ON THE BEHAVIOR NEAR THE CREST
OF WAVES OF EXTREME FORM

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ABSTRACT

The angle \( \phi \) which the free boundary of an extreme wave makes with the horizontal is the solution of a singular, non-linear integral equation that does not fit (as far as we know) into the theory of compact operators on Banach spaces. It has been proved only recently that solutions exist and that (as Stokes suggested in 1880) these solutions represent waves with sharp crests of included angle \( \frac{2\pi}{3} \). In this paper we use the integral equation, known properties of solutions and the technique of the Mellin transform to obtain the asymptotic expansion

\[
\phi(s) = \frac{\pi}{6} + \sum_{n=1}^{k} a_n s^n + o(s^k) \quad \text{as} \quad s \to 0, \quad (*)
\]

to arbitrary order; the co-ordinate \( s \) is related to distance from the crest as measured by the velocity potential rather than by length. The first few (and probably all) of the exponents \( \mu_n \) are transcendental numbers. We are unable to evaluate the coefficients \( a_n \) explicitly, but define some in terms

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of global properties of \( \phi \), and the others in terms of earlier
coefficients. The derivation of (*) includes an assumption about a question
in number theory; if that assumption should be false, logarithmic terms
\( a_{n,j} s^n (\log s)^m n, j \), with \( m, j \) a positive integer, would enter the series at
very large values of \( n \).

Our results confirm the heuristic calculations of Grant [4] and Norman
[9], and, in effect, remove a tentative element from certain statements in
those papers.

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SIGNIFICANCE AND EXPLANATION

This paper concerns waves of permanent form on the free surface of an ideal liquid which is in two-dimensional, irrotational motion under the action of gravity. We consider only extreme waves, often called 'waves of greatest height'; each of these is the end-member of a one-parameter family of waves, and is distinguished from other 'smaller' members of the family by a sharp crest, as in Figure 1(a). Although this corner is physically unrealistic, oceanographers have given such idealized, extreme waves a great deal of attention since Stokes postulated their existence in 1880. (One reason may be the physical importance of the smaller waves, and that scientists like to interpolate.)

The present paper is a contribution to the strict mathematical theory of extreme waves, which has emerged only since 1978. We derive rigorously an asymptotic series that describes the flow near the crest. This confirms and sharpens certain earlier exploratory results due to Grant and Norman. The series should play a useful part in numerical computation of extreme waves.

The responsibility for the wording and views expressed in this descriptive summary lies with MRC, and not with the authors of this report.
ON THE BEHAVIOR NEAR THE CREST OF WAVES OF EXTREME FORM

C. J. Amick¹ and L. E. Fraenkel²

1. Introduction

This paper concerns gravity waves, of permanent and extreme form, on the free surface of an ideal liquid, the flow being two-dimensional, irrotational and in a vertical plane. By a wave of extreme form we mean one that is the 'largest' member of a one-parameter family of such waves, and is characterized by a sharp crest of included angle $2\pi/3$, as shown in Figure 1(a). The existence of such waves was conjectured by Stokes in 1880 and has been proved recently. (For a fuller account, see the Introduction to [1]; for background material, see [8].)

Letting $\theta$ denote the local wave angle (that is, the angle which the free boundary makes with the horizontal; $\theta(x) = \tan^{-1} v(x)$ in the notation of Figure 1), we shall seek the behavior as $\theta \to 0$ of a solution $\theta$ of the equation

$$\theta(s) = \frac{1}{3} \int_0^s K(s,t) \frac{v(t) \sin \theta(t)}{\int_0^t v \sin \theta} \, dt, \quad 0 < s < \pi,$$

where

$$K(s,t) = \frac{1}{\pi} \log \left| \tan \frac{1}{2} s + \tan \frac{1}{2} t \right|$$

and

$$v(t) = \frac{1}{2} (\cos^2 \frac{1}{2} t + b \sin^2 \frac{1}{2} t)^{-1/2}$$

for some $b \in [0,1]$ which we regard as fixed henceforth. Here $\int_0^t v \sin \theta$ stands for $\int_0^t v(u) \sin \theta(u) \, du$; such abbreviations will be used throughout the paper. Equation (1.1) refers only to waves of extreme form; the end-value of a parameter has already been chosen.

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The significance of the weighting function \( v \) is that \( v(s) = -C \partial \phi / \partial s \), where \( C \) is a positive constant and \( \phi \) denotes the velocity potential; the constant \( b \) distinguishes the families mentioned earlier: \( b = 0 \) corresponds to solitary waves (of infinite wavelength and in liquid of finite depth), \( b \in (0,1) \) to periodic waves in liquid of finite depth, and \( b = 1 \) to periodic waves in liquid of infinite depth.

By a solution of (1.1) we mean a function \( \phi \) satisfying the equation pointwise and such that

\[
0 < \phi(s) < \frac{\pi}{2} \quad \text{on} \quad (0, \pi), \quad \phi(\pi) = 0, \quad \phi \text{ is real-analytic}
\]

on \( (0, \pi) \), and \( \phi(s) + \frac{\pi}{6} \) as \( s \to 0 \).

Such solutions are now known to exist ([1], [2], [3], [7], [9]).

There have been many proposals for calculating the shape of the free boundary of an extreme wave by a combination of analytical and numerical methods (long lists of references are given in [5] and [14]). Such calculations were given a new direction by Grant [4], who sought the second term of an asymptotic expansion of \( z(\chi) \) for \( \chi \to 0 \) (here \( z = x + iy \) denotes the complex co-ordinate in the physical plane, and \( \chi = \phi + iy \) the complex potential), the first term being given by Stokes's corner flow. Grant pointed out that the second term must have an exponent that is irrational and 'probably transcendental'; he concluded that 'the structure near the corner is considerably more complicated than has been assumed in the past'. Norman [9] contemplated terms beyond the two considered by Grant; inferred the nature of all the exponents; introduced the assumption that the numbers \( \beta_j \), defined after (1.3) below, are linearly independent over the rationals; and established certain relationships between the coefficients of the series. (However, it seems that no coefficient after the first can be calculated by a merely local analysis.)

Longuet-Higgins and Fox [5], [6] used the exponents and functions arising in the work of Grant and Norman as one part of their theory of waves near an extreme one (the crest being smooth but of very large curvature).

