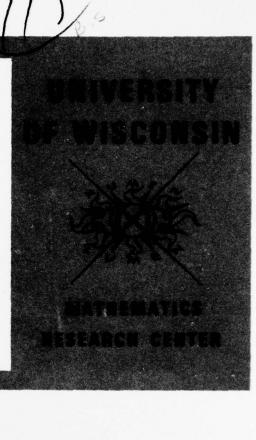


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PERIODIC AND QUASIPERIODIC SOLUTIONS OF  $\Delta u + \lambda u + O(u) = 0$ 

Klaus Kirchgässner and Jürgen Scheurle



Mathematics Research Center University of Wisconsin—Madison 610 Walnut Street Madison, Wisconsin 53706

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# UNIVERSITY OF WISCONSIN - MADISON MATHEMATICS RESEARCH CENTER

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#### ABSTRACT

In this paper the boundary value problem  $\Delta u + \lambda u + f(u, u_x, u_y) = 0$ , u(0,y) = u(1,y) = 0 is studied in the strip  $(0,1) \times R$ , where f is some  $C^2$ -function which, together with its gradient, vanishes at 0,  $\lambda$  is a real parameter. It is shown that, for  $\lambda$  between  $\pi^2$  and  $4\pi^2$ , all small solutions are periodic in y. Moreover, singular solutions exist as local  $H^2$ -limits of periodic solutions with large periods. For values of  $\lambda$  beyond  $4\pi^2$  a formal argument suggests that almost all small solutions are quasiperiodic. The equation is studied as a model for some important but technically cumbersome bifurcation problems in fluid dynamics.

#### EXPLANATION

A nonlinear elliptic boundary-value problem which models aspects of the technically more cumbersome Taylor and Bénard problems of fluid dynamics is considered. It is shown that all small solutions are periodic in one variable when a parameter (corresponding to a bifurcation parameter) is in a suitable range and indications are given that all small solutions are quasiperiodic for larger parameter values when the nonlinearity is analytic.

AMS(MOS) Subject Classification - 35J60, 76DXX

Key Words - Nonlinear elliptic boundary-value problems, periodic solution, quasiperiodic solution, Taylor and Bénard problems

Work Unit Number 1 - Applied Analysis

In this study we try to classify all small solutions of the boundary-value problem

$$\Delta u + \lambda u + \tilde{f}(u, D_{\chi}u, D_{y}u) = 0$$

$$(1.1)$$

$$u(0,y) = u(1,y) = 0$$

in the strip  $\Omega = \{(x,y) \in [0,1] \times \mathbb{R}\}$ . Here,  $\Delta$  is the two-dimensional Laplacean,  $\lambda$  a real parameter and  $\tilde{f}(u,p,q)$  a real valued  $C^2$ -function of its real arguments which, together with its gradient, vanishes for u=p=q=0. The class of solutions considered consists of functions with locally uniform  $H^S$ -norm

$$\sup_{\ell \in \mathbb{Z}} \|u\|_{H^{S}(K_{\ell})} < \varepsilon$$

where  $0 < s \le 2$ ,  $H^s$  the usual Sobolev space, and where  $K_{\ell} = [0,1] \times [(\ell-1)_{\ell},\ell_{\ell}]$  for some  $_{\ell} > 0$  denotes a sequence of compacta covering  $_{\Omega}$ . In view of the physical interpretation given later, our assumption requires small energy input per unit length.

The interest in this question arose from the effort to determine all physically reasonable solutions of the Taylor- and the Bénard problem in fluid dynamics. In both problems a basic (trivial) solution looses its stability to nontrivial solutions which bifurcate at a critical parameter value. In view of the underlying invariance, the set of nontrivial solutions near this bifurcation point is very large [11] but nothing is known about a characterisation of this set. There is an extensive theoretic literature (mathematical or physical) on these problems (see [4] for an excellent survey), however, all approaches assume periodicity in the unbounded variables a priori.

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One might consider this investigation as a first step toward a final answer of those questions, since (1.1) can, to some extent, be regarded as a model for the Navier-Stokes equations - the stationary Burgers equation is a special case of it - .

The conjecture, derived from the linearized equation, that - besides a discrete set of  $\lambda$ -values - all solutions are quasiperiodic is not true without further symmetry restrictions on  $\tilde{f}$  as the one dimensional analogon  $\tilde{u} + \lambda u + \tilde{u}^3 = 0$  already shows. Hence we further impose one of the two following assumptions

(a) 
$$\tilde{f}(u,p,-q) = \tilde{f}(u,p,q)$$
  
(1.2)  
(b)  $\tilde{f}(u,p,-q) = -\tilde{f}(u,p,q)$  and  $\tilde{f}(-u,-p,-q) = \tilde{f}(u,p,q)$ 

Given property (1.2) we are able to settle the question in the  $\lambda$ -interval ( $\pi^2$ ,  $4\pi^2$ ) by proving that all small solutions of (1.1) are periodic in y. The question becomes more delicate for  $\lambda \in (n^2\pi^2, (n+1)^2\pi^2)$  where quasiperiodic solutions are expected. The method applied for the periodic case breaks down since the invertible part of the linearized operator ceases to be continuously invertible. If the nonlinearity  $\tilde{f}$  is real analytic then we can show at least, that formally (considering formal power series) all small solutions are quasiperiodic.

Unfortunately, the analogy to the fluid dynamical situations is rather limited for  $\lambda \in (\pi^2, 4\pi^2)$  since, for the Tayloras well as for the Bénard-problem, one expects two independent frequencies  $\omega_1, \omega_2$ , for  $\lambda$  slightly above the critical value. Hence, the analogous model case is  $\lambda \in (4\pi^2, 9\pi^2)$  and we conjecture therefore that all small solutions in those fluid-dynamical problems are quasiperiodic. For indications of this fact consult [1],[5],[6],[12],

The exceptional set of values  $\lambda = n^2\pi^2$  requires special consideration. We are able to prove for  $\lambda = \pi^2$  the following alternative: either, there exist arbitrary small, nontrivial solutions for  $\lambda = \pi^2$  or, there exist "singular" solutions near  $\lambda = \pi^2$ ,  $\mu = 0$ , whose first Fourier-component is either constant or nonperiodic. The result is natural as analogous ordinary differential equations show. Generalisations to problems where  $\lambda$  is replaced by a strongly elliptic operator of order 2m are possible, but we avoid them here for technical reasons.

To our knowledge, the problem under consideration has not been studied in the literature. There has been some interest in periodic solutions of nonlinear wave equations (see [8] for a recent list of references) but there, the interest is directed towards the existence of special solutions and their regularity.

We start with an investigation of  $\Delta + \lambda$  in a space which is the inductive limit of weighted Sobolev spaces allowing polynomial growth of any order. We show that  $\Delta + \lambda$  is onto for every  $\lambda$ , a fact, which justifies the choice of space and which contributes the corner stone for the reduction of (1.1) to a finite dimensional initial value problem – somewhat in the spirit of Liapunov and Schmidt – (sections 2 and 3). In section 4 we prove uniqueness in Theorem 4.4 . Section 5 contains the existence of just enough periodic resp. quasiperiodic solutions. Section 6 is devoted to the existence of singular solutions as envelopes of periodic ones.

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Let be  $\Omega:=(0,1)\times\mathbb{R}$ ,  $\underline{x}=(x,y)\in\mathbb{R}^2$ , and K a nonempty compact subset of  $\Omega$ . By  $H^S_{10c}(\Omega)$  we denote the set of all funktions  $u:\Omega\to\mathbb{R}$ , with  $u\in H^S(K)$  for all  $K\subset\Omega$ . Here,  $H^S(K)$  is the usual real Sobolev space with the norm  $(\alpha=(\alpha_1,\alpha_2)\in\mathbb{N}_0^2, |\alpha|=\alpha_1+\alpha_2)$ :

$$\|u\|_{H^{s}(K)} = \sup_{\|\alpha\| \le s} (\int_{K} |D^{\alpha}u|^{2} dx)^{1/2}, s \ge 0$$

 $C^0(\overline{\Omega})$  is the Banach space of continuous, real valued functions defined and bounded in  $\overline{\Omega}$  equipped with the norm

$$\|u\|_{C^{0}} = \sup_{\mathbf{x} \in \overline{\Omega}} |u(\underline{\mathbf{x}})| .$$

Furthermore, we need a sequence of Banach spaces  $\mathbf{X}_k^{S}$  which are defined as follows :

$$g_{k}(y) = \max (1, |y|^{k}), k \in \mathbb{N}_{0}$$

$$\lim_{s \to k} \sup_{|\alpha| \le s} (\int_{\Omega} g_{k}^{-2}(y) |19^{\alpha}u(x, y)|^{2} dx)^{1/2}$$

$$x_{k}^{s} = (u \in H_{10c}^{s}(2) | \|u\|_{s, k} < \infty), s \in \mathbb{N}_{0}$$

Observe that  $g_{k+\ell} = g_k g_{\ell}$  holds.

The spaces  $X_k^s$  consist of functions with a given polynomial growth at infinity. Since  $X_k^s \in X_{k+1}^s$  algebraically and topologically, they form a scale of Banach spaces. The boundary conditions u(0,y)=u(1,y)=0 define, via the trace, a subspace of  $X_k^s$  for s > 1/2 and every k, which is denoted by

$$\hat{X}_{k}^{S} := \{ u \in X_{k}^{S} \mid u(0,v) = u(1,v) = 0 \}.$$

Finally let  $X^S$  be the inductive limit of the Bancos spaces  $X_k^S$ ,  $k \in \mathbb{N}_0$ , i.e.  $X^S = \bigcup_{k=0}^{\infty} X_k^S$  is the locally convex topological (lct) vectorspace—equipped with the finest topology which induces on every  $X_k^S$  a topology coarser than the given one ({2},p.429). Hence,  $X^S$  contains exactly those

functions in  $H^S_{loc}(\Omega)$  which , for targe ly!, grow at most like a polynomial, and its lc-topology is such that the injections  $X_k^S \to X^S$  are continuous. The spaces are defined analogously as inductive limits of  $X_k^S$ .

A subset M of a vector space is called absolutely convex (ac) if, for every u,v  $\in$  M,  $\lambda u + \mu u \in$  M holds, whenever  $\|\lambda\| + \|\mu\| \leq 1$ ,  $\lambda, \mu \in \mathbb{R}$ . The absolutely convex hull (ach) of M consists of all finite linear combinations

 $\begin{array}{l} \sum\limits_{j}u_{j}u_{j}\,,\,\sum\limits_{j}|u_{j}|\leq1\,,\,u_{j}\in\mathbb{N}\\ \text{Thus, ach}(\mathbb{M})\text{ is the smallest ac set containing }\mathbb{M}.\text{ The}\\ \text{family of those ac subsets }\mathbb{U}\subset\mathbb{X}^{S},\text{ for which }\mathbb{U}\text{ n }\mathbb{X}_{k}^{S}\text{ is a}\\ \text{neighborhood in }\mathbb{X}_{k}^{S}\text{ for every }k,\text{ forms a basis of neighborhoods in }\mathbb{X}^{S}.\text{ In particular, the sets} \end{array}$ 

$$ach \left( \begin{array}{cc} \omega & \delta_{ck}^{k} \\ k=0 \end{array} \right)$$

form a basis of neighborhoods of 0 in X<sup>S</sup> and hence, generate its topology (cf. [ 0],p. 79). Here,

$$B_{\epsilon_k}^k := \{ u \in X_k^s \mid ||u||_{s,k} < \epsilon_k \}, \epsilon_k > 0$$
.

The set II is called bounded if, for every neighborhood U of 0, there exists a  $\mu = \mu(U)$ , such that  $M \subset \nu U$  for all  $\|\nu\| \ge \mu$ . If X and Y denote lct-vectorspaces and L: X  $\rightarrow$  Y is linear and continuous then, L maps bounded sets into bounded sets; such a map is called bounded. A lct-vector-space X is bornological if and only if every bounded linear map L acting from X into any other lct-vectorspace is continuous ([2],p.177).

As inductive limits of Banach spaces, the  $x^s$  are bornelogical, and even barreled ([9],p.31f). A linear, continuous surjective map  $L: x^s \to x^t$  is a homomorphism, i.e. maps open sets onto open sets [7]. It is not too hard to prove that  $x^s$  is seperated and sequentially complete. Since these properties are not used subsequently we omit the proof.

## LEMMA 2.1

The set  $M \subset X^S$  is bounded if and only if M is contained in some  $X_k^S$  and bounded there.

Proof: If  $U\subseteq X^S$  is a neighborhood of 0 in  $X^S$ , so is  $U\cap X_k^S$  a neighborhood of 0 in  $X_k^S$ . Hence, for  $M\subseteq X_k^S$  bounded, there exists a  $\nu>0$  such that  $M\subseteq \nu(U\cap X_k^S)\subseteq \nu U$  for all  $\|\nu\|\geq \nu$  Thus M is bounded in  $X^S$ .

