delta_{i} \alpha_{i} I_{\tilde{n}}+\tilde{R}\right.  \tag{6.21}\\
& +\sum^{\prime}\left(\delta_{i} / \alpha_{i}\right) \tilde{A}_{i}^{T} P \tilde{A}_{i}+\sum^{\prime \prime} \delta_{i}\left(\tilde{E}_{i}^{T} \tilde{E}_{i}+P \tilde{D}_{i} \tilde{D}_{i}^{T} P\right)
\end{align*}
$$

By introducing multipliers $\lambda \in \mathbb{R}, \lambda \geq 0$, and $Q \in \mathbb{R}^{\bar{n} \times \bar{n}}$, a Lagrangian can be defined as

$$
\begin{equation*}
\mathcal{L}\left(P, A_{c}, B_{c}, C_{c}\right) \triangleq \operatorname{tr}[\lambda P \tilde{V}+\mathcal{Q}(\text { RHS of }(6.21))] \tag{6.22}
\end{equation*}
$$

Setting $\partial \mathcal{L} / \partial P=0$ and utilizing (6.8) implies that $\lambda=1$ without loss of generality, $Q \geq 0$, and $Q$ satisfies

$$
\begin{align*}
0= & \left(\tilde{A}+\frac{1}{2} \sum^{\prime} \delta_{i} \alpha_{i} I_{\tilde{n}}+\sum^{\prime \prime} \tilde{D}_{i} \tilde{D}_{i}^{T} \rho\right) Q+Q\left(\tilde{A}+\frac{1}{2} \sum^{\prime} \delta_{i} \alpha_{i} I_{\tilde{n}}+\sum^{\prime \prime} \tilde{D}_{i} \tilde{D}_{i}^{T} \rho\right)^{T}  \tag{6.23}\\
& +\sum^{\prime}\left(\delta_{i} / \alpha_{i}\right) \tilde{A}_{i} Q \tilde{A}_{i}^{T}+\tilde{V} .
\end{align*}
$$

The remainder of the derivation is exactly parallel to the techniques utilized in $[29,36]$. Briefly, the principal steps are as follows:

Step_1. Compute $\partial \mathcal{L} / \partial A_{c}, \partial \mathcal{L} / \partial B_{c}$ and $\partial \mathcal{L} / \partial C_{c}$;
Step 2. Use (6.9) to show that the lower right $n_{c} \times n_{c}$ blocks of $\mathcal{Q}$ and $\mathcal{P}$ are positive definite;

Step 3. Use $\partial \mathcal{L} / \partial A_{c}=0$ to define a projection $\tau$ and new variables $P, Q, \hat{P}, \hat{Q}, G, \Gamma$;
Step 5. Partition (6.21) and (6.23) into six equations (\#1, $\ldots, \# 6$ ) corresponding to the $n \times n, n \times n_{c}$ and $n_{c} \times n_{c}$ blocks of $P$ and $Q$, respectively;

Step 6. Use equations \#2 and \#3 to solve for $A_{c}$; show that equations \#5 and \#6 also yield $A_{c}$; note that with $A_{c}$ now given, equations \#3 and \#6 are superfluous and can be eliminated;

Step 7. Manipulate equations \#1, \#2, \#4 and \#5 to yield (6.16)-(6.19);
Step 8. Show that Steps 5-7 are reversible so that (6.16)-(6.20) are equivalent to (6.2) or, equivalently, (4.2).

By enforcing the strict inequalities $P>0$ and (6.3), solutions of (6.16)-(6.20) guarantee robust stability with a robust performance bound. The following result follows from Theorem 4.1, Theorem 4.2 and the converse of Theorem 6.1.

Theorem 6.2. Suppose there exist $P, Q, \hat{P}, \hat{Q} \in \mathbb{N}^{n}$ satisfying (6.16)-(6.20), and suppose that (6.3) and $P>0$ are satisfied with ( $P, A_{c}, B_{c}, C_{c}$ ) given by (6.12)-(6.15). Then the compensator $A_{c}, B_{c}, C_{c}$ given by (6.13)-(6.15) solves the Robust Stability Problem and the closed-loop performance (3.7) satisfies the bound

$$
\begin{equation*}
J\left(A_{c}, B_{c}, C_{c}\right) \leq \operatorname{tr} P \tilde{V} . \tag{6.24}
\end{equation*}
$$

The following existence result concerns the solvability of (6.16)-(6.20). Let $n_{u}$ denote the dimension of the unstable subspace of the plant dynamics matrix $A$.

Theorem 6.3. Assume $n_{c} \geq n_{u}, R_{1}>0, V_{1}>0$, suppose the nominal plant, i.e., (3.1), (3.2) with $\delta_{i}=0, i=1, \ldots, p$, is stabilizable and detectable and, in addition, is stabilizable by means of an $n_{c}$ th-order strictly proper dynamic compensator (3.3), (3.4). Then there exist $\bar{\delta}_{1}, \ldots, \bar{\delta}_{p}>0$
such that if $\delta_{i} \in\left[0, \bar{\delta}_{i}\right), i=1, \ldots, p$, then (6.16)-(6.20) have a solution $P, Q, \hat{P}, \hat{Q} \in \mathbb{N}^{n}$ for which ( $A_{c}, B_{c}, C_{c}$ ) given by (6.13)-(6.15) solve the robust stability problem with robust performance bound (6.24).

Proof. From Theorem 3.1 of [37] it follows that there exists a solution to (6.16)-(6.20) which stabilizes the nominal plant. By continuity there exists a neighborhood over which robust stability with performance bound (6.24) holds.

Theorem 6.3 is an existence result which guarantees solvability of the sufficiency conditions over a range of parameter uncertainties. The actual range of uncertainty which can be bounded and the conservatism of the performance bound are problem dependent. To this end we now consider a numerical example.

## 7. Illustrative Numerical Example

To demonstrate the above theory we present an illustrative numerical example. The example chosen was originally used in [2] to illustrate the lack of a guaranteed gain margin for LQG controllers. This example was also considered in [35] for a preliminary robustness study and reconsidered in [46] using $\mu$-analysis. Define

$$
\begin{aligned}
n=n_{u}=2, & m=\ell=p=1, \\
A=\left[\begin{array}{ll}
1 & 1 \\
0 & 1
\end{array}\right], & B=\left[\begin{array}{l}
0 \\
1
\end{array}\right], \quad C=\left[\begin{array}{ll}
1 & 0
\end{array}\right], \\
A_{1}=\left[\begin{array}{ll}
0 & 0 \\
0 & 0
\end{array}\right], & B_{1}=\left[\begin{array}{l}
0 \\
1
\end{array}\right], \quad C_{1}=\left[\begin{array}{ll}
0 & 0
\end{array}\right], \\
R_{1}=V_{1}=\left[\begin{array}{ll}
60 & 60 \\
60 & 60
\end{array}\right], & R_{12}=V_{12}=\left[\begin{array}{l}
0 \\
0
\end{array}\right], \quad R_{2}=V_{2}=1 .
\end{aligned}
$$

Note that the system is open-loop unstable and becomes uncontrollable at $\sigma_{1}=-1$. As can be seen using root locus, a strictly proper stabilizing controller must be of at least second order. Hence we consider (6.16)-(6.20) with $n_{c}=n$ and $\tau_{\perp}=0$. Furthermore, we utilize bound (5.3) and thus set $D=E=0$. Using algorithms described in [38-40], controllers were obtained by solving (6.16)(6.20) for $\left(\delta_{1}, \alpha_{1}\right)=(.1,1),(.2,2)$ and (.4,4). As stated previously, these numerical solutions also verify (4.2) with $P$ given by (6.12). Figure 1 compares the guaranteed robust stability region to the "actual" robust stability region. Tinis robust stability region was evaluated assuming constant $\hat{\sigma}_{1}(\cdot)$ although the theory actually guarantees robustness with respect to time-varying uncertainties.

Thus, the gap between these regions may not be a reliable measure of the conservatism of the results. Note, however, that the design approach appears to provide more stability than is guaranteed a priori. Much of this conservatism may be attributable to the desire for a symmetric stability interval so close to an unstabilizable plant perturbation, i.e., $\sigma_{1}=-1$. Nevertheless, the stability design objectives have been met in accordance with Theorem 6.2. Interestingly, the form of the actual stability region mimics the classical 6 dB downward/infinite dB upward gain margin of full-state-feedback LQR controllers ([1]). Thus, this approach appears to provide an alternative to gain-margin recovery techniques ([9]) which address this specialized form of plant uncertainty. Finally, Figure 2 compares guaranteed closed-loop performance to actual closed-loop performance over the guaranteed closed-loop robust stability region. Again the "actual" region was determined for constant $\hat{\sigma}_{1}(\cdot)$. Controller gains are given in Table 1. Finally, it is interesting to note that higher order robust controllers were obtained for this example in [46] using the $\mu$-function approach.

Acknowledgment. I wish to thank Jill M. Straehla for preparing the manuscript versions of this paper, Scott W. Greeley for carrying out the numerical computations, Dr. Wassim M. Haddad for several helpful suggestions, and the reviewers for numerous helpful comments.


Eigure 1


Figure 2

| $\left(\delta_{1}, \alpha_{1}\right)$ | $\left[\begin{array}{cc}A_{c} & B_{c} \\ \hline-14.917 & 1.0 \\ -85.177 & 3.9657\end{array}\right]$ | $\left[\begin{array}{l}15.917 \\ 79.959\end{array}\right]$ | $\left[\begin{array}{ll}-5.2182 & -4.9657\end{array}\right]$ |
| :---: | :---: | :---: | :---: |
| $(.2 .2)$ | $\left[\begin{array}{ll}-17.963 & 1.0 \\ -133.65 & -4.4614\end{array}\right]$ | $\left[\begin{array}{l}18.963 \\ 127.05\end{array}\right]$ | $\left[\begin{array}{ll}-6.6011 & -5.4614\end{array}\right]$ |
| $(.4 .4)$ | $\left[\begin{array}{ll}-47.813 & 1.0 \\ -1087.3 & -6.5463\end{array}\right]$ | $\left[\begin{array}{ll}48.813 \\ 1073.5\end{array}\right]$ | $\left[\begin{array}{ll}-13.766 & -7.5463\end{array}\right]$ |

Table 1

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## APPENDIX J: $H_{60}$ Theory

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Combined $\mathbf{L}_{\mathbf{2}} / \mathbf{H}_{\infty}$ Model Reduction

## by

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#### Abstract

A model-reduction problem is considered which involves both $L_{2}$ (quadratic) and $H_{\infty}$ (worstcase frequency-domain) aspects. Specifically, the goal of the problem is to minimize an $L_{2}$ modelreduction criterion subject to a prespecifieu $H_{\infty}$ constraint on the model-reduction error. The principal result is a sufficient condition for characterizing reduced-order models with bounded $L_{2}$ and $H_{\infty}$ approximation error. The sufficient condition involves a system of modified Riccati equations coupled by an oblique projection, i.e., idempotent matrix. When the $H_{\infty}$ constraint is absent, the sufficient condition specializes to the $L_{2}$ model-reduction result given in Hyland and Bernstein, 1985.


[^29]
## 1. Introduction

One of the most fundamental problems in dynamic systems theory is to approximate a highorder, complex system with a low-order, relatively simpler model. The resulting reduced-order model can then be used to facilitate the analysis of complex systems as well as the design and implementation of feedback controllers and electronic filters. The model-reduction problem thus refiects the fundamental engineering desire for simplicity of implementation and parsimony of hardware.

In view of the practical motivations for the model-reduction problem, it is not surprising that significant effort has been devoted to this problem in recent years. Indeed, there now exists a welldeveloped theoretical foundation for model reduction under a variety of approximation criteria. Expanding on the original work of Adamjan, Arov and Krein, 1971, progress was achieved in Kung and Lin, 1981, Lin and Kung, 1982, Glover, 1984, Latham and Anderson, 1985, Hung and Glover, 1986, Anderson, 1986, Ball and Ran, 1987, and Parker and Anderson, 1987, for the Hankel-norm approximation criterion. Many of the cited works also present bounds for the closely related $H_{\infty}$ approximation error, although the optimal $H_{\infty}$ model-reduction problem remains open. Alternatively, early progress on the model-reduction problem with a quadratic ( $L_{2}$ ) criterion was achieved in Wilson, 1970, and further explored in Hyland and Bernstein, 1985.

Although the Hankel norm, $H_{\infty}$, and $L_{2}$ model-reduction criteria represent distinct approximation objectives, there exist significant connections. For example, it was shown in Wilson, 1985, that for systems which are either single input or single output, the input and output space topologies can be redefined so that the induced norm of the Hankel operator coincides with the $L_{2}$ system norm. In addition, the optimization technique utilized in Wilson, 1970, was reapplied to the Hilbert-Schmidt Hankel operator topology in Wilson, 1988. In recent work, Wilson, 1988a, has shown that for single-input or single-output systems the quadratic model-reduction criterion is actually an induced norm of the convolution operator itself.

In the present paper we attempt a further unification of the $L_{2}$ and $H_{\infty}$ model-reduction objectives. Specifically, we consider an $L_{2}$ model-reduction problem with a constraint on the $H_{\infty}$ approximation error. The underlying idea involves the suitable application of a frequency-domain inequality due to Willems, 1971, which has recently been applied to $H_{\infty}$ control-design problems in Petersen, 1987, Khargonekar, Petersen and Zhou, 1987, and Bernstein and Haddad, 1988. The principle result of the present paper is a sufficient condition which characterizes reduced-order
models satisfying an optimized $L_{2}$ bound as well as a prespecified $H_{\infty}$ bound. The sufficient condition is a direct generalization of the optimal projection approach developed in Hyland and Bernstein, 1985, for the unconstrained $L_{2}$ problem. While the $L_{2}$-optimal reduced-order model was characterized in Hyland and Bernstein, 1985, by means of a coupled system of two modified Lyapunov equations, the $H_{\infty}$-constrained solution in the present paper involves a coupled system consisting of four modified Riccati equations. As in Hyland and Bernstein, 1985, the coupling is due to the presence of an oblique projection (idempotent matrix) which determines the constrained reduced-order model. When the $H_{\infty}$ constraint is sufficiently relaxed, we show that the conditions given herein specialize directly to those given in Hyland and Bernstein, 1985. Although our result gives sufficient conditions for $H_{\infty}$ approximation, we also state hypotheses under which these conditions are also necessary.

Although numerical algorithms were developed in Hyland and Bernstein, 1985, for the "pure" $L_{2}$ problem, computational methods for the $H_{\infty}$-constrained problem are beyond the scope of the present paper. In view of the additional complexity engendered by the $H_{\infty}$ constraint, more sophisticated algorithms appear necessary. Hence computational methods will focus on the homotopic continuation algorithm developed in Richter, 1987, for reduced-order dynamic compensation.

## Notation and Definitions

$\mathbb{R}, \mathbb{R}^{r \times e}, \mathbb{R}^{r}, \mathbb{E}$
$I_{r},()^{T}, 0_{r \times 0}, 0_{r}$
( ) ${ }^{*}$
tr
$\sigma_{\text {max }}(Z)$
$\lambda_{\text {max }}(Z)$
$\|Z\|_{F}$
$\|h(t)\|_{2}$
$\|H(s)\|_{2}$
$\|H(s)\|_{\infty}$
$\mathbb{S}^{r}, \mathbb{N}^{r}, \mathbb{P}^{r}$
$Z_{1} \leq Z_{2}, Z_{1}<Z_{2}$
$n, m, \ell, n_{m}, q, p ; \tilde{n}$
$x, y, y_{m}, x_{m}, \tilde{y}, \tilde{x}$
real numbers, $r \times s$ real matrices, $\mathbb{R}^{r \times 1}$, expected value $r \times r$ identity matrix, transpose, $r \times s$ zero matrix, $0_{r \times r}$ complex conjugate transpose
trace
largest singular value of matrix Z
largest eigenvalue of matrix $Z$ with real spectrum
[ $\left.\operatorname{tr} \mathrm{ZZ}^{*}\right]^{\frac{1}{2}}$ (Frobenius matrix norm)
$\left[\int_{0}^{\infty}\|h(t)\|_{F}^{2} d t\right]^{\frac{1}{2}}$
$\left[\frac{1}{2 \pi} \int_{-\infty}^{\infty}\|H(j \omega)\|_{F}^{2} d \omega\right]^{\frac{1}{2}}$
$\sup _{\omega \in \mathbb{R}} \sigma_{\text {max }}[H(j \omega)]$
$r \times r$ symmetric, nonnegative-definite, positive-definite matrices
$Z_{2}-Z_{1} \in \mathbb{N}^{r}, Z_{2}-Z_{1} \in \mathbb{P}^{r}, \quad Z_{1}, Z_{2} \in \mathbb{S}^{r}$
positive integers; $n+n_{m}$
$n, \ell, \ell, n_{m}, \ell, \tilde{n}$-dimensional vectors

$$
\begin{aligned}
& \tilde{y}, \tilde{x} \\
& A, B, C \\
& D, E \\
& A_{m}, B_{m}, \\
& \tilde{A}, \tilde{B}, \tilde{C} \\
& \tilde{D}, \tilde{E} \\
& R \\
& w(\cdot) \\
& V \\
& \tilde{R}, \tilde{V} \\
& \boldsymbol{\gamma}
\end{aligned}
$$

$$
y-y_{m}, \quad\left[\begin{array}{c}
x \\
x_{m}
\end{array}\right]
$$

$$
\begin{array}{ll}
A, B, C & n \times n, n \times m, \ell \times n \text { matrices } \\
D, E & m \times p, q \times \ell \text { matrices } \\
A_{m}, B_{m}, C_{m} & n_{m} \times n_{m}, n_{m} \times m, \ell \times n_{m} \text { matrices }
\end{array}
$$

$$
\left[\begin{array}{cc}
A & 0 \\
0 & A_{m}
\end{array}\right],\left[\begin{array}{c}
B \\
B_{m}
\end{array}\right], \quad\left[\begin{array}{ll}
C & -C_{m}
\end{array}\right]
$$

$$
\tilde{B} D=\left[\begin{array}{c}
B D \\
B_{m} D
\end{array}\right], \quad E \tilde{C}=\left[\begin{array}{ll}
E C & -E C_{m}
\end{array}\right]
$$

$$
E^{T} E, \text { model-reduction error-weighting matrix in } \mathbb{P}^{\ell}
$$

p-dimensional standard white noise process
intensity of $D w(\cdot), V=D D^{T} \in \mathbb{P}^{m}$
$\left[\begin{array}{ll}C^{T} R C & -C^{T} R C_{m} \\ -C_{m}^{T} R C & C_{m}^{T} R C_{m}\end{array}\right], \quad\left[\begin{array}{cc}B V B^{T} & B V B_{m}^{T} \\ B_{m} V B^{T} & B_{m} V B_{m}^{T}\end{array}\right]$
positive constant

## 2. Statement of the Problem

In this section we introduce the model-reduction problem with constrained $H_{\infty}$ norm of the model-reduction error. Specifically, we constrain the transfer function of the reduced-order model to lie within a specified $H_{\infty}$ radius of the original system. In this paper we assume that the full-order model is asymptotically stable, i.e., the matrix $A$ is asymptotically stable.
$\mathbf{H}_{\infty}$-Constrained $\mathbf{L}_{2}$ Model-Reduction Problem. Given the $\boldsymbol{n}$ th-order controllable and observable model

$$
\begin{gather*}
\dot{x}(t)=A x(t)+B D w(t)  \tag{2.1}\\
y(t)=C x(t) \tag{2.2}
\end{gather*}
$$

where $t \in[0, \infty)$, determine an $n_{m}$ th-order model

$$
\begin{gather*}
\dot{x}_{m}(t)=A_{m} x_{m}(t)+B_{m} D w(t)  \tag{2.3}\\
y_{m}(t)=C_{m} x_{m}(t) \tag{2.4}
\end{gather*}
$$

which satisfies the following criteria:
(i) $A_{m}$ is asymptotically stable;
(ii) the transfer function of the reduced-order model is within a radius- $\gamma H_{\infty}$ neighborhood of the full-order model, i.e.,

$$
\begin{equation*}
\left\|H(s)-H_{m}(s)\right\|_{\infty} \leq \gamma \tag{2.5}
\end{equation*}
$$

where

$$
\begin{equation*}
H(s) \triangleq E C\left(s I_{n}-A\right)^{-1} B D, \quad H_{m}(s) \triangleq E C_{m}\left(s I_{n_{m}}-A_{m}\right)^{-1} B_{m} D, \tag{2.6}
\end{equation*}
$$

and $\gamma>0$ is a given constant; and
(iii) the $L_{2}$ model-reduction criterion

$$
\begin{equation*}
J\left(A_{m}, B_{m}, C_{m}\right) \triangleq \lim _{t \rightarrow \infty} \mathbb{E}\left\{\left[y(t)-y_{m}(t)\right]^{T} R\left[y(t)-y_{m}(t)\right]\right\} \tag{2.7}
\end{equation*}
$$

is minimized.

Note that the full- and reduced-order systems (2.1)-(2.4) can be written as a single augmented system

$$
\begin{equation*}
\dot{\tilde{x}}(t)=\tilde{A} \tilde{x}(t)+\tilde{D} w(t), \quad t \in[0, \infty) \tag{2.8}
\end{equation*}
$$

so that the $q \times p$ transfer function from $w(t)$ to $E \tilde{y}(t)=\tilde{E} \tilde{x}(t)$ is

$$
\begin{equation*}
\tilde{H}(s)=\tilde{E}\left(s I_{\tilde{n}}-\tilde{A}\right)^{-1} \tilde{D} \tag{2.9}
\end{equation*}
$$

and (2.7) can be written as

$$
\begin{equation*}
J\left(A_{m}, B_{m}, C_{m}\right)=\lim _{t \rightarrow \infty} \mathbb{E}\left\{[E \tilde{y}(t)]^{T}[E \tilde{y}(t)]\right\}=\lim _{t \rightarrow \infty} \mathbb{E}\left[\tilde{x}^{T}(t) \tilde{R} \tilde{x}(t)\right] \tag{2.10}
\end{equation*}
$$

Before continuing it is useful to note that if $A_{m}$ is asymptotically stable then the $L_{2}$ modelreduction criterion (2.7) is given by

$$
\begin{equation*}
J\left(A_{m}, B_{m}, C_{m}\right)=\operatorname{tr} \tilde{Q} \tilde{R} \tag{2.11}
\end{equation*}
$$

where the steady-state covariance

$$
\begin{equation*}
\tilde{Q} \triangleq \lim _{t \rightarrow \infty} \mathbb{E}\left[\tilde{x}(t) \tilde{x}^{T}(t)\right] \tag{2.12}
\end{equation*}
$$

satisfies the augmented Lyapunov equation

$$
\begin{equation*}
0=\tilde{A} \tilde{Q}+\tilde{Q} \tilde{A}^{T}+\tilde{V} \tag{2.13}
\end{equation*}
$$

Using (2.11) and (2.13) it can be shown that the $L_{2}$ criterion (2.7) is an approximation measure involving the full- and reduced-order impulse responses with respect to an $L_{2}$ norm.

Proposition 2.1. The $L_{2}$ model-reduction criterion (2.11) can be written as

$$
\begin{equation*}
J\left(A_{m}, B_{m}, C_{m}\right)=\int_{0}^{\infty}\left\|E C e^{A t} B D-E C_{m} e^{A_{m} t} B_{m} D\right\|_{F}^{2} d t \tag{2.14a}
\end{equation*}
$$

or, equivalently,

$$
\begin{equation*}
J\left(A_{m}, B_{m}, C_{m}\right)=\frac{1}{2 \pi} \int_{-\infty}^{\infty}\left\|H(j \omega)-H_{m}(j \omega)\right\|_{F}^{2} d \omega \tag{2.14b}
\end{equation*}
$$

Proof. It need only be noted that (2.11) is equivalent to

$$
\begin{aligned}
\operatorname{tr} \int_{0}^{\infty} e^{\tilde{A} t} \tilde{V} e^{\tilde{A}^{T} t} d t \tilde{R} & =\operatorname{tr} \int_{0}^{\infty} \tilde{E} e^{\tilde{A} t} \tilde{D} \tilde{D}^{T} e^{\tilde{A}^{T} t} \tilde{E}^{T} d t \\
& =\operatorname{tr} \int_{0}^{\infty}\left(\tilde{E} e^{\tilde{A} t} \tilde{D}\right)\left(\tilde{E} e^{\tilde{A} t} \tilde{D}\right)^{T} d t \\
& =\int_{0}^{\infty}\left\|\tilde{E} e^{\bar{A} t} \tilde{D}\right\|_{F}^{2} d t
\end{aligned}
$$

which is equivalent to (2.14a). Finally, (2.14b) follows from Plancherel's Theorem.
The key step in enforcing (2.5) is to replace the algebraic Lyapunov equation (2.13) by an algebraic Riccati equation. Justification for this technique is provided by the following result.

Lemma 2.1. Let $\left(A_{m}, B_{m}, C_{m}\right)$ be given and assume there exists $Q \in \mathbb{R}^{\bar{n} \times \bar{n}}$ satisfying

$$
\begin{equation*}
Q \in \mathbb{N}^{\bar{n}} \tag{2.15}
\end{equation*}
$$

and

$$
\begin{equation*}
0=\tilde{A} Q+Q \tilde{A}^{T}+\gamma^{-2} Q \tilde{R} Q+\tilde{V} . \tag{2.16}
\end{equation*}
$$

Then

$$
\begin{equation*}
\left(\tilde{A},\left[\gamma^{-2} \odot \tilde{R} Q+\tilde{V}\right]^{\frac{1}{2}}\right) \text { is stabilizable } \tag{2.17}
\end{equation*}
$$

if and only if

$$
\begin{equation*}
A_{m} \text { is asymptotically stable. } \tag{2.18}
\end{equation*}
$$

Furthermore, in this case,

$$
\begin{equation*}
\left\|H(s)-H_{m}(s)\right\|_{\infty} \leq \gamma \tag{2.19}
\end{equation*}
$$

$$
\begin{equation*}
\tilde{Q} \leq \mathbb{O}, \tag{2.20}
\end{equation*}
$$

and

$$
\begin{equation*}
J\left(A_{m}, B_{m}, C_{m}\right) \leq J\left(A_{m}, B_{m}, C_{m}, \mathcal{Q}\right) \tag{2.21}
\end{equation*}
$$

where

$$
\begin{equation*}
J\left(A_{m}, B_{m}, C_{m}, Q\right) \triangleq \operatorname{tr} Q \tilde{R} . \tag{2.22}
\end{equation*}
$$

Proof. Using the assumed existence of a nonnegative-definite solution to (2.16) and the stabilizability condition (2.17), it follows from the dual of Lemma 12.2 of Wonham, 1979, that $\tilde{A}$ is asymptotically stable. Since $\tilde{A}$ is block diagonal, $A_{m}$ is also asymptotically stable. Conversely, since $A$ is assumed to be asymptotically stable, (2.18) implies (2.17). The proof of (2.19) follows from a standard manipulation of (2.16); for details see Lemma 1 of Willems, 1971. To prove (2.20) subtract (2.13) from (2.16) to obtain

$$
\begin{equation*}
0=\tilde{A}(Q-\tilde{Q})+(Q-\tilde{Q}) \tilde{A}^{T}+\gamma^{-2} Q \tilde{R} Q \tag{2.23}
\end{equation*}
$$

which, since $\bar{A}$ is asymptotically stable, is equivalent to

$$
\begin{equation*}
\mathcal{Q}-\tilde{Q}=\int_{0}^{\infty} e^{\tilde{A} t}\left[\gamma^{-2} Q \tilde{R} Q\right] e^{\tilde{A}^{T} t} d t \geq 0 \tag{2.24}
\end{equation*}
$$

Finally, (2.21) follows immediately from (2.20).
Lemma 2.1 shows that the $H_{\infty}$ constraint is automatically enforced when a nonnegative-definite solution to (2.16) is known to exist. Furthermore, the solution $\mathcal{Q}$ provides an upper bound for the actual state covariance $\tilde{Q}$ along with a bound on the $L_{2}$ model-reduction criterion. Next, we present a partial converse of Lemma 2.1 which guarantees the existence of a nonnegative-definite solution to (2.16) when (2.19) is satisfied.

Lemma 2.2. Let ( $A_{m}, B_{m}, C_{m}$ ) be given, suppose $\tilde{A}$ is asymptotically stable, and assume the $H_{\infty}$ approximation constraint (2.19) is satisfied. Then there exists a unique nonnegative-definite solution $Q$ satisfying (2.16). Furthermore, $\left(\tilde{A}+\gamma^{-2} Q \tilde{R}, V^{\frac{1}{2}}\right)$ is stabilizable if and only if $\tilde{A}+\gamma^{-2} Q \tilde{R}$ is asymptotically stable.

Proof. The result is an immediate consequence of Theorems 3 and 2, pp. 150 and 167 of Brockett, 1970, and the dual of Lemma 12.2 of Wonham, 1979.

Finally, we show that the quadratic term $\boldsymbol{\gamma}^{-2} Q \tilde{R} Q$ in (2.16) also constrains the Hankel norm of the approximation error $E \tilde{y}$ when $Q$ is positive definite. To show this let $\tilde{P} \in \mathbb{N}^{\tilde{n}}$ be the
observability Gramian for the augmented system ( $\tilde{A}, \tilde{D} ; \tilde{E})$ which satisfies

$$
\begin{equation*}
0=\tilde{A}^{T} \tilde{P}+\tilde{P} \tilde{A}+\tilde{R} \tag{2.25}
\end{equation*}
$$

Furthermore, note that $\tilde{Q}$ satisfying (2.13) is the dual controllability Gramian.
Proposition 2.2. Let ( $A_{m}, B_{m}, C_{m}$ ) be given and assume there exists $\mathcal{Q} \in \mathbb{P}^{n}$ satisfying (2.16) and (2.17) or, equivalently, (2.18). Then

$$
\begin{equation*}
\lambda_{\text {max }}^{\frac{1}{2}}(\tilde{P} \tilde{Q}) \leq \gamma . \tag{2.26}
\end{equation*}
$$

Proof. Since $Q$ is invertible, (2.16) implies

$$
\begin{equation*}
0=\gamma^{2} \tilde{A}^{T} Q^{-1}+\gamma^{2} Q^{-1} \tilde{A}+\gamma^{2} Q^{-1} \tilde{V} Q^{-1}+\tilde{R} \tag{2.27}
\end{equation*}
$$

Next, subtract (2.25) from (2.27) to obtain

$$
\begin{equation*}
0=\tilde{A}^{T}\left(\gamma^{2} Q^{-1}-\tilde{P}\right)+\left(\gamma^{2} Q^{-1}-\tilde{P}\right) \tilde{A}+\gamma^{2} Q^{-1} \tilde{V} Q^{-1} \tag{2.28}
\end{equation*}
$$

which, since $\tilde{A}$ is asymptotically stable, is equivalent to

$$
\begin{equation*}
\gamma^{2} Q^{-1}-\tilde{P}=\int_{0}^{\infty} e^{\tilde{A}^{T} t}\left[\gamma^{2} Q^{-1} \tilde{V} Q^{-1}\right] e^{\tilde{A} t} d t \geq 0 \tag{2.29}
\end{equation*}
$$

Thus, (2.29) implies $\tilde{P} \leq \gamma^{2} \mathcal{Q}^{-1}$ or, equivalently, $\mathcal{Q}^{\frac{1}{2}} \tilde{P} Q^{\frac{1}{2}} \leq \gamma^{2} I_{\tilde{n}}$. Hence, $\lambda_{\text {max }}^{\frac{1}{2}}(\tilde{P} O) \leq \gamma$. Finally, (2.26) follows immediately from (2.20).
3. The Auxiliary Minimization Problem and Necessary Conditions for Optimality

As discussed in the previous section, the replacement of (2.13) by (2.16) enforces the $H_{\infty}$ approximation constraint between the full- and reduced-order systems and results in an upper bound for the $L_{2}$ model-reduction criterion. That is, if (2.16) is solvable then the reduced-order model $\left(A_{m}, B_{m}, C_{m}\right)$ satisfies the $H_{\infty}$ approximation constraint (2.5) while the actual $L_{2}$ modelreduction criterion is guaranteed to be no worse than the bound given by $J\left(A_{m}, \dot{\Delta}_{m}, C_{m}, \ell\right)$. Hence, $J\left(A_{m}, B_{m}, C_{m}, \mathcal{Q}\right)$ can be interpreted as an auxiliary cost which leads to the following mathematical programming problem.

Auxiliary Minimization Problem. Determine ( $A_{m}, B_{m}, C_{m}, Q$ ) which minimizes $J\left(A_{m}, B_{m}, C_{m}, \mathcal{Q}\right)$ subject to (2.15) and (2.16).

It follows from Lemma 2.1 that the satisfaction of (2.15)-(2.17) leads to 1) $A_{m}$ stable; 2) a bound on the $H_{\infty}$ distance between the full-order and reduced-order systems; and 3) an upper bound for the $L_{2}$ model-reduction criterion. Hence it remains to determine ( $A_{m}, B_{m}, C_{m}$ ) which minimizes $J\left(A_{m}, B_{m}, C_{m}, \mathcal{Q}\right)$ and thus provides an optimized bound for the actual $L_{2}$ criterion $J\left(A_{m}, B_{m}, C_{m}\right)$. Rigorous derivation of the necessary conditions for the Auxiliary Minimization Problem requires additional technical assumptions. Specifically, we restrict ( $A_{m}, B_{m}, C_{m}, \mathcal{Q}$ ) to the open set

$$
\begin{array}{r}
S \triangleq\left\{\left(A_{m}, B_{m}, C_{m}, Q\right): Q \in \mathbb{P}^{\tilde{n}}, \tilde{A}+\gamma^{-2} Q \tilde{R} \text { is asymptotically stable },\right. \\
\text { and } \left.\left(A_{m}, B_{m}, C_{m}\right) \text { is controllable and observable }\right\} . \tag{3.1}
\end{array}
$$

Remark 3.1. The set $S$ constitutes sufficient conditions under which the Lagrange multiplier technique is applicable to the Auxiliary Minimization Problem. Specifically, the requirement that $Q$ be positive definite replaces (2.15) by an open set constraint, the stability of $\tilde{A}+\gamma^{-2} Q \tilde{R}$ serves as a normality condition, and ( $A_{m}, B_{m}, C_{m}$ ) minimal is a nondegeneracy condition.

The following Lemma is needed for the statement of the main result.
Lemma 3.1. Let $\hat{Q}, \hat{P} \in \mathbb{N}^{n}$ and suppose rank $\hat{Q} \hat{P}=n_{m}$. ${ }^{\text {r' }}$ en there exist $n_{m} \times n G, \Gamma$ and $n_{m} \times n_{m}$ invertible $M$, unique except for a change of basis in $\mathbb{R}^{n_{m}}$, such that

$$
\begin{gather*}
\hat{Q} \hat{P}=G^{\boldsymbol{T}} M \Gamma  \tag{3.2}\\
\Gamma G^{\boldsymbol{T}}=I_{n_{m}} \tag{3.3}
\end{gather*}
$$

Furthermore, the $n \times n$ matrices

$$
\begin{gather*}
\tau \triangleq G^{T} \Gamma  \tag{3.4}\\
\tau_{\perp} \triangleq I_{n}-\tau \tag{3.5}
\end{gather*}
$$

are idempotent and have rank $n_{m}$ and $n-n_{m}$, respectively. If, in addition,

$$
\begin{equation*}
\operatorname{rank} \hat{Q}=\operatorname{rank} \hat{P}=n_{m} \tag{3.6}
\end{equation*}
$$

then

$$
\begin{equation*}
\hat{Q}=\tau \hat{Q}, \quad \hat{P}=\hat{P} r . \tag{3.7}
\end{equation*}
$$

Finally, if $P \in \mathbb{N}^{n}$ then the inverse

$$
\begin{equation*}
S \triangleq\left(I_{n}+\gamma^{-2} \hat{Q} P\right)^{-1} \tag{3.9}
\end{equation*}
$$

exists.
Proof. Conditions (3.2)-(3.8) are a direct consequence of Theorem 6.2.5 of Rao and Mitra, 1971. To prove that the inverse in (3.9) exists, note that since the eigenvalues of $\hat{Q} P$ coincide with the eigenvalues of the nonnegative-definite matrix $P^{\frac{1}{2}} \hat{Q} P^{\frac{1}{2}}$, it follows that $\hat{Q} P$ has nonnegative eigenvalues. Thus, the eigenvalues of $I_{n}+\gamma^{-2} \hat{Q} P$ are all greater than one so that the above inverse exists.

Finally, for convenience define

$$
\Sigma \triangleq B V B^{T}, \quad \bar{\Sigma} \triangleq C^{T} R C .
$$

Theorem 3.1. If $\left(A_{m}, B_{m}, C_{m}, \mathcal{Q}\right) \in S$ solves the Auxiliary Minimization Problem then there exist $Q, P, \hat{Q}, \hat{P} \in \mathbb{N}^{n}$ such that

$$
\begin{align*}
A_{m} & =\Gamma\left(A-\gamma^{-4} Q \bar{\Sigma} Q P S\right) G^{T},  \tag{3.10}\\
B_{m} & =\Gamma B  \tag{3.11}\\
C_{m} & =C\left(I_{n}+\gamma^{-2} Q P S\right) G^{T},  \tag{3.12}\\
\mathcal{Q} & =\left[\begin{array}{cc}
Q+\hat{Q} & \hat{Q} \Gamma^{T} \\
\Gamma \hat{Q} & \Gamma \hat{Q} \Gamma^{T}
\end{array}\right], \tag{3.13}
\end{align*}
$$

and such that $Q, P, \hat{Q}, \hat{P}$ satisfy

$$
\begin{align*}
0= & A Q+Q A^{T}+\gamma^{-2} Q \bar{\Sigma} Q+\tau_{\perp} \Sigma \tau_{\perp}^{T},  \tag{3.14}\\
0= & A^{T} P+P A-\gamma^{-4} S^{T} P Q \bar{\Sigma} Q P S+\tau_{\perp}^{T}\left(I_{n}+\gamma^{-2} Q P S\right)^{T} \bar{\Sigma}\left(I_{n}+\gamma^{-2} Q P S\right) \tau_{\perp},  \tag{3.15}\\
0= & \left(A-\gamma^{-4} Q \bar{\Sigma} Q P S\right) \hat{Q}+\hat{Q}\left(A-\gamma^{-4} Q \bar{\Sigma} Q P S\right)^{T}+\gamma^{-4} \hat{Q} S^{T} P Q \bar{Q} Q P S \hat{Q}+\Sigma-\tau_{\perp} \Sigma \tau_{\perp}^{T},  \tag{3.16}\\
0= & \left(A+\gamma^{-2} Q \bar{\Sigma}\right)^{T} \hat{P}+\hat{P}\left(A+\gamma^{-2} Q \bar{\Sigma}\right)+\left(I_{n}+\gamma^{-2} Q P S\right)^{T} \bar{\Sigma}\left(I_{n}+\gamma^{-2} Q P S\right) \\
& \quad-\tau_{\perp}^{T}\left(I_{n}+\gamma^{-2} Q P S\right)^{T} \bar{\Sigma}\left(I_{n}+\gamma^{-2} Q P S\right) \tau_{\perp},  \tag{3.17}\\
& \operatorname{rank} \hat{Q}=\operatorname{rank} \hat{P}=\operatorname{rank} \hat{Q} \hat{P}=n_{m} . \tag{3.18}
\end{align*}
$$

Furthermore, the auxiliary cost is given by

$$
\begin{equation*}
J\left(A_{m}, B_{m}, C_{m}, Q\right)=\operatorname{tr} \tilde{\Sigma}\left(Q+\gamma^{-4} Q P S \hat{Q} S^{T} P Q\right) . \tag{3.19}
\end{equation*}
$$

Conversely, if there exist $Q, P, \hat{Q}, \hat{P} \in \mathbb{N}^{n}$ satisfying (3.14)-(3.18), then ( $A_{m}, B_{m}, C_{m}, \mathcal{Q}$ ) given by (3.10)-(3.13) satisfy (2.15) and (2.16) with the auxiliary cost (2.22) given by (3.19).

Proof. See Appendix A.
Remark 3.1. Theorem 3.1 presents necessary conditions for the Auxiliary Minimization Problem which explicitly synthesize extremal reduced-order models ( $A_{m}, B_{m}, C_{m}$ ). As a check of these conditions, consider the extreme case $n_{m}=n$. Then $G=\Gamma^{-1}$ and thus, without loss of generality, $G=\Gamma=r=I_{n}$ and $\tau_{\perp}=0$. Furthermore, (3.14) implies that $Q=0$ and (3.15) implies that $P=0$. Hence the $H_{\infty}$-constrained full-order model is given (as expected) by ( $A, B, C$ ) regardless of $\boldsymbol{\gamma}$. Furthermore, note that $\mathbb{Q}$ given by (3.13) becomes

$$
Q=\left[\begin{array}{ll}
\hat{Q} & \hat{Q}  \tag{3.20}\\
\hat{Q} & \hat{Q}
\end{array}\right]
$$

so that the quadratic term $\gamma^{-2} Q \tilde{R} Q$ in (2.16) vanishes. Thus (2.16) reduces to (2.13) so that $Q$ coincides with the controllability Gramian $\tilde{Q}$. If, alternatively, the reduced-order constraint is retained but the transfer function approximation constraint (2.5) is sufficiently relaxed, i.e., $\boldsymbol{\gamma} \rightarrow \infty$, then $S=I_{n}$ so that the reduced-order model (3.10)-(3.12) is given by $\left(A_{m}, B_{m}, C_{m}\right)=$ ( $\Gamma A G^{T}, \Gamma B, C G^{T}$ ). In this case (3.14) and (3.15) are superfluous and (3.16) and (3.17) reduce to the optimal projection equations obtained by Hyland and Bernstein, 1985, for the unconstrained $L_{2}$ problem.

## 4. Sufficient Conditions for Combined $\mathrm{L}_{\mathbf{2}} / \mathrm{H}_{\infty}$ Approximation

In this section we combine Lemma 2.1 with the converse of Theorem 3.1 to obtain our main result guaranteeing constrained $H_{\infty}$ approximation along with an optimized $L_{2}$ modej-reduction bound.

Theorem 4.1. Suppose there exist $Q, P, \dot{Q}, \hat{P} \in \mathbb{N}^{n}$ satisfying (3.14)-(3.18) and let $\left(A_{m}, B_{m}, C_{m}, Q\right)$ be given by (3.10)-(3.13). Then $\left(\tilde{A},\left|\mathcal{r}^{-2} Q \tilde{R} Q+\bar{V}\right|^{\frac{1}{2}}\right)$ is stabilizable if and only if $A_{m}$ is asymptotically stable. In this case, the reduced-order transfer function $H_{m}(\rho)$ satisfies the $H_{\infty}$ approximation constraint

$$
\begin{equation*}
\left\|म(s)-H_{m}(s)\right\|_{\infty} \leq \gamma \tag{4.1}
\end{equation*}
$$

and the $L_{2}$ approximation bound

$$
\begin{equation*}
\left\|H(s)-H_{m}(s)\right\|_{2} \leq\left[\operatorname{tr} \dot{\mathcal{L}}\left(Q+\gamma^{-4} Q P S \hat{Q} S^{T} P Q\right)\right]^{\frac{1}{3}} . \tag{4.2}
\end{equation*}
$$

Proof. The converse portion of Theorem 3.1 implies that $Q$ given by (3.13) satisfies (2.15) and (2.16) with auxiliary cost given by (3.19). It now follows from Lemma 2.1 that the stabilizability
condition (2.17) is equivalent to the asymptotic stability of $A_{m}$, the $H_{\infty}$ approximation condition (2.19) holds, and the $L_{3}$ model-reduction criterion satisfies the bound (2.21) which is equivalent to (4.2).

In applying Theorem 4.1 the principal imue concerns conditions on the problem data under which the coupled Riccati equations (3.14)-(3.17) possess nonnegativo-definite molutions. Clearly, for $\gamma$ sufficiently large, (3.14)-(3.17) approximate the "pure" $L_{2}$ solution oblained in Hyland and Bernstein, 1985. In practice, we would aumerically solve (3.14)-(3.17) for succesaively smaller values of $\gamma$ until solutions are no longer obtainable. The important case of interest, however, involvea amall $\boldsymbol{\gamma}$ so that accurate $H_{\infty}$ approximation is enforced. Thus, if (4.1) can be atisfied for a given $\gamma>0$ by a clase of reduced-order models, it is of interest to know whecher one such reduced-order model can be obtained by solving (3.14)-(3.17). Lemma 2.2 guarantoes that ( 2.16 ) pomesses a solution for any model satisfying (4.1). Thus our sufficient conditions will also be necessary so long sa the Auxiliary Minimization Problem possesses at least one extremal over $S$. When this is the case we have the following immediate result.

Propoaition 4.1. Let $\boldsymbol{\gamma}^{*}$ denote the infimum of $\left\|H(a)-H_{m}(s)\right\|_{\infty}$ over all asymptotically stable reduced-order models and suppose that the Auxiliary Minimization Problem has a solution for all $\boldsymbol{\gamma}>\boldsymbol{\gamma}^{\bullet}$. Then for all $\boldsymbol{\gamma}>\boldsymbol{\gamma}^{\boldsymbol{*}}$ there exist $Q, P, \dot{Q}, \dot{P} \in \mathbb{N}^{n}$ satisfying (3.14)-(3.17).

Remark 4.1. As in Hyland and Bernatein, 1985, it can be expected that (3.14)-(3.17) possess multiple solutions. Theorem 4.1 guarantees, however, that the bounds (4.1) and (4.2) are enforced for all such extremals obtained by solving (3.14)-(3.17).

## Appendix A: Proof of Theorem 3.1

To optimize (2.22) over the open eet $S$ subject to the constraint (2.16), form the Lagrangian

$$
\begin{equation*}
\mathcal{C}\left(A_{m}, B_{m}, C_{m}, Q, P, \lambda\right) \triangleq \operatorname{tr}\left\{\lambda Q \tilde{R}+\left|\bar{Q} Q+Q \dot{A}^{T}+\gamma^{-2} Q \dot{R} Q+\bar{V}\right| P\right\} \tag{A.1}
\end{equation*}
$$

where the Lagrange multipliers $\lambda \geq 0$ and $\rho \in \mathbb{R}^{\boldsymbol{n} \times \lambda}$ are not both zero. We thus obtain

$$
\begin{equation*}
\frac{\partial \mathcal{L}}{\partial Q}=\left(\bar{A}+\gamma^{-2} Q \tilde{R}\right)^{T} \rho+\rho\left(\bar{A}+\gamma^{-2} Q \tilde{R}\right)+\lambda \tilde{R} \tag{A.2}
\end{equation*}
$$

Setting $\frac{9}{\partial 8}=0$ yields

$$
\begin{equation*}
0=\left(\bar{A}+\gamma^{-2} Q \bar{R}\right)^{T} P+P\left(\bar{A}+\gamma^{-2} Q \bar{R}\right)+\lambda \bar{R} . \tag{A.3}
\end{equation*}
$$

Since $\tilde{A}+\gamma^{-2} Q \tilde{R}$ is assumed to be atable, $\lambda=0$ implies $P=0$. Hence, it can be asumed withous loss of generality that $\lambda=1$. Furthermore, $P$ is nonnegative defnite.

Now partition $\bar{n} \times \bar{n} Q, P$, into $n \times n, n \times n_{m}$, and $n_{m} \times n_{m}$ subblocks as

$$
Q=\left[\begin{array}{ll}
Q_{1} & Q_{12} \\
Q_{12}^{T} & Q_{2}
\end{array}\right], \quad \mathcal{P}=\left[\begin{array}{cc}
P_{1} & P_{12} \\
P_{12}^{T} & P_{2}
\end{array}\right],
$$

and for notational convenience define

$$
P Q=\left[\begin{array}{ll}
Z_{1} & Z_{12} \\
Z_{21} & Z_{2}
\end{array}\right]
$$

where

$$
\begin{array}{ll}
Z_{1} \triangleq P_{1} Q_{1}+P_{12} Q_{12}^{T}, & Z_{12} \triangleq P_{1} Q_{12}+P_{12} Q_{2} \\
Z_{21} \triangleq P_{12}^{T} Q_{1}+P_{2} Q_{12}^{T}, & Z_{2} \triangleq P_{12}^{T} Q_{12}+P_{2} Q_{2}
\end{array}
$$

Thus, with $\lambda=1$ the stationarity conditions are given by

$$
\begin{align*}
& \frac{\partial \mathcal{L}}{\partial Q}=\left(\tilde{A}+\gamma^{-2} Q \tilde{R}\right)^{T} D+P\left(\tilde{A}+\gamma^{-2} Q \tilde{R}\right)+\tilde{R}=0,  \tag{A.4}\\
& \frac{\partial \mathcal{L}}{\partial A_{m}}=Z_{2}=0,  \tag{A.5}\\
& \frac{\partial \mathcal{L}}{\partial B_{m}}=P_{12}^{T} B V+P_{2} B_{m} V=0,  \tag{A.6}\\
& \frac{\partial \mathcal{L}}{\partial C_{m}}=2 R C_{m} Q_{2}+2 \gamma^{-2} R C_{m} Z_{12}^{T} Q_{12}-2 R C Q_{12}-\gamma^{-2} R C Z_{1}^{T} Q_{12} \\
& \quad-\gamma^{-2} R C Q_{1} Z_{12}-\gamma^{-2} R C Z_{21}^{T} Q_{2}=0 . \tag{A.7}
\end{align*}
$$

Expanding (2.16) and (A.4) yields

$$
\begin{align*}
0= & A Q_{1}+Q_{1} A^{T}+\gamma^{-2}\left(Q_{1} C^{T}-Q_{12} C_{m}^{T}\right) R\left(Q_{1} C^{T}-Q_{12} C_{m}^{T}\right)^{T}+B V B^{T},  \tag{A.8}\\
0= & A Q_{12}+Q_{12} A_{m}^{T}+\gamma^{-2} Q_{1} C^{T} R C Q_{12}-\gamma^{-3} Q_{12} C_{m}^{T} R C Q_{12}-\gamma^{-2} Q_{1} C^{T} R C_{m} Q_{2} \\
& \quad+\gamma^{-2} Q_{12} C_{m}^{T} R C_{m} Q_{2},  \tag{A.9}\\
0= & A_{m} Q_{2}+Q_{3} A_{m}^{T}+\gamma^{-2}\left(Q_{12}^{T} C^{T}-Q_{2} C_{m}^{T}\right) R\left(Q_{12}^{T} C^{T}-Q_{2} C_{m}^{T}\right)^{T}+B_{m} V B_{m}^{T},  \tag{A.10}\\
0= & A^{T} P_{1}+P_{1} A+\gamma^{-2} C^{T} R C Z_{1}^{T}-\gamma^{-3} C^{T} R C_{m} Z_{12}^{T}+\gamma^{-2} Z_{1} C^{T} R C \\
& \quad-\gamma^{-2} Z_{12} C_{m}^{T} R C+C^{T} R C,  \tag{A.11}\\
0= & A^{T} P_{12}+P_{12} A_{m}+\gamma^{-2} C^{T} R C Z_{12}^{T}-\gamma^{-2} Z_{1} C^{T} R C_{m}+\gamma^{-2} Z_{12} C_{m}^{T} R C_{m}-C^{T} R C_{m},  \tag{A.12}\\
0= & A_{m}^{T} P_{2}+P_{2} A_{m}-\gamma^{-2} C_{m}^{T} R C Z_{11}^{T}-\gamma^{-2} Z_{21} C^{T} R C_{m}+C_{m}^{T} R C_{m} . \tag{A.13}
\end{align*}
$$

Now define the $n \times n$ matrices

$$
\begin{gathered}
Q \triangleq Q_{1}-Q_{13} Q_{2}^{-1} Q_{12}^{T}, \quad P \triangleq P_{1}-P_{12} P_{2}^{-1} P_{12}^{T} \\
\dot{Q} \triangleq Q_{12} Q_{2}^{-1} Q_{12}^{T}, \quad \hat{P} \triangleq P_{12} P_{2}^{-1} P_{12}^{T}, \\
T \triangleq-Q_{12} Q_{2}^{-1} P_{2}^{-1} P_{12}^{T}
\end{gathered}
$$

and the $n_{m} \times n_{1} n_{m} \times n_{m}$ and $n_{m} \times n$ matrices

$$
G \triangleq Q_{2}^{-1} Q_{12}^{T}, \quad M \triangleq Q_{2} P_{2}, \quad \Gamma \triangleq-P_{2}^{-1} P_{12}^{T}
$$

The existence of $Q_{2}^{-1}$ and $P_{2}^{-1}$ follows from the fact that ( $A_{m}, B_{m}, C_{m}$ ) is minimal. See Bernstein and Haddad, 1988, and Hyland and Bernstein, 1985, for details. Note that $r=G^{T} \Gamma$. Clearly, $Q, P, \hat{Q}$, and $\hat{P}$ are symmetric and nonnegative definite.

Next note that with the above definitions, (A.5) implies (3.3) and that (3.2) holds. Hence $r=G^{r} \Gamma$ is idempotent, i.e., $r^{2}=r$. Sylvester's inequality yields (3.18). Note also that (3.7) and (3.8) hold.

The components of $\mathcal{Q}$ and $P$ can be written in terms of $Q, P, \hat{Q}, \hat{P}, G$, and $\Gamma$ as

$$
\begin{array}{ll}
Q_{1}=Q+\hat{Q}, & P_{1}=P+\hat{P} \\
Q_{12}=\hat{Q} \Gamma^{T}, & P_{12}=-\hat{P} G^{T} \\
Q_{2}=\Gamma \hat{Q} \Gamma^{T}, & P_{2}=G \hat{P} G^{T} \tag{A.16}
\end{array}
$$

Next note that by using (A.14)-(A.16), (A.7) becomes

$$
C_{m} \hat{S}=C\left[I_{n}+\gamma^{-2}(Q+\hat{Q}) P\right] G^{T}
$$

where

$$
\hat{s} \triangleq I_{n_{-}}+\gamma^{-2} \Gamma \hat{Q} P G^{T}
$$

To prove that $\hat{S}$ is invertible use (3.7) and (3.4) and note that

$$
\begin{aligned}
I_{n_{m}}+\gamma^{-2} \Gamma \hat{Q} P G^{T} & =I_{n_{m}}+\gamma^{-2} \Gamma \hat{Q} r^{T} P G^{T} \\
& =I_{n_{m}}+\gamma^{-2}\left(\Gamma \hat{Q} \Gamma^{T}\right)\left(G P G^{T}\right) .
\end{aligned}
$$

Since $\Gamma \hat{Q} \Gamma^{T}$ and $G P G^{T}$ are nonnegative definite, their product has nonnegative eigenvalues. Thus each eigenvalue of $I_{n_{m}}+\boldsymbol{\gamma}^{-2} \Gamma \hat{Q} P G^{T}$ is real and is greater than unity. Hence $\hat{S}$ is invertible. Now note that by using (3.3) and (3.4) it can be shown that

$$
G^{T} \hat{S}^{-1} \Gamma=S \tau
$$

The expressions (3.11), (3.12) and (3.13) follow from (A.6), (A.7), (3.9) and the definition of $Q$ by using the above identities. Next, computing either $\Gamma$ (A.9)-(A.10) or $G(A .12)+(A .13)$ yields (3.10). Substituting this expression for $A_{m}$ into (A.8)-(A.13) it follows that (A.10) $=\Gamma$ (A.9) and (A.13) $=G(A .12)$. Thus, (A.10) and (A.13) are superfluous and can be omitted. Next, using $(\mathrm{A} .8)+\mathrm{G}^{\boldsymbol{T}} \Gamma(\mathrm{A} .9) \mathrm{G}-(\mathrm{A} .9) \mathrm{G}-[(\mathrm{A} .9) \mathrm{G}]^{\boldsymbol{T}}$ and $\mathrm{G}^{\boldsymbol{T}} \Gamma(\mathrm{A} .9) \mathrm{G}-(\mathrm{A} .9) \mathrm{G}-[(\mathrm{A} .9) \mathrm{G}]^{\boldsymbol{T}}$ yields (3.14) and (3.16). Using (A.11) $+\Gamma^{T} \mathrm{G}(\mathrm{A} .12) \Gamma-(\mathrm{A} .12) \Gamma-[(\mathrm{A} .12) \Gamma]^{T}$ and $\Gamma^{T} \mathrm{G}(\mathrm{A} .12) \Gamma-(\mathrm{A} .12) \Gamma$ $[(A .12) \Gamma]^{T}$ yields (3.15) and (3.17).

Finally, to prove the converse we use (3.10)-(3.18) to obtain (2.16) and (A.4)-(A.7). Let $A_{m}, B_{m}, C_{m}, G, \Gamma, r, Q, P, \hat{Q}, \hat{P}, Q$ be as in the statement of Theorem 3.1 and define $Q_{1}, Q_{12}, Q_{2}$, $P_{1}, P_{12}, P_{2}$ by (A.14)-(A.16). Using (3.3), (3.11) and (3.12) it is easy to verify (A.6) and (A.7). Finally, substitute the definitions of $Q, P, \hat{Q}, \hat{P}, G, \Gamma$, and $r$ into (3.14)-(3.17) along with (3.3), (3.4), (3.7) and (3.8) to obtain (2.16) and (A.4). Finally, note that

$$
Q=\left[\begin{array}{cc}
Q & 0_{n \times n_{m}} \\
0_{n_{m} \times n} & 0_{n_{m}}
\end{array}\right]+\left[\begin{array}{c}
I_{n} \\
\Gamma
\end{array}\right] \hat{Q}\left[I_{n} \quad \Gamma^{T}\right],
$$

which shows that $Q \geq 0$.

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# Steady-State Kalman Filtering with an 

 $\mathbf{H}_{\text {co }}$ Error Bound
## by

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#### Abstract

An estimator design problem is considered which involves both $L_{2}$ (least squares) and $H_{\infty}$ (worst-case frequency-domain) aspects. Specifically, the goal of the problem is to minimize an $L_{2}$ state-estimation error criterion subject to a prespecified $H_{\infty}$ constraint on the state-estimation error. The $H_{\infty}$ estimation-error constraint is embedded within the optimization process by replacing the covariance Lyapunov equation by a Riccati equation whose solution leads to an upper bound on the $L_{2}$ state-estimation error. The principal result is a sufficient condition for characterizing fixed-order (i.e., full- and reduced-order) estimators with bounded $L_{2}$ and $H_{\infty}$ estimation error. The sufficient condition involves a system of modified Riccati equations coupled by an oblique projection, i.e., idempotent matrix. When the $H_{\infty}$ constraint is absent, the sufficient condition specializes to the $L_{2}$ state-estimation result given in [2].


Keywords: Kalman filter, $H_{\infty}$ norm, reduced-order state estimation, optimal projection equations, Hankel norm.

Supported in part by the Air Force Office of Scientific Research under contract F49620-86-C0002.

## 1. Introduction

One of the fundamental problems in dynamic systems theory is the observation of state variables. Although an extensive theoretical foundation has been developed for the quadratic (least squares) error criterion, state estimation with a worst-case frequency-domain design objective has apparently not been considered. In the present paper we thus extend the least squares formulation to include a frequency-domain bound on the state-estimation error. The underlying idea involves the application of state-space techniques which have recently been developed for $H_{\infty}$ control design in [1,4-6]. The results of the present paper are thus complementary to the results obtained in [1].

The principal result of the present paper is a sufficient condition which yields full- and reducedorder estimators satisfying an optimized $L_{2}$ error bound as well as a prespecified $H_{\infty}$ error bound. In the full-order case, the $H_{\infty}$-constrained estimator involves a modified Riccati equation which specializes to the standard steady-state Kalman filter when the $H_{\infty}$ constraint is absent. In the reduced-order case the $H_{\infty}$-constrained result leads to a direct generalization of the optimal projection approach developed in [2] for the unconstrained $L_{2}$ state-estimation problem. While the $L_{2}$-optimal reduced-order state estimator was characterized in [2] by means of a coupled system of one modified Riccati equation and two modified Lyapunov equations, the $H_{\infty}$-constrained solution involves a coupled system consisting of three modified Riccati equations and one modified Lyapunov equation. As in [2], the coupling is due to the presence of an oblique projection (idempotent matrix) with additional coupling now arising from the $H_{\infty}$ constraint. When the $H_{\infty}$ constraint is sufficiently relaxed, these conditions again specialize directly to those given in [2].

We note that the development in the present paper is limited to the case in which the plant is asymptotically stable. These results can also be extended to the unstable plant case, although with additional complexity. This case will thus be treated in a future paper.

The contents of the paper are as follows. After collecting notation in Section 2, the statement of the $H_{\infty}$-Constrained State-Estimation Problem is given in Section 3. The principal result of this section (Lemma 3.1) shows that if the algebraic Lyapunov equation for the covariance is replaced by a modified Riccati equation possessing a nonnegative-definite solution, then the $H_{\infty}$ estimationerror constraint is enforced and the $L_{2}$ state-estimation error criterion is bounded above by an auxiliary cost function. The problem of determining a reduced-order estimator which minimizes this upper bound subject to the Riccati equation constraint is considered in Section 4 as the Auxiliary Minimization Problem. Necessary conditions for the Auxiliary Minimization Problem
(Theorem 4.1) are given in the form of a coupled system of modified algebraic Riccati equations. To develop connections with standard Kalman filter theory the full-order estimator result is also given. In Section 5 the necessary conditions of Theorem 4.1 are combined with Lemma 3.1 to yield sufficient conditions for bounded $H_{\infty}$ and $L_{2}$ estimation error. Although our result gives sufficient conditions for $H_{\infty}$ estimation error, we also state hypotheses under which these conditions are also necessary (Proposition 5.1).

## 2. Notation and Definitions

$\mathbb{R}, \mathbb{R}^{r \times e}, \mathbb{R}^{r}, \mathbb{I}$
$I_{r},()^{\mathrm{T}}, \mathrm{O}_{\mathrm{r} \times \mathrm{a}}, \mathrm{O}_{\mathrm{r}}$
tr
$\sigma_{\max }(Z)$
$\lambda_{\mathrm{max}}(Z)$
$\|Z\|_{\mathrm{F}}$
$\|H(s)\|_{\infty}$
$\mathbb{S}^{r}, \mathbb{N}^{r}, \mathbb{P}^{r}$
$Z_{1} \leq Z_{2}, Z_{1}<Z_{2}$
$n, \ell_{2} n_{e}, p, q, r ; \tilde{n}$
$x, y, y_{e}, x_{e}, \tilde{x}$
$\tilde{x}$
$\tilde{\boldsymbol{x}}$
$A, C$
$D_{1}, D_{2}, E$
$L$
$A_{e}, B_{e}, C_{e}$
$\tilde{A}$
$\tilde{D}, \tilde{E}$
$R$
$w(\cdot)$
$V_{1}, V_{2}$
real numbers, $r \times s$ real matrices, $\mathbb{R}^{r \times 1}$, expected value $r \times r$ identity matrix, transpose, $r \times s$ zero matrix, $0_{r \times r}$
trace
largest singular value of matrix $Z$
largest eigenvalue of matrix $Z$ with real spectrum
$\left[\operatorname{tr} Z Z^{T}\right]^{\frac{1}{2}}$ (Frobenius matrix norm)
$\sup _{\omega \in \mathbb{R}} \sigma_{\max }[H(j \omega)]$
$\boldsymbol{r} \times \mathbf{r}$ symmetric, nonnegative-definite, positive-definite matrices
$Z_{2}-Z_{1} \in \mathbb{N}^{r}, Z_{2}-Z_{1} \in \mathbb{P}^{r}, Z_{1}, Z_{2} \in \mathbb{S}^{r}$
positive integers; $n+n_{e} ; n_{e} \leq n$
$n, \ell, q, n_{e}, \tilde{n}$-dimensional vectors
$\left[\begin{array}{c}x \\ x_{i}\end{array}\right]$
$n \times n, \ell \times n$ matrices
$n \times p, \ell \times p, r \times q$ matrices
$q \times n$ matrix
$n_{\varepsilon} \times n_{e}, n_{e} \times \ell, q \times n_{e}$ matrices
$\left[\begin{array}{cc}A & 0_{n \times n_{e}} \\ B_{e} C & A_{s}\end{array}\right]$
$\left[\begin{array}{c}D_{1} \\ B_{8} D_{2}\end{array}\right], \quad\left[\begin{array}{ll}E L & -E C_{e}\end{array}\right]$
$E^{\mathrm{T}} E$, estimation error weighting in $\mathbb{P}^{a}$
p-dimensional standard white noise process
intensity of $D_{1} w(\cdot), D_{2} w(\cdot) ; \quad V_{1}=D_{1} D_{1}^{T} \in \mathbb{N}^{n}, V_{2}=D_{2} D_{2}^{\mathbf{T}} \in \mathbb{P}^{\ell}$
$V_{12}$
$\tilde{R}, \tilde{V}$ $\boldsymbol{\gamma}$
cross intensity of $D_{2} w(\cdot), D_{2} w(\cdot) ; \quad V_{12}=D_{1} D_{2}^{\mathrm{T}} \in \mathbb{R}^{n \times \ell}$
$\left[\begin{array}{cc}L^{\mathrm{T}} R L & -L^{\mathrm{T}} R C_{e} \\ -C_{e}^{\mathrm{T}} R L & C_{e}^{\mathrm{T}} R C_{e}\end{array}\right], \quad\left[\begin{array}{cc}V_{1} & V_{12} B_{e}^{\mathrm{T}} \\ B_{e} V_{12}^{\mathrm{T}} & B_{e} V_{2} B_{e}^{\mathrm{T}}\end{array}\right]$
positive constant

## 3. Statement of the Problem

In this section we introduce the reduced-order state-estimation problem with a constraint on the $H_{\infty}$ norm of the state-estimation error. Specifically, the transfer function between disturbances and error states is constrained to have $H_{\infty}$ norm less than $\gamma$. In this paper it is assumed that the plant is asymptotically stable, i.e., the eigenvalues of $A$ are in the open left half plane.
$\mathbf{H}_{\infty}$-Constrained State-Estimation Problem. Given the nth-order observed system

$$
\begin{align*}
& \dot{x}(t)=A x(t)+D_{1} w(t)  \tag{3.1}\\
& y(t)=C x(t)+D_{2} w(t) \tag{3.2}
\end{align*}
$$

where $t \in[0, \infty)$, determine an $n_{\varepsilon}$ th-order state estimator

$$
\begin{gather*}
\dot{x}_{e}(t)=A_{\varepsilon} x_{e}(t)+B_{e} y(t),  \tag{3.3}\\
y_{e}(t)=C_{e} x_{e}(t), \tag{3.4}
\end{gather*}
$$

where $n_{\varepsilon} \leq n$, which satisfies the following design criteria:
(i) $A_{e}$ is asymptotically stable;
(ii) the $r \times p$ transfer function

$$
\begin{equation*}
H(s) \triangleq \tilde{E}\left(s I_{\tilde{n}}-\tilde{A}\right)^{-1} \tilde{D} \tag{3.5}
\end{equation*}
$$

from disturbances $w(t)$ to error states $E\left[L x(t)-y_{e}(t)\right]=\tilde{E} \tilde{x}(t)$ satisfies the constraint.

$$
\begin{equation*}
\|H(s)\|_{\infty} \leq \gamma \tag{3.6}
\end{equation*}
$$

where $\gamma>0$ is a given constant; and
(iii) the $L_{2}$ state-estimation error criterion

$$
\begin{equation*}
J\left(A_{e}, B_{e}, C_{e}\right) \triangleq \lim _{t \rightarrow \infty} \mathbb{E}\left\{\left[L x(t)-y_{e}(t)\right]^{\mathrm{T}} R\left[L x(t)-y_{e}(t)\right]\right\} \tag{3.7}
\end{equation*}
$$

is minimized.

It is useful to note that the augmented system (3.1)-(3.4) can be written as

$$
\begin{equation*}
\dot{\tilde{x}}(t)=\tilde{A} \tilde{x}(t)+\tilde{D} w(t), \quad t \in[0, \infty) \tag{3.8}
\end{equation*}
$$

and that (3.7) is equivalent to

$$
\begin{equation*}
J\left(A_{e}, B_{e}, C_{e}\right)=\lim _{t \rightarrow \infty} \mathbb{E}\left\{[\tilde{E} \tilde{x}(t)]^{\mathbf{T}}[\tilde{E} \tilde{x}(t)]\right\}=\lim _{t \rightarrow \infty} \mathbb{E}\left[\tilde{x}^{\mathbf{T}}(t) \tilde{R} \tilde{x}(t)\right] \tag{3.9}
\end{equation*}
$$

Furthermore, if $A_{\varepsilon}$ is asymptotically stable for a given estimator ( $A_{e}, B_{e}, C_{e}$ ) then the $L_{2}$ stateestimation error criterion is given by

$$
\begin{equation*}
J\left(A_{e}, B_{e}, C_{s}\right)=\operatorname{tr} \tilde{Q} \tilde{R}, \tag{3.10}
\end{equation*}
$$

where the steady-state covariance defined by

$$
\begin{equation*}
\tilde{Q} \triangleq \lim _{t \rightarrow \infty} \mathbb{E}\left[\tilde{x}(t) \tilde{x}^{\mathrm{T}}(t)\right] \tag{3.11}
\end{equation*}
$$

satisfies the $\hat{n} \times \tilde{n}$ Lyapunov equation

$$
\begin{equation*}
0=\tilde{A} \tilde{Q}+\tilde{Q} \tilde{A}^{\mathbf{T}}+\tilde{V} . \tag{3.12}
\end{equation*}
$$

Using (3.10) and (3.12) we now show that the criterion (3.7) is an error measure involving the impulse response of (3.8) with respect to an $L_{2}$ norm.

Proposition 3.1. If $A_{e}$ is asymptotically stable then the $L_{2}$ state-estimation error criterion (3.7) can be written as

$$
\begin{equation*}
J\left(A_{e}, B_{e}, C_{e}\right)=\int_{0}^{\infty}\left\|\tilde{E} e^{\tilde{A} t} \tilde{D}\right\|_{F}^{2} d t \tag{3.13}
\end{equation*}
$$

Proof. It need only be noted that (3.10) is equivalent to

$$
\operatorname{tr} \int_{0}^{\infty} e^{\tilde{A} t} \tilde{V} e^{\tilde{A}^{\mathrm{T}} t} d t \tilde{R}=\operatorname{tr} \int_{0}^{\infty}\left(\tilde{E} e^{\tilde{A} t} \tilde{D}\right)\left(\tilde{E} e^{\tilde{A} t} \tilde{D}\right)^{\mathbf{T}} d t
$$

which is equivalent to (3.13).
The key step in enforcing (3.6) is to replace the algebraic Lyapunov equation (3.12) by an algebraic Riccati equation. Justification for this technique is provided by the following result.

Lemma 3.1. Let $\left(A_{e}, B_{e}, C_{e}\right)$ be given and assume there exists $\mathcal{Q} \in \mathbb{R}^{\bar{n} \times \bar{n}}$ satisfying

$$
\begin{equation*}
Q \in \mathbb{N}^{\bar{n}} \tag{3.14}
\end{equation*}
$$

and

$$
\begin{equation*}
0=\tilde{A} Q+Q \tilde{A}^{\mathrm{T}}+\gamma^{-2} Q \tilde{R} Q+\tilde{V} . \tag{3.15}
\end{equation*}
$$

Then

$$
\begin{equation*}
\left(\tilde{A},\left[\gamma^{-2} Q \tilde{R} Q+\tilde{V}\right]^{\frac{1}{2}}\right) \text { is stabilizable } \tag{3.16}
\end{equation*}
$$

if and only if

$$
\begin{equation*}
A_{e} \text { is asymptotically stable. } \tag{3.17}
\end{equation*}
$$

Furthermore, in this case

$$
\begin{gather*}
\|H(s)\|_{\infty} \leq \gamma  \tag{3.18}\\
\tilde{Q} \leq \mathcal{Q} \tag{3.19}
\end{gather*}
$$

and

$$
\begin{equation*}
J\left(A_{e}, B_{e}, C_{e}\right) \leq J\left(A_{e}, B_{e}, C_{e}, \mathcal{Q}\right) \tag{3.20}
\end{equation*}
$$

where

$$
\begin{equation*}
J\left(A_{e}, B_{e}, C_{e}, \mathcal{Q}\right) \triangleq \operatorname{tr} Q \tilde{R} . \tag{3.21}
\end{equation*}
$$

Proof. Using the assumed existence a nonnegative-definite solution to (3.15) and the stabilizability condition (3.16), it follows from the dual of Lemma 12.2 of [9] that $\tilde{A}$ is asymptotically stable. Since $\tilde{A}$ is lower block triangular, $\tilde{A}$ asymptotically stable implies $A_{e}$ is asymptotically stable. Conversely, since $A$ is assumed to be asymptotically stable, (3.17) implies $\tilde{A}$ is asymptotically stable and thus (3.16) holds. The proof of (3.18) follows from a standard manipulation of (3.15); for details see Lemma 1 of [8]. To prove (3.19) subtract (3.12) from (3.15) to obtain

$$
\begin{equation*}
0=\tilde{A}(\mathcal{Q}-\tilde{Q})+(\mathcal{Q}-\tilde{Q}) \tilde{A}^{\mathrm{T}}+\gamma^{-2} Q \tilde{R} Q \tag{3.22}
\end{equation*}
$$

which, since $\tilde{A}$ is asymptotically stable, is equivalent to

$$
\begin{equation*}
\mathcal{Q}-\tilde{Q}=\int_{0}^{\infty} e^{\bar{A}^{t} t}\left[\gamma^{-2} Q \tilde{R} Q\right] e^{\bar{A}^{\boldsymbol{T}} t} d t \geq 0 \tag{3.23}
\end{equation*}
$$

Finally, (3.20) follows immediately from (3.19).

Lemma 3.1 shows that the $H_{\infty}$ constraint is automatically enforced when a nonnegative-definite solution to (3.15) can be shown to exist. Furthermore, the solution $Q$ provides an upper bound for the steady-state covariance $\tilde{Q}$ along with a bound on the $L_{2}$ state-estimation error criterion. Naxt, we present a partial converse of Lemma 3.1 which guarantees the existence of a nonnegative-definite solution to (3.15) when (3.18) is satisfied.

Lemma 3.2. Let ( $A_{s}, B_{s}, C_{s}$ ) be given, suppose $A_{s}$ is asymptotically stable, and assume the $H_{\infty}$ state-estimation error constraint (3.18) is satisfied. Then there exists a unique nonnegativedefinite solution $Q$ satisfying (3.15). Furthermore, $\left(\tilde{A}+\gamma^{-2} Q \tilde{R}, \tilde{D}\right)$ is stabilizable if and only if $\tilde{A}+\boldsymbol{\gamma}^{-2} Q \tilde{R}$ is asymptotically stable.

Proof. The result is an immediate consequence of Theorems 3 and 2 of [3], pp. 150 and 167, along with the dual of Lemma 12.2 of [9].