All these calculations were heuristic rather than rigorous; in particular, Grant [4] and Norman [9] were concerned only to find analytic functions \( z(\chi) \) that satisfy the non-
linear boundary condition of the problem in an asymptotic sense as \( \theta \to 0 \). In this paper we proceed from (1.1) and (1.2) to prove that, subject to an assumption about a question in number theory to be explained presently,

\[
\psi(s) \sim \frac{\pi}{6} + \sum_{n=1}^{\infty} \frac{\mu_n}{s^n} \quad \text{as} \quad s \to 0 .
\]  

(1.3)

The exponents \( \mu_n \) depend on the positive zeros \( \beta_1, \beta_2, \beta_3, \ldots \), with \( \beta_j < \beta_{j+1} \), of \( \sqrt{3}(1+8) - \tan(8\pi/2) \); the \( \beta_j \) are essentially the numbers discovered by Grant [4], and \( \beta_1 \approx 0.8027, \beta_2 \approx 2.9066 \). As Grant thought probable, and as we prove in the Appendix, each \( \beta_j \) is a transcendental number. Each exponent \( \mu_n \) is a finite linear combination, with positive integer coefficients, of numbers in the set \( \{2, \beta_1, \beta_2, \ldots\} \), and contains at least one \( \beta_j \). The \( \mu_n \) are ordered by \( \mu_n < \mu_{n+1} \) (as the symbol \( \sim \) of asymptoticity in (1.3) implies); the first few are \( \beta_1, 2\beta_1, 3\beta_1, \beta_1 + 2, \beta_2, 4\beta_1, 2\beta_1 + 2, \beta_1 + \beta_2, 5\beta_1, 3\beta_1 + 2, 2\beta_1 + \beta_2, \ldots \).

The assumption made in the derivation of (1.3) is stated precisely in Section 4; here we remark that it certainly holds if the set \( \{1, \beta_1, \beta_2, \ldots\} \) is linearly independent over the rationals. Moreover, numerical calculation indicates that the assumption is true for the first hundred of the slightly larger set of exponents arising in the derivation of (1.3). If the assumption should be false, then our method would still be applicable, but logarithmic terms \( a_n s \log s \) with \( a_n, j \) a positive integer, would enter the series at large values of \( n \).

We cannot evaluate the coefficients \( a_n \) in (1.3) when \( \mu_n \notin \{\beta_1, \beta_2, \beta_3, \ldots\} \); we define such coefficients by integrals involving the global behavior of \( \psi \). When \( \mu_n \notin \{\beta_1, \beta_2, \beta_3, \ldots\} \), the corresponding coefficient \( a_n \) is determined by the previous coefficients \( a_1, \ldots, a_{n-1} \); this agrees with Norman’s results in [9].

The expansion (1.3) can be transformed and integrated to yield

\[
z(x) \sim -ig^{-1/3} \left( \frac{3}{2} ix \right)^{2/3} \left( 1 + \sum_{n=1}^{\infty} b_n (ix)^{\mu_n} \right) \quad \text{as} \quad x \to 0 ,
\]  

(1.4)

where the constant \( g \) is the gravitational acceleration, \( \arg(ix) \in [-\pi/2, \pi/2] \), and the coefficients \( b_n \) are real. If the set \( \{1, \beta_1, \beta_2, \ldots\} \) is linearly independent over the rationals, then \( b_n = 0 \) whenever the linear combination defining \( \mu_n \) contains a multiple
of 2. The expansion (1.4) is then of the form proposed by Norman [9], if we interpret liberally certain tentative remarks in that paper (for example, that it is 'possible to consider solutions ... corresponding to combinations of terms from several roots' of \( \sqrt{3}(1+\delta) = \tan(\delta\pi/2) \)).

The plan of the paper is as follows. We begin Section 2 by making the transformation 
\[ \xi = \tan \frac{1}{2}s, \quad \psi(\xi) = \phi(2 \tan^{-1} \xi) \]
in order to obtain a kernel \( k(\xi, \eta) \) that is simpler than \( K(s, t) \). (In effect, we map the unit disk in the plane of \( \zeta = \phi(e^{i\xi}) \) onto the half-plane \( w = \sigma + i\xi : \sigma \geq 0 \) by the conformal transformation \( w = (-1)/(\xi + 1) \).) Since \( \psi(s) = \psi(\tan \frac{1}{2}s) \), an asymptotic expansion of \( \psi(\xi) \) for \( \xi \to 0 \) yields one of \( \psi(s) \) for \( s \to 0 \). The next step is less obvious: we cast the integral equation (2.1) for \( \psi(\xi) \) into a form, (2.10), that contains an elaborate non-linearity but has the virtue of allowing us to construct the expansion of \( \psi(\xi) \) by an inductive process. To begin this process, we show that \( \psi(\xi) - y/6 = O(\xi^9) \) for some \( \alpha \in (\frac{1}{2}, \frac{3}{4}) \).

In Section 3 we combine this preliminary estimate with the use of the Mellin transform to show first that
\[ \psi(\xi) = \frac{\xi}{6} + A_1 \xi^1 + O(\xi^5) \quad \text{as} \quad \xi \to 0 \]  
(1.5)
for any \( \alpha \in (\frac{1}{2}, 1) \), and then that this result implies the improved approximation
\[ \psi(\xi) = \frac{\xi}{6} + A_1 \xi^1 + A_2 \xi^2 + O(\xi^8) \quad \text{as} \quad \xi \to 0 \]  
(1.6)
for any \( m \in (\frac{1}{2}, \frac{1}{2} + 1) \); here \( A_2 \) is a known function of \( A_1 \).

The step from (1.5) to (1.6) points the way to the long inductive proof, in Section 4, of the main results of the paper, which appear in Theorem 4.5 and Corollaries 4.6 and 4.7. The Appendix concerns some properties of the numbers \( B_j \).
2. Preliminary transformations and estimates

Under the transformation \( \xi = \tan \frac{1}{2} s, \eta = \tan \frac{1}{2} t \) and \( \phi(\xi) = \phi(2\tan^{-1}\xi) = \phi(s) \), equation (1.1) becomes

\[
\psi(\xi) = \frac{1}{3} \int_0^\infty k(\xi,\eta) w(\eta) \sin \frac{\eta}{\xi} \, d\eta, \quad 0 < \xi < \infty, \tag{2.1}
\]

where

\[
k(\xi,\eta) = \frac{1}{\pi} \log \frac{\xi + \eta}{|\xi - \eta|} \quad \text{and} \quad w(\eta) = (1 + \eta^2)^{-1/2} (1 + b\eta^2)^{-1/2}.
\]

By a solution of (2.1) we mean a function \( \psi \) satisfying (2.1) pointwise and such that

\[
0 < \psi(\xi) < \frac{\pi}{2} \quad \text{on} \quad (0,\infty), \quad \psi(\xi) = O(\xi^{-1}) \quad \text{as} \quad \xi \to \infty, \quad \psi \text{ is real-analytic}
\]

\[
\text{on} \quad (0,\infty), \quad \text{and} \quad \psi(\xi) + \frac{\pi}{6} \quad \text{as} \quad \xi \to 0.
\]

It is to be understood henceforth that \( \xi \in (0,\infty) \). Occasionally we set \( \xi = 0 \), with the implication that \( \psi(0) = \pi/6 \) and \( \psi \in C[0,\infty) \) (even though in the original problem \( \psi \) is an odd function on \( \mathbb{R} \setminus \{0\} \)).