Now, let M  $\subset$  X be bounded and assume that M is not contained in some  $X_k^S$  or, if so, not bounded there. Then , here is a sequence  $u_k \in M$ ,  $u_k \notin B_{k+1}^k$  for all  $k \in M$ . Hence, have  $\alpha_k \in M_0^2$  and compacta  $K_k \subset \Omega$  such that

(2.1) 
$$\int_{K_{k}} g_{k}^{-2} |D^{\alpha_{k}} u_{k}|^{2} dx > k^{2}$$

Choose a sequence  $(\epsilon_{\ell})$  with  $\beta_{\ell}$   $\epsilon_{\ell}$  < 1, and  $\epsilon_{\ell} \leq \min$  {  $g_{\ell-k}^{-1}(y)$  |  $0 \leq k \leq \ell$ ,  $(x,y) \in \Omega_{k}$  ]. We show that for  $u \in U = ach$  (  $U = B_{\epsilon_{\ell}}^{\ell}$ ) and all k

(2.2) 
$$\int_{K_k} g_k^{-2} |D^{\alpha_k} u|^2 dx < 1$$

holds. Hence,  $u_k$  ¢ kU by (2.1) which yields the contradiction.

To prove (2.2) observe, that for  $0 \le z \le k$ ,  $u \in B_{\epsilon}^{\hat{k}}$ ,

$$\int_{K_k} g_k^{-2} |D^{\alpha_k} u|^2 d\underline{x} \leq \int_{\Omega} g_k^{-2} |D^{\alpha_k} u|^2 d\underline{x} < \varepsilon_k^2 < 1$$

holds and, if a > k,  $u \in B_{\epsilon}^{\ell}$ 

$$\int_{K_{k}} g_{k}^{-2} |D^{\alpha_{k}}u|^{2} d\underline{x} \leq \frac{1}{\min_{\substack{m \text{in } g_{2}^{-2} \\ \underline{x} \in \Omega_{k}}}} \int_{K_{k}} g_{\underline{x}}^{-2} |D^{\alpha_{k}}u|^{2} d\underline{x}$$

$$\leq \frac{2}{k}$$

$$= \frac{2}{k}$$

$$= 1$$

Since U is the ach of  $B_{\varepsilon_{\mathcal{Q}}}^{\varrho}$ , (2.2) follows readily.

## LEMMA 2.2

Let the sequence  $(u_{\ell}) \subset X^S$  be bounded and such that  $(u_{\ell \mid K})$  converges in  $H^S(K)$  for every compact  $K \subset \Omega$ . Then,  $(u_{\ell})$  converges in  $X^S$ .

Proof: According to Lemma 2.1,  $(u_2) \subseteq X_k^s$  holds for some k, and  $(u_2)^l$  bounded there:  $\|u_2\|_{S,k} \le C$ . We show, there exists an  $u \in X_k^s$  such that  $\lim u_2 = u$  in  $X_{k+1}^s$ . Since  $X_{k+1}^s$  lies continuously in  $X^s$ , the assertion follows. Construct an  $u \in H_{loc}^s(\Omega)$  such that  $u_{k|K}$  converges towards  $u_{lK}$ . Hence, for every  $K \subset \Omega$  we have an index  $\ell$  with

$$\sup_{|\alpha| \le s} \left( \int_{K} \sigma_{k}^{-2} |D^{\alpha}(u \cdot u_{\ell})|^{2} d\underline{x} \right)^{1/2} \le 1$$

and thus

$$\sup_{\|\alpha\| \le s} \left( \int_{K} g_{k}^{-2} \|D^{\alpha}u\|^{2} dx \right)^{1/2} \le \|u_{k}\|_{S,k} + 1 \le c + 1 .$$

This yields  $u \in x_k^s$ . For arbitrary  $\delta > 0$  choose  $K \subset \Omega$  so large that  $\lg_1^{-1}(y)! < \delta(4C*2)^{-1}$  for all  $(x,y) \in \mathbb{T} = K$ . Hence

$$\|u - u_{\tilde{x}}\|_{5, k+1} \leq \sup_{\|\alpha\| \leq s} \left( \int_{K} g_{k+1}^{-2} \|D^{\alpha}(u - u_{\tilde{x}})\|^{2} d\underline{x} \right)^{1/2} + \sup_{\|\alpha\| \leq s} \left( \int_{C \setminus K} g_{k+1}^{-2} \|D^{\alpha}(u - u_{\tilde{x}})\|^{2} d\underline{x} \right)^{1/2} \\ \leq \frac{\delta}{2} + \frac{\delta}{2(2C+1)} \|u - u_{\tilde{x}}\|_{5, k} \leq \delta$$

for all sufficiently large &; which proves the assertion.

Subsequently, we need a characterisation of  $X_k^s$  by Fourier components which is an asy consequence of Parseval's equality. We introduce for  $u \in \mathbb{R}^d \setminus \mathbb{R}$  the Banach spaces  $Y_k^s$ :

(2.3) 
$$|u|_{s,k} = \sup_{0 \le \gamma \le s} \int_{\mathbb{R}} g_k^{-2} |D^{\gamma}u|^2 dy ,$$

$$|v|_{s,k} = \{u \in H_{1oc}^s(\mathbb{R}) \mid |u|_{s,k} < \infty \}$$

# LEMMA 2.3

Set  $X_k^0 = X_k^0$  and assume  $u \in X_k^s$  for s = 0,1,or 2. Define  $u_v = \sqrt{2} \int_0^1 u(x,\cdot) \sin v \pi x \, dx$ ,

then we have  $u_{v} \in Y_{k}^{S}$  for all  $v \in \mathbb{N}$  and

(2.4) 
$$\sum_{i=1}^{8} (v\pi)^{2\beta} |D^{Y}u_{i}|^{2} = ||D^{\beta}_{X}D^{Y}_{Y}u_{i}|^{2}$$

for β, γ ε No , β + γ ≤ s.

Conversely, let  $(u_v \mid u_v \in Y_k^s)$  be a sequence for which the left side of (2.4) is finite. Then, the functions  $u_v$  are the Fourier components of some u in  $X_k^s$  which satisfies (2.4).

Proof: 
$$u_y \in Y_k^s$$
 follows from
$$D^Y u_y = \sqrt{2} \int_0^1 D_y^Y u(x, \cdot) \sin v \pi x \, dx$$

and hence

Define

$$φ_β^ν(x) = \sqrt{2}$$

$$\begin{cases}
sin νπx & for β = 0 \\
cos νπx & for β = 1 \\
-sin νπx & for β = 2
\end{cases}$$

and

$$u_y^{\beta,\gamma} = \int_0^1 D_x^{\beta} D_y^{\gamma} u(x,\cdot) \varphi_{\beta}^{\nu}(x) dx$$

Then
$$u_{v}^{\beta, \gamma} = (v\pi)^{\beta} D^{\gamma} u_{v}$$

holds, and Parseval's idendity implies (2.4).

For the proof of the converse set  $u_N^{\beta,\gamma} = \sum_{k=0}^{N} (\nu\pi)^{\beta} D^{\gamma} u_{\nu} \phi_{\beta}^{\nu}$  then  $u_N^{\beta,\gamma}$  converges in  $X_k^0$  towards  $u^{\beta,\gamma}$ , which is the weak  $D_X^{\beta} D_Y^{\gamma}$  derivative of  $u = u^{0,0}$ . Moreover the traces  $u(0,\cdot)$ ,  $u(1,\cdot)$  vanish for s=1,2.

, 4

## § 3 Some Estimates

The solution of equation (1.1) requires some knowledge about the action of the twodimensional Laplacean in the spaces  $X^S$ . We define  $A: X^2 \to X^0$  by  $Au = \Delta u$ . The parameter  $\lambda$  is assumed to lie in  $[n^2\pi^2,(n+1)^2\pi^2)$  for some fixed  $n \in \mathbb{N}_0$ .  $L_v(\lambda) = D_y^2 + (\lambda - v^2\pi^2)$  denotes the operator generated by A on the Fourier components  $u_v$  of u resp.  $f_v$  of f (see Lemma 2.3).

#### LEMMA 3.1

Let be  $f \in X_k^0$ ,  $v \in \mathbb{N}$  and  $v \ge n+1$ . Then there exists a unique solution  $u_v \in Y_k^2$  of  $L_v(\lambda)u_v = f_v$ . It satisfies the inequality

(3.1) 
$$|D^{\Upsilon}u_{v}|_{0,\ell} \leq \frac{c_{1}(k)}{\sqrt{2-\Upsilon}} |f_{v}|_{0,k}, \quad \Upsilon = 0,1,2$$

for some constant  $c_1(k)$ , which is independent of  $\ell$  and  $\nu$ .

Proof: Since  $L_{_{0}}(\lambda)u=0$  has only exponentially growing solutions, the uniqueness is trivial. By means of the Green's function

$$G_{\nu}(y,n) = \frac{1}{2\omega_{\nu}} e^{-\omega_{\nu}(y-n)} \quad \text{for } n < y$$

$$G_{\nu}(y,n) = G_{\nu}(n,y) \quad , \quad \omega_{\nu} = \sqrt{\nu^{2}\pi^{2} - \lambda}$$

one obtains the solutions by the formula

(3.2) 
$$u_{\nu}(y) = \int_{R} G_{\nu}(y,n) f_{\nu}(n) dn$$

Setting  $f_v = g_k h_v$ , one obtains  $h_v \in L_2(IR)$  and, using Cauchy-Schwarz's inequality, y = 0 or 1:

$$\begin{split} \|D^{Y}u_{v}(y)\|^{2} &\leq \omega_{v}^{2\gamma} \int_{\mathbb{R}} G(y,n)g_{k}^{2}(n) dn \int_{\mathbb{R}} G(y,n)h_{v}^{2}(n) dn \\ &\leq \omega_{v}^{2\gamma-2} c_{1}(k) g_{k}^{2}(y) \int_{\mathbb{R}} G(y,n)h_{v}^{2}(n) dn \end{split}$$

Here, we used an estimate for the first integral which can be established by elementary calculations. Now, Fubini's theorem yields:

$$\|D^{\gamma}u_{\nu}\|_{0,\ell}^{2} \le c_{1}(k) \omega_{\nu}^{2\gamma-2} \int_{\mathbb{R}} h_{\nu}^{2}(n) \left(\int_{\mathbb{R}} G(y,n)g_{\ell-k}^{-2}(y) dy\right) dn$$

The inequality  $|g_{\ell-k}^{-2}(y)| \le 1$  for  $\ell \ge k$  and all  $y \in \mathbb{R}$  implies

$$\|D^{\gamma}u_{\nu}\|_{0,\ell}^{2} \leq c_{1}\omega_{\nu}^{2\gamma-4} \int_{\mathbb{R}} h_{\nu}^{2}(\eta) d\eta$$
 for  $\gamma = 0$  or 1.

The case  $\gamma$  = 2 follows from the differential equation. Observing  $\omega_{_{V}} \sim \nu$  completes the proof.

Define for s = 0,1,2 linear projections  $P^S: X^S \to X^S$ ,  $Q^S = id - P^S$  by  $P^Su = \sum_{v=1}^n u_v(y) \sin v\pi x$ .  $P^S$  and  $Q^S$  are continuous and  $P^2$ ,  $Q^2$  commute with A in the sense that  $AP^2 = P^OA$  resp.  $AQ^2 = Q^OA$  holds in  $X^2$ . Moreover,  $Q^SX^S$  is the inductive limit of  $Q^SX^S_k$  having all properties mentioned for  $X^S$  in § 2.

## LEMMA 3.2

Let  $B_2$  denote the restriction of  $A+\lambda$  onto  $Q^2X^0$ . Then,  $B_2: Q^2X^2 \to Q^0X^0$  is a topological isomorphism. Furthermore,

(3.3) 
$$\|u\|_{2,k} \le c_2(k) \|B_2u\|_{0,k}$$
  
holds for  $u \in Q^{2^0}X^2$ .

The proof that  $B_2$  is bijective is a consequence of Lemma 3.1 which, in addition, shows the continuity of  $B_2^{-1}$ . The continuity of  $B_2$  itself is trivial (either direct or by the open mapping principle). Inequality (3.3) follows from (2.4) and (3.1) immediately.

#### COROLLARY 3.3

Let  $f^{\rho}=f(\cdot,\cdot+\rho)$ ,  $\rho\in\mathbb{R}$ , be the  $\rho$ -translation of f. Then, for every  $f\in X^0$ ,  $(B_2^{-1}f)^{\rho}=B_2^{-1}f^{\rho}$ .

For the proof use the representation of  $B_2^{-1}$  given by (3.2) through its action on the Fourier components. Observe that  $G_{\nu}(y,n)$  only depends on the difference y-n.