Finally, we show that the quadratic term $\gamma^{-2} Q \tilde{R} Q$ in (3.15) also constrains the Hanke! norm of the estimation error $E\left[L x(t)-y_{\varepsilon}(t)\right]$ when $Q$ is positive definite. To show this let $\tilde{P} \in \mathbb{N}^{\bar{n}}$ be the observability Gramian for the augmented system $(\tilde{A}, \tilde{D}, \tilde{E})$ which satisfies

$$
\begin{equation*}
0=\tilde{A}^{\mathrm{T}} \tilde{P}+\tilde{P} \tilde{A}+\tilde{R} \tag{3.24}
\end{equation*}
$$

Proposition 3.2. Let ( $A_{e}, B_{e}, C_{e}$ ) be given and assume there exists $Q \in \mathbb{P}^{\bar{n}}$ satisfying (3.15) and (3.16) or, equivalently, (3.17). Then

$$
\begin{equation*}
\lambda_{\text {max }}^{\frac{1}{2}}(\tilde{P} \tilde{Q}) \leq \gamma . \tag{3.25}
\end{equation*}
$$

Proof. Since $Q$ is invertible, (3.15) implies

$$
\begin{equation*}
0=\gamma^{2} \tilde{A}^{\mathrm{T}} Q^{-1}+\gamma^{2} Q^{-1} \tilde{A}+\gamma^{2} Q^{-1} \tilde{V} Q^{-1}+\tilde{R} \tag{3.26}
\end{equation*}
$$

Next, subtract (3.24) from (3.26) to obtain

$$
\begin{equation*}
0=\tilde{A}^{\mathrm{T}}\left(\gamma^{2} Q^{-1}-\tilde{P}\right)+\left(\gamma^{2} Q^{-1}-\tilde{P}\right) \tilde{A}+\gamma^{2} Q^{-1} \tilde{V} Q^{-1} \tag{3.27}
\end{equation*}
$$

which, since $\tilde{A}$ is asymptotically stable, is equivalent to

$$
\begin{equation*}
\gamma^{2} Q^{-1}-\tilde{P}=\int_{0}^{\infty} e^{i^{T} t}\left[\gamma^{2} Q^{-1} \tilde{V} Q^{-1}\right] e^{\bar{A} t} d t \geq 0 \tag{3.28}
\end{equation*}
$$

Thus (3.28) implies $\tilde{P} \leq \gamma^{2} Q^{-1}$ or, equivalently, $Q^{\frac{1}{2}} \tilde{P} Q^{\frac{1}{2}} \leq \gamma^{2} I_{\tilde{n}}$. Hence,

$$
\boldsymbol{\gamma}^{2} \geq \lambda_{\max }\left(Q^{\frac{1}{2}} \tilde{P} Q^{\frac{1}{2}}\right)=\lambda_{\max }\left(\tilde{P}^{\frac{1}{2}} Q \tilde{P}^{\frac{1}{2}}\right) \geq \lambda_{\max }\left(\tilde{P}^{\frac{1}{2}} \tilde{Q} \tilde{P}^{\frac{1}{2}}\right)=\lambda_{\max }(\tilde{P} \tilde{Q})
$$

## 4. The Auxiliary Minimization Problem and Necessary Conditions for Optimality

As discussed in the previous section, the replacement of (3.12) by (3.15) enforces the $H_{\infty}$ state-estimation error constraint and results in an upper bound for the $L_{2}$ state-estimation error criterion. That is, given an estimator ( $A_{e}, B_{e}, C_{e}$ ) satisfying the $H_{\infty}$ estimation constraint, the actual $L_{2}$ state-estimation error criterion is guaranteed to be no worse than the bound given by $J\left(A_{e}, B_{e}, C_{e}, \mathcal{Q}\right)$ if (3.15) is solvable. Hence, $J\left(A_{e}, B_{e}, C_{e}, \mathcal{Q}\right)$ can be interpreted as an auxiliary cost which leads to the following optimization problem.

Auxiliary Minimization Problem. Determine $\left(A_{e}, B_{e}, C_{e}, Q\right)$ which minimizes $J\left(A_{e}, B_{e}\right.$, $C_{e}, \mathcal{Q}$ ) subject to (3.14) and (3.15).

It follows from Lemma 3.1 that the satisfaction of (3.14)-(3.16) leads to 1) $A_{e}$ stable; 2) $H_{\infty}$ estimation error bound $\gamma$; and 3) an upper bound (3.21) for the $L_{2}$ state-estimation error criterion. Hence it remains to determine ( $A_{e}, B_{e}, C_{e}$ ) which minimizes $J\left(A_{e}, B_{e}, C_{e}, Q\right)$ and thus provides an optimized bound for the actual $L_{2}$ criterion $J\left(A_{e}, B_{e}, C_{e}\right)$. Rigorous derivation of the necessary conditions for the Auxiliary Minimization Problem requires additional technical assumptions. Specifically, we restrict ( $A_{e}, B_{e}, C_{e}, \mathcal{Q}$ ) to the open set

$$
\begin{aligned}
S \triangleq\left\{\left(A_{e}, B_{e}, C_{e}, Q\right): Q \in \mathbb{P}^{\tilde{n}}, \quad \tilde{A}+\gamma^{-2} Q \tilde{R}\right. \text { is asymptotically stable, } \\
\text { and } \left.\left(A_{e}, B_{e}, C_{e}\right) \text { is controllable and observable }\right\} .
\end{aligned}
$$

Remark 4.1. The set $S$ constitutes sufficient conditions under which the Lagrange multiplier technique is applicable to the Auxiliary Minimization Problem. Specifically, the requirement that $Q$ be positive definite replaces (3.14) by an open set constraint, the stability of $\tilde{A}+\gamma^{-2} Q \tilde{R}$ serves as a normality condition, and ( $A_{a}, B_{a}, C_{8}$ ) minimal is a nondegeneracy condition.

The following Lemma is needed for the statement of the main result.
Lemma 4.1. Let $\hat{Q}, \hat{P} \in \mathbb{N}^{n}$ and suppose rank $\hat{Q} \hat{P}=n_{e}$. Then there exist $n_{e} \times n G, \Gamma$ and $n_{e} \times n_{e}$ invertible $M$, unique except for a change of basis in $\mathbb{R}^{n_{e}}$, such that

$$
\begin{equation*}
\hat{Q} \hat{P}=G^{\mathrm{T}} M \Gamma \tag{4.1}
\end{equation*}
$$

$$
\begin{equation*}
\Gamma G^{\mathrm{T}}=I_{n_{\mathrm{e}}} . \tag{4.2}
\end{equation*}
$$

Furthermore, the $n \times n$ matrices

$$
\begin{gather*}
\tau \triangleq G^{\mathbf{T}} \Gamma  \tag{4.3}\\
\tau_{\perp} \triangleq I_{n}-\tau \tag{4.4}
\end{gather*}
$$

are idempotent and have rank $n_{e}$ and $n-n_{e}$, respectively. If, in addition,

$$
\begin{equation*}
\operatorname{rark} \hat{Q}=\operatorname{rank} \hat{P}=n_{e} \tag{4.5}
\end{equation*}
$$

then

$$
\begin{equation*}
\hat{Q}=\tau \hat{Q}, \quad \hat{P}=\hat{P} \tau . \tag{4.6}
\end{equation*}
$$

Finally, if $P \in \mathbb{N}^{n}$ then the inverse

$$
\begin{equation*}
S \triangleq\left(I_{n}+\gamma^{-2} \hat{Q} P\right)^{-1} \tag{4.8}
\end{equation*}
$$

exists.
Proof. Conditions (4.1)-(4.7) are a direct consequence of Theorem 6.2.5 of [7]. To prove that the inverse in (4.8) exists, note that since the eigenvalues of $\hat{Q} P$ coincide with the eigenvalues of the nonnegative-definite matrix $P^{\frac{1}{2}} \hat{Q} P^{\frac{1}{2}}$, it follows that $\hat{Q} P$ has nonnegative eigenvalues. Thus, the eigenvalues of $I_{n}+\boldsymbol{\gamma}^{-\mathbf{2}} \hat{Q} P$ are all greater than one so that the above inverse exists.

Finally, for arbitrary $Q \in \mathbb{R}^{n \times n}$ define the notation

$$
\begin{equation*}
Q_{\mathrm{a}} \triangleq Q C^{\mathrm{T}}+V_{12}, \quad \Sigma \triangleq L^{\mathrm{T}} R L . \tag{4.9}
\end{equation*}
$$

Theorem 4.1. If $\left(A_{e}, B_{e}, C_{e}, \mathcal{Q}\right) \in S$ solves the Auxiliary Minimization Problem then there exist $Q, P, \hat{Q}, \hat{P} \in \mathbb{N}^{n}$ such that

$$
\begin{align*}
A_{\mathbf{t}} & =\Gamma\left(A-Q_{\mathrm{a}} V_{2}^{-1} C-\gamma^{-4} Q \Sigma Q P S\right) G^{\mathrm{T}}  \tag{4.10}\\
B_{\mathrm{a}} & =\Gamma Q_{a} V_{2}^{-1}  \tag{4.11}\\
C_{\varepsilon} & =L\left(I_{n}+\gamma^{-2} Q P S\right) G^{\mathrm{T}}  \tag{4.12}\\
\mathcal{Q} & =\left[\begin{array}{cc}
Q+\hat{Q} & \hat{Q} \Gamma^{\mathrm{T}} \\
\Gamma \hat{Q} & \Gamma \hat{Q} \Gamma^{\mathrm{T}}
\end{array}\right] \tag{4.13}
\end{align*}
$$

and such that $Q, P, \hat{Q}, \hat{P}$ satisfy

$$
\begin{align*}
0= & A Q+Q A^{\mathrm{T}}+V_{1}+\gamma^{-2} Q \Sigma Q-Q_{a} V_{2}^{-1} Q_{a}^{\mathrm{T}}+\tau_{\perp} Q_{a} V_{2}^{-1} Q_{a}^{\mathrm{T}} \tau_{\perp}^{\mathrm{T}},  \tag{4.14}\\
0= & A^{\mathrm{T}} P+P A-\gamma^{-4} S^{\mathrm{T}} P Q \Sigma Q P S+\tau_{\perp}^{\mathrm{T}}\left(I_{n}+\gamma^{-2} Q P S\right)^{\mathrm{T}} \Sigma\left(I_{n}+\gamma^{-2} Q P S\right) \tau_{\perp},  \tag{4.15}\\
0= & \left(A-\gamma^{-4} Q \Sigma Q P S\right) \hat{Q}+\hat{Q}\left(A-\gamma^{-4} Q \Sigma Q P S\right)^{\mathrm{T}}+\gamma^{-4} \hat{Q} S^{\mathrm{T}} P Q \Sigma Q P S \hat{Q} \\
& +Q_{a} V_{2}^{-1} Q_{a}^{\mathrm{T}}-\tau_{\perp} Q_{a} V_{2}^{-1} Q_{a}^{\mathrm{T}} \tau_{\perp}^{\mathrm{T}},  \tag{4.16}\\
0= & \left(A-Q_{a} V_{2}^{-1} C+\gamma^{-2} Q \Sigma\right)^{\mathrm{T}} \hat{P}+\hat{P}\left(A-Q_{a} V_{2}^{-1} C+\gamma^{-2} Q \Sigma\right) \\
& +\left(I_{n}+\gamma^{-2} Q P S\right)^{\mathrm{T}} \Sigma\left(I_{n}+\gamma^{-2} Q P S\right)-\tau_{\perp}^{\mathrm{T}}\left(I_{n}+\gamma^{-2} Q P S\right)^{\mathrm{T}} \Sigma\left(I_{n}+\gamma^{-2} Q P S\right) \tau_{\perp},  \tag{4.17}\\
& \operatorname{rank} \hat{Q}=\operatorname{rank} \hat{P}=\operatorname{rank} \hat{Q} \hat{P}=n_{e} . \tag{4.18}
\end{align*}
$$

Furthermore, the auxiliary cost is given by

$$
\begin{equation*}
J\left(A_{e}, B_{\varepsilon}, C_{e}, \mathcal{Q}\right)=\operatorname{tr} L^{\mathrm{T}} R L\left(Q+\gamma^{-4} Q P S \hat{Q} S^{\mathrm{T}} P Q\right) \tag{4.19}
\end{equation*}
$$

Conversely, if there exist $Q, P, \hat{Q}, \hat{P} \in \mathbb{N}^{n}$ satisfying (4.14)-(4.18), then ( $A_{e}, B_{e}, C_{e}, Q$ ) given by (4.10)-(4.13) satisfies (3.14) and (3.15) with the auxiliary cost (3.21) given by (4.19).

Proof. See Appendix A.
Remark 4.2. Theorem 4.1 presents necessary conditions for the Auxiliary Minimization Problem which explicitly synthesize extremal full- and reduced-order estimators ( $A_{e}, B_{e}, C_{e}$ ). If the $H_{\infty}$ estimation constraint is sufficiently relaxed, i.e., $\boldsymbol{\gamma} \rightarrow \infty$, then $S=I_{n}$. In this case equations (4.16) and (4.17) become decoupled from (4.15) and thus (4.15) becomes superfluous. Furthermore, (4.14), (4.16) and (4.17) specialize to the optimal projection equations obtained in [2].

As discussed in [2], in the full-order (Kalman Filter) case $n_{e}=n, G=\Gamma^{-1}$ and thus $G=\Gamma=$ $r=I_{n}$ and $\tau_{\perp}=0$ without loss of generality. To develop further connections with the standard Kalman filter theory assume

$$
\begin{equation*}
V_{12}=0 \tag{4.20}
\end{equation*}
$$

In this case (4.15) implies that $P=0$ so that the gain expressions (4.10)-(4.12) become

$$
\begin{align*}
& A_{\varepsilon}=A-Q C^{\mathrm{T}} V_{2}^{-1} C  \tag{4.21}\\
& B_{\varepsilon}=Q C^{\mathrm{T}} V_{2}^{-1}  \tag{4.22}\\
& C_{\varepsilon}=L \tag{4.23}
\end{align*}
$$

while equations.(4.14)-(4.16) and auxiliary cost (4.19) specialize to

$$
\begin{gather*}
0=A Q+Q A^{\mathrm{T}}+V_{1}+\gamma^{-2} Q L^{\mathrm{T}} R L Q-Q C^{\mathrm{T}} V_{2}^{-1} C Q  \tag{4.24}\\
J\left(A_{e}, B_{e}, C_{e}, Q\right)=\operatorname{tr} L^{\mathrm{T}} R L Q \tag{4.25}
\end{gather*}
$$

Remark 4.3. Note that the necessary conditions for the full-order problem involve one modified Riccati equation. This equation is similar to the observer Riccati equation with the additional quadratic term $\boldsymbol{\gamma}^{-2} Q L^{\mathrm{T}} R L Q$. Finally, note that when the $H_{\infty}$ estimation constraint is sufficiently relaxed, i.e., $\boldsymbol{\gamma} \rightarrow \infty$, (4.24) reduces to the standard observer Riccati equation of steady-state Kalman filter theory.

## 5. Sufficient Conditions for Combined $\mathrm{L}_{\mathbf{2}} / \mathrm{H}_{\infty}$ Estimation

In this section we combine Lemma 3.1 with the converse of Theorem 4.1 to obtain our main result guaranteeing constrained $H_{\infty}$ state-estimation error and an optimized $L_{2}$ state-estimation error bound.

Theorem 5.1. Suppose there exist $Q, P, \hat{Q}, \hat{P} \in \mathbb{N}^{n}$ satisfying (4.14)-(4.18) and let ( $A_{\varepsilon}, B_{e}$, $C_{\varepsilon}, \mathcal{Q}$ ) be given by (4.10)-(4.13). Then $\left(\tilde{A},\left[\gamma^{-2} Q \tilde{R} Q+\tilde{V}\right]^{\frac{1}{2}}\right)$ is stabilizable if and only if $A_{\varepsilon}$ is asymptotically stable. In this case, the transfer function $H(s)$ defined by (3.5) satisfies the $H_{\infty}$ state-estimation error constraint

$$
\begin{equation*}
\|H(s)\|_{\infty} \leq \gamma, \tag{5.1}
\end{equation*}
$$

and the $L_{2}$ state-estimation error criterion (3.7) satisfies the bound

$$
\begin{equation*}
J\left(A_{e}, B_{e}, C_{e}\right) \leq \operatorname{tr} L^{\mathrm{T}} R L\left(Q+\gamma^{-4} Q P S \hat{Q} S^{\mathrm{T}} P Q\right) \tag{5.2}
\end{equation*}
$$

Proof. The converse portion of Theorem 4.1 implies that $\mathcal{Q}$ given by (4.13) satisfies (3.14) and (3.15). It now follows from Lemma 3.1 that the stabilizability condition (3.16) is equivalent to the asymptotic stability of $A_{s}$, the $H_{\infty}$ state-estimation error constraint (3.18) holds, and the $L_{2}$ state-estimation error criterion (3.7) satisfies the bound (3.20) which, by (4.19), is equivalent to (5.2).

In applying Theorem 5.1 the principal issue concerns conditions on the problem data under which the coupled Riccati equations (4.14)-(4.17) possess nonnegative-definite solutions. Clearly,
for $\boldsymbol{\gamma}$ sufficiently large, (4.14)-(4.17) approximate the pure least squares problem considered in [2]. The important case of interest, however, involves small $\boldsymbol{\gamma}$ so that significant $H_{\infty}$ estimation is enforced. Thus, if (5.1) can be satisfied for a given $\gamma>0$, it is of interest to know whether one atch fixed-order estimator can be obtained by solving (4.14)-(4.17). Lemma 3.2 guarantees that (3.15) possesses a solution for any fixed-order estimator satisfying (5.1). Thus our sufficient conditions will also be necessary so long as the Auxiliary Minimization Problem possesses at least one extremal over $S$. When this is the case we have the following result.

Proposition 5.1. Let $\boldsymbol{\gamma}^{*}$ denote the infimum of $\|H(s)\|_{\infty}$ over all asymptotically stable fixedorder estimators and suppose that the Auxiliary Minimization Problem has an extremal for all $\boldsymbol{\gamma}>\boldsymbol{\gamma}^{*}$. Then for all $\boldsymbol{\gamma}>\boldsymbol{\gamma}^{*}$ there exist $Q, P, \hat{Q}, \hat{P} \in \mathbb{N}^{n}$ satisfying (4.14)-(4.17).

## Appendix A: Proof of Theorem 4.1

To optimize (3.21) over the open set $S$ subject to the constraint (3.15), form the Lagrangian

$$
\begin{equation*}
\mathcal{L}\left(A_{e}, B_{e}, C_{e}, Q, P, \lambda\right) \triangleq \operatorname{tr}\left\{\lambda Q \tilde{R}+\left[\tilde{A} Q+Q \tilde{A}^{T}+\gamma^{-2} Q \tilde{R} Q+\tilde{V}\right] P\right\} \tag{A.1}
\end{equation*}
$$

where the Lagrange multipliers $\lambda \geq 0$ and $P \in \mathbb{R}^{\bar{n} \times \bar{n}}$ are not both zero. We thus obtain

$$
\begin{equation*}
\frac{\partial \mathcal{L}}{\partial Q}=\left(\tilde{A}+\gamma^{-2} Q \tilde{R}\right)^{\mathrm{T}} P+P\left(\tilde{A}+\gamma^{-2} Q \tilde{R}\right)+\lambda \tilde{R} \tag{A.2}
\end{equation*}
$$

Setting $\frac{\partial \mathcal{C}}{\partial \mathcal{Q}}=0$ yields

$$
\begin{equation*}
0=\left(\tilde{A}+\gamma^{-2} Q \tilde{R}\right)^{\mathbf{T}} P+P\left(\tilde{A}+\gamma^{-2} Q \tilde{R}\right)+\lambda \tilde{R} \tag{A.3}
\end{equation*}
$$

Since $\tilde{A}+\gamma^{-2} \varrho \tilde{R}$ is assumed to be stable, $\lambda=0$ implies $P=0$. Hence, it can be assumed without loss of generality that $\lambda=1$. Furthermore, $P$ is nonnegative definite.

Now partition $\tilde{n} \times \tilde{n} \mathcal{Q}, \mathcal{P}$ into $n \times n, n \times n_{e}$, and $n_{e} \times n_{e}$ subblocks as

$$
\mathcal{Q}=\left[\begin{array}{cc}
Q_{1} & Q_{12} \\
Q_{12}^{\mathrm{T}} & Q_{2}
\end{array}\right], \quad P=\left[\begin{array}{cc}
P_{1} & P_{12} \\
P_{12}^{\mathrm{T}} & P_{2}
\end{array}\right],
$$

and for notational convenience define

$$
P Q=\left[\begin{array}{ll}
Z_{1} & Z_{12} \\
Z_{21} & Z_{2}
\end{array}\right],
$$

where

$$
\begin{aligned}
Z_{1} \triangleq P_{1} Q_{1}+P_{12} Q_{12}^{\mathrm{T}}, & Z_{12} \triangleq P_{1} Q_{12}+P_{12} Q_{2} \\
Z_{21} \triangleq P_{12}^{\mathrm{T}} Q_{1}+P_{2} Q_{12}^{\mathrm{T}}, & Z_{2} \triangleq P_{12}^{\mathrm{T}} Q_{12}+P_{2} Q_{2}
\end{aligned}
$$

Thus, with $\lambda=1$ the stationarity conditions are given by

$$
\begin{align*}
& \frac{\partial \mathcal{L}}{\partial Q}=\left(\tilde{A}+\gamma^{-2} Q \tilde{R}\right)^{\mathrm{T}} \mathcal{P}+\rho\left(\tilde{A}+\gamma^{-2} Q \tilde{R}\right)+\tilde{R}=0,  \tag{A.4}\\
& \frac{\partial \mathcal{L}}{\partial A_{e}}=Z_{2}=0,  \tag{A.5}\\
& \frac{\partial \mathcal{L}}{\partial B_{e}}=Z_{21} C^{\mathrm{T}}+P_{12}^{\mathrm{T}} V_{12}+P_{2} B_{e} V_{2}=0,  \tag{A.6}\\
& \begin{aligned}
& \frac{\partial \mathcal{L}}{\partial C_{e}}=2 R C_{e} Q_{2}+2 \gamma^{-2} R C_{e} Z_{12}^{\mathrm{T}} Q_{12}-2 R L Q_{12}-\gamma^{-2} R L Z_{1}^{\mathrm{T}} Q_{12} \\
& \quad-\gamma^{-2} R L Q_{1} Z_{12}-\gamma^{-2} R L Z_{21}^{\mathrm{T}} Q_{2}=0 .
\end{aligned}
\end{align*}
$$

Expanding (3.15) and (A.4) yields

$$
\begin{align*}
& 0=A Q_{1}+Q_{1} A^{\mathrm{T}}+V_{1}+\gamma^{-2}\left(Q_{1} L^{\mathrm{T}}-Q_{12} C_{e}^{\mathrm{T}}\right) R\left(Q_{1} L^{\mathrm{T}}-Q_{12} C_{e}^{\mathrm{T}}\right)^{\mathrm{T}},  \tag{A.8}\\
& 0=A Q_{12}+Q_{12} A_{e}^{\mathrm{T}}+Q_{1} C^{\mathrm{T}} B_{d}^{\mathrm{T}}+V_{12} B_{6}^{\mathrm{T}}+\gamma^{-2} Q_{1} L^{\mathrm{T}} R L Q_{12}-\gamma^{-2} Q_{12} C_{e}^{\mathrm{T}} R L Q_{12} \\
& -\gamma^{-2} Q_{1} L^{\mathbf{T}} R C_{e} Q_{2}+\gamma^{-2} Q_{12} C_{e}^{T} R C_{e} Q_{2},  \tag{A.9}\\
& 0=A_{e} Q_{2}+Q_{2} A_{e}^{\mathrm{T}}+B_{e} C Q_{12}+Q_{12}^{\mathrm{T}} C^{\mathrm{T}} B_{e}^{\mathrm{T}}+B_{e} V_{2} B_{e}^{\mathrm{T}} \\
& +\gamma^{-2}\left(Q_{12}^{\mathrm{T}} L^{\mathrm{T}}-Q_{2} C_{e}^{\mathrm{T}}\right) R\left(Q_{12}^{\mathrm{T}} L^{\mathrm{T}}-Q_{2} C_{e}^{\mathrm{T}}\right)^{\mathrm{T}},  \tag{A.10}\\
& 0=A^{\mathrm{T}} P_{1}+P_{1} A+C^{\mathrm{T}} B_{e}^{\mathrm{T}} P_{12}^{\mathrm{T}}+P_{12} B_{\mathrm{e}} C+\gamma^{-2} L^{\mathrm{T}} R L Z_{1}^{\mathrm{T}}+\gamma^{-2} Z_{1} L^{\mathrm{T}} R L \\
& -\gamma^{-2} L^{\mathrm{T}} R C_{e} Z_{12}^{\mathrm{T}}-\gamma^{-2} Z_{12} C_{e}^{\mathrm{T}} R L+L^{\mathrm{T}} R L,  \tag{A.11}\\
& 0=A^{\mathrm{T}} P_{12}+P_{12} A_{\mathrm{e}}+C^{\mathrm{T}} B_{e}^{\mathrm{T}} P_{2}+\gamma^{-2} L^{\mathrm{T}} R L Z_{21}^{\mathrm{T}}-\gamma^{-2} Z_{1} L^{\mathrm{T}} R C_{e} \\
& +\gamma^{-2} Z_{12} C_{e}^{\mathrm{T}} R C_{e}-L^{\mathrm{T}} R C_{e},  \tag{A.12}\\
& 0=A_{6}^{\mathbf{T}} P_{2}+P_{2} A_{6}-\gamma^{-2} C_{6}^{\mathbf{T}} R L Z_{21}^{\mathbf{T}}-\boldsymbol{\gamma}^{-2} Z_{21} L^{\mathbf{T}} R C_{c}+C_{6}^{\mathbf{T}} R C_{c} . \tag{A.13}
\end{align*}
$$

Now define the $n \times n$ matrices

$$
\begin{gathered}
Q \triangleq Q_{1}-Q_{12} Q_{2}^{-1} Q_{12}^{\mathrm{T}}, \quad P \triangleq P_{1}-P_{12} P_{2}^{-1} P_{12}^{\mathrm{T}} \\
\hat{Q} \triangleq Q_{12} Q_{2}^{-1} Q_{12}^{\mathrm{T}}, \quad \hat{P} \triangleq P_{12} P_{2}^{-1} P_{12}^{\mathrm{T}} \\
T \triangleq-Q_{12} Q_{2}^{-1} P_{2}^{-1} P_{12}^{\mathrm{T}}
\end{gathered}
$$

and the $n_{\varepsilon} \times n, n_{8} \times n_{\varepsilon}$, and $n_{e} \times n$ matrices

$$
G \triangleq Q_{2}^{-1} Q_{12}^{\mathrm{T}}, \quad M \triangleq Q_{2} P_{2}, \quad \Gamma \triangleq-P_{2}^{-1} P_{12}^{\mathrm{T}}
$$

The existence of $Q_{2}^{-1}$ and $P_{2}^{-1}$ follows from the fact that ( $A_{e}, B_{e}, C_{e}$ ) is minimal. See $[1,2]$ for details. Note that $\tau=G^{\mathrm{T}} \Gamma$. Clearly, $Q, P, \hat{Q}$, and $\hat{P}$ are symmetric and nonnegative definite.

Next note that with the above definitions, (A.5) implies (4.2) and that (4.1) holds. Hence $\tau=G^{\mathbf{T}} \Gamma$ is idempotent, i.e., $\tau^{2}=r$. Sylvester's inequality yields (4.18). Note also that (4.6) and (4.7) hold.

The components of $\mathcal{Q}$ and $P$ can be written in terms of $Q, P, \hat{Q}, \hat{P}, G$, and $\Gamma$ as

$$
\begin{array}{ll}
Q_{1}=Q+\hat{Q}, & P_{1}=P+\hat{P} \\
Q_{12}=\hat{Q} \Gamma^{\mathrm{T}}, & P_{12}=\hat{P} G^{\mathrm{T}} \\
Q_{2}=\Gamma \hat{Q} \Gamma^{\mathrm{T}}, & P_{2}=G \hat{P} G^{\mathrm{T}} . \tag{A.16}
\end{array}
$$

Next note that by using (A.14)-(A.16), (A.7) becomes

$$
C_{z} \hat{S}=L\left[I_{n}+\gamma^{-2}(Q+\hat{Q}) P\right] G^{T}
$$

where

$$
\hat{S} \triangleq I_{n_{0}}+\gamma^{-2} \Gamma \hat{Q} P G^{\mathrm{T}}
$$

To prove that $\hat{S}$ is invertible use (4.6) and (4.3) and note that

$$
\begin{aligned}
I_{n_{\bullet}}+\gamma^{-2} \Gamma \hat{Q} P G^{\mathrm{T}} & =I_{n_{\epsilon}}+\gamma^{-2} \Gamma \hat{Q} \tau^{\mathrm{T}} P G^{\mathrm{T}} \\
& =I_{n_{\bullet}}+\gamma^{-2}\left(\Gamma \hat{Q} \Gamma^{\mathrm{T}}\right)\left(G P G^{\mathrm{T}}\right)
\end{aligned}
$$

Since $\Gamma \hat{Q} \Gamma^{\mathrm{T}}$ and $G P G^{\mathrm{T}}$ are nonnegative definite, their product has nonnegative eigenvalues. Thus each eigenvalue of $I_{n_{\mathrm{e}}}+\gamma^{-2} \Gamma \hat{Q} P G^{\mathbf{T}}$ is real and is greater than unity. Hence $\hat{S}$ is invertibie. Now note that by using (4.2) and (4.3) it can be shown that

$$
G^{\mathrm{T}} \hat{S}^{-1} \Gamma=S \tau
$$

The expressions (4.11), (4.12) and (4.13) follow from (A.6), (A.7), (4.8) and the definition of $Q$ by using the above identities. Next, computing either $\Gamma$ (A.9) - (A.10) or $G(A .12)+(A .13)$ yields (4.10). Substituting this expression for $A_{e}$ into (A.8) - (A.13) it follows that (A.10) $=\Gamma$ (A.9) and (A.13) $=G$ (A.12). Thus, (A.10) and (A.13) are superfluous and can be omitted. Next, using $(\mathrm{A} .8)+G^{\mathrm{T}} \Gamma(\mathrm{A} .9) G-(\mathrm{A} .9) G-[(\mathrm{A} .9) G]^{\mathrm{T}}$ and $G^{\mathrm{T}} \Gamma(\mathrm{A} .9) G-(\mathrm{A} .9) G-[(\mathrm{A} .9) G]^{\mathrm{T}}$ yields (4.14) and (4.16). Using (A.11) $+\Gamma^{\mathrm{T}} G(\mathrm{~A} .12) \Gamma-(\mathrm{A} .12) \Gamma-[(\mathrm{A} .12) \Gamma]^{\mathrm{T}}$ and $\Gamma^{\mathrm{T}} G(\mathrm{~A} .12) \Gamma-(\mathrm{A} .12) \Gamma-[(\mathrm{A} .12) \Gamma]^{\mathrm{T}}$ yields (4.15) and (4.17).

Finally, to prove the converse we use (4.10)-(4.18) to obtain (3.15) and (A.4) - (A.7). Let, $A_{a}, B_{e}, C_{a}, G, \Gamma, r, Q, P, \hat{Q}, \hat{P}, Q$ be as in the statement of Theorem 4.1 and define $Q_{1}, Q_{12}, Q_{2}, P_{1}$, $P_{12}, P_{2}$ by (A.14)-(A.16). Using (4.4), (4.11) and (4.12) it is easy to verify (A.6) and (A.7). Finally, subsiitute the definitions of $Q, P, \hat{Q}, \hat{P}, G, \Gamma$ and $\tau$ into (4.14)-(4.17) along with (4.2), (4.3), (4.6) and (4.7) to obtain (3.15) and (A.4). Finally, note that

$$
\mathcal{Q}=\left[\begin{array}{cc}
Q & 0_{n \times n_{e}} \\
0_{n_{e} \times n} & 0_{n_{e}}
\end{array}\right]+\left[\begin{array}{c}
I_{n} \\
\Gamma
\end{array}\right] \hat{Q}\left[\begin{array}{ll}
I_{n} & \Gamma^{\mathrm{T}}
\end{array}\right]
$$

which shows that $Q \geq 0$.

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# LQG Control with an $H_{\infty}$ Performance Bound: A Riccati Equation Approach 

by

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#### Abstract

An LQG control-design problem involving a constraint on $H_{\infty}$ disturbance attenuation is considered. The $H_{\infty}$ performance constraint is embedded within the optimization process by replacing the covariance Lyapunov equation by a Riccati equation whose solution leads to an upper bound on $L_{2}$ performance. In contrast to the pair of separated Riccati equations of standard LQG theory, the $H_{\infty}$-constrained gains are given by a coupled system of three nomdified Riccati equations. The coupling illustrates the breakdown of the separation principle for the $H_{\infty}$-constrained problem. Both full- and reduced-order design problems are considered with an $H_{\infty}$ attenuation constraint involving both state and control variables. An algorithm is developed for the full-order design problem and ilustrative numerical results are given.


Supported in part by the Air Force Office of Scientific Research under contracts F49620-86-C0002 and F49620-87-C-0108.

## 1. Introduction

The fundamental differences between Wiener-Hopf-Kalman (WHK) control design (for example, LQG theory [1]) and $H_{\infty}$ control theory [2-4] can be traced to the modeling and treatment of uncertain exogenous disturbances. As explained by Zames in the classic paper [2], LQG design is based upon a stochastic noise disturbance model possessing a fixed covariance (power spectral density), while $H_{\infty}$ theory is predicated on a deterministic disturbance model consisting of bounded power (square-integrable) signals. Since LQG design utilizes a quadratic cost criterion, it follows from Plancherel's theorem that WHK theory strives to minimize the $L_{2}$ norm of the closed-loop frequency response, while $H_{\infty}$ theory seeks to minimize the worst-case attenuation. For systems with poorly modeled disturbances which may possess significant power within arbitrarily small bandwidths, $H_{\infty}$ is clearly appropriate, while, for systems with well-known disturbance power spectral densities, WHK design may be less conservative.

In addition to the fact that $H_{\infty}$ design embodies many classical design objectives [5], it also presents a natural tool for modeling plant uncertainty in terms of normed $H_{\infty}$ plant neighborhoods. In contrast, the $H_{2}$ topology has been shown in [6] to be too weak for a practical robustness theory, while the $H_{\infty}$ norm is not only suitable for robust stabilization but is also conveniently submultiplicative. Within the WHK state space theory, however, the appropriate robustness model appears not to be a nonparametric normed plant neighborhood as in $H_{\infty}$ theory, but rather a parametric uncertainty model. The principal technique for capturing the effects of real parameters within state space models is Lyapunov function theory (see, e.g., [7-16] and the references therein). Such structured uncertainties are difficult to capture nonconservatively within $H_{\infty}$ theory except with specialized refinements ([17]).

In spite of the fundamental differences between WHK design and $H_{\infty}$ theory, a significant connection was discovered in [18]. Specifically, Petersen observed that a modified algebraic Riccati equation developed for parameter-robust full-state-feedback control can also be utilized to yield controllers satisfying $H_{\infty}$ disturbance attenuation bounds. This relationship was further explored in [19] where it was shown that Petersen's modified Riccati equation effectively yields the $H_{\infty}$ optimal full-state-feedback controller. This result is based upon the fact that quadratic stability (i.e., stability with a quadratic Lyapunov function) of the system

$$
\dot{x}=(A+D F E) x, \quad \bar{\sigma}(F) \leq 1,
$$

is quivalent to the stability of $A$ and the small gain condition

$$
\left\|E(s I-A)^{-1} D\right\|_{\infty}<1
$$

The results of [19] thus solve the Standard Problem considered in $[3,4]$ for the static full-statefeedback case.

The extension of these results to the Standard Problem for dynamic output-feedback compensation, however, was not given in $[18,19]$. Within the realm of quadratic robust stabilization, the dynamic output-feedback problem was addressed in [7]. The results of [7] involve a pair of decoupled modified Riccati equations along with an auxiliary inequality. Using different techniques, a mors complete solution was obtained in $[13,14]$ involving a coupled system of three modified Riccati equations for full-order dynamic compensation and a coupled system of four modified Riccati and Lyapunov equations in the fixed-order (i.e., reduced-order) case as in [20]. The results of $[13,14]$ thus raise the following question: What is the relevance of this system of coupled design equations to the problem of $H_{\infty}$ disturbance attenuation via fixed-order compensation?

To begin to address this question the goal of the present paper is to develop a design methodology for fixed-order $L_{2}$ optimal control which includes as a design constraint a bound on $H_{\infty}$ disturbance attenuation. There are three principal motivations for developing such a theory. First, the results of $[18,19]$ present full-state-feedback controllers whose form is directly analogous to the standard LQR solution. However, no $L_{2}$ interpretation was provided in $[18,19]$ to explain this similarity. The present paper, however, provides an $L_{2}$ interpretation within the context of an $H_{\infty}$ design constraint. A novel feature of this mathematical formulation is the dual interpretation of the disturbances. That is, within the context of $L_{2}$ optimality the disturbances are interpreted as white noise signals while, simultaneously, for the purpose of $H_{\infty}$ attenuation the very same disturbance signals have the alternative interpretation of deterministic $L_{2}$ functions. This dual interpretation is unique to the present paper since stochastic modeling does not play a role in $[18,19]$. We also note recent results obtained in [21] which essentially show that the $H_{2}$ plant topology can be induced by imposing $L_{2}$ and $L_{\infty}$ topologies on the disturbance and output spaces, respectively. For further investigation into the relationships between $L_{2}$ and $H_{\infty}$ control, see [22].

The second motivation for our approach is the simultaneous treatment of both $L_{2}$ and $H_{\infty}$ performance criteria which quantitatively demonstrates design tradeoffs. Specifically, in order to enforce the $H_{\infty}$ constraint we derive an upper bound for the $L_{2}$ criterion. Minimization of this upper bound shows that the enforcement of an $H_{\infty}$ disturbance attenuation constraint leads directly
to an increase in the $L_{2}$ performance criterion. Although it would be preferable to minimize the $L_{\mathbf{2}}$ criterion directly, this problem will be considered in a future paper.

The third motivation for our approach is to provide a characterization of fixed-order dynamic output-feedback controllers yielding specified disturbance attenuation. Existing optimal $H_{\infty}$ design methods tend to yield relatively high-order controllers. Intuitively, solving the fixed-order design equations for progressively smaller $H_{\infty}$ disturbance attenuation constraints should, in the limit, yield an $H_{\infty}$ optimal controller over the class of fixed-order stabilizing controllers. Although our main result gives sufficient conditions, we also state hypotheses under which these conditions are also necessary (Proposition 4.1). It should also be noted that the inherent coupling among the modified Riccati equations shows that the classical separation principle of LQG theory is not valid for the $H_{\infty}$-constrained full- and reduced-order design problems.

Of course, since the present paper addresses the problem of fixed-order dynamic compensation, previous full-state feedback results $([18,19])$ are not obtainable as a special case. However, applying the approach of the present paper to the full-state feedback problem yields results which are similar to those of $[18,19]$.

Besides establishing connections with robust stabilizability in state space systems, an immediate benefit of the modified Riccati equation characterization of $H_{\infty}$-optimal controllers is the opportunity to develop novel computational algorithms for controller synthesis. To this end a continuation algorithm has been developed for solving the coupled system of three modified Riccati equations. In a numerical study (see Section 8) we have demonstrated convergence of the algorithm and reasonable computational efficiency. Homotopy methods were suggested for the coupled Riccati equations because of their demonstrated effectiveness in related problems which also involve coupled modified Riccati equations [23-25]. Since $H_{\infty}$ control problems are solvable by established numerical methods [4], it should be stressed that the objective of these numerical studies is not to make direct comparisons with existing $H_{\infty}$ synthesis algorithms, but only to demonstrate solvability of the coupled modified Riccati equations.

The contents of the paper are as follows. After presenting notation at the end of this section, the statement of the $H_{\infty}$-Constrained LQG Control Problem is given in Section 2. The principal result of this section (Lemma 2.1) shows that if the algebraic Lyapunov equation for the closed-loop covariance is replaced by a modified Riccati equation possessing a nonnegative-definite solution, then the closed-loop system is asymptotically stable, the $H_{\infty}$ disturbance attenuation constraint is
satisfied, and the $L_{2}$ performance is bounded above by an "auxiliary" cost function. The problem of determining compensator gains which minimize this upper bound subject to the Riccati equation constraint is considered in Section 3 as the Auxiliary Minimization Problem. Necessary conditions for the Auxiliary Minimization Problem (Theorem 3.1) are given in the form of a coupled system of three modified Riccati equations. In Section 4 the necessary conditions of Theorem 3.1 are combined with Lemma 2.1 to yield sufficient conditions for closed-loop stability, $H_{\infty}$ disturbance attenuation, and bounded $L_{2}$ performance. In Section 5 we specialize the results of Section 3 to the case in which the LQG weights are equal to the $H_{\infty}$ weights. This serves to equalize the $L_{2}$ and $H_{\infty}$ design aspects and, through a series of transformations, the results of Section 3 are recast in a simpler form. These results also clarify connections with [26]. To achieve further design flexibility, the reduced-order control-design problem is considered in Section 6. A simplified qualitative analysis of the full-order design equations is given in Section 7 to highlight important features with regard to existence and multiplicity of solutions. Finally, a numerical algorithm is presented in Section 8 along with illustrative numerical results. A series of designs are obtained to illustrate tradeofs between the $L_{2}$ and $H_{\infty}$ aspects and the conservatism of the $L_{2}$ performance bound. Although in the present paper the numerical results are limited to the case of full-order dynamic compensation, reduced-order designs have been obtained in [27] using Theorem 6.1.

## Notation. Note: All matrices have real entries

$\mathbb{R}, \mathbb{R}^{r \times s}, \mathbb{R}^{r}, \mathbb{E}$
$I_{r},()^{T}, 0_{r \times s}, 0_{r}$
$\operatorname{tr}, \rho()$
$\mathbb{S}^{r}, \mathbb{N}^{r}, \mathbb{P}^{r}$
$Z_{1} \leq Z_{2}, \quad Z_{1}<Z_{2}$
$n, m, \ell, n_{c}, p, q, q_{\infty} ; \tilde{n}$
$x, u, y, x_{c}, \tilde{x}$
$\tilde{\boldsymbol{x}}$
$A, B, C$
$A_{c}, B_{c}, C_{c}$
$\tilde{A}$
$w(\cdot)$
$D_{2}, D_{2}$
real numbers, $r \times s$ real matrices, $\mathbb{R}^{r \times 1}$, expected value
$r \times r$ identity matrix, transpose, $r \times s$ zero matrix, $0_{r \times r}$
trace, spectral radius
$r \times r$ symmetric, nonnegative-definite, positive-definite matrices
$Z_{2}-Z_{1} \in \mathbb{N}^{r}, \quad Z_{2}-Z_{1} \in \mathbb{P}^{r}, \quad Z_{1}, Z_{2} \in \mathbb{S}^{r}$
positive integers; $n+n_{c}\left(n_{c} \leq n\right)$
$n, m, \ell, n_{c}, \tilde{n}$-dimensional vectors
$\left[\begin{array}{c}x \\ x_{c}\end{array}\right]$
$n \times n, n \times m, \ell \times n$ matrices
$n_{c} \times n_{c}, n_{c} \times \ell, m \times n_{c}$ matrices
$\left[\begin{array}{cc}A & B C_{c} \\ B_{c} C & A_{c}\end{array}\right]$
p-dimensional standard white noise
$n \times p, \ell \times p$ matrices; $D_{1} D_{2}^{T}=0$

| $V_{1}, V_{2}$ | $D_{1} D_{1}^{T}, D_{2} D_{2}^{T} ; \quad V_{2} \in \mathbb{P}^{\ell}$ |
| :---: | :---: |
| D | $\left[\begin{array}{c}D_{1} \\ B_{c} D_{2}\end{array}\right]$ |
| $\tilde{V}$ | $\left[\begin{array}{cc}V_{1} & 0_{n \times n_{0}} \\ 0_{n_{0} \times n} & B_{c} V_{2} B_{c}^{T}\end{array}\right]$ |
| $E_{1}, E_{2}$ | $q \times n, q \times m$ matrices; $E_{1}^{T} E_{2}=0$ |
| $\tilde{E}$ | $\left[\begin{array}{ll}E_{2} & E_{2} C_{c}\end{array}\right]$ |
| $R_{1}, R_{2}$ | $E_{1}^{T} E_{1}, E_{2}^{T} E_{2} ; \quad R_{2} \in \mathbb{P}^{m}$ |
| $\tilde{R}$ | $\left[\begin{array}{cc}R_{1} & 0_{n \times n_{c}} \\ 0_{n_{0} \times n} & C_{c}^{T} R_{2} C_{c}\end{array}\right]=\tilde{E}^{T} \tilde{E}$ |
| $E_{1 \infty}, E_{2 \infty}$ | $q_{\infty} \times n, q_{\infty} \times m$ matrices; $\quad E_{1 \infty}^{T} E_{2 \infty}=0$ |
| $\tilde{E}_{\infty}$ | $\left[\begin{array}{lll}E_{1 \infty} & E_{2 \infty} C_{c}\end{array}\right]$ |
| $R_{1 \infty}, R_{2 \infty}$ | $E_{1 \infty}^{T} E_{1 \infty}, E_{2 \infty}^{T} E_{2 \infty}$ |
| $\tilde{R}_{\infty}$ | $\left[\begin{array}{cc}R_{1 \infty} & 0_{n \times n_{0}} \\ 0_{n_{e} \times n} & C_{c}^{T} R_{2 \infty} C_{c}\end{array}\right]=\tilde{E}_{\infty}^{T} \tilde{E}_{\infty}$ |
| $\boldsymbol{\Sigma}, \bar{\Sigma}$ | $B R_{2}^{-1} B^{T}, C^{T} V_{2}^{-1} C$ |
| $\beta, \gamma$ | nonnegative constant, positive constant |

## 2. Statement of the Problem

In this section we introduce the LQG dynamic output-feedback control problem with constrained $H_{\infty}$ disturbance attenuation between the plant and sensor disturbances and the state and control variables. Without the $L_{2}$ performance criterion the problem considered here essentially corresponds to the Standard Problem of [3,4]. For simplicity we restrict our attention to controllers of order $n_{c}=n$ only, i.e., controllers whose order is equal to the dimension of the plant. This constraint is removed in Section 6 where controllers of reduced order are considered. Hence, throughout Sections 2-5 the controller dimension $n_{c}$ and closed-loop plant dimension $\tilde{n} \triangleq n+n_{c}$ should be interpreted as $n$ and $2 n$, respectively. Controllers of order greater than $n$ are generally of less interest in practice and thus are not considered in this paper.
$\mathrm{H}_{\infty}$-Constrained LQG Control Problem. Given the nth-order stabilizable and detectable plant

$$
\begin{gather*}
\dot{x}(t)=A x(t)+B u(t)+D_{1} w(t),  \tag{2.1}\\
y(t)=C x(t)+D_{2} w(t), \tag{2.2}
\end{gather*}
$$

determine an nth-order dynamic compensator

$$
\begin{gather*}
\dot{x}_{c}(t)=A_{c} x_{c}(t)+B_{c} y(t),  \tag{2.3}\\
u(t)=C_{c} x_{c}(t), \tag{2.4}
\end{gather*}
$$

which satisfies the following design criteria:
(i) the closed-loop system (2.1)-(2.4) is asymptotically stable, i.e., $\tilde{A}$ is asymptotically stable;
(ii) the $q_{\infty} \times p$ closed-loop transfer function

$$
\begin{equation*}
H(s) \triangleq \tilde{E}_{\infty}\left(s I_{\tilde{n}}-\tilde{A}\right)^{-1} \tilde{D} \tag{2.5}
\end{equation*}
$$

from $w(t)$ to $E_{1 \infty} x(t)+E_{2 \infty} u(t)$ satisfies the constraint

$$
\begin{equation*}
\|H(s)\|_{\infty} \leq \gamma, \tag{2.6}
\end{equation*}
$$

where $\boldsymbol{\gamma} \boldsymbol{>} \mathbf{0}$ is a given constant; and
(iii) the performance functional

$$
\begin{equation*}
J\left(A_{c}, B_{c}, C_{c}\right) \triangleq \lim _{t \rightarrow \infty} \mathbb{E}\left[x^{T}(t) R_{1} x(t)+u^{T}(t) R_{2} u(t)\right] \tag{2.7}
\end{equation*}
$$

is minimized.
Note that the closed-loop system (2.1)-(2.4) can be written as

$$
\begin{equation*}
\dot{\tilde{x}}(t)=\tilde{A} \tilde{x}(t)+\tilde{D} w(t) \tag{2.8}
\end{equation*}
$$

and that (2.7) becomes

$$
\begin{equation*}
J\left(A_{c}, B_{c}, C_{c}\right)=\lim _{t \rightarrow \infty} \mathbb{E}\left[(\tilde{E} \tilde{x}(t))^{T}(\tilde{E} \tilde{x}(t))\right]=\lim _{t \rightarrow \infty} \mathbb{E}\left[\tilde{x}^{T}(t) \tilde{R} \tilde{x}(t)\right] \tag{2.9}
\end{equation*}
$$

Remark 2.1. Since ( $A, B, C$ ) is assumed to be stabilizable and detectable the set of nth-order stabilizing compensators is nonempty.

Remark 2.2. It is easy to show that the performance functional (2.7) is equivalent to the more familiar expression involving an averaged integral, i.e.,

$$
J\left(A_{c}, B_{c}, C_{c}\right)=\lim _{t \rightarrow \infty} \frac{1}{t} \mathbb{E}\left\{\int_{0}^{t}\left[x^{T}(s) R_{1} x(s)+u^{T}(s) R_{2} u(s)\right] d s\right\} .
$$

Remark 2.3. For convenience we assume $D_{1} D_{2}^{T}=0$, which effectively implies that the plant disturbance and sensor noise are uncorrelated.

Remark 2.4. One may wish to consider a general $L_{2}$ optimization problem of the form $\min \|T-U Q V\|_{2}$, where $Q$ is a parameterization of stabilizing controllers. In this case, without a constraint on the MacMillan degree of $Q$, it may be possible to satisfy (2.6) with smaller values of $\gamma$.

Note that the problem statement involves both $L_{2}$ and $H_{\infty}$ performance weights. In particular, the matrices $R_{1}$ and $R_{2}$ are the $L_{2}$ weights for the state and control variables. By introducing $L_{2}$ weighted variables

$$
z(t)=E_{1} x(t), \quad v(t)=E_{2} u(t),
$$

the cost (2.7) can be written as

$$
J\left(A_{c}, B_{c}, C_{c}\right)=\lim _{t \rightarrow \infty} \mathbb{E}\left[z^{T}(t) z(t)+v^{T}(t) v(t)\right]
$$

For convenience we thus define $R_{1} \triangleq E_{1}^{T} E_{1}$ and $R_{2} \triangleq E_{2}^{T} E_{2}$ which appear in subsequent expressions. Although an $L_{2}$ cross-weighting term of the form $2 x(t)^{T} R_{12} u(t)$ can also be included, we shall not do so here to facilitate the presentation.

For the $H_{\infty}$ performance constraint, the transfer function (2.5) involves weighting matrices $E_{1 \infty}$ and $E_{2 \infty}$ for the state and control variables. The matrices $R_{1 \infty} \triangleq E_{1 \infty}^{T} E_{1 \infty}$ and $R_{2 \infty} \triangleq E_{2 \infty}^{T} E_{2 \infty}$ are thus the $H_{\infty}$ counterparts of the $L_{2}$ weights $R_{1}$ and $R_{2}$. Although we do not require that $R_{1 \infty}$ and $R_{2 \infty}$ be equal to $R_{1}$ and $R_{2}$, we shall require in the next section that $R_{2 \infty}=\beta^{2} R_{2}$, where the nonnegative scalar $\beta$ is a design variable. Finally, the condition $E_{1 \infty}^{T} E_{2 \infty}=0$ precludes an $H_{\infty}$ cross-weighting term which again facilitates the presentation.

Before continuing it is useful to note that if $\tilde{A}$ is asymptotically stable for a given compensator ( $A_{c}, B_{c}, C_{c}$ ) then the performance (2.7) is given by

$$
\begin{equation*}
J\left(A_{c}, B_{c}, C_{c}\right)=\operatorname{tr} \tilde{Q} \tilde{R} \tag{2.10}
\end{equation*}
$$

where the steady-state closed-loop state covariance defined by

$$
\begin{equation*}
\tilde{Q} \triangleq \lim _{t \rightarrow \infty} \mathbb{E}\left[\tilde{x}(t) \tilde{x}^{T}(t)\right] \tag{2.11}
\end{equation*}
$$

satisfies the $\tilde{n} \times \tilde{n}$ algebraic Lyapunov equation

$$
\begin{equation*}
0=\tilde{A} \tilde{Q}+\tilde{Q} \tilde{A}^{T}+\tilde{V} \tag{2.12}
\end{equation*}
$$

Remark 2.5. Using (2.10) and (2.12) it can de shuwn that the $L_{2}$ cost criterion (2.7) can be written in terms of the $L_{2}$ norm of the impulse response of the closed-loop system. Specifically, writing $\tilde{Q}$ satisfying (2.12) as

$$
\tilde{Q}=\int_{0}^{\infty} e^{\tilde{\lambda} t} \tilde{V} e^{\tilde{\Lambda}^{F_{t}}} d t,
$$

(2.10) becomes

$$
J\left(A_{c}, B_{c}, C_{c}\right)=\int_{0}^{\infty}\left\|\tilde{E} e^{\tilde{A} t} \tilde{D}\right\|_{F}^{2} d t
$$

where $\|\cdot\|_{F}$ denotes the Frobenius matrix norm.
The key step in enforcing the disturbance attenuation constraint (2.6) is to replace the algebraic Lyapunov equation (2.12) by an algebraic Riccati equation which overbounds the closed-loop steady-state covariance. Justification for this technique is provided by the following result.

Lemma 2.1. Let ( $A_{c}, B_{c}, C_{c}$ ) be given and assume there exists $\mathcal{Q} \in \mathbb{R}^{\bar{n} \times \bar{n}}$ satisfying

$$
\begin{equation*}
Q \in \mathbb{N}^{n} \tag{2.13}
\end{equation*}
$$

and

$$
\begin{equation*}
0=\tilde{A} Q+Q \tilde{A}^{T}+\gamma^{-2} Q \tilde{R}_{\infty} Q+\tilde{V} \tag{2.14}
\end{equation*}
$$

Then

$$
\begin{equation*}
\left(\tilde{A},\left[\gamma^{-2} Q \tilde{R}_{\infty} Q+\tilde{V}\right]^{\frac{1}{2}}\right) \text { is stabilizable } \tag{2.15}
\end{equation*}
$$

if and only if

$$
\begin{equation*}
\tilde{A} \text { is asymptotically stable. } \tag{2.16}
\end{equation*}
$$

In this case,

$$
\begin{equation*}
\|H(s)\|_{\infty} \leq \boldsymbol{\gamma} \tag{2.17}
\end{equation*}
$$

and

$$
\begin{equation*}
\tilde{Q} \leq 0 . \tag{2.18}
\end{equation*}
$$

Consequently,

$$
\begin{equation*}
J\left(A_{a}, B_{c}, C_{c}\right) \leq J\left(A_{c}, B_{c}, C_{c}, \mathcal{Q}\right) \tag{2.19}
\end{equation*}
$$

where

$$
\begin{equation*}
J\left(A_{c}, B_{c}, C_{c}, Q\right) \triangleq \operatorname{tr} Q \tilde{R} \tag{2.20}
\end{equation*}
$$

Proof. Using the assumed existence of a nonnegative-definite solution to (2.14) and the stabilizability condition (2.15), it follows from the dual of Lemma 12.2 of [28] that $\tilde{A}$ is asymptotically
stable. Conversely, if $\tilde{A}$ is asymptotically stable then (2.15) holds. The proof of (2.17) follows from a standard manipulation of (2.14); for details see Lemma 1 of [29]. To prove (2.18) subtract (2.12) from (2.14) to obtain

$$
\begin{equation*}
0=\tilde{A}(Q-\tilde{Q})+(Q-\tilde{Q}) \tilde{A}^{T}+\gamma^{-2} Q \tilde{R}_{\infty} \mathcal{Q} \tag{2.21}
\end{equation*}
$$

which, since $\tilde{A}$ is asymptotically stable, is equivalent to

$$
\begin{equation*}
Q-\tilde{Q}=\int_{0}^{\infty} e^{\bar{A} t}\left[\gamma^{-2} Q \tilde{R}_{\infty} \mathcal{Q}\right] e^{\bar{T}^{T} t} d t \geq 0 \tag{2.22}
\end{equation*}
$$

Finally, (2.19) follows immediately from (2.18).
Remark 2.6. Note that (2.15) is actually a closed-loop "disturbability" condition which is not concerned with control as such. This condition guarantees that the system does not possess unstable undisturbed modes. Of course, if $\tilde{V}$ is positive definite or ( $\tilde{A}, \tilde{D}$ ) is controllable, then (2.15) is satisfied.

Lemma 2.1 shows that the $H_{\infty}$ disturbance attenuation constraint is automatically enforced when a nonnegative-definite solution to (2.14) is known to exist. Furthermore, the solution $Q$ provides an upper bound for the actual closed-loop state covariance $\bar{Q}$ along with a bound on the $L_{2}$ performance criterion. Next, we present a partial converse of Lemma 2.1 which guarantees the existence of a nonnegative-definite solution to (2.14) when (2.17) is satisfied.

Lemma 2.2. Let ( $A_{c}, B_{c}, C_{c}$ ) be given, suppose $\tilde{A}$ is asymptotically stable, and assume the disturbance attenuation constraint (2.17) is satisfied. Then there exists a unique nonnegativedefinite solution $\mathcal{Q}$ satisfying (2.14). Furthermore, $\left(\tilde{A}+\gamma^{-2} Q \tilde{R}_{\infty}, \tilde{V}^{\frac{1}{2}}\right)$ is stabilizable if and only if $\tilde{A}+\gamma^{-2} \varrho \tilde{R}_{\infty}$ is asymptotically stable.

Proof. The result is an immediate consequence of [30], using Theorems 3 and 2, pp. 150 and 167, along with the dual version of Lemma 12.2 of [28].

Remark 2.7. To further, clarify the relationships between the $L_{2}$ and $H_{\infty}$ aspects of the problem, we note that the closed-loop system can be represented by two possibly different transfer functions. Specifically, with respect to the $L_{2}$ cost criterion, the closed-loop transfer function between disturbances and controlled variables is given by the triple $(\tilde{A}, \tilde{D}, \tilde{E})$ while for the $H_{\infty}$ constraint the closed-loop transfer function (2.5) corresponds to the triple ( $\tilde{A}, \tilde{D}, \tilde{E}_{\infty}$ ).

Finally, it can be shown that the closed-loop Riccati equation (2.14) also enforces a constraint on the norm of the Hankel operator corresponding to the closed-loop system ( $\tilde{A}, \tilde{D}, \tilde{E}_{\infty}$ ) when $Q$
is positive definite. Thus, let $\tilde{P} \in \mathbb{N}^{\tilde{n}}$ denote the solution to

$$
\begin{equation*}
0=\tilde{A}^{T} \tilde{P}+\tilde{P} \tilde{A}+\tilde{R}_{\infty} \tag{2.23}
\end{equation*}
$$

and note that $\tilde{P}$ and $\tilde{Q}$ (satisfying (2.12)) are the observability and controllability Gramians, respectively, of the system ( $\tilde{A}, \tilde{D}, \tilde{E}_{\infty}$ ). As shown in [31], the norm of the Hankel operator corresponding to $\left(\tilde{A}, \tilde{D}, \tilde{E}_{\infty}\right)$ is given by $\lambda_{\text {max }}^{\frac{1}{2}}(\tilde{P} \tilde{Q})$.

Proposition 2.1. Suppose there exists $\mathcal{Q} \in \mathbb{P}^{\bar{n}}$ satisfying (2.14) and such that (2.15) or, equivalently, (2.16) holds. Then

$$
\begin{equation*}
\lambda_{\text {max }}^{\frac{1}{2}}(\tilde{P} \tilde{Q}) \leq \gamma . \tag{2.24}
\end{equation*}
$$

Proof. Since $\mathcal{Q}$ is assumed to be invertible, (2.14) is equivalent to

$$
\begin{equation*}
0=\gamma^{2} \tilde{A}^{T} Q^{-1}+\gamma^{2} Q^{-1} \tilde{A}+\gamma^{2} Q^{-1} \tilde{V} Q^{-1}+\tilde{R}_{\infty} \tag{2.25}
\end{equation*}
$$

Subtracting (2.23) from (2.25) shows that $\gamma^{2} Q^{-1}-\tilde{P} \geq 0$, or, equivalently, $\gamma^{2} I_{\tilde{n}} \geq Q^{\frac{1}{2}} \tilde{P} Q^{\frac{1}{3}}$. Thus,

$$
\gamma^{2} \geq \lambda_{\max }\left(\mathcal{Q}^{\frac{1}{2}} \tilde{P} Q^{\frac{1}{3}}\right)=\lambda_{\max }\left(\tilde{P}^{\frac{1}{2}} Q \tilde{P}^{\frac{1}{2}}\right) \geq \lambda_{\max }\left(\tilde{P}^{\frac{1}{2}} \tilde{Q} \tilde{P}^{\frac{1}{2}}\right)=\lambda_{\max }(\tilde{P} \tilde{Q})
$$

which yields (2.24).

## 3. The Auxiliary Minimization Problem and Necessary Conditions for Optimality

As discussed in the previous section, the replacement of (2.12) by (2.14) enforces the $H_{\infty}$ disturbance attenuation constraint and yields an upper bound for the $L_{2}$ forformance criterion. That is, given a compensator ( $A_{c}, B_{c}, C_{c}$ ) for which there exists a nonnegative-definite solution to (2.14), the actual $L_{2}$ performance $J\left(A_{c}, B_{c}, C_{c}\right)$ of the compensator is guaranteed to be no worse than the bound given by $J\left(A_{c}, B_{c}, C_{c}, \mathcal{Q}\right)$. Hence, $J\left(A_{c}, B_{c}, C_{c}, \mathcal{Q}\right)$ can be interpreted as an auxiliary cost which leads to the following mathematical programming problem.

Auxiliary Minimization Problem. Determine ( $\left.A_{c}, B_{c}, C_{c}, Q\right)$ which minimizes $J\left(A_{c}, B_{c}, C_{c}\right.$, Q) subject to (2.13) and (2.14).

It follows from Lemma 2.1 that the satisfaction of (2.13) and (2.14) along with the generic condition (2.15) leads to 1) closed-loop stability; 2) prespecified $H_{\infty}$ performance attenuation; and 3) an upper bound for the $L_{2}$ performance criterion. Hence it remains to determine ( $A_{c}, B_{c}, C_{c}$ ) which minimizes $J\left(A_{c}, B_{c}, C_{c}, Q\right)$ and thus provides an optimized bound for the actual $L_{2}$ performance
$J\left(A_{c}, B_{c}, C_{c}\right)$. Rigorous derivation of the necessary conditions for the Auxiliary Minimization Problem requires additional technical assumptions. Specifically, we restrict ( $A_{c}, B_{c}, C_{c}, \mathcal{Q}$ ) to the open set

$$
\begin{array}{r}
\mathcal{X} \triangleq\left\{\left(A_{c}, B_{c}, C_{c}, Q\right): Q \in \mathbb{P}^{\bar{n}}, \tilde{A}+\gamma^{-2} Q \tilde{R}_{\infty} \text { is asymptotically stable },\right. \\
\text { and } \left.\left(A_{c}, B_{c}, C_{c}\right) \text { is controllable and observable }\right\} . \tag{3.1}
\end{array}
$$

Remark 3.1. The set $X$ constitutes sufficient conditions under which the Lagrange multiplier technique is applicable to the Auxiliary Minimization Problem. Specifically, the requirement that $Q$ be positive definite replaces (2.13) by an open set constraint, the stability of $\tilde{A}+\gamma^{-2} Q \tilde{R}_{\infty}$ serves as a normality condition, and ( $A_{c}, B_{c}, C_{c}$ ) minimal is a nondegeneracy condition. Note that the stabilizability condition (2.15) and stability condition (2.16) play no role in determining solutions of the Auxiliary Minimization Problem.

The following result presents the necessary conditions for optimality in the Auxiliary Minimization Problem. The proof of this result is given in Appendix A as a special case of the corresponding result for reduced-order dynamic compensation considered in Section 6. As mentioned previously, we assume that $R_{2 \infty}=\beta^{2} R_{2}$. Furthermore, for arbitrary $\hat{Q}, P \in \mathbb{N}^{n}$ define

$$
\begin{equation*}
S \triangleq\left(I_{n}+\beta^{2} \gamma^{-2} \hat{Q} P\right)^{-1} \tag{3.2}
\end{equation*}
$$

Since the eigenvalues of $\hat{Q} P$ coincide with the eigenvalues of the nonnegative-definite matrix $P^{\frac{1}{2}} \hat{Q} P^{\frac{1}{2}}$, it follows that $\hat{Q} P$ has nonnegative eigenvalues. Thus, the eigenvalues of $I_{n}+\beta^{2} \gamma^{-2} \hat{Q} P$ are all greater than one so that the above inverse exists.

Theorem 3.1. If $\left(A_{c}, B_{c}, C_{c}, \mathcal{Q}\right) \in \mathcal{X}$ solves the Auxiliary Minimization Problem then there exist $Q, P, \hat{Q} \in \mathbb{N}^{n}$ such that

$$
\begin{align*}
A_{c} & =A-Q \bar{\Sigma}-\Sigma P S+\gamma^{-2} Q R_{100},  \tag{3.3}\\
B_{c} & =Q C^{T} V_{2}^{-1}  \tag{3.4}\\
C_{c} & =-R_{2}^{-1} B^{T} P S  \tag{3.5}\\
Q & =\left[\begin{array}{cc}
Q+\hat{Q} & \hat{Q} \\
\hat{Q} & \hat{Q}
\end{array}\right], \tag{3.6}
\end{align*}
$$

and such that $Q, P, \hat{Q}$ satisfy

$$
\begin{equation*}
0=A Q+Q A^{T}+V_{1}+\gamma^{-2} Q R_{1 \infty} Q-Q \sum Q, \tag{3.7}
\end{equation*}
$$

$$
\begin{align*}
0= & \left(A+\gamma^{-2}[Q+\hat{Q}] R_{1 \infty}\right)^{T} P+P\left(A+\gamma^{-2}[Q+\hat{Q}] R_{1 \infty}\right)+R_{1}-S^{T} P \Sigma P S,  \tag{3.8}\\
0= & \left(A-\Sigma P S+\gamma^{-2} Q R_{1 \infty}\right) \hat{Q}+\hat{Q}\left(A-\Sigma P S+\gamma^{-2} Q R_{1 \infty}\right)^{T} \\
& +\gamma^{-2} \hat{Q}\left(R_{1 \infty}+\beta^{2} S^{T} P \Sigma P S\right) \hat{Q}+Q \Sigma Q . \tag{3.9}
\end{align*}
$$

Furthermore, the auxiliary cost is given by

$$
\begin{equation*}
J\left(A_{c}, B_{c}, C_{c}, \mathcal{Q}\right)=\operatorname{tr}\left[(Q+\hat{Q}) R_{1}+\hat{Q} S^{T} P \Sigma P S\right] \tag{3.10}
\end{equation*}
$$

Conversely, if there exist $Q, P, \hat{Q} \in \mathbb{N}^{n}$ satisfying (3.7)-(3.9), then ( $A_{c}, B_{c}, C_{c}, \mathcal{Q}$ ) given by (3.3)(3.6) satisfies (2.13) and (2.14) with auxiliary cost (2.20) given by (3.10).

Remark 3.2. If $Q$ and $\hat{Q}$ are nonnegative definite then the fact that (2.13) is satisfied can easily be seen by writing $Q$ as

$$
\mathcal{Q}=\left[\begin{array}{ll}
Q & 0_{n} \\
0_{n} & 0_{n}
\end{array}\right]+\left[\begin{array}{l}
\hat{Q}^{\frac{1}{2}} \\
\hat{Q}^{\frac{1}{2}}
\end{array}\right]\left[\begin{array}{l}
\hat{Q}^{\frac{1}{2}} \\
\hat{Q}^{\frac{1}{2}}
\end{array}\right]^{T} .
$$

Remark 3.3. Setting $\beta=0$ or, equivalently, $E_{2 \infty}=0$, specializes Theorem 3.1 to the "cheap" $H_{\infty}$ control case in which $H_{\infty}$ attenuation between disturbances and controls is not constrained. In this case $S=I_{n}, \mathcal{Q}$ is given by (3.6), and

$$
\begin{align*}
& A_{c}=A-Q \tilde{\Sigma}-\Sigma P+\gamma^{-2} Q R_{1 \infty}  \tag{3.11}\\
& B_{c}=Q C^{T} V_{2}^{-1}  \tag{3.12}\\
& C_{c}=-R_{2}^{-1} B^{T} P \tag{3.13}
\end{align*}
$$

where $Q$ satisfies (3.7), and equations (3.8) and (3.9) become

$$
\begin{align*}
& 0=\left(A+\gamma^{-2}[Q+\hat{Q}] R_{1 \infty}\right)^{T} P+P\left(A+\gamma^{-2}\left[Q+\hat{Q} \mid R_{1 \infty}\right)+R_{1}-P \Sigma P,\right.  \tag{3.14}\\
& 0=\left(A-\Sigma P+\gamma^{-2} Q R_{1 \infty}\right) \hat{Q}+\hat{Q}\left(A-\Sigma P+\gamma^{-2} Q R_{1 \infty}\right)^{T}+\gamma^{-2} \hat{Q} R_{1 \infty} \hat{Q}+Q \hat{\Sigma} Q . \tag{3.15}
\end{align*}
$$

Finally, the auxiliary cost reduces to

$$
\begin{equation*}
J\left(A_{c}, B_{c}, C_{c}, \mathcal{Q}\right)=\operatorname{tr}\left[(Q+\hat{Q}) R_{1}+\hat{Q} P \Sigma P\right] \tag{3.16}
\end{equation*}
$$

Numerical solution of equations (3.7), (3.14) and (3.15) is discussed in Section 8.
Theorem 3.1 presents necessary conditions for the Auxiliary Minimization Problem which explicitly synthesize extremal controllers ( $A_{c}, B_{c}, C_{c}$ ). These necessary conditions comprise a system
of three modified algebraic Riccati equations in variables $Q, P$ and $\hat{Q}$. The $Q$ and $P$ equations are similar to the estimator and regulator Riccati equations of LQG theory, while the $\hat{Q}$ equation has no counterpart in the standard theory. Note that the $Q$ equation is decoupled from the $P$ and $\hat{\boldsymbol{Q}}$ equations and thus can be solved independently. The $P$ equation, however, depends on $Q$. Thus, regulator/estimator separation only holds in one direction which clearly shows that the certainty equivalence principle is no longer valid for the design problem under consideration. Furthermore, since the $P$ and $\hat{Q}$ equations are coupled, they must be solved simultaneously. Finally, note that if the $H_{\infty}$ disturbance attenuation constraint is sufficiently relaxed, i.e., $\boldsymbol{\gamma} \rightarrow \infty$, then the $P$ equation becomes decoupled from the $\hat{Q}$ equation and thus the $\hat{Q}$ equation becomes superfluous. Furthermore, the remaining $Q$ and $P$ equations separate and coincide with the standard LQG result.

## 4. Sufficient Conditions for $H_{\infty}$ Disturbance Attenuation

In this section we combine Lemma 2.1 with the converse of Theorem 3.1 to obtain our main result guaranteeing closed-loop stability, $H_{\infty}$ disturbance attenuation, and an optimized $L_{2}$ performance bound.

Theorem 4.1. Suppose there exist $Q, P, \hat{Q} \in \mathbb{N}^{n}$ satisfying (3.7)-(3.9). Then, with $\left(A_{c}, B_{c}, C_{c}, \mathcal{Q}\right)$ given by (3.3)-(3.6), $\left(\tilde{A},\left[\gamma^{-2} Q \tilde{R}_{\infty} \mathcal{Q}+\tilde{V}\right]^{\frac{1}{2}}\right)$ is stabilizable if and only if $\tilde{A}$ is asymptotically stable. In this case, the closed-loop transfer function $H(s)$ satisfies the $H_{\infty}$ attenuation constraint

$$
\begin{equation*}
\|H(s)\|_{\infty} \leq \boldsymbol{\gamma} \tag{4.1}
\end{equation*}
$$

and the $L_{2}$ performance criterion (2.7) satisfies the bound

$$
\begin{equation*}
J\left(A_{c}, B_{c}, C_{c}\right) \leq \operatorname{tr}\left[(Q+\hat{Q}) R_{1}+\hat{Q} S^{T} P \Sigma P S\right] . \tag{4.2}
\end{equation*}
$$

Proof. The converse portion of Theorem 3.1 implies that $Q$ given by (3.6) satisfies (2.13) and (2.14) with auxiliary cost given by (3.10). It now follows from Lemma 2.1 that the stabilizability condition (2.15) is equivalent to the asymptotic stability of $\tilde{A}$, the $H_{\infty}$ disturbance attenuation constraint (2.17) holds, and the performance bound (2.19), which is equivalent to (4.2), holds.

Remark 4.1. In applying Theorem 4.1 it is not actually necessary to check (2.15) which holds generically. Rather, it suffices to check the stability of $\tilde{A}$ directly which is guaranteed to be equivalent to (2.15).

In applying Theorem 4.1 the principal issue concerns conditions on the problem data under which the coupled Riccati equations (3.7)-(3.9) possess nonnegative-definite solutions. Clearly, for $\gamma$ sufficiently large, (3.7)-(3.9) approximate the standard LQG result so that existence and uniqueness are assured. The important case of interest, however, involves small $\boldsymbol{\gamma}$ so that significant $H_{\infty}$ disturbance attenuation is enforced. Thus, if (4.1) can be satisfied for a given $\boldsymbol{\gamma}>0$, it is of interest to know whether one such controller can be obtained by solving (3.7)-(3.9). Lemma 2.2 guarantees that (2.14) possesses a solution for any controller satisfying (4.1). Thus our sufficient condition will also be necessary so long as the Auxiliary Minimization Problem possesses at least one extremal over $X$. When this is the case we have the following immediate result.

Proposition 4.1. Let $\boldsymbol{\gamma}^{*}$ denote the infimum of $\|H(s)\|_{\infty}$ over all stabilizing nth-order dynamic compensators and suppose that the Auxiliary Minimization Problem has a solution for all $\boldsymbol{\gamma} \boldsymbol{\gamma} \boldsymbol{\gamma}^{*}$. Then for all $\boldsymbol{\gamma}>\boldsymbol{\gamma}^{*}$ there exist $Q, P, \hat{Q} \in \mathbb{N}^{n}$ satisfying (3.7)-(3.9).

Unlike the standard LQG result involving a pair of separated Riccati equations, the new result enforcing $H_{\infty}$ disturbance attenuation involves a nonstandard coupled system of three modifed Riccati equations. The asymmetry of these equations suggests the possibility of a dual result in which the modifications to the standard $P$ and $Q$ Riccati equations are reversed. Such a dual result will generally be different from Theorem 4.1 and thus will yield an improved bound for particular problems. This point was demonstrated in [16] for the problem of robust performance analysis. Due to space limitations, however, we give only a brief outline of the dual $H_{\infty}$ results. Note that $J\left(A_{c}, B_{c}, C_{c}\right)$ given by (2.10) is also given by

$$
\begin{equation*}
J\left(A_{c}, B_{c}, C_{c}\right)=\operatorname{tr} \tilde{P} \tilde{V} \tag{4.3}
\end{equation*}
$$

where $\tilde{P} \in \mathbb{N}^{n}$ is the unique solution to (2.23). Next, utilizing (4.3) in place of (2.10), the $H_{\infty}$ disturbance attenuation constraint (2.6) can now be enforced by replacing (2.23) by the Riccati equation

$$
\begin{equation*}
0=\tilde{A}^{T} P+P \tilde{A}+\gamma^{-2} \rho \tilde{V} P+\tilde{R}_{\infty} \tag{4.4}
\end{equation*}
$$

Note that (4.4) is merely the dual of (2.14). We also require the condition dual to (2.15) given by

$$
\begin{equation*}
\left(\left[\gamma^{-2} P \tilde{V} P+\tilde{R}_{\infty}\right]^{\frac{1}{2}}, \tilde{A}\right) \text { is detectable. } \tag{4.5}
\end{equation*}
$$

Once again, the sufficient conditions for $H_{\infty}$ disturbance attenuation involve a coupled system of three modified Riccati equations dual to (3.7)-(3.9). Similar remarks apply to the reduced-order case given by Theorem 6.1 below.

## 5. Alternative Forms of the Riccati Equations

The purpose of this section is to draw connections with recent results obtained in [26]. As shown in Theorem 4.1, the Riccati equations (3.7)-(3.9) provide sufficient conditions for explicitly synthesizing controllers $\left(A_{c}, B_{c}, C_{c}\right)$ satisfying an $H_{\infty}$ performance bound. In this section we specialize Theorem 3.1 by setting the LQG weights equal to the $H_{\infty}$ weights, i.e., $R_{1}=R_{1 \infty}$ and $\beta=1$. This specialization leads to considerable simplification by equalizing the $L_{2}$ and $H_{\infty}$ design aspects. In this case it turns out that the Riccati equations (3.7)-(3.9) can be transformed to simpler forms which are similar to the results given in [26]. To state the results we require some additional notation and a lemma concerning transformations of a pair of nonnegative-definite matrices.