Combining (2.1) and the formula ([1], p. 197)

\[
\frac{1}{3} \int_0^\infty k(\xi,\eta) \frac{1}{\eta} \, d\eta = \frac{\pi}{6},
\]

we obtain

\[
\psi(\xi) = \frac{\pi}{6} = \frac{1}{3} \int_0^\infty k(\xi,\eta) \frac{d}{d\eta} \log \left( \frac{1}{\eta} \int_0^\eta w \sin \psi \, d\eta \right),
\]

for any constant \( C > 0 \). Define

\[
\gamma(\xi) = \psi(\xi) - \frac{\pi}{6}, \quad (\gamma Y)(\eta) = 2w(\eta) \sin \left( \frac{\eta}{6} + \gamma(\eta) \right), \tag{2.3}
\]

then

\[
\gamma(\xi) = \frac{1}{3} \int_0^\infty k(\xi,\eta) \frac{d}{d\eta} \log \left( \frac{1}{\eta} \int_0^\eta EY \, d\eta \right). \tag{2.4}
\]

Our next transformation of the equation is more elaborate; it involves the integrated kernel (used extensively in [1])

\[
q(\xi,\eta) = \int_0^\xi k(t,\eta) \, dt = \frac{1}{\pi} \left\{ \xi \log \frac{\xi + \eta}{|\xi - \eta|} + \eta \log \frac{|\xi - \eta|^2}{\eta^2} \right\},
\]

with

\[
q_\eta(\xi,\eta) = \frac{1}{\pi} \log \frac{|\xi - \eta|^2}{\eta^2},
\]

and the non-linear operator \( F \) defined by
\[
(FY)(n) = \log \left( \frac{1}{n} \int_0^n EY \right) + \frac{n(EY)(n)}{\int_0^n EY} - 1 \quad (2.5a)
\]
\[
= \frac{d}{dn} \left[ n \log \left( \frac{1}{n} \int_0^n EY \right) \right] . \quad (2.5b)
\]

The function \( FY \) is important throughout the paper; we note its behavior for \( n \to 0 \) and for \( n \to \pm \infty \). By the definition (2.3),
\[
(FY)(n) = 1 + \sqrt{3} Y(n) + O(Y(n)^2 + n^2) \quad \text{as} \quad n \to 0 , \quad (2.6)
\]
and it follows from either form of (2.5) that
\[
(FY)(n) = \sqrt{3} Y(n) + O(Y(n)^2 + n^2) \quad \text{as} \quad n \to 0 . \quad (2.7)
\]

We claim that
\[
0 < \int_0^1 EY < \int_0^n EY < \text{const.}, \quad 1 < n < \infty , \quad (2.8)
\]
where the constant depends on \( Y \). The lower bound follows from (2.6), in which \( Y(n) \to 0 \) as \( n \to 0 \), and because \((FY)(x) > 0\) on \((0, \infty)\) by (2.2). For the upper bound we have
\[
\int_0^n EY < 2 \int_0^\infty \psi(n)(1+n^2 - \frac{1}{2}Y) \text{dn} = \int_0^\infty \psi(t) \sec \frac{t}{2} \text{dt} <
\]
by [3], p. 657, if \( b \in (0,1) \), and by [2], Theorem 4.7(a), for the case \( b = 0 \) of a solitary wave. Also, \( n(EY)(n) \) is bounded because \( w(n) = O(n^{-1}) \) as \( n \to \infty \); hence (2.5a) and (2.8) show that
\[
|(FY)(n)| < \text{const.} \log n, \quad 2 < n < \infty . \quad (2.9)
\]

**Lemma 2.1.** If \( w/6 + Y \) is a solution of (2.1), then
\[
Y(\xi) = \frac{1}{3} \int_0^\infty r(\xi, n)(FY)(n) \text{dn}, \quad 0 < \xi < \infty , \quad (2.10)
\]
where
\[
r(\xi, n) = -\frac{1}{\xi} q_\xi(\xi, n) = -\frac{1}{\xi^2} \log \frac{|\xi - \xi|^2}{n^2} . \quad (2.11)
\]

**Proof.** With the notation (2.3), equation (2.1) becomes
\[
\frac{w}{6} + Y(\xi) = \frac{1}{3} \int_0^\infty k(\xi, n) \frac{(FY)(n)}{\int_0^n EY} \text{dn} .
\]
Differentiation of this with respect to \( \xi \) is legitimate if the resulting integral is written as a Cauchy principal value; noting that \( \xi k_\xi(\xi, n) = -n k_\eta(\xi, n) \), and integrating by parts, we obtain
\[
\xi Y'(\xi) = \frac{1}{3} \int_0^\infty k(\xi, n) \frac{d}{dn} \left[ \frac{n(EY)(n)}{\int_0^n EY} \right] \text{dn} , \quad (2.12)
\]
where the integral exists by the properties (2.2) of $\psi \equiv \pi/6 + \gamma$ and because

$\psi'(n) = O(n^{-1})$ as $n \to 0$ (see [7]). We add (2.4) and (2.12), refer to the definition (2.5a) of $F_Y$, and integrate with respect to $\xi$; there results

$$
\xi_Y(\xi) = \frac{1}{3} \int_0^\infty q(\xi, n) \frac{d}{dn} (FY)(n) dn .
$$

We may integrate by parts because (for fixed $\xi \in (0, \infty)$, as elsewhere) $q(\xi, n)$ is $O(\log 1/n)$ as $n \to 0$, and is $O(n^{-1})$ as $n \to \infty$, and because we have the estimates (2.7) and (2.9) for $F_Y$. Accordingly,

$$
\xi_Y(\xi) = -\frac{1}{3} \int_0^\infty q_1(\xi, n)(FY)(n) dn ,
$$

and this is (2.10).

**Lemma 2.2.** If $\pi/6 + \gamma$ is a solution of (2.1), then there exists an exponent $a \in (\frac{1}{2}, \frac{3}{4})$ such that $\gamma(\xi) = O(\xi^a)$ as $\xi \to 0$.