#### LEMMA 3.4

Let  $\rho \ge 1$  be fixed and define  $K_{\ell} = [0,1] \times [(\ell-1)\rho,\ell\rho]$  for  $\ell=1,2,\ldots$  and  $K_{\ell} = [0,1] \times [\ell\rho,(\ell+1)\rho]$  for  $\ell=-1,-2,\ldots,\mathbb{Z}' = \mathbb{Z} \setminus \{0\}$ . Assume

for some  $f \in X^0$ .

Then there exists a  $\rho$ -independent constant  $\gamma_1$  such that

$$\sup_{\boldsymbol{\ell} \in \mathbb{Z}^*} \|B_2 Q^0 f\|_{H^2(K_{\ell})} \leq \gamma_1 \sup_{\boldsymbol{\ell}} \|f\|_{H^0(K_{\ell})}$$

is satisfied.

Proof: Observe that  $g_2^{-2} \ge \rho^2$  holds in  $K_1$  and use Lemma 3.2 to obtain

$$\begin{split} \|B_{2}^{-1}Q_{0}f^{\ell\rho}\|_{H^{2}(K_{1})} & \leq \rho^{2}\|B_{2}^{-1}Q_{0}f^{\ell\rho}\|_{2,2} \leq c_{2}(2)\rho^{2}\|f^{\ell\rho}\|_{0,2} \\ & = c_{2}(2)\rho^{2}\sum_{|k|=1}^{\infty}\|g_{2}^{-1}f^{\ell\rho}\|_{H^{0}(K_{k})} \end{split}$$

Since  $\lg_2^{-1}(y) \leq \min (1,1/(k-1)^2 \rho^2)$  for all  $y \in K_k$  and in view of

$$\sup_{k \in \mathbb{Z}^{+}} \|f^{\ell \rho}\|_{H^{0}(K_{k})} = \sup_{\ell \in \mathbb{Z}^{+}} \|f^{\ell \rho}\|_{H^{0}(K_{1})}$$

the inequality follows immediately.

It is useful subsequently to introduce

$$f|_{\ell} = f \times K_{\ell}$$
 ,  $\ell \in \mathbb{Z}'$ 

where  $x_{K_{\ell}}$  denotes the characteristic function of  $K_{\ell}$ , defined in Lemma 3.4. Instead of  $\|f\|_{\mathcal{L}^{\parallel}H(K_{2})}$  we use the less correct but unambigous notation  $\|f\|_{\mathcal{H}(K_{2})}$ .

#### LEMMA 3.5

Let be  $f \in Q^0 X^0$ , f(x,y) = 0 for  $(x,y) \in S_p = [0,1] \times (-\infty,p)$ . Then ,  $u \in B_2^{-1} f$  restricted to  $S_p$  has the form :

$$B_2^{-1}f|_{S_p} = \sum_{v=n+1}^{\infty} u_v(o) \sin v\pi x e^{\omega_v y}$$

Furthermore, the following inequality holds:

$$\|B_2^{-1}f\|_{H^2(K_{-1})} \le e^{-\rho \omega_{n+1}} \|B_2^{-1}f\|_{H^2(K_1)}$$

An analogous result is valid, when  $S_p$  is replaced by  $[0,1] \times (-\rho,\infty)$  and the roles of  $K_{-1}$  and  $K_1$  are interchanged.

Proof: The representation of u follows by a straightforward computation. To obtain the inequality use Lemma 2.3 and the following estimates:

$$\begin{split} \|D_{x}^{\beta}D_{y}^{\gamma}u\|_{H^{0}(K_{1})}^{2} &= \sum_{\nu=n+1}^{\infty} (\nu\pi)^{2\beta}\omega_{\nu}^{2\gamma} u_{\nu}(\rho)^{2} \int_{-\infty}^{\rho} e^{2\omega}\nu^{y} dy \\ &\geq \frac{e^{2\omega}n+1^{\rho}}{2} \sum_{\nu=n+1}^{\infty} (\nu\pi)^{2\beta}\omega_{\nu}^{2\gamma-1} u_{\nu}(\rho)^{2} (1-e^{-2\omega\nu^{\rho}}) \\ &= e^{2\omega}n+1^{\rho} \|D_{x}^{\beta}D_{y}^{\gamma}u\|_{H^{0}(K_{-1})}^{2}, \quad \text{qed.} \end{split}$$

Now, we study the action of  $A+\lambda$  on  $P^2X^2$ . The operator is still onto but an estimate of the form (3.2) cannot be expected. Allowing a higher order growth we arrive at estimates which are sufficient for the uniqueness proof in the next section.

#### LEMMA 3.6

Let be  $f \in X_k^0$  for some  $k \in \mathbb{N}_0$ , and assume  $v \le n$ . Then,  $L_v(\lambda)u_v = f_v$  possesses a solution which satisfies

$$10^{\gamma}u_{\nu}l_{0,2} \leq c_{3}(k) lf_{\nu}l_{0,k}$$
,  $\gamma = 0,1,2$ 

where  $\ell > k + \frac{3}{2}$  if  $\lambda = n^2 \pi^2$  and  $\nu = n$ , and  $\ell > k+1$  otherwise.

Proof: Define  $\phi$  by  $L_{\nu}(\lambda)\phi=0$ ,  $\phi(0)=0$ ,  $\phi'(0)=1$ . Consider first the case  $\nu< n$  or  $\lambda>n^2\pi^2$ . Then  $\phi$  and its derivatives are bounded and

$$u_{v}(y) = \int_{0}^{y} \varphi(y-\eta) f_{v}(\eta) d\eta$$

is the desired solution. With  $f_v=g_kh_v$  , one obtains  $h_v\in H^0(IR)$  and, if  $\gamma=0$  or 1 :

$$\begin{split} \|D^{Y}u_{v}\|_{0, 2}^{2} & \leq C(\phi) \int_{\mathbb{R}} g_{2}^{-2}((\int_{0}^{|y|} g_{k}^{2} dn) \int_{0}^{|y|} h_{v}^{2}(n) dn) dy \\ & \leq 2C(\phi) \|f_{v}\|_{0, k}^{2} \int_{\mathbb{R}} g_{2}^{-2}(k+\frac{1}{2}) dy \\ & \leq c_{3}(k) \|f_{v}\|_{0, k}^{2} \text{ if } 2 > k+1 . \end{split}$$

For  $\nu=n$  and  $\lambda=n^2\pi^2$  one has  $\phi(y)=y$  . The modified estimation yields

$$|D^{Y}u_{v}|_{0,k}^{2} \le c_{1} |f_{v}|_{0,k}^{2} \int_{\mathbb{R}} g_{\ell-(k+1)}^{2} dy \le c_{3}(k) |f_{v}|_{0,k}^{2}$$

if  $\ell > k + \frac{3}{2}$ . The inequalities for  $\gamma = 2$  follow from the differential equation, qed.

The proof shows that  $u_{\nu}$  is uniquely determined by the initial condition  $u_{\nu}(0) = u_{\nu}^{*}(0) = 0$ . To put this into the right framework define

$$Z_{1} = P^{2}X^{2}$$

$$\tilde{P} : X^{2} \rightarrow \ker(A+\lambda) , \quad \tilde{Q} = id - \tilde{P}$$

$$\tilde{P} = \sqrt{2} \int_{v=1}^{n} \sin v\pi x \left( \int_{0}^{1} u(x,0) \sin v\pi x \, dx \, \cos\sqrt{\lambda-v^{2}\pi^{2}} \, y \right)$$

$$+ \frac{1}{\sqrt{\lambda-v^{2}\pi^{2}}} \int_{0}^{1} D_{y}u(x,0) \sin v\pi x \, dx \, \sin\sqrt{\lambda-v^{2}\pi^{2}} \, y$$

if  $\lambda > n^2\pi^2$ ; otherwise, for  $\lambda = n^2\pi^2$ , we write  $\tilde{P}_{n-1}$  for the above expression when n is replaced by n-1 and define

$$\tilde{P}u = \tilde{P}_{n-1}u + \sqrt{2} \sin n\pi x \left( \int_{0}^{1} u(x,0) \sin n\pi x \, dx + y \int_{0}^{1} D_{y}u(x,0) \sin n\pi x \, dx \right)$$

Note that the trace  $D_y u(x,o)$  is well defined since  $u \in X^2$ . It is easily seen that  $\tilde{P}$  and  $\tilde{Q}$  are continuous projections,  $\tilde{P}_{|\ker(A+\lambda)} = id$  and  $\tilde{P}u = 0$  if and only if  $u_v(x,o) = D_y u_v(x,o)$  holds for  $0 \le v \le n$ . Hence, according to our previous Lemma

(3.5) 
$$B_{1} = (A+\lambda) |\tilde{Q}Z_{1}| : \tilde{Q}Z_{1} + P^{O}X^{O}$$

$$A_{1} = (A+\lambda) |\tilde{P}Z_{1}| : \tilde{P}Z_{1} \rightarrow \{0\}$$

 $\mathsf{B}_1$  is a topological isomorphism. Therefore, the decomposition

$$\tilde{X}^2 = \tilde{P}Z_1 \oplus \tilde{Q}Z_1 \oplus Q^2\tilde{X}^2$$

defines a decomposition of  $A+\lambda$  in the form

(3.6) 
$$A + \lambda = A_1 + B$$
,  $B = B_1 + B_2$ ,

where B is a top isomorphism between  $QZ_1 \oplus Q^2X^2$  and  $X^0$ .

#### THEOREM 3.7

Assume that  $n^2\pi^2 \le \lambda < (n+1)^2\pi^2$  holds for some  $n \in \mathbb{N}_0$ . Then,

(i) dim ker 
$$(A + \lambda) = 2n$$

(ii) 
$$A + \lambda = A_1 \oplus B$$

where  $A_1$  and B are given through (3.5), B:  $\tilde{QZ}_1 \oplus \tilde{Q}_1^2 \to \chi^0$  is a topological isomorphism and the following estimate holds

$$\|u\|_{2,\ell} \le c(k) \|Bu\|_{0,k}$$
 if  $\ell > k + \frac{3}{2}$ .

The following two Lemmas estimate the influence of the initial conditions  $u_v(0)$ ,  $u_v'(0)$  on the  $H^2(K_{\pm 1})$ -norm of  $u_v(0)$ , which will be crucial for a continuation argument in the uniqueness proof.

## LEMMA 3.8

Adopt the notations of Lemma 3.4 assume  $\ker(A+\lambda) \neq \{0\}$ . Then, there exists a p-independent constant  $\gamma_2$  such that

$$\|u\|_{H^{2}(K_{\pm 1})} \leq \gamma_{2}^{p^{k}}(\|f\|_{H^{0}(K_{\pm 1})} + \|\tilde{P}u\|_{2,\ell})$$

holds for  $k = \max(2, \ell)$ ,  $\ell \ge 1$ , and  $f \in P^{O}X^{O}$ ,  $(A+\lambda)u = f$ .

Proof: Observe that

$$(A+\lambda)v = f, \tilde{P}v = 0, f_{|1} = 0$$

yields  $v_{|1} = 0$ , since  $v \in P^{2^0_X 2}$  and  $v_{v}(0) = v_{v}(0) = 0$ ,  $1 \le v \le n$ . Hence

$$(B_1^{-1}f)_{|1} = (B_1^{-1}(f_{|1}))_{|1}$$
,

and we conclude from Lemma 3.6 :

$$\| u - P u \|_{H^{2}(K_{1})} = \| B_{1}^{-1}(f_{1}) \|_{H^{2}(K_{1})} \le C_{1}^{\rho} \| B_{1}^{-1}f_{1} \|_{2,2}$$

$$\le C_{2}^{\rho} \| f_{1} \|_{0,0} = C_{2}^{\rho} \| f \|_{H^{0}(K_{1})}$$

In view of  $u \in P^{202}$  we have

$$\|Pu\|_{H^{2}(K_{1})} \leq C_{3} \rho^{2} \|Pu\|_{2,2}$$

for every  $\ell \ge 1$ ,  $C_3$  independent of  $\ell$ ; whence

$$\|\|u\|_{H^{2}(K_{1})} \le C_{2^{\rho}}^{2} \|\|f\|_{H^{0}(K_{1})} + C_{3^{\rho}}^{\ell} \|\|Pu\|_{2,\ell}$$

which implies the assertion ( $\rho \geq 1$ ).

#### LEMMA 3.9

The following inequality is valid for all  $u \in P^{2^{\circ}X^{2}}$  and  $\ell > \frac{3}{2}$ :

$$\|\tilde{P}u\|_{2,\ell} \leq \gamma(\ell) \|u\|_{H^2(K_{\pm 1})}$$

Proof: Equation (3.4) immediately implies  $(0 \le \beta + \gamma \le 2)$ :

$$\|D_{x}^{\beta}D_{y}^{\gamma}\tilde{P}u\|_{0,2}^{2} \leq C_{1}\sum_{v=1}^{n}(u_{v}^{2}(0)+u_{v}^{2}(0))\int_{\mathbb{R}}(y^{2}+1)g_{2}^{-2}(y) dy$$

From Sobolev's imbedding theorem one obtains  $u(\cdot,0)$ ,  $D_y u(\cdot,0) \in L_2(0,1)$  for every  $u \in H^2(K_{-1})$ . Hence

$$\|u_{v}'(0)\|^{2} \le \int_{0}^{1} \|D_{y}u(x,0)\|^{2} dx \le C_{2}\|u\|_{H^{2}(K_{-1})}^{2}$$

and a similar inequality for  $[u_v(0)]^2$ . Thus, by Lemma 2.3 the assertion follows.