Lemma 5.1. Let $P, Q \in \mathbb{N}^{n}$. Then the following statements hold.
(i) $P Q$ has nonnegative eigenvalues.
(ii) $\rho(P Q) \leq(<) \gamma^{-2}$ if and only if $P^{\frac{1}{2}} Q P^{\frac{1}{2}} \leq(<) \gamma^{-2} I_{n}$. Furthermore, if $P$ is positive definite then $\rho(P Q) \leq(<) \gamma^{-2}$ if and only if $Q \leq(<) \gamma^{-2} P^{-1}$.
(iii) Define $Z \triangleq P\left(I_{n}+\gamma^{-2} P Q\right)^{-1}$. Then $Z$ is nonnegative definite,

$$
\begin{equation*}
P\left(I_{n}+\gamma^{-2} Q P\right)^{-1}=\left(I_{n}+\gamma^{-2} P Q\right)^{-1} P \tag{5.1}
\end{equation*}
$$

and $\rho(Z Q)<\boldsymbol{\gamma}^{-2}$ or, equivalently, $Q^{\frac{1}{2}} Z Q^{\frac{1}{2}}<\boldsymbol{\gamma}^{-2} I_{n}$. Furthermore, $Z$ is positive definite if and only if $P$ is positive definite. In this case,

$$
\begin{equation*}
Z=\left(P^{-1}+\gamma^{-2} Q\right)^{-1} \tag{5.2}
\end{equation*}
$$

(iv) Suppose $\rho(P Q)<\gamma^{-2}$ and define $Y \triangleq P\left(I_{n}-\gamma^{-2} Q P\right)^{-1}$. Then $Y$ is nonnegative definite,

$$
\begin{equation*}
P\left(I_{n}-\gamma^{-2} Q P\right)^{-1}=\left(I_{n}-\gamma^{-2} P Q\right)^{-1} P \tag{5.3}
\end{equation*}
$$

and

$$
\begin{equation*}
Q P\left(I_{n}-\gamma^{-2} Q P\right)^{-1}=\left(I_{n}-\gamma^{-2} Q P\right)^{-1} Q P . \tag{5.4}
\end{equation*}
$$

Furthermore, $P$ is positive definite if and only if $Y$ is positive definite. In this case,

$$
\begin{equation*}
Y=\left(P^{-1}-\gamma^{-2} Q\right)^{-1} \tag{5.5}
\end{equation*}
$$

Proof. The results are easily proven and thus the details are omitted.
Proposition 5.1. Let $R_{1}=R_{1 \infty}, \beta=1$, suppose there exist $Q, \hat{Q} \in \mathbb{N}^{n}$ and $Z \in \mathbb{P}^{n}$ satisfying

$$
\begin{align*}
& 0=A Q+Q A^{T}+V_{1}+\gamma^{-2} Q R_{1} Q-Q \bar{\Sigma} Q,  \tag{5.6}\\
& 0=\left(A+\gamma^{-2} Q R_{1}\right)^{T} Z+Z\left(A+\gamma^{-2} Q R_{1}\right)+R_{1}-Z \Sigma Z+\gamma^{-2} Z Q \bar{\Sigma} Q Z,  \tag{5.7}\\
& 0=\left(A-\Sigma Z+\gamma^{-2} Q R_{1}\right) \hat{Q}+\hat{Q}\left(A-\Sigma Z+\gamma^{-2} Q R_{1}\right)^{T}+\gamma^{-2} \hat{Q}\left(R_{1}+Z \Sigma Z\right) \hat{Q}+Q \bar{\Sigma} Q,  \tag{5.8}\\
& \qquad \rho(Z \hat{Q})<\gamma^{-2}, \tag{5.9}
\end{align*}
$$

and let $\left(A_{c}, B_{c}, C_{c}, \mathcal{Q}\right)$ be given by

$$
\begin{align*}
A_{c} & =A-Q \Sigma-\Sigma Z+\gamma^{-2} Q R_{2},  \tag{5.10}\\
B_{c} & =Q C^{T} V_{2}^{-1}  \tag{5.11}\\
C_{c} & =-R_{2}^{-1} B^{T} Z  \tag{5.12}\\
\mathcal{Q} & =\left[\begin{array}{cc}
Q+\hat{Q} & \hat{Q} \\
\hat{Q} & \hat{Q}
\end{array}\right] . \tag{5.13}
\end{align*}
$$

Then $\left(\tilde{A},\left[\gamma^{-2} Q \tilde{R} Q+\tilde{V}\right]^{\frac{1}{2}}\right)$ is stabilizable if and only if $\tilde{A}$ is asymptotically stable. In this case, the closed-loop transfer function $H(s)$ satisfies the $H_{\infty}$ disturbance attenuation constraint

$$
\begin{equation*}
\|H(s)\|_{\infty} \leq \gamma \tag{5.14}
\end{equation*}
$$

and the $L_{2}$ performance criterion (2.7) satisfies the bound

$$
\begin{equation*}
J\left(A_{c}, B_{c}, C_{c}\right) \leq \operatorname{tr}\left[(Q+\hat{Q}) R_{1}+\hat{Q} Z \Sigma Z\right] . \tag{5.15}
\end{equation*}
$$

Proof. We need only show that (5.6)-(5.9) imply (3.7)-(3.9). Clearly, (3.7) is a restatement of (5.6). Since $Z$ is positive definite and $\rho(Z \hat{Q})<\boldsymbol{\gamma}^{-3}$, it is possible by Lemma 5.1 (iv) to define $P \triangleq Z\left(I_{n}-\gamma^{-2} \hat{Q} Z\right)^{-1}=\left(Z^{-1}-\gamma^{-2} \hat{Q}\right)^{-1}$. Furthermore, $P$ is positive definite and it follows that $Z=P S$. Thus (5.8) implies (3.9). To obtain (3.8), form $P\left[Z^{-1}(5.7) Z^{-1}-\gamma^{-2}(5.8)\right] P$ and combine terms.

Remark 5.1. It is easy to see that the proof of Proposition 5.1 can be reversed. That is, when $P$ satisfying (3.8) is positive definite define $Z \triangleq\left(P^{-1}+\gamma^{-2} \hat{Q}\right)^{-1}$ which leads to (5.7) and (5.8) using Lemma 5.1 (iii). Thus, when $P$ is positive definite, the form of the equations (5.6)-(5.9) represents no loss of generality.

Proposition 5.2. Let $R_{1}=R_{1 \infty}, \beta=1$, suppose there exist $Q, \hat{Q} \in \mathbb{N}^{n}$ and $P_{0} \in \mathbb{P}^{n}$ satisfying

$$
\begin{align*}
0= & A Q+Q A^{T}+V_{1}+\gamma^{-2} Q R_{1} Q-Q \bar{\Sigma} Q,  \tag{5.16}\\
0= & A^{T} P_{0}+P_{0} A+R_{1}+\gamma^{-2} P_{0} V_{1} P_{0}-P_{0} \Sigma P_{0},  \tag{5.17}\\
0= & {\left[A-\Sigma P_{0}\left(I_{n}-\gamma^{-2} Q P_{0}\right)^{-1}+\gamma^{-2} Q R_{1}\right] \hat{Q}+\hat{Q}\left[A-\Sigma P_{0}\left(I_{n}-\gamma^{-2} Q P_{0}\right)^{-1}+\gamma^{-2} Q R_{1}\right]^{T} } \\
& +\gamma^{-2} \hat{Q}\left[R_{1}+P_{0}\left(I_{n}-\gamma^{-2} Q P_{0}\right)^{-1} \Sigma\left(I_{n}-\gamma^{-2} P_{0} Q\right)^{-1} P_{0}\right] \hat{Q}+Q \bar{\Sigma} Q,  \tag{5.18}\\
& \rho\left[(Q+\hat{Q}) P_{0}\right]<\gamma^{-2}, \tag{5.19}
\end{align*}
$$

and let $\left(A_{c}, B_{c}, C_{c}, \mathcal{Q}\right)$ be given by

$$
\begin{align*}
A_{c} & =A-Q \bar{\Sigma}-\Sigma P_{0}\left(I_{n}-\gamma^{-2} Q P_{0}\right)^{-1}+\gamma^{-2} Q R_{1}  \tag{5.20}\\
B_{c} & =Q C^{T} V_{2}^{-1}  \tag{5.21}\\
C_{c} & =-R_{2}^{-1} B^{T} P_{0}\left(I_{n}-\gamma^{-2} Q P_{0}\right)^{-1}  \tag{5.22}\\
Q & =\left[\begin{array}{cc}
Q+\hat{Q} & \hat{Q} \\
\hat{Q} & \hat{Q}
\end{array}\right] . \tag{5.23}
\end{align*}
$$

Then $\left(\tilde{A},\left[\gamma^{-2} Q \tilde{R} Q+\tilde{V}\right]^{\frac{1}{2}}\right)$ is stabilizable if and only if $\tilde{A}$ is asymptotically stable. In this case, the closed-loop transfer function $H(s)$ satisfies the $H_{\infty}$ disturbance attenuation constraint

$$
\begin{equation*}
\|H(s)\|_{\infty} \leq \gamma, \tag{5.24}
\end{equation*}
$$

and the $L_{2}$ performance criterion (2.7) satisfies the bound

$$
\begin{equation*}
J\left(A_{c}, B_{c}, C_{c}\right) \leq \operatorname{tr}\left[(Q+\hat{Q}) R_{1}+\hat{Q} P_{0}\left(I_{n}-\gamma^{-2} Q P_{0}\right)^{-1} \Sigma\left(I_{n}-\gamma^{-2} P_{0} Q\right)^{-1} P_{0}\right] \tag{5.25}
\end{equation*}
$$

Proof. As in the proof of Proposition 5.1, it need only be shown that (5.16)-(5.19) imply (3.7)-(3.9). Hence define $P=\left[P_{0}^{-1}-\gamma^{-2}(Q+\hat{Q})\right]^{-1}$ and form the equation $P\left[P_{0}^{-1}(5.17) P_{0}^{-1}-\right.$ $\boldsymbol{\gamma}^{-2}(5.16)-\gamma^{-2}(5.18) \mid P$ to obtain (3.8). Equations (3.7) and (3.9) are immediate.

Remark 5.2. To clarify the relationships among (3.7)-(3.9), (5.6)-(5.8), and (5.16)-(5.18), we tabulate the transformations involving $P, Z$ and $P_{0}$ :

$$
\begin{align*}
& P=\left(Z^{-1}-\gamma^{-2} \hat{Q}\right)^{-1}=\left[P_{0}^{-1}-\gamma^{-2}(Q+\hat{Q})\right]^{-1}  \tag{5.26}\\
& Z=\left(P^{-1}+\gamma^{-2} \hat{Q}\right)^{-1}=\left(P_{0}^{-1}-\gamma^{-2} Q\right)^{-1}  \tag{5.27}\\
& P_{0}=\left(Z^{-1}+\gamma^{-2} Q\right)^{-1}=\left[P^{-1}+\gamma^{-2}(Q+\hat{Q})\right]^{-1} \tag{5.28}
\end{align*}
$$

- Remark 5.3. It is important to note that numerically solving equations (5.6)-(5.8) and (5.16)(5.18) does not require that $Z$ and $P_{0}$ be invertible. As shown in the Proofs of Propositions 5.1 and 5.2 the positive definiteness assumptions are used to construct equation (3.8).

Remark 5.4. Note that the gains (5.10)-(5.12) and (5.16)-(5.18) of Propositions 5.1 and 5.2 are independent of the matrix $\hat{Q}$ satisfying equations (5.8) and (5.18). Nevertheless, equations (5.8) and (5.18) must have a solution in order to enforce the solvability of (2.14) which implies that the $H_{\infty}$ constraint is satisfied. Thus our result does not yield a guarantee of $H_{\infty}$ performance unless (5.8) and (5.18) can be solved numerically. Of course, $\hat{Q}$ is also required to evaluate the $L_{2}$ performance bound (5.15) or (5.25). Note that the solutions of $Q$ and $P_{0}$ of (5.16) and (5.17) are analogous to the matrices $X_{\infty}$ and $Y_{\infty}$ of [26]. Finally, note that (5.19) implies that $\rho\left(Q P_{0}\right)<\gamma^{-2}$ which is essentially condition 5.2 (iii) of [26].

Remark 5.5. The transformations (5.26)-(5.28) used to obtain the form of the equations given by Propositions 5.1 and 5.2 depend strongly upon the assumptions $R_{1 \infty}=R_{1}$ and $\beta=1$. That is, if either $R_{1 \infty} \neq R_{1}$ or $\beta \neq 1$ then these transformations cannot be carried out. Thus, although (5.6)-(5.8) and (5.16) -(5.18) possess numerical advantages over (3.7)-(3.9), these alternative forms exist only in the very special case in which the $L_{2}$ and $H_{\infty}$ weights are equalized. Moreover, in the presence of parameter uncertainties ([13,15]) or nonstrictly proper controller design ([32]), such transformations seem to be precluded.

Remark 5.6. It is interesting to note that equations (5.6) and (5.7) with controller gains (5.10)-(5.12) are already known since they are identical to the optimality conditions of the exponential-quadratic-Gaussian problem treated in [33]. Specifically, see equations (3.1) and (4.1) on pages 603 and 609, respectively. As shown in [33], minimizing the criterion

$$
J=\lim _{t \rightarrow \infty} \mathbb{E} \mu e^{\frac{u}{2}\left(x^{T} R_{1} x+u^{T} R_{2} u\right)}
$$

leads to the pair of modified Riccati equations (5.6) and (5.7) with $\boldsymbol{\gamma}^{-2}$ replaced by $\mu$. This implies that with equation (5.8) the exponential-of-quadratic design problem effectively enforces a bound of $\mu^{-\frac{1}{2}}$ on the $H_{\infty}$ norm of the closed-loop transfer function. What remains to be achieved then is a deeper understanding of this connection. For related references, see $[34,35]$.

## 6. Extensions to Reduced-Order Dynamic Compensation

In this section we extend Theorem 4.1 by expanding the formulation of Section 3 to allow the compensator to be of fixed dimension $n_{c}$ which may be less than the plant order $n$. Hence, in this section define $\tilde{n}=n+n_{c}$, where $n_{c} \leq n$. As in [20] this constraint leads to an oblique projection which introduces additional coupling in the design equations along with an additional equation. The following lemma is required.

Lemma 6.1. Let $\hat{Q}, \hat{P} \in \mathbb{N}^{n}$ and suppose rank $\hat{Q} \hat{P}=n_{c}$. Then there exist $n_{c} \times n G, \Gamma$ and $n_{c} \times n_{c}$ invertible $M$, unique except for a change of basis in $\mathbb{R}^{n_{c}}$, such that

$$
\begin{equation*}
\hat{Q} \hat{P}=G^{T} M \Gamma, \quad \Gamma G^{T}=I_{n_{0}} \tag{6.1}
\end{equation*}
$$

Furthermore, the $n \times n$ matrices

$$
\begin{equation*}
\tau \triangleq G^{\boldsymbol{T}} \Gamma, \quad \tau_{\perp} \triangleq I_{n}-\tau \tag{6.3}
\end{equation*}
$$

are idempotent and have rank $n_{c}$ and $n-n_{c}$, respectively.
Proof. Conditions (6.1)-(6.4) are a direct consequence of Theorem 6.2.5 of [36].
Theorem 6.1. Let $n_{c} \leq n$, suppose there exist $Q, P, \hat{Q}, \hat{P} \in \mathbb{N}^{n}$ satisfying

$$
\begin{align*}
0= & A Q+Q A^{T}+V_{1}+\gamma^{-2} Q R_{1 \infty} Q-Q \bar{\Sigma} Q+\tau_{\perp} Q \tilde{\Sigma} Q \tau_{\perp}^{T},  \tag{6.5}\\
0= & \left(A+\gamma^{-2}[Q+\hat{Q}] R_{1 \infty}\right)^{T} P+P\left(A+\gamma^{-2}[Q+\hat{Q}] R_{1 \infty}\right)+R_{1} \\
& \quad-S^{T} P \Sigma P S+\tau_{\perp}^{T} S^{T} P \Sigma P S \tau_{\perp},  \tag{6.6}\\
0= & \left(A-\Sigma P S+\gamma^{-2} Q R_{1 \infty}\right) \hat{Q}+\hat{Q}\left(A-\Sigma P S+\gamma^{-2} Q R_{1 \infty}\right)^{T} \\
& +\gamma^{-2} \hat{Q}\left(R_{1 \infty}+\beta^{2} S^{T} P \Sigma P S\right) \hat{Q}+Q \bar{\Sigma} Q-\tau_{\perp} Q \bar{\Sigma} Q \tau_{\perp}^{T},  \tag{6.7}\\
0= & \left(A-Q \bar{\Sigma}+\gamma^{-2} Q R_{1 \infty}\right)^{T} \hat{P}+\hat{P}\left(A-Q \bar{\Sigma}+\gamma^{-2} Q R_{1 \infty}\right) \\
& +S^{T} P \Sigma P S-\tau_{\perp}^{T} S^{T} P \Sigma P S \tau_{\perp},  \tag{6.8}\\
& \operatorname{rank} \hat{Q}=\operatorname{rank} \hat{P}=\operatorname{rank} \hat{Q} \hat{P}=n_{c}, \tag{6.9}
\end{align*}
$$

and let $\left(A_{c}, B_{c}, C_{c}, Q\right)$ be given by

$$
\begin{align*}
A_{c} & =\Gamma\left(A-Q \bar{\Sigma}-\Sigma P S+\gamma^{-2} Q R_{1 \infty}\right) G^{T},  \tag{6.10}\\
B_{c} & =\Gamma Q C^{T} V_{2}^{-1},  \tag{6.11}\\
C_{c} & =-R_{2}^{-1} B^{T} P S G^{T},  \tag{6.12}\\
Q & =\left[\begin{array}{cc}
Q+\hat{Q} & \hat{Q} \Gamma^{T} \\
\Gamma \hat{Q} & \Gamma \hat{Q} \Gamma^{T}
\end{array}\right] . \tag{6.13}
\end{align*}
$$

Then, $\left(\tilde{A},\left[\gamma^{-2} Q \tilde{R}_{\infty} Q+\tilde{V}\right]^{\frac{1}{2}}\right)$ is stabilizable if and only if $\tilde{A}$ is asymptotically stable. In this case, the closed-loop transfer function $H(s)$ satisfies the $H_{\infty}$ disturbance attenuation constraint

$$
\begin{equation*}
\|H(s)\|_{\infty} \leq \gamma \tag{6.14}
\end{equation*}
$$

and the $L_{2}$ performance criterion (2.7) satisfies the bound

$$
\begin{equation*}
J\left(A_{c}, B_{c}, C_{c}\right) \leq \operatorname{tr}\left[(Q+\hat{Q}) R_{1}+\hat{Q} S^{T} P \Sigma P S\right] \tag{6.15}
\end{equation*}
$$

Remark 6.1. It is easy to see that Theorem 6.1 is a direct generalization of Theorem 4.1. To recover Theorem 4.1, set $n_{c}=n$ so that $\tau=G=\Gamma=I_{n}$ and $\tau_{\perp}=0$. In this case the last term in each of (6.5)-(6.8) can be deleted and equation (6.8) becomes superfluous. Furthermore, (6.5)-(6.7) now reduce to (3.7)-(3.9), as expected. If, furthermore, $\beta=0$ then $S=I_{n}$ so that equations (6.5)-(6.7) now reduce to the "cheap" $H_{\infty}$ control case given by (3.7), (3.14) and (3.15). Alternatively, setting $\boldsymbol{\gamma}=\infty$ and retaining the reduced-order constraint $n_{c}<n$ yields the result of [20].

Remark 6.2. Consider the case $R_{1 \infty}=R_{1}$ and $\beta=1$. By introducing a new variable $Z=P S=\left(P^{-1}+\gamma^{-2} \hat{Q}\right)^{-1}$ as in Section 5, equation (6.6) becomes

$$
\begin{align*}
0= & \left(A+\gamma^{-2} Q R_{1}\right)^{T} Z+Z\left(A+\gamma^{-2} Q R_{1}\right)+R_{1}  \tag{6.16}\\
& -Z \Sigma Z+\tau_{\perp}^{T} Z \Sigma Z \tau_{\perp}+\gamma^{-2} Z\left[Q \Sigma Q-\tau_{\perp} Q \bar{\Sigma} Q \tau_{\perp}^{T}\right] Z .
\end{align*}
$$

Note that (6.16) specializes to (5.7) when $\tau_{\perp}=0$ (i.e., $n_{c}=n$ ). Furthermore, $P S$ can be replaced by $Z$ in (6.7)-(6.12). Next, to generalize (5.17), define $P_{0}$ as in (5.28) so that (5.17) becomes

$$
\begin{align*}
0= & A^{T} P_{0}+P_{0} A+R_{1}+\gamma^{-2} P_{0} V_{1} P_{0}-P_{0} \Sigma P_{0} \\
& +\left(I_{n}-\gamma^{-2} P_{0} Q\right) \tau_{\perp}^{T}\left(I_{n}-\gamma^{-2} P_{0} Q\right)^{-1} P_{0} \Sigma P_{0}\left(I_{n}-\gamma^{-2} Q P_{0}\right)^{-1} r_{\perp}\left(I_{n}-\gamma^{-2} Q P_{0}\right) . \tag{6.17}
\end{align*}
$$

Again, (6.17) specializes to (5.17) when $\tau_{\perp}=0$.

## 7. Analysis of the Design Equations

Before developing numerical algorithms, it is instructive to analyze the design equations to determine existence and multiplicity of nonnegative-definite solutions. Although a detailed theoretical analysis remains an area for future research, in this section we present a simplified treatment which highlights important asymptotic properties of the equations. It turns out that several key properties are discernible by considering the scalar case $n=1$. Although many of these properties
can be developed for general $n$, the simplified scalar case suffices for obtaining a useful qualitative analysis. Here we consider only (3.7), (3.14) and (3.15).

Since the $Q$ equation (3.7) is decoupled from (3.14) and (3.15), it can be analyzed separately. It is easy to see that there exists a unique nonnegative solution for $\gamma>\left(R_{1} / \bar{\Sigma}\right)^{\frac{1}{2}}$ as in the case of a standard Riccati equation with stabilizability and detectability hypothesis. Furthermore, it can be seen that for

$$
\left(R_{1} /\left[\bar{\Sigma}+\left(A^{2} / V_{1}\right)\right]\right)^{\frac{1}{2}}<\gamma<\left(R_{1} / \bar{\Sigma}\right)^{\frac{1}{2}}
$$

there exist two nonnegative solutions when $A$ is stable and zero nonnegative solutions when $A$ is unstable. Below this lower bound for $\boldsymbol{\gamma}$ nonnegative solutions $Q$ do not exist. This result thus indicates (as in LQG theory [42]) a lower bound to the achievable $H_{\infty}$ disturbance attenuation as determined by the sensor noise intensity $V_{2}$ appearing in $\bar{\Sigma}$.

Since the $P$ and $\hat{Q}$ equations (3.14) and (3.15) are coupled they must be analyzed jointly. Since (3.15) is a standard Riccati equation, it follows under generic hypotheses that it possesses exactly one nonnegative-definite solution for all values of $Q$ and $\hat{Q}$. The analysis of the $\hat{Q}$ equation is, however, more involved. It can be shown that the existence of real solutions is a complicated function of $\gamma, Q$, and $P$. When real solutions do exist, it follows that there exist either zero or two nonnegative-definite solutions. To obtain further qualitative insight into the solutions $P$ and $\hat{Q}$ we fix $\gamma$ and allow $R_{2} \rightarrow 0$, that is, the cheap $L_{2}$ control case. It thus follows that $P \sim\left(R_{1} / \Sigma\right)^{\frac{1}{2}}$ and that either $\hat{Q} \sim 2 \gamma^{2}\left(\Sigma / R_{1}\right)^{\frac{1}{2}}$ or $\hat{Q} \sim \frac{1}{2} \bar{\Sigma} Q^{2}\left(\Sigma R_{1}\right)^{-\frac{1}{2}}$, which correspond to the previously mentioned pair of solutions satisfying (3.15). This result thus indicates that an arbitrarily small $H_{\infty}$ disturbance attenuation constraint $\gamma$ can be achieved (subject to the solvability of (3.7)) by sufficiently increasing the $L_{2}$ controller authority. That is, since solutions exist in the cheap $L_{2}$ control case, the $H_{\infty}$ disturbance attenuation constraint is achievable. The ability to achieve small $\boldsymbol{\gamma}$ is also attributable to the fact that since $\beta=0, H_{\infty}$ disturbance attenuation to the control variables is not limited in (3.7), (3.14) and (3.15) as in Theorems 3.1 and 6.1. Of course, as is well known, it is not possible to make $\boldsymbol{\gamma} \rightarrow 0$ by letting $\Sigma \rightarrow \infty$ and $\bar{\Sigma} \rightarrow \infty$ when the system possesses nonminimum phase zeros. Also, note that both of the asymptotic solutions to (3.15) are guaranteed to yield the bound (4.1). The solution of interest, however, is $\hat{Q}=O\left(\Sigma^{-\frac{1}{2}}\right)$ since it clearly yields a lower value of $J\left(A_{c}, B_{c}, C_{c}, Q\right)$ than $\hat{Q}=O\left(\Sigma^{\frac{1}{2}}\right)$. Finally, similar analysis can be applied to (5.6)-(5.8) and (5.16)-(5.18).

## 8. Numerical Algorithm and Illustrative Results

In this section we describe a numerical algorithm which has been developed and implemented for solving the coupled Riccati equations (3.7), (3.14) and (3.15). We also present numerical results for an illustrative example.

Coupled modified Riccati equations arise in a variety of problems and homotopic continuation methods have been shown to be particularly successful [23-25]. To solve (3.7), (3.14) and (3.15) we have implemented a simplified continuation method involving the constraint constant $\gamma$. The idea is to exploit the fact that for large $\gamma$ the problem is approximated by LQG which provides a reliable starting solution. The continuation parameter $\gamma$ is then successively decreased until either a desired value of $\boldsymbol{\gamma}$ is achieved or no further decrease is possible. This algorithm is now summarized. Let $\varepsilon>0$ denote a convergence criterion.

Algorithm 8.1. To solve (3.7), (3.14) and (3.15), perform the following steps:

Step 1: Initialize $\boldsymbol{\gamma}>\mathbf{0}$;
Step 2: Solve (3.7) for $Q$;
Step 3: Let $k=0, \hat{Q}_{0}=0$;
Step 4: Solve (3.14) for $P_{k+1}=P$ with $\hat{Q}=\hat{Q}_{k}$;
Step 5: Solve (3.15) for $\hat{Q}_{k+1}=\hat{Q}$ with $P=P_{k+1}$;
Step 6: If $k \geq 1$ check for $\left\|P_{k+1}-P_{k}\right\|<\varepsilon$ and $\left\|\hat{Q}_{k+1}-\hat{Q}_{k}\right\|<\varepsilon$;
Step 7: If convergence is not achieved in Step 6 (or $k=0$ ) let $k \leftarrow k+1$ and go to Step 4; otherwise decrease $\boldsymbol{\gamma}$ and go to Step 2.

Steps 2, 4 and 5 were carried out using a standard Riccati solver [37] which proved to be reliable even when the quadratic term was indefinite or nonnegative definite. For instance, for the example considered below the term $\boldsymbol{\gamma}^{-\mathbf{3}} \boldsymbol{R}_{1}-\boldsymbol{\Sigma}$ was indefinite for all finite $\boldsymbol{\gamma}$. The crucial step in the algorithm is the decreasing of $\boldsymbol{\gamma}$ in Step 7. Significant effort was devoted to providing a smooth transition to smaller values of $\boldsymbol{\gamma}$ without sacrificing computational efficiency. The development of more sophisticated continuation algorithms remains an area for future research.

The example considered was formulated in [38] and was considered extensively in [24,25,39] to
compare reduced-order design methods. The example is interesting since it possesses a complex pair of nonminimum phase zeros due to the fact that the physical system (coupled rotating disks) has noncolocated sensors and actuators. The plant is of eighth order and has two neutrally stable poles. The problem data are as follows:

$$
n=n_{c}=8, \quad m=\ell=1, \quad q=p=2,
$$

$$
\begin{gathered}
A=\left[\begin{array}{cccccccc}
-0.161 & 1 & 0 & 0 & 0 & 0 & 0 & 0 \\
-6.004 & 0 & 1 & 0 & 0 & 0 & 0 & 0 \\
-0.5822 & 0 & 0 & 1 & 0 & 0 & 0 & 0 \\
-9.9835 & 0 & 0 & 0 & 1 & 0 & 0 & 0 \\
-0.4073 & 0 & 0 & 0 & 0 & 1 & 0 & 0 \\
-3.982 & 0 & 0 & 0 & 0 & 0 & 1 & 0 \\
0 & 0 & 0 & 0 & 0 & 0 & 0 & 1 \\
0 & 0 & 0 & 0 & 0 & 0 & 0 & 0
\end{array}\right] \quad B=\left[\begin{array}{c}
0 \\
0 \\
.0064 \\
.00235 \\
.0713 \\
1.0002 \\
.1045 \\
.9955
\end{array}\right] \quad C=\left[\begin{array}{ll}
1 & 0_{1 \times 7}
\end{array}\right] \\
E_{1}=E_{1 \infty}=10^{-3}\left[\begin{array}{llllllll}
0 & 0 & 0 & 0 & .55 & 11 & 1.32 & 18 \\
0 & 0 & 0 & 0 & 0 & 0 & 0 & 0
\end{array}\right], \\
E_{2}=\left[\begin{array}{l}
0 \\
1
\end{array}\right], \quad E_{2 \infty}=\left[\begin{array}{l}
0 \\
0
\end{array}\right], \quad \beta=0, \\
D_{1}=\left[\begin{array}{lll}
B & 0_{8 \times 1}
\end{array}\right], \quad D_{2}=\left[\begin{array}{ll}
0 & 1
\end{array}\right] .
\end{gathered}
$$

With the problem data as given, the LQG controller was found to yield a closed-loop $H_{\infty}$ performance of 1.39 (i.e., 2.87 dB above unity gain). Using Algorithm 8.1 we obtained a solution for $\gamma=.52$ for a net $H_{\infty}$ performance improvement of 8.7 dB (see Figure 1). Note that this result is consistent with Proposition 8.1 of [3] which implies that the maximum ratio of the $H_{\infty}$ performance of the optimal $L_{2}$ controller to the $H_{\infty}$ performance of the optimal $H_{\infty}$ controller can be no more than twice the number of right-half-plane zeros, which for the present problem with two nonminimum phase zeros corresponds to a factor of 4 (i.e., 12 dB ).

Our numerical experience revealed two interesting features. First, the loop between Steps 4 and 6 converged reliably. However, a critical value $\boldsymbol{\gamma}_{\text {min }}$ of $\boldsymbol{\gamma}$ was invariably found below which solutions could not be computed. This value $\gamma_{\text {min }}$ appears to represent the best achievable $H_{\infty}$ performance for the given $L_{2}$ weights. Second, for each value of $\boldsymbol{\gamma} \geq \boldsymbol{\gamma}_{\min }$ for which a solution
was computed, the actual $H_{\infty}$ performance was close to this value revealing that the $H_{\infty}$ bound is tight. For example, the actual worst-case attenuation of the $\gamma=.52$ design shown in Figure 1 is .511 . Controller characteristics are given in Table 1. Note that in each case the $L_{2}$ performance bound is within $30 \%$ of the actual $L_{2}$ performance.

| $H_{\infty}$ Attenuation | Actual $H_{\infty}$ | $L_{2}$ Performance | Actual $L_{2}$ |
| :---: | :---: | :---: | :---: |
| Constraint $\boldsymbol{\gamma}$ | Attenuation <br>  <br>  <br>  <br> $\boldsymbol{H}(s) \\|_{\infty}$ | Bound <br> $J\left(A_{c}, B_{c}, C_{c}, Q\right)$ | Performance <br> $J\left(A_{c}, B_{c}, C_{c}\right)$ |
| 2 | 1.39 | - | .143 |
| 1.5 | 1.18 | .159 | .146 |
| 1.0 | 1.06 | .171 | .151 |
| .9 | .855 | .204 | .168 |
| .8 | .797 | .217 | .176 |
| .7 | .732 | .236 | .187 |
| .52 | .661 | .262 | .203 |

Table 1

Acknowledgments. We wish to thank Professor P. P. Khargonekar for several helpful discussions and for providing a preprint of [19], : Straehla for transforming the original manuscript into $T_{E} X$, A. Daubendiek, S. Greeley, S. Richter, and A. Tellez for developing the numerical algorithm and performing the calculations of Section 8, D. Hyland, E. Collins, and L. Davis for helpful discussions and suggestions, A. N. Madiwale for providing simplifications of (6.6)-(6.8), and the reviewers for several helpful comments.

## Appendix A: Proof of Theorem 6.1

To optimize (2.20) over the open set $X$ subject to the constraint (2.14), form the Lagrangian

$$
\begin{equation*}
\mathcal{L}\left(A_{c}, B_{c}, C_{c}, Q, P, \lambda\right) \triangleq \operatorname{tr}\left\{\lambda Q \tilde{R}+\left[\tilde{A} Q+Q \tilde{A}^{T}+\gamma^{-2} Q \tilde{R}_{\infty} Q+\tilde{V}\right] P\right\} \tag{A.1}
\end{equation*}
$$

where the Lagrange multipliers $\lambda \geq 0$ and $P \in \mathbb{R}^{\bar{n} \times \bar{n}}$ are not both zero. We thus obtain

$$
\begin{equation*}
\frac{\partial \mathcal{L}}{\partial Q}=\left(\tilde{A}+\gamma^{-2} \varrho \tilde{R}_{\infty}\right)^{T} P+\mathcal{P}\left(\tilde{A}+\gamma^{-2} Q \tilde{R}_{\infty}\right)+\lambda \tilde{R} \tag{A.2}
\end{equation*}
$$



Figure 1

Setting $\partial \mathcal{L} / \partial Q=0$ yields

$$
\begin{equation*}
0=\left(\tilde{A}+\gamma^{-2} Q \tilde{R}_{\infty}\right)^{T} \rho+P\left(\tilde{A}+\gamma^{-2} Q \tilde{R}_{\infty}\right)+\lambda \tilde{R} \tag{A.3}
\end{equation*}
$$

Since $\tilde{A}+\gamma^{-2} Q \tilde{R}_{\infty}$ is assumed to be stable, $\lambda=0$ implies $P=0$. Hence, it can be assumed without loss of generality that $\lambda=1$. Furthermore, $P$ is nonnegative definite.

Now partition $\tilde{n} \times \tilde{n} \rho, P$ into $n \times n, n \times n_{c}$, and $n_{c} \times n_{c}$ subblocks as

$$
\mathcal{Q}=\left[\begin{array}{cc}
Q_{1} & Q_{12} \\
Q_{12}^{T} & Q_{2}
\end{array}\right], \quad \mathcal{P}=\left[\begin{array}{cc}
P_{1} & P_{12} \\
P_{12}^{T} & P_{2}
\end{array}\right]
$$

Thus, with $\lambda=1$ the stationarity conditions are given by

$$
\begin{align*}
& \frac{\partial \mathcal{L}}{\partial Q}=\left(\tilde{A}+\gamma^{-2} Q \tilde{R}_{\infty}\right)^{T} P+P\left(\tilde{A}+\gamma^{-2} Q \tilde{R}_{\infty}\right)+\tilde{R}=0,  \tag{A.4}\\
& \frac{\partial \mathcal{L}}{\partial A_{c}}=P_{12}^{T} Q_{12}+P_{2} Q_{2}=0,  \tag{A.5}\\
& \frac{\partial \mathcal{L}}{\partial B_{c}}=P_{2} B_{c} V_{2}+\left(P_{12}^{T} Q_{1}+P_{2} Q_{12}^{T}\right) C^{T}=0,  \tag{A.6}\\
& \frac{\partial \mathcal{L}}{\partial C_{c}}=R_{2} C_{c} Q_{2}+\beta^{2} \gamma^{-2} R_{2} C_{c}\left(P_{1} Q_{12}+P_{12} Q_{2}\right)^{T} Q_{12}+B^{T}\left(P_{1} Q_{12}+P_{12} Q_{2}\right)=0 \tag{A.7}
\end{align*}
$$

Expanding (2.14) and (A.4) yields

$$
\begin{align*}
0= & A Q_{1}+Q_{1} A^{T}+B C_{c} Q_{12}^{T}+Q_{12} C_{c}^{T} B^{T}+\gamma^{-2} Q_{1} R_{1 \infty} Q_{1} \\
& +\beta^{2} \gamma^{-2} Q_{12} C_{c}^{T} R_{2} C_{c} Q_{12}^{T}+V_{1},  \tag{A.8}\\
0= & A Q_{12}+Q_{12} A_{c}^{T}+B C_{c} Q_{2}+Q_{1} C^{T} B_{c}^{T}+\gamma^{-2} Q_{1} R_{1 \infty} Q_{12} \\
& +\beta^{2} \gamma^{-2} Q_{12} C_{c}^{T} R_{2} C_{c} Q_{2},  \tag{A.9}\\
0= & A_{c} Q_{2}+Q_{2} A_{c}^{T}+B_{c} C Q_{12}+Q_{12}^{T} C^{T} B_{c}^{T}+\gamma^{-2} Q_{12}^{T} R_{1 \infty} Q_{12} . \\
& +\beta^{2} \gamma^{-2} Q_{2} C_{c}^{T} R_{2} C_{c} Q_{2}+B_{c} V_{2} B_{c}^{T},  \tag{A.10}\\
0= & A^{T} P_{1}+P_{1} A+C^{T} B_{c}^{T} P_{12}^{T}+P_{12} B_{c} C+\gamma^{-2} R_{1 \infty}\left(P_{1} Q_{1}+P_{12} Q_{12}^{T}\right)^{T} \\
& +\gamma^{-2}\left(P_{1} Q_{1}+P_{12} Q_{12}^{T}\right) R_{1 \infty}+R_{1},  \tag{A.11}\\
0= & A^{T} P_{12}+P_{12} A_{c}+C^{T} B_{c}^{T} P_{2}+P_{1} B C_{c}+\gamma^{-2} R_{1 \infty}\left(P_{12}^{T} Q_{1}+P_{2} Q_{12}^{T}\right)^{T} \\
& +\beta^{2} \gamma^{-2}\left(P_{1} Q_{12}+P_{12} Q_{2}\right) C_{c}^{T} R_{2} C_{c},  \tag{A.12}\\
0= & A_{c}^{T} P_{2}+P_{2} A_{c}+P_{12}^{T} B C_{c}+C_{c}^{T} B^{T} P_{12}+C_{c}^{T} R_{2} C_{c} . \tag{A.13}
\end{align*}
$$

Lemma A.1. $Q_{2}$ and $P_{2}$ are positive definite.
Proof. By a minor extension of results from [40], (A.10) can be rewritten as

$$
0=\left(A_{c}+B_{c} C Q_{12} Q_{2}^{+}\right) Q_{2}+Q_{2}\left(A_{c}+B_{c} C Q_{12} Q_{2}^{+}\right)^{T}+\Psi,
$$

where

$$
\Psi \triangleq \gamma^{-2} Q_{12}^{T} R_{1 \infty} Q_{12}+\beta^{2} \gamma^{-2} Q_{2} C_{c}^{T} R_{2} C_{c} Q_{2}+B_{c} V_{2} B_{c}^{T}
$$

and $Q_{2}^{+}$is the Moore-Penrose or Drazin generalized inverse of $Q_{2}$. Next note that since ( $A_{c}, B_{c}$ ) is controllable it follows from Lemma 2.1 and Theorem 3.6 of [28] that $\left(A_{c}+B_{c} C Q_{12} Q_{2}^{+}, \Psi^{\frac{1}{2}}\right)$ is also controllable. Now, since $Q_{2}$ and $\Psi$ are nonnegative definite, Lemma 12.2 of [28] implies that $Q_{2}$ is positive definite. Using (A.13), similar arguments show that $P_{2}$ is positive definite.

Since $R_{2}, V_{2}, Q_{2}, P_{2}$ are invertible, (A.5)-(A.7) can be written as

$$
\begin{gather*}
-P_{2}^{-1} P_{12}^{T} Q_{12} Q_{2}^{-1}=I_{n_{0}},  \tag{A.14}\\
B_{c}=-P_{2}^{-1}\left(P_{12}^{T} Q_{1}+P_{2} Q_{12}^{T}\right) C^{T} V_{2}^{-1},  \tag{A.15}\\
C_{c}\left[I_{n_{0}}+\beta^{2} \gamma^{-2}\left(Q_{12}^{T} P_{1}+Q_{2} P_{12}^{T}\right) Q_{12} Q_{2}^{-1}\right]=-R_{2}^{-1} B^{T}\left(P_{1} Q_{12}+P_{12} Q_{2}\right) Q_{2}^{-1} . \tag{A.16}
\end{gather*}
$$

Now define the $n \times n$ matrices

$$
\begin{gathered}
Q \triangleq Q_{1}-Q_{12} Q_{2}^{-1} Q_{12}^{T}, \quad P \triangleq P_{1}-P_{12} P_{2}^{-1} P_{12}^{T} \\
\hat{Q} \triangleq Q_{12} Q_{2}^{-1} Q_{12}^{T}, \quad \hat{P} \triangleq P_{12} P_{2}^{-1} F_{12}^{T} \\
\tau \triangleq-Q_{12} Q_{2}^{-1} P_{2}^{-1} P_{12}^{T}
\end{gathered}
$$

and the $n_{c} \times n, n_{c} \times n_{c}$ and $n_{c} \times n$ matrices

$$
G \triangleq Q_{2}^{-1} Q_{12}^{T}, \quad M \triangleq Q_{2} P_{2}, \quad \Gamma \triangleq-P_{2}^{-1} P_{12}^{T}
$$

Note that $\tau=G^{\boldsymbol{T}} \Gamma$.
Clearly, $Q, P, \hat{Q}$ and $\hat{P}$ are symmetric and $\hat{Q}$ and $\hat{P}$ are nonnegative definite. To show that $Q$ and $P$ are also nonnegative definite, note that $Q$ is the upper left-hand block of the nonnegativedefinite matrix $\tilde{\mathscr{Q}} \mathcal{Q} \tilde{\mathscr{Q}}^{T}$, where

$$
\tilde{Q} \triangleq\left[\begin{array}{cc}
I_{n} & -Q_{12} Q_{2}^{-1} \\
O_{n_{0} \times n} & I_{n_{0}}
\end{array}\right]
$$

Similarly, $P$ is nonnegative definite.

Next note that with the above definitions (A.14) is equivalent to (6.2) and that (6.1) holds. Hence $\tau=G^{T} \Gamma$ is idempotent, i.e., $\boldsymbol{r}^{\mathbf{2}}=\boldsymbol{\tau}$.

It is helpful to note the identities

$$
\begin{gather*}
\hat{Q}=Q_{12} G=G^{T} Q_{12}^{T}=G^{T} Q_{2} G, \quad \hat{P}=-P_{12} \Gamma=-\Gamma^{T} P_{12}^{T}=\Gamma^{T} P_{2} \Gamma  \tag{A.17}\\
\tau G^{T}=G^{T}, \quad \Gamma \tau=\Gamma  \tag{A.18}\\
\hat{Q}=\tau \hat{Q}, \quad \hat{P}=\hat{P} \tau  \tag{A.19}\\
\hat{Q} \hat{P}=-Q_{12} P_{12}^{T} \tag{A.20}
\end{gather*}
$$

Using (A.14) and Sylvester's inequality, it follows that

$$
\operatorname{rank} G=\operatorname{rank} \Gamma=\operatorname{rank} Q_{12}=\operatorname{rank} P_{12}=n_{c}
$$

Now using (A.17) and Sylvester's inequality yields

$$
n_{c}=\operatorname{rank} Q_{12}+\operatorname{rank} G-n_{c} \leq \operatorname{rank} \hat{Q} \leq \operatorname{rank} Q_{12}=n_{c},
$$

which implies that rank $\hat{Q}=n_{c}$. Similarly, rank $\hat{P}=n_{c}$, and rank $\hat{Q} \hat{P}=n_{c}$ follows from (A.20).
The components of $Q$ and $P$ can be written in terms of $Q, P, \hat{Q}, \hat{P}, G$ and $\Gamma$ as

$$
\begin{array}{ll}
Q_{1}=Q+\hat{Q}, & P_{1}=P \div \hat{P} \\
Q_{12}=\hat{Q} \Gamma^{T}, & P_{12}=-\hat{P} G^{T} \\
Q_{2}=\Gamma \hat{Q} \Gamma^{T}, & P_{2}=G \hat{P} G^{T} \tag{A.23}
\end{array}
$$

Next note that by using (A.21)-(A.23) it can be shown that the right-hand coefficient of $C_{c}$ in (A.16) is given by

$$
\hat{S} \triangleq I_{n_{0}}+\beta^{2} \gamma^{-2} \Gamma \hat{Q} P G^{T}
$$

To prove that $\hat{S}$ is invertible use (A.19) and (6.3) and note that

$$
\begin{aligned}
I_{n_{\&}}+\beta^{2} \gamma^{-2} \Gamma \hat{Q} P G^{T} & =I_{n_{+}}+\beta^{2} \gamma^{-2} \Gamma \hat{Q} \tau^{T} P G^{T} \\
& =I_{n_{\&}}+\beta^{2} \gamma^{-2}\left(\Gamma \hat{Q} \Gamma^{T}\right)\left(G P G^{T}\right)
\end{aligned}
$$

Since $\Gamma \hat{Q} \Gamma^{T}$ and $G P G^{T}$ are nonnegative definite, their product has nonnegative eigenvalues (see Lemma 5.1). Thus each eigenvalue of $I_{n_{e}}+\beta^{2} \gamma^{-2} \Gamma \hat{Q} P G^{T}$ is real and is greater than unity. Hence $\hat{S}$ is invertible. Now note that by using (6.2) and (6.3) it can be shown that

$$
G^{T} \hat{S}^{-1} \Gamma=S \tau
$$

The expressions (6.11), (6.12) and (6.13) follow from (A.15), (A.16), and the definition of $\mathcal{Q}$. Next, computing either $\Gamma$ (A.9)-(A.10) or $G^{\prime}(A .12)+(A .12)$ yields (6.10). Substituting (A.21)-(A.23) into (A.S)-(A.13) and the expression for $A_{c}$ into (A.9), (A.10), (A.12) and (A.13) it follows that (A.10) $=\Gamma($ A.9 ) and (A.13) $=G$ (A.12). Thus, (A.10) and (A.13) are superfluous and can be omitted. Thus, (A.8)-(A.13) reduce to

$$
\begin{align*}
0= & A Q+Q A^{T}+V_{1}+\gamma^{-2}(Q+\hat{Q}) R_{1 \infty}(Q+\hat{Q})+\beta^{2} \gamma^{-2} \hat{Q} S^{T} P \Sigma P S \hat{Q} \\
& +(A-\Sigma P S) \hat{Q}+\hat{Q}(A-\Sigma P S)^{T},  \tag{A.24}\\
0= & {\left[(A-\Sigma P S) \hat{Q}+\hat{Q}(A-\Sigma P S)^{T}+Q \bar{\Sigma} Q+\gamma^{-2}(Q+\hat{Q}) R_{1 \infty}(Q+\hat{Q})-\gamma^{-2} Q R_{1 \infty} Q\right.} \\
& \left.+\beta^{2} \gamma^{-2} \hat{Q} S^{T} P \Sigma P S \hat{Q}\right] \Gamma^{T},  \tag{A.25}\\
0= & \left(A+\gamma^{-2}[Q+\hat{Q}] R_{1 \infty}\right)^{T} P+P\left(A+\gamma^{-2}[Q+\hat{Q}] R_{1 \infty}\right)+R_{1} \\
& +\left(A-Q \bar{\Sigma}+\gamma^{-2} Q R_{1 \infty}\right)^{T} \hat{P}+\hat{P}\left(A-Q \bar{\Sigma}+\gamma^{-2} Q R_{1 \infty}\right),  \tag{A.26}\\
0= & {\left[\left(A-Q \bar{\Sigma}+\gamma^{-2} Q R_{1 \infty}\right)^{T} \hat{P}+\hat{P}\left(A-Q \bar{\Sigma}+\gamma^{-2} Q R_{1 \infty}\right)+S^{T} P \Sigma P S\right] G^{T} } \tag{A.27}
\end{align*}
$$

Next, using (A.24) $+G^{T} \Gamma$ (A.25) $G-(\mathrm{A} .25) G-[(\mathrm{A} .25) G]^{T}$ and $G^{T} \Gamma$ (A.25) $G-(\mathrm{A} .25) G-[(\mathrm{A} .25) G]^{T}$ yields (6.5) and (6.7). Similarly, using (A.26) $+\Gamma^{T} G(A .27) \Gamma$-(A.27) $\Gamma$ - $[(A .27) \Gamma]^{T}$ and $\Gamma^{T} G(A .27) \Gamma-(A .27) \Gamma-[(A .27) \Gamma]^{T}$ yields (6.6) and (6.8).

Finally, to prove the converse we use (6.5)-(6.13) to obtain (2.14) and (A.4)-(A.7). Let $A_{c}, B_{c}, C_{c}, G, \Gamma, \tau, Q, P, \hat{Q}, \hat{P}, \mathcal{Q}$ be as in the statement of Theorem 6.1 and define $Q_{1}, Q_{12}, Q_{2}$, $P_{1}, P_{12}, P_{2}$ by (A.21)-(A.23). Using (6.2), (6.11) and (6.12) it is easy to verify (A.6) and (A.7). Finally, substitute the definitions of $Q, P, \hat{Q}, \hat{P}, G, \Gamma$, and $\tau$ into (6.5)-(6.8) using (6.2), (6.3), and (A.19) to obtain (2.14) and (A.4). Finally, note that

$$
\mathcal{Q}=\left[\begin{array}{cc}
Q & 0_{n \times n_{0}} \\
0_{n_{0} \times n} & 0_{n_{0}}
\end{array}\right]+\left[\begin{array}{c}
I_{n} \\
\Gamma
\end{array}\right] \hat{Q}\left[\begin{array}{ll}
I_{n} & \Gamma^{T}
\end{array}\right]
$$

which shows that $Q \geq 0$.

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## Abstract

In a recent paper a unification of the $L_{2}$ (LQG) and $H_{\infty}$ control design problems was obtained in terms of modified algebraic Riccati equations. In the present paper these results are extended to guarantee robust $L_{2}$ and $\mathrm{H}_{\infty}$ performance in the presence of structured realvalued parameter variations ( $\Delta \mathrm{A}, \Delta \mathrm{B}, \Delta \mathrm{C}$ ) in the state space model. For design flexibility the paper considers two distinct types of uncertainty bounds for both full- and reduced-order dynamic compensation. An important special case of these results generates $\underline{L}_{2} / \mathrm{H}_{\infty}$ controller designs with guaranteed gain margins.

## I. Introduction

It has recently been shown that the solution to the optimal $H_{\infty}$ disturbance attenuation problem can be expressed in terms of a pair of modified Riccati equations $([1,2])$ Furthermore, it was shown in [1] that $L_{2} / H_{\infty}$ design tradeoffs can be achieved by solving a coupled system consisting of three modified Riccati equations. As is well known, the disturbance attenuation problem can be used to guarantee robustness with respect to unstructured plant uncertainties. However, if plant uncertainty is present in the form of parametric variations of the state space model, then alternative bounding techniques are required. The goal of the present paper is thus to extend the results of [1] to include bounds on the effect of real-valued structured parameter variations.

In the absence of an $\mathrm{H}_{\mathrm{x}}$ design constraint, robust stability and $L_{2}$ performance for dynamic compensator design were guaranteed in [3.4] by incorporating quadratic lyapunov bounds within LQG design theory. Two distinct bounds were considered. In [3] a quadratic bound was used while in [4] a linear bound was employed. In each case futl- and reduced-order dynamic compensators were characterized by means of coupled systems of modified Riccati and Lyapunov equations.

To design $\mathrm{Hl}_{\infty}$ controllers which are robust with respect to structured real-valued parameter variations we proceed by combining the results of [1] with those of [3.4]. That is, we derive coupled systems of modified Riccati and Lyapunov equations whose solutions yield controllers which are guaranteed to satisfy a prespecified lloo attenuation constraint for all variations ( $\Delta \Lambda$, $\Delta \mathrm{Al}, \Delta \mathrm{C})$ belonging to a given uncertainty set. If the uncertainty is absent (i.e., $\Delta A=0$, etc.) then the results of [1] are recovered, while if the $H_{\infty}$ constraint is relaxed, then the results of $[3,4]$ are obtained. Tlus the results of [1] can be viewed as a specializa$t$ ion of a broader design theory which accounts for structured real-valued parameter uncertainty. Finally, we state all results for the case of a fixed-order (i.e., reduced-order) controller for maximal design

[^30]flexibility. Extensions to even more general design problems are mentioned in Section 9 but omitted here for lack of space.

Notation. Note: All matrices have real entries.
$\mathbb{R}, \mathbb{R}^{r x s}, \mathbb{R}, E \quad$ real numbers, $r \times s$ real matrices, $\mathbb{R r} \times 1$, expected value

| $\mathrm{I}_{\mathrm{r}},()^{\top}, \mathbf{0}_{\text {r x }}, \mathrm{O}_{\mathrm{r}}$ | $\mathrm{r} x \mathrm{r}$ identity matrix. transpose. $r \times s$ zero matrix, $0_{r x r}$ |
| :---: | :---: |
| $\mathrm{Sr}, \mathrm{Nr}, \mathrm{Pr}$ | r x r symmetric, nonnegative definite, positive-definite matrices |
| $\mathrm{Z}_{1} \leq \mathrm{Z}_{2}, \mathrm{Z}_{1}<\mathrm{Z}_{2}$ | $\begin{aligned} & Z_{2}-Z_{1} \epsilon \mathbb{M r}, Z_{2}-Z_{1} \in \mathbb{P r}, \\ & Z_{1} ; Z_{2} \in S^{r} . \end{aligned}$ |
| $n, m, \ell, n_{s}$ | positive integers |
| $\mathrm{p}, \mathrm{dx}, \mathrm{q}, \mathrm{qx} ; \tilde{\mathrm{n}}$ | positive integers: $n+n_{c}\left(n_{c} \leq n\right)$ |
| $\mathrm{x}, \mathrm{u}, \mathrm{y}, \mathrm{x}_{\mathrm{c}}, \mathrm{x}^{\text {en }}$ | $\mathrm{n}, \mathrm{m}, \ell, \mathrm{n}_{\mathrm{c}}, \tilde{\mathrm{n}}$-dimensional vectors |
| $A, \Delta A ; B, \Delta B ; C, \Delta C$ | $n \times n ; n \times m ; \ell \times n$ matrices |
| $A_{c}, B_{c}, C_{c}$ | $\mathrm{n}_{\mathrm{c}} \times \mathrm{n}_{\mathrm{c}}, \mathrm{n}_{\mathrm{c}} \times \ell, m \times n_{c}$ matrices |
| $\overline{\mathrm{x}}, \overline{\mathrm{A}}, \Delta \overline{\mathrm{A}}$ | $\left[\begin{array}{l}x \\ x_{c}\end{array}\right],\left[\begin{array}{cc}A & B C_{c} \\ B_{c} C & A_{c}\end{array}\right],\left[\begin{array}{lll}\Delta A & \Delta B C_{c} \\ B_{c} \Delta C & 0\end{array}\right]$ |

d-dimensional standard white noise
$D_{1}, D_{2}$
$\mathrm{n} \times \mathrm{d}, f \times \mathrm{d}$ matrices: $\mathrm{D}_{1} \mathrm{D}_{2}^{\top}=0$
$V_{1} . V_{2}$
D.
$E_{1}, E_{2}$
$D_{1} D_{1}^{\top}, D_{2} D_{2}^{\top}: V_{2} \in \mathbb{P}^{\ell}$
$\left[\begin{array}{c}D_{1} \\ B_{c} D_{2}\end{array}\right],\left[\begin{array}{cc}V_{1} & 0_{n \times n_{c}} \\ 0_{n_{c} \times n} & B_{c} \gamma_{2} B_{c}^{\top}\end{array}\right]$
$\mathrm{L}_{1}, \mathrm{E}_{2}$
$q \times n, q \times m$ matrices: $E_{1}^{\top} E_{2}=0$
$\dot{E}, R_{1}, R_{2}$
$\stackrel{R}{R}$
$E_{1} x, E_{2} x \quad q x \times n, q x \times m$ matrices:
$E_{1}^{\top} x E_{2} x=0$
$\tilde{E}_{\infty}, R_{1} \infty, R_{2 \infty}$
$\left[E_{1} E_{2} C_{c}\right], E_{1}^{\top} E_{1}, E_{2}^{\top} E_{2} ; R_{2} \in \mathbb{P}_{m}$

$$
\left[\begin{array}{cc}
R_{1} & 0_{n \times n_{c}} \\
0_{n_{c} \times n} & C_{c}^{\top} R_{2} C_{c}
\end{array}\right]=\tilde{E}^{\top} \tilde{E}
$$

$\left[\begin{array}{ll}E_{1} \infty & E_{2} x C_{c}\end{array}\right] . E_{1}^{\top} x E_{1} \infty, E_{2}^{\top} \times E_{2} \infty$
$\tilde{R}_{\infty}$

$$
\begin{aligned}
& D_{1} \infty, D_{2 \infty} \\
& \text { n } \times \mathrm{d} \infty, \ell \times \text { d } \infty \text { matrices } ; D_{1} \infty D_{2}{ }^{\top}=0 \\
& V_{1} \infty, V_{2 \infty} \\
& D_{1} \infty D_{1}^{\top} \infty, D_{2} \infty D_{2}^{\top} \\
& \tilde{D}_{\infty}, \bar{V}_{\infty} \\
& {\left[\begin{array}{c}
D_{1 \infty} \\
B_{c} D_{2 \infty}
\end{array}\right],\left[\begin{array}{cc}
V_{1 \infty} & 0_{n \times n_{c}} \\
0_{n_{c} \times n} & B_{c} V_{2 \infty B_{c}}^{\top}
\end{array}\right]} \\
& \beta, \gamma, a \quad \text { nonnegative constant; positive } \\
& \text { constants } \\
& A_{a}, A_{c_{a}} \\
& A+\frac{\alpha}{2} I_{n}, A_{c}+\frac{\alpha}{2} I_{n} .
\end{aligned}
$$

In this section we state the Robust stability and $L_{2}$ performance problem with robust Hoo disturbance attenuation constraint. Specifically, we consider a fixed-order dynamic output-feedback control-design problem with structured real-valued plant parameter uncertainties and constrained $H_{o}$ disturbance attenuation. This problem involves a set $U \subset \mathbb{R n} \times \boldsymbol{n} \times \mathbb{R} \times \times m \times$ $\mathbb{R}^{l \times m}$ of uncertain perturbations ( $\Delta A, \Delta B, \Delta C$ ) of the nominal system matrices ( $A, B, C$ ). The goal of the problem is to determine a fixed-order, strictly proper dynamic compensator ( $A_{c}, B_{c}, C_{c}$ ) which (a) stabilizes the plant for all variations in $U$, (b) satisfies an Ho constraint on disturbance rejection for all variations in U, and (c) minimizes the worst-case value over the uncertainty set $U$ of a steady-state $L_{2}$ performance criterion. In this and the following section no explicit structure is assumed for the elements of $U$. In Sections 4 and 7 , two specific structures of variations in $U$ will be introduced.

Hoc - Constrained Robust Dynamic Compensation Problem. Given the nth-order stabilizable and detectable plant with structured real-valued plant parameter variations

$$
\begin{gather*}
\dot{x}(t)=(A+\Delta A) x(t)+(B+\Delta B) u(t)+D_{1} w(t) .  \tag{2.1}\\
y(t)=(C+\Delta C) x(t)+D_{2} w(t) . \tag{2.2}
\end{gather*}
$$

determine an $n_{c}$ th-order dynamic compensator

$$
\begin{align*}
& \dot{x}_{c}(t)=A_{c} x_{c}(t)+B_{c} y(t)  \tag{2.3}\\
& u(t)=C_{c} x_{c}(t), \tag{2.4}
\end{align*}
$$

which satisfies the following design criteria:
(i) the closed-loop system (2.1) - (2.4) is asymptotically stable for all ( $\Delta A . \Delta B, \Delta C ; \in U$, i.e. $\bar{A}+\Delta \bar{i}$ is asymptotically stable for all ( $\Delta \hat{A}, \Delta B, \Delta C$ ) $\epsilon \mathrm{U}$ :
(ii) the $q \times x d$ closed-loop transfer function

$$
\begin{equation*}
H_{\Delta \tilde{A}}(s) \triangleq \tilde{E}_{\infty}\left[s I_{\tilde{n}}-(\tilde{A}+\Delta \tilde{A})\right]^{-1} \tilde{D} \tag{2.5}
\end{equation*}
$$

from $w(t)$ to $E_{1} x x(t)+E_{2 x u}(t)$ satisfies the constraint

$$
\begin{equation*}
\left\|\left\|_{\Delta A}^{-(s)}\right\|_{\infty} \leq \gamma,(\Delta A, \Delta B, \Delta C) \epsilon U,\right. \tag{2.6}
\end{equation*}
$$

(iii) the performance functional

$$
J\left(A_{c}, B_{c}, C_{C}\right) \triangleq \sup _{(\Delta A, \Delta B, \Delta C) \in U} \limsup _{t \rightarrow \infty}
$$

$$
\begin{equation*}
E\left[x^{\top}(t) R_{1} x(t)+u^{\top}(t) R_{2} u(t)\right] \tag{2.7}
\end{equation*}
$$

is minimized.
Note that for each uncertain variation $(\Delta A, \Delta B$, $\Delta C) \in U$, the closed-loop system can be written as

$$
\begin{equation*}
\dot{\bar{x}}(t)=(\tilde{A}+\Delta \tilde{A}) \tilde{x}(t)+\tilde{D} w(t), t \in[0, \infty) \tag{2.8}
\end{equation*}
$$

and that (2.7) becomes

$$
\begin{align*}
J\left(A_{c}, B_{c}, C_{c}\right)= & \sup _{(\Delta A, \Delta B, \Delta C) \epsilon U} \underset{t \rightarrow \infty}{\lim \sup } \\
& E\left[\tilde{x}^{\top}(t) \tilde{R} \tilde{x}(t)\right] \tag{2.9}
\end{align*}
$$

Note that the problem statement involves both $L_{2}$ and $H_{o o}$ performance weights. In particular, the matrices $R_{1}$ and $R_{2}$ are the $l_{2}$ weights for the state and control variables. By introducing the variables

$$
z(t)=E_{1} x(t), v(t)=E_{2} u(t)
$$

the cost (2.7) can be written as

$$
\begin{aligned}
J\left(A_{c}, B_{c}, C_{c}\right)= & \sup _{(\Delta A, \Delta B, \Delta C) \epsilon U} \quad \limsup _{t \rightarrow \infty} \\
& E\left[z^{\top}(t) z(t)+v^{\top}(t) v(t)\right]
\end{aligned}
$$

For convenience we thus define $R_{1} \triangleq E_{1}^{\top} E_{1}$ and $R_{2} \triangleq E_{2}^{\top} E_{2}$ which appear in subsequent expressions. Although an $L_{2}$ cross-weighting term of the form $2 x^{\top}(t) R_{12} u(t)$ can also be included. we shall not do so here to facilitate the presentation.

For the Hx performance constraint, the transfer function (2.5) involves weighting matrices $E_{1} x$ and $E_{2} x$ for the state and control variables. The matrices $R_{1} \infty$ $\triangleq E_{1} \times E_{1} \infty$ and $R_{2} x \triangleq E_{2}{ }^{\top} x E_{2} \infty$ are thus the $H_{x}$ counterparts of the $L_{2}$ weights $R_{1}$ and $R_{2}$. Although we do not require that $R_{1} x$ and $R_{2} x$ be equal to $R_{1}$ and $R_{2}$. we shall require that $R_{2} x=B^{2} R_{2}$, where the nonnegative scalar 3 is a design variable. he further note that the assumption $E_{1}{ }^{\top} E_{2} \infty=0$ precludes an $H x$ crossweighting term which again facilitates the presentation. Finally, similar remarks apply to the disturbance and sensor noise intensities $V_{1} \triangleq D_{1} D_{1}^{\top}, V_{2} \triangleq$ $D_{2} D_{2}^{\top}, V_{1_{\infty}} \triangleq D_{1_{x}} D_{1}^{\top}$ and $V_{2} \triangleq D_{2_{x}} D_{2_{x}}^{\top}$ for the $L_{2}$ and $H_{x}$ designs respectively. We note that $h(t)$ can be interpreted as white noise for the $L_{2}$ design aspect and as an $L_{2}$ signal for the $H_{x}$ design aspect [1].

Before continuing it is useful to note that if $(\tilde{A}+\Delta \tilde{i})$ is asymptotically stable for all ( $\Delta A, \Delta B, \Delta C$ ) $\epsilon U$ for a given compensator ( $A_{c}, B_{c}, C_{c}$ ) then the performance (2.7) is given by

$$
\begin{equation*}
J\left(A_{c}, B_{c}, C_{c}\right)=\sup _{(\Delta A, \Delta B, \Delta C) \in U} \operatorname{tr} \tilde{Q}_{\Delta X} \tilde{A}^{\tilde{R}}, \tag{2.10}
\end{equation*}
$$

where the steady-state closed-loop state covariance defined by

$$
\begin{equation*}
\overline{\mathrm{Q}}_{\Delta \tilde{A}} \triangleq \lim _{\mathrm{t} \rightarrow \infty} \mathrm{E}\left[\tilde{\mathrm{x}}(\mathrm{t}) \bar{x}^{\top}(\mathrm{t})\right] \tag{2.11}
\end{equation*}
$$

satisfies the $\tilde{\mathbf{n}} \mathbf{x} \overline{\mathrm{n}}$ algebraic Lyapunov equation

$$
\begin{equation*}
0=(\tilde{A}+\Delta \tilde{A}) \tilde{Q}_{\Delta \tilde{A}}+\tilde{Q}_{\Delta \tilde{A}}(\tilde{A}+\Delta \tilde{A})^{\top}+\bar{V} \tag{2.12}
\end{equation*}
$$

The key step in guaranteeing robust stability and performance is to replace the uncertain terms in the covariance lyapunov equation (2.12) by a bounding function $\Omega$. Note that since $\Delta \tilde{A}$ is independent of $A_{c}$, the bounding function $\Omega$ need only depend on $B_{c}$ and $C_{c}$. Furthermore, the $H_{\infty}$ constraint (2.6) on the disturbance attenuation is enforced for all ( $\Delta A, \Delta D, \Delta C$ ) $\in U$ by replacing the modified algebraic Lyapunov equation (2.12) by an algebraic Riccati equation which overbounds the closed-loop steady-state covariance. Justification for this technique is provided by the following result.

Lemma 2.1. Let $\Omega: \quad \mathbb{R}^{n_{c} \times \mathbb{I}} \times \mathbb{R}^{m \times n_{c}} \times \mathbb{N}^{\tilde{n}} \longrightarrow S^{\bar{n}}$ be such that

$$
\begin{align*}
\Delta \tilde{A} \mathbb{Q}+\emptyset \Delta \tilde{A}^{\top} \leq & \Omega\left(B_{c}, C_{c}, 母\right) ;(\Delta A, \Delta B, \Delta C) \epsilon U \\
& \left(B_{c}, C_{c}, 母\right) \in \mathbb{R}^{n_{c} \times \ell^{\prime}} \times \mathbb{R}^{m \times n_{c}} \times \mathbb{N}^{\tilde{n}}, \tag{2.13}
\end{align*}
$$

and, for a given ( $A_{c}, B_{c}, C_{c}$ ), suppose there exists Q $\in \mathbb{N n}^{\text {ñ }}$ satisfying

$$
\begin{equation*}
0=\tilde{A} Q+Q \tilde{A}^{\top}+\gamma^{-2} Q \tilde{R}_{x} Q+\Omega\left(B_{c}, C_{c}, Q\right)+\bar{V} \tag{2.14}
\end{equation*}
$$

then
$\left(\bar{A}+\Delta \tilde{A},\left[\bar{V}+\gamma^{2} \underline{Q} \bar{R}_{x} \mathbf{Q}+\Omega\left(B_{c}, C_{c}, \mathbf{Q}\right)-\left(\Delta \tilde{A} \mathbf{Q}+Q \Delta \tilde{A}^{\top}\right)\right]^{\frac{1}{2}}\right)$ is stabilizable for all $(\Delta A, \Delta B, \Delta C) \epsilon U$.
if and only if
$\tilde{A}+\Delta \tilde{A}$ is asymptotically stable, $(\Delta A, \Delta B, \Delta C) \in U$
In this case,

$$
\begin{equation*}
\left\|H_{\Delta \tilde{A}}(\mathrm{~s})\right\|_{\infty} \leq r, \quad(\Delta A, \Delta \mathrm{~B}, \Delta \mathrm{C}) \in \mathrm{U} \tag{2.17}
\end{equation*}
$$

and

$$
\begin{equation*}
\tilde{\mathbb{Q}}_{\Delta \tilde{A}} \leq \mathbf{Q}, \quad(\Delta A, \Delta B, \Delta C) \epsilon U \tag{2.18}
\end{equation*}
$$

where $\tilde{\mathrm{a}}_{\Delta \tilde{A}}$ is given by (2.12). Consequently,

$$
\begin{equation*}
J\left(A_{c}, B_{c}, C_{c}\right) \leq J\left(A_{c}, B_{c}, C_{c}, \mathbf{Q}\right), \tag{2.19}
\end{equation*}
$$

where

$$
\begin{equation*}
J\left(A_{c}, B_{c}, C_{c}, \mathbf{q}\right) \stackrel{\Delta}{=} \operatorname{tr} \mathbf{q} \tilde{R} \tag{2.20}
\end{equation*}
$$

Proof. First note for clarity that in (2.13) $\mathbb{Q}$ denotes an arbitrary element of $\mathrm{Nn}^{\mathrm{n}}$ since (2.13) holds for all $\mathbf{Q} \in \mathbb{N n}$. while in (2.14) $\mathbf{Q}$ donotes a specific
solution to (2.14). Now for $(\Delta A, \Delta B, \Delta C) \in U$, (2.14) is equivalent to

$$
\begin{align*}
0= & (\tilde{A}+\Delta \tilde{A}) \mathbf{Q}+\mathbf{Q}(\tilde{A}+\Delta \tilde{A})^{\top}+\gamma^{-2} \mathbf{Q} \tilde{R} \mathbf{Q} \\
& +\Omega\left(B_{c}, C_{c}, \tilde{Q}\right)-\left(\Delta \tilde{\tilde{Q}}+\mathbb{Q}+\tilde{A}^{T}\right)+\tilde{V} . \tag{2.21}
\end{align*}
$$

Hence, by assumption, (2.21) has a solution $\square \in \mathbb{N a n}^{\bar{n}}$ for all $(\Delta A, \Delta B, \Delta C) \in U$ and, by (2.13), $\Omega\left(B_{c}, C_{c}, Q\right)$ - $\left(\Delta \bar{A} \cap+Q \Delta \tilde{A}^{\top}\right)$ is nonnegative definite. Now if the stabilizability condition (2.15) holds for all ( $\Delta A, \Delta B$, $\Delta C) \in U$, it follows from Lemna 12.2 of [5] that $\bar{A}+\Delta \bar{A}$ is asymptotically stable for all ( $\Delta A, \Delta B, \Delta C$ ) $\in U$.
Conversely, if $\tilde{A}+\Delta \tilde{A}$ is asymptotically stable for all ( $\Delta A, \Delta B, \Delta C$ ) $\epsilon U$, then ( 2.15 ) holds. The proof of (2.17) now follows from a standard manipulation of (2.14). Next, subtracting (2.12) from (2.20) yields

$$
\begin{aligned}
& 0=(\tilde{A}+\Delta \tilde{A})\left(\mathbf{Q}-\tilde{Q}_{\Delta \tilde{A}}\right)+\left(\mathbf{Q}-\tilde{Q}_{\Delta \tilde{A}}\right)(\tilde{A}+\Delta \bar{A})^{\top}+\gamma^{-2} \mathbf{Q} \tilde{R}_{x} \mathbf{Q} \\
& +\Omega\left(B_{c}, C_{c}, \mathbf{Q}\right)-\left(\Delta \bar{Q} \mathbf{Q}+\mathbf{Q} \Delta \tilde{A}^{\top}\right) .
\end{aligned}
$$

or, equivalently, since $\tilde{A}+\Delta \tilde{A}$ is asymptotically stable for all $(\Delta A, \Delta B, \Delta C) \in U$,

$$
\begin{gathered}
\mathbf{Q}-\tilde{Q}_{\Delta \tilde{A}}==_{0}^{\infty} e^{(\tilde{A}+\Delta \tilde{A}) t}\left[\gamma^{-2} \mathbf{Q} \tilde{R}_{x} \mathbf{Q}+\Omega\left(B_{c}, \mathbf{C}_{c}, \mathbf{Q}\right)\right. \\
-\left(\Delta \tilde{\left.\left.\tilde{A} \mathbf{Q}+\mathbf{Q} \tilde{A}^{\top}\right)\right] \mathrm{e}^{(\overline{\mathrm{A}}+\Delta \tilde{\mathrm{A}})^{\top} \mathrm{t}} \mathrm{dt} \geq 0}\right.
\end{gathered}
$$

which implies (2.18). The performance bound (2.19) is now an immediate consequence of (2.18).

Remark 2.2. Note that (2.15) is actually a closed-loop "disturbability" condition which is not concerned with control.as such. This condition guarantees that the closed-loop system does not possess unstable undisturbed modes. In applying Lemma 2.1 it may be convenient to replace condition (2.15) with a stronger condition which is easier to verify in practice. Clearly. (2.15) is satisfied if $\left[\tilde{\hat{V}}+\gamma^{2} \mathbf{Q R}_{\mathrm{R}} \mathrm{Q}\right.$ $\left.+\Omega\left(B_{c}, C_{c} . Q\right)-\left(\Delta \tilde{A} Q+Q \Delta A^{\top}\right)\right]$ is positive definite for all ( $\Delta A, \Delta B, \Delta C) \in U$. This will be the case, for example, if either $\bar{V}$ is positive definite or strict inequality holds in (2.13) Also, it follows from Theorem 3.6 of [5] that (2.15) is implied by the stronger condition that
$(\tilde{A}+\Delta \tilde{A}, \tilde{D})$ is stabilizable. $(\Delta A, \Delta B, \Delta C) \in U$.

## III. The Aluxiliarv Minimization Problem

As shown in the previous section, the replacement of (2.12) by (2.14) enforces the $H_{x}$ disturbance attenuation constraint and yields an upper bound for the worst case $L_{2}$ performance criterion. That is, given a compensator ( $A_{c}, B_{c}, C_{c}$ ) for which there exists a non-negative-definite solution to (2.14), the actual worst case $L_{2}$ performance ( $A_{c}, B_{c}, C_{c}$ ) of the compensator is guaranteed to be no torse than the bound given by $\boldsymbol{J}\left(A_{c}, B_{c}, C_{c}, \mathbb{Q}\right)$. Hence, $J\left(A_{c}, B_{c}, C_{c}, \mathbf{q}\right)$ can be interpreted as an auxiliary cost which leads to the following optimization problem.

Auxiliarv Vinimization Prohlem. Determine $\left(A_{c}\right.$, $B_{c}, C_{c}$, Q) which minimizes $J\left(A_{c}, B_{c}, C_{c}, ~ Q\right)$ subject to (2.14) with $\ \in \mathbb{N}^{\mathrm{N}}$.

It follows from Lemma 2.1 that the satisfaction of (2.14) for $Q \in \mathbb{N}^{n}$ along with the generic condition
(2.15) leads to 1) closed-loop stability for all ( $\Delta A$, $\Delta B, \Delta C) \epsilon U$; 2) prespecified $H_{\infty}$ performance attenuation for all ( $\Delta A, \Delta B, \Delta C$ ) $\epsilon U$; and 3 ) an upper bound for the worst case $L_{2}$ performance criterion. Hence, it remains to determine ( $A_{c}, B_{c}, C_{c}$ ) which minimizes $J$ ( $A_{c}, B_{c}$, $\left.C_{c}, ~ q\right)$ and thus provides an optimized bound for the actual worst case $L_{2}$ performance $J\left(A_{c}, B_{c}, C_{c}\right)$ over all $(\Delta A, \Delta B, \Delta C) \in U$.