**Proof.** In this proof we abbreviate $(FY)(n)$ to $F_Y(n)$, and similarly for other functions of the same kind.

(i) For any $c \in (0, 1)$, let $\xi \in (0, c]$ and rewrite (2.10) as

$$
\gamma(\xi) = \frac{1}{3} \int_0^c r(\xi, n)F_Y(n) dn + R_Y(\xi, c), \quad 0 < \xi < c \ ,
$$

where, in view of (2.7) and (2.9),

$$
|R_Y(\xi, c)| = \left| \frac{1}{3} \int_c^\infty r(\xi, n)F_Y(n) dn \right|
$$

$$
< \text{const.} \left\{ \int_c^\infty \frac{1}{\xi^2} \left[ \log \frac{(\pi^2 - n^2)}{n^2} \right] dn + \int_c^\infty \frac{\xi - \log n}{n^2} dn \right\}
$$

$$
< \text{const.} \left\{ \int_c^\infty \frac{u^2}{u^2 - 1} du + \xi \int_2^\infty \frac{\log n}{n^2} dn \right\}
$$

$$
< \text{const.} \frac{\xi}{c} \ ,
$$

where the constant depends on $\gamma$ but is independent of $c$. It is natural to define

$$
G_\gamma(n) = 2 \sin \frac{\pi}{6} \gamma(n) = 1 ;
$$

then $G_\gamma(n) \sim \sqrt{3} \gamma(n)$ as $n \to 0$ and $\gamma(n) \to 0$, and, by the definition (2.3),

$$
E_Y(n) = w(n)[1 + G_\gamma(n)] = 1 + G_\gamma(n) + O(n^2) \quad \text{as} \quad n \to 0 .
$$

Moreover, we can so define $U_Y$, $V_Y$ and $W_Y$ that

-7-
\[
\frac{1}{n} \log \int_0^n G_y = \log \left( 1 + \frac{1}{n} \int_0^n G_y + O(n^2) \right) = 1 - \frac{1}{n} \int_0^n G_y + O(n^2),
\]

\[
G_y(n) = \cos \gamma(n) - 1 + \sqrt{3} \sin \gamma(n) = \sqrt{3} \sin \gamma(n) - 1;
\]

then

\[
G_y(n) + 1, \ U_y(n) + 1, \ V_y(n) + 1 \quad \text{and} \quad W_y(n) + 1 \quad \text{as} \quad n \to 0,
\]

and

\[
F_y(n) = \{ V_y(n) - U_y(n) - G_y(n)U_y(n) \} \frac{1}{n} \int_0^n G_y + G_y(n) + \delta_y(n)
\]

\[
= \sqrt{3} \{ V_y(n) - U_y(n) - G_y(n)U_y(n) \} \frac{1}{n} \int_0^n W_y + \sqrt{3} W_y(n)\gamma(n) + \delta_y(n),
\]

where \( \delta_y(n) = O(n^2) \) as \( n \to 0 \).

Finally, define a linear operator \( L \) on the space \( C[0,c] \) by

\[
(Lf)(\xi) = \frac{1}{\sqrt{3}} \int_0^c x(\xi, n) \left[ (V_y(n) - U_y(n) - G_y(n)U_y(n)) \frac{1}{n} \int_0^n G_y + G_y(n) + \delta_y(n) \right] + \frac{3}{\sqrt{3}} \int_0^c W_y(n) f(n) \, dn,
\]

and the integral equation (2.13) becomes

\[
\gamma(\xi) = (Ly)(\xi) + S_y(\xi, c), \quad 0 < \xi < c,
\]

where

\[
S_y(\xi, c) = F_y(\xi, c) + \frac{1}{3} \int_0^c x(\xi, n)\delta_y(n) \, dn = O(\xi/c),
\]

because

\[
\left| \int_0^c x(\xi, n)\delta_y(n) \, dn \right| \leq \text{const.} \left| \int_0^c \frac{1}{u^2} \left| \log \frac{1-u}{u} \right| u^2 \, du \right| < \text{const.} \cdot c \xi.
\]

(ii) For every \( \varepsilon \in (0, \frac{1}{4}] \), choose \( c = c(\varepsilon) \) so small that

\[
|V_y(n) - U_y(n) - G_y(n)U_y(n)| + |1 - W_y(n)| < \varepsilon \quad \text{for all} \quad n \in [0, c(\varepsilon)],
\]

this is possible by (2.14). For any \( a \in [0,1) \) and any \( \varepsilon \in (0, \frac{1}{4}] \), define the Banach space

\[
X_a = X_a(c(\varepsilon)) = \{ f \in C[0, c(\varepsilon)] : \| f \|_a < \infty \},
\]

where

\[
\| f \|_a = \sup_{0 < \xi < c(\varepsilon)} \left| \xi^{-a} |f(\xi)| \right|.
\]

Consider the linear operator \( L \) defined by (2.15), we wish to show that \( L \) maps \( X_a \) into itself, for \( \varepsilon \) sufficiently small and for some \( \alpha > \frac{1}{2} \), and that
It is easy to show that $L_\alpha \in C(0, c(\epsilon))$ when $f \in X_\alpha$ and we estimate $\|L_\alpha f\|_\alpha$ as follows.

\[
| (L_\alpha) (\xi) | < \frac{1}{\sqrt{3}} c(\xi, n) | x(\xi, n) | \left( \frac{c(\xi, n)}{n} \right) f + (1+\epsilon) | f | \frac{\partial}{\partial n} d n \\
< \frac{1}{\sqrt{3}} (1+\epsilon) \|L_\alpha f\|_\alpha c(\xi, n) | x(\xi, n) | \left( \frac{1+\epsilon}{1+\alpha} n^\alpha + n \right) \frac{\partial}{\partial n} d n \\
< \frac{1}{\sqrt{3}} (1+\epsilon) (1+\frac{\epsilon}{1+\alpha}) \|L_\alpha f\|_\alpha \int_0^\infty \left| \log \frac{t^2-n^2}{n^2} \right| n^\alpha d n \\
= \frac{1}{\sqrt{3}} (1+\epsilon) (1+\frac{\epsilon}{1+\alpha}) \|L_\alpha f\|_\alpha \int_0^\infty \left| \log \frac{1-u^2}{u^2} \right| u^\alpha d u ,
\]
where $0 < \xi < c(\epsilon)$. It follows that

\[
\|L_\alpha f\|_\alpha < (1+\epsilon)(1+\frac{\epsilon}{1+\alpha}) u(\alpha) ,
\]

where

\[
u(\alpha) = \frac{1}{\sqrt{3}} \int_0^\infty \left| \log \frac{1-u^2}{u^2} \right| u^\alpha d u \\
= \frac{1}{\sqrt{3}} \left( \int_0^\infty \log \frac{u^2}{1-u^2} u^\alpha d u + 2 \int_0^\infty \log \frac{1-u^2}{u^2} u^\alpha d u \right) .
\]