# § 4 Uniqueness

Our aim, to give a complete description of all small solutions of (1.1), requires a proof for the fact that there exists at most one "small" solution with given projection  $\tilde{P}u$ . The content of this section is a precise formulation of this statement.

We use the notations of the preceding section, in particular the constants  $\gamma_1, \gamma_2, \gamma_3$  appearing in various Lemmas. Equation (1.1) is now written in operator notation

(4.1) 
$$(A + \lambda)u + f(u) = 0$$
,  $u \in X^2$ 

where  $\lambda$  satisfies  $n^2\pi^2 \le \lambda \le (n+1)^2\pi^2$ .

## 4.1 Assumptions for f

 $f: H_{1oc}^{s}(\Omega) \to H_{1oc}^{o}(\Omega)$  for some s,  $0 \le s \le 2$ .

$$f_1 = f_{|H^S(K_1)} : H^S(K_1) \rightarrow H^O(K_1)$$

is supposed to be differentiable in a neighborhood of 0, and its derivative  $f_1'$  should be continuous at 0. Furthermore,  $f_1(0) = 0$ ,  $f_1'(0) = 0$ ,  $f(u^{\sigma}) = f^{\sigma}(u)$  for all  $\sigma \in \mathbb{R}$  is assumed.

Setting

$$p(u,v) = P^{O}(f(u)-f(v))$$

$$q(u,v) = Q^{O}(f(u)-f(v)) ,$$

We obtain from 4.1 quite easily, that for every  $\delta > 0$  there is a n > 0 - independent of 2 - such that

for all  $\ell \in Z'$  and all  $u,v \in H_{loc}^S(\Omega)$  with  $\|u\|_{H^S(K_\ell)}^{-1} \cdot \|v\|_{H^S(K_\ell)}^{-1} \le \eta$ . Since  $p^0,Q^0$  commute with  $u \mapsto u^0$  we have  $p(u^\sigma,v^\sigma) = p^\sigma(u,v)$  and similarly for  $q,\sigma \in \mathbb{R}$ .

For  $\ell \in \mathbb{N}$ ,  $k = \max(2, \ell)$  define constants  $\rho$ , a,  $\delta$ , e<sub>m</sub> satisfying the conditions :

$$a = \frac{1}{2} \min(1, (\gamma_2 \gamma_3 \rho^k)^{-1})$$

$$\rho \ge \max(-\frac{1}{\omega_{n+1}} \ln \frac{a}{2}, 1)$$

$$\delta = \min(a(2\gamma_2 \rho^k)^{-1}, a^2(12\gamma_1)^{-1})$$

$$\epsilon_m = a^m \epsilon_0, \quad \epsilon_0 \ge 2n$$

with n taken from (4.2). Observe that the sequence  $\boldsymbol{\epsilon}_m$  tends to 0.

# 4.2 Property I(m)

The pair  $u, v \in X^S$  is said to have property  $\pi(m)$ ,  $m \in \mathbb{N}_0$ , if

$$\sup_{\sigma} \| \mathbf{u}^{\sigma} \|_{H^{S}(K_{1})} \leq n , \quad \sup_{\sigma} \| \mathbf{v}^{\sigma} \|_{H^{S}(K_{1})} \leq n$$

$$\| \mathbf{u} - \mathbf{v} \|_{H^{2}(K_{j})} \leq \varepsilon_{m+1-|j|}$$

holds for  $j = \pm 1, \ldots, \pm m$ .

# LEMMA 4.1

Let  $u,v \in X^S$  have property  $\Pi(m)$  for all  $m \in \mathbb{N}_0$ . Then, for  $\tau = +1$  or  $\tau = -1$ , one has

$$\begin{split} \| \, B_2^{-1} \, & \sum_{j=2\tau}^{\tau^\infty} \, q_{j,j}(u,v) \|_{H^2(K_\tau)} \leq \frac{1}{6} \, \varepsilon_{m+1} \\ \| \, B_2^{-1} \, & \sum_{j=-\tau}^{\tau^\infty} \, q_{j,j}(u,v) \|_{H^2(K_\tau)} \leq \frac{1}{6} \, \varepsilon_{m+1} \, , \quad m \in \mathbb{N}_0 \, . \end{split}$$

Proof: We proceed by induction.  $\Pi(0)$  implies  $\|q(u,v)\|_{H^0(K_j)} \le \varepsilon_0 \delta$  for all  $j \in Z'$ . In view of Lemma 3.4 and (4.3) the inequalities hold for m = 0. Let them be satisfied for m-1. Since u,v have property  $\Pi(m)$  one obtains

 $\|(q(u,v))\|_{H^0(K_j)} < \delta \in m+1-|j|$  for |j| = 1,2,...,m.

Now, define  $\tilde{u} = \sum_{j=3}^{\infty} u_{j}$  and similarly v. Observe that  $\tilde{u}^{\rho}$ ,  $\tilde{v}^{\rho}$  possess property  $\pi(m-1)$ , whence

$$\|B_{2}^{-1}\int_{j=2}^{\infty} q_{jj}(\tilde{u}^{\rho},\tilde{v}^{\rho})\|_{H^{2}(K_{1})} \leq \frac{1}{6} \epsilon_{m}$$

Since we have  $(\sum_{j=2}^{\infty} q_{j}(\tilde{u}^{p},\tilde{v}^{p}))_{k} = 0$  for  $k = 1,0,-1,-2,\ldots$ , it follows, using Lemma 3.5, by assumption :

$$\begin{split} &\| B_{2}^{-1} \sum_{j=3}^{\infty} q_{|j}(u,v) \|_{H^{2}(K_{1})} = \| B_{2}^{-1} \sum_{j=2}^{\infty} q_{|j}^{\rho}(u,v) \|_{H^{2}(K_{-1})} \\ & \leq e^{-\rho \omega_{n+1}} \| B_{2}^{-1} \sum_{j=2}^{\infty} q_{|j}(\tilde{u}^{\rho}, \tilde{v}^{\rho}) \|_{H^{2}(K_{1})} \leq \frac{1}{6} \varepsilon_{m} e^{-\rho \omega_{n+1}} \end{split}$$

Moreover n(m-1) for u,v and Lemma 3.4 yield :

Hence, by (4.3), one obtains :

$$\|B_2^{-1}\sum_{j=2}^{\infty} q_{j,j}(u,v)\|_{H^2(K_1)} \leq \gamma_1 \delta \varepsilon_{m-1} + \frac{1}{6} \varepsilon_m \ \mathrm{e}^{-\rho \omega_{n+1}} \leq \frac{1}{6} \varepsilon_{m+1} \ .$$

Similarly one shows

$$\|B_{2}^{-1}\int_{j=-2}^{-\infty} q_{|j}(u,v)\|_{H^{2}(K_{-1})} \leq \frac{1}{6} \epsilon_{m+1}$$
.

Let  $\chi_{-}$  denote the characteristic function  $\chi_{[0,1]\times\mathbb{R}^{-}}$ . Note that  $u_{-}^{\rho}=u^{\rho}\chi_{-}$ ,  $v_{-}^{\rho}=v^{\rho}\chi_{-}$  have property  $\pi(m)$ , implying

$$\|B_{2}^{-1} \sum_{j=-1}^{\infty} q_{j}(u,v)\|_{H^{2}(K_{1})} = \|B_{2}^{-1} \sum_{j=-2}^{\infty} q_{j}(u_{-}^{p},v_{-}^{p})\|_{H^{2}(K_{-1})}$$

$$\leq \frac{1}{6} \varepsilon_{m+1} .$$

The last inequality for  $\tau = -1$  follows quite similarly, qed.

#### LEMMA 4.2

If  $u, v \in X^S$  have property  $\pi(m)$ , for all  $m \in \mathbb{N}_0$ , then

$$\|B_2^{-1}q(u,v)\|_{H^2(K_{+1})} \le \frac{1}{2} \epsilon_{m+1}$$

holds for every m & IN ..

Proof: Property  $\Pi(m)$  for u,v and (4.2) imply  $\|q\|_1\|_{H^0(K_1)} \le \delta \varepsilon_m$ . Thus, by Lemma 3.4, one has:  $\|B_2^{-1}q\|_1\|_{H^2(K_1)} \le \delta \gamma_1 \varepsilon_m$ . Therefore, in view of Lemma 4.1, we obtain

$$\begin{split} \|B_{2}^{-1}q\|_{H^{2}(K_{1})} & \leq \|B_{2}^{-1}q\|_{1}\|_{H^{2}(K_{1})} + \|B_{2}^{-1}\sum_{j=2}^{\infty} q\|_{j}\|_{H^{2}(K_{1})} \\ & + \|B_{2}^{-1}\sum_{j=-1}^{\infty} q\|_{j}\|_{H^{2}(K_{1})} \\ & \leq \delta \gamma_{1} \varepsilon_{m} + \frac{1}{3} \varepsilon_{m+1} \leq \frac{1}{2} \varepsilon_{m+1} \end{split}$$

# LEMMA 4.3

Let  $m \in \mathbb{N}_0$  be fixed. Assume two solutions  $u, v \in X^2$  of (4.1) have property  $\Pi(m)$  and satisfy Pu = Pv then,

$$\|P^2(u-v)\|_{H^2(K_j)} \leq \frac{\epsilon_{m+1-|j|}}{2}$$

holds for |j| = 1, 2, ..., m+1.

Proof: Since  $P^{0}(\Lambda+\lambda) = (\Lambda+\lambda)P^{2}$  one obtains

$$(A+\lambda)p^{2}(u-v) = p(u,v)$$

Now, Lemma 3.8 yields

(4.4) 
$$\|p^2(u-v)\|_{H^2(K_1)} \le r_2 e^{k} (\delta \|u-v\|_{H^2(K_1)} + \|p(u-v)\|_{2,k})$$

where k = max(2, x),  $k \ge 1$ . The same 'emma 3.9 implies

(4.5) 
$$\|\tilde{P}(u-v)\|_{2,\epsilon} \leq \gamma_3(\epsilon) \|P^2(u-v)\|_{2,\epsilon}$$

We fix & and show inductively that

(4.6) 
$$\|P^{2}(u^{j\rho}-v^{j\rho})\|_{H^{2}(K_{1})} < \frac{1}{2}(\epsilon_{m+1-j})$$

holds for  $0 \le j \le m$ . For j = 0 we obtain from (4.6) and  $\tilde{P}u = \tilde{P}v$  :

$$\|P^{2}(u-v)\|_{H^{2}(K_{1})} \leq \gamma_{2} \rho^{k} \delta \varepsilon_{m} < \frac{1}{2} \varepsilon_{m+1}$$

Assume the validity of (4.6) for j-1 < m. As we know,  $u^{j\rho}$ ,  $v^{j\rho}$  are solutions of (4.1); they have property  $\pi(m-1)$ . Hence, in view of the assumptions, one has:

Combining this inequality with (4.4) and (4.5) yields

$$\begin{split} \|P^{2}(u-v)\|_{H^{2}(K_{j})} &= \|P^{2}(u^{j\rho}-v^{j\rho})\|_{H^{2}(K_{1})} \\ &\leq \gamma_{2}\rho^{k}(\delta\varepsilon_{m-j} + \gamma_{3}(z)\frac{1}{2}(\varepsilon_{m+2-j})) \\ &< \frac{1}{2}(\varepsilon_{m+1-j}) \end{split}$$

Thus, (4.6) is valid for all n,  $0 \le j \le m$ . A similar argument holds for negative indices j, and the Lemma is proved.

#### THEOREM 4.4

Let be  $n^2\pi^2 \le \lambda \le (n+1)^2\pi^2$ ,  $n \in \mathbb{N}_0$ . Define A:  $\chi^2 + \chi^0$  by Au =  $\Delta u$  and assume 4.1 for f.

Then, there exists a positive  $\epsilon$  such that, for every two solutions u,v of equation (4.1), the conditions  $\tilde{P}u = \tilde{P}_V$  and  $\sup_{\sigma} \|u^{\sigma}\|_{H^{S}(K_1)} < \epsilon$ ,  $\sup_{\sigma} \|v^{\sigma}\|_{H^{S}(K_1)} < \epsilon$  imply u = v.

If  $ker(A+\lambda) = \{0\}$ ,  $\tilde{P}$  is the null operator, and Theorem 4.4 asserts local uniqueness of the trivial solution u=0.