## IV. Uncertainty Structure: Linear Bound

Having established the theoretical basis for our approach, we now assign explicit structure to the set $U$ and bounding function $\Omega$. Specifically, the uncertainty set $U$ is assumed to be of the form

$$
\begin{gather*}
\mathrm{U}=\left\{(\Delta \mathrm{A}, \Delta \mathrm{~B}, \Delta \mathrm{C}): \Delta \mathrm{A}=\sum_{i=1}^{p} \sigma_{i} \mathrm{~A}_{\mathrm{i}}, \Delta \mathrm{~B}=\sum_{i=1}^{p} \sigma_{i} \mathrm{~B}_{\mathrm{i}},\right. \\
\left.\Delta \mathrm{C}=\sum_{\mathrm{i}=1}^{p} \sigma_{\mathrm{i}} \mathrm{C}_{\mathrm{i}}, \quad \sum_{i=1}^{p} \sigma_{\mathrm{i}}^{2} / a_{i}^{2} \leq 1\right\}, \tag{4.1}
\end{gather*}
$$

where, for $i=1, \ldots, p: A_{i} \in \mathbb{R n \times n}, B_{i} \in \mathbb{R n \times m}$, and $C_{i}$ $\epsilon \mathbb{R} \ell \times n$ are fixed matrices denoting the structure of the parametric uncertainty; $a_{i}$ is a given positive number; and $\sigma_{\mathrm{i}}$ is an uncertain real parameter. Note that the uncertain parameters $\sigma_{\mathrm{i}}$ are assumed to lie in a speci-
fied ellipsoidal region in $\mathbb{R}^{p}$. The closed- loop system (2.8) thus has structured uncertainty of the form

$$
\begin{equation*}
\Delta \tilde{A}=\sum_{i=1}^{p} \sigma_{i} \tilde{A}_{i}, \tag{4.2}
\end{equation*}
$$

where

$$
\tilde{X}_{i} \triangleq\left[\begin{array}{lr}
A_{i} & B_{i} C_{c} \\
B_{c} C_{i} & 0
\end{array}\right] ; i=1, \ldots, p
$$

Note that the symmetry of the uncertainty set entails no loss of generality by requiring only a redefinition of the nominal plant matrices.

In order to obtain explicit gain expressions for ( $A_{c}, B_{c}, C_{c}$ ) in Sections 5 and 6 , we shall require that at most one of the perturbations $\Delta B$ and $\Delta C$ is nonzero. he thus consider the cases ( $\Delta A, \Delta C$ ) $\epsilon \mathrm{U}$ or $(\Delta A, \Delta B) \epsilon$ U. If this assumption is not imposed, then optimality conditions can still be derived, but at the expense of closed-form gain expressions. In this section, and Section 5 we will assume that $\Delta B=0$ (i.e., $B_{i}=0$, $i=$ $1, \ldots, p)$ and $\Omega\left(B_{c}, C_{c}, \mathbf{q}\right)$ becomes $\Omega\left(B_{c}, \mathbf{q}\right)$. The dual case $\Delta B \neq 0, \Delta C=0$ (i.e., $C_{i}=0, i=1, \ldots, p$ ) will be considered in Section 6.

For the structure of $U$ specified by (4.1), the bound $\Omega$ satisfying (2.13) can now be given a concrete form.

Proposition 4.1. Let $\alpha$ be an arbitrary positive scalar. Then the function
satisfies (2.13) with $U$ given by (4.1) and $\Delta B=0$.
Proof. See [4]. व
Remark 4.1. Note that the bound $\Omega$ given by (4.3) consists of two distinct terms. The first term a can be thought of as arising from an exponential time weighting of the cost, or, equivalently, from a uniform
right shift of the open-loop dynamics. The second term $\alpha^{-1} \sum_{i=1}^{f} a_{i}^{2} \bar{A}_{i} \eta \tilde{A}_{i}^{\top}$ arises naturally from a multiplicative white noise model. Such interpretations have no bearing on the results obtained here since, only the bound $\Omega$ defined by (4.3) is required. We call (4.3) the linear bound since it is linear in 9 . For a more detailed discussion on (4.3) see [4].

With $\Omega$ defined by (4.3), the modified Riccati equation (2.14) becomes

$$
\begin{equation*}
0=\tilde{A} Q+Q \tilde{A}^{\top}+\gamma^{2} \mathbf{Q} \tilde{R}_{\infty} Q+a \mathbb{Q}+\alpha^{-1} \sum_{i=1}^{\mathcal{L}} a_{i}^{2} \tilde{A}_{i} Q \tilde{A}_{i}^{\top}+\bar{V} \tag{4.4}
\end{equation*}
$$

or, equivalently,

$$
\begin{equation*}
0=\tilde{A}_{\alpha} \mathbf{Q}+\mathbf{Q} \tilde{A}_{\alpha}^{\top}+\gamma^{2} \mathbf{Q} \tilde{R}_{\alpha} \mathbf{Q}+{\underset{\mathrm{i}}{1}}_{\mathcal{F}}^{=_{1}} \delta_{\mathrm{i}} \tilde{\lambda}_{\mathrm{i}} \mathbf{Q} \tilde{A}_{\mathrm{i}}^{\top}+\tilde{V}, \tag{4.5}
\end{equation*}
$$

where $\delta_{i} \triangleq a_{1}^{2} / a$ and
$\tilde{A}_{a} \triangleq \tilde{A}+\frac{\alpha}{2} I_{\tilde{n}}=\left[\begin{array}{ll}A & B C_{c}^{c} \\ B_{c} C & A_{c}\end{array}\right]$.
V. Sufficient Conditions for Robust Stability and Performance with Robust Ho Disturbance Attentuation: Linear Bound

In this section we state sufficient conditions for characterizing fixed-order (i.e., full- and reducedorder) controllers guaranteeing closed-loop stability for all $(\Delta A, \Delta C) \in U$, constrained $H_{x}$ disturbance attenuation for all ( $\Delta A, \Delta C$ ) $\epsilon U$, and an optimized worst case $\mathrm{L}_{2}$ performance bound.

In order to state the main results we require some additional notation and a factorization lemma.

Lemma 5.1. Let $\dot{\mathrm{Q}}, \hat{\mathrm{P}} \in \mathbb{N}^{\mathrm{n}}$ and suppose rank $\dot{\mathrm{Q}} \dot{\mathrm{P}}=\mathrm{n}_{\mathrm{c}}$ then there exist $n_{c} \times n \quad G, \Gamma$ and $n_{c} \times n_{c}$ invertible $M$, unique except for a change of basis in $\mathbb{R}^{n} c$. such that

$$
\begin{equation*}
\hat{Q} \hat{P}=G^{\top} W \Gamma . \quad \Gamma G^{\top}=I_{n_{c}} . \tag{5.1}
\end{equation*}
$$

Furthernore, the nxn matrices

$$
\begin{equation*}
\tau \triangleq G^{\top} \Gamma, \quad \tau_{1} \triangleq I_{n}-\tau \tag{5.3}
\end{equation*}
$$

are idempotent and rank $n_{c}$ and $n-n_{c}$, respectively. Finally, if $P \in \mathbb{N}^{n}$ and $B \geq 0$ then the inverse

$$
\begin{equation*}
S \triangleq\left(I_{n}+32 \gamma^{-2} \hat{\mathbf{Q}}\right)^{-1} \tag{5.5}
\end{equation*}
$$

exists.
For arbitrary $\mathbb{Q}, \dot{\mathbf{Q}} \in \mathbb{R}^{n \times n}$ and $a>0$ define the following notation:

$$
\begin{aligned}
& V_{2 s} \triangleq V_{2}+\sum_{i=1}^{p} \delta_{i} C_{i}(Q+\dot{Q}) C_{i}^{\top}, \\
& Q_{\mathrm{s}} \triangleq Q C^{\top}+\sum_{i=1}^{p} \delta_{i} A_{i}(Q+\dot{Q}) C_{i}^{\top}, \quad \Sigma \triangleq B R_{2}^{-1} B^{\top} .
\end{aligned}
$$

Theorem 5.1. Suppose there exist $\mathbb{Q}, \mathbf{P}, \hat{\mathbb{Q}}, \hat{\mathrm{P}} \in \mathbb{N}$ satisfying

$$
\begin{align*}
& 0=A_{a} \emptyset+Q A_{\alpha}^{\top}+\gamma^{-2} Q R_{1 \infty} Q+V_{1} \\
& +\sum_{i=1}^{p} \delta_{i} A_{i}(Q+\hat{Q}) A_{i}^{\top}-Q_{S} V_{2 s}^{-1} Q_{s}^{\top}+\tau_{1} Q_{S} V_{2 s}^{-1} Q_{s}^{\top} \tau_{1}^{\top},  \tag{5.6}\\
& U=\left(A_{a}+\gamma^{-2}[Q+\hat{Q}] R_{i_{\infty}}\right)^{\top} P+P\left(A_{\alpha}+\gamma^{-2}[Q+\hat{Q}] R_{i_{\infty}}\right)+R_{i} \\
& +\sum_{i=1}^{p} \delta_{i}\left[A_{i}^{\top} P A_{i}+\left(A_{i}-Q_{b} V_{2 s}^{-1} C_{i}\right)^{\top} \dot{P}\left(A_{i}-Q_{s} V_{2 s}^{-1} C_{i}\right)\right] \\
& -S^{\top} P E P S+\tau_{1}^{\top} S^{\top} P E P S \tau_{1},  \tag{5.7}\\
& 0=\left(A_{a}-\Sigma P S+\gamma^{-2} Q R_{1_{\infty}}\right) \hat{Q}+\hat{Q}\left(A_{a}-\Sigma P S+\gamma^{-2} Q R_{1_{\infty}}\right)^{\top}+ \\
& \gamma^{-2} \dot{Q}\left(R_{1_{\infty}}+\beta^{2} S^{\top} P \Sigma P S\right) \hat{Q}+Q_{s} V_{2 \mathrm{~s}}^{-1} Q_{s}^{\top}-\tau_{\perp} Q_{S} V_{2 \mathrm{~s}}^{-1} \mathrm{Q}_{\mathrm{s}}^{\top} \tau_{\perp}^{\top},  \tag{5.8}\\
& 0=\left(A_{a}-Q_{s} V_{2 s}{ }^{-1} C+\gamma^{-2} Q_{1_{1}}\right)^{\top} \dot{P}+\dot{P}\left(A_{\alpha}-Q_{s} V_{2 s}^{-1} C\right. \\
& \left.+\gamma^{-2} \mathrm{QR}_{1_{\infty}}\right)+\mathrm{S}^{\top} \mathrm{PEPS}-\tau_{+}^{\top} \mathrm{S}^{\top} \mathrm{PEPS} \tau_{1}, \tag{5.9}
\end{align*}
$$

$$
\begin{equation*}
\operatorname{rank} \hat{\mathbf{Q}}=\operatorname{rank} \hat{\mathbf{P}}=\operatorname{rank} \hat{\mathbf{Q}} \hat{\mathbf{P}}=\mathrm{n}_{\mathbf{c}}, \tag{5.10}
\end{equation*}
$$

and let ( $A_{c}, B_{c}, C_{c}, \mathbf{Q}$ ) be given by

$$
\begin{align*}
& A_{c}=\Gamma\left(A-\Sigma P S-Q_{S} V_{2 s}^{-1} C+\gamma^{-2} Q R_{1_{\infty}}\right) G^{\top}  \tag{5.11}\\
& B_{c}=\Gamma Q_{s} V_{2 s}^{-1},  \tag{5.12}\\
& C_{c}=-R_{2}^{-1} B^{\top} P S G^{\top},  \tag{5.13}\\
& Q=\left[\begin{array}{cc}
Q+\tilde{Q} & \dot{Q} \Gamma^{\top} \\
\Gamma Q & \Gamma \hat{Q} \Gamma^{\top}
\end{array}\right] .
\end{align*}
$$

Then. $(\bar{A}+\Delta \bar{A}, \bar{D})$ is stabilizable if and only if $\bar{X}+\Delta \bar{A}$ is asymptotically stable for all ( $\Delta d, \Delta C$ ) $\epsilon \mathrm{U}$. In this case, the closed-loop transfer function $H_{\Delta i}(\mathrm{~s})$ satisfies the $H_{x}$ disturbance attenuation constraint

$$
\begin{equation*}
\left\|\left\|_{\Delta \tilde{A}}(s)\right\|_{\infty} \leq \gamma,(\Delta A, \Delta C) \in U\right. \tag{5.15}
\end{equation*}
$$

and the worst case $L_{2}$ performance criterion (2.10) satisfies the bound

$$
\begin{equation*}
J\left(A_{c}, B_{c}, C_{c}\right) \leq \operatorname{tr}\left[(\mathbb{Q}+\hat{\mathbf{Q}}) R_{1}+\hat{\mathbf{Q}} S^{\top} P \tilde{C} P S \hat{Q}\right] . \tag{5.16}
\end{equation*}
$$

Proof. The proof follows from Lemma 2.1 by combining the proofs of Theorem 6.1 of [1] and Theorem 6.1 of [4]. व

Remark 5.1: Theorem 5.1 presents sufficient conditions for designing controllers yielding robust stability and performance with a constraint on the How norm of the closed-loop transfer function for a state-space system with structured real-valued plant parameter variations. These sufficient conditions comprise a system of three modified Riccati equations and one modified Lyapunov equation coupled by the optimal projection $r$, the linear uncertainty bound, and the Hoo constraint. If the uncertainty bound is deleted, then the results of [1] are recovered. If, alternatively, the uncertainty terms are retained but the How constraint is sufficiently relaxed, i.e., $\gamma \rightarrow \infty$, then the results of [4] are recovered for the case $B_{i}=0, i=1, \ldots, p$.

Remark 5.2. To specialize Theorem 5.1 to the full order case $n_{c}=n$, it is only necessary to set $G^{\top}=[-1$ so that $G=\Gamma=\tau=I_{n}$ and $\tau_{\perp}=0$ without loss of generality. Now the last term in each of (5.6) - (5.9) can be deleted and $G$ and $\Gamma$ in (5.11) - (5.14) can be taken to be the identity. It is interesting to note that in the full-order case the Hoo design problem with structured parameter variations is comprised of four coupled Riccati/Lyapunov equations. This coupling illustrates the breakdown of regulator/est imator sejaration and shows that the certainty equivalence principle is no longer valid. This is not surprising since separation also breaks down for the full-order $L_{2}$ result with parameter uncertainties [4].

Remark 5.3. When solving (5.6) - (5.10) numerically, the uncertainty terms and the $H_{\infty}$ constraint can be adjusted to examine tradeoffs among performance, robustness, and disturbance rejection. Specifically, the uncertainty range $\alpha_{i}$ and the structure matrices $A_{i}$. $C_{i}$ appearing in $Q_{s}$ and $V_{2 s}$ along with $\gamma$ can be varied systematically to determine the region of solvability of the design equations (5.6) - (5.9).

Remark 5.4. Although equations (5.6) - (5.10) appear formidable, they are, in fact, quite numerically tractable. For related problems involving coupled Riccati equations, homotopic continuation methods have been shown to be effective (see [1] and the References therein).

Remark 5.5. We point out that if $\beta=0$ or, equivalently, $E_{2_{\infty}}=0$, which corresponds to the "cheap" $\mathrm{H}_{\infty}$ control case (i.e., How attenuation between disturbances and controls is not constrained), it is possible to obtain closed-form gains ( $A_{c}, B_{c}, C_{c}$ ) given by a modified set of design equations when all three of $\Delta A, \Delta B$, and $\Delta C$ are nonzero. Because of space limitations this result will be given in a future paper.

Bemark 5.6. An important special case of the results of Section 5 is obtained by setting $\Delta t=0$, $\Delta B=0, \Delta C=\sigma_{1} C_{1}$, and $C_{1}=C$. The resulting $L_{2} / H_{\infty}$ design possesses guaranteed gain margin of $\pm 100 \sigma_{1}$ percent at the sensor output.

## VI. The Dual Case: Linear Bound

Cnlike the standard LQG result involving a pair of uncoupled Riccati equations, the new result guaranteeing robust stability, robust performance, and disturbance rejection involves a coupled system of four modified Riccat $1 /$ Lyapunov equations. The asymmetry of these equations suggests the possibility of a dual result in which the modifications to the standard Riccati equations are reversed. One motivation for developing such dual results is that for certain problems the dual bounds may be sharper than the primal bound introduced in Section 4 . Furthermore, the dual theory permits distinct $\|_{\infty}$ disturbance weights ( $V_{1_{\infty}}$ and $V_{2_{\infty}}$ ), although ne now requre $R_{1_{\infty}}=R_{2}$. Finally, the dual theory allows for uncertainty in the control matrix $B$ (i.e., $\Delta B \neq 0$ ). Although we now require $\Delta C=0$, (i.e., $C_{i}=$ $0, i=1, \ldots, p$ ) to obtain closed-form gain expressions for ( $A_{c}, B_{c}, C_{c}$ ). We begin with the following lemma:

Lemma 6.1. Suppose the system (2.8) is asymptotically stable for all ( $\Delta A, \Delta B, \Delta C$ ) $\in \mathrm{U}$ for a given ( $A_{c}$. $B_{c}, C_{c}$ ). Then

$$
\begin{equation*}
J\left(A_{c}, B_{c}, C_{c}\right)=\quad \sup \quad \operatorname{tr} \tilde{P}_{\Delta A^{-}} \bar{V} \tag{6.1}
\end{equation*}
$$

where $\bar{P}_{\Delta \tilde{A}} \in \operatorname{Man}^{n}$ is the unique solution to

$$
\begin{equation*}
0=(\tilde{A}+\Delta \tilde{A})^{\top} \tilde{P}_{\Delta \tilde{A}}+\tilde{P}_{\Delta \tilde{A}}(\tilde{A}+\Delta \tilde{A})+\tilde{R} . \tag{6.2}
\end{equation*}
$$

Proof. See [3]. $\quad$ a
Utilizing (6.1) in place of (2.10), the $\mathrm{H}_{\infty}$ disturbance attenuation constraint from plant and sensor disturbances to the state and control variables given by

$$
\begin{equation*}
\left\|\tilde{\bar{A}}_{\Delta \tilde{A}}(s)\right\|_{\infty}=\left\|\tilde{E}\left[\mathrm{SI}_{\tilde{\mathrm{A}}}-(\tilde{A}+\Delta \tilde{\mathrm{A}})\right]^{-1} \tilde{\mathrm{D}}_{\infty}\right\|_{\infty} \leq \gamma \tag{6.3}
\end{equation*}
$$

can now be enforced by replacing (2.14) by the modified Riccati equation

$$
\begin{align*}
0= & \tilde{A}^{\top} \mathbf{P}+\mathbf{P} \tilde{A}+\gamma^{-2} \mathbf{P} \tilde{V}_{\infty} \mathbf{P}+\hat{\Omega}\left(\mathrm{C}_{\mathbf{c}}, P\right)+\tilde{\mathrm{R}}, \\
& (\Delta \mathrm{~A}, \Delta \mathrm{~B}, \Delta \mathrm{C}) \epsilon \mathrm{U}, \tag{6.4}
\end{align*}
$$

where

$$
\begin{equation*}
\Delta \tilde{A}^{\top} P+P \Delta \bar{A} \leq \hat{\Omega}\left(C_{c}, P\right),(\Delta A, \Delta B, \Delta C) \epsilon U . \tag{6.5}
\end{equation*}
$$

Note that (6.4) is merely the dual of (2.14). We also require the condition dual to (2.15) given by
$\left(\left[\tilde{\mathbf{R}}+\gamma^{2} \mathbf{P} \hat{V}_{\infty} \mathbf{P}+\hat{\mathbf{R}}\left(\mathrm{C}_{\mathrm{c}}, \mathbf{P}\right) \cdot\left(\Delta \tilde{A}^{\top} \mathbf{P}+\mathbf{P} \Delta \tilde{\mathrm{A}}\right)\right]^{\frac{1}{2}}, \tilde{\mathrm{~A}}+\Delta \tilde{\mathrm{A}}\right)$
is detectable for all $(\Delta A, \Delta B) \epsilon U$.
For the structure of $U$ as specified by (4.1) with $\Delta C=0$, the bound $\hat{\Omega}$ satisfying (6.5) can now be given a concrete form.

Proposition 6.1. Let a be an arbitrary positive scalar. Then the function

$$
\begin{equation*}
\bar{\Omega}\left(C_{c}, P\right) \triangleq a P+a^{\prime} \cdot \sum_{i=1}^{p} a_{i}^{2} \bar{A}_{i}^{\top} P \Lambda_{i}^{\top} \tag{6.7}
\end{equation*}
$$

satisfies (6.5) with $U$ given by (4.1) and $\Delta C=0$. with $\dot{\Omega}$ defined by ( $6 . i$ ), the modified dual Riccati equation (6.4) becomes

$$
\begin{equation*}
0=\tilde{A}_{a}^{\top} P+P \tilde{A}_{a}+\gamma^{-2} P \tilde{V}_{\alpha} P+\sum_{i=1}^{P} \delta_{i} \tilde{X}_{i}^{\top} P \tilde{X}_{i}+\tilde{R} . \tag{6.8}
\end{equation*}
$$

We can now state sufficient conditions for robust stability, robust $L_{2}$ performance, and robust disturbance attenuation for the dual problem. For arbitrary $\mathbb{Q}, P, \dot{P} \in \mathbb{R}^{n \times n}$ and $\alpha>0$ define the following notation: $R_{2 s} \triangleq R_{2}+\sum_{i=1}^{p} \delta_{i} B_{i}^{\top}(P+\dot{P}) B_{i}, P_{s} \stackrel{\Delta}{\Delta}{ }^{\top} P+\sum_{i=1}^{P} \delta_{i} B_{i}^{\top}(P+\dot{P}) A_{i}$, $\hat{S} \triangleq\left(I_{n}+\gamma^{-2} \beta^{2}(\hat{P})^{-1}, \quad \bar{\Sigma} \triangleq C^{\top} V_{2}{ }^{1} C\right.$.

Theorem 6.1. Suppose there exist $P, \mathbb{Q}, \dot{P}, \dot{\mathbf{0}} \in \mathbb{N}^{n}$ satisfying (5.10) and

$$
\begin{align*}
0= & A_{a}^{\top} P+P A_{a}+\gamma^{-2} P V_{t_{\infty}} P+R_{1}+\sum_{i=1}^{P} \delta_{i} A_{i}^{\top}(P+\hat{P}) A_{i} \\
& -P_{s}^{\top} R_{2} \frac{1}{s} P_{s}+\tau_{L}^{\top} P_{s}^{\top} R_{2} \frac{1}{s} P_{s} \tau_{L}^{\top}, \tag{6.9}
\end{align*}
$$

$$
\begin{align*}
0= & \left(A_{a}+\gamma^{-2} V_{1_{\infty}}[P+\hat{P}]\right) Q+Q\left(A_{\alpha}+\gamma-2 V_{1_{\infty}}[P+\hat{P}]\right)^{\top} \\
& +V_{1}+\sum_{i=1}^{p} \delta_{i}\left[A_{i} Q A_{i}^{\top}+\left(A_{i}-B_{i} R_{2}^{1} P_{s}\right) \tilde{Q}\left(A_{i}-B_{i} R_{2}{ }_{2}^{1} P_{s}\right)^{\top}\right] \\
& -\hat{S} Q \bar{\Sigma} Q \dot{S}^{\top}+\tau_{\perp} \hat{S} Q \bar{\Sigma} Q \dot{S}^{\top} \tau_{\perp}^{\top}, \tag{6.10}
\end{align*}
$$

$0=\left(A_{\alpha}-\dot{S} Q \bar{\Sigma}+\gamma^{-2} V_{1_{\infty}} P\right)^{\top} \dot{P}+\dot{P}\left(A_{\alpha}-\tilde{S} Q \bar{Q}+\gamma^{-2} V_{t_{\infty}} P\right)+$

$$
\begin{equation*}
\gamma^{-2} \hat{P}\left(Y_{1_{\infty}}+\beta \approx \hat{S} Q \bar{E} Q \hat{S}^{\top}\right) \hat{P}+P_{s}^{\top} R_{2}^{1}{ }_{\delta}^{1} P_{s}-\tau_{\perp}^{\top} P_{s}^{\top} R_{2}^{1} P_{s} \tau_{\perp} \tag{6.11}
\end{equation*}
$$

$0=\left(\mathrm{A}_{\alpha}-\mathrm{BR}_{2}^{1}{ }_{\mathrm{S}}^{1} \mathrm{P}_{\mathrm{s}}+\gamma^{-2} \mathrm{~V}_{1_{\infty}} \mathrm{P}\right) \dot{\mathrm{Q}}+\hat{\mathrm{Q}}\left(\mathrm{A}_{\alpha}-\mathrm{BR}_{2}{ }_{5}^{1} \mathrm{P}_{\mathrm{s}}+\gamma^{\left.-2 \mathrm{~V}_{1}{ }_{\infty} \mathrm{P}\right)^{\top}, ~}\right.$
$+\hat{S} Q \tilde{\Sigma} Q \hat{S}^{\top}-\tau_{\perp}^{\top} \mathbf{S} Q \bar{\Sigma} Q \hat{S}^{\top} \tau_{\perp}^{\top}$,
and let $\left(A_{c}, B_{c}, C_{c}, P\right)$ be given by

$$
\begin{align*}
& A_{c}=\Gamma\left(A-\hat{S} Q \bar{\Sigma}-B R_{2}^{1}{ }_{s}^{1} P_{s}+\gamma^{-2} V_{1 \infty} P\right) G^{\top},  \tag{6.13}\\
& B_{c}=\Gamma \hat{S} Q C^{\top} V_{2}^{-1},  \tag{6.14}\\
& C_{c}=-R_{2}{ }_{5}^{1} P_{s} G^{\top},  \tag{6.15}\\
& P=\left[\begin{array}{cc}
P+\hat{P} & -\dot{P} G^{\top} \\
-G \dot{P} & G P G^{\top}
\end{array}\right] \tag{6.16}
\end{align*}
$$

Then, $(\tilde{E}, \tilde{A}+\Delta \tilde{A})$ is detectable if and only if $\bar{A}+\Delta \bar{A}$ is asymptotically stable for all ( $\Delta \mathrm{A}, \Delta \mathrm{B}$ ) $\in \mathrm{U}$. In this case. the closed-loop transfer function $\dot{H}_{\Delta \tilde{A}}(s)$ satisfies the $H_{\infty}$ disturbance attenuation constraint,

$$
\begin{equation*}
\left\|\tilde{I}_{\Delta \tilde{A}^{-}}(\mathrm{s})\right\|_{\infty} \leq \gamma,(\Delta \mathrm{A}, \Delta \mathrm{~B}) \in \mathrm{U}, \tag{6.17}
\end{equation*}
$$

and the worst case $L_{2}$ performance criterion (6.1) satisfies the bound

$$
\begin{equation*}
\left(A_{c}, B_{c}, C_{c}\right) \leq \operatorname{tr}\left[(P+\hat{P}) V_{t}+\hat{P} \bar{S} Q \tilde{\Sigma} Q \dot{S}^{\top} \dot{P}\right] . \tag{6.18}
\end{equation*}
$$

Remark 6.1. The dual case of Remark 5.6 is obtained by setting $\Delta A=0, \Delta B=\sigma_{1} B_{1}, \Delta C=0$, and $B_{1}=B$. The resulting $\mathrm{L}_{2} / \mathrm{H}_{0}$ design possesses guaranteed gain margin of $\pm 100 \sigma_{1}$ percent at the input.

## VII. Uncertaintv Structure and Sufficient Conditions for Robust Stability and Performance with How Disturbance Attenuation: Quadratic Bound

We now assign a different structure to the uncertainty set $U$ and the bounding function $\Omega$. Specifically, the uncertainty set $U$ is assumed to be of the form

$$
\begin{align*}
& U=\left\{(\Delta A, \Delta B, \Delta C): \Delta A=\sum_{i=1}^{p} F_{i} M_{i} N_{i} G_{i}\right. \\
& \Delta B=\sum_{i=1}^{p} F_{i} M_{i} N_{i} H_{i}, \Delta C=\sum_{i=1}^{p} K_{i} M_{i} N_{i} G_{i}, \\
&\left.M_{i} M_{i}^{\top} \leq \bar{M}_{i}, N_{i}^{\top} N_{i} \leq \bar{N}_{i}, i=1, \ldots, p\right\}, \tag{7.1}
\end{align*}
$$

where, for $i=1, \ldots, p: F_{i} \in \mathbb{R}^{n \times r_{i}}, G_{i} \in \mathbb{R}^{t_{i} \times n}$, $H_{i} \in \mathbb{R}^{t_{i} \times m}$, and $K_{i} \in \mathbb{R}^{l \times r_{i}}$ are fixed matrices denoting the structure of the uncertainty; $\bar{\pi}_{i} \in \mathbb{N}^{\mathbf{r}}$ and $\bar{N}_{i} \in \mathbb{N}^{\mathbf{t}}$ :
are given uncertainty bounds; and $M_{i} \in \mathbb{R}^{r}{ }^{1 \times s} \mathbf{i}$ and $N_{i} \in \mathbb{R}^{\mathbf{j}}{ }^{\mathrm{xt}}{ }^{\mathrm{i}}$ are uncertain matrices.

In order to obtain explicit gain expressions ( $A_{c}$, $B_{c}, C_{c}$ ) we again consider two cases, 1$)(\Delta A, \Delta C) \in U$ with $\Delta B=0$ and 2) $(\Delta A, \Delta B) \in U$ with $\Delta C=0$. When $\Delta B=0$ the closed-loop system has structured uncertainty of the form

$$
\begin{equation*}
\Delta \bar{A}=\sum_{i=1}^{p} \tilde{F}_{i} M_{i} N_{i} \tilde{G}_{i} \tag{7.2}
\end{equation*}
$$

where,

$$
\tilde{F}_{i} \triangleq\left[\begin{array}{c}
F_{i} \\
B_{c} K_{i}
\end{array}\right], \tilde{G}_{i} \triangleq\left[\begin{array}{ll}
G_{i} & 0
\end{array}\right] .
$$

In this case the quadratic bound $\Omega$ satisfying (2.13) can now be given a concrete form.

Proposition 7.1. The function

$$
\begin{equation*}
\Omega\left(B_{c}, \mathbf{Q}\right) \triangleq \sum_{i=1}^{P} \tilde{\mathrm{~F}}_{\mathrm{i}} \bar{M}_{\mathrm{i}} \tilde{\mathrm{~F}}_{\mathrm{i}}^{\top}+\mathbb{Q} \tilde{G}_{\mathrm{i}}^{\top} \bar{N}_{\mathrm{i}} \tilde{G}_{1} \mathbf{Q} \tag{7.3}
\end{equation*}
$$

satifies (2.13) with $U$ given by (7.1) and $\Delta B=0$.
Proof. See [3]. व
Thus, with $\Omega$ defined by (7.3), the modified Riccati equation (7.3) becomes

$$
\begin{align*}
0= & \tilde{A} Q+q \bar{A}^{\top}+\gamma^{2} \mathbf{Q} \tilde{R}_{\infty} Q+\tilde{V} \\
& +\sum_{i=1}^{p}\left[\tilde{F}_{i} \bar{H}_{i} \tilde{F}_{i}^{\top}+q \tilde{G}_{i}^{\top} \bar{N}_{i} \tilde{G}_{1} Q\right] . \tag{7.4}
\end{align*}
$$

For arbitrary $\mathbb{q} \in \mathbb{R}^{n \times n}$ define:

$$
\begin{array}{ll}
Q_{a} \triangleq Q C^{\top}+\sum_{i=1}^{P} F_{i} \bar{M}_{i} K_{i} . & D \triangleq \sum_{i=1}^{P} F_{i} \bar{M}_{i} F_{i}^{\top}, \\
V_{2 a} \triangleq V_{2}+\sum_{i=1}^{P} K_{i} \bar{M}_{i} K_{i}^{\top}, & E \triangleq \sum_{i=1}^{P} G_{i}^{\top} \bar{N}_{i} G_{i} .
\end{array}
$$

Theorem $7 . \frac{1}{2}$. Suppose these exist $\mathbf{O}, \mathbf{P}, \hat{\mathbf{Q}}, \hat{\mathrm{P}} \in \mathbb{N}^{n}$ satisfying (5.10) and

$$
\begin{align*}
0= & A Q+Q A^{\top}+\gamma^{2} Q_{R_{1}} Q+V_{1}+Q E Q+D \\
& -Q_{a} V_{2}^{1} \frac{1}{a} Q_{a}^{\top}+\tau_{\perp} Q_{a} V_{2} \frac{1}{a} Q_{a}^{\top} \tau_{1}^{\top},  \tag{7.5}\\
0= & \left(A+[Q+\hat{Q}]\left[\gamma^{-2} R_{1_{\infty}}+E\right]\right)^{\top} P \\
& +P\left(A+[Q+\hat{Q}]\left[\gamma^{-2} R_{1_{\infty}}+E\right]\right) \\
& +R_{1}-S^{\top} P E P S+\tau_{\perp}^{\top} S^{\top} P \Sigma P S \tau_{\perp}, \tag{7.6}
\end{align*}
$$

$$
0=\left(A-\Sigma P S+Q\left[\gamma^{2} R_{1_{\infty}}+E\right]\right) \hat{Q}+\tilde{Q}(A-\Sigma P S
$$

$$
\left.+\mathbf{Q}\left[\gamma^{-2} R_{1_{\infty}}+E\right]\right)^{\top}+\hat{\mathbb{Q}}\left(\gamma^{-2}\left[R_{1_{\infty}}+3^{2} S^{\top} P \Sigma P S\right]+E\right) \tilde{Q}
$$

$$
\begin{equation*}
+Q_{a} V_{2}^{1} Q_{a}^{\top}-\tau_{1} Q_{a} V_{2}^{\frac{1}{a}} Q_{a}^{\top} \tau_{\perp}^{\top}, \tag{7.7}
\end{equation*}
$$

$$
0=\left(A-Q_{a} V_{2}^{1} \frac{1}{C} C+\mathbb{Q}\left[\gamma^{-2} R_{1_{\infty}}+E\right]\right)^{\top} \dot{P}+\dot{P}\left(A-Q_{a} V_{2}^{-1} C\right.
$$

$$
\begin{equation*}
\left.+Q\left[\gamma^{-2} R_{1_{\infty}}+E\right]\right)+S^{\top} P \Sigma P S-\tau_{\perp}^{\top} S^{\top} P \Sigma P S_{\tau_{\perp}} \tag{7.8}
\end{equation*}
$$

and let $Q$ be given by (5.14) and ( $A_{c}, B_{c}, C_{c}$ ) by

$$
\begin{align*}
& A_{c}=\Gamma\left(A-\Sigma P S-Q_{a} V_{2}^{1} \frac{1}{C}+Q\left[\gamma^{-2} R_{1 \infty}+E\right]\right) G^{\top},  \tag{7.9}\\
& B_{c}=\Gamma Q_{a} V_{2 \frac{1}{a}},  \tag{7.10}\\
& C_{c}=-R_{2}^{1} B^{\top} P S G^{\top} . \tag{7.11}
\end{align*}
$$

Then, $(\tilde{A}+\Delta \tilde{A} . \tilde{D})$ is stabilizable if and $n n l y$ if $\tilde{A}+\Delta \tilde{X}$ is asymptotically stable for all ( $\Delta \mathrm{A}, \Delta \mathrm{C}$ ) $\epsilon \mathrm{U}$. In this case, the closed-loop transfer function $\|_{\Delta A}(\mathrm{~s})$ sat isfies the $H_{\infty}$ disturbance attenuation constraint

$$
\begin{equation*}
\left\|H_{\Delta \tilde{A}}(s)\right\| \infty \leq \gamma,(\Delta A, \Delta C) \in U, \tag{7.12}
\end{equation*}
$$

and the worst case $L_{2}$ performance criterion (2.10) satisfies the bound

$$
\begin{equation*}
J\left(A_{c}, B_{c}, C_{c}\right) \leq \operatorname{tr}\left[(Q+\hat{Q}) R_{1}+\hat{Q} S^{\top} P \Sigma P S Q ̀\right] . \tag{7.13}
\end{equation*}
$$

Proof. The proof follows by combining the proofs of Theorems 6.1 of [1] and Theorem 8.1 of [3]. a

Remark i.2. It is interesting to note that the full-order case $n_{c}=n$ with $G=\Gamma=\tau=I_{n}$ and $\tau_{\perp}=0$
(see Remark 5.1), $\dot{\mathrm{P}}$ plays no role so that ( 7.8 ) is superfluous. Thus, unlike the full-order result for the linear bound involving four equations, the fullorder quadratic bound involves three modified Riccati equations coupled by the uncertainty term and the $\mathrm{H}_{\mathrm{o}}$ constraint. If, alternatively, the reduced-order constraint is retained, but the uncertainty terms are deleted, then the results of [1] are recovered. If, furthermore, the uncertainty terms are retained, but the $H_{\infty}$ constraint is sufficiently relaxed, i.e., $\gamma \rightarrow \infty$, the results of [3] are recovered.

## VIII. The Dual Case: Onadratic Bound

For the structure of $U$ as specified by (7.1) with $\Delta C=0$, the closed-loop system has structured uncertainty of the form

$$
\begin{equation*}
\Delta \tilde{A}=\sum_{i=1}^{p} \tilde{F}_{i} Y_{i} X_{i} \tilde{\mathbf{G}}_{i} \tag{8.1}
\end{equation*}
$$

where
$\tilde{\mathbf{F}} \triangleq\left[\begin{array}{c}\mathrm{F}_{\mathrm{i}} \\ 0\end{array}\right], \quad \tilde{\mathrm{G}}_{\mathrm{i}} \triangleq\left[\begin{array}{lll}\mathrm{G}_{\mathbf{i}} & H_{i} \mathrm{C}_{\mathrm{c}}\end{array}\right]$.
Pronosition 8.1. The function

$$
\dot{\Omega}\left(\mathrm{C}_{\mathrm{C}}, \mathbf{P}\right) \triangleq \sum_{\mathrm{i}=1}^{\mathrm{P}} \tilde{\mathrm{G}}_{\mathrm{i}}^{\top} \overline{\mathrm{N}}_{\mathrm{i}} \tilde{\mathrm{G}}_{\mathrm{i}}+\mathrm{P} \tilde{F}_{\mathrm{i}} \overline{\mathrm{M}}_{\mathrm{i}} \tilde{\mathrm{~F}}_{\mathrm{i}} \mathbf{P}
$$

satisfies (6.5) with $U$ given by (7.1) and $\Delta C=0$.
With $\hat{\Omega}$ defined by (8.2), the modified dual equation ( 6.4 ) becomes

$$
\begin{align*}
0= & \tilde{A}^{\top} P+P \bar{A}+\gamma^{-2 P} \tilde{V}_{x x} P+\tilde{R} \\
& +\sum_{i=1}^{P}\left[\tilde{G}_{i}^{\top} \bar{N}_{i} \tilde{G}_{i}+P \tilde{F}_{i} \bar{M}_{i} \tilde{F}_{i} P\right] . \tag{8.3}
\end{align*}
$$

For arbitrary $\mathbf{P} \in \operatorname{Raxn}$ define:

$$
P_{a} \triangleq B^{\top} P+\sum_{i=1}^{\mathcal{E}} H_{i}^{\top} \bar{N}_{i} G_{i}, \quad R_{2 a} \triangleq R_{2}+\sum_{i=1}^{P} H_{i}^{\top} \bar{N}_{i} H_{i}
$$

Theorem 8.1. Suppose there exists $\mathbf{P}, \mathbf{Q}, \dot{\mathrm{P}}, \dot{\mathbf{Q}} \in \mathbb{N}$ satisfying (5.10) and

$$
\begin{align*}
& 0=A^{\top} P+P A+\gamma^{-2} P V_{1_{\infty}} P+R_{1}+E+P D P \\
& -P_{a}^{\top} R_{2}{ }_{a}^{1} P_{a}+r_{\perp}^{\top} P_{a}^{\top} R_{2}^{1} P_{a} \tau_{\perp}^{\top},  \tag{8.4}\\
& 0=\left(A+\left[\gamma^{-2} V_{1_{\infty}}+D\right][P+\hat{P}]\right) \mathbf{C}+\mathbf{Q}\left(A+\left[\gamma^{-2} V_{1_{\infty}}+D\right]\right. \\
& [P+\hat{P}])^{\top}+V_{1}-\hat{S} Q \tilde{\Sigma} Q \hat{S}^{\top}+\tau_{\perp} \hat{S} Q \bar{\Sigma} \underline{S^{\top}} r_{+}^{\top},  \tag{8.5}\\
& 0=\left(A-\hat{S} Q \bar{\Sigma}+\left[\gamma^{-2} V_{1_{\infty}}+D\right] P\right)^{\top} \hat{\mathrm{P}}+\hat{\mathrm{P}}(\mathrm{~A}-\hat{\mathrm{S}} \mathrm{Q} \overline{\mathrm{~L}} \\
& \left.+\left[\gamma^{-2} V_{1_{\infty}}+D\right] P\right)+\dot{P}\left(\gamma^{-2}\left[V_{1_{\infty}}+\beta^{2} \hat{S} Q \bar{\Sigma} Q \hat{S}^{\top}\right]+D\right) \hat{P} \\
& +P_{a}^{\top} R_{2}^{1} P_{a} P_{d} \tau_{d}^{\top} P_{a}^{\top} R_{2}^{1} P_{a} \tau_{1},  \tag{8.6}\\
& 0=\left(A-B R_{2}{ }_{2}^{1} P_{a}+\left[\gamma^{-2} V_{1_{\infty}}+D\right] P\right) \hat{\mathbf{Q}}+\hat{\mathbf{Q}}\left(A-B R_{2}^{1} \mathrm{P}_{\mathrm{a}}\right. \\
& \left.+\left[\gamma^{-2} V_{1_{\infty}}+D\right] P\right)+\hat{S} Q \bar{\Sigma} \underline{Q} \hat{S}^{\top}-\tau_{\perp} \mathrm{S} Q \bar{Q} Q \hat{S}^{\top} \tau_{\perp}^{\top}, \tag{8.7}
\end{align*}
$$

and let $P$ be given by ( 6.16 ) and ( $A_{c}, B_{c}, C_{c}$ ) by

$$
\begin{align*}
& A_{c}=\Gamma\left(A-\hat{S} Q \bar{\Sigma}-B R_{2}^{-1} P_{a}+\left[\gamma^{-2} V_{1 \infty}+D\right] P\right) G^{\top},  \tag{8.8}\\
& B_{c}=\Gamma \hat{S} Q C^{\top} V_{2}^{1},  \tag{8.9}\\
& C_{c}=-R_{2}^{-1} P_{a} G^{\top} . \tag{8.10}
\end{align*}
$$

Then $(\tilde{E}, \tilde{A}+\Delta \tilde{A})$ is detectable if and only if $\tilde{A}+\Delta \tilde{A}$ is asymptotically stable for all $(\Delta A, \Delta B) \in U$. In this case, the closed-loop transfer function $\hat{H}_{\Delta \hat{A}}(s)$ satisfies the $H_{l o \infty}$ disturbance attenuation constraint

$$
\begin{equation*}
\left\|\dot{H}_{\Delta \tilde{A}^{(s)}}\right\|_{\infty} \leq \gamma, \quad(\Delta A, \Delta B) \epsilon U, \tag{8.11}
\end{equation*}
$$

and the worst case $L_{2}$ performance criterion (6.1) satisfies the bound

$$
\begin{equation*}
J\left(A_{c}, B_{c}, C_{c}\right) \leq \operatorname{tr}\left[(P+\hat{P}) V_{1}+\hat{P} \hat{S} Q \bar{Q} Q \hat{S}^{\top} \hat{P}\right] \tag{8.12}
\end{equation*}
$$

## VIII. Further Extensions

The results of this paper can be readily extended in several directions:

1) Mixed bounds, i.e., letting $\Delta A=\Delta A_{1}+\Delta A_{2}$ and bounding $\Delta A_{1}$ with the linear bound and $\Delta A_{2}$ with the quadratic bound (this would unify the linear and quadratic bound results)
2) $L_{2}$ and $H_{\infty}$ cross weighting terms (e.q.. $x^{\top} R_{12} u$ ) as well as correlated plant disturbance and sensor noise
3) nonstrictly proper plant model, i.e., (2.2) replaced by
$y(t)=(C+\Delta C) x(t)+(D+\Delta D) u(t)+D_{2} w(t)$
4) nonstrictly proper controller, i.e., (2.4) replaced by

$$
\begin{equation*}
u(t)=C_{c} X_{c}(t)+D_{c} y(t) \tag{2.4}
\end{equation*}
$$

and the related problems of singular control weighting ( $R_{2} \geq 0$ ) and singular measurement noise ( $V_{2} \geq 0$ )
5) discrete-time and sampled data design.

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## Optimal Output Feedback for Nonzero Set Point Regulation

## DENNIS S. BERNSTEIN AND WASSIM M. HADDAD

Abstract-Motivated by the results of Artstein and Leizarowitz [2] on steady-state periodic tracking, amtinuous-time monzero sef point regulation problem is considered which involves 1) moisy and aonnoisy measurements, 2) weighted and unwelghted controls, 3) correlated plant/ measurement molse and cross weighting, 4) nonzero-menn disturbances, and 5) state-, control-, and measurement-dependent white moise. It is shown that in the sbsence of multiplicatlve disturbances the closed-loop control can be designed independently of the opentoop controf. Ualike [2], the results are obtained without using the overtaking criterion.

## I. Introduction

The quadratic performance criterion

$$
\begin{equation*}
J \leqq \int_{0}^{t_{3}} x^{T}(t) Q x(t)+u^{T}(t) R u(t) d t \tag{1.1}
\end{equation*}
$$

expresses the desire to minimize deviations of the state $x(t)$ of the system

$$
\begin{equation*}
\dot{x}(t)=A x(t)+B u(t) \tag{1.2}
\end{equation*}
$$

from the regulation point $x=0$. As is well known [1, pp. 270-276], the nonzero sel point criterion

$$
\begin{equation*}
J_{x}=\int_{0}^{t_{1}}\left[x(t)-x^{T}\right]^{T} Q[x(t)-x]+u^{T}(t) R u(t) d t \tag{1.3}
\end{equation*}
$$

[^31]presents no additional difficulty as long as $x(t)$ and $u(t)$ are replaced by $x(t)-\bar{x}$ and $u(t)-\tilde{u}$, where $u$ satisfies
\[

$$
\begin{equation*}
0=A x+B u \tag{1.4}
\end{equation*}
$$

\]

Closer inspection, however, reveals that this approach is suboptimal. Specifically, the offset $u$ in the control may correspond to an unacceptably high level of control effort when $u^{\top} R u$ is large. Hence, this approach overlooks design tradeoffs concerning the control effort required for maintaining the nonzero regulation point $\bar{x}$. Moreover, such an approach is impossible when 4 satisfying (1.4) does not exist.

A significant advance in extending the LQR formulation to steady-state tracking problems (and, hence, to nonzero set point regulation) was given by Artstein and Leizarowitz in [2]. They consider the performance criterion

$$
\begin{equation*}
J_{\infty} \leq \int_{0}^{\infty}[x(t)-\Gamma(t)]^{r} Q(x(t)-\Gamma(t)]+u^{r}(t) R u(t) d t \tag{1.5}
\end{equation*}
$$

where $\Gamma(\cdot)$ is periodic on $[0, \infty)$ and the minimization of $J_{\infty}$ is performed in the sense of the overtaking criterion. For the nonzero set point problem ( $\Gamma(t)=\bar{x}$ ) with full-state feedback plus constant offset control law

$$
\begin{equation*}
u(t)=K x(t)+\alpha \tag{1.6}
\end{equation*}
$$

it follows from [2, Theorem 2] that $K$ and $\alpha$ are given by

$$
\begin{gather*}
K=-R^{-1} B^{T} P  \tag{1.7}\\
\alpha=-R^{-1} B^{T}(A-\Sigma P)^{-T} Q x \tag{1.8}
\end{gather*}
$$

where $P$ satisfies the Riccati equation

$$
\begin{equation*}
0=A^{r} P+P A+Q-P \Sigma P \tag{1.9}
\end{equation*}
$$

with $\leq \hat{\text { 全 }} B R^{-1} B^{T}$.
Two features of the control law (1.6)-(1.8) are noteworthy. First, (1.6) consists of both a closed-loop feedback component $K x(t)$ and an openloop component $\alpha$ depending upon the regulation point (Fig. 1). Second (and more important), is the observation that the closed-loop control component is independent of the open-loop component. From a practical point of view this feature is quite useful since it implies that the feedback gain $K$ can be determined without regard to the set point. Hence, a change in the desired set point $\vec{x}$ during on-line operation does not necessitate resolving the Riccati equation in real time; only $\alpha$ requires updating. For a new value of $\dot{x}^{*}, \alpha$ can readily be recomputed on-line via the matrix multiplication operation (1.8).

The contribution of the present note is an extension of the result of [2] as applied to the nonzero set point regulation problem without using the overtaking criterion. We extend this result in the following different ways.

1) Outpur Feedback with Noisy and Nonnoisy Measurements: To obrain a more realistic problem setting, we consider the case in which the full state is not available, but rather only measured linear combinations of states. Moreover, we consider the possibility that some of the measurements are comupted by white noise while others are noise free. Note that the noise-free case was considered in [3] while the fully noisy case is the standard assumption in LQG theory As in [4]-[6] we express the solution in terms of a projection corresponding to the noise-free measurements.
2) Singular Control Weighting: As noted in [6], [7] static continu-ous-time feedback of noise-corrupted measurements results in unbounded cost uniess the corresponding controls are unweighted. Hence, we allow for both weighted and unweighted controls to which the noise-free and noisy measurements are fed, respectively. This setting leads to an additional projection dual to the projection arising from the noise-free measurements [6].
3) Correlated Plant and Measurement Noise and Crass Weight. ing: To allow greater design flexibility we allow the possibility that the plant and measurement noise are correlated. In addition, we consider the dual design feature, namely, cross weighting in the performance criterion.
4) Nonzero-Mean Disturbances: In addition to the presence of zero-


Fig. 1 ,
mean white plant disturbances we allow for the possibility of a nonzero constant disturbance offset. In contrast to [1, pp. 277-281], our result shows that the presence of a constant disturbance offset leads to an additional offset in the open-loop component of the control.
5) Multiplicative White Noise: In addition to the above generalizations we allow for the presence of multiplicative disturbances in the plant. The control law thus generalizes previous results involving state-, controland measurement-dependent noise [8]-[11]. As shown in [12]-[14], the multiplicative white noise model can be used to guarantee robustness with respect to deterministic plant parameter variations.

## II. Notation and Definitions

| $\mathfrak{J}, \mathbb{B}^{r \times s}, \mathbb{R}^{r}$ 它 | real numbers, $r \times s$ real matrices, $\sqrt[3]{ } 3^{\prime \times 1}$ expectation |
| :---: | :---: |
| $I_{r},()^{T}$ | $r \times r$ identity, transpose |
| ( ${ }^{(8)}$ | Kronecker sum. Kronecker product |
| asymptotically stable matrix | matrix with eigenvalues in open lefthalf plane |
| $n, m_{1}, m_{2}, l_{1}, l_{2}, p, r$ | positive integers |
| $x, u_{1}, u_{2}, y_{1}, y_{2}$ | $n, m_{1}, m_{2}, J_{1}, l_{2}$-dimensional vectors |
| $A, A_{i} ; B_{1}, B_{1 i} ; C_{1}, C_{1 i}$ | $n \times n$ matrices; $n \times m_{1}$ matrices; $l_{1} \times$ $n$ matrices, $i=1, \cdots, p$ |
| $B_{2}, C_{2}, K_{1}, K_{2}, L$ | $n \times m_{2}, l_{2} \times n, m_{1} \times l_{2}, m_{2} \times l_{1}, r \times$ $n$ matrices |
| $\alpha_{1}, \alpha_{2}, \gamma, \delta$ | $m_{1}, m_{2,} n, r$-dimensional vectors |
| $v_{i}(t)$ | unit variance white noise, $i=1, \cdots$, $p$ |
| $w_{0}(t), w_{1}(t)$ | $n$-dimensional, $l_{1}$-dimensional white noise |
| $V_{0}, V_{1}$ | intensities of $w_{0}, w_{1} ; V_{0} \geq 0, V_{1}>0$ |
| V01 | $n \times l_{1}$ cross intensity of $w_{0}, w_{1}$ |
| $R_{0}, R_{1}$ | $r \times r$ and $m_{1} \times m_{1}$ state and control weightings: $R_{0} \geq 0, R_{1}>0$ |
| $\boldsymbol{R}_{01}$ | $r \times m_{1}$ cross weighting; |
|  | $R_{0}-R_{01} R_{1}^{-1} R_{01}^{T} \geq 0$ |
| $\overline{\mathcal{A}}, \bar{A}_{1}$ | $A+B_{1} K_{1} C_{2}+B_{1} K_{2} C_{1}, A_{1}+B_{11} K_{1} C_{2}$ |
|  | $+B_{2} K_{2} C_{1}, i=1, \cdots, p$ |
| $\boldsymbol{B}, \boldsymbol{B}_{1}$ | $B_{1} \alpha_{1}+B_{2} \alpha_{2}+\gamma, B_{1}, \alpha_{1}, i=1, \cdots, p$ |
| $\dot{W}(t)$ | $w_{0}(t)+B_{2} K_{2} w_{1}(t)$ |
| $\bar{V}$ | $V_{0}+V_{01} K_{2}^{T} B_{2}^{T}+B_{2} K_{2} V_{01}^{T}+$ |
|  | $B_{2} K_{2} V_{1} K_{2}^{\top} B_{2}^{T}$ |
| $\boldsymbol{R}$ | $L^{\top} \boldsymbol{R}_{0} L+L^{\top}{ }^{\top} \boldsymbol{R}_{01} K_{1} C_{2}+$ |
|  | $C_{2}^{T} K_{1}^{T} R_{01}^{T} L+C_{2}^{T} K_{1}^{T} R_{1} K_{1} C_{2}$ |

For arbitrary $n \times n Q, P$ such that the indicated inverses exist, define:

$$
Q \Leftrightarrow Q C_{1}^{T}+V_{01}+\sum_{1=1}^{p} A_{1}\left(Q+A_{1} A_{s}^{T}\right) C_{11}^{r}
$$

$$
\sigma_{s} \triangleq B_{1}^{T} P+R_{01}^{T} L+\sum_{i=1}^{f} B_{11}^{T} P A_{1}
$$

$$
\begin{aligned}
& \tau_{1} \triangleq Q C_{2}^{T}\left(C_{2} Q C_{2}^{r}\right)^{-1} C_{2}, \quad \tau_{2} \triangleq B_{2}\left(B_{2}^{\top} P B_{2}\right)^{-1} B_{2}^{T} P, \\
& T_{1,} \triangleq I_{n}-T_{1}, T_{2 \perp} \triangleq I_{n}-T_{2}, \\
& V_{11} \triangleq V_{1}+\sum_{i=1}^{p} C_{1 i}\left(Q+m_{1} m_{3}^{T}\right) C_{1,}^{T}, \quad R_{11} \hat{10} R_{1}+\sum_{i=1}^{p} B_{11}^{T} P B_{1,},
\end{aligned}
$$

$A_{s} \triangleq A-B_{1} R_{s}^{-1} O_{s} r_{1}-r_{2} Q V_{1 s}^{-1} C_{1}, \quad m_{s} \triangleq-A_{s}^{-1} B$,

$\varphi_{i} \triangleq R_{1,}-\left(\otimes_{s} T_{1} A_{s}^{-1} B_{1}+B_{1}^{T} A_{z}^{-T} T_{1,}^{T} \otimes_{s}^{T}\right)$,
$\Delta_{1 s} \triangleq P_{s} T_{1} A_{i}^{-1}+B_{1}^{T} \bar{A}_{j}^{-r} P_{i}, \quad \Delta_{2 s} \triangleq B_{2}^{T}\left(P \bar{A}_{z}^{-1}+\bar{A}_{j}^{-T} P\right)$,

$$
a_{1} \triangleq \Delta_{1}, B_{2}, \quad a_{25} \triangleq \Delta_{2 .} B_{21}
$$

$\Lambda_{s} \triangleq \hat{A}_{s}^{-\Gamma}\left(L^{\top} R_{0}-T_{1}^{\top} Q_{5}^{\top} R_{1 s}^{-1} R_{01}^{T}\right)$,
$\Lambda_{1 s} \triangleq B_{1}^{T} \Lambda_{s}-R_{01}^{T}, \quad \Lambda_{2 \mathrm{~s}} \triangleq B_{2}^{T} \Lambda_{1}$,
$Q \triangleq Q C_{1}^{T}+V_{01}, \quad Q \triangleq B_{1}^{T} P+R_{01}^{T}$,
$A \triangleq A-B_{1} R_{1}^{-1} Q r_{1}-\tau_{2} Q V_{1}^{-1} C_{1}, \quad i n \triangleq-A^{-1} B$,

$$
\phi \triangleq R_{1}-\left(F T_{1}, \hat{A}^{-1} B_{1}+B_{1}^{r} \hat{A}^{-T} T_{1}^{r} Q^{r}\right),
$$

$\Delta_{1} \triangleq Q_{r_{1}} \hat{A}^{-1}+B_{1}^{r} \hat{A}^{-r} P, \quad \Delta_{2} \triangleq B_{2}^{r}\left(P \hat{A}^{-1}+\hat{A}^{-} T^{P}\right)$.

$$
\Omega_{1} \triangleq \Delta_{1} B_{2}, \quad \Omega_{2} \triangleq \Delta_{2} B_{2}
$$

$\Delta \triangleq A^{-T}\left(L^{\top} R_{0}-T_{1} \ominus^{T} R_{1}^{-1} R_{01}^{T}\right)$,
$\Lambda_{1} \triangleq B_{1}^{T} \Lambda-R_{01}^{r}, \quad \Lambda_{2} \triangleq B_{2}^{T} \Lambda$.
 $\Delta_{10} \triangleq P r_{1,} A_{0}^{-1}+B_{1}^{r} \hat{A}_{0}^{-r} P$,

$$
\begin{aligned}
& \Lambda_{10} \triangleq B_{1}^{r} A_{0}^{-T}\left(L^{T} R_{0}-T_{1} \otimes^{T} R_{1}^{-1} R_{01}^{T}\right)-R_{01}^{T}, \\
& A_{1}=A-B_{1} R_{1}^{-1} P, \quad \sum \triangleq B_{1} R_{1}^{-1} B, \quad A_{P}^{T} \triangleq A-\Sigma P .
\end{aligned}
$$

## ili. nonzero Set pont regulation

Nonzero Set Point Problem: Given the controlled system
$x(t)=\left(A+\sum_{i=1}^{n} v_{1}(t) A_{i}\right) x(t)$

$$
\begin{equation*}
+\left(B_{1}+\sum_{i=1}^{p} u_{1}(t) B_{1}\right) u_{1}(t)+B_{2} u_{2}(t)+w_{0}(t)+\gamma \tag{3.1}
\end{equation*}
$$

with measurements

$$
\begin{gather*}
y_{1}(t)=\left(C_{1}+\sum_{i=1}^{f} v_{1}(t) C_{11}\right) x(t)+w_{1}(t)  \tag{3.2}\\
y_{2}(t)=C_{2} x(t) \tag{3.3}
\end{gather*}
$$

where $t \in[0, \infty)$, determine $K_{1}, K_{2}, \alpha_{1}$, and $\alpha_{2}$ such that the static output feedback law

$$
\begin{align*}
& u_{1}(t)=K_{1} y_{2}(t)+\alpha_{1}  \tag{3.4}\\
& u_{2}(t)=K_{2} y_{1}(t)+\alpha_{2} \tag{3.5}
\end{align*}
$$

minimizes the performance criterion

$$
\begin{align*}
J\left(K_{1}, K_{2}, \alpha_{1}, \alpha_{2}\right)=\lim _{t \rightarrow \infty} & \mathbb{T}\left[(L x(t)-\delta)^{r} R_{0}(L x(t)-\delta)\right. \\
& \left.+2(L x(t)-\delta)^{r} R_{01} u_{1}(t)+u_{1}^{r}(t) R_{1} u_{1}(t)\right] \tag{3.6}
\end{align*}
$$

The closed-loop system (3.1)-(3.5) can be writien as

$$
\begin{equation*}
x(t)=\left(A+\sum_{i=1} v_{1}(t) A_{i}\right) x(t)+B+\sum_{i=1}^{L} v_{i}(t) B_{i}+\hat{w}(t) . \tag{3.7}
\end{equation*}
$$

To analyze (3.7) define the second-moment and covariance matrices

$$
Q(t) \triangleq \mathbb{E}\left[x(t) x^{\top}(t)\right], \quad Q(t) \triangleq Q(t)-m(t) m^{T}(t)
$$

where $m(t) \triangleq \mathbb{E}[x(t)]$. It follows from [15, p. 142], that $Q(t), Q(t)$, and $m(t)$ satisfy

$$
\begin{aligned}
\dot{Q}(t)=A Q(t) & +Q(t) A^{T}+Q_{+}+B_{m}^{T}(t)+m(t) B^{r} \\
& +\sum_{i=1}^{n}\left[A_{i} Q(t) A_{i}^{T}+A_{i} m(t) B_{i}^{T}+B_{i} m^{T}(t) A_{i}^{T}+B_{i} B_{i}^{T}\right],
\end{aligned}
$$

$$
Q(t)=A Q(t)+Q(t) A^{r}+\bar{V}+\sum_{i=1}^{n}\left[A, Q(t) A_{i}^{T}+A_{i} m(t) m^{T}(t) A_{i}^{T}\right.
$$

$$
\left.+\bar{A}_{1} m(t) B_{1}^{T}+\bar{B}_{1} \ldots T^{T}(t) A_{i}^{T}+\bar{B}_{1} B_{1}^{T}\right],
$$

$$
\begin{equation*}
\dot{m}(t)=\bar{A} m(t)+\bar{B} . \tag{3.10}
\end{equation*}
$$

To consider the steady state, we restrict our consideration to the set of second-moment stabilizing gains

$$
\delta_{s} \triangleq\left\{\left(K_{1}, K_{2}\right): A \oplus A+\sum_{i=1}^{p} A_{i} \otimes A_{i} \text { is asymptotically stable }\right\} .
$$

It follows from fundamental properties of Lyapunov equations that if ( $\boldsymbol{K}_{1}$, $\left.K_{2}\right) \in \delta_{5}$, then $\bar{A}$ is also asymprotically stable. Hence, for $\left(K_{1}, K_{2}\right) \in \delta_{s}$, $Q \triangleq \lim _{t \rightarrow \infty} Q(t), Q \triangleq \lim _{t \rightarrow-Q(t)}$ and $m \triangleq \lim _{t \rightarrow \infty} m(t)$ exist and satisfy
$0=A Q+Q A^{T}+B m+m B^{T}+\nabla+\sum_{i=1}^{p}\left(A_{i} Q A_{i}^{T}\right.$

$$
\begin{equation*}
\left.+A_{i} m B_{i}^{T}+B_{i} m^{T} A_{i}^{T}+B_{i} B_{i}^{T}\right] . \tag{3.11}
\end{equation*}
$$

$0=\bar{A} Q+Q \tilde{A}+\bar{Q}+\sum_{i=1}^{D}\left[A_{i} Q \tilde{A}_{i}^{T}+A_{i} m m^{T} A_{i}^{T}\right.$
$\left.+A_{i} m B_{1}^{T}+B_{i} m^{T} A_{1}^{T}+B_{i} B_{1}^{T}\right]$,
$0=\bar{A} m+B$.
Now $J\left(K_{1}, K_{2}, \alpha_{1}, \alpha_{2}\right)$ is given by

$$
\begin{align*}
J\left(K_{1}, K_{2}, \alpha_{1}, \alpha_{i}\right)= & \llbracket r\left[\left(Q+m m^{\top}\right) R_{j}-2 m^{\top} L^{\top} R_{0} \delta\right. \\
& +\delta^{\top} R_{0} \delta+2 m^{\top} L^{\top} R_{01} \alpha_{1} \\
& -2 \delta^{\top} R_{01} K_{1} C_{2} m-2 \delta^{\top} R_{01} \alpha_{1} \\
& +2 m^{\top} C_{2}^{\top} K_{1}^{\top} R_{1} \alpha_{1}+\alpha_{1}^{\top} R_{1} \alpha_{1} . \tag{3.14}
\end{align*}
$$

Associated with $Q$ is its dual $P \geq 0$ which is the unique solution of

$$
\begin{equation*}
0=\bar{A}^{T} P+P \bar{A}+\bar{R}+\sum_{i=1}^{p} \bar{A}_{i}^{T} P \bar{A}_{i} \tag{3.15}
\end{equation*}
$$

To obtain closed-form expressions for the feedback gains we further restrict consideration to the set
$\delta_{1}^{+} \triangleq\left\{\left(K_{1}, K_{2}\right) \in \delta_{1}: C_{2} Q C_{2}^{T}, B_{2}^{\top} P B_{2}, \Phi_{1}\right.$
and $\Omega_{1, ~}^{T} \Phi_{5}^{-1} \square_{1,}+\Omega_{2}$ are inverible),
and assume

$$
\begin{equation*}
\left[B_{11} \neq 0=C_{11}=0\right], \quad i=i, \cdots, p \tag{3.16}
\end{equation*}
$$

Optimizing (3.14) subject to (3.12) and (3.13) yields the following reauk illustrated in Fig. 2.

Theorem 3.1: Suppose $K_{1}, K_{2}, \alpha_{1}, \alpha_{2}$ wolve the nonzero point


Fig. 2.
problem with $\left(K_{1}, K_{2}\right) \in \delta_{s^{+}}$Then there exist $n \times n Q, P \geq 0$ such that

$$
\begin{gather*}
K_{1}=-R_{1 s}^{-1} Q_{s} Q C_{2}^{T}\left(C_{2} Q C_{2}^{T}\right)^{-1},  \tag{3.17}\\
K_{2}=-\left(B_{2}^{T} P B_{2}\right)^{-1} B_{2}^{T} P Q_{s} V_{1 s}^{-1}  \tag{3.18}\\
\alpha_{1}=\Phi_{s}^{-1}\left(\Omega_{1 s} \alpha_{2}+\Delta_{1 s} \gamma-\Lambda_{1 s} \delta\right),  \tag{3.19}\\
\alpha_{2}=\left(\Omega_{1 s}^{T} \Phi_{s}^{-1} \Omega_{1 s}+\Omega_{2 s}\right)^{-1}\left[\left(\Omega_{1 s}^{T} \Phi_{s}^{-1} \Delta_{1 s}+\Delta_{2 s}\right) \gamma-\left(\Omega_{1 s}^{T} \Phi_{s}^{-1} \Lambda_{1 s}+\Lambda_{2 s}\right) \delta\right], \tag{3.20}
\end{gather*}
$$

and such that $Q$ and $P$ satisfy

$$
\begin{align*}
& 0=\left(A-B_{1} R_{1 ;}^{-1} \nabla_{s} \tau_{1}\right) Q+Q\left(A-B_{1} R_{1:}^{-1} \nabla_{s} \tau_{1}\right)^{T}+V_{0} \\
& +\sum_{i=1}^{p}\left[\left(A_{i}-B_{1 i} R_{1 ;}^{-1} \mathcal{O}_{3} \tau_{1}\right) Q\left(A_{i}-B_{1,} R_{1 ;}^{-1} \mathcal{O}_{s} T_{1}\right)^{T}\right. \\
& \left.+\hat{A}_{i s} \hat{H}_{s} \tilde{f}_{s}^{T} \hat{A}_{i s}^{T}+\hat{A}_{1 s} f_{s} \bar{B}_{i}^{T}+\bar{B}_{i} \hat{m}_{s}^{T} \hat{A}_{i s}^{T}+\hat{B}_{s} B_{i}^{T}\right] \\
& -Q, V_{1 s}^{-1} Q_{s}^{T}+r_{2,} Q_{s} V_{15}^{-1} Q_{s}^{r} r_{2 \perp}^{T} \text {, }  \tag{3.21}\\
& 0=\left(A-\tau_{2} Q_{3} V_{1:}^{-1} C_{1}\right)^{\top} P+P\left(A-r_{2} Q_{1} V_{15}^{-1} C_{1}\right)+R_{0} \\
& +\sum_{i=1}^{p}\left(A_{i}-\tau_{2} Q V_{15}^{-1} C_{11}\right)^{T} P\left(A_{i}-\tau_{2} Q, V_{15}^{-1} C_{1 j}\right) \\
& -\nabla_{s}^{T} R_{1 s}^{-1} \nabla_{s}+\tau_{1+}^{T} Q_{s}^{T} R_{1 s}^{-1} \nabla_{s} \tau_{11} . \tag{3.22}
\end{align*}
$$

Outline of Proof: As in [16] the result is obtained by forming the Lagrangian while accounting for (3.12) and (3.13). Define
$\mathcal{L}\left(K_{1}, K_{2}, \alpha_{1}, \alpha_{2}, Q, m\right)=\operatorname{tr}\left[\lambda_{0} J\left(K_{1}, K_{2}, \alpha_{1}, \alpha_{2}\right)\right.$ $+P[$ RHS of $\left.(3.12)]+\lambda^{T}(A m+B)\right]$
where $\lambda_{0} \geq 0$ and $\lambda \in \mathbb{R}^{n}$. Setting $\partial \mathcal{L} / \partial Q=0$ and using the secondmoment stability assumption it follows that $\lambda_{0}=1$ without loss of generality. The derivation now follows by setting the partial derivatives of $\mathcal{L}$ with respect to $K_{1}, K_{2}, \alpha_{1}, \alpha_{2}$, and $m$ to zero and solving for the gains. To assist the reader in carrying out the details we note that $\lambda$ is given by

$$
\begin{aligned}
& \lambda=-2 A^{-T}\left[-P B-\sum_{i=1}^{P} A_{i}^{T} P B_{i}+L^{T} R_{0} \delta\right. \\
&\left.-L^{\top} R_{01} \alpha_{1}+C_{2}^{T} K_{1} R_{01}^{T} \delta-C_{2}^{T} K_{1}^{T} R_{1} \alpha_{1}\right]
\end{aligned}
$$

Remark 3.I: Because of the presence of $\delta$ in (3.21) via $\alpha_{1}$ and $\alpha_{2}$ in $B$ (ese the definition of $\hat{m}_{s}$ ) and $\alpha_{1}$ in $B_{i}$, the closed-loop component of the control law (3.17)-(3.20) cannor be designed independently of the openloop component. As now shown, independence is recovered when the multiplicative noise terms are absent.

## IV. Speciallzations of Theorem 3.1

To draw connections with the previous literature, a series of specializations of Theorem 3.1 is now given. We begin by deleting all multiplicative white noise terms, i.e.,

$$
\begin{equation*}
A_{i}, B_{11}, C_{11}=0, \quad i=1, \cdots, p \tag{4.1}
\end{equation*}
$$

In this case the stabilizing set $\delta_{s}$ can be characterized by

$$
\delta=\left\{\left(K_{1}, K_{2}\right): A \text { is asymptotically stable }\right\}
$$

and, furthermore, $\delta_{s}{ }^{+}$becomes

$$
\delta^{+} \triangleq\left\{\left(K_{1}, K_{2}\right) \in \delta: C_{2} Q C_{2}^{T}, B_{2}^{T} P B_{2},\right.
$$

$\Phi$ and $\Omega_{1}^{T} \Phi^{-1} \Omega_{1}+\Omega_{2}$ are invertible $\}$.
Corollary 4.I: Assume (4.1) is satisfied and suppose $K_{1}, K_{2}, \alpha_{1}, \alpha_{2}$ solve the nonzero set point problem with $\left(K_{1}, K_{2}\right) \in \mathcal{S}^{+}$. Then there exist $n \times n Q, P \geq 0$ such that

$$
\begin{gather*}
K_{1}=-R_{1}^{-1} \mathscr{P} Q C_{2}^{T}\left(C_{2} Q C_{2}^{T}\right)^{-1},  \tag{4.2}\\
K_{2}=-\left(B_{2}^{T} P B_{2}\right)^{-1} B_{2}^{T} P Q V_{1}^{-1},  \tag{4.3}\\
\alpha_{1}=\Phi^{-1}\left(\Omega_{1} \alpha_{2}+\Delta_{1} \gamma-\Lambda_{1} \delta\right) \tag{4.4}
\end{gather*}
$$

$$
\begin{equation*}
\alpha_{2}=\left(\Omega_{1}^{\tau} \Phi^{-1} \Omega_{1}+\Omega_{2}\right)^{-1}\left[\left(\Omega_{1}^{\tau} \Phi^{-1} \Delta_{1}+\Delta_{2}\right) \gamma-\left(\Omega_{1}^{\tau} \Phi^{-1} \Lambda_{1}+\Lambda_{2}\right) \delta\right] \tag{4.5}
\end{equation*}
$$

and such that $Q$ and $P$ satisfy

$$
\begin{align*}
& 0=\left(A-B_{1} R_{1}^{-1} Q \tau_{1}\right) Q+Q\left(A-B_{1} R_{1}^{-1} Q \tau_{1}\right)^{T}+V_{0} \\
&-Q V_{1}^{-1} Q^{T}+\tau_{2 A} Q V_{1}^{-1} Q^{T} \tau_{2 A}^{T},  \tag{4.6}\\
& 0=\left(A-T_{2} Q V_{1}^{-1} C_{1}\right)^{T} P+P\left(A-\tau_{2} Q V_{1}^{-1} C_{1}\right)+R_{0} \\
&-Q^{T} R_{1}^{-1} P+\tau_{1}^{T} Q^{T} R_{1}^{-1} Q_{\tau_{1}} . \tag{4.7}
\end{align*}
$$

In addition to (4.1) we assume further that all controls are weighted and all measurements are noise free, i.e.,

$$
\begin{equation*}
B_{2}=0, \quad C_{1}=0 \tag{4.8}
\end{equation*}
$$

This corresponds to the setting considered in [3]. It follows from the proof of Theorem 3.1 that the assumption $B_{2}=0$ leads to $T_{2}=0$, and $C_{1}=0$ corresponds to deleting (3.5). Hence $\delta$ and $\mathcal{S}^{+}$are now given by

$$
\begin{aligned}
& \mathcal{S}_{0} \triangleq\left\{K_{1}: A+B_{1} K_{1} C_{2} \text { is asymptotically stable }\right\} \\
& \delta_{0}^{+} \triangleq\left\{K_{1} \in \delta_{0}: C_{2} Q C_{2}^{T} \text { and } \Phi_{0} \text { are invertible }\right\}
\end{aligned}
$$

Corollary 4.2: Assume (4.1) and (4.8) are satisfied and suppose $K_{1}$ and $\alpha_{1}$ solve the nonzero set point problem with $K_{1} \in \delta_{0}$. Then there exist $n \times n Q, P \geq 0$ such that

$$
\begin{gather*}
K_{1}=-R_{1}^{-1} \odot Q C_{2}^{T}\left(C_{2} Q C_{2}^{T}\right)^{-1},  \tag{4.9}\\
\alpha_{1}=\Phi_{0}^{-1}\left(\Delta_{10} \gamma-\Lambda_{10} \delta\right), \tag{4.10}
\end{gather*}
$$

and such that $Q$ and $P$ satisfy

$$
\begin{align*}
& 0=\left(A-B_{1} R_{1}^{-1} P T_{1}\right) Q+Q\left(A-B_{1} R_{1}^{-1} P r_{1}\right)^{T}+V_{0}  \tag{4.11}\\
& 0=A^{T} P+P A+R_{0}-Q^{T} R_{1}^{-1} Q+r_{1}^{T} \otimes^{T} R_{1}^{-1} \odot r_{1+} \tag{4.12}
\end{align*}
$$

We now specialize further to the full-state feedback case, i.e.,

$$
\begin{equation*}
C_{2}=I_{n} \tag{4.13}
\end{equation*}
$$

and hence $T_{1}=T_{n}$ and $\tau_{1+}=0$. Now $\delta_{0}$ and $\delta_{0}^{+}$become
$\delta_{1} \triangleq\left\{K_{1}: A+B_{1} K_{1}\right.$ is asymptotically stable $\}$.
$\delta_{1} \Delta\left\{X_{1} \in \delta_{1}: Q\right.$ is invertible $\}$.