The first of these two integrals is evaluated in Section 3, see (3.11a); in the second, we set $u = 1/x$ and integrate by parts; then

\[
u(\alpha) = \frac{1}{\sqrt{3} (1+\alpha)} \tan \frac{\pi}{2} + \frac{4}{1+\alpha} \int_0^\infty \frac{dx}{\sqrt{2} x^{1+\alpha}}.
\] (2.18)

To obtain a simple majorant, replace $x^\alpha$ by 1 in the last integral, which is then easily evaluated; thus

\[
u(\alpha) < \frac{1}{\sqrt{3} (1+\alpha)} \left( \frac{\pi}{2} + \frac{4}{1+\alpha} \log (1 + \sqrt{2}) \right) = \tau(\alpha), 
\]
say.

Now $\tau(\frac{1}{2}) < 0.82$, $\tau(\frac{3}{4}) > 1.16$ and $\tau'(\alpha) > 0$ on $(0,1)$; hence there exists a number $\alpha \in (\frac{1}{2}, \frac{3}{4})$ such that $\nu(\alpha) < \tau(\alpha) < \tau(\alpha) < 1$ for all $\alpha \in [0,\alpha]$. If we choose $\epsilon$ sufficiently small, then $\|L_\alpha f\|_\alpha < 1$ for all $\alpha \in [0,\alpha]$.

(iii) With $\epsilon$ and $c = c(\epsilon)$ now fixed, the estimate (2.16b) of $S_\alpha$ may be written $S_\alpha(\xi) = O(\xi)$ as $\xi \to 0$. Hence $S_\alpha \in X_\alpha$ for all $\alpha \in (0,1)$, and so the equation

\[f(\xi) = (L_\alpha)(\xi) + S_\alpha(\xi)\]
has a unique solution $f_\alpha \in X_\alpha$ for each $\alpha \in [0,\alpha]$. Since
the solutions are identical; in other words, \( f_\alpha = f_0 \) for all \( \alpha \in [0, \tilde{\alpha}] \) because each \( f_\alpha \in X_0 \). We know that \( Y \in X_0 \) and satisfies (2.16a), whence \( Y = f_0 \in X_\alpha \) for all \( \alpha \in [0, \tilde{\alpha}] \). In particular, \( Y \in X_{\tilde{\alpha}} \), so that \( \tilde{\alpha} \) may be chosen as the exponent in the statement of the lemma.

Remark. We shall see from Lemmas 3.1 and 3.2 that the best possible exponent is the number \( \beta_1 \) introduced after (1.3), so that

\[
\frac{1}{\sqrt{3(1+\beta_1^2)}} \tan \frac{\beta_1^2}{2} = 1.
\]

(\( \beta_1 \) is best possible if \( A_1 \neq 0 \) in Lemma 3.2.) This corresponds to only the first term of \( \mu(\alpha) \) in (2.18), so that precise treatment of the second term is not worth while.
3. The asymptotic expansion of $\psi(\xi)$ to three terms

Our observation in (2.7) that $(FY)(n) \sim \sqrt{3} \gamma(n)$ as $n \to 0$ suggests that the integral equation (2.10) be written

$$
\gamma(\xi) = \frac{1}{\sqrt{3}} \int_0^1 x(\xi,n)\gamma(n)dn + \rho_\gamma(\xi), \quad 0 < \xi < 1,
$$

where

$$
\rho_\gamma(\xi) = \frac{1}{3} \int_0^1 x(\xi,n)((FY)(n) - \sqrt{3} \gamma(n))dn + \frac{1}{3} \int_1^\infty x(\xi,n)(FY)(n)dn.
$$

From (2.7) and Lemma 2.2 we have

$$(FY)(n) - \sqrt{3} \gamma(n) = O(n^{-\alpha}) \quad \text{as} \quad n \to 0, \quad \text{with} \quad \alpha \in \left(\frac{1}{2}, \frac{3}{4}\right),$$

and it follows (by an estimate like (2.17) for the first integral in (3.2), and an estimate like (2.13b) for the second) that $\rho_\gamma(\xi) = O(1)$ as $\xi \to 0$.

All the results to come will be consequences of (3.1) and (3.2). Our plan is to bootstrap from some expansion of $\psi(\xi)$ for $\xi \to 0$ (at present, $\psi/6 + O(\xi^2)$) to an estimate of $\rho_\gamma(\xi)$ (at present, $O(\xi)$), and then to derive a more complete expansion of $\psi(\xi)$ from (3.1), regarded as a linear equation for $\gamma$ in which $\rho_\gamma$ is 'known'. The proof by induction in Section 4 will follow this program to exhaustion; without the more explicit first steps in this section, it would probably be incomprehensible.

The Mellin transform ([12], p. 7) will be our main tool for analysis of (3.1). Let $f$ be piecewise continuous on $(0,\infty)$, define $\phi(x,s) = x^{-1}f(x)$ for $x > 0$, where $s = \sigma + it \in C$, and assume that $\phi(x,s) \in L^1_1(0,\infty)$ for $a < \sigma < b$. Then, for $s$ in this strip, we define the Mellin transform $\hat{f}$ of $f$ by

$$
\hat{f}(s) = \int_0^\infty x^{-1}f(x)dx,
$$

we also write this as $\hat{f}(s) = \hat{f}(x)$. Note that $\hat{f}$ is analytic in the strip $\text{Re} \ s \in (a,b)$. The inversion formula ([12], p. 46) is

$$
f(x) = \frac{1}{2\pi i} \int_{c-i\infty}^{c+i\infty} x^{-s} \hat{f}(s)ds, \quad a < c < b,
$$

at points of continuity of $f$. The product formula ([12], p. 54) is

$$
\hat{f}(-s+A)g(s+B) := \hat{B}^s \int_0^\infty u^{A-1}f(u)g(0)du,
$$

provided that $\hat{f}(-s+A)$ and $g(s+B)$ have a common strip of convergence. Finally, we record the following property ([12], p. 118) of the Mellin transform.
If \( \int_0^a |x|^p f(x) \frac{dx}{x} < \) for some \( a \in (a,b) \) and
\[ \text{some } p \in (1,2), \text{ then } f(0+i) \in L_p((-\infty,\infty)), \text{ where } \frac{1}{p} + \frac{1}{p'} = 1. \] (3.6)

To evaluate certain contour integrals arising from (3.4), and to state our theorems, we shall need the following lemma, a variant of which has already been given in [5] for a different purpose. The significance of the exponents \( \beta_j \) in the lemma is evident from the observation (a particular case of Lemma 4.2) that, if \( \sqrt{5}(1+\beta) = \tan(8\pi/2) \) and \( \beta > -1, \)
then
\[ \frac{1}{\sqrt{5}} \int_0^1 r(\xi,\eta) \eta^\beta d\eta = \xi^\beta + a(\xi), \quad 0 < \xi < 1, \] (3.7)
where \( a(\xi) \) is \( O(\xi) \) and real-analytic on \([0,1)\), and contains only odd powers of \( \xi \) in its Taylor series about the origin. In other words, the linear integral operator in (3.1) leaves the functions \( \xi^{-\beta} \) almost invariant, merely adding to them such functions \( a \), and these latter turn out to be unimportant.