Proof: We show inductively that u,v has property  $\mathbb{B}(m)$  for all  $m \in \mathbb{N}_0$ , yielding  $\|u-v\|_{H^2(K_j)} = 0$  for  $|j| = 1, 2, \ldots$  and thus u = v.

Set  $\varepsilon = n \le \varepsilon_0/2$ , then  $\Pi(0)$  holds for u,v. Assume, they have property  $\Pi(m-1)$ . Then,  $u^{j\rho}$ ,  $v^{j\rho}$  have property  $\Pi(m-1-ljl)$  for  $ljl = 0,1, \ldots, m-1$ . Since  $u^{j\rho}$ ,  $v^{j\rho}$  solve (4.1) one has:

$$\mathbb{Q}^2(u^{j\rho} - v^{j\rho}) = B_2^{-1} q(u^{j\rho}, v^{j\rho})$$

whence, via Lemma 4.2, one obtains

$$\begin{split} \| \varrho^2(u-v) \|_{H^2(K_{j+1})} &= \| \varrho^2(u^{j\rho}-v^{j\rho}) \|_{H^2(K_1)} \\ &\leq \frac{1}{2} \, \epsilon_{m-1\,j\,l} \quad \text{for } 0 \leq j \leq m-1 \\ \| \varrho^2(u-v) \|_{H^2(K_{j-1})} &\leq \| \varrho^2(u^{j\rho}-v^{j\rho}) \|_{H^2(K_{-1})} \end{split}$$

$$\leq \frac{1}{2} \epsilon_{m-ij} \quad \text{for } 0 \geq j \geq -(m-1)$$

Furth rmore, if  $ker(A+\lambda) + \{0\}$ , Lemma 4.3 implies

$$\|P^2(u-1)\|_{H^2(K_{j+1})} \leq \frac{1}{2} \epsilon_{m-1j1}$$

for  $j = 0, \pm 1, \dots, \pm (m-1)$ . If ker  $(\Lambda + \lambda) = (0), n^2$  is the idendity on  $\chi^2$ . Hence, in both cases, one obtains

 $\|\mathbf{u} - \mathbf{v}\|_{H^{2}(K_{j})} < \epsilon_{m-1,j}$  for  $\|\mathbf{j}\| = 1, \dots, m$ 

Therefore u, y have property I(m) and the Theorem is proved.

## COROLLARY 4.5

Assume Theorem 4.4. Then for every  $\mu \in \mathbb{R}$  there exist positive numbers  $\lambda_0(\mu)$ ,  $\eta_0(\mu)$  such that for every  $\lambda \in [\mu-\lambda_0, \mu+\lambda_0]$ ,  $\tilde{\mu} \in \ker(A+\mu)$ , there is at most one solution of (5.1) satisfying  $\tilde{P}u = \hat{u}$ ,  $\sup_{\sigma} \|u^{\sigma}\|_{H^{\infty}(K_1)} \leq \eta_0$ .

Proof: Consider in a neighborhood of  $(\mu,0)$  the mapping

 $T: (\lambda, u) \rightarrow Au + \lambda u + f(u)$ ,  $T: \mathbb{R} \times \mathbb{X}^2 + \mathbb{X}^0$ 

The derivative of T at  $(\mu,0)$  is given by  $T'(0)(\lambda,u) = Au + \mu u$  which has the kernel ker  $T'(0) = \mathbb{R} \times \ker (A+\mu)$ . Moreover,  $(\lambda,u) \mapsto (\lambda,Pu)$  defines a continuous, linear projection from  $\mathbb{R} \times X^2$  into ker T'(0). Now the proof proceeds literally as in Theorem 4.4.

While the last theorem proves that, for given  $\tilde{u} \in \tilde{P}_X^{0.2}$ , there is at most one solution u of (4.1) with  $\tilde{P}u = \tilde{u}$ , which is uniformly small in  $H_{loc}^{S}(\Omega)$ , it remains to be shown that there exists such a solution. The nature of this existence problem changes dramatically between the  $\lambda$  - intervals  $(\pi^2, 4\pi^2)$  and  $(n^2\pi^2, (n+1)^2\pi^2)$ ,  $n \ge 2$ .

For the first case we are able to prove, for arbitrary  $C^p$  - mappings f satisfying certain symmetry requirements, that all small solutions are periodic in y. Moreover, if a certain non-degeneracy condition is met, "singular" solutions exist which are locally uniform limits of solutions having arbitrarily large irreducible periods (see next section). In the case  $\lambda \in (n^2\pi^2, (n+1)^2\pi^2)$ ,  $n \geq 2$ , the question of existence is more delicate since problems of small denominations are involved. For real analytic f we are able to show that (4.1) can be solved by a formal power series. Convergence and generalisations to  $C^p$  - functions f are still open.

We are going to concrete the abstract properties of f, assumed in the preceding section. Let U be some neighborhood of 0 in IR and  $\tilde{f}: U \times IR^2 \to IR$  be a  $C^p$ -map,  $p \ge 2$ , satisfying  $\tilde{f}(0) = 0$  and  $D\tilde{f}(0) = 0$  for its gradient. We consider the equation

$$\Delta u + \lambda u + \tilde{f}(u, D_{x}u, D_{y}u) = 0$$
(5.1)
$$u(0,y) = u(1,y) = 0$$

for  $x \in \mathbb{R}$ . Since  $H_{1oc}^2(\Omega) \subset W_{\infty,loc}^1(\Omega)$ ,  $u,D_xu,D_yu$  are locally  $L_{\infty}$ -functions. It follows immediately that  $f: u \mapsto \tilde{f}(u,D_xu,D_yu)$  defines a  $C^D$  - mapping from the real space  $H_{1oc}^2(\Omega)$  into  $H_{1oc}^0(\Omega)$  which, in particular,

satisfies assumption 4.1 . Similarly, if f is real analytic for lul <  $\delta$ , lpl <  $\delta$ , lql <  $\delta$  then f, as a mapping from  $H^2(K)$  into  $H^0(K)$ , is analytic in the sense of [4], p.112 for every compact  $K \subset \Omega$ , if  $\|u\|_{H^2(K)} < \delta'$ .

The case  $\lambda \in (\pi^2, 4\pi^2)$  is considered first.  $\tilde{f}$  is assumed to be a  $C^p$ -function,  $p \geq 2$ . Periodic solutions are constructed via the ansatz  $u(x,y) = v(x,\omega y)$ , where v(x,z) is  $2\pi$ -periodic in z and where  $\omega$  varies near  $\omega^0 := (\lambda - \pi^2)^{1/2}$ . One obtains the boundary-value problem for v:

$$(\mathcal{L}(\omega) + \lambda)v + F(\omega, v) = 0$$
  
(5.2)  
 $v(0,z) = v(1,z) = 0, v 2\pi$ -periodic in z

where the following notations have been used

$$\mathcal{L}(\omega) = D_{xx}^2 + \omega^2 D_{zz}^2$$
(5.2a)
$$F(\omega, v) = \tilde{f}(v, D_x v, \omega D_z v)$$

Although only real solutions are of interest we work, for technical reasons, in complex Hilbert-spaces  $H^S_\#$ , the real subspaces of which are denoted by  $H^S_\#$ . Note that  $\mathcal L$  and F are defined on real spaces; therefore we have always to assure that they act on real elements.

Set  $\Omega_1 = [0,1] \times IR$  and define

$$\mathbb{H}_{\#}^{0} \equiv \{ v \in \mathbb{H}_{10C}^{0}(\Omega_{1}) / v \ 2\pi \text{-periodic in } z \}$$

with the scalar product

$$(v^1, v^2)_0 = \int_{(0,1)\times T_1} v^{1} \overline{v^2} d\mu$$
,

 $\nu$  denoting the Lebesque measure and  $T_1$  the interval (0,2 $\pi$ ).

Furthermore, we introduce :

$$\mathbb{H}^{02}_{\#} = \{ v \in \mathbb{H}^{2}_{loc}(\Omega_{1}) / v \in \mathbb{H}^{0}_{\#} \text{ and } v(0,\cdot) = v(1,\cdot) = 0 \}$$

with the scalar product

$$(v^{1}, v^{2})_{2} = \sum_{|\alpha| \le 2} (D^{\alpha}v^{1}, D^{\alpha}v^{2})_{0}$$

Let V denote a suitable neighbourhood of  $\omega^0$  in  $\mathbb{R}$ . Then,  $\mathcal{L}$  and F map V  $\times$   $\mathbb{H}^2_\#$  into  $\mathbb{H}^0_\#$ ;  $\mathcal{L}$  is smooth in  $\omega$  and linear in v, F is a  $C^p$ -mapping satisfying  $F(\omega,0) = D_vF(\omega,0) = 0$ .

For fixed  $\lambda \in (\pi^2, 4\pi^2)$ , the kernel of  $\mathcal{L}(\omega^0) + \lambda$  is spanned by the  $\mathbb{H}^0_\#$ -orthonormal system

(5.3) 
$$\varphi^{j}(x,z) = \pi^{-1/2} \sin|j|\pi x e^{izj}, j = \pm 1$$

where  $z_{-j} = -z_{j}$  per definition. Define  $P^{2}$  resp.  $P^{c}$  in  $H^{o}_{\#}$  resp. in  $H^{o}_{\#}$  by

(5.4) 
$$P^{s}v = \sum_{j = 1}^{s} (v, \varphi^{j})_{o}\varphi^{j}$$
,  $s = 0 \text{ or } 2$ .

 $\mathsf{P}^\mathsf{S}$  is a projection commuting with  $\mathcal Z$  in the following sense

$$P^{0}Z(\omega) = Z(\omega)P^{2}$$

and which, restricted to the real subspaces  $H_\#^0$  resp.  $H_\#^0$  acts as a real projection. Set  $Q^S \equiv id - P^S$  then, for sufficiently small  $|\omega - \omega^0|$ 

$$\mathcal{L}_{2}(\omega) = (\mathcal{L}(\omega) + \lambda) |_{\Omega^{2}H^{2}_{\#}}$$

is a topological isomorphism between  $\Omega^{2}H_{\#}^{2}$  and  $\Omega^{0}H_{\#}^{0}$ . Setting

$$c_{j} = (v, \varphi^{j})_{0}, w = \Omega^{2}v_{0}, c_{-j} = \overline{c_{j}}$$

( denoting complex conjugation), one obtains

(5.5a) 
$$\mathcal{L}_{2}(\omega)w = Q^{\circ}F(\omega, \sum_{\substack{j,j=1 \\ j \neq j}} c_{j} \varphi^{j} + w)$$

Since  $\|F(\omega,v)\|_0 = o(\|v\|_2)$ , equation (5.5a) yields, via the implicit funktion theorem, a unique solution  $w(\omega,\underline{c}) \in \mathbb{H}_{\#}^2$  as a real  $\mathbb{C}^D$ -funktion of  $\omega \in V(\omega^0)$ ,  $\underline{c} \in U(\underline{0})$ , where V and U are suitably chosen neighborhoods of  $\omega^0$  and  $\underline{0} \in \mathbb{C}^2$ . Furthermore, w satisfies

uniformly for  $\omega \in V(\omega^0)$ . Hence, near v=0,  $\omega=\omega^0$ , equation (5.2) is equivalent to

(5.5b) 
$$(\mathcal{Z}(\omega) + \lambda) \sum_{\substack{j \mid j \mid = 1}} c_j \varphi^j = P^0 F(\omega, \sum_{\substack{j \mid j \mid = 1}} c_j \varphi^j + w(\omega, \underline{c}))$$

For general F, existence of nontrivial solutions of (5.2), and thus of (5.5), cannot be expected without further assumptions Therefore, we require

a) 
$$F(\omega, v)(x, -z) = F(\omega, v(x, -z))$$
  
(5.6) or  
b)  $F(\omega, v)(x, -z) = -F(\omega, v(x, -z))$ ,  $F(\omega, -v) = F(\omega, v)$   
a.e. in  $\Omega_1$ 

## LEMMA 5.1

Assume  $c_{-1} = \overline{c_1}$ , and define  $s_j(\underline{c}) \equiv c_j e^{i\theta j}$ ,  $j = \pm 1$ ,  $\theta_{-j} = -\theta_j$ , then, the following idendities hold:

(i) 
$$Pv(x,z,s(\underline{c})) = Pv(x,z+\theta_1,\underline{c})$$

$$w(x,z,s(\underline{c})) = w(x,z+\theta_1,\underline{c})$$

(ii) 
$$Pv(x,z,\overline{c}) = Pv(x,-z,\underline{c})$$

$$w(x,z,\overline{c}) = w(x,-z,\underline{c}) , \text{ if } (5.6a) \text{ holds}$$

(iii) 
$$Pv(x,z,-\overline{c}) = -Pv(x,-z,\underline{c})$$

$$w(x,z,-\overline{c}) = -w(x,-z,\underline{c}) , \text{ if } (5.6b) \text{ holds}$$

Here, the  $\omega$ -dependence of w has been suppressed, P stands for any P<sup>S</sup>, and  $\overline{c}$  denotes the vector  $(\overline{c_1},\overline{c_{-1}})$ .