Corollary 4.3: Assume (4.1), (4.8), and (4.13) are satisfied and suppose $K_{1}$ and $\alpha_{1}$ solve the nonzero set point problem with $K_{1} \in \delta_{1}^{+}$. Then there exist $n \times n Q, P \geq 0$ such that

$$
\begin{equation*}
K_{1}=-R_{1}^{-1} P, \tag{4.14}
\end{equation*}
$$

$$
\begin{equation*}
\alpha_{1}=R_{1}^{-1} B_{1} A_{1}^{-T} P_{\gamma}-R_{1}^{-1}\left[B_{1}^{T} A_{1}^{-T}\left(L^{\top} R_{0}-\phi^{T} R_{1}^{-1} R_{01}^{T}\right)-R_{01}^{T}\right] \delta \tag{4.15}
\end{equation*}
$$

and such that $Q$ and $P$ satisfy

$$
\begin{gather*}
0=\left(A-B_{1} R_{1}^{-1} \mathcal{P}\right) Q+Q\left(A-B_{1} R_{1}^{-1} \mathcal{P}\right)^{T}+V_{0}  \tag{4.16}\\
0=A^{T} P+P A+R_{0}-\mathcal{P}^{\top} R_{1} \mathcal{P} . \tag{4.17}
\end{gather*}
$$

Finally, setting

$$
\begin{equation*}
\gamma=0, \quad R_{01}=0, \quad L=I_{n} \tag{4.18}
\end{equation*}
$$

we obtain the result of [2].
Corollary 4.4: Assume (4.1), (4.8), (4.13), and (4.18) are satisfied and suppose $K_{1}$ and $\alpha_{1}$ solve the nonzero set point problem with $K_{1} \in \delta_{1}$. Then there exists $n \times n P \geq 0$ such that

$$
\begin{gather*}
K_{1}=-R_{1}^{-1} B_{1}^{\top} P,  \tag{4.19}\\
\alpha_{1}=-R_{1}^{-1} B_{1}^{T} A_{j}^{\top} R_{0} \delta \tag{4.20}
\end{gather*}
$$

and such that $P$ satisfies

$$
\begin{equation*}
0=A^{\top} P+P A+R_{0}-P \Sigma P \tag{4.21}
\end{equation*}
$$

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## Analysis of Time-Varying Scaled Systems Via General Orthogonal Polynomials

## TSU TLAN LEE AND YIH FONG CHANG

## Abstract-Gemeral orthosonal polynomials are introduced to amalyze

 and approximate the solution of a class of scaled systems. Uaiag the operational matrix of integration, together with the operntional mintix of Herar transformation, the dynamical equation of a scaled system is reduced to a set of simultaneous linear algebraic equations. The coeflicient vectors of the general orthoganal polynomials can be determined recursively by the derived algorithm. An iliustritive example is given to demonstrate the validity and applicability of the orthogonal polynomial approximations.
## I. INTRODUCTION

An investigation of the dynamics of an overhead current collection mechanism for an electric locomotive by Ockendon and Taylor [12] revealed that under certain conditions, the dynamics of the systems is characterized by a differential equation containing terms with a scaled argument of the form

$$
\begin{gathered}
\dot{X}(t)=A X(\lambda r)+B X(t) \\
X(0)=X_{0}
\end{gathered}
$$

where $X(\lambda t)$ and $X(t)$ are $n$-vectors and $A$ and $B$ are $n \times n$ matrices and the constant $0<\lambda<1$. This type of differential equation also plays an important role in several chemical processes [3], [13]. This equation was first saudied by Fox et al. [11] with the introduction of a finite difference method for $0<\lambda<1$. Recently, the solution of such a scaled system has been obtained by several different orthogonal functions, such as blockpulse functions [14], [2], [3], Walsh functions [1], delayed unit step functions [4], Laguerre polynomials [5], Chebyshev polynomials [6], [7], and Legender polynomials [15]. The common approach of these methods is the use of the operational matrix of integration together with the operational matrix of scaling to reduce the differential equation to a set of linear algebraic equations, which is more suitable for computer programming

In this note we will employ the operational matrix of integration and product operational matrix of the general orthogonal polynomials, together with the operational matrix of linear transformation, which will be derived later, to obtain the solution of the scaled system. The operational matrix of linear transformation is derived based on the following properties, namely, the pure recurrence relation

$$
\begin{equation*}
\phi_{1+1}(z)=\left(a_{1} z+b_{1}\right) \phi_{i}(z)-c_{1} \phi_{t-1}(z) \tag{1}
\end{equation*}
$$

with

$$
\phi_{0}(z)=1 ; \phi_{1}(z)=a_{0} z+b_{0}
$$

and the differential recurrence relation

$$
\begin{equation*}
\phi_{1}(z)=A_{i} \dot{\phi}_{1+1}(z)+B_{i} \dot{\phi}_{1}(z)+C_{1} \dot{\phi}_{1-1}(z) \tag{2}
\end{equation*}
$$

where recurrence coefficients $a_{i}, b_{i}, c_{i}$ and differential recurrence coefficients $A_{i}, B_{i}$, and $C_{i}$, are specified by the particular orthogonal polynomials under consideration and some are listed in 19]. The aim of this paper is twofold: 1) to derive an operational matrix of linear transformation for general orthogonal polynomials so that the scaled

[^32]
## Optimal output feedback for non-zero set point regulation: the discrete-time case

WASSIM M. HADDAD $\dagger$ and DENNIS S. BERNSTEIN ${ }_{+}^{+}$
Optimal discrete-time static output feedback is considered for a non-zero set point problem with non-zero mean disturbances. The optimal control law consists of a closed-loop component for feeding back the measurements and a constant open-loop component which accounts for the non-zero set point and non-zero disturbance mean. An additional feature is the presence of state-, control- and measurement-dependent white noise. It is shown that in the absence of multiplicative disturbances, the closed-loop controller can be designed independently of the open-loop control.

## Notation and definitions

$\mathbf{R}, \mathbf{R}^{r \times s}, \mathbf{R}^{r}, \mathbf{E}$ real numbers, $r \times s$ real matrices, $\mathbf{R}^{r \times 1}$, expectation
$I_{n},()^{T} \quad n \times n$ identity, transpose
$\otimes$ Kronecker product
$\operatorname{tr} Z$ trace of square matrix $Z$
asymptotically
stable matrix $n, m, l, p$ positive integers
$x \quad n$-dimensional vector
$u, y \quad m$-, $l$-dimensional vectors
$A, A_{i} ; B, B_{i} ; C, C_{i} \quad n \times n$ matrices, $n \times m$ matrices, $l \times n$ matrices, $i=1, \ldots, p$
$L, K \quad r \times n$ matrix, $m \times l$ matrix
$\delta, \gamma, \alpha \quad r$-, $n$-, $m$-dimensional vectors
$k$ discrete-time index $1,2, \ldots$
$v_{i}(k)$ unit variance white noise, $i=1, \ldots, p$
$w_{1}(k), w_{2}(k) \quad n$-dimensional, $l$-dimensional white noise processes
$V_{1}, V_{2} n \times n$ covariance of $w_{1}, l \times l$ covariance of $w_{2} ; V_{1} \geqslant 0, V_{2} \geqslant 0$
$V_{12} n \times l$ cross-covariance of $w_{1}, w_{2}$
$R_{1}, R_{2} \quad r \times r$ and $m \times m$ state and control weightings; $R_{1} \geqslant 0, R_{2} \geqslant 0$
$R_{12} \quad r \times m$ cross weighting; $R_{1}-R_{12} R_{2}^{-1} R_{12}^{\top} \geqslant 0$
$\tilde{A}, \bar{A}_{i} \quad A+B K C, A_{i}+B_{i} K C+B K C_{i}, i=1, \ldots, p$
A $I_{n}-\tilde{A}$
$\widetilde{B} \quad B \alpha+\gamma$
$\tilde{B}_{i} \quad B_{i} \alpha$

$$
\tilde{w}=w_{1}+B K w_{2}+\sum_{i=1}^{p} B_{i} K w_{2}
$$

Received 9 March 1987.

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$$
\begin{array}{ll}
\tilde{V} & V_{1}+V_{12} K^{\mathrm{T}} B^{\mathrm{T}}+B K V_{12}^{\mathrm{T}}+B K V_{2} K^{\mathrm{T}} B^{\mathrm{T}}+\sum_{i=1}^{p} B_{i} K V_{2} K^{\mathrm{T}} B_{i}^{\mathrm{T}} \\
\tilde{R} & L^{\mathrm{T}} R_{1} L+L^{\mathrm{T}} R_{12} K C+C^{\mathrm{T}} K^{\mathrm{T}} R_{12}^{\mathrm{T}} L+C^{\mathrm{T}} K^{\mathrm{T}} R_{2} K C \\
& +\sum_{i=1}^{p} C_{i}^{\mathrm{T}} K^{\mathrm{\top}} R_{2} K C_{i}
\end{array}
$$

For arbitrary $m \in \mathbf{R}^{n}$ and $Q, P \in \mathbf{R}^{n \times n}$ define:

$$
\begin{array}{ll}
R_{2 s} \triangleq R_{2}+B^{\mathrm{T}} P B+\sum_{i=1}^{p} B_{i}^{\mathrm{T}} P B_{i}, & V_{2 s} \triangleq V_{2}+C Q C^{\mathrm{T}}+\sum_{i=1}^{p} C_{i}\left(Q+m m^{\mathrm{T}}\right) C_{i}^{\mathrm{T}} \\
R_{2 a} \triangleq R_{2}+B^{\mathrm{T}} P B, & V_{2 a} \triangleq V_{2}+C Q C^{\mathrm{T}} \\
P_{s} \triangleq B^{\mathrm{T}} P A+R_{12}^{\mathrm{T}}+\sum_{i=1}^{p} B_{i}^{\mathrm{T}} P A_{i}, & Q_{s} \triangleq A Q C^{\mathrm{T}}+V_{12}+\sum_{i=1}^{p} A_{i}\left(Q+m m^{\mathrm{T}}\right) C_{i}^{\mathrm{T}} \\
P_{a} \triangleq B^{\mathrm{T}} P A+R_{12}^{\mathrm{T}}, & Q_{a} \triangleq A Q C^{\mathrm{T}}+V_{12} \\
P_{s 1} \triangleq R_{12}^{\mathrm{T}}+\sum_{i=1}^{p} B_{i}^{\mathrm{T} P A_{i},} & Q_{s 1} \triangleq V_{12}+\sum_{i=1}^{p} A_{i}\left(Q+m m^{\mathrm{T}}\right) C_{i}^{\mathrm{T}}
\end{array}
$$

## 1. Introduction

The quadratic performance criterion

$$
\begin{equation*}
J=\sum_{k=0}^{N} x^{\mathrm{T}}(k) R_{1} x(k)+u^{\mathrm{T}}(k) R_{2} u(k) \tag{1}
\end{equation*}
$$

expresses the desire to minimize deviations of the state $x(k)$ of the system

$$
\begin{equation*}
x(k+1)=A x(k)+B u(k)+w(k) \tag{2}
\end{equation*}
$$

from the regulation point $x=0$. As is well known (Kwakernaak and Sivan, 1972, pp. 504-509), the non-zero set point criterion

$$
\begin{equation*}
J_{r}=\sum_{k=0}^{N}[x(k)-\bar{x}]^{\top} R_{1}[x(k)-\bar{x}]+u^{\top}(k) R_{2} u(k) \tag{3}
\end{equation*}
$$

presents no additional difficulty so long as $x(k)$ and $u(k)$ are replaced by $x(k)-\bar{x}$ and $u(k)-\bar{u}$. where $\bar{u}$ satisfies

$$
\begin{equation*}
\bar{x}=A \bar{x}+B \bar{u} \tag{4}
\end{equation*}
$$

Closer inspection, however, reveals that this approach is suboptimal. Specifically, the offset $\bar{u}$ in the control may correspond to an unacceptably high level of control effort when $\bar{u}^{\top} R_{2} \bar{u}$ is large. Hence (3) overlooks design tradeoffs concerning the control effort required for maintaining the non-zero regulation point $\bar{x}$. Moreover, such an approach is impossible when $\bar{u}$ satisfying (4) does not exist.

A significant advance in extending the full-state-feedback $L Q R$ formulation to steady-state periodic tracking problems (and hence to the special case of non-zero set point regulation) was given by Artstein and Leizarowitz (1985). Bernstein and Haddad ( 1987 b) generalize the results of Artstein and Leizarowitz (1985) for the nonzero set point regulation problem to include noisy and non-noisy measurements, weighted and unweighted controls, correlated plant/measurement noise, cross weighting, non-zero mean disturbances, and state-, control- and measurement-dependent multiplicative white noise. They consider the steady-state performance criterion

$$
\begin{equation*}
J=\lim _{t \rightarrow x} \mathrm{E}\left[(L x(t)-\delta)^{\mathrm{\top}} R_{1}(L x(t)-\delta)+2(L x(t)-\delta)^{\top} R_{12} u(t)+u^{\top}(t) R_{2} u(t)\right] \tag{5}
\end{equation*}
$$

where $\delta$ is the non-zero regulation point. For full-state feedback with $R_{12}=0$ and $L=$ identity, Artstein and Leizarowitz (1985) show that for a constant offset control law

$$
\begin{equation*}
u(t)=K x(t)+\alpha \tag{6}
\end{equation*}
$$

$K$ and $\alpha$ are given bv

$$
\begin{gather*}
K=-R_{2}^{-1} B^{\top} P  \tag{7}\\
\alpha=-R_{2}^{-1} B^{\top}(A-\Sigma P)^{-\top} R_{1} \delta \tag{8}
\end{gather*}
$$

where $P$ satisfies the Riccati equation

$$
0=A^{\mathrm{T}} P+P A+R_{1}-P \Sigma P
$$

with

$$
\Sigma \triangleq B R_{2}^{-1} B^{\top}
$$

Two features of the control law (6)-(8) are noteworthy. First, (6) consists of both a closed-loop feedback component $K x(t)$ and an open-loop component $\alpha$ depending upon the regulation point $\delta$. And, second (and more important), is the observation that the closed-loop control component is independent of the open-loop component. From a practical point of view this feature is quite useful since it implies that the feedback gain $K$ can be determined without regard to the set point. Hence a change in the desired set point $\delta$ during on-line operation does not necessitate re-solving the Riccati equation in real time; only $x$ requires updating. For a new value of $\delta, x$ can readily be recomputed on-line via the matrix multiplication operation (8). In the presence of multiplicative disturbances, however, the independence of the closed-loop component from the open-loop component is lost.

The purpose of the present paper is to provide a self-contained derivation of the optimality conditions for the non-zero set point problem in the discrete-time case. To obtain a realistic problem setting, we consider the case in which the full state is not available, but rather only noise-corrupted measurements of linear combinations of states For greater design flexibility, we also allow the possibility for correlated plant and measurement noise. In addition, we consider the dual design feature, namely, cross weighting in the performance criterion. The presence of a non-zero constant plant disturbance in conjunction with zero-mean white plant disturbances, i.e. a nonzero mean disturbance, is also considered. Our results show that the presence of a non-zero constant disturbance component leads to an additional offset in the openloop component of the control. Finally, in addition to the above generalizations we allow for the presence of multiplicative disturbances in the plant. The control law thus generalizes previous results involving state-, control- and measurement-dependent noise (Bernstein and Haddad 1987). As shown in Bernstein and Greeley (1986) and Haddad (1987), the multiplicative white noise model can be used for robustness with respect to plant parameter variations.

## 2. Non-zero set point regulation

2.1. Non-zero set point problem

Given the $\boldsymbol{n}$ th-order controlled system

$$
\begin{equation*}
x(k+1)=\left(A+\sum_{i=1}^{p} v_{i}(k) A_{i}\right) x(k)+\left(B+\sum_{i=1}^{p} v_{i}(k) B_{i}\right) u(k)+w_{1}(k)+\gamma \tag{9}
\end{equation*}
$$

with measurements

$$
\begin{equation*}
y(k)=\left(C+\sum_{i=1}^{p} v_{i}(k) C_{i}\right) x(k)+w_{2}(k) \tag{10}
\end{equation*}
$$

where $k=1,2, \ldots$, determine $K$ and $\alpha$ such that the static output feedback controller

$$
\begin{equation*}
u(k)=K y(k)+\alpha \tag{11}
\end{equation*}
$$

minimizes the steady-state performance criterion

$$
\begin{equation*}
J(K, \alpha) \triangleq \lim _{k \rightarrow \infty} \mathrm{E}\left[(L x(k)-\delta)^{\mathrm{T}} R_{1}(L x(k)-\delta)+2(L x(k)-\delta)^{\mathrm{\top}} R_{12} u(k)+u^{\mathrm{T}}(k) R_{2} u(k)\right] \tag{12}
\end{equation*}
$$

Using the notation of § 1 the closed-loop system (9)-(11) can be written as

$$
\begin{equation*}
x(k+1)=\left(\tilde{A}+\sum_{i=1}^{p} v_{i}(k) \tilde{A}_{i}\right) x(k)+\tilde{B}+\sum_{i=1}^{p} v_{i}(k) \tilde{B}_{i}+\tilde{w}(k) \tag{13}
\end{equation*}
$$

To analyse (13) define the second-moment and covariance matrices

$$
\tilde{Q}(k) \triangleq \mathrm{E}\left[x(k) x^{\mathrm{\top}}(k)\right], \quad Q(k) \triangleq \tilde{Q}(k)-m(k) m^{\mathrm{\top}}(k)
$$

where $m(k) \triangleq E[x(k)]$. It follows from (13) that $\tilde{Q}(k), Q(k)$ and $m(k)$ satisfy

$$
\begin{align*}
& \tilde{Q}(k+1)=\tilde{A} \tilde{Q}(k) \tilde{A}^{\mathrm{T}}+\tilde{A} m(k) \tilde{B}^{\mathrm{T}}+\tilde{B} m^{\mathrm{T}}(k) \tilde{A}^{\mathrm{T}}+\tilde{B} \tilde{B}^{\top} \\
& +\sum_{i=1}^{\rho}\left[\tilde{A_{i}} \tilde{Q}(k) \tilde{A}_{i}^{\mathrm{T}}+\tilde{B}_{i} m^{\mathrm{T}}(k) \tilde{A}_{i}^{\mathrm{T}}+\tilde{A}_{i} m(k) \tilde{B}_{i}^{\mathrm{T}}+\tilde{B}_{i} \tilde{B}_{i}^{\mathrm{T}}\right]+\tilde{V}  \tag{14}\\
& Q(k+1)=\tilde{A} Q(k) \tilde{A}^{\mathrm{T}}+\sum_{i=1}^{p}\left[\tilde{A_{i} Q(k) \tilde{A}_{i}^{\mathrm{T}}+\tilde{A}_{i} m(k) m^{\mathrm{T}}(k) \tilde{A}_{i}^{\mathrm{T}}}\right. \\
& \left.\quad+\tilde{B}_{i} m^{\mathrm{T}}(k) \tilde{A}_{i}^{\mathrm{T}}+\tilde{A}_{i} m(k) \tilde{B}_{i}^{\mathrm{T}}\right]+\tilde{V}  \tag{15}\\
& \quad m(k+1)=\tilde{A} m(k)+\tilde{B} \tag{16}
\end{align*}
$$

To consider the steady state, we restrict our consideration to the set of closed-loop second-moment stabilizing gains

$$
\mathbf{S}_{3} \triangleq\left\{K: \tilde{A} \otimes \tilde{A}+\sum_{i=1}^{p} \tilde{A}_{i} \otimes \tilde{A}_{i} \text { is asymptotically stable }\right\}
$$

It follows from fundamental properties of Lyapunov equations that if $K \in \mathbf{S}_{s}$, then $\tilde{A}$ is also asymptotically stable. Hence, for $K \in \mathbf{S}_{s}, \tilde{Q} \triangleq \lim _{k \rightarrow \infty} \tilde{Q}(k), Q \triangleq \lim _{k \rightarrow \infty} Q(k)$ and $m \triangleq \lim _{k \rightarrow \infty} m(k)$ exist and satisfy

$$
\begin{gather*}
\tilde{Q}=\tilde{A} \tilde{Q} \tilde{A}^{\mathrm{T}}+\tilde{A} m \tilde{B}^{\mathrm{T}}+\tilde{B} m^{\mathrm{T}} \tilde{A}^{\mathrm{T}}+\tilde{B} \tilde{B}^{\mathrm{T}} \\
+\sum_{i=1}^{R}\left[\tilde{A_{i}} \tilde{Q} \tilde{A}_{i}^{\mathrm{T}}+\tilde{B}_{i} m^{\mathrm{T}} \tilde{A}_{i}^{\mathrm{T}}+\tilde{A}_{i} m \tilde{B}_{i}^{\mathrm{T}}+\tilde{B}_{i} \tilde{B}_{i}^{\mathrm{T}}\right]+\tilde{V}  \tag{17}\\
Q=\tilde{A} Q \tilde{A}^{\mathrm{T}}+\sum_{i=1}^{R}\left[\tilde{A}_{i} Q \tilde{A}_{i}^{\mathrm{T}}+\tilde{A}_{i} m m^{\mathrm{T}} \tilde{A}_{i}^{\mathrm{T}}+\tilde{B}_{i} m^{\mathrm{T}} \tilde{A}_{i}^{\mathrm{T}}+\tilde{A}_{i} m \tilde{B}_{i}^{\mathrm{T}}+\tilde{B}_{i} \tilde{B}_{i}^{\mathrm{T}}\right]+\tilde{V}  \tag{18}\\
m=\mathbf{A}^{-1} \tilde{B} \tag{19}
\end{gather*}
$$

Note that since $\tilde{A}$ is asymptotically stable, the inverse in (19) exists. For $K \in \mathbf{S}_{\mathbf{s}}$, it now
follows that $J(K, \alpha)$ is given by

$$
\begin{align*}
J(K, \alpha)= & \operatorname{tr}\left[\left(Q+m m^{\mathrm{T}}\right) \tilde{R}\right]+\operatorname{tr}\left[K^{\mathrm{T}} R_{2} K V_{2}\right]+\delta^{\mathrm{T}} R_{1} \delta-2 m^{\mathrm{T}} L^{\mathrm{T}} R_{1} L \delta \\
& +2 m^{\mathrm{T}} L^{\mathrm{T}} R_{12} \alpha-2 \delta^{\mathrm{T}} R_{12} K C m-2 \delta^{\mathrm{T}} R_{12} \alpha+2 m^{\mathrm{T}} C^{\mathrm{T}} K^{\mathrm{T}} R_{2} \alpha+\alpha^{\mathrm{T}} R_{2} \alpha \tag{20}
\end{align*}
$$

Associated with $Q$ is its dual $P \geqslant 0$ which is the unique solution of

$$
\begin{equation*}
P=\tilde{A}^{\mathrm{T}} P \tilde{A}+\sum_{i=1}^{p} \tilde{A}_{i}^{\mathrm{\top}} P \tilde{A}_{i}+\tilde{R} \tag{21}
\end{equation*}
$$

To obtain closed-form expressions for the feedback gain $K$, we further restrict consideration to the set

$$
\mathbf{S}_{s}^{+} \triangleq\left\{K \in \mathbf{S}_{s}: R_{2 s}>0, V_{2 s}>0 \text { and } \Psi_{s} \text { is invertible }\right\}
$$

where

$$
\begin{aligned}
\Psi_{s} \triangleq & B^{\mathrm{T}} A^{-\mathrm{T}} L^{\mathrm{T}} R_{1} L A^{-1} B+B^{\mathrm{T}} \mathbf{A}^{-\mathrm{T}} L^{\mathrm{T}} R_{12}\left(I_{m}+K C A^{-1} B\right) \\
& +\left(I_{m}+K C A^{-1} B\right)^{\mathrm{T}} R_{12}^{\mathrm{T}} L A^{-1} B+\left(I_{m}+K C A^{-1} B\right)^{\mathrm{T}} R_{2}\left(I_{m}+K C A^{-1} B\right) \\
& +\sum_{i=1}^{p}\left[B^{\mathrm{T}} \mathbf{A}^{-\mathrm{T}} \tilde{A}_{i}^{\mathrm{T}} P \tilde{A}_{i} \mathbf{A}^{-1} B+B_{i}^{\mathrm{T}} P \tilde{A}_{i} \mathbf{A}^{-1} B+B^{\mathrm{T}} \mathbf{A}^{-\mathrm{T}} A_{i}^{\mathrm{T}} P B_{i}\right. \\
& \left.+B_{i}^{\mathrm{T}} P B_{i}+B^{\mathrm{T}} \mathbf{A}^{-\mathrm{T}} C_{i}^{\mathrm{T}} K^{\mathrm{T}} R_{2} K C_{i} A^{-1} B\right]
\end{aligned}
$$

Furthermore, we assume that

$$
\begin{equation*}
\left[B_{i} \neq 0 \Rightarrow C_{i}=0\right], \quad i=1, \ldots, p \tag{22}
\end{equation*}
$$

i.e. for each $i \in\{1, \ldots, p\}, B_{i}$ and $C_{i}$ are not both non-zero. Essentially, (22) expresses the condition that the control-dependent and measurement-dependent disturbances are independent. There are no constraints, however, on correlation with the statedependent noise. For the statement of the main theorem define

$$
\begin{gathered}
\Lambda_{s} \triangleq B^{\mathrm{T}} \mathbf{A}^{-\mathrm{T}} L^{\mathrm{T}}\left(R_{1} L+R_{12} K C\right) A^{-1}+\left(I_{m}+K C A^{-1} B\right)^{\mathrm{T}}\left(R_{12}^{\mathrm{T}} L+R_{2} K C\right) \mathbf{A}^{-1} \\
+\sum_{i=1}^{p}\left(\left[A_{i} A^{-1} B+B_{i}\right]^{\mathrm{T}} P \tilde{A}_{i} A^{-1}+B^{\mathrm{T}} A^{-\mathrm{T}} C_{i}^{\mathrm{T}} K^{\mathrm{T}} R_{2} K C_{i} A^{-1}\right) \\
\Omega \triangleq B^{\mathrm{T}} \mathbf{A}^{-\mathrm{T}} L^{\mathrm{T}} R_{1} L+R_{12}^{\mathrm{T}} L+B^{\mathrm{T}} \mathbf{A}^{\mathrm{T}} C^{\mathrm{T}} K^{\mathrm{T}} R_{12}^{\mathrm{T}} L
\end{gathered}
$$

Theorem 2.1
Suppose $K$ and $x$ solve the non-zero set point problem with $K \in \mathbf{S}_{s}^{+}$. Then there exist $n \times n Q, P \geqslant 0$ such that

$$
\begin{gather*}
K=-R_{2 s}^{-1}\left(B^{\mathrm{T}} P A Q C^{\mathrm{T}}+P_{s 1} Q C^{\mathrm{T}}+B^{\mathrm{T}} P Q_{s 1}\right) V_{2 s}^{-1}  \tag{23}\\
x=-\Psi_{s}^{-1}\left[\Lambda_{s i}+\Omega \delta\right] \tag{24}
\end{gather*}
$$

and such that $Q$ and $P$ satisfy

$$
\begin{align*}
& Q=A Q A^{\top}+V_{1}+\sum_{i=1}^{p}\left[\left(A_{i}+B_{i} K C\right) Q\left(A_{i}+B_{i} K C\right)^{\top}+B_{i} K V_{2} K^{\mathrm{T}} B_{i}^{\mathrm{T}}+\tilde{A}_{i} m m^{\mathrm{T}} \tilde{A}_{i}^{\mathrm{T}}\right. \\
& \left.+\bar{B} m^{\mathrm{T}} \tilde{A}_{i}^{\top}+\tilde{A}_{i} m \tilde{B}^{\top}+\tilde{B}_{i} \tilde{B}_{i}{ }^{\mathrm{T}}\right] \\
& +\left(Q_{s}+B K V_{2 s}\right) V_{2 s}^{-1}\left(Q_{s}+B K V_{2 s}\right)^{\top}-Q_{s} V_{25} Q_{s}^{\top}  \tag{25}\\
& P=4^{\mathrm{T}} P A+R_{1}+\sum_{i=1}^{R}\left[\left(A_{i}+B K C_{i}\right)^{\top} P\left(A_{i}+B K C_{i}\right)+C_{i}^{\top} K^{\top} R_{2} K C_{i}\right. \\
& \left.+\left(P_{s}+R_{2 s} K C\right)^{\top} R_{2 s}^{-1}\left(P_{s}+R_{2 s} K C\right)-P_{s}^{\top} R_{2 s} P_{s}\right] \tag{26}
\end{align*}
$$

Proof
The derivation of the necessary conditions is a straightforward application of the Lagrange multiplier technique. To optimize (20) over $\mathbf{S}_{s}^{+}$subject to the constraints (18) and (19), form the lagrangian

$$
\begin{aligned}
L(K, \alpha, Q, P, m) \triangleq \operatorname{tr} & {\left[\lambda_{0} J(K, x)+\left(\tilde{A} Q \tilde{A}^{\top}+\sum_{i=1}^{\rho}\left[\tilde{A}_{i} Q \tilde{A}_{i}^{\top}+\tilde{A}_{i} m m^{\top} \tilde{A}_{i}^{\top}+\tilde{B}_{i} m^{\top} \tilde{A}_{i}^{\top}\right.\right.\right.} \\
& \left.\left.\left.+\tilde{A}_{i} m \tilde{B}_{i}^{\top}+\tilde{B}_{i} \tilde{B}_{i}^{\top}\right]+\tilde{V}-Q\right) P+\dot{\lambda}^{\top}(\tilde{A} m+\tilde{B}-m)\right]
\end{aligned}
$$

where the Lagrange multipliers $\lambda_{0} \geqslant 0, \lambda \in \mathbf{R}^{n}$ and $P \in \mathbf{R}^{n \times n}$ are not all zero. Setting $\partial \mathrm{L} / \partial Q=0$ and using the second-moment stability assumption it follows that $\lambda_{0}=1$ without loss of generality. Thus the stationarity conditions are given by

$$
\begin{align*}
& \frac{\partial \mathrm{L}}{\partial Q}=\tilde{A}^{\mathrm{\top}} P \tilde{A}+\sum_{i=1}^{R} \tilde{A}_{i}^{\mathrm{T}} P \tilde{A}_{i}+\tilde{R}-P=0  \tag{27}\\
& \frac{\partial \mathrm{~L}}{\partial P}=\tilde{A} Q \tilde{A}^{\top}+\sum_{i=1}^{p}\left[\tilde{A}_{i} Q \tilde{A}_{i}^{\top}+\tilde{A}_{i} m m^{\top} \tilde{A}_{i}^{\top}+\tilde{B}_{i} m^{\top} \tilde{A}_{i}^{\top}+\tilde{A}_{i} m \tilde{B}_{i}^{\top}+\tilde{B}_{i} \tilde{B}_{i}^{\top}\right] \\
& +\tilde{V}-Q=0  \tag{28}\\
& \frac{\partial \mathbf{L}}{\partial \hat{\lambda}^{\top}}=\tilde{A} m+\tilde{B}-m=0  \tag{29}\\
& \frac{\partial \mathrm{~L}}{\partial \mathrm{~K}}=R_{2 s} K V_{2 s}+R_{12}^{\top} Q C^{\top}+B^{\mathrm{T}} P A Q C^{\top} \\
& +\sum_{i=1}^{p}\left[B_{i} P A_{i} Q C+B^{\mathrm{T}} P A_{i}\left(Q+m m^{\mathrm{T}}\right) C_{i}^{\mathrm{T}}\right]+B^{\mathrm{T}} P V_{12}=0  \tag{30}\\
& \frac{\partial \mathrm{~L}}{\partial \alpha}=\sum_{i=1}^{\mathrm{g}}\left[B_{i}^{\top} P A_{i} m+B_{i}^{\mathrm{T}} P B_{i} K C+B_{i}^{\top} P B_{i} x\right]+\frac{1}{2} B^{\mathrm{T}} \dot{i}+R_{12}^{\mathrm{T}} L m \\
& -R_{12}^{\mathrm{T}} \delta+R_{2} \mathrm{KCm}+\mathrm{R}_{2} \alpha=0  \tag{31}\\
& \frac{\partial \mathrm{~L}}{\partial m}=\tilde{R} m+\sum_{i=1}^{p}\left[\tilde{A}_{i}^{\top} P B_{i}+\tilde{A}_{i}^{\top} P \tilde{A}_{i} m\right]-\frac{1}{2} A^{\top} \dot{\lambda}-L^{\top} R_{1} \delta+L^{\top} R_{12} \chi \\
& -C^{\top} K^{\top} R_{12}^{\top} \delta+C^{\top} K^{\top} R_{2} x=0  \tag{32}\\
& \lambda=2 A^{-\mathrm{T}}\left(\tilde{R} m+\sum_{i=1}^{p}\left[\tilde{A}_{i}^{\top} P \tilde{B}_{i}+\tilde{A}_{i} P \tilde{A}_{i}^{\top} m\right]-L^{\mathrm{T}} R_{1} \delta+L^{\mathrm{T}} R_{12} \alpha\right. \\
& \left.-C^{\top} K^{\top} R_{12}^{\top} \delta+C^{\top} K^{\top} R_{2} x\right) \tag{33}
\end{align*}
$$

Using the definitions for $Q_{s 1}$ and $P_{s 1}$ along with (33), we obtain (23) and (24). Substituting the expressions for the optimal gains into (27) and (28) yields (25) and (26).

## Remark 1

Because of the presence of $\delta$ in (25) via $m$ in both $Q_{s 1}$ and $V_{2 s}$ and in (25) via $\tilde{B}$ (in $m$ ) and $\bar{B}_{i}$, the closed-loop component of the control law (23) cannot be
determined independently of the open-loop component. As shown in the following section, independence is recovered when the multiplicative noise terms are absent.

## Remark 2

To specialize Theorem 2.1 to the standard regulation problem, set $\delta=0$ and $;=0$ yielding Theorem 2.1 of Bernstein and Haddad (1987 a).

## 3. Specializations of Theorem 2.1

A series of specializations of Theorem 2.1 is now given. We begin by deleting all multiplicative white noise terms, i.e.

$$
\begin{equation*}
A_{i}, B_{i}, C_{i}=0, \quad i=1, \ldots, p \tag{34}
\end{equation*}
$$

In this case the stabilizing set $\mathrm{S}_{\mathrm{s}}$ can be characterized by

$$
\mathbf{S}=\{K: \tilde{A} \text { is asymptotically stable }\}
$$

and, furthermore, $\mathbf{S}_{\mathbf{s}}^{+}$becomes

$$
\mathbf{S}^{+} \triangleq\left\{K \in \mathbf{S}: \boldsymbol{R}_{2 a}>0, V_{2 a}>0 \text { and } \Psi_{a} \text { is invertible }\right\}
$$

where

$$
\begin{aligned}
\Psi_{a} \triangleq & B^{\top} \mathbf{A}^{-\top} L^{\mathrm{T}} R_{1} L \mathbf{A}^{-1} B+B^{\mathrm{\top}} \mathbf{A}^{-\mathrm{T}} L^{\mathrm{T}} R_{12}\left(I_{m}+K C A^{-1} B\right) \\
& +\left(I_{m}+K C A^{-1} B\right)^{\mathrm{T}} R_{12}^{\mathrm{T}} L A^{-1} B+\left(I_{m}+K C A^{-1} B\right)^{\mathrm{T}} R_{2}\left(I_{m}+K C A^{-1} B\right)
\end{aligned}
$$

For the statement of Corollary 3.1 define

$$
\Lambda_{a} \triangleq B^{\mathrm{T}} \mathbf{A}^{-\mathrm{T}} L^{\top}\left(R_{1} L+R_{12} K C\right) \mathbf{A}^{-1}+\left(I_{m}+K C A^{-1} B\right)^{\mathrm{\top}}\left(R_{12}^{\mathrm{T}} L+R_{2} K C\right) \mathrm{A}^{-1}
$$

## Corollary 3.1

Assume (34) is satisfied and suppose $K$ and $x$ solve the non-zero set point problem with $K \in \mathbf{S}^{+}$. Then there exist $n \times n Q, P \geqslant 0$ such that

$$
\begin{gather*}
K=-R_{2 a}^{-1}\left(B^{\mathrm{T}} P A Q C^{\mathrm{T}}+R_{12}^{\mathrm{T}} Q C^{\mathrm{T}}+B^{\mathrm{T}} P V_{12}\right) V_{2 a}^{-1}  \tag{35}\\
x=-\Psi_{a}^{-1}\left[\Lambda_{a} i^{\prime}+\Omega \delta\right] \tag{36}
\end{gather*}
$$

and such that $Q$ and $P$ satisfy

$$
\begin{align*}
& Q=A Q A^{\mathrm{T}}+V_{1}+\left(Q_{a}+B K V_{2 a}\right) V_{2 a}^{-1}\left(Q_{a}+B K V_{2 a}\right)^{\mathrm{T}}-Q_{a} V_{2 a} Q_{a}^{\mathrm{T}}  \tag{37}\\
& P=A^{\mathrm{T}} P A+R_{1}+\left(P_{a}+R_{2 a} K C\right)^{\mathrm{\top}} R_{2 a}^{-1}\left(P_{a}+R_{2 a} K C\right)-P_{a}^{\mathrm{T}} R_{2 a} P_{a} \tag{38}
\end{align*}
$$

Finally, setting

$$
\begin{equation*}
\because=0, \quad R_{12}=0, \quad V_{12}=0, \quad r=n, \quad L=I_{n} \tag{39}
\end{equation*}
$$

we obtain the discrete-time version of Artstein and Leizarowitz (1985) for the case of output feedback. Define

$$
\mathbf{S}_{1}^{+} \triangleq\left\{K \in \mathbf{S}: R_{2 a}>0, V_{2 a}>0 \text { and } \Psi_{1}>0\right\}
$$

where

$$
\Psi_{1} \triangleq B^{\top} \mathbf{A}^{-\top} R_{1} A^{-1} B+\left(I_{m}+K C A^{-1} B\right)^{\top} R_{2}\left(I_{m}+K C A^{-1} B\right)
$$

Corollary 3.2
Assume (34) and (39) are satisfied and suppose $K$ and $x$ solve the non-zero set
point problem with $K \in \mathbf{S}_{1}^{+}$. Then there exist $n \times n Q, P \geqslant 0$ such that

$$
\begin{gather*}
K=-R_{2 a}^{-1} B^{\mathrm{T}} P_{A Q C} V_{2 a}^{-1}  \tag{40}\\
\alpha=-\Psi_{1}^{-1} B^{\mathrm{T}} A^{-\mathrm{T}} R_{1} \delta \tag{41}
\end{gather*}
$$

and such that $Q$ and $P$ satisfy

$$
\begin{align*}
Q= & A Q A^{\mathrm{T}}+V_{1}+\left(A Q C^{\mathrm{T}}+B K V_{2 a}\right) V_{2 a}^{-1}\left(A Q C^{\mathrm{T}}+B K V_{2 a}\right)^{\mathrm{T}} \\
& -A Q C^{\mathrm{T}} V_{2 a} C Q A^{\mathrm{T}}  \tag{42}\\
P= & A^{\mathrm{T}} P A+R_{1}+\left(B^{\mathrm{T}} P A+R_{2 a} K C\right)^{\mathrm{T}} R_{2 a}^{-1}\left(B^{\mathrm{T}} P A+R_{2 a} K C\right) \\
& -A^{\mathrm{T}} P B R_{2 a} B^{\mathrm{T}} P A \tag{43}
\end{align*}
$$

## 4. Directions for further research

The extension to fixed-order dynamic compensation for non-zero set point regulation appears possible using the approach of Hyland and Bernstein (1984) and Haddad (1987). A generalization of Theorem 2.1 to design periodic tracking controllers (either static or dynamic) via the parameter optimization approach is being developed.

## ACKNOWLEDGMENT

Dr Bernstein was supported in part by the Air Office of Scientific Research under contract F49620-86-C-0002.

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## Optimal Nonzero Set Point Regulation Via Fixed-Order Dynamic Compensation

WASSIM M. HADDAD and DENNIS S. BERNSTEIN


#### Abstract

Standard LOG control theory is generalized to a regulation problem involving specified nonzero set points for the state and control variables and monzero-mean disturbarices. For generality, the results are ebtained for the probtem of fixed-order (i.e., not necescarily full-order) dyamic compensation. When the state, control, and disturbance offsets are set to zero and the compensitor order is set equal to the plant dimension. the standard LOK result is recovered. These results provide the dynamic counterpart for the nonzero set point regulation results obtained in (1) via static controllers.


## 1. Introduction

As discussed in [1], the standard quadratic performance criterion expresses the desire to maintain the state and control variables in the neighborhood of the origin. If regulation is desired about nonzero state and control offsets. then, in special cases, the set points can be translated to the origin and standard theory can be applied (see. e.g., [2. pp. 270276]). In general. however, (see [1]) such a translation may either be suboptimal or impossible. The latter situation may occur, for example, if the number of state components with specified nonzero set points is greater than the number of controls, while the former is the case when the control offser is particularly costly.

Motivated by the work of Leizarowitz and Artstein [3], [4] on the more general problems of periodic and nonperiodic tracking, the nonzero set point problem was addressed in [1] for the case of static output-feedback controllers. The goal of the present note is to derive analogous results for the case of dynamic compensation considered by Leizarowitz in [5]. As in [1], the solution we obtain has the satisfying feature that the closed-loop dynamic-feedback-compensation gains are independent of the open-loop control components which arise from the state and control set points. Thus, if the state set point is changed during operation, then only the open-loop control components require updating. Consequently, there is no need to recalculate the closed-loop gains by solving Riccati equations in real time. The overall theory thus permits the treatment of step commands within standard LQG thenry.

For generality the development herein incorporates several special features which provide additional flexibility in applications. These include: 1 ) conscant disturbance vectors in addition to zero-mean additive plant and measurement noise (i.e.. nonzero-mean disturbances): 2) correlated plant and measurement noise; 3) state/control performance cross-weighting. 4) arbitrary set points for selected linear combinations of the state and control variables (sec $L_{1}$ and $L_{\text {: }}$ in the problem statement in Section III); and 5) fixed-order (i.e., full- or reduced-order) compensation. Because of the last feature, the results obtuined in the present note also generalize the results of $\{0 \mid$. For clarits, we specialize the main result to the usual full-order LQG case.

## iI. Notation and Definitions

8. $n^{\prime \times s}, 3^{\prime}$, E

Real numbers. $r \times s$ real matrices. $f^{r 1}$. expectation.

Manuscript received September 15. 1987; revised November 23, 1987. This work was supported in part by the Air Force Office of Scientific Research under Contract F49620-86-C-0002.
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IEEE Los Number 8821570.

| $I_{n},()^{\top},()^{\prime}$ | $n \times n$ identity, transpose. group generalized inverse. |
| :---: | :---: |
| asymprotically stable matrix | Matrix with eigenvalues in open lefi-inalf plane. |
| $n, m, l, q, r, n_{c}$ | Positive integers. |
| $\boldsymbol{\pi}$ | $n+n_{r}$. |
| $x, u, y, x_{s}, \tilde{x}$ | $n, m, l, n_{r}, n$-dimensional vectors. |
| $A, B, C, D$ | $n \times n, n \times m, l \times n, l \times m$ matrices. |
| $A_{r}, B_{r}, C_{r}$ | $n_{c} \times n_{c}, n_{c} \times l, m \times n_{c}$ matrices. |
| $L_{1}, L_{2}$ | $q \times n, r \times m$ matrices. |
| $\delta_{1}, \delta_{2}$ | $q, r$-dimensional set point vectors. |
| $\gamma_{1}, \gamma_{2}$ | $n, 1$-dimensional constant disturbance vectors. |
| $\alpha, \alpha_{c}$ | $m, n_{c}$-dimensional control vectors. |
| $\bar{\delta}, \bar{\gamma}, \tilde{\alpha}$ | $\left[\begin{array}{l}\delta_{1} \\ \delta_{2}\end{array}\right],\left[\begin{array}{c}\gamma_{1} \\ B_{c} \gamma_{2}\end{array}\right],\left[\begin{array}{c}\alpha \\ \alpha_{r}\end{array}\right]$. |
| $w_{1}(t), w_{2}(t)$ | $n$, 1-dimensional zero-mean white noise processes. |
| $v_{1}, v_{2}$ | Intensities of $w_{1}, w_{2} ; V_{1} \geq 0, V_{2}>0$. |
| $V_{12}$ | $n \times 1$ cross intensity of $w_{1}, w_{2}$. |
| $\bar{w}(t), \bar{V}$ | $\left[\begin{array}{c}w_{1}(t) \\ B_{c} w_{2}(t)\end{array}\right] \cdot\left[\begin{array}{cc}V_{1} & V_{12} B_{c}^{T} \\ B_{c} V_{12}^{T} & B_{c} V_{2} B_{s}^{r}\end{array}\right]$ |
| $\boldsymbol{R}_{1}, \boldsymbol{R}_{\mathbf{2}}$ | $q \times q$ and $r \times r$ state and control weightings: $R_{1} \geq 0, R_{2} \geq 0, L_{2}^{T} R_{2} L_{2}>0$ |
| $\boldsymbol{R}_{12}$ | $q \times r$ cross weighting: $L_{1}^{T} R_{1} L_{1}-$ $L_{1}^{T} R_{12} L_{2}\left(L_{2}^{T} R_{2} L_{2}\right)^{-1} L_{2}^{T} R_{i 2}^{T} L_{1} \geq 0$. |
| $\bar{R}$ | $\left[\begin{array}{cc}L_{1}^{T} R_{1} L_{1} & L_{1}^{T} R_{12} L_{2} C_{c} \\ C_{c}^{T} \sum_{2}^{\top} R_{12}^{T} L_{1} & C_{c}^{\top} L_{2}^{T} R_{2} L_{2} C_{c}\end{array}\right]$. |
| $\bar{A}, \bar{B}$ | $\left[\begin{array}{cc}A & B C_{c} \\ B_{c} C & A_{c}+B_{c} D C_{r}\end{array}\right],\left[\begin{array}{cc}B & 0 \\ B_{c} D & I_{n_{c}}\end{array}\right]$ |
| $m, m_{c}$ | $n, n_{r}$-dimensional vectors. |
| $\bar{m}$ | $\left[\begin{array}{l}m \\ m\end{array}\right]$ |
| $\tilde{R}_{1}$ | $\left[\begin{array}{ccc}L_{1}^{\top} R_{1} L_{1}-L_{1}^{\top} R_{12} L_{2}\left(L_{2}^{\top} R_{2} L_{2}\right)^{-1} L_{:}^{\dagger} R_{12}^{\top} L_{1} & 0 \\ 0 & 0\end{array}\right]$ |
| $\overline{R_{12}}, \overline{R_{2}}, \bar{R}!$ | $\left[\begin{array}{cc}L L_{:}^{T} R_{1}^{T} L_{1} & L_{2}^{T} R_{2} L_{2} C_{1} \\ 0 & 0\end{array}\right] \cdot\left[\begin{array}{cc}L_{:}^{T} R_{2} L_{2} & 0 \\ 0 & 0\end{array}\right]$. |
|  | $\left[\begin{array}{ccc}\left(L T R_{2} L_{2}\right) & 1 & 0 \\ 0 & 0\end{array}\right]$ |
| $\bar{N}, 5$ | $\left[\begin{array}{cc}L!R_{1}^{\prime} & L!R_{:}^{\prime} \\ 0 & 0\end{array}\right],\left[\begin{array}{cc}L_{1}^{T} R_{1} & L_{1}^{T} R_{1:} \\ C_{1}^{\prime} L_{2}^{T} R_{12}^{T} & C_{i}^{r} L_{2}^{\prime} R_{2}\end{array}\right]$ |

For arbitrary $n \times n Q, P$ define:

$$
\begin{gathered}
Q_{0} \hat{\underline{E}} Q C^{\top}+V_{12}, \quad P_{4} \triangleq B^{\top} P+L_{2}^{\top} R_{12}^{\top} L_{1}, \\
A_{Q} \triangleq A-Q V_{2}^{-1} C^{\top}, \quad A_{p} \triangleq A-B\left(L_{2}^{T} R_{2} L_{2}\right)^{-1} P_{8} .
\end{gathered}
$$

## II. Dynamic Compensation for Nonzero Set Point regulation

## A. Nonzero Set Point Problem

Given the mh-order stabilizable and detectable plant

$$
\begin{gather*}
\dot{x}(t)=A x(t)+B u(t)+w_{1}(t)+\gamma_{1}, \quad t \in[0, \infty)  \tag{3.1}\\
y(t)=C x(t)+D u(t)+w_{2}(t)+\gamma_{2} \tag{3.2}
\end{gather*}
$$

design a fixed-order dynamic compensator

$$
\begin{gather*}
\dot{x}_{c}(t)=A_{c} x_{c}(t)+B_{c} y(t)+\alpha_{c}  \tag{3.3}\\
u(t)=C_{c} x_{r}(t)+\alpha \tag{3.4}
\end{gather*}
$$

which minimizes the steady-state performance criterion

$$
\begin{align*}
J\left(A_{r}, B_{c}, C_{i}, \alpha_{,}, \alpha_{f}\right) \leqq & \lim _{1 \rightarrow \infty} \Xi\left[\left(L_{1} x(t)-\delta_{1}\right)^{T} R_{1}\left(L_{1} x(t)-\delta_{1}\right)\right. \\
& +2\left(L_{1} x(t)-\delta_{1}\right)^{\top} R_{12}\left(L_{2} u(t)-\delta_{2}\right) \\
& \left.+\left(L_{2} u(t)-\delta_{2}\right)^{\top} R_{2}\left(L_{2} u(t)-\delta_{2}\right)\right] \tag{3.5}
\end{align*}
$$

Remark 3.1: The cost functional (3.5) is identical to the LQG criterion (usually stated in terms of an averaged integral) with the exception of the shifted set points $\delta_{1}$ and $\delta_{2}$ and matrices $L_{1}$ and $L_{2}$ for selecting linear combinations of components of $x$ and $u$.

The closed-loop system (3.1)-(3.4) can be writen as

$$
\begin{equation*}
\dot{\dot{x}}(t)=\tilde{A} \tilde{x}(t)+\bar{B} \bar{\alpha}+\tilde{x}(t)+\dot{\gamma}, \quad t \in[0, \infty) \tag{3.6}
\end{equation*}
$$

where $\bar{x}(t) \doteq\left\{x^{T}(t),\left.x_{r}^{T}(t)\right|^{T}\right.$ and the closed-loop disturbance $\tilde{w}(t)$ has nonnegative-definite intensity $\mathcal{V}$. To analyze (3.6) define the covariance matrix

$$
Q(t)=E\left[(\bar{x}(t)-\tilde{m}(t))(\bar{x}(t)-\tilde{m}(t))^{T}\right]=E\left[\bar{x}(t) \bar{x}^{\top}(t)\right]-\tilde{m}(t) \tilde{m}^{\top}(t)
$$

where $\bar{m}(t) \doteq \equiv \mid \bar{x}(t)]$ As shown in [1]. $\bar{Q}(t)$ and $\bar{m}(t)$ satisfy

$$
\begin{gather*}
\dot{Q}(t)=\bar{A} Q(t)+\bar{Q}(t) \bar{A}^{\prime}+\bar{V},  \tag{3.7}\\
\dot{m}(t)=\bar{A} \bar{m}(t)+\bar{B} \bar{\alpha}+\dot{\gamma} . \tag{3.8}
\end{gather*}
$$

To guarantec that $f$ is finite and independent of initial conditions, we restrict our attention to the set of admissible stabilizing compensators

$$
S \doteq\left\{\left(A_{1}, B_{1}, C_{1}\right): \bar{A} \text { is asymptotically stable }\right\}
$$

 exist and satisfy

$$
\begin{align*}
& 0=A \boldsymbol{Q}+Q A^{T}+D  \tag{3.9}\\
& 0=A \hat{f}+\boldsymbol{B} \dot{\alpha}+\dot{\gamma} \tag{3.10}
\end{align*}
$$

Since the value of $J$ is independent of the imernal realization of the transfer function corresponding to (3.3) and (3.4). without loss of
generality we further restrict our attention to the set
$\mathcal{S}^{\circ} \triangleq\left\{\left(A_{c}, B_{c}, C_{c}\right) \in \mathcal{S}:\left(A_{c}, B_{c}\right)\right.$ is controllable
and $\left(A_{c}, C_{r}\right)$ is observable\}.
Now $J\left(A_{c}, B_{c}, C_{c}, \alpha, \alpha_{c}\right)$ is given by

$$
\begin{align*}
J\left(A_{c}, B_{c}, C_{c}, \alpha_{1} \alpha_{c}\right)= & \operatorname{tr}\left[Q+\bar{m} \bar{m}^{T}\right] R-2 m^{T} L_{1}^{T} R_{1} \delta_{1}+\delta_{1}^{T} R_{1} \delta_{1} \\
& +2 m^{T} L_{1}^{T} R_{12} L_{2} \alpha-2 m^{T} L_{1}^{T} R_{12} \delta_{2} \\
& -2 \delta_{1}^{T} R_{12} L_{2} C_{c} m_{c}-2 \delta_{1}^{T} R_{12} L_{2} \alpha \\
& +2 \delta_{1}^{T} R_{12} \delta_{2}+2 m_{c}^{T} C_{c}^{T} L_{2}^{T} R_{2} L_{2} \alpha-2 m_{r}^{T} C_{c}^{T} L_{2}^{T} R_{2} \delta_{2} \\
& -2 \alpha^{T} L_{2}^{T} R_{2} \delta_{2}+\alpha^{T} L_{2}^{T} R_{2} L_{2} \alpha+\delta_{2}^{T} R_{2} \delta_{2} . \tag{3.11}
\end{align*}
$$

To obtain closed-form expressions for the feedback gains we further restrict consideration to the set

$$
\delta^{*} \hat{\equiv}\left\{\left(A_{c}, B_{c}, C_{c}\right) \in \delta^{\prime}: \Omega>0\right\}
$$

where

The following factorization lemma is needed for the statement of the main result.
Lemmo 3.1: Suppose $n \times n Q, P$ are nonnegative definite and rank $Q \mathcal{P}=n_{c}$. Then there exist $n_{c} \times n G, \Gamma$ and $n_{c} \times n_{c}$ invertible $M$ such that

$$
\begin{gather*}
\phi \mathcal{P}=G^{T} M \Gamma  \tag{3.12}\\
\Gamma G^{T}=I_{n_{4}} \tag{3.13}
\end{gather*}
$$

Furthermore, $G, M$, and $\Gamma$ are unique except for a change of basis in $\|^{n} c$. Proof: See [6].
As shown in [6]. $Q P$ has a group generalized inverse $(Q P)^{\prime \prime}=$ $G^{T} M^{\prime} \Gamma$, and the matrix

$$
\begin{equation*}
T \leqq Q P(Q P)^{\prime}=G^{T} \Gamma \tag{3.14}
\end{equation*}
$$

is an oblique projection. A triple ( $G, M, \Gamma$ ) satisfying (3.12) and (3.13)
 projecrive factorization of $Q \dot{P}$. Furthermore, define the coinplementary projection $T_{\perp} \leqslant I_{n}-\tau$. Optimizing (3.11) subject to (3.9) and (3.10) yields the following result illustrated in Fig. 1.

Theorem 3.1: Suppose $\left(A_{c}, B_{c}, C_{c}, \alpha, \alpha_{c}\right)$ solves the nonzero set point problem with $\left(A_{c}, B_{c}, C_{c}\right) \in \delta^{\prime \prime}$. Then there exist $n \times n$ nonnegativedefinite matrices $Q, P, Q, P$ such that, for some projective factorization $(G, M, \Gamma)$ of $Q \hat{P}, A_{c}, B_{c}, C_{c}, \alpha$, and $\alpha_{c}$ are given by

$$
\begin{equation*}
A_{c}=\Gamma\left|A-B\left(L_{2}^{r} R_{2} L_{2}\right)^{-1} P_{a}-Q_{a} V_{2}^{-1} C+Q_{a} V_{2}^{-1} D\left(L_{2}^{r} R_{2} L_{2}\right)^{-1} P_{a}\right| G^{T} \tag{3.15}
\end{equation*}
$$

$$
\left[\begin{array}{c}
a  \tag{3.17}\\
a_{c}
\end{array}\right]=\Omega^{-1}\left[\left(R_{12}-B^{\top} A^{-\top} R\right) A^{-1} \dot{\gamma}+\left(N-B^{\top} A-\Gamma S\right) \delta\right]
$$



Fig. 1.
and such that $Q, P, Q$, and $\bar{P}$ satisfy

$$
\begin{align*}
& 0=A^{T} P+P A+L_{1}^{T} R_{1} L_{1}-P_{j}^{T}\left(L_{2}^{T} R_{3} L_{2}\right)^{-1} P_{g}  \tag{3.19}\\
& +\tau_{1}^{T} P_{d}^{T}\left(L_{2}^{T} R_{2} L_{2}\right)^{-1} P_{A} T_{\perp}, \tag{3.20}
\end{align*}
$$

$$
\begin{align*}
& \operatorname{rank} \boldsymbol{Q}=\operatorname{rank} \boldsymbol{P}=\operatorname{rank} \boldsymbol{Q} \boldsymbol{P}=\boldsymbol{n}_{\mathrm{f}} . \tag{3.22}
\end{align*}
$$

Proof: See Section IV.
Remark 3.2: The results of [6] are a special case of Theorem 3.1. To see this set $\delta_{1}=\gamma_{1}=0, \delta_{2}=0, \gamma_{2}=0, L_{1}=I_{n}$, and $L_{2}=I_{m}$, which yields the results of 16$]$ with the added features of correlated plant/ measurement noise ( $V_{12}$ ). cross weighting ( $R_{12}$ ), and a direct transmission term ( $D$ ) in the plant dynamics

As discussed in |6|, in the full-order (LQG) case $n_{c}=n$ the Lyapunov equations (3.21) and (3.22) for $\phi$ and $\boldsymbol{P}$ are superfluous. In this case $\boldsymbol{G}=$ $\Gamma{ }^{\prime}$ and thus $G=\Gamma=\tau=I_{n}$ without loss of generality. To develop further connections with standard LQG theory, assume

$$
\begin{equation*}
L_{1}=I_{n}, L_{2}=I_{m}, R_{12}=0, V_{12}=0 \tag{3.24}
\end{equation*}
$$

and definc

$$
\begin{aligned}
& \AA=\left[\begin{array}{cc}
R_{1} & 0 \\
0 & C_{r}^{T} R_{1} C
\end{array}\right] \text {, 苗, }=\left[\begin{array}{cc}
R_{1} & 0 \\
0 & 0
\end{array}\right] \text {. } \\
& \mathbb{R}_{12} \doteq\left[\begin{array}{cc}
0 & R_{2} C \\
0 & 0
\end{array}\right] \cdot \mathbb{R}_{2}=\left[\begin{array}{cc}
R_{2} & 0 \\
0 & 0
\end{array}\right] \text {. } \\
& \hat{\pi} \cong\left[\begin{array}{cc}
0 & R_{3} \\
0 & 0
\end{array}\right], \tilde{\delta} \doteq\left[\begin{array}{cc}
R_{1} & 0 \\
0 & C_{r}^{r} R_{3}
\end{array}\right] .
\end{aligned}
$$

In this case $S$ " becomes

$$
\delta^{\prime} \leq\left\{\left(A_{c}, B_{c}, C_{c}\right) \in \delta^{\prime}: \Omega>0\right\}
$$

where

Corollary 3.1: Let $n_{r}=n$, assume (3.24) is satisfied. and suppose ( $A_{c}, B_{c}, C_{c}, \alpha_{,} \alpha_{c}$ ) solves the full-crder nonzero set point problem with $\left(A_{c}, B_{c}, C_{c}\right) \in \mathcal{S}^{n}$. Then there exist $n \times n$ nonnegative-definite matrices $Q, P$ such that $A_{c}, B_{c}, C_{c}, \alpha$, and $\alpha_{c}$ are given by

$$
\begin{gathered}
A_{r}=A-B R_{2}^{-1} B^{T} P-Q C^{T} V_{2}^{-1} C+Q C^{T} V_{2}^{-1} D R_{2}^{-1} B^{T} P, \\
B_{c}=Q C^{T} V_{2}^{-1} \\
C_{r}=-R_{2}^{-1} B^{T} P, \\
{\left[\begin{array}{c}
\alpha \\
\alpha_{c}
\end{array}\right]=\Omega^{-1}\left[\left(\delta_{12}-B^{T} \bar{A}-T \bar{R}\right) \bar{A}^{-1} \tilde{\gamma}+\left(\bar{\delta}-B^{T} \bar{A}^{-T} \delta\right) \delta\right]}
\end{gathered}
$$

and such that $Q, P$ satisfy

$$
\begin{aligned}
& 0=A Q+Q A^{T}+V_{1}-Q C^{T} V_{2}^{-1} C Q \\
& 0=A^{T} P+P A+R_{1}-P B R_{2}^{-1} B^{\prime} P
\end{aligned}
$$

Remark 3.3: Note that by setting $\delta_{1}=\gamma_{1}=0, \delta_{2}=C, \gamma_{2}=0$, and $D$ $=0$. Corollary 31 yields the standard LQG result.

Remark 3.4: It is easy to see that in the full-order case $n_{\mathrm{s}}=n$ a solution to the nonzero set point problem exists as long as $\bar{\Omega}$ is positive definite. In the reduced-order case, however, the situation is more complex. For details, see 181 .

## IV. Proof of Theorem 3.1

To optamize ( 3.11 ) over the open set $S$ " subject to the constraints (3.9) and (3.10). form the Lagrangian

$$
\begin{aligned}
& \mathcal{L}\left(A_{c}, B_{r}, C_{r}, \alpha, \alpha_{c}\right) \doteq \operatorname{tr}\left\{\lambda_{r} J\left(A_{c}, B_{r}, C_{c}, \alpha, \alpha_{c}\right)\right. \\
& \\
& \left.\quad+\left(\tilde{A} \bar{Q}+\bar{Q} \bar{A}^{r}+\tilde{V}\right) \bar{P}+\bar{\lambda}^{r}(\bar{A} \bar{m}+\tilde{B} \tilde{\alpha}+\tilde{\gamma})\right\}
\end{aligned}
$$

where the Lagrange multipliers $\lambda_{0} \geq 0, \bar{\lambda} \in s^{n}$, and $\mathcal{F} \in j 3^{i \times n}$ are not all zero. Setting $\partial \mathscr{L} / \partial Q=0$ and using the fact that $\bar{A}$ is asymprotically stable, it follows that $\lambda_{n}=1$ without loss of generality.

Now partition $\tilde{n} \times \dot{n} \oint, \dot{P}$ into $n \times n, n \times n_{c}, n_{6} \times n_{c}$ subblocks and $\bar{\lambda} \in s^{n}$ into $\Omega^{n}$ and $s^{n} c$ components as

$$
Q=\left[\begin{array}{ll}
Q_{1} & Q_{12} \\
Q_{12}^{r} & Q_{2}
\end{array}\right], P=\left[\begin{array}{ll}
P_{1} & P_{12} \\
P_{12}^{r} & P_{2}
\end{array}\right], \bar{\lambda}=\left[\begin{array}{c}
\lambda_{1} \\
\lambda_{2}
\end{array}\right] .
$$

Thus, the stationarity conditions are given by

$$
\begin{align*}
& \frac{\partial \mathcal{L}}{\partial P}=A Q+Q A^{r}+Q=0 .  \tag{4.1}\\
& \frac{\partial \mathcal{L}}{\partial Q}=A^{T} P+P A+R=0, \tag{4.2}
\end{align*}
$$

$$
\begin{align*}
& \frac{\partial \mathcal{L}}{\partial \tilde{f}}=\hat{R} \tilde{m}+\frac{1}{2} A^{r} \bar{\lambda}+R_{12}^{r} \dot{\alpha}-S \delta=0,  \tag{4.3}\\
& \frac{\partial \mathcal{L}}{\partial \dot{\alpha}}=\left[\begin{array}{c}
{\left[L_{2}^{T} R_{2} L_{2} 0\right] \bar{\alpha}+\left[L{ }_{2}^{r} R_{12}^{T} L_{1} L{ }_{2}^{T} R_{2} L_{2} C_{c}\right] \bar{m}-\left[L_{2}^{T} R_{12} L_{2}^{T} R_{2}\right] \delta+\frac{1}{2} B^{T} \lambda_{1}+\frac{1}{2} D^{T} B_{c}^{T} \lambda_{2}} \\
\frac{1}{2} \lambda_{2}
\end{array}\right]=0,  \tag{4.4}\\
& \frac{\partial \mathcal{L}}{\partial A_{c}}=P_{12}^{r} Q_{12}+P_{2} Q_{2}+\frac{1}{2} \lambda_{2} m_{c}^{r}=0,  \tag{4.5}\\
& \frac{\partial \mathcal{L}}{\partial B_{c}}=P_{12}^{T} V_{12}+P_{2} B_{c} V_{2}+\left(P_{12}^{T} Q_{1}+P_{2} Q_{12}^{T}\right) C^{T}+\left(P_{12}^{T} Q_{12}+P_{2} Q_{2}\right) C_{c}^{T} D^{T}+\frac{1}{2} \lambda_{2} m^{T} C^{T}+\frac{1}{2} \lambda_{2} \gamma_{2}^{T}+\frac{1}{2} \lambda_{2} \alpha^{T} D^{T}=0,  \tag{4.6}\\
& \frac{\partial \mathcal{L}}{\partial C_{c}}=L_{2}^{T} R_{12}^{T} L_{1} Q_{12}+L_{2}^{T} R_{2} L_{2} C_{c} Q_{2}+L_{2}^{T} R_{12}^{T} L_{1} m m_{c}^{T}+L_{2}^{T} R_{2} L_{2} C_{c} m_{c} m_{c}^{T}+B^{T}\left(P_{1} Q_{12}+P_{12} Q_{2}\right)+D^{T} B_{c}^{T}\left(P_{12}^{T} Q_{12}+P_{2} Q_{2}\right)-L_{2}^{T} R_{12}^{T} \delta_{1} m_{c}^{T} \\
& +L_{2}^{T} R_{2} L_{2} \alpha m_{c}^{T}-L_{2}^{T} R_{2} \delta_{2} m_{c}^{T}+\frac{1}{2} B^{T} \lambda_{1} m_{c}^{T}=0 . \tag{4.7}
\end{align*}
$$

Expanding (4.1) and (4.2) yields

$$
\begin{gather*}
0=A Q_{1}+Q_{1} A^{T}+V_{1}+B C_{c} Q_{12}^{T}+Q_{12} C_{r}^{T} B^{T}  \tag{4.8}\\
0=A Q_{12}+Q_{12} A_{c}^{T}+B C_{r} Q_{2}+Q_{1} C^{\top} B_{c}^{T}+V_{12} B_{c}^{T}+Q_{12} C_{c}^{T} D^{T} B_{c}^{r}  \tag{4.9}\\
0=A_{c} Q_{2}+Q_{2} A_{c}^{r}+B_{c} C Q_{12}+Q_{12}^{T} C^{T} B_{c}^{T}+B_{c} V_{2} B_{c}^{T} \\
+B_{c} D C_{c} Q_{2}+Q_{2} C_{c}^{T} D^{T} B_{c}^{T}  \tag{4.10}\\
0=A^{T} P_{12}+P_{12} A_{c}+C^{T} B_{c}^{T} P_{2}+P_{1} B C_{c}+L_{1}^{T} R_{12} L_{2} C_{c}+P_{12} B_{c} D C_{c}  \tag{4.11}\\
0=A_{c}^{T} P_{2}+P_{2} A_{c}+C_{c}^{T} B^{T} P_{12}+P_{12}^{T} B C_{c}+C_{c}^{T} L_{2}^{T} R_{2} L_{2} C_{c}  \tag{4.12}\\
+ \\
+C_{c}^{T} D^{T} B_{c}^{T} P_{2}+P_{2} B_{c} D C_{r} . \tag{4.13}
\end{gather*}
$$

Next, note that (4.4) implies that $\lambda_{2}=0$, and thus (4.5) can be written as

$$
\begin{equation*}
-P_{2}^{-1} P_{12}^{T} Q_{12} Q_{2}^{-1}=I_{n_{c}} \tag{4.14}
\end{equation*}
$$

The existence of $Q_{2}^{-1}$ and $P_{2}^{-1}$ follows from the fact that $\left(A_{c}, B_{c}, C_{c}\right)$ is minimal. See [6] for details. Now define the $n \times n$ matrices

$$
\begin{gathered}
Q \triangleq Q_{1}-Q_{12} Q_{2}^{-1} Q_{12}^{T} P \triangleq P_{1}-P_{12} P_{2}^{-1} P_{12}^{T} \\
Q \triangleq Q_{12} Q_{2}^{-1} Q_{12}^{T} P \triangleq P_{12} P_{2}^{-1} P_{12}^{T} \\
T \triangleq-Q_{12} Q_{2}^{-1} P_{2}^{-1} P_{12}^{T}
\end{gathered}
$$

and the $n_{c} \times n, n_{c} \times n_{c}$, and $n_{c} \times n$ matrices

$$
G \cong Q_{2}^{\prime} Q_{12}^{T}, \quad M \triangleq Q: P_{2}, \quad \Gamma \triangleq-P_{2}^{-1} P_{12}^{T}
$$

Note that $\tau=G^{T} \Gamma$. Clearly. $Q, P, Q$, and $P$ are symmetric and nonnegative definite.

Next note that with the above definitions, (4.14) is equivalent to (3.13) and that (3.12) holds. Hence. $T=G^{T} T$ is idempotent, i.e., $r^{2}=T$. Sylvester's inequality yields (3.23). Note also that

$$
\phi=\tau Q, \quad \beta=\beta_{\tau}
$$

The components of $Q$ and $\bar{P}$ can be written in terms of $Q, P, Q, P, G$, and $\Gamma$ as

$$
\begin{array}{ll}
Q_{1}=Q+Q, \quad P_{1}=P+P \\
Q_{12}=Q \Gamma^{\top}, & P_{12}=-\beta G^{T} \\
Q_{2}=\Gamma Q \Gamma^{T} & P_{2}=G P G^{T}
\end{array}
$$

The expressions (3.16) and (3.17) follow from (4.6) and (4.7) by using the $n_{c}$ and $n$ components of (4.4), respectively, and the above identities. Next, computing either $\Gamma(4.9)-(4.10)$ or $G(4.12)+(4.13)$ yields (3.15).

Substituting this expression for $\boldsymbol{A}_{c}$ into (4.8)-(4.13) it follows that (4.10) $=\Gamma(4.9)$ and $(4.13)=G(4.12)$. Thus. (4.10) and (4.13) are superfluous and can be omitted. Next. using (4.8) $+G T(4.9) G-(4.9) G$ $[(4.9) G]^{T}$ and $G^{T} T(4.9) G-(4.9) G-[(4.9) G]^{T}$ yields (3.19) and (3.21). Using (4.11) $+\Gamma^{T} G(4.12) \Gamma-(4.12) \Gamma-[(4.12) \Gamma]^{T}$ and $\Gamma^{T} G(4.12) \Gamma-(4.12) \Gamma-[(4.12) \Gamma]^{T}$ yields (3.20) and (3.22).

To obrain (3.18) note that (4.4) can be rewritten as

$$
\begin{equation*}
R_{2} \tilde{\alpha}+\bar{R}_{12} \tilde{m}-\tilde{N} \delta+\frac{1}{2} B^{r} \bar{\lambda}=0 \tag{4.15}
\end{equation*}
$$

Next, note that (4.3) is equivalent to

$$
\begin{equation*}
\frac{1}{2} \bar{\lambda}=-\bar{A}-T \bar{R} \bar{m}+\bar{A}-T \bar{S} \bar{S}-\bar{A}-\tau \bar{R}_{12}^{T} \bar{\alpha} . \tag{4.16}
\end{equation*}
$$

Substituting (4.16) into (4.15) yields

Next, note that (3.10) is equivalent to

$$
\begin{equation*}
\bar{m}=-\bar{A}^{-1} \dot{B} \tilde{\alpha}-\bar{A}^{-1} \dot{\gamma} \tag{4.18}
\end{equation*}
$$

Now, substituting (4.18) into (4.17) yields

$$
\begin{align*}
& +\left(\bar{R}_{12}-B^{\top} \bar{A}-T \bar{R}\right) \bar{A}^{-1} \dot{\gamma} \tag{4.20}
\end{align*}
$$

Finally, note that the coefficient of $\bar{\alpha}$ in $(4.20)$ is equivalent to $\Omega$ and thus (4.20) yields (3.18).

## V. CONClUDING REMARKS

The results of the present note can be combined with the results of [1] to obtain nonstrictly proper controllers leading to a gencralization of [7]. Current research is focused on extending the results of the present note to larger classes of command and disturbance signals.

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## A Study of Controllability and Time-Optimal Control of a Robot Model with Drive Train Compliances and Actuator Dynamics

A. AILON AND G. LANGHOLZ


#### Abstract

The problems of robot controllability and time-optimal control where drive train compliances and actuator dynamics are incorporated in the mathematical model is the subject of this note. This study demonstrates the conditions that ensure the existence of a timeoptimal control, and establishes controllability of the augmented model (robot and actuator) in open- and closed-loop form. This note describes a procedure for the derivation of easily computable functional inequalities which represent upper bounds on the norm of the augmented system's time response.


## I. INTRODUCTION

To obtain the control strategy of mechanical manipulators, various control schemes are presented in the available literature. A few examples are resolved control [1], inverse problems technique [2], and resolved acceleration control [3]. In most cases, the control scheme involves the computation of the appropriate generalized forces by the equation

$$
H(\theta) \theta^{(2)}+K(\theta, \theta)+R(\theta)=q
$$

where $\theta$ and $q$ are the vectors of the generalized coordinates and forces, respectively. $H$ is the moment of inertia matrix. $K$ is a vector specifying centrifugal and Coriolis effects, and $\boldsymbol{R}$ is a vector specifying gravitational effects.

In much of the literature the actuators providing the drive torques are modeled as pure torque sources. However, this approach is in most cases a simplification of the realistic models of the system $[4]-[8]$.

The objective of this note is to study controllability and to investigate the conditions which ensure the existence of a control function that transfers the augmented model of the mechanical system, the actuator's dy namics. and the drive train's compliances. from a given initial position to a desired target in a minimum time The $m$ on: and the approach are useful for the design of a linear contr...ter and can be used as a point of departure for a more gencral model of a robst arm.

## iI. The Mathematical Model.

The Lagrange formulation of a multilink mechanical system is given by

$$
\begin{equation*}
d\left(\partial L / \partial \dot{\theta}_{1}\right) / d t-\partial L / \partial \theta_{1}=q_{1} . \quad i=1,2 \cdots n \tag{1}
\end{equation*}
$$

Manuscnpt received April 18. 1985, revised July 11. 1986. April 10, 1987. and Seprember 23, 1987
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where $L=T-V . T$ and $V$ are the kinetic and potential energies of the system, respectively.