**Lemma 3.1.** The only zeros of \( \sqrt{5}(1-s) + \tan(s\pi/2) \) in the half-plane \( \Re s < 1 \) are simple zeros on the negative real axis. We denote such points by \( s = -\beta_j, \ j = 1,2,3, \ldots, \)
with \( \beta_j < \beta_{j+1} \). Then \( \beta_j \in (2j-2,2j-1) \) for all \( j \), and \( \beta_j = 2j-1 + O(j^{-1}) \) as \( j \to \infty \) (the \( 0 \)-term is negative). Also,
\[ \beta_1 \approx 0.8027, \quad \beta_2 \approx 2.9066, \quad \beta_3 \approx 4.9383. \]

**Proof.** Set \( s = 1 - 2z/\pi \) and \( z = x+iy \); we have to solve
\[ \frac{2\sqrt{3}}{\pi} z + \cot z = 0, \quad x > 0, \]
or equivalently, since no zero or pole of \( \tan z \) is a solution,
\[ z \tan z = \frac{-\pi}{2\sqrt{3}}, \quad x > 0. \]

If this equation has a solution with \( y \neq 0 \), then the imaginary part of the equation gives
\[ \frac{\sinh 2y}{2y} = -\frac{\sin 2x}{2x}, \]
which is impossible (because \( \sinh 2y/2y > 1 \), while \( -\sin 2x/2x < 0 \) for \( 0 < x < \pi/2 \)
and \( -\sin 2x/2x < 1/2 \) for \( x > \pi/2 \)). The remaining assertions now follow from elementary analysis of the points where the graphs of \( \tan x \) and \( \cot \) intersect, for \( x > 0 \).
LEMMA 3.2. If \( \phi \) is a solution of (2.1), then
\[
\psi(\xi) = \frac{1}{6} + A_1 \xi + o(\xi^2) \quad \text{as} \quad \xi \to 0,
\]
for some constant \( A_1 \), and any \( \xi \in (\beta, 1) \). Here \( \beta_1 \) is as in Lemma 3.1.

Proof. (i) Define
\[
g(\xi) = \begin{cases} 
\gamma(\xi), & 0 < \xi < 1, \\
0, & 1 < \xi,
\end{cases}
\]
and
\[
h(\xi) = \begin{cases} 
p_\gamma(\xi), & 0 < \xi < 1, \\
\frac{1}{\sqrt{3}} \int_0^1 r(\xi, n) \gamma(n) \, dn, & 1 < \xi,
\end{cases}
\]
so that (3.1) becomes
\[
g(\xi) = \frac{1}{\sqrt{3}} \int_0^1 r(\xi, n) \gamma(n) \, dn + h(\xi), \quad 0 < \xi < 1,
\]
since the integrand is zero for \( n > 1 \), and both sides are zero for \( \xi > 1 \). Setting \( n = \xi u \), we arrive at a form suitable for the Mellin transform:
\[
g(\xi) = \frac{1}{\sqrt{3}} \int_0^\infty R(u) r(\xi u) \, du + h(\xi), \quad 0 < \xi < 1,
\]
where
\[
R(u) = \frac{\gamma(1, u) - \gamma(u)}{u^2}.
\]
Because of the plethora of notation, we emphasize that \( g \) is a truncated form of \( \gamma = \phi - \psi/6 \), that \( h(\xi) \) is presumably of smaller order than \( g(\xi) \) as \( \xi \to 0 \), and that the dependence of \( h \) on \( \gamma \) is now implicit.

(ii) The term \( h(\xi) \) may be estimated as follows. Let \( \alpha \in \left( \frac{1}{2}, \frac{3}{2} \right) \) be the exponent in Lemma 2.2; then for \( \xi > 2 \) (and hence \( u < \frac{1}{2} \) when we set \( n = \xi u \))
\[
|h(\xi)| < \frac{1}{\sqrt{3}} \int_0^1 |r(\xi, n)| |\gamma(n)| \, dn
\]
\[
< \text{const.} \xi^{\alpha} \int_0^{1/\xi} |\log \left( \frac{1-u^2}{u^2} \right)| u^\alpha \, du
\]
\[
< \text{const.} \xi^{\alpha} \int_0^{1/\xi} \log \left( \frac{1-u^2}{u^2} \right) u^\alpha \, du
\]
\[
< \text{const.} \xi^{-1} \log \xi,
\]
-13-
and for $1 < \xi < 2$ the integral is bounded by a constant. Since $\rho_y(\xi) = O(\xi)$ as $\xi \to 0$, we have
\[ |h(\xi)| < \begin{cases} \text{const. } \xi, & 0 < \xi < 1, \\ \text{const. } \xi^{-1}(1 + \log \xi), & 1 < \xi. \end{cases} \] (3.9)

It follows that $h(s)$ exists and is analytic for $-1 < \sigma < 1$, where $s = \sigma + it$, and from (3.6) that
\[ h(\sigma+i.) \in L_{p'}(\mathbb{R}) \text{ for all } \sigma \in (-1,1) \text{ and all } p' > 2. \] (3.10)