<u>Proof</u>: The proof for Pv follows by inspection. To prove (i) for w one shows that  $w_{\theta_1} \equiv w(\cdot, \cdot + \theta_1, \underline{c})$  satisfies (5.5a) if  $\underline{c}$  is replaced by  $\underline{s}(\underline{c})$ :

$$\mathcal{L}_{2}(\omega)w_{\theta_{1}} = (\Omega^{\circ}F(\omega,Pv+w))_{\theta_{1}} = \Omega^{\circ}F(\omega,Pv(s(\underline{c}))+w_{\theta_{1}})$$

The unique solvability of this equation yields (i). Assertion (ii) is shown quite similarly. To prove (iii) we set  $w^- = w(\cdot, -\cdot, \underline{c})$  and obtain

$$\mathcal{L}_{2}(\omega)(-w^{-}) = -(0^{\circ}F(\omega,Pv+w))^{-} = 0^{\circ}F(\omega,(Pv)^{-}+w^{-})$$

$$= 0^{\circ}F(\omega,-Pv(-\overline{\underline{c}}) - (-w^{-}))$$

$$= 0^{\circ}F(\omega,Pv(-\overline{\underline{c}}) + w^{-})$$

Again, the assertion follows, since (5.5a) is uniquely solvable.

Setting  $\tau \equiv \omega^2 - (\omega^0)^2$  and  $G(\tau,\underline{c}) \equiv F(\omega,(Pv+w))(\underline{c})$  equation (5.5b) can be written as follows:

$$(5.7)_{j}$$
  $-\tau c_{j} + (6(\tau, \underline{c}), \varphi^{j})_{0} = 0$ ,  $j = \pm 1$ 

First, we solve  $(5.7)_1$  in the subspace  $c_{-1} = c_1 = \overline{c_1}$ . Since  $G(\tau, \underline{0}) = 0$ , we have in this case,

$$N(\tau,\underline{c}) = \begin{cases} \frac{(G(\tau,\underline{c}), \varphi^1)}{c_1} & \text{if } c_1 = 0 \\ \\ D_{c_1}(G(\tau,c),\varphi^1) & \text{if } c_1 = 0 \end{cases}$$

as a  $C^{p-1}$ -function near  $\tau=0$ ,  $\underline{c}=\underline{0}$ . Then,  $(5.7)_1$  yields, via the implicit function theorem, a unique nontrivial solution  $\tau(c_1)$  which, for small  $|c_1|$  is a  $C^{p-1}$ -curve.

Observe that, in view of Lemma 5.1(i),

$$\frac{\left(G(\tau,\underline{c}),\varphi^{1}\right)_{0}}{c_{1}} \mid_{S(\underline{c})} = \frac{\left(G(\tau,\underline{c}),\varphi^{1}\right)_{0}}{c_{1}} \mid_{\underline{c}}$$

holds if  $c_1 \neq 0$ . Hence, by rotation, we obtain, from the solution constructed above, a solution  $\tau(\underline{c})$  of  $(5.7)_1$  for all small  $|\underline{c}|$ ,  $\underline{c} \in \mathbb{C}^2$  satisfying  $c_1 = \overline{c_1}$ .

It remains to be shown that the solution  $\tau(\underline{c})$  of  $(5.7)_1$  solves  $(5.7)_{-1}$  as well.Property (5.6a) and its implication in Lemma 5.1(ii) yields :

$$\frac{(G(\tau,\underline{c}),\varphi^{1})}{c_{1}} \mid_{\underline{c}} = \frac{(F(\omega,Pv+w),\varphi^{1})(\underline{c})}{\overline{c_{1}}}$$

$$= \frac{(F(\omega,Pv+w)(-z),\varphi^{1})(\underline{c})}{\overline{c_{1}}}$$

$$= \frac{(G(\tau,\underline{c}),\varphi^{1})}{\overline{c_{1}}}$$

The same equality is true if (5.6b) and hence Lemma 5.1(iii) holds. Since we have  $\tau(\overline{c}) = \tau(\underline{c})$ , equation (5.7)<sub>-1</sub> is satisfied.

#### THEOREM 5.2

Let be  $\lambda \in (\pi^2, 4\pi^2)$  and let  $F: V(\omega^0) \times H_\#^2 \to H_\#^0$  be a real  $C^p$ -map,  $p \ge 2$ , which satisfies (5.6a) or (5.6b). Then, for every sufficiently small  $|\underline{c}| \ \underline{c} \in \mathbb{C}^2, \ \overline{c_1} = c_{-1}$ , there exists a  $C^{p-1}$ -solution  $\omega(\underline{c})$ ,  $V(\underline{c})$  of (5.2) of the form

$$w(\underline{c}) = \omega^{0} + O(|\underline{c}|)$$

$$v(\underline{c}) = \sum_{j,j=1}^{n} c_{j} \varphi^{j} + o(|\underline{c}|)$$

We are now able to complete the characterisation of all small solutions by an existence result which complements the uniqueness theorem 4.4.

## THEOREM 5.3

Let be  $\lambda \in (\pi^2, 4\pi^2)$ , and let  $\tilde{\mathbb{P}}$  be defined by (3.4) (n=1). Assume that  $\tilde{\mathbf{f}}: \mathbb{U} \times \mathbb{R}^2 \to \mathbb{R}$  - U some neighborhood of 0 in  $\mathbb{R}$  - is a  $\mathbb{C}^p$ -map,  $p \geq 2$ , satisfying  $f(0,\underline{0}) = 0$ ,  $Df(0,\underline{0}) = 0$ . Suppose further that one of the following two conditions holds:

(a) 
$$\tilde{f}(u,p,-q) = \tilde{f}(u,p,q)$$

(b) 
$$\tilde{f}(u,p,-q) = -\tilde{f}(u,p,q)$$
 and  $\tilde{f}(-u,-p,-q) = \tilde{f}(u,p,q)$ 

Then, for every  $\tilde{u} \in \ker(A+\lambda)$  with sufficiently small norm there exists a y-periodic solution u of (5.1) with  $\tilde{P}u = \tilde{u}$ .

<u>Proof</u>: Note that the assumptions for  $\tilde{f}$  guarantee that  $F(\omega,v) \equiv \tilde{f}(v,D_Xv,\omega D_Zv)$  has property (5.6a) or (5.6b) and fulfills the regularity requirements of the preceding theorem. Let  $\omega(\underline{c})$ ,  $v(\underline{c})$  be a solution of (5.2). Then  $u(x,y) = v(x,\omega y)$  satisfies (5.1). Moreover,  $\tilde{P}u$  has the form:

$$\tilde{P}u(x,y) = \frac{2}{\sqrt{\pi}} (\alpha_1 \cos \omega^0 y - \beta_1 \sin \omega^0 y) \sin \pi x + o(|\underline{c}|),$$

$$c_1 = \alpha_1 + i\beta_1$$

where the remainder  $o(|\underline{c}|)$  is continuously differentiable with respect to  $\alpha$ ,8. Hence Pu = u can be solved using simply the implicit function theorem again, q.e.d. .

Now, we consider the general case  $\lambda \in (n^2\pi^2, (n+1)^2\pi^2)$  with  $n \in \mathbb{N}$  satisfying  $n \geq 2$ . Setting  $\omega_j^0 := (\lambda - j^2\pi^2)^{1/2}$  we construct quasiperiodic solutions via the ansatz  $u(x,y) := v(x,\omega_1y,\ldots,\omega_ny)$  where  $\underline{\omega} \equiv (\omega_1,\ldots,\omega_n)$  is close to  $\underline{\omega}^0$  in  $\mathbb{R}^n$ , and where v is  $2\pi$ -periodic in every  $z_j$  ("v is  $2\pi$ -periodic in  $\underline{z}$ "). With the notations

$$\mathcal{L}(\underline{\omega}) = D_{xx}^{2} + \int_{j,k=1}^{n} \omega_{j} \omega_{k} D^{2}_{z_{j}^{z_{k}}}$$

$$F(\underline{\omega},v) = f(v,D_{x}v,\underline{\omega},\nabla_{z}v)$$

we obtain the equation

$$(\mathcal{L}(\underline{\omega}) + \lambda)v + F(\underline{\omega}, v) = 0$$

$$(5.3)$$

$$v(0,\underline{z}) = v(1,\underline{z}) = 0 , v 2\pi\text{-periodic in }\underline{z}.$$

The formal similarity with (5.2) is not accidental; rather the preceding analysis carries over to the present case if  $\omega$ , z,  $\alpha_1$ ,  $\tau_1$ , and |j|=1 are replaced by  $\underline{\omega}$ ,  $\underline{z}$ ,  $\alpha_n\equiv [0,1]\times \mathbb{R}^n$ ,  $\tau_n\equiv (0,2\pi)^n$ , and |j|=1, ..., n. The functions

(5.9) 
$$\varphi^{j}(x,\underline{z}) = 2^{-\frac{n-1}{2}} \pi^{-\frac{n}{2}} \sin|j|\pi x \cdot e^{iz_{j}}, |j| = 1, ..., n$$

form a  $\mathbb{H}^0_{\#}$ -orthonormal basis of the 2n-dimensional kernel of  $\mathcal{L}(\underline{\omega}^0)$  +  $\lambda$ , a space which, quite obviously, consists of quasiperiodic solutions only. The projections  $P^S,Q^S$  as well as  $\mathcal{L}_2(\underline{\omega})$  are defined similarly to the case n=1. The fundamental difference here lies in the fact that  $\mathcal{L}_2(\underline{\omega})$ , though invertible, has no continuous inverse. Hence, the analogue of (5.5a)

(5.10) 
$$\mathcal{Z}_{2}(\underline{\omega})W = \Omega^{0}F(\underline{\omega}, \sum_{\substack{j=1 \ j \in I}}^{n} c_{j}\varphi^{j} + W)$$

cannot be solved as simply as before.

Spire gain

Let us assume that  $\tilde{f}$  is a real analytic function of its arguments. Then, since F then is a real analytic mapping from  $H^2_{\#}$  into  $H^0_{\#}$ , one can solve (5.10) by a formal power series

$$W = \sum_{\substack{|\alpha| \ge 0 \\ |\beta| \ge 2}} \underline{\tau}^{\alpha} \underline{c}^{\beta} W_{\alpha\beta}$$

-  $\alpha, \beta \in \mathbb{N}_0^n$ , [a], [B] their lengths,  $\underline{\tau} = (\tau_1, \dots, \tau_n)$ ,  $\tau_j = \omega_j^2 - \omega_j^2$ . Observe that, in view of the injectivity of  $\mathcal{X}_2(\underline{\omega}^0)$ , the coefficients  $N_{\alpha\beta}$  can be determined recursively. The proof of convergence by standard arguments fails however since, with increasing  $\beta$ , the inverse of  $\mathcal{X}_2(\underline{\omega}^0)$  grows indefinitely in norm.

While the proof of unique solvability of (5.10) seens to be a formidable task – and we are not yet able to accomplish it – the solution of the system of bifurcation equations analogous to (5.5b) is not. To show this, let us assume that (5.10) has a unique real analytic solution for sufficiently small |c| and |t|, |c| = |c| . The following Lemma is an immediate consequence of this assumption.

# LEHHA 5.4

Assume  $c_{-j} = \overline{c_j}$  for  $j = 1, \ldots, n$  and denote by  $c^{\ell} \in \mathbb{C}^{2n}$  a vector with  $c_{\ell} = \overline{c_{\ell}} = 0$ . Then, the solution  $w(\underline{\omega},\underline{c})(x,\underline{z})$  of (5.10) is independent of  $z_{\ell}$ .