Let $p_{i}$ be the ith generalized momentum [9]. Using Legender's dual transformation

$$
\begin{equation*}
p_{i}=\partial L / \partial \dot{\theta}_{i}, \quad i=1,2, \cdots n \tag{2}
\end{equation*}
$$

Since $L$ is a quadratic function in $\dot{\theta}_{1}, \boldsymbol{p}$ is linear in for any given $\theta$, i.e.,

$$
\begin{equation*}
p_{i}=\sum_{j=1}^{n} a_{i j}(\theta) \dot{\theta}_{j}, \quad i=1,2, \cdots n \tag{3}
\end{equation*}
$$

with $a_{i j}(\theta)=\partial^{2} L / \partial \dot{\theta}_{i} \partial \dot{\theta}_{j}$,
The inertial matrix is $H=\left\{\partial^{2} L / \partial \dot{\theta}, \partial \dot{\theta}\right\}_{n \times n}$ with $\operatorname{det}(H)=h(\theta)>0$, $\forall \theta$, where det $(\cdot)$ is the determinant of $(\cdot)$. Now, from (1). (2), and (3) we have

$$
\begin{array}{ll}
\dot{p}_{1}=\partial L / \partial \theta_{1}+q_{1}, & i=1,2, \cdots n \\
\dot{\theta}_{1}=\sum b_{v}(\theta) p_{j}, & i=1,2, \cdots, n \tag{5}
\end{array}
$$

Using (5) one obtains

$$
\begin{equation*}
\partial L / \partial \theta_{s}=\left[\sum_{i=1}^{n} \sum_{i=1}^{n} c_{s y}(\theta) p_{1} p_{k}\right] /(\operatorname{det}(H))^{2} . \tag{6}
\end{equation*}
$$

Equations (4)-(6) constitute the state equations of the $n$-link mechanical system which can be written as

$$
\begin{equation*}
\dot{z}(t)=F(z(t))+B q(t), z\left(t_{0}\right)=z_{0} \tag{7}
\end{equation*}
$$

where the vectors $z=\left[p^{\top} \theta^{\top}\right]^{T}, q=\left[q_{1} q_{2} \cdots q_{n}\right]^{\top}$, and $F=\left[F_{1} F_{2} \cdots\right.$ $\left.F_{2 n}\right]^{r}$ are in Euclidean vector space with the usual norm $\|z\|^{2}=\sum_{i=1}^{2 n}$ $\left(z_{1}\right)^{2}$. We also have

$$
\begin{aligned}
F_{s} & =\left[\sum_{i=1}^{n} \sum_{i=1}^{n} c_{* 1} p_{1} p_{j}\right] /[\operatorname{det}(H)]^{2}, \quad s=1,2, \cdots, n \\
& =\left[\sum_{i=1}^{n} d_{s,} p_{1}\right] / \operatorname{det}(H) . \quad s=n+1, n+2, \cdots, 2 n
\end{aligned}
$$

and $B=\left[\frac{1}{8}\right]$. where $I$ is the $n \times n$ identity matrix.
As an example, the exact equations for the two-link mechanical system which is confined to move in the vertical plane are given by

$$
\begin{align*}
& \dot{p}_{1}= {\left[p_{1} p_{2} l_{1} l_{c 2} m_{2} \sin \left(\theta_{2}-\theta_{1}\right) \operatorname{det}(H)-\left[0.5 p_{1}^{2} I_{2}+0.5 p_{2}^{2}\left(I_{1}+m_{2} I_{1}\right)\right.\right.} \\
&\left.\left.-p_{1} p_{2} E\right] 2 E l_{1} l_{c 2} m_{2} \sin \left(\theta_{2}-\theta_{1}\right)\right] /(\operatorname{det}(H))^{2} \\
&-\left(m_{1} g l_{1}+m_{2} g l_{1}\right) \sin \theta_{1}+q_{1}=F_{1}\left(p_{1}, p_{2}, \theta_{1}, \theta_{2}\right)+q_{1} \\
& \dot{p}_{2}=-\left[p_{1} p_{2} l_{1} l_{c 2} m_{2} \sin \left(\theta_{2}-\theta_{1}\right) \operatorname{det}(H)-\left\{0.5 p_{1}^{2} I_{2}+0.5 p_{2}^{2}\left(I_{1}+m_{2} I_{1}^{2}\right)\right.\right. \\
&\left.\left.-p_{1} p_{2} E\right] 2 E l_{1} l_{c 2} m_{2} \sin \left(\theta_{2}-\theta_{1}\right)\right] /(\operatorname{det}(H))^{2}-m_{2} g l_{e 2} \sin \theta_{2}+q_{2} \\
&= F_{2}\left(p_{1}, p_{2}, \theta_{1}, \theta_{2}\right)+q_{2} \\
& \dot{\theta}_{1}=\left(p_{1} I_{2}-p_{2} E\right) ; \operatorname{det}(H)=F_{3}\left(p_{1}, p_{2}, \theta_{1}, \theta_{2}\right) \\
& \dot{\theta}_{2}=\left(p_{2}\left(I_{1}+m_{2} I_{1}^{2}\right)-p_{1} E\right) / \operatorname{det}(H)=F_{0}\left(p_{1}, p_{2}, \theta_{1}, \theta_{2}\right) \tag{8}
\end{align*}
$$

where $E=l_{1} l_{c 2} m_{2} \cos \left(\theta_{2}-\theta_{1}\right), m_{1}$ and $l_{1}$ are the mass and the length of the ith link, respectively, $I_{1}$ is the moment of inertia of the ith link with respect to the ith joint, and $l_{c i}$ is the distance from the $i$ th joint to the center of gravity of the ith link.

The term det $(\boldsymbol{H})$ is a trigonometric function of, and periodical in. $\theta_{1}$. This function attains its minimum in the interval $0 \leq \theta_{i} \leq 2 \pi, i=1,2$; $\cdots n$, and therefore

$$
\begin{equation*}
\operatorname{det}(H) \geq k>0, \quad \forall z \in R^{2 n} \tag{9}
\end{equation*}
$$

We turn now to the dynamics of the robot's drivers. The robot is

# Optimal Reduced-Order State Estimation 

for
Unstable Plants

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#### Abstract

The problem of optimal reduced-order steady-state state estimation is considered for the case in which the plant has unstable poles. In contrast to the standard full-order estimation problem involving a single algebraic Riccati equation, the solution to the reduced-order problem involves one modified Riccati equation and one Lyapunov equation coupled by a projection matrix. This projection is completely distinct from the projection obtained in Bernstein and Hyland, 1985, for stable plants.


Supported in part by the Air Force Office of Scientific Research under contract F49620-86-C0002.

## 1. Introduction

It has recently been shown that optimal reduced-order, steady-state state estimators can be characterized by means of an algebraic system of equations consisting of one modified Riccati equation and two modified Lyapunov equations coupled by a projection matrix. The solution given in Bernstein and Hyland, 1985, however, was confined to problems in which the plant is asymptotically stable, while in practice it is often necessary to obtain estimators for plants with unstable modes. The purpose of the present paper is to obtain results similar to those of Bernstein and Hyland, 1985, for unstable plants.

Intuitively, it is clear that finite, steady-state state-estimation error for unstable plants is achievable only when the estimator retains, or duplicates in some sense, the unstable modes. Roughly speaking, the solution given in Bernstein and Hyland, 1985, is inapplicable to the unstable problem for the simple reason that the range of the projection matrix may not fully encompass the unstable subspace. Hence, in the present paper we derive a new reduced-order solution which is constrained to estimate all of the unstable states. Specifically, for a plant with an unstable subspace of dimension $n_{u}$, we characterize the optimal estimator of order $n_{u}$ which observes all of the unstable states.

As in Bernstein and Hyland, 1985, the solution is given in terms of an oblique projection (denoted in the present paper by $\mu$ ) which characterizes the optimal estimator gains. Again in contrast to the lone observer Riccati equation of the standard full-order theory, the optimal reducedorder estimator gains for an unstable plant are given by an algebraic system which, in the present case, consists of one modified Riccati equation and one Lyapunov equation coupled by the projection matrix $\mu$.

It is important to stress that the solution derived in the present paper is fundamentally different " from the solution obtained in Bernstein and Hyland, 1985, for two reasons. First, the estimator obtained in Bernstein and Hyland, 1985, was characterized by three matrix equations (in variables $Q, \hat{Q}$, and $\hat{P}$ ) while the solution obtained herein involves twn matrix equations (in variables $Q$ and $P$ ). And, second, since the projecticin $\mu$ arising in the present paper depends upon $P$, it is completely distinct from the projection $\tau$ appearing in Bernstein and Hyland, 1985, which depends upon $\hat{Q}$ and $\hat{P}$. Hence the results of the present paper neither generalize, nor are a special case of, the results of Bernstein and Hyland, 1985.

In applying the results of the present paper we note that the solution is applicable to problems in which the unstable subspace also includes additional stable modes. Indeed, the only constraint in applying the theory is that the unstable subspace include all of the unstable poles. To clarify this point (see Sections 2 and 3 for notation) we note that all unstable poles of $A$ must be contained in $A_{u}$, but $A_{u}$ may also contain an arbitrary number of selected stable poles. Thus, the estimator derived in the present paper can be viewed as a subspace-constrained observer-estimator.

Finally, the result given herein is only a partial solution to the reduced-order estimation problem. Specifically, a reduced-order estimator which includes all of the unstable modes and optimal combinations of a fixed number of stable modes should involve both projections $\tau$ and $\mu$ and four matrix equations in variables $Q, P, \hat{Q}$, and $\hat{P}$. When the result is specialized to the full-order case we expect the two projections to merge to form the identity and the four matrix equations to collapse to the single observer Riccati equation. A third projection $\nu$ due to singular measurement noise and static estimation can also be incorporated (Haddad and Bernstein, 1987, Halevi, 1988). This general solution remains the subject of current research.

After introducing notation, we consider the reduced-order estimation problem for continuoustime plants in Section 2. In Section 3 the corresponding discrete-time problem is considered. For stable plants the reduced-order discrete-time solution was given in Bernstein, Davis, and Hyland, 1986.

## Notation and Definitions

Note: All matrices have real entries

| $\mathbb{R}, \mathbb{R}^{r \times}, \mathbb{R}^{r}, \mathbb{E}$ | real numbers, $r \times s$ real matrices, $\mathbb{R}^{r \times 1}$, expected value |
| :--- | :--- |
| $I_{n},()^{T}, 0_{r \times e}, 0_{r}$ | $n \times n$ identity matrix, transpose, $r \times s$ zero matrix, $0_{r \times r}$ |
| $n, \ell, n_{e}, n_{u}, n_{e}, q$ | positive integers |
| $x, y, x_{e}, x_{u}, x_{e}, y_{e}$ | $n, \ell, n_{e}, n_{u}, n_{e}, q$-dimensional vectors |
| $A, C$ | $n \times n, \ell \times n$ matrices |
| $A_{u}, A_{u e}, A_{e}$ | $n_{u} \times n_{u}, n_{u} \times n_{e}, n_{e} \times n_{s}$ matrices |
| $C_{u}, C_{e}$ | $\ell \times n_{u}, \ell \times n_{s}$ matrices |
| $L, L_{u}, L_{e}$ | $q \times n, q \times n_{u}, q \times n_{e}$ matrices |
| $R$ | $q \times q$ positive-definite matrix |
| $A_{e}, B_{e}, C_{e}, D_{e}$ | $n_{e} \times n_{e}, n_{e} \times \ell, q \times n_{e}, q \times \ell$ matrices |
| $t, k$ | $t \in[0, \infty)$, discrete-time index $1,2,3, \ldots$ |

$\tilde{A}$
$w_{1}(\cdot), w_{2}(\cdot)$
$V_{1}$
$V_{2}$
$V_{12}$
$\tilde{w}(\cdot)$
$\tilde{V}$
$A-\left[\begin{array}{c}I_{n_{e}} \\ O_{n_{e} \times n_{e}}\end{array}\right] B_{e} C$
$n$, $\ell$-dimensional continuous-time or discrete-time white noise processes
$n \times n$ nonnegative-definite intensity or covariance of $w_{1}(\cdot)$
$\ell \times \ell$ positive-definite intensity or covariance of $w_{2}(\cdot)$
$n \times \ell$ cross intensity or cross-covariance of $w_{1}(\cdot), w_{2}(\cdot)$
$w_{1}(\cdot)-\left[\begin{array}{c}I_{n_{4}} \\ 0_{n, \times n_{4}}\end{array}\right] B_{\mathrm{e}} w_{2}(\cdot)$
$V_{1}-V_{12} B_{e}^{T}\left[\begin{array}{ll}I_{n_{\psi}} & 0_{n_{u} \times n_{\rho}}\end{array}\right]-\left[\begin{array}{c}I_{n_{\psi}} \\ 0_{n_{v} \times n_{u}}\end{array}\right] B_{\varepsilon} V_{12}^{T}+\left[\begin{array}{c}I_{n_{u}} \\ 0_{n_{\varepsilon} \times n_{u}}\end{array}\right] B_{\varepsilon} V_{2} B_{e}^{T}\left[\begin{array}{ll}I_{n_{u}} & 0_{n_{u} \times n_{e}}\end{array}\right]$

## 2. Problem Statement and Main Theorem

Reduced-Order State-Estimation Problem. Given the $\boldsymbol{n}$ th-order observed system

$$
\begin{align*}
& \dot{x}(t)=A x(t)+w_{1}(t),  \tag{2.1}\\
& y(t)=C x(t)+w_{2}(t), \tag{2.2}
\end{align*}
$$

design an $n_{\varepsilon}$ th-order state estimator

$$
\begin{gather*}
\dot{x}_{e}(t)=A_{e} x_{e}(t)+B_{e} y(t),  \tag{2.3}\\
y_{e}(t)=C_{e} x_{e}(t), \tag{2.4}
\end{gather*}
$$

which minimizes the state-estimation error criterion

$$
\begin{equation*}
J\left(A_{e}, B_{e}, C_{e}\right) \triangleq \lim _{t \rightarrow \infty} \mathbb{E}\left[L x(t)-y_{e}(t)\right]^{T} R\left[L x(t)-y_{e}(t)\right] \tag{2.5}
\end{equation*}
$$

In this formulation the plant is partitioned into possibly unstable and stable subsystems. Thus, letting $x(t)=\left[x_{u}^{T}(t), x_{i}^{T}(t)\right]^{T}$ and $w_{1}(t)=\left[w_{1 u}^{T}(t), w_{1,}^{T}(t)\right]^{T},(2.1)$ can be written as

$$
\left[\begin{array}{l}
\dot{x}_{u}(t)  \tag{2.6}\\
\dot{x}_{s}(t)
\end{array}\right]=\left[\begin{array}{cc}
A_{u} & A_{u s} \\
0_{n_{\bullet} \times n_{u}} & A_{s}
\end{array}\right]\left[\begin{array}{l}
x_{u}(t) \\
x_{s}(t)
\end{array}\right]+\left[\begin{array}{l}
w_{1 u}(t) \\
w_{1!}(t)
\end{array}\right],
$$

where $A_{u} \in \mathbb{R}^{n_{u} \times n_{u}}$ is possibly unstable, $A_{s} \in \mathbb{R}^{n_{0} \times n_{0}}$ is asymptotically stable, and the measurement equation (2.2) becomes

$$
y(t)=\left[C_{u} C_{0}\right]\left[\begin{array}{l}
x_{u}(t)  \tag{2.7}\\
x_{s}(t)
\end{array}\right]+w_{2}(t) .
$$

Furthermore, the matrix $L$, which is partitioned as

$$
L=\left[\begin{array}{ll}
L_{u} & L_{d} \tag{2.8}
\end{array}\right]
$$

identifies the states or linear combinations of states whose estimates are desired. The dimension $n_{e}$ of the estimator state $x_{c}$ is fixed to be equal to the order of the unstable part of the system, i.e., $n_{c}=n_{u}$. Thus, the goal of the Reduced-Order State-Estimation Problem is to design an estimator of order $n_{u}$ which yields quadratically optimal estimates of specified linear combinations of states of the system. As mentioned in Section 1, $A_{u}$ includes all unstable modes of $A$ as well as an arbitrary number of selected stable modes of $A$.

Since $A_{u}$ may contain unstable modes, define the error state $z(t) \triangleq x_{u}(t)-x_{e}(t)$ satisfying

$$
\begin{equation*}
\dot{i}(t)=\left(A_{u}-B_{e} C_{u}\right) x_{u}(t)-A_{e} x_{e}(t)+\left(A_{u s}-B_{e} C_{e}\right) x_{s}(t)+w_{1 u}(t)-B_{e} w_{2}(t) . \tag{2.9}
\end{equation*}
$$

Note that the explicit dependence of the error states $z(t)$ on the unstable states $x_{u}(t)$ can be eliminated by constraining

$$
\begin{equation*}
A_{\mathrm{e}}=A_{u}-B_{\mathrm{e}} C_{u} \tag{2.10}
\end{equation*}
$$

so that (2.9) becomes

$$
\begin{equation*}
\dot{z}(t)=\left(A_{u}-B_{e} C_{u}\right) z(t)+\left(A_{u e}-B_{e} C_{e}\right) x_{e}(t)+w_{1 u}(t)-B_{e} w_{2}(t) . \tag{2.11}
\end{equation*}
$$

Similarly, the explicit dependence of the estimation error (2.5) on the unstable states $x_{u}(t)$ can be eliminated by setting

$$
\begin{equation*}
C_{e}=L_{u} . \tag{2.12}
\end{equation*}
$$

Now (2.9)-(2.11) yield

$$
\begin{equation*}
\dot{\tilde{x}}(t)=\tilde{A} \tilde{x}(t)+\tilde{w}(t), \tag{2.13}
\end{equation*}
$$

where

$$
\tilde{x}(t) \triangleq\left[\begin{array}{c}
z(t) \\
x_{s}(t)
\end{array}\right], \quad \tilde{A} \triangleq\left[\begin{array}{cc}
A_{u}-B_{e} C_{u} & A_{u s}-B_{e} C_{s} \\
0_{n_{i} \times n_{u}} & A_{\mathrm{e}}
\end{array}\right],
$$

and $\tilde{w}(t)$ and its intensity $\tilde{V}$ are given in Section 1.
To guarantee that $J$ is finite, consider the set of asymptotically stable reduced-order estimators

$$
S \triangleq\left\{\left(A_{e}, B_{e}, C_{e}\right): \quad A_{e}=A_{u}-B_{e} C_{u} \text { is asymptotically stable }\right\}
$$

so that $\tilde{A}$ is asymptotically stable. Of course, $S$ is nonempty if $\left(A_{u}, C_{u}\right)$ is detectable. Furthermore, for nondegeneracy we restrict our attention to the set of admissible estimators

$$
S^{+} \triangleq\left\{\left(A_{e}, B_{e}, C_{e}\right) \in S: \quad\left(A_{e}, C_{e}\right) \text { is observable }\right\}
$$

where $A_{8}$ and $C_{6}$ are given by（2．10）and（2．12）．Also，for arbitrary $Q \in \mathbb{R}^{n \times n}$ define the notation

$$
Q_{a} \triangleq Q C^{T}+V_{12}
$$

Theorem 2．1．Suppose $\left(A_{e}, B_{e}, C_{e}\right) \in S^{+}$solves the Reduced－Order State－Estimation Prob－ lem with constraints（2．10）and（2．12）．Then there exist $n \times n$ nonnegative－definite matrices $Q, P$ such that $A_{s}, B_{e}, C_{e}$ are given by

$$
\begin{align*}
& A_{\varepsilon}=\Phi\left(A-Q_{a} V_{2}^{-1} C\right) F^{T}  \tag{2.14}\\
& B_{\varepsilon}=\Phi Q_{a} V_{2}^{-1}  \tag{2.15}\\
& C_{e}=L F^{T} \tag{2.16}
\end{align*}
$$

and such that $Q, P$ satisfy

$$
\begin{align*}
& 0=A Q+Q A^{T}+V_{1}-Q_{a} V_{2}^{-1} Q_{a}^{T}+\mu_{\perp} Q_{a} V_{2}^{-1} Q_{a}^{T} \mu_{\perp}^{T}  \tag{2.17}\\
& 0=\left(A-\mu Q_{a} V_{2}^{-1} C\right)^{T} P+P\left(A-\mu Q_{a} V_{2}^{-1} C\right)+L^{T} R L \tag{2.18}
\end{align*}
$$

where

$$
\begin{align*}
& P=\left[\begin{array}{cc}
P_{u} & P_{u_{e}} \\
P_{u \varepsilon}^{T} & P_{\varepsilon}
\end{array}\right] \in \mathbb{R}^{\left(n_{\psi}+n_{0}\right) \times\left(n_{*}+n_{e}\right)},  \tag{2.19}\\
& F \triangleq\left[\begin{array}{ll}
I_{n_{\star}} & 0_{n_{u} \times n_{\star}}
\end{array}\right], \quad \Phi \triangleq\left[I_{n_{\star}} P_{u}^{-1} P_{u}\right],  \tag{2.20}\\
& \mu \triangleq F^{T} \Phi=\left[\begin{array}{cc}
I_{n_{凶}} & P_{u}^{-1} P_{u \bullet} \\
0_{n_{0} \times n_{u}} & 0_{n_{\varphi}}
\end{array}\right], \quad \mu_{\perp} \triangleq I_{n}-\mu . \tag{2.21}
\end{align*}
$$

Furthermore，the minimal cost is given by

$$
\begin{equation*}
J\left(A_{e}, B_{e}, C_{e}\right)=\operatorname{tr} Q L^{T} R L \tag{2.22}
\end{equation*}
$$

Proof．See Appendix A．
Remark 2．1．Note that since $\Phi F^{T}=I_{n_{凶}}$ the $n \times n$ matrix $\mu$ which couples the modified Riccati equation（2．17）and the Lyapunov equation（2．18）is idempotent，i．e．，$\mu^{2}=\mu$ ．Note also that rank $\mu=n_{u}$ ．This projection is completely distinct from the projection $\tau$ appearing in Bernstein and Hyland， 1985.

Remark 2.2. In the full-order case $n_{u}=n$, Theorem 2.1 corresponds to the standard steadystate Kalman filter result. To see this, formally set $\Phi=F=\mu=I_{n}$ and $\mu_{\perp}=0_{n}$ so that (2.18) is superfluous and (2.17) specializes to the standard observer Riccati equation.

Remark 2.3. Note that (2.14) and (2.16) are merely restatements of (2.10) and (2.12). Furthermore, (2.15) implies that $\tilde{A}=A-\mu Q_{a} V_{2}^{-1} C$ so that the coefficient of $P$ in (2.18) is asymptotically stable.

## 3. Discrete-Time Formulation

Discrete-Time Reduced-Order State-Estimation Problem. Given the nth-order observed system

$$
\begin{gather*}
x(k+1)=A x(k)+w_{1}(k)  \tag{3.1}\\
y(k)=C x(k)+w_{2}(k) \tag{3.2}
\end{gather*}
$$

design an $n_{e}$ th-order state estimator

$$
\begin{gather*}
x_{\epsilon}(k+1)=A_{e} x_{e}(k)+B_{e} y(k)  \tag{3.3}\\
y_{e}(k)=C_{e} x_{e}(k)+D_{e} y(k) \tag{3.4}
\end{gather*}
$$

which minimizes the discrete-time state-estimation error criterion

$$
\begin{equation*}
\hat{J}\left(A_{e}, B_{e}, C_{\varepsilon}, D_{\varepsilon}\right) \triangleq \lim _{k \rightarrow \infty} \mathbb{E}\left[L x(k)-y_{\epsilon}(k)\right]^{T} R\left[L x(k)-y_{e}(k)\right] \tag{3.5}
\end{equation*}
$$

Because of the discrete-time setting it is now possible as in Bernstein, Davis, and Hyland, 1986, to permit a static feedthrough term $D_{e}$ in the estimator design. The gain $D_{e}$ represents a static least squares estimator in conjunction with the dynamic estimator ( $A_{e}, B_{e}, C_{e}$ ).

As in the continuous-time case, the plant is partitioned into stable and possibly unstable subsystems according to (2.6). Furthermore, an error state $z(k) \triangleq x_{u}(k)-x_{e}(k)$ is defined, $A_{e}$ is constrained as in (2.10), and $C_{a}$ is constrained to be $L_{u}-D_{e} C_{u}$. Thus, the augmented system consisting of the error states $z(k)$ and the stable states $x_{d}(k)$ becomes

$$
\begin{equation*}
\tilde{x}(k+1)=\tilde{A} \tilde{x}(k)+\tilde{w}(k) \tag{3.6}
\end{equation*}
$$

where $\tilde{x}(k) \triangleq\left[z^{T}(k), x_{\theta}^{T}(k)\right]^{T}$.

To guarantee that $J$ is finite and to obtain closed-form expressions for the estimator gains we restrict our attention to the sets

$$
\begin{gathered}
\hat{S} \triangleq\left\{\left(A_{\mathrm{e}}, B_{\mathrm{e}}, C_{\mathrm{a}}, D_{\mathrm{e}}\right): \quad A_{\mathrm{e}}=A_{u}-B_{\mathrm{e}} C_{u} \text { is asymptotically stable }\right\} \\
\hat{S}^{+} \triangleq\left\{\left(A_{\mathrm{e}}, B_{\mathrm{e}}, C_{e}, D_{\mathrm{e}}\right) \in \hat{S}: \quad\left(A_{e}, C_{\mathrm{e}}\right) \text { is observable }\right\}
\end{gathered}
$$

Also, for arbitrary $Q \in \mathbb{R}^{n \times n}$ define the notation

$$
\hat{Q}_{a} \triangleq A Q C^{T}+V_{12}, \quad \hat{V}_{2} \triangleq V_{2}+C Q C^{T} .
$$

Theorem 3.1. Suppose $\left(A_{\varepsilon}, B_{\varepsilon}, C_{e}, D_{e}\right) \in \hat{S}^{+}$solves the Discrete-Time Reduced-Order StateEstimation Problem. Then there exist $n \times n$ nonnegative-definite $Q, P$ such that $A_{e}, B_{e}, C_{e}, D_{e}$ are given by

$$
\begin{align*}
& A_{e}=\Phi\left(A-\hat{Q}_{a} \hat{V}_{2}^{-1} C\right) F^{T},  \tag{3.7}\\
& B_{e}=\Phi \hat{Q}_{a} \hat{V}_{2}^{-1}  \tag{3.8}\\
& C_{e}=\left(L-D_{e} C\right) F^{T},  \tag{3.9}\\
& D_{e}=L Q C^{T} \hat{V}_{2}^{-1} \tag{3.10}
\end{align*}
$$

and such that $Q, P$ satisfy

$$
\begin{gather*}
Q=A Q A^{T}+V_{1}-\hat{Q}_{a} \hat{V}_{2}^{-1} \hat{Q}_{a}^{T}+\mu_{\perp} \hat{Q}_{a} \hat{V}_{2}^{-1} \hat{Q}_{a}^{T} \mu_{\perp}^{T}  \tag{3.11}\\
P=\left(A-\mu \hat{Q}_{a} \hat{V}_{2}^{-1} C\right)^{T} P\left(A-\mu \hat{Q}_{a} \hat{V}_{2}^{-1} C\right)+\left(L-D_{e} C\right)^{T} R\left(L-D_{e} C\right), \tag{3.12}
\end{gather*}
$$

where $F, \Phi, \mu$, and $\mu_{\perp}$ are defined by (2.19)-(2.21). Furthermore, the minimal cost is given by

$$
\begin{equation*}
\hat{J}\left(A_{e}, B_{e}, C_{e}, D_{e}\right)=\operatorname{tr}\left[\left(L Q L^{T}-D_{e} V_{2} D_{e}^{T}\right) R\right] \tag{3.13}
\end{equation*}
$$

Proof. See Appendix A.
Remark 3.1. If a strictly proper estimator is desired, then delete $D_{e}$ in (3.9), (3.12), and (3.13).

## Appendix A: Proof of Theorems 2.1 and 3.1

To analyze (2.13) define the second-moment matrix

$$
\begin{equation*}
Q(t) \triangleq \mathbb{E}\left[\tilde{x}(t) \tilde{x}^{T}(t)\right] \tag{A.1}
\end{equation*}
$$

which satisfies

$$
\begin{equation*}
\dot{Q}(t)=\tilde{A} Q(t)+Q(t) \tilde{A}^{T}+\tilde{V}, \quad t \geq 0 . \tag{A.2}
\end{equation*}
$$

Since $\left(A_{e}, B_{a}, C_{c}\right) \in S, \tilde{A}$ is asymptotically stable and

$$
Q \triangleq \lim _{t \rightarrow \infty} \mathbb{E}\left[\tilde{x}(t) \tilde{x}^{T}(t)\right]
$$

exists and satisfies

$$
\begin{equation*}
0=\tilde{A} Q+Q \tilde{A}^{T}+\tilde{V} . \tag{A.3}
\end{equation*}
$$

Next note that (2.5) can be written as

$$
\begin{equation*}
J\left(A_{e}, B_{e}, C_{e}\right)=\operatorname{tr} Q L^{T} R L . \tag{A.4}
\end{equation*}
$$

To minimize (A.4) over the open set $S^{+}$subject to the constraint (A.3), form the Lagrangian

$$
\begin{equation*}
\mathcal{L}\left(A_{a}, B_{a}, C_{c}, Q, P, \lambda\right) \triangleq \operatorname{tr}\left[\lambda Q L^{T} R L+\left(\tilde{A} Q+Q \tilde{A}^{T}+\tilde{V}\right) P\right] \tag{A.5}
\end{equation*}
$$

where the Lagrange multipliers $\lambda \geq 0$ and $P \in \mathbb{R}^{n \times n}$ are not both zero. Setting $\partial \mathcal{L} / \partial Q=0, \lambda=0$ implies $P=0$ since $\tilde{A}$ is asymptotically stable. Hence, without loss of generality set $\lambda=1$.

Now partition $n \times n P$ into $n_{u} \times n_{u}, n_{u} \times n_{d}$, and $n_{s} \times n_{s}$ subblocks as

$$
P=\left[\begin{array}{cc}
P_{u} & P_{u \bullet}  \tag{A.6}\\
P_{u \bullet}^{T} & P_{\bullet}
\end{array}\right]
$$

Thus the stationarity conditions are given by

$$
\begin{gather*}
\frac{\partial \mathcal{L}}{\partial Q}=\tilde{A}^{T} P+P \tilde{A}+L^{T} R L=0  \tag{A.7}\\
\frac{\partial \mathcal{L}}{\partial B_{z}}=P_{u} B_{\varepsilon} V_{z}-\left[P_{u} P_{u \varepsilon}\right]\left(Q C^{T}+V_{12}\right)=0 . \tag{A.8}
\end{gather*}
$$

Expanding the $n_{u} \times n_{u}$ subblock of (A.7) yields

$$
\begin{equation*}
0=\left(A_{u}-B_{e} C_{u}\right)^{T} P_{u}+P_{u}\left(A_{u}-B_{e} C_{u}\right)+L_{u}^{T} R L_{u} \tag{A.9}
\end{equation*}
$$

which, using (2.10) and (2.12), is equivalent to

$$
\begin{equation*}
0=A_{a}^{T} P_{u}+P_{u} A_{\varepsilon}+C_{e}^{T} R C_{\mathrm{e}} \tag{A.10}
\end{equation*}
$$

Thus, since $\left(A_{6}, B_{e}, C_{6}\right) \in S^{+},\left(A_{\varepsilon}, C_{6}\right)$ is observable and it follows from (A.10) that $P_{u}$ is positive definite. Since $P_{u}$ is thus invertible, define the $n_{u} \times n$ matrices

$$
F \triangleq\left[\begin{array}{ll}
I_{n_{u}} & 0_{n_{u} \times n_{\delta}} \tag{A.11}
\end{array}\right], \quad \Phi \triangleq\left[I_{n_{u}} P_{u}^{-1} P_{u s}\right],
$$

and the $n \times n$ matrix $\mu \triangleq F^{T} \Phi$. Note that since $\Phi F^{T}=I_{n_{u}}, \mu$ is idempotent, i.e., $\mu^{2}=\mu$.
Next note that (A.8) and (A.11) imply (2.15). Similarly, (2.14) is equivalent to (2.10) with $B_{e}$ given by (2.15). Finally, (2.16) is a restatement of (2.12). Now, using the expression for $B_{e}, \tilde{A}$ and $\tilde{V}$ become

$$
\begin{gather*}
\tilde{A}=A-\mu Q_{a} V_{2}^{-1} C,  \tag{A.12}\\
\tilde{V}=V_{1}-V_{12} V_{2}^{-1} Q_{a}^{T} \mu^{T}-\mu Q_{a} V_{2}^{-1} V_{12}^{T}+\mu Q_{a} V_{2}^{-1} Q_{a}^{T} \mu^{T} . \tag{A.13}
\end{gather*}
$$

Finally, (2.17) and (2.18) follow from (A.3) and (A.7) using (A.12) and (A.13).
For the discrete-time problem define the second-moment matrix

$$
Q(k) \triangleq \mathbb{E}\left[\tilde{x}(k) \tilde{x}^{T}(k)\right],
$$

which satisfies

$$
\begin{equation*}
Q(k+1)=\tilde{A} Q(k) \tilde{A}^{r}+\tilde{V} . \tag{A.14}
\end{equation*}
$$

Since $\tilde{A}$ is asymptotically stable,

$$
Q \triangleq \lim _{k \rightarrow \infty} \mathbb{E}\left[\tilde{x}(k) \tilde{x}^{T}(k)\right]
$$

exists and satisfies

$$
\begin{equation*}
Q=\tilde{A} Q \tilde{A}^{T}+\tilde{V} . \tag{A.15}
\end{equation*}
$$

The remainder of the proof follows as above for the continuous-time case.
Acknowledgment. We wish to thank David C. Hyland for several helpful suggestions.

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# The Optimal Projection Equations for Reduced-Order, Discrete-Time Modeling, Estimation, and Control 

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#### Abstract

The optimal projection equations derived previously for reduced-order, continuous-time modeling. estimation. and control are developed for the discrete-time case. The design equations are presented in a concise. unified manner to facilitate their accessibility for the development of numerical algorithms for practical applications. As in the continuous-time case, the standard Kalman filter and linear-quadratic-Gaussian results are immediately obtained as special cases of the estimation and control results.


## Nomenclature

| $A, B, C=n \times n, n \times m, l \times n$ matrices |  |
| :---: | :---: |
| $A_{m,}, B_{m,}, C_{m}$ | $=n_{m} \times n_{m}, n_{m} \times m, l \times n_{m}$ matrices |
| $A_{e}, B_{e}, C_{c}, D_{e}=n_{e} \times n_{e}, n_{e} \times \ell, p \times n_{e}, p \times \ell$ matrices |  |
| $A_{1}, B_{c}, C_{l}, D_{l}=n_{c} \times n_{c}, n_{c} \times \ell, m \times n_{c}, m \times \ell$ matrices |  |
| $E_{1}$ | $=$ matrux with unity in the ( $i, i$ ) position and zeros elsewhere |
| $\boldsymbol{E}$ | - expected value |
| $I_{\text {r }}$ | $=r \times r$ identity matrix |
| $\boldsymbol{k}$ | = discrete-time index 1,2,3,.. |
| $L$ | $=p \times n$ matrix |
| $n, m, l_{,} n_{m}$, |  |
| $n_{\sim}, n_{i}, p$ | $=$ positive integers, $1 \leq n_{m}, n_{c}, n_{c} \leq \boldsymbol{n}$ |
| $R, N, R_{2}=\ell \times \ell, p \times p, m \times m$ positive-definite matrices |  |
| $R_{1} \quad=n \times n$ nonnegative-definite matrix |  |
| $R_{12}$ | $=n \times m$ matrix such that $R_{1}-R_{12} R_{2}^{-1} R_{12}^{\top}$ |
| $\boldsymbol{R}, \boldsymbol{R}^{\prime \times}=$ real number |  |
| $\bar{R} \quad=\left[\begin{array}{ll}R_{1} & R_{12} \\ R_{12}^{T} & R_{2}\end{array}\right]$ |  |
|  |  |
| ir $Z \quad=$ trace of square matix $Z$ |  |
| $u_{V}, y, y_{\text {l }} \quad=m-, t$ - $p$-dimensional vectors |  |
|  | $=m \times m$ positive-definite covariance of $w$ |
| $\nu_{1} \quad=n \times n$ nonnegative-definite covariance of $w_{1}$ | $=n \times n$ nonnegative-definite covariance of $w_{i}$ |
| $V_{2} \quad=\ell \times \ell$ positive-definite covariance of $n_{2}$ |  |
| $V_{1}$ : $\quad=n \times \rho$ cross-covariance of $w_{1}, w_{2}$ |  |
|  |  |
| $x_{1}, x_{m}, x_{r}, x_{1}=n_{-}, n_{m} \cdot n_{e}, n_{c}$-dimensional vectors |  |
| $Z_{0,1}=(i, j)$ element of matrix $Z$ |  |
| $Z^{\dagger} \quad=$ transpose of vector or matrix $Z$ |  |
| $Z^{-T}=\left(Z^{T}\right)^{-1}$ or $\left(Z^{-1}\right)^{T}$ |  |
| П(*) | $=\Psi E \Psi^{-1}$ (unit-rank eigenprojection ${ }^{47}$ ) |
| $p(Z)$ | = rank of matrix $Z$ |

[^33]
## 1. Introduction

IN a recent series of papers ${ }^{1 / 11}$ it has been shown that the first-order necessary conditions for quadratically optimal, continuous-time. reduced-order modeling. estimation, and control can be transformed into coupled systems of two, three. and four matrix equations, respectively. This coupling, due to the presence of an oblique projection (idempotent matrix). arises as a rigorous consequence of optimality, bence suggesting the name optimal projection. For the estimation and control problems, this formulation provides a direct generalization of classical steady-state Kalman filter and linear-quadratic-Gaussian (LQG) control theory. In the full-order case the projection becomes the identity matrix, the additional two modified Lyapunor equations drop out, and the remaining modified Ricrati equations become the usual Riccati equations.

Coupling via the optimal projection supports the view that sequential reduced-order design procedures consisting of either 1) model reduction followed by estimator (controller) design or 2 ) estimator (controller) design followed by estimator (controller) reduction are generally not optimal. Furthermore, for the control problem the coupled structure of the equations yields the further insight that in the reduced-order case there is no longer separation between the operations of state estimation and state-estimate feedback. i.e., the certainty equivalence principle breaks down

For practical applications. the optimal projection equations permit the development of altemative numerical algorithms that operate through successive iteration of the optimal projection ${ }^{4.6}$ rather than by gradient search techniques. ${ }^{12}$ By recognizing that each local extremalcorresponds to $n$, possible choices out of $n$ rank-l eigenprojections of the product of a pair of pseudogramians, it is possible to efficiently idenufy the global minimum ${ }^{10}$ This idea is philosophically similar to Skelton's component-cost analysis. ${ }^{15-15}$

The purpose of the present paper is to develop the optimal projection equations for reduced-order modeling, estimation, and control in the discrete-lime case. Since the underlying theory has been discussed previously, ${ }^{1.11}$ the presentation berein is geared toward a clear and concise statement of the main results to facilitate numerical developments and practical application. For example, by expressing the optimal projection in terms of eigenprojections, a variety of novel algorithms are immediately suggested. For illustrative purposes we apply the results on reduced-order state estimation to a third-order problein to obtain reduced-order estimators and the results on reduced-order dynamic compensation to a tenth-order problem to obtain reduced-order controllers.

Because of the discrete-time setting it is now possible to permit static feedthrough gains in the estimator and controller designs. As previously noted,' nonsingular control weighting and measurement noise in the continuous-time case permit only a purely dynamic (strictly proper) controller. Note that this is precisely the case in continuous-time LQG theory, which always yields a strictly proper feedback controller. The static gains in the discrete-time estimator problem permit simultaneous, unified treatment of nondynamic least-squares estimation along with dynamic (Kalman filter-type) estimation.

The references include a representative sampling of papers on quadratically optimal reduced-order modeling. ${ }^{\text {B-25 }}$ estimation. ${ }^{3533}$ and control, ${ }^{34-46}$ along with closely related approaches. For emphasis on the discrete-time problem, see Refs. 18, 30, 41, 42, 44, and 45.

## II. Problem Statement and Main Results

We now state the reduced-order modeling. estimation, and control problems. The object of the model-reduction problem is to determine a model of reduced state-space dimension whose steady-state response to white noise inputs (or, equivalently, impulse response) best approximates, in a quadratic (least-squares) sense, the response of a given high-order system. In the reduction process the order of the reduced model is fixed and the optimization is performed over the wodel parameters.

## Reduced-Order Modeling Problem

Given the model

$$
\begin{gather*}
x(k+1)=A x(k)+B w(k)  \tag{1}\\
y(k)=C x(k) \tag{2}
\end{gather*}
$$

design a redv' ed-order model

$$
\begin{gather*}
x_{m}(k+1)=A_{m} x_{m}(k)+B_{m} w(k)  \tag{3}\\
y_{m}(k)=C_{m} x(k) \tag{4}
\end{gather*}
$$

which minimizes the model-reduction criterion

$$
\begin{align*}
& J_{m}\left(A_{m} \cdot B_{m}, C_{m}\right) \\
& \quad \triangleq \lim _{k \rightarrow x} E\left[y_{m}(k)-y(k)\right] R\left[y_{m}(k)-y(k)\right] \tag{5}
\end{align*}
$$

The goal of the reduced-order state-estimation problem is to design an estimator of given order which yields quadratically optimal (least squares) estimates of specified linear combinations $L x$ of states $x$ In practice, the order of the estimator may be determined by implementation constraints, such as reaj-time computing capability. Note that the feedthrough term $D$, permits the utilization of a static least-squares estimator in conjunction with the d.namic estimator ( $A_{e}, B_{e}, C_{e}$ ).

## Reduced-Order State-Estimation Problem

Given the observed system

$$
\begin{gather*}
x(k+1)=A x(k)+w_{1}(k)  \tag{6}\\
y(k)=C x(k)+w_{2}(k) \tag{7}
\end{gather*}
$$

design a reduced-order state estimator

$$
\begin{gather*}
x_{r}(k+1)=A_{e} x_{r}(k)+B_{e} y(k)  \tag{8}\\
y_{r}(k)=C_{e} x_{r}(k)+D_{e} y(k) \tag{9}
\end{gather*}
$$

which ummimes the staie-estumation criterion

$$
\begin{align*}
& J_{f}\left(A_{,}, B_{e}, C_{e}, D_{e}\right) \\
& \quad \& \lim _{k \rightarrow \infty} E\left[v_{e}(k)-L x(k)\right]^{r} N\left[y_{e}(k)-L x(k)\right] \tag{10}
\end{align*}
$$

For the fixed-order dynamic-compensation problem, a static feedthrough term is included, i.e., the controller may be nonstrictly proper.

## Reduced-Order Dynamic-Compensation Problem

Given the controlled system

$$
\begin{gather*}
x(k+1)=A x(k)+B u(k)+w_{1}(k)  \tag{1i}\\
y(k)=C x(k)+w_{2}(k) \tag{12}
\end{gather*}
$$

design a reduced-order dynamic compensator

$$
\begin{gather*}
x_{t}(k+1)=A_{1} x_{l}(k)+B_{y} y(k)  \tag{13}\\
u(k)=C_{،} x_{l}(k)+D_{i} y(k) \tag{14}
\end{gather*}
$$

which minimizes the dynamic-compensation criterion

$$
\begin{gather*}
J_{1}\left(A_{1}, B_{1}, C_{1}, D_{1}\right) \triangleq \lim _{k \rightarrow \infty} E\left[x(k)^{T} R_{1} x(k)\right. \\
\left.+2 x(k)^{T} R_{12} u(k)+u(k)^{T} R_{2} u(k)\right] \tag{15}
\end{gather*}
$$

To guarantee that $J_{m}, J_{e}$, and $J_{1}$ are finite and independent of initial conditions, consideration is restricted to the following (open) sets. [A triple $(\hat{A}, \hat{B}, \hat{C})$ is min ral if $(\hat{A}, \hat{B})$ is controllable and ( $\hat{A}, \hat{C}$ ) is observable.]

$$
\mathscr{S}_{m_{1}} \triangleq\left\{\left(A_{m}, B_{m 1}, C_{m}\right):\right.
$$

$A_{m}$ is stable and $\left(A_{m}, B_{m}, C_{m}\right)$ is minimal $\}$

$$
\mathscr{S}_{e} \triangleq\left\{\left(A_{e}, B_{e}, C_{e}, D_{e}\right):\right.
$$

$$
\left.A_{e} \text { is stable and }\left(A_{e}, B_{e}, C_{e}\right) \text { is minimal }\right)
$$

$$
\mathscr{P}_{1} \triangleq\left\{\left(A_{1}, B_{1}, C_{1}, D_{1}\right):\right.
$$

$$
\left[\begin{array}{cc}
A+B D_{c} C & B C_{c} \\
B_{1} C & A_{c}
\end{array}\right] \text { is stable and }\left(A_{4}, B_{c}, C_{c}\right) \text { is minimal }
$$

 factorization lemma aill be needed for the main results
Lemma 2.1 Let $\tau \in \boldsymbol{R}^{n \times n}$. Then

$$
\begin{gather*}
\tau^{2}=\tau  \tag{16}\\
\rho(\tau)=n_{r} \tag{17}
\end{gather*}
$$

if, and only if, there exist $G . \Gamma \in R^{\prime \prime \times n}$ such that

$$
\begin{align*}
G^{T} \Gamma & =\tau  \tag{18}\\
\Gamma G^{T} & =I_{n}, \tag{19}
\end{align*}
$$

Furthermore, $G$ an:- are unique to a change of basis in $\boldsymbol{R}^{n}$.
Proof Sufficienc) is obvious. To prove necessity, first note that due to Eq. (16) the eigenvalues of + are cither 0 or 1 Further, it is easy to see that thas a diagonal Jordan canonical form Hence, the result follows from

$$
T=S\left[\begin{array}{cc}
I_{n}, & 0 \\
0 & 0
\end{array}\right] S^{-1}=G^{T} \Gamma
$$

where $G=\left[\phi^{\top} \quad 0\left|S^{\top} . I=\left|\phi^{-1} \quad 0\right| S^{-1}\right.\right.$, and $\phi \in R^{n, \times n}$. $\quad$. For convenience. call $G$ and $\Gamma$ satisfying Eqs. (18) and (19) a profectue facturisation of $\tau$. Furthermore, for $n \times n$ non-negative-definite matrices (i.e., symmetric matrices with nonnegative eigenvalues) $\mathcal{P}$ and $\mathscr{F}$. define the set of contragredi-
ently diagor alizing transformations

$$
\begin{aligned}
& \mathscr{S}(\mathscr{G}, \mathscr{P}) \\
& \quad\left\{\left\{\Psi \in R^{n \times n}: \Psi^{-1} \mathscr{P} \Psi^{-T} \text { and } \Psi^{T} \mathscr{P} \Psi \text { are diagonal }\right\}\right.
\end{aligned}
$$

It follows from Ref. 48, p. 123, Theorem 6.2.5, that $\mathscr{S}(2,9)$ is always nonempty. This set does not, however, have a unique element since basis rearrangements and sign transpositions may be incorporated into $\Psi$. Further nonuniqueness arises if 29 has repeated eigenvalues.

Theorem 2.1. Suppose $A$ is stable and $\left(A_{m}, B_{m}, C_{m}\right) \in \mathscr{S}_{m}$ solves the reduced-order modeling problem. Then there exist $n \times n$ nonnegative-definite matrices $\hat{Q}$ and $\hat{P}$ such that $A_{m}$, $B_{m}$, and $C_{m}$ are given by

$$
\begin{align*}
A_{m} & =\Gamma A G^{\top}  \tag{20}\\
B_{m} & =\Gamma B  \tag{21}\\
C_{m} & =C G^{T} \tag{22}
\end{align*}
$$

and such that $\hat{Q}$ and $\hat{P}$ satisfy

$$
\begin{align*}
& \hat{Q}=A \tau \hat{Q}_{\tau}^{\top} A^{\tau}+B V B^{T}  \tag{23}\\
& \hat{P}=A^{T} \tau^{\top} \hat{P}_{\tau} A+C^{\tau} R C \tag{24}
\end{align*}
$$

where

$$
\begin{equation*}
\tau \pm \sum_{i=1}^{n_{m}} \Pi_{i}(\Psi) \tag{25}
\end{equation*}
$$

for some $\Psi \in \mathscr{Q}(\hat{Q}, \hat{P})$ such that $\left(\Psi^{-1} \hat{Q} \hat{P} \Psi\right)_{\langle i n} \neq 0, i=$ $1, \ldots, n_{m}$, and some projective factorization $G, f$ of $\tau$. Furthermore, the minimal cost is given by

$$
\begin{equation*}
J_{m}\left(A_{m}, B_{m}, C_{m}\right)=\operatorname{tr}\left[\left(W_{c}-\tau \hat{Q}^{\top} r^{\top}\right) C^{\top} R C\right] \tag{26}
\end{equation*}
$$

where $W_{i}$ is the unique (nonnegative-definite) solution of

$$
W_{r}=A W_{r} A^{T}+B V B^{T}
$$

For convenience in stating the estimator result, define the notation

$$
\begin{aligned}
& \Sigma_{Q} \triangleq\left(A Q C^{\top}+V_{12}\right) \dot{V}_{2}^{-1}\left(A Q C^{\top}+V_{12}\right)^{\top} \\
& \Sigma_{P} \triangleq\left(L-D_{R} C\right)^{T} N\left(L-D_{R} C\right) \\
& A_{Q} \pm A-\left(A Q C^{\top}+V_{12}\right) \dot{V}_{2}^{-1} C \\
& \hat{V}_{2} \pm V_{2}+C Q C^{\top} \\
& T_{2} \triangleq I_{n}-T
\end{aligned}
$$

Theorem 2.2. Suppose $A$ is stable and ( $\left.A_{r}, B_{c}, C_{c}, D_{e}\right) \in \mathscr{S}_{e}$ solves the reduced-order state-estimation problem. Then there exist $n \times n$ nonnegative-definite matrices $Q, \hat{Q}$, and $\hat{P}$ such that $A_{e}, B_{e}, C_{e}$, and $D_{e}$ are given by

$$
\begin{align*}
& A_{e}=\Gamma\left[A-\left(A Q C^{T}+V_{12}\right) \dot{V}_{2}^{1} C\right] G^{T}  \tag{27}\\
& B_{c}=\Gamma\left(A Q C^{T}+V_{12}\right) \dot{V}_{2}^{1}  \tag{28}\\
& C_{e}=\left(L-D_{f} C\right) G^{T}  \tag{29}\\
& D_{e}=L Q C^{\top} \tilde{V}_{2}^{-1} \tag{30}
\end{align*}
$$

and such that $\boldsymbol{Q}, \hat{Q}$, and $\hat{P}$ satisfy

$$
\begin{align*}
Q & =A Q A^{T}-\left(A Q C^{T}+V_{12}\right) \dot{V}_{2}^{-1}\left(A Q C^{T}+V_{12}\right)^{T} \\
& +V_{1}+\tau_{1} \dot{Q}_{1}^{T}  \tag{31}\\
\hat{Q} & =A \tau \dot{Q}^{T} A^{T}+\Sigma_{Q}  \tag{32}\\
\hat{P} & =A_{Q}^{T} \tau r_{P} \tau A_{Q}+\Sigma_{P} \tag{33}
\end{align*}
$$

where

$$
\begin{equation*}
\tau \triangleq \sum_{i=1}^{n_{c}} \Pi_{i}(\Psi) \tag{34}
\end{equation*}
$$

for some $\Psi \in \mathscr{G}(\hat{Q}, \hat{P})$ such that $\left(\Psi^{-1} \hat{Q} \hat{P} \psi\right)_{(0,1)} \neq 0, i=$ $1, \ldots, n_{e}$, and some projective factorization $G, \Gamma$ of $\tau$. Furthermore, the minimal cost is given by

$$
\begin{equation*}
J_{e}\left(A_{e}, B_{e}, C_{e}, D_{e}\right)=\overleftarrow{t}\left[\left(L Q L^{T}-D_{e} \bar{V}_{2} D_{e}^{T}\right) N\right] \tag{35}
\end{equation*}
$$

For the control result, define the additional notation

$$
\begin{aligned}
& \Sigma_{Q}^{\prime} \triangleq\left(A Q C^{T}+V_{12}+B D_{c} \bar{V}_{2}\right) \tilde{V}_{2}^{1}\left(A Q C^{T}+V_{12}+B D_{c} \tilde{V}_{2}\right)^{T} \\
& \Sigma_{P}^{\prime} \triangleq\left(B^{T} P A+R_{12}^{T}+\tilde{R}_{2} D_{c} C\right)^{T_{R_{2}}^{-1}\left(B^{T} P A+R_{12}^{T}+\tilde{R}_{2} D_{c} C\right)} \\
& A_{P} \triangleq A-B \bar{R}_{2}^{-1}\left(B^{T} P A+R_{12}^{T}\right) \\
& \tilde{R}_{2} \triangleq R_{2}+B^{r} P B \\
& M \triangleq\left[\begin{array}{c}
I_{n} \\
D_{c} C
\end{array}\right], \quad \hat{M} \triangleq\left[\begin{array}{c}
I_{n} \\
-\tilde{R}_{2}^{-1}\left(B^{T} P A+R_{12}^{T}\right)
\end{array}\right]
\end{aligned}
$$

Theorem 2.3. Suppose $\left(A_{e}, B_{c}, C_{c}, D_{c}\right) \in \mathscr{S}_{c}$ solves the re-duced-order dynamic-compensation problem. Then there exist $n \times n$ nonnegative-definite matrices $\boldsymbol{Q}, P, \hat{Q}$, and $\hat{P}$ such that $A_{i}, B_{c}, C_{c}$, and $D_{c}$ are given by

$$
\begin{align*}
A_{1}= & \Gamma\left[A-\left(A Q C^{T}+V_{12}\right) \bar{V}_{2}^{-1} C\right. \\
& \left.-B \tilde{R}_{2}^{-1}\left(B^{T} P A+R_{12}^{T}\right)-B D_{c} C\right] G^{T}  \tag{36}\\
B_{1}= & \Gamma\left[\left(A Q C^{T}+V_{12}\right) \dot{V}_{2}^{-1}+B D_{1}\right]  \tag{37}\\
C_{1}= & -\left[\tilde{R}_{2}^{-1}\left(B^{T} P A+R_{12}^{T}\right)+D_{c} C\right] G^{T} \\
D_{c}= & -\tilde{R}_{2}^{-1}\left(B^{T} P A Q C^{T}+R_{12}^{T} Q C^{T}+B^{T} P V_{12}\right) \bar{V}_{2}^{-1} \tag{39}
\end{align*}
$$

and such that $Q, P, \hat{Q}$, and $\dot{P}$ satisty

$$
\begin{align*}
& Q=A Q A^{T}-\left(A Q C^{T}+V_{12}\right) \dot{V}_{2}^{-1}\left(A Q C^{T}+V_{12}\right)^{T} \\
& +V_{1}+\tau_{\perp} \hat{Q} \tau_{\perp}^{T}  \tag{40}\\
& P=A^{T} P A-\left(B^{T} P A+R_{12}^{T}\right)^{T} \bar{R}_{2}^{-1}\left(B^{T} P A+R_{12}^{T}\right) \\
& +R_{1}+\tau_{\perp}^{T} \hat{P}_{r_{\perp}}  \tag{41}\\
& \hat{\boldsymbol{Q}}=A_{P} \boldsymbol{T}_{\boldsymbol{Q}}{ }_{T}^{\top} A_{P}^{T}+\boldsymbol{\Sigma}_{\boldsymbol{Q}}^{\prime}  \tag{42}\\
& \hat{P}=A_{Q}^{T}{ }^{\top} T_{P_{\tau}} A_{\rho}+\Sigma_{P}^{\prime} \tag{43}
\end{align*}
$$

where

$$
\begin{equation*}
\tau \Perp \sum_{i=1}^{n_{i}} n_{i}(\Psi) \tag{44}
\end{equation*}
$$

for some $\Psi \in \mathscr{G}(\hat{Q}, \hat{P})$ such that $\left(\Psi^{-1} \hat{Q}^{\hat{P}} \Psi\right)_{(1,1)} \neq 0, i=$ $1, \ldots, n_{1}$, and some projective factorization $G, \Gamma$ of $r$. Further.
more, the minimal cost is given by

$$
\begin{equation*}
J_{1}\left(A_{1}, B_{1}, C_{1}, D_{1}\right)=\operatorname{tr}\left[\left(M Q M^{T}+\hat{M} r \grave{Q}_{\tau^{T}} \dot{M}^{T}\right) \tilde{R}\right] \tag{4S}
\end{equation*}
$$

Remark 2.1. To specialize the estimation and control results to the strictly proper (no-feedthrough) case, merely ignore Eqs. (30) and (39) and set $D_{c}=0$ and $D_{1}=0$ wherever they appear.

Remark 2.2. In the full-order cases $n_{c}=\boldsymbol{n}$ and $n_{\boldsymbol{c}}=\boldsymbol{n}$ in Theorems 2.2 and 2.3, the projection $\tau$ becomes the identity and Eqs. (32). (33). (42), and (43) play no role. In this case G and $\Gamma$ are also the identity. Specializing further to the purely dynamic case $D_{r}=0, D_{t}=0$ as in the previous remark yields the standard Kalman filter and LQG results.

Remark 2.3. As previously noted.? ${ }^{10.11}$ the indeterminacy in specifying the projective factorization $G$. $\Gamma$ satisfying Eqs. (18) and (19) corresponds to nothing more than an arbitrary choice of internal state-space basis for the design systems $\left(A_{m}, B_{m}, C_{m i}\right),\left(A_{i}, B_{e} C_{e}\right)$, and $\left(A_{1}, B_{i}, C_{i}\right)$.

Remark 2.4. Since $\hat{Q}$ and $\hat{P}$ are balanced by means of the transformation $\Psi \in \mathscr{D}(\hat{Q}, \hat{P})$, it follows that $\Psi^{-1} \hat{Q} \hat{P} \Psi$ is diagonal. Hence, $\hat{Q} \hat{P}$ is semisimple and thus $\Pi_{,}(\Psi)$ is a rank-1 eigenprojection of $\hat{Q} \hat{P}$. (A semisimple matix possesses a diagonal Jordan form. ${ }^{47.48}$ ) Although the optimal projection $\tau$ is characterized in Eqs. (25). (34), and (44) as the sum of rank-1 eigenprojections of $\hat{Q} \hat{P}$. because of the nonuniqueness in $\boldsymbol{O}(\hat{Q}, \hat{P})$. the theorems do not specify which eigenprojections actually comprise $T$. From analytical examples ${ }^{10}$ it can be seen that each of the $\binom{n_{1}}{n_{r}}$ possible projections may correspond to a local extremal in the optimizauon problem.

Remark 2.5. The proofs of Theorems 2.1-2.3 are similar to the continuous-time results and, hence, have been omitted. To help the reader reconstruct the lengthy manipulations, the key details differing from the continuous-time case are pointed out For the control problem. an $\left(n+n_{t}\right) \times\left(n+n_{c}\right)$ discretetime algebraic Lyapunov equation is obtained for the steady state covariance of the closed-loop system. Regarding this equation as a side constraint. the Lagrange multiplier technique is used to compute stationarity conditions that yield explicit expressions for $A_{1}, B_{1}, C_{1}$. and $D_{\text {. }}$. The piojection arises when these expressions are substituted into the original


Fia. 1 Roor-mean-square performance vs controiler order for five-mode beam evample
augrented Lyapunov equation and its dual. The interesting aspect is that the explicit gain expressions and the definition of the optimal projection arise in the reverse order as compared to the continuo تix :im derivation. Similar remarks apply to the reduced-order modeling and estimation problems.

## III. Examples

As an application of Theorem 2.2 on reduced-order state estimation, the stirred-tank example from Ref. 36, pp. 107. 473, and 531, is considered. Igooring the undisturbed volume state, the remaining states are the incremental tank concentration and variations in the feed concentrations. The problem data are as follows:

$$
\begin{aligned}
& A=\left[\begin{array}{lll}
0.9048 & 0.06702 & 0.02262 \\
0 & 0.8825 & 0 \\
0 & 0 & 0.9048
\end{array}\right] \\
& C=\left[\begin{array}{lll}
1 & 0 & 0
\end{array}\right], V_{2}=10^{-6}, V_{12}=0, N=L=I_{3} \\
& V_{1}=\left[\begin{array}{lll}
5.399 \times 10^{-6} & 8.015 \times 10^{-5} & 8.762 \times 10^{-5} \\
8.015 \times 10^{-5} & 2.212 \times 10^{-3} & 0 \\
8.762 \times 10^{-5} & 0 & 7.251 \times 10^{-3}
\end{array}\right]
\end{aligned}
$$

The standard Kalman filter result is

$$
\begin{aligned}
& A_{e}=\left[\begin{array}{lll}
-0.7959 & 0.06702 & 0.2262 \\
-8.205 & 0.8825 & 0 \\
-10.29 & 0 & 0.9048
\end{array}\right] \\
& B_{e}=\left[\begin{array}{c}
1.701 \\
8.205 \\
10.29
\end{array}\right] \\
& C_{e}=I_{3}, \quad D_{e}=0
\end{aligned}
$$

with performance $J_{c}=0.0358515$
Permitting nonzero feedthrough $D_{c}$. yields

$$
\begin{aligned}
& A_{e}=\left[\begin{array}{clc}
0.09885 & 0.1061 & -0.0167 \\
-0.6851 & 0.02176 & 0.09015 \\
0.1503 & 0.2652 & 0.8707
\end{array}\right], B_{e}=\left[\begin{array}{c}
-13.34 \\
-10.08 \\
2.237
\end{array}\right] \\
& C_{r}=\left[\begin{array}{clc}
-0.006261 & -0.007282 & 0.0007445 \\
0.5716 & -0.0002166 & 0.007202 \\
0.6822 & -0.006021 & -0.005440
\end{array}\right] \\
& D_{r}=\left[\begin{array}{c}
0.9068 \\
9.298 \\
11.37
\end{array}\right]
\end{aligned}
$$

where the (improved) performance is $J_{e}=0.032401049$
For the reduced-order results. an algorithm for solving all three equations (31-33) is described briefly. Begin by setting $\tau=I$ and solving Eqs. (31-33) for the "full-order" values of $\boldsymbol{Q}$. $\hat{Q}$, and $\hat{P}$. Choose $n_{e}$ eigenprojections of $\hat{Q} \hat{P}$ in diagonalizing coordinates and iterate the modified Lyapunov equations (32) and (33) until convergence of $r, \hat{Q}$, and $\hat{P}$ is obtained Return to Eq. (31) and solve for $Q$ with $V_{1}+\tau_{\perp} Q T_{1}^{T}$ as the new nonhomogeneous term in the Riccati equation. Repeat the above steps until convergence is reached

In applying this algorithm to the present example, the eigenprojections were chosen for convenience in accordance with the largest eigenvalues of $\hat{Q} \hat{P}$. The results indicate attainment of the globai minimum. For the optimal second-order filter, the gains are given by

$$
\begin{aligned}
& A_{c}=\left[\begin{array}{ll}
0.09898 & -0.1137 \\
0.6632 & -0.02285
\end{array}\right], \quad B_{e}=\left[\begin{array}{c}
-13.33 \\
9.762
\end{array}\right] \\
& C_{e}=\left[\begin{array}{cc}
-0.006259 & 0.007796 \\
0.5716 & -0.0002737 \\
0.6823 & 0.006121
\end{array}\right], \quad D_{r}=\left[\begin{array}{c}
0.9068 \\
9.298 \\
11.37
\end{array}\right]
\end{aligned}
$$

with $J_{e}=0.032401094$, and, for the first-order filter,

$$
\begin{aligned}
& A_{e}=0.1498, \quad B_{e}=-14.82 \\
& C_{e}=\left[\begin{array}{c}
-0.001661 \\
0.5748 \\
0.6897
\end{array}\right], \quad D_{e}=\left[\begin{array}{c}
0.9749 \\
9.489 \\
11.65
\end{array}\right]
\end{aligned}
$$

with $J_{e}=0.03240418$. Convergence to this accuracy was obtained with 7 iterations of Eqs. (31-33) for the second-order filter and 10 iterations for the first-order filter. Note that the performance degrades only slighly with reduced order, and the static gain term gives the first-order filter better performance than the standard (full-order) Kalman filter.
To illustrate Theorem 2.3 for designing reduced-order dynamic compensators, consider a simply supported beam with two colocated sensor/actuator pairs. Assuming the beam has length 2 and that the sensor/actuator pairs are placed at coordinates $a=55 / 172$ and $b=46 / 43$, a continuous-time model of the following form is obtained:

$$
\dot{x}=\hat{A} x+\hat{B} u+\hat{w}_{1} \quad y=\hat{C} x+\hat{w}_{2}
$$

where, retaining the first five modes,

$$
\begin{aligned}
& \hat{A}=\text { block-diag }\left(\left[\begin{array}{cc}
0 & 1 \\
-\omega_{1}^{2} & -2 \zeta \omega_{1}
\end{array}\right], \cdots,\left[\begin{array}{cc}
0 & 1 \\
-\omega_{5}^{2} & -2 \zeta \omega_{5}
\end{array}\right]\right) \\
& \omega_{1}=i^{2} ; \quad i=1, \ldots, 5, \quad \zeta=0.005 \\
& \hat{B}_{(, 1)}=0.5\left(1+(-1)^{\prime}\right) \sin (i \pi b / 2), \quad i=1, \ldots, 10 \\
& \hat{B}_{(1,2)}=-0.5\left(1+(-1)^{\prime}\right) \sin (i \pi a / 2), \quad i=1, \ldots, 10 \\
& \hat{C}=\hat{B}^{T}
\end{aligned}
$$

The intensities $\hat{\boldsymbol{V}}_{1}$ and $\hat{\boldsymbol{V}}_{2}$ of $\hat{\boldsymbol{w}}_{1}$ and $\hat{\boldsymbol{w}}_{2}$ are chosen to be

$$
\hat{V}_{1}=0.1 I_{10}, \quad \hat{V}_{2}=0.01 I_{2}
$$

and it is assumed that $\hat{w}_{1}$ and $\hat{w}_{2}$ are uncorrelated. For the continuous-time cost

$$
J=\lim _{1 \rightarrow \infty} E\left[x^{\top} R_{1} x+2 x^{\top} \hat{R}_{12} u+u^{\top} \hat{R}_{2} u\right]
$$

set

$$
\begin{gathered}
\hat{R}_{1}=\operatorname{block} \cdot \operatorname{diag}\left(\left[\begin{array}{cc}
1 & 0 \\
0 & 1 / \omega_{1}
\end{array}\right] \cdots,\left[\begin{array}{cc}
1 & 0 \\
0 & 1 / \omega_{5}
\end{array}\right]\right) \\
\hat{R}_{12}=0, \quad \hat{R}_{2}=I_{2}
\end{gathered}
$$

To convert to the discrete-time problem with discretization interval $h$, let ${ }^{36,49}$

$$
\begin{aligned}
& A=e^{i n}, \quad B=\int_{0}^{n} e^{i t} \hat{B} d t, \quad C=\hat{C} \\
& V_{1}=\int_{0}^{n} e^{i t} \hat{V}_{1} e^{\lambda^{\top} t} d t, \quad V_{2}=\hat{V}_{2} \\
& R_{1}=\hat{R}_{1}, \quad R_{12}=0, \quad R_{2}=\hat{R}_{2}
\end{aligned}
$$

The design equations (40-43) for the control problem can be solved using exactly the same techniques as in the previous example for the estimation problem. For the strictly proper case ( $D_{f}=0$ ), a series of controllers was designed with $n_{c}=$ $1.2, \ldots, 10$. where the $n_{c}=10$ result is the LQG solution. The
gains for the case $n_{1}=4$, for example, are given by

$$
\begin{aligned}
& A_{i}=\left[\begin{array}{llll}
0.9317 & 0.1572 & -0.2130 & -0.005038 \\
0.0137 & 0.6879 & 0.2519 & 0.4085 \\
0.3330 & -0.0580 & 0.7713 & -0.2602 \\
0.05980 & -0.3297 & 0.3918 & 0.3005
\end{array}\right] \\
& B_{i}=\left[\begin{array}{ccc}
-0.4920 & -0.2166 \\
0.6179 & -0.5959 \\
0.2253 & -0.02572 \\
0.07221 & -0.4863
\end{array}\right] \\
& C_{i}=\left[\begin{array}{cccc}
0.05864 & -0.3094 & -0.01815 & 0.2409 \\
-0.1301 & 0.1463 & -0.1945 & -0.07192
\end{array}\right]
\end{aligned}
$$

Figure 1 summarizes the results for each order, where the rms controller performance is given by

$$
J_{i}^{\prime}=E\left[\lim _{k \rightarrow x} x(k)^{T} R_{1} x(k)\right]^{\frac{1}{2}}
$$

These rest is provide a tradeoff study of performance versus controller urder that can be used to assess processor requirements.

## IV. Conclusion

Optimality conditions have been obtained for the problems of least-squares, reduced-order (i.e., fixed-order), discrete-time modeling. estimation, and control. These conditions comprise systems of two. three, and four matrix equations, respectively, coupled by an oblique projection which determines the optimal system gains. When the order of the estimator or controller is equal to the order of the plant, the oblique projection becomes the identity matrix and the estimation and control results specialize to the standard discrete-time Kalman filter and linear-quadratic-Gaussian results. The design results are applied to two illustrative examples. For a third-order stirredtank problem, filters of first and second order are obtained, and, for a simply supported Euler beam example with five flexible modes (i.e., 10 states), a series of reduced-order controllers with $1,2, \ldots 9$ poles is obtained. The latter results illustrate the tradeoff between control-system performance and controller order.

## Acknowledgment

The authors wish to thank Scott W. Greeley for carrying out the computations in Sec. III.

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# The Optimal Projection Equations for Fixed-Order Sampled-Data Dynamic Compensation with Computation Delay 

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#### Abstract

For an LQG-iype sampled-data regulator problem which accounts for computational delay and utilizes an averaging $\mathbf{A} / \mathbf{D}$ device, the equivalent discrete-time problem is shown to be of increased order due to the inclusion of delayed measurement states. The optimal projection equations for reduced-order, discrete-time compensation are applied to the augmented problem to characterize low-order controllers. The design results are illustrated on a tenth-order flexible beam example.


## I. INTRODUCTION

Classical sampled-data control theory has been extensively developed [1]-[7] and is widely used in practival applications. Sampled-data design based upon modern optimal control theory has also been developed, although to a considerably lesser extent [8]-[14]. The goals of the present note are twofold. First, for an LQG-type sampled-data regulator problem which explicitly accounts for computational delay, we obtain an equivalent discrete-time problem (Theorem 2.1 and Corollary 2.1). The timing diagram in Fig. 1 illustrates the unavoidable delay in the feedback loop (see Section II for notation).' A salient feature of this problem is that rather than reflace the continuous-time white noise measurement model by a discrete-time version (which is often done in the literature since continuous-time white noise cannot be sampled). we employ an averag-ing-type A/D device as in [8, p. 82] [see (2.5)].
The second goal of the note is to present a novel design procedure which is applicable to the equivalent discrete-time problem, and which thus directly accounts for the delay effects. Since the discrete-time model is of augmented order $n+l(n=$ number of plant states, $l=$ number of measurements), it seems natural to seek dynamic feedback of reduced order. To this end. we apply the optimal projection equations for discretetime dynamic compensation to the equivalent discrete-time problem to characterize optimal controllers of order $n_{c} \leq n+l$. These equations, which were previously derived in [17] for the continuous-time case, are discussed in [15]. [16]. Note that, in practice, the computational delay (and. hence, sample interval) in real-time controller implementation depends directly upon the controller order $n_{c}$. For example, by reducing $n_{c}$ the sample rate can effectively be increased. Thus, the engineering tradeoffs of performance versus controller order and sample interval can be investigated using the results of this note.

This note also includes formulas for integrals of matrix exponentials arising in the sampled-data/discrete-time conversion, along with an algorithm for solving the optimal projection equations. The results are applied to a tenth-order flexible beam example.

## II. Sampled-Data Problem and Equivalent Discrete-Time Formulation

The following notation and definitions will be used throughout.

$$
\begin{array}{ll}
I_{r}, 0_{r}, \ldots, & r \times r \text { identity matrix, } r \times s \text { zero matrix, } 0_{r \times r} \\
Z_{1 .}, 1 & (i, j) \text { element of matrix } Z
\end{array}
$$

Manuscript received Octoher 4, 1985; revised May 12. 1986. This work was supported in part by the Aur Force Office of Scientific Research under Contracts AFOSR F49620. 84C-0015. F49620-K6-C-0002, and F49620-86-C-0038

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IEEE Log Number K609472.
${ }^{1}$ For simplicity, the timing diagram Fig 1 applies to the case in which the compensator is strictl) proper. In the noxe the results are suated for the more general case in whech a direct (sutic) feedthrough $\operatorname{eem} D_{i} \hat{y}(k)$ is imiluded.


Fig. 1. Timing diagram for sampled-data controller.

| $Z^{T}, Z^{-T}$ | transpose of vector or matrix $Z,\left(Z^{r}\right)^{-1}$ |
| :---: | :---: |
| $E_{i}$ | matrix with unity in the ( $i, i$ ) position and zeros elsewhere |
| $E, R, R^{\prime \times s}$ | expected value, real numbers, $r \times s$ real matrices |
| stable matrix | matrix with eigenvalues in open unit disk |
| nonnegativesemisimple matrix | diagonalizable matrix with nonnegative eigenvalues [18]. [19] |
| $n, m, l, n_{c}$ | positive integers, $1 \leq n_{c} \leq n+l$ |
| $x, u, y, x_{\text {c }}$ | $n, m, l, n_{c}$-dimensional vectors |
| $A, B, C$ | $n \times n, n \times m, l \times n$ matrices |
| $A_{c}, B_{c}, C_{c}, D_{c}$ | $n_{c} \times n_{c}, n_{c} \times l, m \times n_{c}, m \times 1$ matrices |
| $w_{1}, w_{2}$ | $n$, $l$-dimensional zero-mean continuous-time white noise processes |
| $V_{1}$ | $n \times n$ nonnegative-definite intensity of $w_{1}$ |
| $V_{2}$ | $l \times 1$ positive-definite intensity of $w_{2}$ |
| $V_{12}$ | $n \times 1$ cross intensity of $w_{1}, w_{2}$ |
| $\boldsymbol{R}_{1}$ | $n \times n$ nonnegative-definite state weighting matrix |
| $\boldsymbol{R}_{2}$ | $m \times m$ positive-definite control weighting matrix |
| $R_{12}$ | $n \times m$ cross weighting matrix such that $R_{1}-$ $R_{12} R_{2}^{-1} R_{1}^{T}$ is nonnegative definite |
| $k$ | discrete-time index, 1, 2, 3, ${ }^{\text {a }}$, |

in the statement of the sampled-data control problem the sample interval $h$ and the controller order $n_{c}$ are fixed and the optimization is performed over the controller parameters $\left(A_{c}, B_{c}, C_{c}, D_{c}\right)$. For design tradeoff studies $h$ and $n_{c}$ can be varied and the problem can be solved for each pair of values of interest.