(iii) To find the Mellin transform of the function $R$ in (3.8), we recall that our original kernel in (2.1) was
\[ k(\xi, \eta) = \frac{1}{\pi} \log \frac{\xi + \eta}{|\xi - \eta|}, \]
and note from [12], p. 192, that (with $s = \sigma + it$)
\[ \int_0^\infty u^{s-1} k(1, u) du = \frac{1}{s} \tan \frac{\pi s}{2}, \quad -1 < \sigma < 1. \]
Setting $u = n/\xi$ for fixed $\xi \in (0, \infty)$, multiplying both sides by $\xi^s$, and then integrating with respect to $\xi$, we obtain
\[ \int_0^\infty n^{s-1} q(\xi, n) d\xi = \frac{\xi^{s+1}}{s(s+1)} \tan \frac{\pi s}{2}, \quad -1 < \sigma < 1. \]
We may integrate by parts because $n^s q(\xi, n)$ is $O(n^{\sigma+1} \log 1/n)$ as $n \to 0$, and is $O(n^{\sigma-1})$ as $n \to \infty$, and there results
\[ -\int_0^\infty n^{s-1} q(\xi, n) d\xi = \frac{\xi^{s+1}}{s+1} \tan \frac{\pi s}{2}, \quad -1 < \sigma < 1. \]
Set $n = \xi u$ and $s = z-1$, then
\[ -\frac{1}{u} \int_0^\infty u^{s-1} \log \frac{1-u^2}{u^2} du = \frac{1}{z} \tan \frac{(z-1)\pi}{2}, \quad 0 < \Re z < 2, \] (3.11a)
that is,
\[ R(z) = \frac{1}{z} \tan \frac{(z-1)\pi}{2}, \quad 0 < \Re z < 2. \] (3.11b)

(iv) We are now in a position to take the Mellin transform of the integral equation (3.8a). By Lemma 2.2, $g(\xi) = O(\xi^\alpha)$ as $\xi \to 0$, with $\alpha \in \left(\frac{1}{2}, \frac{3}{4}\right)$; therefore $\hat{g}(s)$ exists and is analytic for $\sigma > -\alpha$. We have just shown that $\hat{h}(s)$ and $\hat{R}(-s+1)$ exist and are analytic for $-1 < \sigma < 1$. Hence, if $-\alpha < \sigma < 1$, we may apply the Mellin transform to (3.8a) and use the product formula (3.5) with $A = 1$ and $B = 0$:
\[ \hat{g}(s) = \frac{1}{\sqrt{3}} \hat{R}(1-s)\hat{g}(s) + \hat{h}(s), \quad -\alpha < \sigma < 1, \]
or, equivalently,

$$g(s) - h(s) = Q(s)\hat{h}(s), \quad -\alpha < \sigma < 1,$$

where

$$Q(s) = \frac{R(1-s)}{\sqrt{3-3R(1-s)}} - \frac{\tan(\pi s/2)}{\sqrt{3(1-s) + \tan(\pi s/2)}}. \quad (3.13)$$

The inversion formula (3.4) gives, for all \( \xi \in (0,\infty) \setminus \{1\} \),

$$g(\xi) - h(\xi) = \frac{1}{2\pi i} \int_{c-i\infty}^{c+i\infty} \xi^{-s} Q(s)\hat{h}(s) ds, \quad -\alpha < \sigma < 1. \quad (3.14)$$

(v) Let \( \xi \in (0,1) \), so that \( g(\xi) = \gamma(\xi) \). Since \( h(\xi) = O(\xi) \) as \( \xi \to 0 \), it suffices to show that the right-hand side of (3.14) behaves like \( A_1 \xi^{-1} + O(\xi^\beta) \) for some constant \( A_1 \) and any \( \xi \in (\beta_1,1) \). We prove this by moving the path of integration as far to the left as our knowledge of \( h \) allows, and taking the residue at any pole of the integrand between the path in (3.14) and the new one. The details are as follows.

Given \( c \in (-\alpha,1) \) and \( \xi \in (\beta_1,1) \), define \( \Gamma_M \) to be the (positively directed) rectangular contour with corners at \( c \pm iM \) and \( -c \pm iM \) for some large \( M > 0 \), and consider the contour integral

$$\frac{1}{2\pi i} \int_{\Gamma_M} \xi^{-s} Q(s)\hat{h}(s) ds.$$

The contribution of the horizontal parts \((-\xi < c < c, t = \pm\infty)\) of \( \Gamma_M \) tends to zero as \( M \to \infty \), because (a) \( |\xi^{-s}| = \xi^{-\sigma} \) with \( 0 < \sigma < 1 \); (b) the fact that \( \tan(\pi s/2) + i \) as \( t \to \pm\infty \) with \( \sigma \) fixed implies, in view of (3.13), that \( Q(s) = O(\xi^{-1}) \); (c) the Riemann-Lebesgue lemma, applied to the definition integral (3.3) of \( \hat{h}(s) \), shows that \( \hat{h}(s) \to 0 \) as \( t \to \pm\infty \) with \( \sigma \) fixed.

Since \( \hat{h}(s) \) is analytic for \(-1 < \sigma < 1\), we conclude from (3.13) and Lemma 3.1 that the only singularity of \( \xi^{-s} Q(s)\hat{h}(s) \), with \(-1 < \sigma < 1\), is a simple pole at \( s = -\beta_1 \).

Accordingly,

$$\frac{1}{2\pi i} \left[ \int_{c-i\infty}^{c+i\infty} - \int_{c+i\infty}^{c-i\infty} \xi^{-s} Q(s)\hat{h}(s) ds \right]$$

$$= \lim_{M \to \infty} \frac{1}{2\pi i} \int_{\Gamma_M} \xi^{-s} Q(s)\hat{h}(s) ds$$

$$= \xi^{-\beta_1} \hat{h}(s) \bigg|_{s = -\beta_1} \hat{h}(\beta_1) = A_1 \xi^{-1}, \quad \text{say}, \quad (3.15)$$

and so
We bound this last integral (without the factor $\xi^{1/2\pi}$) by means of Hölder's inequality, noting that $\mathcal{O}(-i+1.) \in L^2(\mathbb{R}^\infty)$ because $\mathcal{O}(s) = O(e^{-s})$ as $s \to \infty$, and that $h(-i+1.) \in L^2(\mathbb{R}^\infty)$ by (3.10). As we remarked earlier, the desired result now follows, because $g(\xi) = \gamma(\xi)$ for $\xi < 1$ and $h(\xi) = O(\xi)$ as $\xi \to 0$.

Remarks. 1. As was to be expected from (3.1) and (3.7), we cannot evaluate the coefficient $A_1$ in Lemma 3.2; in (3.15) we have defined it in terms of the global behavior of $\gamma$.

2. Stokes conjectured (see [1], pp. 194 and 199) that the profile of an extreme wave is convex between the crest and the trough; in other words, that $\psi(\xi)$ is non-increasing on $(0, \infty)$. Now Lemma 3.2 and analysis of the equation which results from differentiation of (2.1) imply that $\xi^2 \psi'(\xi) = A_1 \xi \beta_1 \xi^{-1} + O(\xi^0)$. Therefore we expect that $A_1 < 0$, and a calculation shows that this corresponds to $b_1 > 0$ in (1.4). Numerical evidence ([6], equation (2.10) with the value $B = 0.131$ on p. 776; [14], Table 1) suggests that $b_1 > 0$ and hence that $A_1 < 0$.