The bifurcation equations, corresponding to (5.5b), are as follows:

$$(5.11) \qquad -\tau_{\ell} c_{\ell} + (G(\underline{\tau},\underline{c}),\varphi^{\ell})_{0} = 0 , \quad \tau_{-\ell} = \tau_{\ell} ,$$

$$|\ell| = 1, \ldots, n .$$

where

$$G(\underline{\tau},\underline{c}) \equiv F(\underline{\omega}, \sum_{j=1}^{n} c_j \varphi^j + w(\underline{\omega},\underline{c}))$$

Observe that Lemma 5.1 still holds for  $|j|=1,\ldots,n$ . As in (5.7), we solve (5.11) first for  $\ell=1,\ldots,n$  and  $c_{-j}=c_{j}=\overline{c_{j}},\ j=1,\ldots,n$ . G $(\underline{\tau},\underline{c}^{\ell})$  is independent of  $z_{\ell}$  according to the preceding Lemma. Hence  $(G(\underline{\tau},\underline{c}),\varphi^{\ell})_{0}=0$  if  $c_{\ell}=0$ . Thus

$$N_{\ell}(\underline{\tau},\underline{c}) = \begin{cases} \frac{(G(\underline{\tau},\underline{c}),\phi^{\ell})_{0}}{c_{\ell}} & \text{if } c_{\ell} \neq 0 \\ \\ D_{c_{\ell}}(G(\underline{\tau},\underline{c}),\phi^{\ell})_{0} & \text{if } c_{\ell} = 0 \end{cases}$$

is an analytic map near  $\underline{\tau} = \underline{c} = \underline{0}$  and  $(5.11)_{\ell}$ ,  $\ell = 1, \ldots, n$ , yields locally a unique analytic solution  $\underline{\tau}(\underline{c})$ . Again it follows, via Lemma 5.1, that

$$\frac{\left(G\left(\underline{\tau},\underline{c}\right),\varphi^{\ell}\right)_{0}}{c_{1}}\Big|_{S\left(\underline{c}\right)} = \frac{\left(G\left(\underline{\tau},\underline{c}\right),\varphi^{1}\right)_{0}}{c_{1}}\Big|_{\underline{c}}$$

for every  $c_1 \neq 0$ . Therefore, the system  $(5.11)_{\ell}$ ,  $\ell=1,\ldots,n$ , has a unique analytic solution  $\underline{\tau}(\underline{c})$  in some neighborhood of  $\underline{0} \in \boldsymbol{t}^{2n}$  satisfying  $c_{-j} = \overline{c_{j}}$ ,  $j=1,\ldots,n$ . The proof, that  $\tau_{-\ell} = \tau_{\ell}$  fulfil as well the equations  $(5.11)_{\ell}$ ,  $\ell=-1,\ldots,-n$ , proceeds literally as in the periodic case.

# PROPOSITION 5.5

Let be  $\lambda \in (n^2\pi^2, (n+1)^2\pi^2, 2 \le n \in \mathbb{N}$  and let  $\tilde{P}$  be defined by (3.4). Assume that  $f: U \times \mathbb{R}^2 \to \mathbb{R}$  be real analytic near the origin, satisfying the conditions of Theorem 5.3. Furthermore, suppose that equation (5.10) has a unique real analytic solution near  $\underline{c} = \underline{0}$  and  $\underline{\omega} = \underline{\omega}^0$ . Then there exists, for every  $\tilde{u} \in \ker(\Lambda + \lambda)$  with sufficiently small norm, a solution u of (5.1) with  $\tilde{P}u = \tilde{u}$ , which is quasiperiodic in y.

## § 6 Existence of Singular Solutions

In this section we consider the set of solutions near u=0,  $\lambda=\pi^2$ . Either there exist nontrivial solutions of (5.1) for  $\lambda=\pi^2$  or there are solutions for  $\lambda>\pi^2$ , which are local  $H^2$ -limits of periodic solutions with infinitely growing irreducible periods.

We restrict  $\lambda$  to the interval  $(\pi^2, 4\pi^2)$ . According to the preceding section the set of solutions with small  $H^S(K_g)$ -norm consists of periodic functions exclusively which form a two-dimensional manifold over ker  $(A+\lambda)$ . Let be  $\omega_0 = (\lambda - \pi^2)^{1/2}$ ,  $u(x,y) = v(x,\omega y)$  and assume v(x,z) to be  $2\pi$ -periodic in z. Define the real Hilbert-space  $H_e^S$  for arbitrary real  $s \ge 0$  as follows:

$$H_e^s(\Omega_1) = \{v \in H_+^s(\Omega_1) \mid v(x,-z) = v(x,z)\}$$

Similarly  $H_e^{0s}(\Omega_1)$ , s>1, denotes the corresponding subspace of  $H_{\#}^{s}$ .

#### 6.1 HYPOTHESIS

Assume that  $\tilde{\mathbf{f}}$  in (5.1) has the properties given there and that  $\tilde{\mathbf{f}}(\mathbf{u},\mathbf{p},\mathbf{-q})=\tilde{\mathbf{f}}(\mathbf{u},\mathbf{p},\mathbf{q})$  holds. In addition we suppose that the mapping  $\mathbf{f}:\mathbf{u}\mapsto\mathbf{f}(\mathbf{u},\mathbf{D}_{\mathbf{x}}\mathbf{u},\mathbf{D}_{\mathbf{y}}\mathbf{u})$  from  $\mathbf{H}_{\mathrm{loc}}^{\mathbf{S}}(\mathbf{g})$  into  $\mathbf{H}_{\mathrm{loc}}^{\mathbf{O}}(\mathbf{g})$  is continuously differentiable near  $\mathbf{u}=0$  for some real  $\mathbf{s}<2$ .

Under this hypothesis  $F(\omega, v) = f(v, D_X v, \omega D_Z v)$  is continuously differentiable near v = 0 and satisfies (5.7a) as well as  $F(\omega, 0) = 0$ ,  $D_v F(\omega, 0) = 0$ . Observe that every polynomial in u with  $\tilde{f}(0) = D\tilde{f}(0) = 0$  satisfies the assumption 6.1 for  $\tilde{f}$ .

Note that  $\mathcal{Z}(\omega)$  - defined in (5.2a) - has a continuous inverse  $\mathcal{Z}^{-1}(\omega)$ :  $H_e^0 \to H_e^2$  which is compact as a mapping from  $H_{\ell}^0$  into  $H_{\ell}^s$ . Hence,  $\mathcal{Z}^{-1}(\omega)F(\omega,\cdot)$  defines a completely

continuous operator in  $\mathbb{H}_e^s$  ahich satisfies  $\mathbb{H}_e^{-1}(\omega)F(\omega,v)\mathbb{H}_s = o(\mathbb{H}v\mathbb{H}_s)$  uniformly for  $\omega$  in compact subsets of  $\mathbb{R}^+$ . Equation (5.2) can be written as follows

(6.1) 
$$v = -\lambda x^{-1}(\omega)v - x^{-1}(\omega)F(\omega,v)$$

We intend to apply the global result of Rabinowitz [10] for fixed  $\lambda$  and variable  $\omega$ . Although our  $\omega$ -dependence is somewhat more general than in [10], the global existence still carries over since, at  $\omega = \omega_1^0 = (\lambda - \pi^2)^{1/2}$ , a simple eigenvalue of  $-\lambda^{-1}(\omega)$  in  $H_e^S$  crosses at the point 1 the unique circle with nonvanishing velocity.

Denote by  $C_{\lambda}^{+}$  resp.  $C_{\lambda}^{-}$  the connected components constructed in [13]. Theorem 1.40, which neet the point  $(\omega_{1}^{\circ 2},0)$  in  $\mathbb{R}^{+} \times \mathbb{H}_{2}^{s}$ . Then each of them has one of the following properties :

(i) 
$$C_{\lambda}^{\pm}$$
 is unbounded

(ii) 
$$C_{\lambda}^{\pm}$$
 meets another point  $(\tilde{\omega}^2,0)$ 

(iii) 
$$C_{\lambda}^{\pm}$$
 meets the subspace (0)  $\times H_{e}^{S}$ 

Subsequently we will show that alternative (iii) holds. In view of equation (6.1) we may consider  $\mathbb{C}^{\pm}_{\lambda}$  as subsets of  $\mathbb{R}^{+} \times \mathbb{H}^{2}_{e}$  having the same properties as in  $\mathbb{R}^{+} \times \mathbb{H}^{s}_{e}$ .

Define 
$$S^{+} \equiv \{v \in H_{2}^{2} \mid v_{1} = \sqrt{2} \quad \begin{cases} v(x, \cdot) \sin \pi x \, dx \, has \, exactly \\ \text{two simple zeroes in } [0, 2\pi) \, and \, v_{1}(0) > 0 \ \end{cases}$$

similarly S with  $v_1(0) < 0$ . Since we consider functions which are even in z and since, for  $v \in H_e^2$ ,  $v_1$  is continuously differentiable in R, the sets S are open in  $H_e^2$ . Moreover, near  $(\omega_0^2,0)$ ,  $C_\lambda^+ = (\omega_0^2,0)$  is a subset of  $R^+ \times S^+$  and similarly for  $C_\lambda^-$ . Observe that  $v \in S^+$  and v periodic in z implies that v has the irreducible period  $2\pi$ , i.e.  $2\pi$  is the largest possible period. If  $u(x,y) = v(x,\omega y)$  then u has the irreducible period  $2\pi$ .

Define

$$r_{\lambda} = \{u \in \mathbb{N}_{10c}^2 \mid u(x,y) = v(x,y), \frac{2}{\lambda}, v \in C_{\lambda}^{\pm}\}$$

r a r u r a

It is easily seen that  $\Gamma_{\lambda}$  as well as  $\Gamma_{\lambda}^{\pm}$  are connected subsets of  $H_{loc}^{2}(\Omega)$ . Subsequently we use the following notation :

$$E(u) = \sup_{\sigma \in \mathbb{R}} \|u^{\sigma}\|_{H^{2}}(\kappa_{1})$$

We have two alternatives which are treated seperat ly:

I An  $\varepsilon_0$  > 0 can be found such that for every  $\varepsilon \in (0,\varepsilon^{-1}]$  there exists a sequence  $(\lambda_n,u_n)$  of solutions of (5.1) with  $\lambda_n + \pi^2 + 0$ ,  $u_n \in \Gamma_{\lambda_n}$  and  $E(u_n) = \varepsilon$ .

II For every  $\epsilon$  > 0 there exists a  $\delta(\epsilon)$  > 0 such that

$$\sup_{u \in \Gamma_{\lambda}} E(u) < \varepsilon$$

$$0 < |\lambda - \pi^{2}| < \delta(\varepsilon)$$

holds.

#### THEOREM 6.2

Suppose hypothesis 6.1 and I hold. Then, for  $\lambda=\pi^2$ , and for all  $\epsilon$  in  $[0,\epsilon_0]$  there exists a solution  $u_\epsilon$  of (5.1) with  $E(u_\epsilon)=\epsilon$ .

Proof: Choose the sequence  $(\lambda_n, u_n)$  from I. We may assume that  $E(u_n) \leq \|u_n\|_{H^2(K_1)}$  ( $u_n^\sigma$  is again a solution). Since  $u_n^{(2)}(K_j)$  lies compactly in  $u_n^{(3)}(K_j)$ ,  $u_n^{(3)}(K_j)$ ,  $u_n^{(3)}(K_j)$ . Therefore,  $u_n^{(3)}(K_j)$  and  $u_n^{(3)}(K_j)$  converge in  $u_n^{(3)}(K_j)$  and are bounded in  $u_n^{(3)}(K_j)$  uniformly in j. Hence, by Lemma 2.2, they converge in  $u_n^{(3)}(K_j)$ . Using the equation

(6.2) 
$$u_{n_i} = -\lambda_{n_i} A^{-1} u_{n_i} - A^{-1} f(u_{n_i})$$

where A, defined in § 3 , is a topological isomorphism from  $x^0$  into  $x^0$ , we conclude that  $(u_n)$  converges in  $x^0$ . Thus  $u \in x^0$  is a solution of (5.1) which, in view of  $E(u_n) \leq \|u_n\|_{H^2(K_1)}$  satisfies  $E(u) = \varepsilon$ .

Now we turn to case II. Let be  $J\subset\mathbb{R}^+$  be compact set; then there are positive constants  $C_1,C_2$  - independent of u and w - such that

(6.3) 
$$C_1 \| \mathbf{v} \|_2 \leq \mathbf{E}(\mathbf{u}) \leq C_2 \| \mathbf{v} \|_2$$

holds for  $v(x, \omega y) = u(x, y)$ ,  $\omega \in J$ .

### LEMMA 6.3

Assume II and  $\lambda - \pi^2$  sufficiently small. Then

$$C_{\lambda}^{\pm} \sim \{(\omega_0^2, 0)\} \subset \mathbb{R}^+ \times S_{\lambda}^{\pm}$$

Proof : Take  $\lambda=\pi^2<\min(\delta(\eta_0),\lambda_0)$  where  $\eta_0,\lambda_0$  are from Corollary 4.5.Then, using II , we obtain :

(6.4) 
$$\sup_{u \in \Gamma_{\lambda}} E(u) < \eta_0$$

Furthermore, near  $(\omega^2,0)$  we have  $(\omega^2,v)\in\mathbb{R}^+\times S^+$  if v+0 and  $(\omega^2,v)\in C_\lambda^+$ .

If  $C_{\lambda}^{+} = \{(\omega_{0}^{2},0)\} \notin \mathbb{R}^{+} \times S_{\lambda}^{+}$  we show that  $C_{\lambda}^{+}$  contains a point  $(\tilde{\omega}^{2},0)$ , with  $\tilde{\omega}^{2} \neq \omega_{0}^{2}$ . Indeed, one would have a noint  $(\tilde{\omega}^{2},v) \in C_{\lambda}^{+}$  with  $v \in \partial S_{\lambda}^{+}$  (boundary of  $S_{\lambda}^{+}$ ). Since v is even in z,  $v_{1}$  and thus  $u_{1}$  has a double zero in  $[0,2\pi)$ . We may sume  $u_{1}(0) = u_{1}^{*}(0) = 0$ . Corollary 4.5 yields u = 0 and here v = 0.