## Fixed-Order, Sampled-Data Dynamic-Compensation Problem

Given the $\boldsymbol{n}$ th-order continuous-time system

$$
\begin{equation*}
\dot{x}(t)=A x(t)+B u(t)+w_{1}(t) \tag{2.1}
\end{equation*}
$$

with continuous-time measurements

$$
\begin{equation*}
y(t)=C x(t)+D u(t)+w_{2}(t) \tag{2.2}
\end{equation*}
$$

design an $\boldsymbol{n}_{\mathrm{f}}$ th-order discrete-time compensator

$$
\begin{gather*}
x_{r}(k+1)=A_{r} x_{r}(k)+B_{r} y(k),  \tag{2.3}\\
\hat{u}(k)=C_{r} x_{r}(k)+D_{r} f(k) \tag{2.4}
\end{gather*}
$$

which, with A/D averaged measurements

$$
\begin{equation*}
f(k) \leq \frac{1}{h} \int_{(\alpha-1)^{n}}^{a n} y(t) d t \tag{2.5}
\end{equation*}
$$

and D/A wero-order-hold controls

$$
\begin{equation*}
u(t)=\Delta(k), \quad t \in[k h,(k+1) h), \tag{2.6}
\end{equation*}
$$

minimizes the performance criterion

$$
\begin{align*}
J\left(A_{r}, B_{c},\right. & \left.C_{c}, D_{c}\right) \underline{\underline{i}} \lim _{i \rightarrow \infty} E \frac{1}{t} \\
& \cdot \int_{0}^{1}\left[x(s)^{T} R_{1} x(s)+2 x(s)^{\tau} R_{12} u(s)+u(s)^{\tau} R_{2} u(s)\right] d s \tag{2.7}
\end{align*}
$$

The main result of this section concerns propagation of the plant and digitized measurements over one time step. For notational convenience, define

$$
H(s) \triangleq \int_{0}^{s} e^{t r} d r
$$

Theorem 2.1: For the fixed-order, sampled-data dynamic-compensation problem, the plant dynamics (2.1), averaged measurements (2.5) and performance criterion (2.7) have the equivalent discrete-time representations

$$
\begin{gather*}
x^{\prime}(k+1)=A^{\prime} x^{\prime}(k)+B^{\prime} \hat{u}(k)+w^{\prime}(k)  \tag{2.8}\\
\hat{y}(k)=C^{\prime} x^{\prime}(k-1)+D^{\prime} \hat{u}(k-1)+w_{2}^{\prime}(k-1) \tag{2.9}
\end{gather*}
$$

$J\left(A_{c}, B_{c}, C_{c}, D_{c}\right)=\delta+\lim _{k \rightarrow \infty} E\left[x^{\prime}(k)^{\top} R_{i} x^{\prime}(k)\right.$

$$
\begin{equation*}
\left.+2 x^{\prime}(k)^{T} R_{i_{2}}^{\prime} \dot{u}(k)+\hat{u}(k)^{T} R_{2}^{\prime} \hat{u}(k)\right] \tag{2.10}
\end{equation*}
$$

where

$$
x^{\prime}(k) \leqq x(k h) . \delta \triangleq \frac{1}{h} \operatorname{tr} \int_{0}^{A} \int_{0}^{s} e^{A} V_{1} e^{A^{T}} R_{1} d r d s
$$

$$
A^{\prime} \triangleq e^{A \star}, B^{\prime} \triangleq H(h) B, C^{\prime} \triangleq \frac{1}{h} C H(h), D^{\prime} \triangleq \frac{1}{h} C \int_{0}^{h} H(s) d s B+D
$$

$w_{i}(k)$ and $w_{2}^{\prime}(k)$ are zero-mean, white noise processes with

$$
E\left\{\left[\begin{array}{l}
w_{1}^{\prime}(k) \\
w_{2}^{\prime}(k)
\end{array}\right]\left[\begin{array}{l}
w_{1}^{\prime}(k) \\
w_{2}^{\prime}(k)
\end{array}\right] \tau\right\}=\left[\begin{array}{cc}
V_{i}^{\prime} & V_{12}^{\prime} \\
V_{12}^{\prime T} & V_{i}^{\prime}
\end{array}\right]
$$

where

$$
\begin{aligned}
& V_{i} \triangleq \int_{0}^{n} e^{A s} V_{1} e^{A} T_{s} d s, V_{12}^{\prime} \triangleq \frac{1}{h} \int_{0}^{h} e^{A s} V_{1} H^{T}(s) d s C^{T}+\frac{1}{h} H(h) V_{12} \\
& V_{2}^{\prime} \triangleq \frac{1}{h} V_{2}+\frac{1}{h^{2}} C \int_{0}^{n} H(s) V_{1} H^{T}(s) d s C^{T} \\
& \\
& +\frac{1}{h^{2}} C \int_{0}^{n} H(s) d s V_{12}+\frac{1}{h^{2}} V_{12}^{T} \int_{n}^{n} H^{T}(s) d s C^{T}
\end{aligned}
$$

and

$$
R_{i} \leqslant \frac{1}{h} \int_{0}^{h} e^{A_{2}} R_{1} e^{4 s} d s, R_{i 2} \triangleq \frac{1}{h} \int_{0}^{A} e^{T_{3}} R_{1} H(s) d s B+\frac{1}{h} H^{r}(h) R_{i 2}
$$

$R_{i} \triangleq R_{2}+\frac{1}{h} B^{r} \int_{0}^{n} H^{r}(s) R_{1} H(s) d s B$

$$
+\frac{1}{h} B^{r} \int_{0}^{n} H^{r}(s) d s R_{12}+\frac{1}{h} R_{12}^{r} \int_{0}^{h} H(s) d s B
$$

The proof of this theorem is atraightforward calculation involving standard techniques, and hence is omitted. The result is more comprehensive than previous work, however, and includes several results as special cases. For example, the expressions for $A^{\prime}$ and $B^{\prime}$ are standard; $C^{\prime}$ is
given by (10.9), [8, p. 83]; $R_{i}^{\prime}, R_{12}^{\prime}$, and $R_{i}$ are given in [10], [12]; and $V_{1}^{\prime}, V_{i 2}^{\prime}$, and $V_{2}^{\prime}$ can be found in $[8, p, 85]$. The expressions for $\delta$ and $D^{\prime}$ appear to be new.

Note that the averaged measurements depend upon delayed samples of the state. By augmenting the discretized state equation (2.8) to include these measurements, it is possible to state the original sampled-data problem as a discrete-time problem.

Corollary 2.1. With the nocation

$$
\begin{gathered}
f(k) \triangleq\left[\begin{array}{l}
x^{\prime}(k) \\
f(k)
\end{array}\right], A \triangleq\left[\begin{array}{cc}
A^{\prime} & 0_{n \times 1} \\
C^{\prime} & 0_{i}
\end{array}\right], \hat{B}\left[\begin{array}{l}
B^{\prime} \\
D^{\prime}
\end{array}\right], \\
C \triangleq\left[0_{1 \times n} \quad l_{1}\right\}, \hat{w}(k) \triangleq\left[\begin{array}{l}
w_{i}^{\prime}(k) \\
w_{2}^{\prime}(k)
\end{array}\right], \hat{V}\left[\begin{array}{cc}
V_{i}^{\prime} & V_{i 2} \\
V_{12}^{\prime} & V_{2}
\end{array}\right], \\
R_{1} \triangleq\left[\begin{array}{cc}
R_{1}^{\prime} & 0_{n \times 1} \\
0_{1 \times n} & 0_{1}
\end{array}\right], R_{12} \triangleq\left[\begin{array}{l}
R_{i 2} \\
0_{1 \times m}
\end{array}\right], R_{2} \triangleq R_{2}
\end{gathered}
$$

the fixed-order, sampled-data dynamic-compensation problem is equivalent to the following discrete-time problem. Given the $(n+l)$ th-order discrete-time system

$$
\begin{equation*}
\hat{x}(k+1)=\bar{A} \tilde{x}(k)+\bar{B} \dot{u}(k)+\hat{w}(k) \tag{2.11}
\end{equation*}
$$

with discrete-time measurements

$$
\begin{equation*}
f(k)=C \hat{x}(k) \tag{2.12}
\end{equation*}
$$

design an $n_{c}$ th-order discrete-ime compensator of the form (2.3), (2.4), which minimizes

$$
\begin{align*}
& J\left(A_{c}, B_{c}, C_{c}, D_{c}\right)=\delta+\lim _{k \rightarrow-\infty} E\left[\hat{x}(k)^{\top} R_{1} \hat{x}(k)\right. \\
&\left.+2 \hat{x}(k)^{T} R_{12} \tilde{u}(k)+\hat{u}(k)^{T} R_{2} \hat{u}(k)\right] . \tag{2.13}
\end{align*}
$$

Remark 2.1: The equivalent cost (2.13) involves a constant offset $\delta$ which serves as a lower bound on the sampled-data performance, i.e., a "discretization floor." Note that

$$
\delta=\frac{h}{2} \operatorname{tr} V_{1} R_{1}+O\left(h^{2}\right)
$$

Remark 2.2: Although the measurements $\hat{\varphi}(k)$ are noise free, the singularity is not so serious as singular measurement noise in the continuous-time case where the Kalman filter gains are expressed in terms of the inverse of the measurement noise intensity. In the discrete-time case, rather, it is required that $r+\overline{C Q C}^{T}$ be invertible, where $\bar{V}$ is the measurement noise covariance (see [11, p. 530], or [15], [16]).

Remark 2.3: The increase in plant order from $n$ to $n+l$ is due to the computational delay and A/D process. Since discrete-time LQG theory yields a possibly unwieldy ( $n+l$ )th-order controller, we seek "reducedorder" controllers. Note that in this context an $n$ th-order controller can be regarded as being of reduced order.

Remark 2.4: As pointed out in [10], particular choices of the sample interval $h$ may result in a loss of controllability and observability for the equivalent discrete-time problem. Hence, these properties must be verified before applying control design methods.

## III Application of the Optimal projection Equations to the. Equivalent Discrete-Time Problem

We now apply the optimal projection equations for discrete-time dynamic compensation to the equivatent discrete-time problem. The following easily proved lemma will be needed.

Lemma 3.1: Let $t \in \boldsymbol{R}^{(n+1) \times(n+1)}$. Then

$$
\begin{equation*}
T^{2}=T \tag{3.1}
\end{equation*}
$$

if and only if there exist $G, r \in R^{n \times(n+1)}$ such that

$$
\begin{align*}
& G^{T} \Gamma=r,  \tag{3.3}\\
& r G^{r}=I_{r q} . \tag{3.4}
\end{align*}
$$

Furthennore, $G$ and $\Gamma$ are unique to a change of basis in $R^{n} c$.
Call $G$ and $\Gamma$ satisfying (3.3), (3.4) a projective factorization of $\tau$. Furthermore. for $n \times n$ nonnegative-definite matrices $Q$ and $P$, define the set of contragrediently diagonalizing transformations

$$
\mathcal{C}(Q . \mathbb{Q}) \leqq\left\{\Psi \in R^{n \times n}: \Psi^{-1} Q \Psi^{-T} \text { and } \Psi^{\top} \mathcal{P} \Psi \text { are diagonal }\right\} .
$$

It follows from [19, Theorem 6.2.5, p. 123] that $\mathcal{C}(\mathcal{Q}, \mathcal{Q})$ is always nonemply. This set does not, however. have a unique eiement since basis rearrangements and sign transpositions may be incorporated into $\Psi$. Further nonuniqueness arises if $Q \odot$ has repeated eigenvalues.

To guarantee that $J$ is finite and independent of initial conditions, we restrict our consideration to the (opea) set

$$
\begin{aligned}
\delta \doteq\left\{\left(A_{c}, B_{c}, C_{c}, D_{c}\right):\right. & {\left[\begin{array}{cc}
\bar{A}+B D_{c} C & B C_{c} \\
B_{c} C & A_{c}
\end{array}\right] } \\
& \text { is stable and } \left.\left(A_{c}, B_{c}, C_{c}\right) \text { is minimal }\right\} .
\end{aligned}
$$

For the design problem it is required that $\delta$ be nonempry, i.e., that the augmented system be stabilizable. We also require the notation

$$
R_{2} \triangleq R_{2}+B^{T} P B_{1}, \hat{V}_{2} \triangleq C Q C^{\top}, T_{+} \triangleq I_{n+1}-T
$$

$$
A_{Q} \triangleq \hat{A}-\hat{A} Q \mathcal{C}^{\top} \hat{D}_{2}^{-1} C, A_{p} \triangleq \hat{A}-B R_{2}^{-1}\left(B^{\top} P \hat{A}+R_{12}^{T}\right)
$$

$$
\Sigma_{Q} \triangleq\left(\tilde{A} Q C^{\top}+B D_{1} \hat{V}_{2}\right) \tilde{V}_{2}^{-1}\left(\hat{A} Q C^{T}+\hat{B} D_{c} \tilde{V}_{2}\right)^{r}
$$

$$
\Sigma_{\rho} \triangleq\left(\hat{B}^{\top} P \hat{A}+R_{12}^{T}+\hat{R}_{2} D_{c} \hat{C}\right)^{\top} \hat{R}_{2}^{-1}\left(B^{\top} P \hat{A}+R_{12}^{T}+\hat{R}_{2} D_{c} \hat{C}\right)
$$

$$
M \triangleq\left[\begin{array}{l}
I_{n+1} \\
D_{c} \mathcal{C}
\end{array}\right], \bar{M} \triangleq\left[\begin{array}{c}
I_{n+1} \\
-\hat{R}_{2}^{-1}\left(\dot{B}^{r} P \bar{A}+\kappa_{12}^{T}\right)
\end{array}\right]
$$

$$
R \equiv\left[\begin{array}{ll}
R_{1} & R_{12} \\
R_{12}^{T} & R_{2}
\end{array}\right]
$$

Theorem 3. J: Suppose ( $A_{c}, B_{c}, C_{c}, D_{c}$ ) $\in S$ solves the fixed-order. sampled-data dynamic-compensation problem. Then there exist $(n+l)$ $x(n+l)$ nonnegative-definite matrices $Q, P, Q$, and $P$ such that $A_{c}$, $B_{c}, C_{c}$, and $D_{c}$ are given by
and such that $Q, P, Q$, and $\boldsymbol{P}$ satisfy
where

$$
\begin{equation*}
r=\sum_{i=1}^{n_{r}} \pm E_{i} \pm-1 \tag{3.13}
\end{equation*}
$$

for some $\psi \in \mathbb{C}(Q, \beta)$ such that $\left(\boldsymbol{\psi}^{-1} Q^{\beta} \Psi\right)_{(0, i)} \neq 0, i=1, \cdots, n_{c}$, and some projective factorization $G, \Gamma$ of $r$. Furthermore, the minimal cost is

$$
\begin{align*}
& \boldsymbol{Q}=\hat{A}_{P}{ }^{T} \boldsymbol{Q}_{T}{ }^{T} \hat{A}_{p}^{T}+\Sigma_{Q} . \tag{3.10}
\end{align*}
$$

$$
\begin{align*}
& A_{\mathrm{c}}=\Gamma\left[\hat{A}-\hat{A} Q \mathcal{C}^{\top} \bar{D}_{2}^{-1} \bar{C}-B \hat{R}_{2}^{-1}\left(\mathcal{B}^{\top} P \hat{A}+R_{1}^{T}\right)-B D_{c} C\right] G^{T},  \tag{3.5}\\
& \left.B_{c}=\Gamma \mid \hat{A} Q \mathcal{C}^{\top} \hat{V}_{i}^{-1}+\hat{B} D_{c}\right] . \tag{3.6}
\end{align*}
$$

$$
\begin{align*}
& D_{c}=-\hat{R}_{2}^{-1}\left(\bar{B}^{\top} P \hat{A} Q C^{\top}+\hat{R}_{12}^{\top} Q C^{\top}\right) \bar{V}_{2}^{-1} . \tag{3.7}
\end{align*}
$$

given by

$$
\begin{equation*}
J\left(A_{c}, B_{c}, C_{f}, D_{c}\right)=\delta+\operatorname{tr}\left\{\left(M Q M^{\top}+\hat{M}_{T} Q_{T}{ }^{T} \dot{M}^{T}\right) \dot{R}\right\} . \tag{3.14}
\end{equation*}
$$

Remark 3.1: Theorem 3.1 can immediately be specialized to the more restrictive probiem in which the compensator is strictly proper. This can be done in both the full- and reduced-order cases by ignoring (3.8) and setting $D_{c}=0$ wherever it appears. See [15]. [16].

## IV. Numerical Evaluation of Integrals Involving Matrix Exponentials

To evaluate the exponential/integral expressions appearing in Theorem 2.1, we utilize the approach of [20]. The idea is to eliminate the need for integration by computing the matrix exponential of appropriate block matrices. Numerical matrix exponentiation is discussed in [21].

Proposition 4.1: Consider the following paritioned matrix exponentials of order $(3 n+l) \times(3 n+l) .(3 n+m) \times(3 n+m),(2 n+m)$ $\times(2 n+m)$, and $(3 n) \times(3 n)$, respectively:
$\left[\begin{array}{cccc}F_{1} & F_{2} & F_{3} & F_{4} \\ 0_{n} & F_{5} & F_{6} & F_{1} \\ 0_{n} & 0_{n} & F_{2} & F_{9} \\ 0_{1 \times n} & 0_{1 \times n} & 0_{1 \times n} & I_{1}\end{array}\right] . \Leftrightarrow \exp \left[\begin{array}{cccc}-A & I_{n} & 0_{n} & 0_{n \times 1} \\ 0_{n} & -A & V_{1} & V_{12} \\ 0_{n} & 0_{n} & A^{T} & C^{T} \\ 0_{1 \times n} & 0_{1 \times n} & 0_{1 \times n} & 0_{1}\end{array}\right] h$,
$\left[\begin{array}{cccc}F_{10} & F_{11} & F_{12} & F_{13} \\ 0_{n} & F_{14} & F_{1 s} & F_{10} \\ 0_{n} & 0_{n} & F_{19} & F_{10} \\ 0_{m \times n} & 0_{m \times n} & 0_{m \times n} & I_{m}\end{array}\right] \triangleq \exp \left[\begin{array}{cccc}-A^{r} & I_{n} & 0_{n} & 0_{n \times m} \\ 0_{n} & -A^{T} & R_{1} & R_{12} \\ 0_{n} & 0_{n} & A & B \\ 0_{m \times n} & 0_{m \times n} & 0_{m \times n} & 0_{m}\end{array}\right] h$,

$$
\begin{aligned}
& {\left[\begin{array}{ccc}
I_{n} & F_{19} & F_{20} \\
0_{n} & F_{21} & F_{22} \\
0_{m \times n} & 0_{m \times n} & I_{m}
\end{array}\right] \triangleq \exp \left[\begin{array}{ccc}
0_{n} & I_{n} & 0_{n \times m} \\
0_{n} & A & B \\
0_{m \times n} & 0_{m \times n} & 0_{m}
\end{array}\right] h,} \\
& {\left[\begin{array}{ccc}
F_{23} & F_{24} & F_{23} \\
0_{n} & F_{22} & F_{21} \\
0_{n} & 0_{n} & F_{23}
\end{array}\right] \triangleq \exp \left[\begin{array}{ccc}
-A & I_{n} & 0_{n} \\
0_{n} & -A & V_{1} \\
0_{n} & 0_{n} & A^{T}
\end{array}\right] h .}
\end{aligned}
$$

Then

$$
\begin{aligned}
& A^{\prime}=F_{13}, B^{\prime}=F_{13}, C^{\prime}=\frac{1}{h} F_{9}^{T}, D^{\prime}=\frac{1}{h} C F_{20}+D, \delta=\frac{1}{h} \operatorname{tr} R_{1} F_{24}^{T} F_{23}, \\
& V_{i}^{\prime}=F_{1}^{\tau} F_{60}, V_{12}^{\prime}=\frac{1}{h} F_{1}^{T} F_{1,}, V_{2}^{\prime}=\frac{1}{h}\left(V_{2}+\frac{1}{h} C F_{1}^{T} F_{4}+\frac{1}{h} F_{4}^{\tau} F_{3} C^{r}\right), \\
& R_{1}^{\prime}=\frac{1}{h} F_{12}^{T} F_{13}, R_{12}^{\prime}=\frac{1}{h} F_{17}^{T} F_{16}, R_{2}^{\prime}=R_{2}+\frac{1}{h}\left(B^{T} F_{17}^{T} F_{13}+F_{13}^{T} F_{17} B\right) .
\end{aligned}
$$

## V. Numerical Solution of the discrete-Time Optimal Projection Equations

The following algorithm is proposed for solving (3.9)-(3.12).

## Algorithm S.I:

Step 1: Initialize $k=0$ and $\tau^{(0)}=I_{n, 1}$.
Step 2: With $T_{\underline{1}} \boldsymbol{T}^{(k)}$ solve (3.9)-(3.12) for $Q^{(k)} \triangleq Q . P^{(1)} \triangleq P, Q^{(*)}$ $\leqq Q$. and $\boldsymbol{P}^{(k)} \leqq \boldsymbol{P}$.
Step 3: If $k \geq 1$ check for convergence: If $\|\left(Q^{(k)}, P^{(k)}, Q^{(k)}, P^{(k)}\right)-$ $\left(Q^{(k-1)}, P^{(k-11}, Q^{(k-1)}, P^{(k-1)}\right) \|>$ tol then continue; else go to Step 6.
Step 4: Select $\Psi^{(k)} \in \mathbb{C}\left(\phi^{(k)}, \beta^{(t)}\right)$ and update $\tau^{(k+1)}=\sum^{n c}$, $\boldsymbol{\psi}^{(k)} \boldsymbol{E}_{i}\left(\boldsymbol{\Psi}^{(k)}\right)^{-1}$.
Step S: Increment $k$ and go to Step 2.
Step 6: Evaluate (3.5)-(3.8) with $Q=Q^{(t)}, P=P^{(t)}, Q=Q^{(t)}, \rho$ $=\beta^{(k)}, G^{\top} \Gamma=r^{(k)}, \Gamma G^{\top}=I_{n_{c}}$.
Remark S.I: In solving the Riccati equation (3.9), the nonhomogeneous term is taken to be $D+r_{\Lambda} Q r_{1}$, which is nonnegative definite. SI uilar remarks apply to (3.10).

TABLE 1
COSt $J$ for beam Example

|  | 0 | . 1 | . 5 |
| :---: | :---: | :---: | :---: |
| 12 | -- | 1.3715 | 3.0134 |
| 10 | 1.1677 | 1.3723 | 3.0134 |
| 8 | 1.1883 | 1.3823 | 3.0162 |
| 6 | 1.2086 | 1.4421 | 3.0195 |
| 4 | 1.3330 | 1.5752 | 3.0812 |
| 2 | 1.4789 | 2.0425 | 3.3406 |
| Open-Loop Cost (ume) is 101.73 |  |  |  |

Remork 5.2: The critical step of Algorithm 5.1 is the choice of $\Psi^{(k)}$ for constructing the projection $\tau^{(k+1)}$. Since $\Psi^{(k)}$ can include basis rearrangements, the choice of $\Psi^{(k)}$ essentially corresponds to a selection of $n_{c}$ rank1 eigenprojections out of $n+$ /possible eigenprojections. This selection is discussed at length in [22] where it is pointed out that the choice of eigenprojections determines which local extremal will be reached by the algorithm. Component-cost methods have thus been utilized as a promising selection criteria. Because of the eigenprojection structure of the necessary conditions, Algorithm 5.1 is fundamentally different from gradient search methods.

## Vi. Illustrative Example: Control of a Flexible Beam

Consider a simply supported beam of length two with two colocated sensor/actuator pairs placed at coordinates $a_{1}=55 / 172$ and $a_{2}=45 / 43$. Define

$$
\begin{aligned}
& B_{(1,1}=0.5(-1)^{+1}\left(1+(-1)^{\prime}\right) \sin (i \pi a, / 4) \text {, } \\
& i=1, \cdots, 10, j=1,2, C=B^{r} . \\
& \begin{array}{l}
v_{1}=\underset{1=1 .}{\text { block-diag }}\left[\begin{array}{ll}
0 & 0 \\
0 & 1
\end{array}\right], v_{12}=0, v_{2}=0.1 /_{2}, \\
R_{1}=\text { block-diag } \\
\cdots=1
\end{array}\left[\begin{array}{cc}
1 & 0 \\
0 & \omega_{1}^{-2}
\end{array}\right], R_{12}=0, R_{2}=0.1 I_{2} .
\end{aligned}
$$

For $n_{c}=10,8,6,4,2$ continuous-time controllers were designed using the results of [17] and, for $n_{c}=12,10,8,6,4,2$ and $h=0.1,0.5$, strictly proper ( $D_{\mathrm{c}}=0$ ) discrete-time controllers were obbained from Theorem 3.1. The results are summarized in Table I.

## Vil. Directions for further Developments

The following extensions and related developments immediately suggest themselves: reduced-order, discrete-time modeling/estimation of
continuous-time systems [22], [23]: robust sampled-data control of uncerain systems with multiplicative noise [24]-[27]; multirate sampling [28]. [29]: alternative A/D and D/A devices and asynchronous sampling/ control update; infinite-dimensional systems [30], [31].

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# The optimal projection equations for reduced-order, discrete-time state estimation for linear systems with multiplicative white noise 

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Received 2 July 1986


#### Abstract

The optimal projection equations obtained in [2,3] for reduced-order, discrete-time state estimation are generalized to include the effects of state- and measurement-dependent noise to provide a model of parameter uncertainty. In contrast to the single matrix Riccati equation arising in the full-order (Kalman filter) case, the optimal steady-state reduced-order discrete-time estimator is characterized by three matrix equations (one modified Riccati equation and two modified Lyapunov equations) coupled by both an oblique projection and stochastic effects.


Keywords: Reduced-order Kalman filter, Robust estimation.

## 1. Introduction

In a recent series of papers [1-3] it has been shown that the first-order necessary conditions for optimal continuous and discrete-time reduced-order state-estimation can be transformed into coupled systems of three matrix equations (one modified Riccati equation and two modified Lyapunov equations). The coupling is due to the presence of an oblique projection (idempotent matrix) which arises as a rigorous consequence of the stationarity conditions. This formulation provides a direct generalization of the classical steady-state Kalman filter theory. Specifically, in the full-order case, the projection becomes the identity matrix, the additional two modified Lyapunov equations drop out, and the remaining modified Riccati equation reduces to the standard observer Riccati equation for the Kalman filter expression. Related results in reduced-order estimator design can be found in [4-17].

An acuutional extension of classical state estimation involves the inclusion of state- and measurementdependent disturbances [18-24]. One motivation for such a model is to design estimators which are desensitized, i.e., robustified, to actual parameter variations [25-31]. For the continuous-time control problem this has been justified in [32-38].

As shown in [36] for the continuous-time case, applying the optimal projection approach to the multiplicative white noise model yields an extended formulation of the optimality conditions for reducedorder state estimation. Specifically, the system of three matrix equations characterizing the optimal estimator are now coupled by both the oblique projection and stochastic effects.

[^34]The purpose of the present paper is to provide a self-contained derivation of the optimality conditions for reduced-order state estimation in the presence of both state- and measurement-dependent white noise in the discrete-time case. The goal of the development is to present the optimality conditions in a clear, concise manner to facilitate the development of numerical algorithms for practical application.

## 2. Notation and definitions

$\mathbf{R}, \mathbf{R}^{r \times s}, \mathbf{R}^{\prime}, \mathrm{E}$ real numbers, $r \times s$ real matrices, $\mathbf{R}^{r \times 1}$, expectation. $I_{n},()^{\mathrm{T}}, \otimes \quad n \times n$ identity, transpose, Kronecker product [39].
$\tau_{\perp} \quad I_{n}-\tau, \tau \in \mathbf{R}^{n \times n}$.
$n, m, l, n_{e}, q$ positive integers, $1 \leq n_{c} \leq n$.
$x, x_{e} \quad n, n_{e}$-dimensional vectors.
$y, y_{e} \quad l, q$-dimensional vectors.
$A, A_{i} ; C, C_{i} \quad n \times n$ matrices, $l \times n$ matrices, $i=1, \ldots, p$.
$A_{\mathrm{e}}, B_{\mathrm{e}}, C_{\mathrm{e}}, D_{\mathrm{e}} \quad n_{\mathrm{e}} \times n_{\mathrm{e}}, n_{\mathrm{e}} \times l, q \times n_{\mathrm{e}}, q \times l$ matrices.
$k \quad$ discrete-time index $1,2,3, \ldots$.
$v_{i}(k) \quad$ unit variance white noise, $i=1, \ldots, p$.
$w_{1}(k), w_{2}(k) \quad n$-dimensional, $l$-dimensional white noise processes.
$V_{1} \quad n \times n$ nonnegative-definite covariance of $w_{1}(k)$.
$V_{2} \quad l \times l$ positive-definite covariance of $w_{2}(k)$.
$V_{12} \quad n \times l$ cross-covariance of $w_{1}(k), w_{2}(k)$.
$R \quad q \times q$ positive-definite matrix.
$L \quad q \times n$ matrix.

$$
\begin{aligned}
& \tilde{A}=\left[\begin{array}{cc}
A & 0 \\
B_{\mathrm{e}} C & A_{\mathrm{e}}
\end{array}\right], \quad \tilde{A}_{i}=\left[\begin{array}{cc}
A_{i} & 0 \\
B_{\mathrm{e}} C_{i} & 0
\end{array}\right], \quad i=1, \ldots, p, \\
& \tilde{w}(k)=\left[\begin{array}{c}
w_{1}(k) \\
B_{e} w_{2}(k)
\end{array}\right], \quad \tilde{V}=\left[\begin{array}{cc}
V_{1} & V_{12} B_{e}^{\mathrm{T}} \\
B_{e} V_{12}^{\mathrm{T}} & B_{\mathrm{e}} V_{2} B_{\mathrm{e}}^{\mathrm{T}}
\end{array}\right], \\
& \tilde{R}=\left[\begin{array}{cc}
L^{\top} R L-L^{\top} R D_{e} C-C^{\top} D_{e}^{\top} R L+C^{\top} D_{e}^{\top} R D_{e} C+\sum_{i=1}^{p} C_{i}^{\top} D_{e}^{\top} R D_{e} C_{i} & -L^{\top} R C_{e}+C^{\top} D_{e}^{\top} R C_{i} \\
-C_{e}^{\top} R L+C_{e}^{\top} R D_{e} C & C_{e}^{\top} R C_{e}
\end{array}\right] .
\end{aligned}
$$

$Z_{(t, j)} \quad(i, j)$ element of matrix $Z$.
$\rho(\bar{Z}) \quad$ rank of matrix $Z$.
$\operatorname{tr} Z \quad$ trace of a square matrix $Z$.
$E_{i} \quad$ square matrix with unity in the ( $i, i$ ) position and zeros elsewhere.
$\pi_{i}(\psi) \quad \psi E_{i} \psi^{-1}$ (unit-rank eigenprojection).
$\mathcal{N}(Z), \mathscr{R}(Z)$ null space, range of matrix $Z$.
An asymptotically stable matrix is a matrix with eigenvalues in the open unit disk: a nonnegative-definite matrix is a symmetric matrix with nonnegative eigenvalues; and a positive-definite matrix is a symmetric matrix with positive eigenvalues.

For arbitrary $n \times n, Q, \hat{Q}, \tau$, define

$$
\begin{aligned}
& V_{2 s} \triangleq V_{2}+C Q C^{\mathrm{T}}+\sum_{i=1}^{p} C_{i}\left(Q+\tau \hat{Q} \tau^{\mathrm{T}}\right) C_{i}^{\mathrm{T}}, \quad Q_{\mathrm{s}} \triangleq A Q C^{\mathrm{T}}+V_{12}+\sum_{i=1}^{p} A_{i}\left(Q+\tau \hat{Q}^{\mathrm{T}}\right) C_{i}^{\mathrm{T}}, \\
& \hat{V}_{2 \mathrm{~s}} \triangleq V_{2}+C Q C^{\mathrm{T}}+\sum_{i=1}^{p} C_{i}(Q+\hat{Q}) C_{i}^{\mathrm{T}}, \quad \hat{Q}_{\mathrm{s}} \triangleq A Q C^{\mathrm{T}}+V_{12}+\sum_{i=1}^{p} A_{i}(Q+\hat{Q}) C_{i}^{\mathrm{T}}
\end{aligned}
$$

## 3. Problem statement and main theorem

Reduced-Order State-Estimation Problem. Given the $n$-th-order observed system

$$
\begin{align*}
& x(k+1)=\left(A+\sum_{i=1}^{p} v_{i}(k) A_{i}\right) x(k)+w_{1}(k)  \tag{3.1}\\
& y(k)=\left(C+\sum_{i=1}^{p} v_{i}(k) C_{i}\right) x(k)+w_{2}(k) \tag{3.2}
\end{align*}
$$

design an $n_{e}$-th reduced-order state estimator

$$
\begin{align*}
& x_{e}(k+1)=A_{e} x_{\mathrm{e}}(k)+B_{\mathrm{e}} y(k),  \tag{3.3}\\
& y_{e}(k)=C_{\mathrm{e}} x_{e}(k)+D_{e} y(k) \tag{3.4}
\end{align*}
$$

which minimizes the state-estimation error criterion

$$
\begin{equation*}
J\left(A_{e}, B_{e}, C_{e}, D_{e}\right) \triangleq \lim _{k \rightarrow \infty} \mathrm{E}\left[L x(k)-y_{e}(k)\right]^{\mathrm{T}} R\left[L x(k)-y_{e}(k)\right] \tag{3.5}
\end{equation*}
$$

In this formulation the matrix $L$ identifies the states, or linear combinations of states, whose estimates are desired. The order $n_{e}$ of the estimator state $x_{e}$ is determined by implementation constraints, i.e., by the computing capability available for realizing (3.3), (3.4) in real time. Note that the feedthrough term $D_{\mathrm{e}}$ permits the utilization of a static least squares estimator in conjunction with the dynamic estimator ( $A_{e}, B_{e}, C_{e}$ ). Thus, the goal of the Reduced-Order State-Estimation Problem is to design an estimator of given order that yields quadratically optimal (least squares) estimates of specified linear combinations of states.

To guarantee that $J$ is finite assume that $A$ is asymptotically stable and consider the set of asymptotically stable reduced-order (i.e., fixed-order) estimators

$$
\mathscr{S} \triangleq\left\{\left(A_{\mathrm{e}}, B_{e}, C_{\mathrm{e}}, D_{\mathrm{e}}\right): A_{\mathrm{e}} \text { is asymptotically stable }\right\} .
$$

Since the value of $J$ is independent of the internal realization of the transfer function corresponding to (3.3) and (3.4), without loss of generality we further restrict our attention to the set of admissible estimators

$$
\mathscr{S}^{+} \triangleq\left\{\left(A_{e}, B_{e}, C_{e}, D_{e}\right) \in \mathscr{S}:\left(A_{e}, B_{e}\right) \text { is controllable and }\left(A_{e}, C_{e}\right) \text { is observable }\right\} .
$$

The following factorization lemma is needed for the statement of the main result.

Lemma 3.1. Let $\tau \in \mathbb{R}^{n \times n}$. Then

$$
\begin{align*}
& \tau^{2 \rho}=\tau,  \tag{3.6}\\
& \rho(\tau)=n_{\mathrm{e}} . \tag{3.7}
\end{align*}
$$

if and only if there exist $G, \Gamma \in \mathbb{R}^{n_{e} \times n}$ such that

$$
\begin{align*}
& G^{\boldsymbol{\top}} \Gamma=\tau  \tag{3.8}\\
& \Gamma G^{\boldsymbol{T}}=I_{n_{e}} . \tag{3.9}
\end{align*}
$$

Furthermore, $G$ and $\Gamma$ are unique to a change of basis in $\mathbf{R}^{\boldsymbol{n}_{c}}$.

Proof. See [3].
For convenience, call $G$ and $\Gamma$ satisfying (3.8) and (3.9) a projective factorization of $\tau$. Furthermore, for $n \times n$ nonnegative-definite matrices $\hat{\mathcal{Q}}$ and $\hat{P}$, define the set of contragrediently diagonalizing transformations

$$
\mathscr{D}(\hat{Q}, \hat{P}) \triangleq\left\{\psi \in \mathbf{R}^{n \times n}: \psi^{-1} \hat{Q} \psi^{-\top} \text { and } \psi^{\top} \hat{P} \psi \text { are diagonal }\right\} .
$$

It follows from Theorem 6.2.5, p. 123 of [40], that $\mathscr{D}(\hat{Q}, \hat{P})$ is always nonempty. This set does not, however, have a unique element since basis rearrangements and sign transpositions may be incorporated into $\psi$. Further nonuniqueness arises if $\hat{Q} \hat{P}$ has repeated eigenvalues.

Theorem 3.1. Suppose $A$ is asymptotically stable and $\left(A_{e}, B_{e}, C_{e}, D_{e}\right) \in \mathscr{S}^{+}$solves the Optimal ReducedOrder State-Estimation Problem. Then there exist $n \times n$ nonnegative-definite matrices $Q, \hat{Q}$ and $\hat{P}$ such that $A_{e}, B_{e}, C_{-}$and $D_{e}$ are given by

$$
\begin{align*}
& A_{\mathrm{e}}=\Gamma\left(A-Q_{\mathrm{s}} V_{2 \mathrm{~s}}^{-1} C\right) G^{\mathrm{T}},  \tag{3.10}\\
& B_{\mathrm{e}}=\Gamma Q_{\mathrm{s}} V_{2 \mathrm{~s}}^{-1},  \tag{3.11}\\
& C_{\mathrm{e}}=\left(L-D_{\mathrm{e}} C\right) G^{\mathrm{T}},  \tag{3.12}\\
& D_{\mathrm{e}}=L Q C^{\mathrm{T}} V_{2 \mathrm{~s}}^{-1}, \tag{3.13}
\end{align*}
$$

and such that $Q, \hat{Q}$ and $\hat{P}$ satisfy

$$
\begin{align*}
& Q=A Q A^{\mathrm{T}}+\sum_{i=1}^{p} A_{i}\left(Q+\tau \hat{Q}^{\mathrm{T}}\right) A_{i}^{\mathrm{T}}+V_{1}-Q_{\mathrm{s}} V_{2 \mathrm{~s}}^{-1} Q_{\mathrm{s}}^{\mathrm{T}}+\tau_{\perp} \hat{Q} \tau_{\perp}^{\mathrm{T}},  \tag{3.14}\\
& \hat{Q}=A \tau \hat{Q} \tau^{\mathrm{T}} A^{\mathrm{T}}+Q_{\mathrm{s}} V_{2 \mathrm{~s}}^{-1} Q_{\mathrm{s}}^{\mathrm{T}},  \tag{3.15}\\
& \hat{P}=\left(A-Q_{\mathrm{s}} V_{2 \mathrm{~s}}^{-1} C\right)^{\mathrm{T}} \tau^{\mathrm{T}} \hat{P} \tau\left(A-Q_{\mathrm{s}} V_{2 \mathrm{~s}}^{-1} C\right)+\left(L-D_{\mathrm{e}} C\right)^{\mathrm{T}} R\left(L-D_{\mathrm{e}} C\right), \tag{3.16}
\end{align*}
$$

where

$$
\begin{equation*}
\tau \triangleq \sum_{i=1}^{n_{e}} \pi_{i}(\psi) \tag{3.17}
\end{equation*}
$$

for some $\psi \in \mathscr{D}(\hat{Q}, \hat{P})$ such that $\left(\psi^{-1} \hat{Q} \hat{P} \psi\right)_{(, i, i)} \neq 0, i=1, \ldots, n_{\mathrm{e}}$, and some projective factorization $G, \Gamma$ of $\tau$. Furthermore, the minimal cost is given by

$$
\begin{equation*}
J\left(A_{\mathrm{e}}, B_{\mathrm{e}}, C_{\mathrm{e}}, D_{\mathrm{e}}\right)=\operatorname{tr}\left[\left(L Q L^{\mathrm{T}}-D_{\mathrm{e}} V_{2 \mathrm{~s}} D_{\mathrm{e}}^{\mathrm{T}}\right) R\right] . \tag{3.18}
\end{equation*}
$$

Remark 3.1. It is useful to note that (3.10) can be replaced by

$$
\begin{equation*}
A_{\mathrm{e}}=\Gamma A G^{\mathrm{T}}-B_{\mathrm{e}} C G^{\mathrm{T}} \tag{3.10}
\end{equation*}
$$

Remark 3.2. Because of (3.9) the $n \times n$ matrix $\tau$ which couples the three equations (3.14)-(3.16) is idempotent, i.e., $\tau^{2}=\tau$. In general, this 'optimal projection' is an oblique projection (as opposed to an orthogonal projection) since it is not necessarily symmetric. It should be stressed that the form of the optimal reduced-order estimator (3.10)-(3.13) is a direct consequence of optimality and not the result of an a priori assumption on the structure of the reduced-order estimator.

Remark 3.3. To specialize the result to the strictly proper (no feedthrough) case, merely ignore (3.13) and set $D_{\mathrm{e}}=0$ wherever it appears.

Remark 3.4. Replacing $x_{e}$ by $S x_{e}$, where $S$ is invertible, yields the 'equivalent' estimator $\left(S A_{e} S^{-1}, S B_{e}, C_{e} S^{-1}, D_{e}\right)$ with $J\left(S A_{e} S^{-1}, S B_{e}, C_{e} S^{-1}, D_{e}\right)=J\left(A_{e}, B_{e}, C_{e}, D_{e}\right)$. Note that transformation of the estimator state basis corresponds to the alternative factorization $\tau=\left(S^{-\mathrm{T}} G\right)^{\mathrm{T}}(S \Gamma)$.

Remark 3.5. Note that for the optimal values of $A_{e}$ and $B_{e}$ the estimator dynamics (3.3) assume the usual observer form

$$
\begin{equation*}
x_{\mathrm{e}}(k+1)=\Gamma A G^{\mathrm{T}} x_{e}+\Gamma Q_{s} V_{2 s}^{-1}\left(y-C G^{\mathrm{T}} x_{\mathrm{e}}\right) \tag{3.19}
\end{equation*}
$$

By introducing the quasi-full-state estimate $\hat{x} \triangleq G^{\top} x_{e} \in \mathbb{R}^{n}$ so that $\tau \hat{x}=\hat{x}$ and $x_{e}=\Gamma \hat{x} \in \mathbb{P}^{n c}$, (3.19) can be written as

$$
\begin{equation*}
\hat{x}(k+1)=\tau A \tau \hat{x}+\tau Q_{\mathrm{s}} V_{2 \mathrm{~s}}^{-1}(y-C \hat{x}) . \tag{3.20}
\end{equation*}
$$

Although the implemented estimator (3.19) has the state $x_{e} \in \boldsymbol{R}^{n_{e}}$, (3.19) can be viewed as a quasi-full-order estimator whose geometric structure is entirely dictated by the projection $\tau$ and the stochastic effects. Specifically, error inputs $Q_{\mathrm{s}} V_{2 \mathrm{~s}}^{-1}(y-C \hat{x})$ are annihilated unless they are contained in $[\mathscr{N}(\tau)]^{\perp}=\mathscr{R}\left(\tau^{\mathrm{T}}\right)$. Hence, the observation subspace of the estimator is precisely $\mathscr{R}\left(\tau^{\mathrm{T}}\right)$.

Specializing Theorem 3.1 to the noise-free case $A_{i}=0, C_{i}=0, i=1, \ldots, p$, yields Theorem 2.2 of [2,3]. Alternatively, specializing Theorem 3.1 to the full-order case $n_{e}=n$ reveals that the Lyapunov equation for $\hat{\boldsymbol{P}}$ is superfluous. In this case it follows from Remark 3.4 that $G=\Gamma=I_{n}$ without loss of generality.

Corollary 3.1. Assume $n_{e}=n, A$ is asymptotically stable and $\left(A_{e}, B_{e}, C_{e}, D_{e}\right) \in \mathscr{S}^{+}$solves the Optimal Full-Order State-Estimation Problem. Then there exist $n \times n$ nonnegative-definite matrices $Q$ and $\hat{Q}$ such that $A_{e}, B_{e}, C_{e}$ and $D_{e}$ are given by

$$
\begin{align*}
& A_{\mathrm{e}}=A-\hat{Q}_{\mathrm{s}} \hat{V}_{2 \mathrm{~s}}^{-1} C  \tag{3.21}\\
& B_{\mathrm{e}}=\hat{Q}_{\mathrm{s}} \hat{V}_{2 \mathrm{~s}}^{1}  \tag{3.22}\\
& C_{\mathrm{e}}=L-D_{e} C  \tag{3.23}\\
& D_{\mathrm{e}}=L Q C^{\mathrm{T}} \hat{V}_{2 \mathrm{~s}}^{-1} \tag{3.24}
\end{align*}
$$

and such that $Q$ and $\hat{Q}$ satisfy

$$
\begin{align*}
& Q=A Q A^{\mathrm{T}}+\sum_{i=1}^{p} A_{i}(Q+\hat{Q}) A_{i}^{\mathrm{T}}+V_{1}-\hat{Q}_{\mathrm{s}} \hat{V}_{2 \mathrm{~s}}^{-1} \hat{Q}_{\mathrm{s}}^{\mathrm{T}}  \tag{3.25}\\
& \hat{Q}=A \hat{Q} A^{\mathrm{T}}+\hat{Q}_{\mathrm{s}} \hat{V}_{2 \mathrm{~s}}^{-1} \hat{Q}_{\mathrm{s}}^{\mathrm{T}} \tag{3.26}
\end{align*}
$$

Furthermore, the minimal cost is given by

$$
\begin{equation*}
J\left(A_{e}, B_{e}, C_{e}, D_{e}\right)=\operatorname{tr}\left[\left(L Q L^{\mathrm{T}}-D_{e} \hat{V}_{2 \mathrm{~s}} D_{e}^{\mathrm{T}}\right) R\right] . \tag{3.27}
\end{equation*}
$$

Remark 3.6. To recover the standard Kalman filter result from Corollary 3.1 set $A_{i}=0, C_{i}=0$. $i=1, \ldots, p$, so that (3.25) and (3.26) are decoupled and (3.26) is superfluous. Since the standard Kalman filter is strictly proper, set $D_{\mathrm{e}}=0$ as in Remark 3.3.

## 4. Proof of the main theorem

Using the notation of Section 2 the augmented system (3.1) and (3.3) can be written as

$$
\begin{equation*}
\tilde{x}(k+1)=\left(\tilde{A}+\sum_{i=1}^{P} v_{i}(k) \tilde{A_{i}}\right) \tilde{x}(k)+\tilde{w}(k) \tag{4.1}
\end{equation*}
$$

where $\left.\tilde{x}(k) \Delta x^{\mathrm{T}}(k), x_{\mathrm{e}}^{\mathrm{T}}(k)\right]^{\mathrm{T}}$. To analyze (4.1) it is useful to define the second-moment matrix

$$
\begin{equation*}
\tilde{Q}(k)=\mathrm{E}\left[\tilde{x}(k) \tilde{x}^{\mathrm{T}}(k)\right] . \tag{4.2}
\end{equation*}
$$

It follows from (4.1) and (4.2) that $\tilde{Q}(k)$ must satisfy

$$
\begin{equation*}
\tilde{Q}(k+1)=\tilde{A} \tilde{Q}(k) \tilde{A}^{\mathrm{T}}+\sum_{i=1}^{p} \tilde{A}_{i} \tilde{Q}(k) \tilde{A}_{i}^{\mathrm{T}}+\tilde{V} \tag{4.3}
\end{equation*}
$$

Lemma 4.1. $A_{e}$ is asymptotically stable if and only if

$$
\mathscr{A} \triangleq \tilde{A} \otimes \tilde{A}+\sum_{i=1}^{p} \tilde{A}_{i} \otimes \tilde{A}_{i}
$$

is asymptotically stable.
Proof. The result follows from properties of the Kronecker product applied to partitioned matrices. See [36] for details.

Hence $\mathscr{A}$ stable assures

$$
\tilde{Q}^{\triangleq} \lim _{k \rightarrow \infty} \mathrm{E}\left[\tilde{x}(k) \tilde{x}^{\mathrm{T}}(k)\right]
$$

exists. Furthermore, $\tilde{Q}$ and its nonnegative-definite dual $\tilde{P}$ are the unique solutions of the modified Lyapunov equations

$$
\begin{align*}
& \tilde{Q}=\tilde{A} \tilde{Q} \bar{A}^{\top}+\sum_{i=1}^{p} \tilde{A}_{i} \tilde{Q} \bar{A}_{i}^{\top}+\tilde{V},  \tag{4.4}\\
& \tilde{P}=\tilde{A}^{\top} \tilde{P} \tilde{A}+\sum_{i=1}^{p} \tilde{A}_{i}^{\tilde{P} \tilde{P} \tilde{A}_{i}+\tilde{R} .} \tag{4.5}
\end{align*}
$$

Partition $\left(n+n_{e}\right) \times\left(n+n_{e}\right) \tilde{Q}, \tilde{P}$ into $n \times n, n \times n_{e}$, and $n_{e} \times n_{e}$ subblocks as

$$
\tilde{Q}=\left[\begin{array}{cc}
Q_{1} & Q_{12} \\
Q_{12}^{\top} & Q_{2}
\end{array}\right], \quad \tilde{P}=\left[\begin{array}{cc}
P_{1} & P_{12} \\
P_{12}^{\top} & P_{2}
\end{array}\right]
$$

and define the $n \times n$ nonnegative-definite matrices

$$
\begin{array}{ll}
Q \triangleq Q_{1}-Q_{12} Q_{2}^{-1} Q_{12}^{\mathrm{T}}, & \hat{\hat{Q}} \triangleq Q_{12} Q_{2}^{-1} Q_{12}^{\mathrm{T}}, \quad P \triangleq P_{1}-P_{12} P_{2}^{-1} P_{12}^{\mathrm{T}}, \quad \hat{\hat{P}} \triangleq P_{12} P_{2}^{-1} P_{12}^{\mathrm{T}}, \\
\hat{Q} \triangleq A \hat{\hat{Q}} A^{\top}+Q_{s} V_{2 \mathrm{~s}}^{-1} Q_{\mathrm{s}}^{\mathrm{T}}, & \hat{P} \triangleq\left(A-Q_{\mathrm{s}} V_{2 \mathrm{~s}}^{-1} C\right)^{\mathrm{T}} \hat{\hat{P}}\left(A-Q_{\mathrm{s}} V_{2 \mathrm{~s}}^{-1} C\right)+\left(L-D_{\mathrm{c}} C\right)^{\top} R\left(L-D_{\mathrm{e}} C\right),
\end{array}
$$

where $\tau \hat{\hat{Q}} \tau^{\top}$ is replaced by $\hat{\hat{Q}}$ in $Q_{\mathrm{s}}$ and $V_{2 \mathrm{~s}}$ and the $n_{\mathrm{e}} \times n, n_{e} \times n_{\mathrm{e}}, n_{\mathrm{e}} \times n$ matrices

$$
G \triangleq Q_{2}^{-1} Q_{12}^{\mathrm{T}}, \quad M \triangleq Q_{2} P_{2}, \quad \Gamma \triangleq-P_{2}^{-1} P_{12}^{\mathrm{T}} .
$$

To minimize (3.5) subject to the constraint (4.4), form the Lagrangian

$$
\mathscr{L}\left(A_{\mathrm{e}}, B_{\mathrm{e}}, C_{\mathrm{e}}, D_{\mathrm{e}}, \tilde{Q}, \tilde{P}, \lambda\right) \triangleq \operatorname{tr}\left[\lambda J\left(A_{\mathrm{e}}, B_{\mathrm{e}}, C_{\mathrm{e}}, D_{e}\right)+\left(\tilde{A} \tilde{Q} \tilde{A}^{\mathrm{T}}+\sum_{i=1}^{p} \tilde{A}_{i} \tilde{Q} \tilde{A}_{i}^{\top}+\tilde{V}-\tilde{Q}\right) \tilde{P}\right]
$$

where the Lagrange multipliers $\lambda \geqslant 0$ and $\tilde{P} \in \mathbf{R}^{\left(n+n_{0}\right) \times\left(n+n_{0}\right)}$ are not both zero. Setting $\partial \mathscr{Q} / \partial \tilde{Q}=0, \lambda=0$
implies $\tilde{P}=0$ since $\left(A_{e}, B_{e}, C_{e}, D_{e}\right) \in \mathscr{P}^{+}$. Hence, without loss of generality set $\lambda=1$. Thus the stationarity conditions are given by

$$
\begin{align*}
& \frac{\partial \mathscr{L}}{\partial \tilde{P}}=\tilde{A} \tilde{Q} \tilde{A}^{\mathrm{T}}+\sum_{i=1}^{p} \tilde{A}_{i} \tilde{Q} \tilde{A}_{i}^{\mathrm{T}}+\tilde{V}-\tilde{Q}=0,  \tag{4.6}\\
& \frac{\partial \mathscr{L}}{\partial \tilde{Q}}=\tilde{A}^{\mathrm{T}} \tilde{P} \tilde{A}+\sum_{i=1}^{p} \tilde{A}_{i}^{\mathrm{T}} \tilde{P} \tilde{A}_{i}+\tilde{R}-\tilde{P}=0,  \tag{4.7}\\
& \frac{\partial \mathscr{L}}{\partial A_{e}}=P_{12}^{\mathrm{T}} A Q_{12}+P_{2} B_{e} C Q_{12}+P_{2} A_{e} Q_{2}=0,  \tag{4.8}\\
& \frac{\partial \mathscr{L}}{\partial B_{e}}=P_{12}^{\mathrm{T}} Q_{s}+P_{2} B_{e} V_{2 \mathrm{~s}}=0,  \tag{4.9}\\
& \frac{\partial \mathscr{L}}{\partial C_{e}}=-R L Q_{12}+R D_{e} C Q_{12}+R C_{e} Q_{2}=0,  \tag{4.10}\\
& \frac{\partial \mathscr{L}}{\partial D_{e}}=D_{e} V_{2 s}-L Q C^{\mathrm{T}}=0 . \tag{4.11}
\end{align*}
$$

Expanding (4.6) and (4.7) yields

$$
\begin{align*}
& A(Q+\hat{\hat{Q}}) A^{\mathrm{T}}+\sum_{i=1}^{p} A_{i}(Q+\hat{\hat{Q}}) A_{i}^{\mathrm{T}}+V_{1}-Q-\hat{\hat{Q}}=0,  \tag{4.12}\\
& {\left[A \hat{\hat{Q}} A^{\mathrm{T}}+Q_{s} V_{2 s}^{-1} Q_{s}^{\mathrm{T}}-\hat{\hat{Q}}\right] \Gamma^{\mathrm{T}}=0,}  \tag{4.13}\\
& \Gamma\left[A \hat{\hat{Q}} A^{\mathrm{T}}+Q_{s} V_{2 s}^{-1} Q_{s}^{\mathrm{T}}-\hat{\hat{Q}}\right] \Gamma^{\mathrm{T}}=0,  \tag{4.14}\\
& {\left[\left(A-Q_{s} V_{2 s}^{-1} C\right)^{\mathrm{T}} \hat{\hat{P}}\left(A-Q_{s} V_{2 s}^{-1} C\right)+\left(L-D_{e} C\right)^{\mathrm{T}} R\left(L-D_{e} C\right)-\hat{\hat{P}}\right] G^{\mathrm{T}}=0,}  \tag{4.15}\\
& G\left[\left(A-Q_{s} V_{2 s}^{-1} C\right)^{\mathrm{T}} \hat{\tilde{P}}\left(A-Q_{s} V_{2 s}^{-1} C\right)+\left(L-D_{e} C\right)^{\mathrm{T}} R\left(L-D_{e} C\right)-\hat{\hat{P}}\right] G^{\mathrm{T}}=0 . \tag{4.16}
\end{align*}
$$

Note that the $(1,1)$ subblock of equation $(4.7)$ which characterizes $P_{1}$ has been omitted from the above equations since the estimator gains are independent of $P_{1}$.

Using (4.8)-(4.11) we obtain (3.10)-(3.13). Using (4.12) $\left.+G^{\top} \Gamma(4.13) G-(4.13) G-((4.13) G)\right)^{\mathrm{T}}$ and $\left.G^{\mathrm{T}} \Gamma(4.13) G-(4.13) G-((4.13) G)\right)^{\mathrm{T}}+(4.13)-(4.13)$ yields (3.14) and (3.15). Using $\Gamma^{\mathrm{T}} G(4.15) \Gamma$ (4.15) $\Gamma$ $-((4.15) \Gamma))^{T}+(4.15)-(4.15)$ yields (3.16). Finally, $\Gamma(4.13)-(4.14)$ or $G(4.15)-(4.16)$ yields $\Gamma G^{\top}=I_{n_{\theta}}$.

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# Optimal projection equations for discrete-time fixed-order dynamic compensation of linear systems with multiplicative white noise 

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The optimal projection equations for discrete-time reduced-order dynamic compensation are generalized to include the effects of state-, control- and measurementdependent noise. In addition, the discrete-time static output feedback problem with multiplicative disturbances is considered. For both problems, the design equations are presented in a concise, unified manner to facilitate their accessibility for developing numerical algorithms for practical applications.

## Notation and definitions

$R, R^{r \times s}, R^{r}, E$ real numbers, $r \times s$ real matrices, $\mathbf{R}^{r \times 1}$, expectation
$I_{n},()^{\top} \quad n \times n$ identity, transpose
$\otimes$ Kronecker product
$\tau_{\perp} \quad I_{n}-\tau, \tau \in \mathbf{R}^{n \times n}$
asymptotically
stable matrix
non-negative-
semisimple matrix
non-negative-
definite matrix
positive-definite
matrix symmetric matrix with positive eigenvalues
$n, m, l, n_{c}, p$ positive integers, $l \leqslant n_{c} \leqslant n$
$x, x_{c} \quad n$-, $n_{c}$-dimensional vectors
$u, y \quad m$-, $l$-dimensional vectors
$A, A_{i} ; B, B_{i} ; C, C_{i} \quad n \times n$ matrices, $n \times m$ matrices, $l \times n$ matrices, $i=1, \ldots, p$
$A_{\mathrm{c}}, B_{\mathrm{c}}, C_{\mathrm{c}}, D_{\mathrm{c}} n_{\mathrm{c}} \times n_{\mathrm{c}}, n_{\mathrm{c}} \times l, m \times n_{\mathrm{c}}, m \times l$ matrices
$k$ discrete-time index $1,2, \ldots$
$v_{i}(k)$ unit variance white noise, $i=1, \ldots, p$
$w_{1}(k), w_{2}(k) \quad n$-dimensional, l-dimensional white noise processes
$V_{1}, V_{2} n \times n$ covariance of $w_{1}, l \times l$ covariance of $w_{2} ; V_{1} \geqslant 0, V_{2} \geqslant 0$
$V_{12} n \times l$ cross-covariance of $w_{1}, w_{2}$
$R_{1}, R_{2}$ state and control weightings; $R_{1} \geqslant 0, R_{2} \geqslant 0$
$R_{12} n \times m$ cross weighting; $R_{1}-R_{12} R_{2}^{-1} R_{12}^{\mathrm{T}} \geqslant 0$
$\tilde{A}, \tilde{A}_{i} \quad A+B D_{\mathrm{c}} C, A_{i}+B_{i} D_{\mathrm{c}} C+B D_{\mathrm{c}} C_{i}, i=1, \ldots, p$
$\tilde{w} \quad w_{1}+B D_{c} w_{2}+\sum_{i=1}^{p} B_{i} D_{c} n_{2}$

[^35]| $\tilde{\boldsymbol{V}}$ | $V_{1}+V_{12} D_{c}^{\mathrm{T}} B^{\mathrm{T}}+B D_{c} V_{12}^{\mathrm{T}}+B D_{c} V_{2} D_{c}^{\top} B^{\top}+\sum_{i=1}^{p} B_{i} D_{c} V_{2} D_{c}^{\mathrm{T}} B_{i}^{\mathrm{T}}$ |
| :---: | :---: |
| $\tilde{R}$ | $R_{1}+R_{12} D_{c} C+C^{\top} D_{c}^{\top} R_{12}^{\top}+C^{\top} D_{c}^{\top} R_{2} D_{c} C+\sum_{i=1}^{p} C_{i}^{\top} D_{c}^{\top} R_{2} D_{c} C_{i}$ |
|  | $\left[\begin{array}{cc}\tilde{A} & B C_{\mathrm{c}} \\ B_{\mathrm{c}} C & A_{\mathrm{c}}\end{array}\right],\left[\begin{array}{cc}\bar{A}_{i} & B_{i} C_{\mathrm{c}} \\ B_{\mathrm{c}} C_{i} & 0\end{array}\right]$ |
| \% | $\left[\begin{array}{c}\check{w} \\ B_{c} w_{2}\end{array}\right]$ |
| $\underline{V}$ | $\left[\begin{array}{cc}\tilde{V} & V_{12} B_{c}^{\top}+B D_{c} V_{2} B_{c}^{\top} \\ B_{c} V_{12}^{\mathrm{T}}+B_{c} V_{2} D_{\mathrm{c}}^{\mathrm{T}} B^{\mathrm{T}} & B_{\mathrm{c}} V_{2} B_{\mathrm{c}}^{\mathrm{T}}\end{array}\right]$ |
| \% | $\left[\begin{array}{cc}\tilde{R} & R_{12} C_{\mathrm{c}}+C^{\mathrm{T}} D_{\mathrm{c}}^{\mathrm{T}} R_{2} C_{\mathrm{c}} \\ C_{\mathrm{c}}^{\mathrm{T}} R_{12}^{\mathrm{T}}+C_{\mathrm{c}}^{\mathrm{T}} R_{2} D_{\mathrm{c}} C & C_{\mathrm{c}}^{\mathrm{T}} R_{2} C_{\mathrm{c}}\end{array}\right]$ |
| $\begin{gathered} Z_{(i . j)} \\ \operatorname{tr} Z \\ \rho(Z) \end{gathered}$ | (i, $j$ ) element of matrix $Z$ trace of square matrix $Z$ rank of matrix $Z$ |
| $E_{i}$ | matrix with unity in the ( $i, i$ ) position and zeros elsewhere |
| $\Pi_{i}(\psi)$ | $\psi E_{i} \psi^{-1}$ (unit-rank eigenprojection) |

For arbitrary $n \times n Q, P, \hat{Q}, \hat{P}, \tau$ define:

$$
\begin{array}{ll}
V_{2 s} \triangleq V_{2}+C Q C^{\mathrm{T}}+\sum_{i=1}^{p} C_{i} Q C_{i}^{\top}, & R_{2 s} \triangleq R_{2}+B^{\mathrm{T}} P B+\sum_{i=1}^{p} B_{i}^{\mathrm{T}} P B_{i} \\
Q_{s} \triangleq A Q C^{\mathrm{T}}+V_{12}+\sum_{i=1}^{p} A_{i} Q C_{i}^{\mathrm{T}}, & P_{s} \triangleq B^{\mathrm{T}} P A+R_{12}^{\mathrm{T}}+\sum_{i=1}^{p} B_{i}^{\top} P A_{i} \\
Q_{s 1} \triangleq V_{12}+\sum_{i=1}^{p} A_{i} Q C_{i}^{\top}, & P_{s 1} \triangleq R_{12}^{\mathrm{T}}+\sum_{i=1}^{p} B_{i}^{\top} P A_{i} \\
\hat{V}_{2 s} \triangleq V_{2}+C Q C^{\mathrm{T}}+\sum_{i=1}^{p} C_{i}\left(Q+\tau \hat{Q} \tau^{\top}\right) C_{i}^{\top}, & \hat{R}_{2 s} \triangleq R_{2}+B^{\top} P B+\sum_{i=1}^{p} B_{i}^{\top}\left(P+\tau^{\top} \hat{P}_{\tau}\right) B_{i} \\
\hat{Q}_{5} \triangleq A Q C^{\top}+V_{12}+\sum_{i=1}^{p} A_{i}\left(Q+\tau \hat{Q} \tau^{\top}\right) C_{i}^{\top}, & \hat{P}_{s} \triangleq B^{\top} P A+R_{12}^{\top}+\sum_{i=1}^{p} B_{i}^{\top}\left(P+\tau^{\top} \hat{P}_{\tau}\right) A_{i} \\
\hat{Q}_{s 1} \triangleq V_{12}+\sum_{i=1}^{\nu} A_{i}\left(Q+\tau \hat{Q}^{\mathrm{T}}\right) C_{i}^{\top}, & \hat{P}_{s 1} \triangleq R_{12}^{\top}+\sum_{i=1}^{p} B_{i}^{\top}\left(P+\tau^{\top} \hat{P}_{\tau}\right) A_{i}
\end{array}
$$

## 1. Introduction

Hyland and Bernstein (1984) showed that the first-order necessary conditions for quadratically optimal continuous-time fixed-order dynamic compensation can be transformed into a coupled system of four matrix equations (two modified Riccati equations and two modified Lyapunov equations). The coupling is due to the presence of an oblique projection (idempotent matrix) which arises as a rigurous consequence
of optimality. This formulation provides a generalization of classical LQG control theory, since in the full-order case the projection becomes the identity matrix, the modified Lyapunov equations drop out, and the modified Riccati equations reduce to the usual LQG equations. Coupling via the optimal projection implies that sequential reduced-order design procedures consisting of either model reduction followed by controller design or controller design followed by controller reduction are generally suboptimal. Furthermore, the coupled structure of the equations yields the insight that in the reduced-order case there is no longer separation between the operations of state estimation and state-estimate feedback, i.e. the certainty equivalence principle breaks down.

The above developments for the continuous-time problem have, morcover, been carried out by Bernstein, Davis and Hyland (1986) in a discrete-time setting. As in the continuous-time case, the optimal reduced-order compensator is characterized by a pair of modified Riccati equations and a pair of modified Lyapunov equations coupled by an oblique projection. Furthermore, because of the discrete-time setting it is now possible to permit static feedthrough gains in both the full- and reduced-order controller designs. As pointed out by Hyland and Bernstein (1984), non-singular control weighting and measurement noise in the continuous-time case permit only a purely dynamic (strictly proper) controller. Note that this is precisely the case in continuous-time LQG theory, which yields strictly proper feedback controllers.

An immediate application of the discrete-time results is a rigorous treatment of the linear-quadratic sampled-data reduced-order dynamic-compensation problem (Bernstein, Davis and Greeley 1986). By explicitly accounting for real-time computational delay in the feedback loop, the sampled-data control-design problem can be transformed into an equivalent discrete-time problem. The dimension of the equivalent discrete-time system, however, is augmented by the available measurements which are treated as delay states. The optimal projection equations for discrete-time fixed-order dynamic compensation can thus be used to obtain controllers of tractably low dimension in spite of dimension augmentation.

Design considerations concerning stability and performance robustness with respect to unknown parameter variations can also be incorporated into the fixedorder dynamic-compensation design process. This can be accomplished by introducing white noise inte the plant via the imperfectly known parameters (Bernstein and Hyland 1985. Bernstein and Greeley 1986 a). Intuitively speaking, the quadratically optimal feedback controller designed in the presence of such disturbances is automatically desensitized to actual parameter variations. As shown by Bernstein and Greeley (1986 b), the modification of the closed-loop covariance equation due to multiplicative noise can be used to guarantee robust stability and performance by means of a Lyapunov function and a performance bound.

An interesting aspect of the design equations for the multiplicative noise model is the breakdown of the separation principle even in the full-order case. That is, even when coupling due to the oblique projection is absent, coupling due to stochastic effects remains. This is a graphic portrayal of observations made previously, e.g. by Gustafson and Speyer (1975). An alternative, apparently suboptimal approach involving certainty equivalent controllers for guaranteeing stochastic stability was developed by Yaz (1986).

The purpose of the present paper is to extend the optimal projection equations for fixed-order discrete-time dynamic compensation given by Bernstein, Davis and Hyland (1985) to include the effects of state-, control- and measurement-dependent
white noise. The main result (Theorem 3.1) presents the necessary conditions for optimality as a system of four matrix equations (two modified discrete-time Riccati equations and two modified discrete-time Lyapunov equations) coupled by both the optimal projection and stochastic effects. For the sake of completeness, the optimality conditions for discrete-time static output feedback are given by Theorem 2.1.

## 2. Static output feedback

2.1. Discrete-time static output-feedback problem

Given the controlled system

$$
\begin{gather*}
x(k+1)=\left(A+\sum_{i=1}^{p} v_{i}(k) A_{i}\right) x(k)+\left(B+\sum_{i=1}^{p} v_{i}(k) B_{i}\right) u(k)+w_{1}(k)  \tag{1}\\
y(k)=\left(C+\sum_{i=1}^{p} v_{i}(k) C_{i}\right) x(k)+w_{2}(k) \tag{2}
\end{gather*}
$$

where $k=1,2, \ldots$, determine $D_{c}$ such that the static output feedback law

$$
\begin{equation*}
u(k)=D_{c} y(k) \tag{3}
\end{equation*}
$$

minimizes the performance criterion

$$
\begin{equation*}
J \triangleq \lim _{k \rightarrow \infty} \mathrm{E}\left[x^{\mathrm{T}}(k) R_{1} x(k)+2 x^{\mathrm{T}}(k) R_{12} u(k)+u^{\mathrm{T}}(k) R_{2} u(k)\right] \tag{4}
\end{equation*}
$$

Using the notation given at the beginning of this paper, the closed-loop system (1)-(3) can be written as

$$
\begin{equation*}
x(k+1)=\left(\tilde{A}+\sum_{i=1}^{p} v_{i}(k) \tilde{A}_{i}\right) x(k)+\tilde{w}(k) \tag{5}
\end{equation*}
$$

Define the second-moment matrix

$$
\begin{equation*}
Q(k)=E\left[x(k) x^{\top}(k)\right] \tag{6}
\end{equation*}
$$

satisfying

$$
\begin{equation*}
Q(k+1)=\tilde{A} Q(k) \bar{A}^{\top}+\sum_{i=1}^{p} \bar{A}_{i} Q(k) \bar{A}_{i}^{\top}+\tilde{V} \tag{7}
\end{equation*}
$$

To consider the steady state, we restrict our consideration to the set of second-moment stabilizing gains

$$
\mathbf{S} \triangleq\left\{D_{\mathrm{c}}: \tilde{A} \otimes \tilde{A}+\sum_{i=1}^{p} \tilde{A}_{i} \otimes \tilde{A}_{i} \text { is asymptotically stable }\right\}
$$

The requirement $D_{c} \in S$ implies the existence of the steady-state closed-loop state covariance $Q \triangleq \lim _{k \rightarrow \infty} Q(k)$. Furthermore, $Q$ and its non-negative-definite dual $P$ are unique solutions of the modified discrete-time Lyapunov equations

$$
\begin{align*}
& Q=\tilde{A} Q \tilde{A}^{\mathrm{T}}+\sum_{i=1}^{p} \tilde{A}_{i} Q \tilde{A}_{i}^{\mathrm{T}}+\tilde{V}  \tag{8}\\
& P=\tilde{A}^{\mathrm{T}} P \tilde{A}+\sum_{i=1}^{p} \tilde{A}_{i}^{\mathrm{T}} P \tilde{A}_{i}+\tilde{R} \tag{9}
\end{align*}
$$

An additional technical requirement is that $D_{\mathrm{c}}$ be confined to the set

$$
\mathbf{S}^{+} \triangleq\left\{D_{c} \in \mathbf{S}: V_{2 s}>0 \quad \text { and } \quad R_{2 s}>0\right\}
$$

In order to obtain closed-form expressions for extremal values of the closed-loop control gains, the static- and dynamic-compensation problems require the technical assumption

$$
\begin{equation*}
\left[B_{i} \neq 0 \Rightarrow C_{i}=0\right], \quad i=1, \ldots, p \tag{10}
\end{equation*}
$$

i.e. for each $i \in\{1, \ldots, p\}, B_{i}$ and $C_{i}$ are not both non-zero. Essentially, (10) expresses the condition that the control-dependent and measurement-dependent noises are independent. There are no constraints, however, on correlation with the statedependent noise. By optimizing (4) with respect to $D_{c}$ and manipulating (8) and (9), we obtain the following result.

## Theorem 2.1

Suppose $D_{c} \in \mathbf{S}^{+}$solves the discrete-time static output-feedback problem. Then there exist $n \times n Q, P \geqslant 0$ such that

$$
\begin{equation*}
D_{c}=-R_{2 s}^{-1}\left[B^{\mathrm{T}} P A Q C^{\mathrm{T}}+P_{s 1} Q C^{\mathrm{T}}+B^{\mathrm{T}} P Q_{\mathrm{s} 1}\right] V_{2 s}^{-1} \tag{11}
\end{equation*}
$$

and such that $Q$ and $P$ satisfy

$$
\begin{align*}
Q= & A Q A^{\mathrm{T}}+V_{1}+\sum_{i=1}^{p}\left[\left(A_{i}+B_{i} D_{c} C\right) Q\left(A_{i}+B_{i} D_{c} C\right)^{\mathrm{T}}+B_{i} D_{\mathrm{c}} V_{2} D_{\mathrm{c}}^{\mathrm{T}} B_{i}^{\mathrm{T}}\right] \\
& +\left(Q_{\mathrm{s}}+B D_{\mathrm{c}} V_{2 s}\right) V_{2 s}^{-1}\left(Q_{s}+B D_{c} V_{2 s}\right)^{\mathrm{T}}-Q_{s} V_{2 s} Q_{s}^{\mathrm{T}}  \tag{12}\\
P= & A^{\mathrm{T}} P A+R_{1}+\sum_{i=1}^{p}\left[\left(A_{i}+B D_{c} \dot{C}_{i}\right)^{\mathrm{T}} P\left(A_{i}+B D_{c} C_{i}\right)+C_{i}^{\mathrm{T}} D_{\mathrm{c}}^{\mathrm{T}} R_{2} D_{\mathrm{c}} C_{i}\right] \\
& +\left(P_{s}+R_{2 s} D_{c} C\right)^{\mathrm{T}} R_{2 s}^{-1}\left(P_{s}+R_{2 s} D_{c} C\right)-P_{s}^{\mathrm{T}} R_{2 s} P_{s} \tag{13}
\end{align*}
$$

## 3. Dynamic output feedback

We now expand the formulation of the static problem to include a dynamic compensator.

### 3.1. Discrete-time dynamic output-feedback problem

Given the controlled system (1), (2), determine $A_{\mathfrak{c}}, B_{c}, C_{c}, D_{c}$ such that the dynamic output-feedback law

$$
\begin{gather*}
x_{\mathbf{c}}(k+1)=A_{\mathbf{c}} x_{\mathrm{c}}(k)+B_{\mathrm{c}} y(k)  \tag{14}\\
u(k)=C_{\mathrm{c}} x_{\mathbf{c}}(k)+D_{\mathbf{c}} y(k) \tag{15}
\end{gather*}
$$

minimizes the performance criterion (4).
We restrict our attention to the second-moment-stabilizing controllers

$$
\begin{aligned}
& \mathbf{S} \triangleq\left\{\left(A_{c}, B_{c}, C_{c}, D_{c}\right): A^{*} \otimes{ }^{z}+\sum_{i=1}^{p} A_{i} \otimes \tilde{A}_{i}\right. \text { is asymptotically stable and } \\
& \left.\quad\left(A_{c}, B_{c}, C_{c}\right) \text { is minimal }\right\}
\end{aligned}
$$

which implies the existence of $\tilde{Q} \triangleq \lim _{k \rightarrow x} E\left[\tilde{x}(k) \tilde{x}^{\top}(k)\right]$, where $\tilde{x}(k) \triangleq\left[x^{\top}(k), x_{c}^{\top}(k)\right]^{\top}$.

Furthermore, $\bar{Q}$ and its non-negative-definite dual $\tilde{P}$ are the unique solutions to the modified discrete-time Lyapu.iov equations

$$
\begin{align*}
& \tilde{Q}=\bar{A} \bar{Q} \tilde{A}^{\boldsymbol{z}}{ }^{\top}+\sum_{i=1}^{p} \tilde{A}_{i} \tilde{Q}^{\tilde{z}} \tilde{A}_{i}^{\top}+\tilde{V} \tag{16}
\end{align*}
$$

An additional technical assumption is that $\left(A_{\mathrm{c}}, B_{\mathrm{c}}, C_{\mathrm{c}}, D_{\mathrm{c}}\right)$ be confined to the set

$$
\mathbf{S}^{+\Delta} \triangleq\left\{\left(A_{\mathrm{c}}, B_{\mathrm{c}}, C_{\mathrm{c}}, D_{\mathrm{c}}\right) \in \mathbf{S}: \hat{R}_{2 s}>0 \quad \text { and } \quad \hat{V}_{2 s}>0\right\}
$$

The following lemma is required.
Lemma 3.1
Let $\tau \in \mathbf{R}^{n \times n}$. Then

$$
\begin{align*}
\tau^{2} & =\tau  \tag{18}\\
\rho(\tau) & =n_{c} \tag{19}
\end{align*}
$$

if and only if there exist $G, \Gamma \in \mathbf{R}^{n \times n}$ such that

$$
\begin{align*}
G^{\mathrm{T}} \Gamma & =\tau  \tag{20}\\
\Gamma G^{\top} & =I_{n_{c}} \tag{21}
\end{align*}
$$

Proof
See Bernstein, Davis and Hyland (1986).
For convenience call $G$ and $\Gamma$ satisfying (20) and (21) a projectice factorization of $\tau$. Furthermore, for $n \times n$ non-negative-definite matrices $\hat{Q}, \hat{P}$. define the set of contragrediently diagonalizing transformations (see Rao and Mitra 1971. p. 123)

$$
\mathbf{D}(\hat{Q}, \hat{P}) \geqslant\left\{\psi \in \mathbf{R}^{n \times m}: \psi^{-1} \hat{Q} \psi^{-\top} \text { and } \psi^{\top} \hat{P} \psi \text { are diagonal }\right\}_{;}
$$

Theorem 3.1
Suppose $\left(A_{\mathrm{c}}, B_{\mathrm{c}}, C_{\mathrm{c}}, D_{\mathrm{c}}\right) \in \mathbf{S}^{+}$solves the discrete-time dynamic output-feedback problem. Then there exist $n \times n Q, P, \hat{Q}, \hat{P} \geqslant 0$ such that

$$
\begin{align*}
& A_{\mathrm{c}}=\Gamma\left[A-B \hat{R}_{2 s}^{-1} \hat{P}_{s}-\hat{Q}_{s} \hat{V}_{2 s}^{-1} C-B D_{\mathrm{c}} C\right] G^{\top}  \tag{22}\\
& B_{\mathrm{c}}=\Gamma\left[\hat{Q}_{s} \hat{V}_{2 s}^{-1}+B D_{\mathrm{c}}\right]  \tag{23}\\
& C_{\mathrm{c}}=-\left[\hat{R}_{2 s}^{-1} \hat{P}_{s}+D_{\mathrm{c}} C\right] G^{\top}  \tag{24}\\
& D_{\mathrm{c}}=-\hat{R}_{2 s}^{-1}\left[B^{\top} P A Q C^{\top}+\hat{P}_{s 1} Q C^{\top}+B^{\top} P \hat{Q}_{s 1}\right] \hat{V}_{2 s}^{-1} \tag{25}
\end{align*}
$$

and such that $Q, P, \hat{Q}, \hat{P}$ satisfy

$$
\begin{align*}
Q= & A Q A^{\top}+V_{1}+\tau_{1} \hat{Q}_{\tau}^{\top}+\sum_{i=1}^{\rho}\left[\left(A_{i}-B_{i} \hat{R}_{2 s}^{-1} \hat{P}_{s}\right) \tau \hat{Q}^{\top} \tau^{\top}\left(A_{i}-B_{i} \hat{R}_{2 s}^{-1} \hat{P}_{s}\right)^{\top}\right. \\
& \left.+\left(A_{i}+B_{i} D_{c} C\right) Q\left(A_{i}+B_{i} D_{c} C\right)^{\top}+B_{i} D_{c} V_{2} D_{c}^{\top} B_{i}^{\top}\right]-\hat{Q}_{s} \hat{V}_{2 s}^{-1} \hat{Q}_{s}^{\top} \tag{26}
\end{align*}
$$

$$
\begin{align*}
P= & A^{\mathrm{T}} P A+R_{1}+\tau_{1}^{\mathrm{T}} \hat{P}_{\tau_{1}}+\sum_{i=1}^{p}\left[\left(A_{i}-\hat{Q}_{s} \hat{V}_{2 s}^{-1} C_{i}\right)^{\mathrm{T}} \tau^{\mathrm{T}} \hat{P}^{2}\left(A_{i}-\hat{Q}_{s} \hat{V}_{2 s}^{-1} C_{i}\right)\right. \\
& \left.+\left(A_{i}+B D_{c} C_{i}\right)^{\mathrm{T}} P\left(A_{i}+B D_{c} C_{i}\right)+C_{i}^{\mathrm{T}} D_{c}^{\mathrm{T}} R_{2} D_{c} C_{i}\right]-\hat{P}_{s}^{\mathrm{T}} \hat{R}_{2 s}^{-1} \hat{P}_{s}  \tag{27}\\
\hat{Q}= & \left(A-B \hat{R}_{2 s}^{-1} \hat{P}_{s}\right) \tau \hat{Q}^{\mathrm{T}}\left(A-B \hat{R}_{2 s}^{-1} \hat{P}_{s}\right)^{\mathrm{T}}+\left(\hat{Q}_{s}+B D_{c} \hat{V}_{2 s}\right) \hat{V}_{2 s}^{-1}\left(\hat{Q}_{s}+B D_{c} \hat{V}_{2 s}\right)^{\mathrm{T}}  \tag{28}\\
\hat{P}= & \left(A-\hat{Q}_{s} \hat{V}_{2 s}^{-1} C\right)^{\mathrm{T}} \tau^{\mathrm{T}} \hat{P}_{\tau}\left(A-\hat{Q}_{s} \hat{V}_{2 s}^{-1} C\right)+\left(\hat{P}_{s}+\hat{R}_{2 s} D_{c} C\right)^{\mathrm{T}} \hat{R}_{2 s}^{-1}\left(\hat{P}_{s}+\hat{R}_{2 s} D_{c} C\right) \tag{29}
\end{align*}
$$

where

$$
\begin{equation*}
\tau \Delta \sum_{i=1}^{n_{c}} \Pi_{i}(\psi)=G^{\top} \Gamma \tag{30}
\end{equation*}
$$

for some $\psi \in \mathrm{D}(\hat{Q}, \hat{P})$ such that $\left(\psi^{-1} \hat{Q} \hat{P} \psi\right)_{(i, i)} \neq 0, i=1, \ldots, n_{\mathrm{c}}$, and some projective factorization $G, \Gamma$ of $\tau$.

## Remark 3.1

To specialize the result to the strictly proper (no feedthrough) case, merely ignore (25) and set $D_{c}=0$ wherever it appears.

## Remark 3.2

As previously pointed out by Bernstein. Davis and Hyland (1986), the indeterminacy in specifying the projective factorization $G, \Gamma$ satisfying (20) and (21) corresponds to an arbitrary choice of internal state-space basis for the design system ( $A_{\mathrm{c}}, B_{\mathrm{c}}, C_{\mathrm{c}}$ ).

## Remark 3.3

In the full-order case $n_{c}=n$, the projection $\tau$ becomes the identity and (28) and (29) play no role. In this case $G^{\top} \Gamma=\Gamma G^{\mathrm{T}}=I_{n}$ and thus $G$ and $\Gamma$ can be chosen to be the identity. Deleting all multiplicative white noise terms corresponding to state-, controland measurement-dependent disturbances, i.e. $A_{i}, B_{i}, C_{i}=0, i=1, \ldots, p$, and specializing further to the purely dynamic case ( $D_{c}=0$ ) yields the standard LQG result. Alternatively, setting $n_{c}<n$ and deleting the multiplicative noise terms yields the results of Bernstein, Davis and Hyland (1986).

## 4. Proof of Theorem 3.1

Partition $\left(n+n_{c}\right) \times\left(n+n_{c}\right) \bar{Q}, \tilde{P}$ into $n \times n . n \times n_{c}$, and $n_{c} \times n_{c}$ sub-blocks as

$$
\tilde{Q}=\left[\begin{array}{cc}
Q_{1} & Q_{12} \\
Q_{12}^{\top} & Q_{2}
\end{array}\right] . \quad \tilde{P}=\left[\begin{array}{cc}
P_{1} & P_{12} \\
P_{12}^{\top} & P_{2}
\end{array}\right]
$$

and define the $n \times n$ non-negative-definite matrices

$$
\begin{gathered}
Q \stackrel{\Delta}{=} Q_{1}-Q_{12} Q_{2}^{-1} Q_{12}^{\top}, \quad P \triangleq P_{1}-P_{12} P_{2}^{-1} P_{12}^{\top} \\
\hat{\hat{Q}} \triangleq Q_{12} Q_{2}^{-1} Q_{12}^{\top}, \quad \hat{\tilde{P}} \triangleq P_{12} P_{2}^{-1} P_{12}^{\top} \\
\hat{Q} \triangleq\left(A-B \hat{R}_{2 s}^{-1} \hat{P}_{s}\right) \hat{\hat{Q}}\left(A-B \hat{R}_{2 s}^{-1} \hat{P}_{s}\right)^{\top}+\left(\hat{Q}_{s}+B D_{c} \hat{V}_{2 s}\right)^{-1} \hat{V}_{2 s}^{-1}\left(\hat{Q}_{s}+B D_{c} \hat{V}_{2 s}\right)^{\top} \\
\hat{P} \triangleq\left(A-\hat{Q}_{s} \hat{V}_{2 s}^{-1} C\right)^{\top} \hat{P}\left(A-\hat{Q}_{s} \hat{V}_{2 s}^{-1} C\right)+\left(\hat{P}_{s}+\hat{R}_{2 s} D_{c} C\right)^{\top} \hat{R}_{2 s}^{-1}\left(\hat{P}_{s}+\hat{R}_{2 s} D_{c} C\right)
\end{gathered}
$$

where $\tau \hat{Q}^{\top}{ }^{\top}$ and $\tau^{\top} \hat{P}_{\tau}$ in $\hat{Q}_{s}, \hat{P}_{s}, \hat{V}_{2 s}$, and $\hat{R}_{2 s}$ are replaced by $\hat{Q}$ and $\hat{\hat{P}}$, and the $n_{c} \times n$,
$n_{c} \times n_{c}$, and $n_{c} \times n$ matrices

$$
G \triangleq Q_{2}^{-1} Q_{12}^{\top}, \quad M \triangleq Q_{2} P_{2}, \quad \Gamma \triangleq-P_{2}^{-1} P_{12}^{\top}
$$

Define the lagrangian
$\left\langle A_{\mathrm{c}}, B_{\mathrm{c}}, C_{\mathrm{c}}, \mathrm{D}_{\mathrm{c}}, \tilde{Q}, \tilde{P}, \lambda\right) \triangleq \operatorname{tr}\left[\dot{\lambda} J\left(A_{\mathrm{c}}, B_{\mathrm{c}}, C_{\mathrm{c}}, D_{\mathrm{c}}\right)+\left(\tilde{A} \tilde{Q}^{\boldsymbol{z}}{ }^{\mathrm{T}}+\sum_{i=1}^{p} \tilde{A}_{i} \tilde{Q}_{\hat{Q}}^{\boldsymbol{A}}{ }_{i}^{\mathrm{T}}+\tilde{\tilde{V}}-\tilde{Q}\right) \tilde{P}\right]$
where the Lagrange multipliers $\lambda \geqslant 0$ and $\tilde{P} \in \mathbf{R}^{\left(n+n_{c}\right) \times\left(n+n_{c}\right)}$ are not both zero. Setting $\partial L / \bar{Q}=0, \lambda=0$ implies $\tilde{P}=0$ since $\left(A_{c}, B_{c}, C_{c}, D_{c}\right) \in \mathbf{S}^{+}$. Hence, without loss of generality, set $i=1$. Thus the stationarity conditions are given by

$$
\begin{align*}
& \frac{\partial L}{\partial A_{\mathrm{c}}}=P_{12}^{\mathrm{T}} A Q_{12}+P_{12}^{\mathrm{T}} B D_{\mathrm{c}} C Q_{12}+P_{12}^{\mathrm{T}} B C_{\mathrm{c}} Q_{2}+P_{2} A_{\mathrm{c}} Q_{2}+P_{2} B_{\mathrm{c}} C Q_{12}=0  \tag{33}\\
& \frac{\partial L}{\partial B_{\mathrm{c}}}=P_{2} B_{\mathrm{c}} \hat{V}_{2 s}+P_{12}^{\mathrm{T}} \hat{Q}_{\mathrm{s}}+P_{12}^{\mathrm{T}} B D_{\mathrm{c}} \hat{V}_{2 \mathrm{~s}}=0  \tag{34}\\
& \frac{\partial L}{\partial C_{c}}=\hat{P}_{s} Q_{12}+\hat{R}_{2 s} C_{c} Q_{2}+\hat{R}_{2 s} D_{c} C Q_{12}=0  \tag{35}\\
& \frac{\partial L}{\partial D_{c}}=\hat{R}_{2 s} D_{c} \hat{V}_{2 s}+B^{\mathrm{T}} P A Q C^{\mathrm{T}}+\hat{P}_{s 1} Q C^{\mathrm{T}}+B^{\mathrm{T}} P \hat{Q}_{s 1}=0
\end{align*}
$$

Expanding (31) and (32) yields

$$
\begin{align*}
& A Q A^{\mathrm{T}}+\left(\hat{Q}_{\mathrm{s}}+B D_{\mathrm{c}} \hat{V}_{2 \mathrm{~s}}\right) \hat{V}_{2 \mathrm{~s}}^{-1}\left(\hat{Q}_{\mathrm{s}}+B D_{\mathrm{c}} \hat{V}_{2 s}\right)^{\mathrm{T}}-\hat{Q}_{s} \hat{V}_{2 \mathrm{~s}}^{-1} \hat{Q}_{\mathrm{s}}^{\mathrm{T}} \\
& +\left(A-B \hat{R}_{2 s}^{-1} \hat{P}_{s}\right) \hat{Q}\left(A-B \hat{R}_{2 s}^{-1} \hat{P}_{s}\right)^{\top} \\
& +V_{1}-Q-\hat{\hat{Q}}+\sum_{i=1}^{p}\left[\left(A_{i}-B_{i} \hat{R}_{2 s}^{-1} \hat{P}_{s}\right) \hat{\hat{Q}}\left(A_{i}-B_{i} \hat{R}_{2 s}^{-1} \hat{P}_{s}\right)^{\top}\right. \\
& \left.+\left(A_{i}+B_{i} D_{c} C\right) Q\left(A_{i}+B_{i} D_{c} C\right)^{\top}+B_{i} D_{c} V_{2} D_{c}^{\top} B_{i}^{\top}\right]=0  \tag{37}\\
& {\left[\left(\hat{Q}_{s}+B D_{c} \hat{V}_{2 s}\right) \hat{V}_{2 s}^{-1}\left(\hat{Q},+B D_{c} \hat{V}_{2 s}\right)^{\mathrm{T}}+\left(A-B \hat{R}_{2 s}^{-1} \hat{P}_{s}\right) \hat{Q}\left(A-B \hat{R}_{2 s}^{-1} \hat{P}_{s}\right)^{\mathrm{T}}-\hat{Q}\right] \Gamma^{\mathrm{T}}=0} \\
& \Gamma\left[\left(\hat{Q}_{s}+B D_{c} \hat{V}_{2 s}\right) \hat{V}_{2 s}^{-1}\left(\hat{Q}_{s}+B D_{c}+\hat{V}_{2 s}\right)^{\top}\right.  \tag{38}\\
& \left.+\left(A-B \hat{R}_{2 s}^{-1} \hat{P}_{s}\right) \hat{\hat{Q}}\left(A-B \hat{R}_{2 s}^{-1} \hat{P}_{s}\right)^{\mathrm{T}}-\hat{\hat{Q}}\right] \Gamma^{\mathrm{T}}=0  \tag{39}\\
& A^{\mathrm{T}} P A+\left(\hat{P}_{\mathrm{s}}+\hat{R}_{2 \mathrm{~s}} D_{\mathrm{c}} C\right)^{\mathrm{T}} \hat{R}_{2 \mathrm{~s}}^{-1}\left(\hat{P}_{\mathrm{s}}+\hat{R}_{2 \mathrm{~s}} D_{\mathrm{c}} C\right)-\hat{P}_{\mathrm{s}}^{\mathrm{T}} \hat{R}_{2 \mathrm{~s}}^{-1} \hat{P}_{\mathrm{s}} \\
& +\left(A-\hat{Q}_{s} \hat{V}_{2 s}^{-1} C\right)^{\mathrm{T}} \hat{\tilde{P}}\left(A-\hat{Q}_{s} \hat{V}_{2 s}^{-1} C\right) \\
& +R_{1}-P-\hat{\hat{P}}+\sum_{i=1}^{p}\left[\left(A_{i}-\hat{Q}_{s} \hat{b}_{2 s}^{-1} C_{i}\right)^{\top} \hat{\tilde{P}}\left(A_{i}-\hat{Q}_{s} \hat{V}_{2 s}^{-1} C_{i}\right)\right. \\
& \left.+\left(A_{i}+B D_{c} C_{i}\right)^{\mathrm{T}} P\left(A_{i}+B D_{c} C_{i}\right)+C_{i}^{\mathrm{T}} D_{c}^{\mathrm{T}} K_{2} D_{c} C_{i}\right]=0  \tag{40}\\
& {\left[\left(\hat{P}_{s}+\hat{R}_{2 s} D_{c} C\right)^{\top} \hat{R}_{2 s}^{-1}\left(\hat{P}_{s}+\hat{R}_{2 s} D_{c} C\right)+\left(A-\hat{Q}_{s} \hat{V}_{2 s}^{-1} C\right)^{\top} \hat{P}\left(A-\hat{Q}_{s} \hat{V}_{2 s}^{-1} C\right)-\hat{P}\right] G^{\top}=0} \tag{41}
\end{align*}
$$

$$
\begin{equation*}
G\left[\left(\hat{P}_{s}+\hat{R}_{2 s} D_{\mathrm{c}} C\right)^{T} \hat{R}_{2 s}^{-1}\left(\hat{P}_{s}+\hat{R}_{2 s} D_{\mathrm{c}} C\right)+\left(A-\hat{Q}_{s} \hat{V}_{2 s}^{-1} C\right)^{\mathrm{T}} \hat{\tilde{P}}\left(A-\hat{Q}_{s} \hat{V}_{2 s}^{-1} C\right)-\hat{P}\right] G^{\mathrm{T}}=0 \tag{42}
\end{equation*}
$$

Using (33)-(36) we obtain (22)-(25). Using (37) $+G^{T} \Gamma(38) G-(38) G-(38 G)^{T}$ and $G^{\mathrm{T}} \Gamma(38) G-(38) G-(38 G)^{\mathrm{T}}$ yields (26) and (28). Similarly, using (40) $+\Gamma^{\mathrm{T}} G(41) \Gamma$ $-(41) \Gamma-(41 \Gamma)^{\top}$ and $\Gamma^{\top} G(41) \Gamma-(41) \Gamma-(41 \Gamma)^{\top}$ we obtain (27) and (29). Also. $\Gamma(38)-(39)$ or $G(41)-(42)$ yields $\Gamma G^{\top}=I_{n_{c}}$ so that $\tau=G^{\top} \Gamma=\tau^{2}$. Finally. (39) and (42) imply $\hat{Q}=\tau \hat{Q} \tau^{\top}$ and $\hat{P}=\tau^{\top} \hat{P}_{\tau}$.

## Remark 4.1

An interesting differer cetween the above discrete-time derivation and the continuous-time derivation of Hyland and Bernstein (1984) is that the explicit gain expressions and the definition of the optimal projection arise in the reverse order.

## 5. Directions for further research

The principal application of Theorems 2.1 and 3.1 is the sampled-data problem with parameter uncertainties. Although generalization of the results of Bernstein et al. (1986) is possible, there appear to be a number of mathematical issues which arise. A related development appears in Tiedemann and De Koning (1984). A more extensive treatment of the results of the present paper can be found in Haddad (1987).

## ACKnowledgment

Dr. Bernstein was supported in part by the Air Force Office of Scientific Research under contract AFOSR F49620-86-C-0002.

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## INEQUALITIES FOR THE TRACE OF MATRIX EXPONENTIALS*

## DENNIS S. BERNSTEIN +


#### Abstract

Several inequalities involving the trace of matrix exponentials are derived. The Golden-Thompon inequality $\operatorname{tr} c^{4+B} \leq \operatorname{tr} \varepsilon^{4,} e^{B}$ for symmetric $A$ and $B$ is obtained as a special case along with the new inequalit: tr $e^{4} e^{e^{T}} \leqq \operatorname{tr} e^{4 \cdot A^{T}}$ for nonnormal $A$.


Key words. matris expunential, inequalits trace
ANIS(MOS) subject classification. 15

1. Introduction. For $n \times n$ real symmetric matrices $A$ and $b$. the Golden-Thompson inequality $[1]-[5]$ states that

$$
\begin{equation*}
\operatorname{tr} c^{t+A} \leqq \operatorname{tr} c^{4} e^{t} \tag{1.1}
\end{equation*}
$$

Reference [5] generalizes (1.1) to allow arbitrary spectral functions in place of the trace operator and provides an overview of applications in which these inequalities arise.

In contrast to (1.1). problems in linear-quadratic optimal feedback control [6] typically involve a performance functional $J$ of the form

$$
\begin{equation*}
J=\operatorname{tr} \int_{0}^{x} e^{-4 t} L^{\prime} e^{4^{2} t} R d t \tag{1.2}
\end{equation*}
$$

where $I$ and $R$ denote noise intensity and performance weighting matrices. respectively. and $A$ denotes the linear system dynamics matrix. The form of (1.2) thus suggests inequalities of the form (1.1) involving .1 and $.4^{\prime}$. Where $A$ is nonnormal, in place of symmetric 4 and $B$. Such inequalities are motivated by robust sampled-data controldesign problems which require performance bounds for uncertain system models.

The main result of the present note is the inequalits

$$
\begin{equation*}
\operatorname{tr} ،^{-1} \varepsilon^{\prime \prime} \leqq \operatorname{tr} \boldsymbol{c}^{1 \cdot} \mathbf{1}^{\prime} \tag{1.3}
\end{equation*}
$$

Rather surprisingly, the sign of the inequalit? (1.3) is opposite to the sign of (1.1). To understand whe this is the case, we derice a series of inequalitics which. upon appropriate specialization. yield hoth (1.i) and (1.3).
2. Inequalities. The following femma is required. LLef (" denote the transpose of a matrix (.)
 (2.1)

$$
\operatorname{tr} c^{\prime} \doteq \operatorname{tr} c^{\prime \prime} C^{\prime \prime}=\operatorname{tr}\left(C^{-1}\right)^{\prime}
$$

 second follows from a result of $k$. Fan (sce [4. pr. 234. 516]).

Thilortm 2.1. If.t. $B \in R^{\prime \prime}$ ", then
(2.2)

1-2)

$$
\left.\begin{array}{l}
\operatorname{tr} e^{-1} e^{H}  \tag{2.3}\\
\frac{1}{2} \operatorname{tr}\left(e^{2 A}+e^{2 B}\right)
\end{array}\right\} \leqq \frac{1}{2} \operatorname{tr}\left(e^{\prime A} e^{\prime \prime}+e^{H} e^{H^{\prime}}\right) \leqq \frac{1}{2} \operatorname{tr}\left(e^{H \cdot A^{\prime}}+e^{H+B^{\prime}}\right)
$$

[^36]Proof. Defining $C=e^{A / 2 r} e^{B / 2 r}$, (2.1) becomes

$$
\operatorname{tr}\left(e^{A / 2 r} e^{B / 2 r}\right)^{2 r} \leqq \operatorname{tr}\left(e^{A / 2 r} e^{B / 2 r}\right)^{\prime}\left(e^{B^{T} / 2 r} e^{A^{T} / 2 r}\right)^{r} \leqq \operatorname{tr}\left(e^{A / 2 r} e^{B / 2 r} e^{B^{T} / 2 r} e^{4^{T} / 2 r}\right)^{r}
$$

Letting $r \rightarrow \alpha$, the exponential product formula [5, p. 60] and its generalization [7. p. 97 ] yield the first two inequalities of (2.2). The third inequality of (2.2) follows from Corollary 3 of [5] while the fourth inequality of (2.2) follows from

$$
0 \leqq \operatorname{tr}\left[e^{\left(\cdot 1+A^{7} / 1 / 2\right.}-e^{\left(B+B^{T} / / 2\right.}\right]^{2}
$$

To prove (2.3) note that the upper leftmost inequality follows from $0 \leqq$ $\operatorname{tr}\left(e^{-4}-e^{B}\right)\left(c^{-4}-c^{B}\right)^{7}$. The remaining inequalities in (2.3) follow from $\operatorname{tr} \boldsymbol{c}^{-4} \leqq$ tr $e^{4} e^{4^{T}} \leqq \operatorname{tr} e^{A+A^{\top}}$. which is a consequence of (2.2) with $B=.4$.

Corollary 2.1. If $A \in R^{\prime \cdot n}$, then

$$
\begin{align*}
& \operatorname{tr} e^{2-t} \leqq \operatorname{tr} e^{-4} e^{4^{7}} \leqq \operatorname{tr} e^{-4+A^{7}} \leqq \frac{n}{2}+\frac{1}{2} \operatorname{tr} e^{2\left(t+A^{7}\right)},  \tag{2.4}\\
& \operatorname{tr} e^{24} \leqq \frac{n}{2}+\frac{1}{2} \operatorname{tr} e^{2.4} e^{24^{7}} \leqq \frac{n}{2}+\frac{1}{2} \operatorname{tr} e^{2\left(.4+4^{T}\right)} \tag{2.5}
\end{align*}
$$

If $A, B \in R^{n \times n}$ arc symmetric, then

$$
\begin{equation*}
\operatorname{tr} e^{4+B} \leqq \operatorname{tr} e^{A} e^{B} \leqq \frac{1}{\operatorname{tr}}\left(e^{24}+e^{2 B}\right) \tag{2.6}
\end{equation*}
$$

Proof. The first two inequalities of (2.4) follow from the first two inequalities of (2.2) with $B=A$. The last inequality of (2.4) follows from the last inequality of (2.2) with $B=0$ and $A$ replaced by 2.4. Inequalities (2.5) follow from (2.3) with $B=0$ and $A$ replaced by 2.4 while ignoring the lower leftmost term in (2.3). Finally, (2.6) follows from (2.2).

Remark. The second inequality in (2.4) and the first inequality in (2.6) correspond io (1.3) and (1.1). respectively:
3. Additional inequalities. The question immediatel arises as to whether any additional inequalities insolsing the expressions annearing in (2.4) and (2.5) are true. Nole that tr $e^{4} e^{H}$ in (2.3) cannot be merged with (2.2) because of the sign reversal between (1.1) and (1.3). It can readily the seen that the only remaining possibilities are

$$
\begin{equation*}
\operatorname{tr} e^{\prime \prime} \cdot e^{\prime} e^{\prime \mu} \cdot R^{\prime} \because \frac{1}{2} \operatorname{tr}\left(e^{\prime 2}+e^{2 H}\right) . \tag{3.1}
\end{equation*}
$$

$$
\begin{equation*}
\operatorname{tr} a^{+A} a^{\prime \prime}=\frac{1}{2} \operatorname{tr}\left(a^{2}+a^{\prime N}\right) \tag{3.2}
\end{equation*}
$$

B: randomly generating $A$ and $B$. (3.1) was shoun whe false Since (3.2) implies (3.1). (3.2) must also be false. Furthermore in the case $B^{\prime}=-B$. inequality (3.1), which becomes

$$
\begin{equation*}
\operatorname{tr} c^{\left.1+A^{2}\right)}: \frac{n}{2}+\frac{1}{2} \operatorname{tr} c^{\prime} c^{\prime \prime} \tag{3.4}
\end{equation*}
$$

was also shown to be false. Hence (2.4) and (2.5) cannot be merged. Finally, the remaining inequality (3.3) was also shown to be false even when $B=0$.

Remark. The results of this paper can be generalized to the case in which $A$ and $B$ are complex matrices. Generalization to arbitrary spectral functions [5] remains an area for further research.

Acknowledgment. I wish to thank Scott W. Greeley for carrying out numerical calculations which suggested the results of this paper.

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# PROBLEMS AND SOLUTIONS 

Edited ay MURRAY S. KLAMKIN
Collaborating, Editors: CECIL C. Rousseau otto G. RUEHR
All problems und solutions should he sent, lypewritten in duplicate, to Murray. S. Klamkin. Deparıment of Mathematics; University of Alberta, Edmonton, Alherra, Canada T6G 2GI. An asterisk placed beside' a problem number indicates that the prohlem was submitted without solution. Proposers and solvers whese' solutions are published will receive 5 reprints of the corresponding problem section. Other solvers will recerive just one reprint provided a self-addressed stamped (U.S.A. or Canada) envelope is enclosed. Proposers and solvers desiring acknowledgment of their contrihutioms should inchude a self-addressed stamped postcard (no stamps necessary ouside the L'S.A. and Canada). Solutions should be received hy. June 30, 1988.

## PROBLEMS

## Commuting Matrix Exponentials

Problem 88-1*, by Dennis S. Bernstein (Harris Corporation, Melbourne, Florida).
In feedback control theory for sampled-data systems, the equivalent discrete-time dynamics matrix is given by $e^{A h}$, where $h$ is the sample interval and $A$ is the dynamics matrix for the original continuous-time system. When $A$ is perturbed by $A_{0}$ (due possibly to some modeling uncertainty), then it is necessary to consider $e^{(1)+A_{0} h_{1}}$. For robust control system design it thus may be of interest to know when $e^{\left(1+\pi_{0}\right) h}$ can be decomposed into a nominal part involving $A$ and a perturbed part involving $A_{0}$. Analogous questions arise in the study of bilinear control systems of the form

$$
\dot{x}=A x+u B x
$$

where $u$ is a scalar control. In this case the Lie group generated by $A$ and $B$ plays a central role. Again, it is of interest to know how $\boldsymbol{c}^{4+B}$ is related to $e^{4}$ and $e^{H}$. the principal result being the Baker-Campbell-Hausdorff formula. Of course, it is well known that when

$$
\begin{equation*}
A B=B A . \tag{1}
\end{equation*}
$$

where $A, B$ are real $n \times n$ matrices, then both

$$
\begin{equation*}
c^{-4} c^{H}=e^{H} c^{\prime} \tag{2}
\end{equation*}
$$

and

$$
\begin{equation*}
c^{\prime} e^{\prime \prime}=c^{\prime+1 \prime} \tag{3}
\end{equation*}
$$

hold. It is less well known that the converse is not true. Specifically, examples are given in [1] which show that (2) may hold while (1) and (3) are violated, and that (3) may hold while (1) is violated. Interestingly, it is stated without proof in [1] that (3) implies (2). Prove this claim or find a counterexample. A copy of [1] is availabie from the proposer.

## REFERENCE

[1] M. Frechet, Les solutions non commutables de léquation matricielle $e^{x} e^{r}=\rho^{\lambda+r}$. Rend. Circ. Mat. Palermo 2, (1952), pp. 11-27.

## ROBUST STABILITY FOR SAMPLED-DATA CONTROL SYSTEMS

by

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## ABSTRACT

In this paper we consider the robust stability of a continuous-time system under computer control. The uncertainty is modeled as additive perturbations to the matrices in a continuous-time state space description of the plant. Our methods exploit the resulting exponential-like uncertainty structure in the sampled-data control system and we develop sufficient conditions for such a system to be robustly stable.

## 1. Introduction

A sampled-data control system consists of a continuous-time plant under computer control; see Figure 1. Generally speaking, if the matrices in a state space description of a continuous-time plant are uncertain, then the resulting closed-loop, discrete-time system possesses an exponential-like uncertainty structure. This is true even if the continuous-time plant has linear uncertainty. Existing methods; e.g., see [1]-[4], are inadequate in analyzing such uncertain discrete-time systems since they do not directly handle these exponential-like structures. Notable exceptions include the conic sector approach in [5], and the stochastic parameter formulation in [6]. Jndeed, the present paper was motivated by the approach of [6] which, as shown in [12] can be reinterpreted to yield conditions for deterministic robust stability. Our objective is thus to develop a robust stability test which exploits the specific nonlinear uncertainty structures occurring in sampled-data control systems.

In the sequel, the following notation will be used. For $X \varepsilon R^{n \times n}, X$ ' denotes the transpose of $X$, while $X>0(X \geq 0)$ means that $X$ is positive definite (positive semi-definite). The spectral radius of $X$ is given by $\rho(X)$. Let $\otimes, \oplus$ and "vec" denote the Kronecker product, Kronecker sum and column stacking operators respectively; see [9]. In addition we shall use "vec ${ }^{-1}$ " to denote the operation of forming a (usually square) matrix from a column vector.

## 2. Problem Formulation

In this section we state the robust analysis problem for sampled-data control systems using static feedback. To begin, consider the n-dimensional continuous-time plant

$$
\begin{align*}
& \dot{x}(t)=(A+\Delta A) x(t)+(B+\Delta B) u(t) \\
& y(t)=C x(t)
\end{align*}
$$

where $A \varepsilon R^{n \times n}, B \varepsilon R^{n \times m}$ and $C \varepsilon R^{l \times n}$ denote nominal state space data and where $\Delta A$ and $\Delta B$ represent perturbations in $A$ and $B$ respectively. The pair
of suitably dimensioned matrices ( $\Delta A, \Delta B$ ) belongs to an uncertainty set $U$ given by

$$
\begin{equation*}
\mathbf{U}=\left\{(\Delta A, \Delta B): \quad \Delta A=\sum_{i=1}^{p} \sigma_{i} A_{i} \quad \text { and } \quad \Delta B=\sum_{i=1}^{p} \sigma_{i} B_{i}, \quad \sum_{i=1}^{p} \sigma_{i}^{2} \leq 1\right\} \tag{2.1}
\end{equation*}
$$

where $A_{i}$ and $B_{i}$ reflect the "structure" of the uncertainty and where $\sigma_{i}$ is an uncertain real parameter; note that an uncertain parameter $\sigma_{i}$ may appear in both $\Delta A$ and $\Delta B$, and it's possible to have $A_{i}=0$ and $B_{i} \neq 0$ or vice-versa.

Now, consider the sampled-data system in Figure 1 with a sampling period of $h$ seconds. We assume perfect synchronization between the A/D (sampler) and the (D/A) (zero-order hold) and ignore finite word-length effects and computational delays. We also assume that a static control law

$$
\begin{equation*}
u(k h)=K y(k h) ; \quad k=0,1,2, \ldots, \tag{2.2}
\end{equation*}
$$

is implemented for some given gain $K \varepsilon R^{m \times l}$.
Our purpose is to analyze the robust stability of this closed-loop system, and to this end we consider its evolution at the sample instances kh by forming the associated discrete-time system

$$
\begin{equation*}
x(k+1)=\left[e^{(A+\Delta A) h}+\int_{0}^{h} e^{(A+\Delta A) \tau} d \tau(B+\Delta B) K C\right] x(k) \tag{2.3}
\end{equation*}
$$

In the above $x(k)$ denotes $x(k h)$; we have abused notation for the sake of conciseness. Given arbitrary ( $\Delta A, \Delta B$ ) $\varepsilon U$, system (2.3) is discrete-time stable if all the eigenvalues of

$$
e^{(A+\Delta A) h}+\int_{0}^{h} e^{(A+\Delta A) \tau} d \tau\langle B+\Delta B) K C
$$

lie within the open unit disk. Additionally, (2.3) is said to be robustly discrete-time stable if it is discrete-time stable for all ( $\Delta A, \Delta B$ ) $\varepsilon U$.

## 3. Main Result

We now develop a sufficient condition for the robust stability of (2.3). This condition exploits the exponential structure of the uncertainty; i.e., rather than "overbounding" the uncertainty with an additive model of the form

$$
x(k+1)=(\Phi+\Delta \Phi) x(k),
$$

our methods treat the exponential structure of the nonlinearity in (2.3) more directly.

To show robust stability we will construct a parameter-independent (independent of the uncertain matrices ( $\Delta A, \Delta B$ )) quadratic Lyapunov function $V(x)$ for (2.3). Thus, let $P$ be some free, positive-definite symmetric matrix and consider the quadratic Lyapunov candidate

$$
V(x)=x^{\prime} P x, \quad x \in R^{n}
$$

System (2.3) is robustly discrete-time stable if
$\Delta V(\Delta A, \Delta B, P) \triangleq$
$\left[e^{(A+\Delta A) h}+\int_{0}^{h} e^{(A+\Delta A) \tau} d \tau(B+\Delta B) K C\right]^{\prime} P\left[e^{(A+\Delta A) h}+\int_{0}^{h} e^{(A+\Delta A) \tau} d \tau(B+\Delta B) K C\right]-P$
is negative definite for all $(\Delta A, \Delta B) \varepsilon U$. This follows since $x^{\prime} \Delta V x$ is the Lyapunov difference associated with $V=x^{\prime} P x$ and (2.3); i.e.,

$$
x^{\prime} \Delta V(\Delta A, \Delta B, P) x=V(x(k+1))-V(x(k)) .
$$

System (2.3) is robustly discrete-time stable if this difference is negative for all $x \in R^{n}$ and all $(\Delta A, \Delta B) \varepsilon U$; e.g., see [7].

A critical step in our development is to express the uncertainty in (3.1) solely in terins of the matrix exponential. Indeed, using the identity (see [8])

$$
\exp \left(\left[\begin{array}{cc}
(A+\Delta A) & (B+\Delta B)  \tag{3.2}\\
0 & 0
\end{array}\right] h\right)=\left[\begin{array}{cc}
e^{(A+\Delta A) h} & \int_{0}^{h} e^{(A+\Delta A) \tau} d \tau(B+\Delta B) \\
0 & I
\end{array}\right]
$$

(3.1) can be rewritten
$\Delta V(\Delta A, \Delta B, P)=$
$\left[\begin{array}{ll}I & 0\end{array}\right] \exp \left(\left[\begin{array}{cc}(A+\Delta A) & (B+\Delta B) \\ 0 & 0\end{array}\right] h\right)\left[\begin{array}{l}I \\ K C\end{array}\right] P\left[I \quad C^{\prime} K^{\prime}\right] \exp \left(\left[\begin{array}{cc}(A+\Delta A) & (B+\Delta B) \\ 0 & 0\end{array}\right]^{\prime} h\right)\left[\begin{array}{l}I \\ 0\end{array}\right]-P$

Our next result provides a parameter-independent upper bound to $\Delta V$. To give this bound, let $\alpha>0$ be given and define

$$
\bar{A}_{\alpha} \triangleq\left(\left[\begin{array}{ll}
A & B  \tag{3.4}\\
0 & 0
\end{array}\right]+\frac{\alpha}{2} I\right) \biguplus\left(\left[\begin{array}{ll}
A & B \\
0 & 0
\end{array}\right]+\frac{\alpha}{2} I\right)+\frac{1}{\alpha} \underset{i=1}{p}\left[\begin{array}{ll}
A_{i} & B_{i} \\
0 & 0
\end{array}\right] \otimes\left[\begin{array}{ll}
A_{i} & B_{i} \\
0 & 0
\end{array}\right]
$$

and

$$
\pi_{p} \triangleq\left[\begin{array}{l}
I  \tag{3.5}\\
K C
\end{array}\right] P\left[\begin{array}{ll}
I & C^{\prime} K^{\prime}
\end{array}\right]
$$

Lemma 3.1 (See Appendix for proof): If $P \varepsilon R^{n \times n}$ is positive definite and symmetric and $\alpha>0$, then
$\Delta V(\Delta A, \Delta B, P)=$
$\left[\begin{array}{ll}I & 0\end{array}\right] \exp \left(\left[\begin{array}{cc}(A+\Delta A) & (B+\Delta B) \\ 0 & 0\end{array}\right] h\right)\left[\begin{array}{l}I \\ K C\end{array}\right] P\left[I \quad C^{\prime} K^{\prime}\right] \exp \left(\left[\begin{array}{cc}(A+\Delta A) & (B+\Delta B) \\ 0 & 0\end{array}\right]^{\prime} h\right)\left[\begin{array}{l}I \\ 0\end{array}\right]-P^{\prime}$

$$
\leq\left[\begin{array}{ll}
I & 0
\end{array}\right] \operatorname{vec}^{-1}\left[e^{\bar{A}_{\alpha} h} \operatorname{vec}\left(\Pi_{P}\right)\right]\left[\begin{array}{l}
I  \tag{3.6}\\
0
\end{array}\right]-P
$$

for all $(\Delta A, \Delta B) \varepsilon U$.

$$
\left[\begin{array}{ll}
I & 0
\end{array}\right] \operatorname{vec}^{-1}\left[e^{\bar{A}_{\alpha} h} \operatorname{vec}\left(\Pi_{P}\right)\right]\left[\begin{array}{l}
I  \tag{3.7}\\
0
\end{array}\right]-P=-I
$$

which is equivalent to

$$
\begin{equation*}
\left(M_{\alpha}-I\right) \operatorname{vec}(P)=-\operatorname{vec}(I) \tag{3.8}
\end{equation*}
$$

where

$$
M_{\alpha} \triangleq\left[\begin{array}{ll}
I & 0
\end{array}\right] e^{\bar{A}_{\alpha} h}\left(\left[\begin{array}{l}
I \\
K C
\end{array}\right] \otimes\left[\begin{array}{l}
I \\
K C
\end{array}\right]\right)
$$

When does (3.8) have a positive-definite solution? Our next lemma gives a sufficient condition.

Lemma 3.2 (See Appendix for proof): If there exists an $\alpha>0$ such that

$$
\begin{equation*}
\rho\left(M_{\alpha}\right)<1, \tag{3.10}
\end{equation*}
$$

then (3.8) has a positive-definite solution $P$.


Now, assume (3.10) holds for some $\alpha>0$. From Lemma 3.2, equation (3.7) has a positive-definite solution $P$; hence, using (3.6) and (3.7) it follows that

$$
\begin{equation*}
\Delta V(\Delta A, \Delta B, P) \leq-I \tag{3.11}
\end{equation*}
$$

for all ( $\Delta A, \Delta B$ ) $\varepsilon U$. We have thus proven the following main result.

Theorem 3.1: If there exists an $\alpha>0$ such that

$$
\rho\left(M_{\alpha}\right)<1,
$$

then the sampled-data system (2.3) is robustly discrete-time stable.

It's important to note that the condition of Theorem 3.1, $\rho\left(M_{\alpha}\right)<1$, is always satisfied if the "nominal" system is stable and there is no uncertainty; i.e., $\Delta A=0$ and $\Delta B=0$. Indeed, taking $\alpha \simeq 0$, a straightforward manipulation using (3.2), (3.4), (3.9) and identities (A.1) - (A.4) in the Appendix gives

$$
\begin{equation*}
M_{\alpha}=\left(e^{A h}+\int_{0}^{h} e^{A \tau} B K C d \tau\right) \otimes\left(e^{A h}+\int_{0}^{h} e^{A \tau} B K C d \tau\right) \tag{3.12}
\end{equation*}
$$

Since the nominal system is stable, then all the eigenvalues of

$$
e^{A h}+\int_{0}^{h} e^{A \tau_{B K C d} \tau}
$$

lie within the unit disk. This implies, together with (3.12) and the fact that the eigenvalues of a Kronecker product of two matrices are the products of the eigenvalues of these two matrices (see [9]), that all the eigenvalues of $M_{\alpha}$ lie within the unit disk. Thus, $\rho\left(M_{\alpha}\right)<1$.

The stability result of Theorem 3.1 is also valid for the case when (2.3) has time-varying uncertain parameters $\sigma_{i}(t)$. This is a consequence of having established stability via a parameter-independent lyapunov function. Finally, we remark that a dual result holds when $\Delta B=0$ and one allows uncertainty in $C$; i.e, $C \rightarrow C+\Delta C$.

## 4. Example

In Soroka and Shaked [13], the robustness of a continuous-time system under "cheap" LQ regulation is studied. The resulting closed-loop system is described by

$$
\begin{align*}
& \dot{x}(t)=\left[A+(B+\Delta B) L_{r}\right] x(t) \\
& y(t)=C x(t) \tag{4.1}
\end{align*}
$$

where

$$
\begin{gather*}
A=\left[\begin{array}{rr}
-1 & 0 \\
0 & -2
\end{array}\right] ; \quad B=\left[\begin{array}{l}
1 \\
1
\end{array}\right] ; \quad C=\left[\begin{array}{r}
-1 \\
1
\end{array}\right]^{\prime} ; \\
\Delta A=0 ; \quad \Delta B=\sigma_{1} B_{1} ; \quad B_{1}=\left[\begin{array}{c}
\bar{\sigma} \\
0
\end{array}\right] ; \quad \sigma_{1}^{2} \leq 1 ; \quad \bar{\sigma}>0 \tag{4.2}
\end{gather*}
$$

and where the $L Q$ regulator gains are given by

$$
\begin{equation*}
L_{r}=[1+q-\sqrt{5+2 q} \quad 2 \sqrt{5+2 q}-q-4]^{\prime} ; \quad q=\sqrt{4+1 / r} \tag{4.3}
\end{equation*}
$$

with $r>0$ being the control weighting in the quadratic performance index

$$
\int_{0}^{h}\left[y^{2}(t)+r u^{2}(t)\right] d t
$$

Soroka and Shaked showed that stability robustness decreased as the control became more cheap ( $r \rightarrow 0$ ).

Now, suppose this LQ regulator gain is to be implemented in a computer. What is the robustness of this sampled-data control system? We shall use the sufficient condition in Theorem 3.1 to help answer this question. First, however, we translate the continuous-time LQ gain to one suitable for sampled-data control since the nominal discrete-time system is unstable if we implement gain $K=L_{r}$ in the computer. Following [pp. 189-191, Astrom and WittenmarkI, we take

$$
\begin{equation*}
K=L_{r}\left[I+\left(A-B L_{r}\right) h / 2\right] \tag{4.4}
\end{equation*}
$$

For given sampling period $h$, control weighting $r$ and uncertainty bound $\bar{\sigma}$, we are now in a position to determine if $\rho\left(M_{\alpha}\right)<1$ for some $\alpha>0$. For example, for $h=.1, r=.05$ and $\bar{\sigma}=.5$, we plot $\rho\left(M_{\alpha}\right)$ versus $\alpha$ in Figure 2 and observe that $\rho\left(M_{\alpha}\right)<1$ for $\alpha>$. 3. Hence, the sampled-data system is stable for $\bar{\sigma}=.5$. Also, for $h=.1$ we determine, for various $r$, the largest $\bar{\sigma}$ for which $\rho\left(M_{\alpha}\right)<1$ for some $\alpha>0$. We compare these results to the actual
discrete-time stability boundary and to Soroka and Shaked's results [13]; see Figure 3. To compute the the actual discrete-time stability boundary, we assume time-invariant uncertainty $\sigma_{1}$ and find the largest $\bar{\sigma}$ for which

$$
\rho\left(e^{A h}+\int_{0}^{h} e^{A \tau} d \tau\left[\begin{array}{c}
1+\bar{\sigma}_{1}  \tag{4.5}\\
1
\end{array}\right] K\right)<1 .
$$

for all $\left|\sigma_{1}\right| \leq \bar{\sigma}$. Recall that the stability result in Theorem 3.1 is valid for time-varying uncertainties $\sigma_{1}(t)$ as well; hence, it's natural to expect that the resuts using $\rho\left(M_{\alpha}\right)$ will be conservative compared to those using (4.5). This was indeed true except for values $1 / r \simeq 10$. For these values, the stability criterion of Theorem 3.1 predicts ranges of stable $\sigma_{1}$ which are larger than those computed using (4.5). This is clearly incorrect and we contribute these discrepencies to roundoff errors in the computations. Finally, we note that the actual discrete-time stability region is larger than its continuous-time counterpart.

ACKNOWLEDGEMENT

We thank Mr. Mujdat Arabacioglu for his assistance in preparing the example.

## APPENDIX: PROOF OF LEMMAS 3.1 AND 3.2

To prove Lemmas 3.1 and 3.2 we first need some identities and observations; see [7] for details.

Identities: If $X \varepsilon \mathbb{R}^{\mathbf{r X r}}$ and $Y \varepsilon \mathbb{R}^{S \times S}$, then by definition the Kronecker sum is given by

$$
\begin{equation*}
X \oplus Y=X \Theta I+I \Theta Y \tag{A.1}
\end{equation*}
$$

Next, if the indicated products exist, then

$$
\begin{equation*}
\operatorname{vec}(X Y Z)=\left(Z^{\prime} X X\right) \operatorname{vec}(Y) \tag{A.2}
\end{equation*}
$$

$$
\begin{equation*}
(X \otimes Y)(Z \otimes W)=(X Z) \otimes(Y W) \tag{A.3}
\end{equation*}
$$

Finally, if $X$ is square, then

$$
\begin{equation*}
e^{X \oplus X}=e^{X} \otimes e^{X} \tag{A.4}
\end{equation*}
$$

Observation 1: Given arbitrary $(\Delta A, \Delta B) \varepsilon U$, the solution to the matrix differential equation

$$
\dot{Y}(t)=\left[\begin{array}{cc}
(A+\Delta A) & (B+\Delta B)  \tag{A.5}\\
0 & 0
\end{array}\right] Y(t)+Y(t)\left[\begin{array}{cc}
(A+\Delta A) & (B+\Delta B) \\
0 & 0
\end{array}\right]^{\prime}, \quad t \geq 0, \quad Y(0)=Y_{0}
$$

is

$$
Y(t)=\exp \left(\left[\begin{array}{cc}
(A+\Delta A) & (B+\Delta B)  \tag{A.6}\\
0 & 0
\end{array}\right] t\right) Y_{0} \exp \left(\left[\begin{array}{cc}
(A+\Delta A) & (B+\Delta B) \\
0 & 0
\end{array}\right]^{\prime} t\right)
$$

This is a well-known result.

Observation 2: Let $\alpha>0$, then the solution to the matrix differential equation
$\dot{\bar{Y}}(t)=\left(\left[\begin{array}{ll}A & B \\ 0 & 0\end{array}\right]+\frac{\alpha}{2} I\right) \bar{Y}(t)+\bar{Y}(t)\left(\left[\begin{array}{ll}A & B \\ 0 & 0\end{array}\right]^{\prime}+\frac{\alpha}{2} I\right)+\frac{1}{\alpha_{i}} \sum_{i=1}^{p}\left[\begin{array}{ll}A_{i} & B_{i} \\ 0 & 0\end{array}\right] \bar{Y}(t)\left[\begin{array}{ll}A_{i} & B_{i} \\ 0 & 0\end{array}\right]^{\prime}$,

$$
\begin{equation*}
t \geq 0, \quad \bar{Y}(0)=Y_{0} \tag{A.7}
\end{equation*}
$$

is

$$
\begin{equation*}
\bar{Y}(t)=\operatorname{vec}^{-1}\left[e^{\bar{A}_{\alpha}^{t}} \operatorname{vec}\left(\bar{Y}_{0}\right)\right] \tag{A.8}
\end{equation*}
$$

To show (A.8), apply the "vec" operation to both sides of (A.7) and use identities (A.1) and (A.2) to get
$\operatorname{vec}(\dot{\bar{Y}}(t))=$
$\left\{\left(\left[\begin{array}{ll}A & B \\ 0 & 0\end{array}\right]+\frac{\alpha}{2} I\right) \oplus\left(\left[\begin{array}{cc}A & B \\ 0 & 0\end{array}\right]+\frac{\alpha}{2} I\right)+\frac{1}{\alpha} \underset{i=1}{p}\left(\left[\begin{array}{ll}A_{i} & B_{i} \\ 0 & 0\end{array}\right] \otimes\left[\begin{array}{ll}A_{i} & B_{i} \\ 0 & 0\end{array}\right]\right)\right\} \operatorname{vec}(\bar{Y}(t))$.
The solution to (A.9) is

$$
\operatorname{vec}(\bar{Y}(t))=\left[\mathrm{e}^{\bar{A}_{\alpha} t} \operatorname{vec}\left(\bar{Y}_{0}\right)\right]
$$

from which (A.8) follows.
observation 3: If $\alpha>0$ and if $\overline{\mathrm{Y}}_{0}$ is positive semi-definite, then the solution in (A.8) satisfies

$$
\begin{equation*}
\bar{Y}(t)=\operatorname{vec}^{-1}\left[e^{\bar{A}_{\alpha}} \operatorname{vec}\left(\bar{Y}_{0}\right)\right] \geq 0, \quad t \geq 0 \tag{A.9}
\end{equation*}
$$

To prove this observation, let $S, N \varepsilon R^{r \times r}$ with $N$ positive semi-definite. From (A.2)

$$
(S \otimes S) \operatorname{vec}(N)=\operatorname{vec}\left(S N S^{\prime}\right)
$$

which implies

$$
\begin{equation*}
\operatorname{vec}^{-1}[(S \otimes S) \operatorname{vec}(N)]=S N S^{\prime} \geq 0 \tag{A.10}
\end{equation*}
$$

Using (A.3)

$$
\begin{equation*}
(S \otimes s)^{i}=\left(s^{i} \otimes s^{i}\right) \tag{A.11}
\end{equation*}
$$

so that

$$
\begin{align*}
\operatorname{vec}^{-1}\left[e^{S \otimes S} \operatorname{vec}(N)\right] & =\operatorname{vec}^{-1}\left[\sum_{i=0}^{\infty}(i!)^{-1}(S \otimes S)^{i} \operatorname{vec}(N)\right] \\
& =\operatorname{vec}^{-1}\left[\sum_{i=0}^{\infty}(i!)^{-1}\left(S^{i} \otimes S^{i}\right) \operatorname{vec}(N)\right] \\
& =\sum_{i=0}^{\infty} \operatorname{vec}^{-1}\left[(i!)^{-1}\left(S^{i} \otimes S^{i}\right) \operatorname{vec}(N)\right] \\
& =\sum_{i=0}^{\infty}(i!)^{-1} S^{i} N S^{i} \\
& \geq 0 . \tag{A.12}
\end{align*}
$$

Furthermore, from (A.4) and (A.10)

$$
\begin{align*}
\operatorname{vec}^{-1}\left[e^{S} \oplus S_{\operatorname{vec}(N)}\right] & =\operatorname{vec}^{-1}\left[\left(e^{S} \otimes e^{S}\right) \operatorname{vec}(N)\right] \\
& =e^{S} \mathrm{Ne}^{S} \\
& \geq 0 \tag{A.13}
\end{align*}
$$

Now, using the exponential product formula (see [10, pg. 97]) we write

$$
\begin{equation*}
e^{\bar{A}_{\alpha} t}=\lim _{j \rightarrow \infty}\left\{\exp \left[\frac{1}{j}\left(\tilde{A}_{\alpha} \oplus \tilde{A}_{\alpha}\right) t\right] \prod_{i=1}^{p} \exp \left[\frac{1}{\alpha j}\left(\tilde{A}_{i} \otimes \tilde{A}_{i}\right) t\right]\right\}^{j} \tag{A.14}
\end{equation*}
$$

where

$$
\tilde{A}_{\alpha} \triangleq\left[\begin{array}{cc}
A & B  \tag{A.15}\\
0 & 0
\end{array}\right]+\frac{\alpha}{2} I ; \quad \tilde{A}_{i} \triangleq\left[\begin{array}{cc}
A_{i} & B_{i} \\
0 & 0
\end{array}\right] .
$$

Consequently,

$$
\begin{align*}
& \operatorname{vec}^{-1}\left[e \bar{A}_{\alpha}^{t} \operatorname{vec}\left(\bar{Y}_{0}\right)\right] \\
&=\operatorname{vec}^{-1}\left\{\lim _{j \rightarrow \infty}\left\{\exp \left[\frac{1}{j}\left(\tilde{A}_{\alpha} \oplus \tilde{A}_{\alpha}\right) t\right] \prod_{i=1}^{p} \exp \left[\frac{1}{\alpha j}\left(\tilde{A}_{i} \otimes \tilde{A}_{i}\right) t\right]\right\}^{j} \operatorname{vec}\left(\bar{Y}_{0}\right)\right] \\
&=\lim _{j \rightarrow \infty}\left\{\operatorname{vec}^{-1}\left\{\left(\exp \left[\frac{1}{j}\left(\bar{A}_{\alpha} \oplus \bar{A}_{\alpha}\right) t\right] \prod_{i=1}^{p} \exp \left[\frac{1}{\alpha j}\left(\tilde{A}_{i} \otimes \bar{A}_{i}\right) t\right]\right\}^{j} \operatorname{vec}\left(\bar{Y}_{0}\right)\right]\right\} . \tag{A.16}
\end{align*}
$$

We'll now show that the expression in the limit brackets in (A.16) is positive semi-definite for all positive integers $j$. Indeed, this is sufficient to prove that (A.9) holds. For simplicity take $j=p=1$ and let $\tilde{N}$ satisfy

$$
\begin{equation*}
\operatorname{vec}(\tilde{N})=\exp \left[\frac{1}{\alpha}\left(\tilde{A}_{i} \otimes \tilde{A}_{i}\right) t\right] \operatorname{vec}\left(\bar{Y}_{0}\right) \tag{A.17}
\end{equation*}
$$

or equivalently

$$
\begin{equation*}
\tilde{N}=\operatorname{vec}^{-1}\left\{\exp \left[\frac{1}{\alpha}\left(\tilde{A}_{i} \otimes \tilde{A}_{i}\right) t\right] \operatorname{vec}\left(\bar{Y}_{0}\right)\right\} \tag{A.18}
\end{equation*}
$$

Since $\bar{Y}_{0}$ is assumed positive semi-definite, then, from (A.12), $\tilde{N}$ is positive semi-definite. From (A.13), (A.17) and this fact, the expression in the limit brackets of (A.16) satisfies

$$
\exp \left[\left(\tilde{A}_{\alpha} \oplus \tilde{A}_{\alpha}\right) t\right] \exp \left[\frac{1}{\alpha}\left(\tilde{A}_{i} \otimes \tilde{A}_{i}\right) t\right] \operatorname{vec}\left(\bar{Y}_{0}\right)=\exp \left[\left(\tilde{A}_{\alpha} \oplus \tilde{A}_{\alpha}\right) t\right] \operatorname{vec}(\tilde{N})
$$

$$
\geq 0
$$

A similar argument, using (A.12) and (A.13) alternately, shows that the expression in the limit brackets of (A.16) is positive semi-definite for arbitrary positive integers $j$ and $p$. This completes the proof of Observation 3.

Proof of Lemma 3.1: Assume $P$ positive definite and $\alpha>0$. Furthermore, consider the matrix differential equations in (A.5) and (A.7) with

$$
\begin{equation*}
Y_{0}=\bar{Y}_{0}=\Pi_{P} \tag{A.19}
\end{equation*}
$$

where $\Pi_{p}$ is given in (3.5). Subtracting (A.5) and (A.7) gives

$$
\begin{gather*}
\dot{\bar{Y}}(t)-\dot{Y}(t)=\left[\begin{array}{cc}
(A+\Delta A) & (B+\Delta B) \\
0 & 0
\end{array}\right](\bar{Y}(t)-Y(t))+(\bar{Y}(t)-Y(t))\left[\begin{array}{cc}
(A+\Delta A) & (B+\Delta B) \\
0 & 0
\end{array}\right]^{\prime}+\psi_{\alpha}(t) \\
t \geq 0, \quad \bar{Y}(0)-Y(0)=0 \tag{A.20}
\end{gather*}
$$

where

$$
\begin{align*}
\Psi_{\alpha}(t) \triangleq \tilde{A}_{\alpha} \bar{Y}(t) & +\bar{Y}(t) \tilde{A}_{\alpha}^{\prime}+\sum_{\alpha}^{1} \sum_{i=1}^{p} \tilde{A}_{i} \bar{Y}(t) \tilde{A}_{i}^{\prime} \\
& -\left(\left[\begin{array}{cc}
(A+\Delta A) & (B+\Delta B) \\
0 & 0
\end{array}\right] \bar{Y}(t)+\bar{Y}(t)\left[\begin{array}{cc}
(A+\Delta A) & (B+\Delta B) \\
0 & 0
\end{array}\right]^{\prime}\right) \tag{A.21}
\end{align*}
$$

and where $\tilde{A}_{\alpha}$ and $\tilde{A}_{i}$ are defined in (A.15).

Claim 1: $\quad \psi_{\alpha}(t)$ is positive semi-definite for all $t \geq 0$.
Proof of Claim 1: From (A.15) and (A.21)

$$
\begin{align*}
& \Psi_{\alpha}(t)=\left(\left[\begin{array}{ll}
A & B \\
0 & 0
\end{array}\right]+\frac{\alpha}{2} I\right) \bar{Y}(t)+\bar{Y}(t)\left(\left[\begin{array}{ll}
A & B \\
0 & 0
\end{array}\right]^{\prime}+\frac{\alpha}{2} I\right) \\
& -\left(\left[\begin{array}{cc}
(A+\Delta A) & (B+\Delta B) \\
0 & 0
\end{array}\right] \bar{Y}(t)+\bar{Y}(t)\left[\begin{array}{cc}
(A+\Delta A) & (B+\Delta B) \\
0 & 0
\end{array}\right]^{\prime}\right)+\frac{1}{\alpha} \sum_{i=1}^{p} \tilde{A}_{i} \bar{Y}^{(t)} \tilde{A}_{i}^{\prime} \\
& =\alpha \bar{Y}(t)-\left(\left[\begin{array}{cc}
(A+\Delta A) & (B+\Delta B) \\
0 & 0
\end{array}\right] \bar{Y}(t)+\bar{Y}(t)\left[\begin{array}{cc}
(A+\Delta A) & (B+\Delta B) \\
0 & 0
\end{array}\right]^{\prime}\right)+\frac{1}{\alpha} \sum_{i=1}^{P} \tilde{A}_{i} \bar{Y}(t) \tilde{A}_{i}^{\prime} \\
& \geq \alpha \sum_{i=1}^{p} \sigma_{i}^{2} \bar{Y}(t)-\sum_{i=1}^{p} \sigma_{i}\left(\tilde{A}_{i} \bar{Y}(t)+\bar{Y}(t) \tilde{A}_{i}^{\prime}\right)+\frac{1}{\alpha} \sum_{i=1}^{p} \tilde{A}_{i} \bar{Y}(t) \tilde{A}_{i}^{\prime} \\
& =\sum_{i=1}^{p}\left[\sqrt{\alpha} \sigma_{i} I-\frac{1}{\sqrt{\alpha}} \tilde{A}_{i}\right] \bar{Y}(t)\left[\sqrt{\alpha} \sigma_{i} I-\frac{1}{\sqrt{\alpha}} \tilde{A}_{i}\right]^{\prime} . \tag{A.22}
\end{align*}
$$

From Observation $3, \bar{Y}(t)$ is positive semi-definite for all $t \geq 0$. It thus follows from (A.22) that $\Psi_{\alpha}(t)$ is likewise positive semi-definite for all $t \geq 0$. This proves the claim.

Now, since $\bar{Y}(0)-Y(0)=0$, the solution to (A.20) is

$$
\left.\left.\bar{Y}(t)-Y(t)=\int_{0}^{t} \exp \left[\begin{array}{cc}
(A+\Delta A) & (B+\Delta B)  \tag{A.23}\\
0 & 0
\end{array}\right](t-s)\right] \psi_{\alpha}(s) \exp \left[\begin{array}{cc}
(A+\Delta A) & (B+\Delta B) \\
0 & 0
\end{array}\right]^{\prime}(t-s)\right] d s
$$

From (A.23) and Claim 1 it follows that

$$
\begin{equation*}
\bar{Y}(t)-Y(t) \geq 0, \quad t \geq 0 \tag{A.24}
\end{equation*}
$$

Combining (A.6) and (A.8) with (A.24) gives

$$
\begin{aligned}
0 & \leq \bar{Y}(h)-Y(h) \\
& =\operatorname{vec}^{-1}\left[e^{\bar{A}_{\alpha} h} \operatorname{vec}\left(\Pi_{p}\right)\right]-\exp \left(\left[\begin{array}{cc}
(A+\Delta A) & (B+\Delta B) \\
0 & 0
\end{array}\right] h\right) \Pi_{P} \exp \left(\left[\begin{array}{cc}
(A+\Delta A) & (B+\Delta B) \\
0 & 0
\end{array}\right]^{\prime} h\right)
\end{aligned}
$$

which implies that
$\left[\begin{array}{ll}I & 0\end{array}\right] \exp \left(\left[\begin{array}{cc}(A+\Delta A) & (B+\Delta B) \\ 0 & 0\end{array}\right] h\right) \quad \Pi_{p} \exp \left(\left[\begin{array}{cc}(A+\Delta A) & (B+\Delta B) \\ 0 & 0\end{array}\right]^{\prime} h\right)\left[\begin{array}{l}I \\ 0\end{array}\right]-P$

$$
\leq\left[\begin{array}{ll}
I & 0
\end{array}\right] \operatorname{vec}^{-1}\left[e^{\bar{A}_{\alpha} h} \operatorname{vec}\left(\Pi_{P}\right)\right]\left[\begin{array}{l}
I \\
0
\end{array}\right]-P
$$

which is the desired result. The proof of Lemma 3.1 is complete.
$\nabla \nabla \nabla$

Proof of Lemma 3.2: Assume $\alpha>0$ such that (3.10) holds. We must show that (3.8) has a positive-definite solution $P$. Indeed, since $\rho\left(M_{\alpha}\right)<1$, then $\left(M_{\alpha}-I\right)$ is invertible and (3.8) has a unique solution

$$
\operatorname{vec}(P)=\left(I-M_{\alpha}\right)^{-1} \operatorname{vec}(I)
$$

or

$$
P=\operatorname{vec}^{-1}\left[\left(I-M_{\alpha}\right)^{-1} \operatorname{vec}(I)\right] .
$$

Now, with $M_{\alpha}^{0} \triangleq I$, it follows from [11, Theorem 6.7.1] that $\left(I-M_{\alpha}\right)^{-1}=$ $\sum_{i=0}^{\infty} M_{\alpha}^{i}$. Consequently,

$$
\begin{equation*}
P=\operatorname{vec}^{-1}\left[\sum_{i=0}^{\infty} M_{\alpha}^{i} \operatorname{vec}(I)\right]=\sum_{i=0}^{\infty} \operatorname{vec}^{-1}\left[M_{\alpha}^{i} \operatorname{vec}(I)\right] . \tag{A.25}
\end{equation*}
$$

Claim 2: If i is a positive integer, then

$$
\begin{equation*}
\operatorname{vec}^{-1}\left[M_{\alpha}^{i} \operatorname{vec}(I)\right] \geq 0 \tag{A.26}
\end{equation*}
$$

Proof of Claim 2: The proof proceeds by induction. From the definition of $M_{\alpha}$ in (3.9) and identity (A.2), we have for $i=1$

$$
\begin{align*}
\operatorname{vec}^{-1}\left[M_{\alpha} \operatorname{vec}(I)\right] & =\operatorname{vec}^{-1}\left\{\left[\begin{array}{ll}
I & 0
\end{array}\right] e^{\bar{A}_{\alpha} h}\left(\left[\begin{array}{l}
I \\
K C
\end{array}\right] \otimes\left[\begin{array}{l}
I \\
K C
\end{array}\right]\right) \operatorname{vec}(I)\right\} \\
& =\operatorname{vec}^{-1}\left\{[ \begin{array} { l l } 
{ I } & { 0 }
\end{array} ] e ^ { \overline { A } _ { \alpha } h } \operatorname { v e c } \left([ \begin{array} { l } 
{ I } \\
{ K C }
\end{array} ] \left[\begin{array}{ll}
I & \left.\left.\left.C^{\prime} K^{\prime}\right]\right)\right\}
\end{array}\right.\right.\right. \\
& =\left[\begin{array}{ll}
I & 0
\end{array}\right] \operatorname{vec}^{-1}\left\{e^{\bar{A}_{\alpha} h} \operatorname{vec}\left(\left[\begin{array}{l}
I \\
K C
\end{array}\right]\left[\begin{array}{ll}
I & \left.\left.\left.C^{\prime} K^{\prime}\right]\right)\right\}
\end{array}\right] \begin{array}{l}
I \\
0
\end{array}\right] .\right. \tag{A.27}
\end{align*}
$$

It follows from (A.9), (A.27) and Observation 3 with

$$
\bar{Y}(0)=\left[\begin{array}{l}
I \\
K C
\end{array}\right]\left[\begin{array}{ll}
I & C^{\prime} K^{\prime}
\end{array}\right], \quad t=h
$$

that $\operatorname{vec}^{-1}\left[M_{\alpha} \operatorname{vec}(I)\right] \geq 0$. Now, for the induction step, assume $\operatorname{vec}^{-1}\left[M_{\alpha}^{i} \operatorname{vec}(I)\right] \geq 0$. Then

$$
\begin{align*}
\operatorname{vec}^{-1}\left[M_{\alpha}^{i+1} \operatorname{vec}(I)\right] & =\operatorname{vec}^{-1}\left[M_{\alpha} M_{\alpha}^{i} \operatorname{vec}(i)\right] \\
& =\operatorname{vec}^{-1}\left\{M_{\alpha} \operatorname{vec}\left[\operatorname{vec}^{-1}\left(M_{\alpha}^{i} \operatorname{vec}(I)\right)\right]\right] . \tag{A.28}
\end{align*}
$$

From (A.28) and Observation 3 with

$$
\bar{Y}(0)=\operatorname{vec}^{-1}\left[M_{\alpha}^{i} \operatorname{vec}(I)\right], \quad t=h,
$$

it follows that $\operatorname{vec}^{-1}\left[M_{\alpha}^{i+1} \operatorname{vec}(I)\right] \geq 0$. This proves the lemma.

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Figure 1 : Sampled-Data Control System


Figure 2 : Plot of $\rho\left(M_{\alpha}\right)$ versus $\alpha$ for $h=0.1$, $r=0.05, \quad \bar{\sigma}=0.5$



[^0]:    Eigure 25. The Matrix Majorant Is a Bound for the Matrix Block Norm. i.e.. the Nonnegative Natrix Each of Whose Elamente Is che Narm of the Corresponding Block of a Given Matrix

[^1]:    Manuscripe received October 2, 1984; revised February 1. 1985. This paper is based on a prior submission of March 14, 1984. Paper recommended by Past Associate Editor, B. R. Barmish. This work was supported in part by the Air Force at Lincoln Laboratory/M.I.T., Lexington, MA.

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    ' The quadratic efror criterion has been chosen for consideration in the present paper because of its relation to the standard engineering practice of stating specifications in terms of ims deviation.

[^2]:    ${ }^{2}$ The projection was. however, pointed out in [28, p. 29$]$.

[^3]:    ' J will occasionally be referred to as the "model-reduction error" or. simply, as the "cost."

[^4]:    Manuscript received August 15. 1983: revind February 14. 1983. This paper is based on a prior submission of March 10. 1983. Paper recommended hy P. R Kumar. Past Chairman of the Stochastic Control Committec This work was suppuned ty the Deparment of the Air Force and was perfurmed at Lincoln Lethratory. A1 IT
    The authors are with the Controls Analy mes and Sy nthesis Grimp. Harris Corp. GASD, Melbourne. FL 12902.

[^5]:    * Received by the editors December 6, 1983, and in revised form September 15, 1984. This work was performed at Lincoin Laboratory/MIT and was sponsored by the Department of the Air Force.
    † Harris Corporation, Government Aerospace Systems Division, Controls Analysis and Synthesis Group, Melbourne, Florida 32901.
    ' Examples of such components include tapped delay lines and surface acoustic wave devices. Although acoustoelectric convolvers [3, p. 465] can perform continuous-time integration, synthesis of the desired impulse-response kernel can incorporate only finitely many specified parameters. The obvious fact should also be noted that physical limitations impose an upper bound on the number of design parameters that can be incorporated in the construction of any device. For an extensive treatment of this subject, see [72].

[^6]:    ${ }^{2}$ Alternatively, we could have adopted the white noise formulation of [4]. The main difierence between the two white noise formalisms is that Balakrishnan works with finitely additive rather than countably additive measures. Strictly speaking, then, even in finite dimensions Balakrishnan's white noise is different from the standard notion (see [6, pp. 307, 315]).

[^7]:    ${ }^{3}$ This assumption and its analogue, the lack of a cross-weighting term $x(t)^{\top} R_{12} u(1)$ in (3.5), are for convenience only. See $\$ 5$.
    ${ }^{4}$ We must require that $R_{1}$ and $V_{1}$ be nuclear since covariance operators in the white noise formulation of [6] are not necessarily trace class as they are in the formulation of [4].

[^8]:    ${ }^{3}(3.14 \mathrm{a})$ refers to $\rho(\hat{Q})=n_{c}$ etc.

[^9]:    * This research supported in part by the Air Force Office of Scientific Research under contract F49620-86-C-0002
    ** This research supported in part by the Air Force Office of Scientific Research under grant AFOSR-87-0351. Part of this research was carried out while this author was a visiting scientist at the Institute for Computer Applications in Science and Engineering (ICASE) at the NASA Langley Research Center in Hampton, VA which is operated under NASA contracts NAS1-17070 and NAS1-18107.

[^10]:    76. 

    D. S. Bernstein, "Sequential Design of Decentralized Dymamic Compensators Using the Optimal Projection Equations," Int. J. Contr., Vol. 46, pp. 1569-1577, 1987.
    121. D. S. Bernstein and W. M. Haddad, "Kobust Decentralized Static Output Feedback," submitted to Sys, Contr. Lett.

[^11]:    Received 15 December 1986.
    $\dagger$ Harris Corporation, Government Aerospace Systems Division, P.O. Box 94000, Melbourne, Florida 32902, U.S.A.

[^12]:    * The Kuhn-Tucker theorem requires a priori verification of a constraint qualification which is difficult to confirm in the present context. The Fritz John version is less restrictive and hence more suitable.

[^13]:    1 As shown to [29], the formula for the derivative of a meslar function whith roapect so symancric ergumens $Q$ and $\beta$ ecails a modificasion of (5.6) and (5.7). Since thene framients ate being set to zero, however, the fleal rewlt is idexaical. Ahernatively. $Q$ and $\beta$ cen be viewod (as we are doing bere) as astiorary matrix variables. Symmetry is maposed only a poorteriori by the form of (5.4) and (5.5) and the sebility of $\lambda$. Honce. memeamicully, the reauh of (29) is bot required.

[^14]:    Manuscript received February 6, 1986; revised November 25. 1986 and May 5. 1997. This work was aupported in part by the Air Force Office of Sciemetfic Research under Contrmists AFOSR F49620-84-C-0015 and F49620-86-C-0002.
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    IEEE LOE Number $\$ 716159$.

[^15]:    $\overline{1}_{\text {This research was supported }}$ in part by the Air Force Office of Scientific Research under contracts F49620-84-C-0015, F49620-86-C-0002 and P49620-86-C-0038.
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[^16]:    Manuscript received July 23. 1985: revised October 24. 1985. This work wes supported in part by the Air Force Office of Scientific Reseanch under Contrict AFOSR F49620-84-C.0015.
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    IEEE LOE Number 8407268.

[^17]:    Manuscript received April 23, 1986: revised April 13, 1987. Paper recommended by Associatc Editor, D. P. Looze. This work was supported in part by the Air Force Office of Scientific Research under Contracts F49620. 86-C-0002 and F49620-86-C-0038.
    The authur is with the Government Aerospace Systems Division, Melbourne. FL 32902.

    IEEE Log Number 8716818
    Several exceptions to these remarks should be noted: necessary and wufficient conditions have been given in [12]. [21]: a quadratic cost functional is utilized in [10]. [13]: and controller effort is considered in [20].

[^18]:    Manuscript received August 8. 1986; revised May 6, 1987. Paper recommended by Associate Editor. M. G. Safonov. This work was supporied in pari by the Air Force Office of Scientific Research under Contracts F49620-86-C-0002 and F49620-86-C-0038.
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    IEEE Log Number 8716541.

[^19]:    ' Uncertainties in a single subsystem can also be regarded as interaction uncerainties. To see this, write $\dot{x}=(A+G) x$ twice so that the uncernainty $G$ is represented by $\left[\begin{array}{cc}0 \\ c & 0 \\ 0\end{array}\right]$.

[^20]:    ${ }^{2}$ Upon first reading the uncertainty represented by (3.3) can be ignored since the principal contribution concerns the treatment of (3.2).

[^21]:    Manuscript received June 8. 1987: revined September 25. 1987. This work was supported in part by the U.S. Air Force Office of Scientific Research under Contract F49620-86-C-0002.

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[^23]:    Manukipa received Augusp 4. 1987; revised Ocwber 16. 1987. This' work wa mpponed in part by the Air Force Onfice of Scientific Revearch ander Comerser F49620-86-C0002
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[^24]:    Manuscripm received Fehruary 13. 1987. revised Seprember 28. 1987. Paper recommended by Associate Editor, B. H. Krogh.
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    IEEE Log Number 8821013.

[^25]:    * Supported in part by the Air Force Office of Scientific Research under contract F49620-86-C-0002.

[^26]:    This research was supported in part by the Air Force Office of Scientific Research under contract F49620-86-C-0002.

[^27]:    * It is also interesting to note that in Hamilton-Jacobi-Bellman sufficiency theory the performance functional is expressed in terms of a value function which also serves as a Lyapunov function for the closed-loop system. These connections will be explored in a future paper.

[^28]:    * Asymptotic stability for a nonautonomous system is defined in the standard way. See, e.g., [42].

[^29]:    Supported in part by the Air Force Office of Scientific Research under contract F49620-86-C0002

[^30]:    This work is sponsored by the Dept. of the Air Force. The views expressed are those of the author and do not reflect the official policy or position of the U.S. Government.
    ** Supported in part by the Air Force Office of Scientific Research under contract F49620-86-C-0002.

[^31]:    Manuscript received July 17, 1986: revised March 4, 1987. This work was supported in part by the Air Force Office of Scientific Research under Contracts F49620-86-C-0002 and P49620-86-C- -038
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[^33]:    Received April 1, 1985; revision received Aug. 26, 1985. Copyrighe C American Institute of Aeronautics and Astronautics, Inc., 1986. All rights reserved.
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[^34]:    * Supported in part by the Air Force Office of Scientific Research under contract F49620-86-C-0002.

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[^36]:    - Received by the editors January 21. 1487, accepted for publication May 11. 198? This research was supported in part by the Aur Force Ottice of Scientutic Research under contract F49620-86-C-0002.
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