If  $(\omega_0^2,0)$  is the only point in  $C_\lambda^+$  with v=0, the assertion is proved. Since, for some neighborhood U of  $(\omega_0^2,0)$   $C_\lambda^+$  O U and  $C_\lambda^-$  O U are disjoint,  $C_\lambda^+$  cannot leave  $R_\lambda^+$  x S near  $(\omega_0^2,0)$  and hence everywhere. Up to here, the proof for  $C_\lambda^-$  is the same.

It remains to show that  $(\omega_0^2,0)$  is the only point in  $C_\lambda = C_\lambda^+ \cap C_\lambda^-$  with v=0. Assume the contrary: Then,there exists a sequence  $(\omega_n^2,v_n) \in C_\lambda$   $\cap (\mathbb{R}^+ \times S)$ ,  $S=S^+ \cup S^-$ , such that  $\|v_n\|_2 \to 0$ ,  $\omega_n^2 \to \omega^2$  with  $\tilde{\omega}^2 \neq \omega_0^2$ . Therefore  $v_{n,1}(0) \to 0$ ,  $v_{n,1}(0) = 0$  holds which implies  $u_{n,1}(0) \to 0$ ,  $u_{n,1}(0) = 0$ . In view of (3.4) and Theorem 4.4 we conclude that  $u_n$  must coincide with the solutions constructed in Theorem 5.5 for large n. But those solutions have period near  $\frac{2\pi}{\omega_0}$  whereas  $u_n$  has period near  $\frac{2\pi}{\omega}$ . Since  $v_n \in S$ , the periods are irreducible and the contradiction follows.

Subsequently, the constants  $\eta_0$  and  $\lambda_0$  are taken from Corollary 4.5 for  $\mu=\pi^2$ .

## LEMMA 6.4

Suppose that case II holds. Then, for sufficiently small  $\lambda = \pi^2$ , the projection of  $C_{\lambda}^{\pm} \subset \mathbb{R}^+ \times \overset{02}{\text{e}}$  into  $\mathbb{R}^+$  forms an interval (0,b], b > 0.

Proof: We show first that  $C_{\lambda}^{+}$  is contained in some interval  $(0,b_{0}]$ . Otherwise, since  $\omega_{0}^{2} < 3\pi^{2}$  for  $\pi^{2} < \lambda < 4\pi^{2}$ , there is a sequence  $C_{\lambda_{n}}^{+}$ ,  $\lambda_{n} \rightarrow \pi^{2}$ , such that  $C_{\lambda_{n}}^{+}$  intersects  $\{3\pi^{2}\}$  x  $\mathring{\mathbb{H}}_{e}^{2}$  for every n. Hence, we have a sequence  $u_{n} \in \Gamma_{\lambda_{n}}$  with period  $2\pi/\sqrt{3}\pi$  in y satisfying - in view of II -  $E(u_{n}) \rightarrow 0$ . Dividing (6.2) by  $E(u_{n})$  yields:

$$\frac{u_n}{E(u_n)} = -\lambda_n \Lambda^{-1} \left( \frac{u_n}{E(u_n)} + \frac{f(u_n)}{E(u_n)} \right)$$

As in the proof of Theorem 6.2 one obtains a subsequence  $u_i = u_{n_i}/E(u_{n_i})$  converging towards some u in  $X_1^0$ . Furthermore,  $f(u_{n_i})/E(u_{n_i}) + 0$  in  $H^0_{1oc}(\Omega)$ , hence in  $X_1^0$ . The above equation

implies  $u_i \rightarrow u$  in  $X^2$  and thus  $(\Lambda + \pi^2)u = 0$ . Since all  $u_i$  have the irreducible period  $2\pi/\sqrt{3}\pi$  and since  $E(u_i) = 1$  holds, u possesses the same properties. However, there is no such solution in ker  $(\Lambda + \pi^2)$ . Therefore, the projection of  $C_{\lambda}^+$  into  $R^+$  is contained in some interval  $(0,b_0]$ .

Now, let be  $\lambda-\pi^2<\min(\delta(n_0),\lambda_0)$ . In view of II we have sup E(u) <  $n_0$  which, by (6.3) yields:  $u\in r_\lambda$ 

$$(\omega_{2}^{2 \operatorname{sup}}, v) \in C_{\lambda}^{+} \| v \|_{2} \leq \frac{1}{C_{1}} n_{0}$$

$$\omega_{\epsilon}^{2} [a, b_{0}]$$

for every a  $\in$  (0,b<sub>o</sub>), where the constant C<sub>1</sub> only depends on a. Therefore, C<sub>\lambda</sub><sup>+</sup> possesses property P<sub>3</sub>. Since C<sub>\lambda</sub><sup>+</sup> is connected and closed, the same is true for its projection on  $\mathbb{R}^+$  which proves the assertion for C<sub>\lambda</sub><sup>+</sup>. An identical argument holds for C<sub>\lambda</sub><sup>-</sup>.

We call a nontrivial solution of (5.1) singular, if it is the  $H^2_{loc}(\Omega)$ -limit of functions in  $\Gamma_{\lambda}$  but does not belong to  $\Gamma_{\lambda}$ . Note that such a singular solution u is nonperiodic in y or  $u_1$  is constant.

## THEOREM 6.5

Let II and hypothesis 6.1 be valid and assume  $\lambda - \pi^2 > 0$  be sufficiently small. Then, there is an interval  $J = (\alpha, \beta) \subset \mathbb{R}$  containing 0 such that, for every  $a \in J$ , there exists an unique solution u of (5.1) with  $E(u) < \eta_0$  and  $u_1(0) = a$ ,  $u_1'(0) = 0$ . Among these solutions, there are those of arbitrary large period.

Furthermore, there exist two singular solutions  $u^j$  satisfying  $u_1^1(0) = \alpha$ ,  $u_1^2(0) = \beta$ ,  $u^j(0) = 0$ ,  $E(u) \leq n_0$ . The map  $a \leftrightarrow u$  from  $[\alpha,\beta]$  into  $X^2$  is continuous.

Proof: According to II we have (6.4). The linear map  $\phi: u \to u_1(0)$  from  $H^2_{\log 1}(\Omega)$  into R is continuous and injective by Corollary 4.5. Hence,  $\phi^{-1}$  is continuous. Since  $\Gamma_{\lambda}$  is connected and bounded in  $H^2_{\log 1}(\Omega)$ ,  $J = \phi \Gamma_{\lambda}$  is an interval with  $0 \in J$ . We show:  $\delta$  is open.

Take a curve  $s_1$  in  $C_\lambda^+$  connecting  $(\omega_1^2,v^1)\in C_\lambda^+$  and  $(\omega_0^2,0)$  and let  $\sigma_1\equiv\{u\in\Gamma_\lambda/u(x,y)=v(x,\omega y),(\omega^2,v)\in s_1\}$  be the trace in  $\Gamma_\lambda$ . Then  $\Phi\sigma_1$  is an interval  $[0,\beta_1]$ . If  $(\omega_2^2,v^2)\notin s_1$ , we obtain, in view of Corollary 4.5, for the corresponding interval  $[0,\beta_2]$ ,  $\beta_2>\beta_1$ . Therefore  $\Phi(u^1)\in J$  (  $u^j(x,y)=v^j(x,\omega^j y)$ ) and, since a similar argument holds for  $C_\lambda^-$ , J is open.

In view of Lemma 6.4., to every  $(\omega_1^2, v^1) \in C_\lambda^+$  there exists a  $(\omega_2^2, v^2) \in C_\lambda^+$ ,  $(\omega_2^2, v^2) \notin s_1$  and with  $0 < \omega_2^2 < \omega_1^2$ . Thus we obtain a sequence  $u_n \in \Gamma_{\lambda_n}^+$  with  $u_{n,1}(0) \neq \alpha$  having period  $2\pi/\omega_n$  which increase indefinitely. We conclude — as in the proof of Theorem 6.2 — that a subsequence  $u_{n_1}$  converges towards a solution  $u^1$  of (5.1) in  $H_{10c}^2(\Omega)$  satisfying

(6.5a) 
$$u_1^1(0) = \alpha$$
,  $u_1^1'(0) = 0$ ,  $E(u^1) \le \eta_0$ .

Similarly one constructs a solution u<sup>2</sup> with

(6.5b) 
$$u_1^2(0) = \beta$$
,  $u_1^2(0) = 0$ ,  $E(u^2) \le \eta_0$ .

The continuity of the map

$$a \mapsto \begin{cases} \phi^{-1}a & \text{for } a \in J \\ u^j & \text{on the boundary of } J \end{cases}$$

acting from  $[\alpha,\beta]$  into  $\overset{0}{X^2}$ , follows in  $\overset{0}{J}$  from the continuity of  $\phi^{-1}$ . At the boundary we argue as follows: Every sequence  $u_n$  satisfying  $u_{n,1}(0) \to \alpha$  (or  $\beta$ ), converges towards  $u^1$  (or  $u^2$ ) in  $H^2_{loc}(\Omega)$  since - by Corollary 4.5 - the  $u^j$ 's are uniquely determined through (6.5).

To prove that the solutions  $u^j$  are singular, i.e.  $u^j \notin \Gamma_{\lambda}$ , we show that  $u^j_1$  is either constant or not periodic. Note that the orbits  $\ell_u = \{(u_1(y), u_1'(y))/ y \in |R|\}$  of periodic functions  $u \in \Gamma_{\lambda}$  in  $|R|^2$  are simply closed curves enclosing the origin. Since, according to Corollary 4.5, the  $\ell_u$ 's do not intersect and since

 $\Gamma_{\lambda}$  is connected in  $H^2_{loc}(\Omega)$ , the set  $U = \bigcup_{u \in \Gamma_{\lambda}} \ell_u$  forms an open neighborhood of 0. The orbits of  $u^j$ , j=1,2, must belong to the boundary  $\partial U$ . Hence, if  $u^j$  is not constant but periodic, it possesses at least two simple zeros per period. Since  $u^j_1$  is the local  $C^1$ -limit of functions with arbitrary large irreducible period which have exactly two simple zeros per period, a contradiction follows.

If hypothesis 6.1 b) is valid a similar result holds in the space of odd functions :

$$H_{\sigma}^{s}(\Omega_{1}) \equiv \{v \in H_{\#}^{s}(\Omega_{1}) / v(x,-z) = -v(x,z) \}$$

Choose for  $S^{\pm}$  the set of those  $v \in H^{0}_{\sigma}$  which have exactly two simple zeroes in  $[0,2\pi)$  and satisfy  $v_{1}'(0) > 0$  (< 0). In this case one obtains solutions u of (5.1) with  $u_{1}(0) = 0$ ,  $u_{1}'(0) = a$  for all a in the closure of J.

### COROLLARY 6.6

Let  $\tilde{f}$  satisfy the hypothesis 6.1 a) or 6.1 b). Moreover, assume II to hold and take  $\lambda-\pi^2>0$  sufficiently small. Then, there exists an open neighborhood U of 0 in ker  $(A+\lambda)$  and a continuous map  $\Psi: cl\ U \to X^2$  such that  $u=\Psi(\tilde{u}),\ \tilde{u}\in cl\ U$  is a solution of (5.1) with  $E(u)\leq \eta_0$ . For  $\tilde{u}\in U,\ \Psi(\tilde{u})$  is periodic, for  $\tilde{u}\in\partial U,\ \Psi(\tilde{u})$  is singular.

Proof: According to Theorem 6.5 and the following remark a component  $\Gamma_{\lambda} \subset H^2_{loc}(\Omega)$  of periodic solutions u exist in each case, satisfying  $E(u) < \eta_o$ . Their orbits in the  $(u_1,u_1')$ -plane form an open neighborhood of 0 in  $\mathbb{R}^2$ . Since  $u^{\sigma}$  is also a solution for all  $\sigma \in \mathbb{R}$  we obtain solutions u of (5.1) satisfying  $u_1(0) = a$ ,  $u_1'(0) = b$ ,  $E(u) < \eta_o$  for arbitrary  $(a,b) \in \mathbb{U}$ .

Now, consider the case (a,b)  $\in$  3U. As in the proof of Theorem 6.5 we conclude, for every sequence  $(a_n,b_n)\in U$  converging toward (a,b), that the corresponding periodic solutions  $u_n$  converge in  $H^2_{loc}(\Omega)$  towards a solution u of (5.1) which, by Corollary 4.5, is uniquely determined, and which

satisfies  $u_1(0) = a, u_1'(0) = b$ . Hence, the mapping  $(a,b) \mapsto u$  from  $\mathbb{R}^2$  into  $H^2_{loc}(\Omega)$  is injective and continuous. The orbit of u in the  $(u_1, u_1')$ -plane lies in  $\mathfrak{d}U$ ; hence u is singular, qed.

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