3. The term $A_2 \xi^{-1}$ in our next result has a character quite different from that of $A_1 \xi^{-1}$; instead of arising from (3.7), it arises from terms that are essentially squares of the first perturbation, and so $A_2$ is determined uniquely by $A_1$.

**Theorem 3.3.** If $\psi$ is a solution of (2.1), then
\[ \psi(\xi) = \frac{\beta_1}{\xi} + \frac{1}{\xi} + A_2 \xi^{-1} + O(\xi^{-1}) \quad \text{as} \quad \xi \to 0, \]
for any $m \in (2, 3, \beta_1, \xi')$. Here $A_1$, $\beta_1$ and $\xi$ are as in Lemma 3.2, so that $\xi \in (\beta_1, 1)$, and
\[ A_2 = -\frac{(4+8s_i^2)(\tan \beta_1 \pi)A_1^2}{2\sqrt{\gamma(1+8s_i^2)(\tan \beta_1 \pi) - \sqrt{3}(1+2s_i)}} \approx 0.1341 A_1^2. \]

Proof. We shall use the notation in the proof of Lemma 3.2, and shall begin by showing that the previous estimate (3.9) of $h$ can be sharpened, by means of the result of Lemma 3.2, to
\[ h(\xi) = C_1 \xi + C_2 \xi^{-1} + O(\xi^{-1}), \quad 0 < \xi < 1, \quad (3.17a) \]

\[ |h(\xi)| \leq \text{const.} \xi^{-1}(1 + \log \xi), \quad 1 < \xi, \quad (3.17b) \]

where \( C_1 \) is a constant of no interest, because ultimately it will drop out of the integral equation, and \( C_2 \) is a constant that is known in terms of \( A_1 \). With this estimate in hand, we can move the path of integration in (3.14) further to the left than was possible before. Taking due account of the poles of the integrand between the path in (3.14) and the new one, we shall obtain the result of the theorem.

(i) Let \( \xi \in (0,1] \). A straightforward calculation gives the following.

\[ (EY)(n) = 2w(n) \sin\left(\frac{\xi}{6} + \gamma(n)\right) \]

\[ = 1 + \sqrt{3} \gamma(n) - \frac{1}{2} A_1^2 n^{28} + O(n \log n) \]

\[ (FY)(n) = \sqrt{3} \gamma(n) = \log\left(\frac{1}{n} EY\right) + \frac{\partial(EY)(n)}{\partial EY} - 1 - \sqrt{3} \gamma(n) \]

\[ = \frac{4+8B_1^2}{2(1+8B_1^2)} A_1^2 n^{28} + O(n \log n) \]

Since \( 28_1 \) and \( B_1 + 1 \) are not odd positive integers, Lemmas 4.2 and 4.3 yield

\[ \frac{1}{3} \int_0^1 x(\xi,n)((FY)(n) - \sqrt{3} \gamma(n)) d\xi = C_0 \xi + C_2 \xi^{-1} + O(\xi^{-1}) \], \quad (3.18)

where the constant \( C_0 \) will be a part of \( C_1 \) (and is therefore of no interest), and

\[ C_2 = \frac{(4+8B_1^2)(\tan B_1)^2 A_1^2}{6(1+8B_1^2)^2} \], \quad (3.19)

The other term in the definition (3.2) of \( \rho_Y(n) \) is

\[ \frac{1}{3} \int_1^\infty x(\xi,n)(FY)(n) = -\frac{1}{3\pi} \int_1^\infty \log\left(1 - \frac{E^2}{n^2}\right)(FY)(n) d\xi \]

\[ = \frac{1}{5 \log \xi} \int_1^\xi \frac{2k-1}{k} \frac{(FY)(n)}{n^2k} d\xi \], \quad (3.20)
where the integral is $O(k^{-1})$ as $k \to \infty$, so that the series converges for $\xi \leq 1$.

Equations (3.18) and (3.20) imply that

$$
\rho_1(\xi) = C_1\xi + C_2\xi^{-1} + O(\xi^{-1})
$$

and thus prove (3.17a), while (3.17b) is a previous estimate.

(ii) Define

$$
h_1(\xi) = \begin{cases} 
C_1\xi + C_2\xi^{-1}, & 0 < \xi < 1, \\
0, & 1 < \xi ,
\end{cases}
$$

and $h_2 = h - h_1$. Then, for $\sigma > -1$,

$$
h_1(s) = \frac{C_1}{s+1} + \frac{C_2}{s+28},
$$

while $\hat{h}_2(s)$ exists and is analytic for $-\beta - \xi < \sigma < 1$, and $\hat{h}_2(\sigma+i.) \in L_p'(-\infty,\infty)$ for all $\sigma \in (-\beta - \xi, 1)$ and all $p' > 2$. The obvious analytic (more precisely, meromorphic) continuations of $\hat{h}_1$ and $\hat{h}$ into the strip of convergence of $\hat{h}_2$ will also be denoted by $\hat{h}_1$ and $\hat{h}$, respectively.

We can now proceed from (3.14) as in the proof of Lemma 3.2, step (v), except that now we move the path of integration to $\sigma = -m$ for any $m \in (28,\beta + \xi)$, and collect the residues of the poles of $\xi^{-\xi}Q(s)h(s)$ at $s = -\beta - \xi$ and $-28$. Since $Q(-1) = -1$, there results

$$
g(\xi) - h(\xi) = \xi^{-\xi} - C_1\xi + C_2\xi^{-1} + \int_{-m}^{28} \xi^{-\xi}Q(-m+i)\hat{h}(-m+i)dt ;
$$

the last integral is bounded as before since $\hat{h}_1(-m+i.)$ and $\hat{h}_2(-m+i.)$ are both in $L_2(-\infty,\infty)$. Substituting for $h(\xi)$ from (3.17a), we have

$$
\gamma(\xi) = \xi^{-\xi} + \xi^{-1} + C_1(28) + 1)\xi^{-1} + O(\xi^{-1})
$$

as $\xi \to 0$,

where $C_2$ is given by (3.19) and $Q(-28)$ by (3.13).
4. The asymptotic expansion of $\Psi(\xi)$ to any number of terms

4.1. Notation

Let $\mathbb{N}_0 = \{0, 1, 2, \ldots\}$, $\mathbb{N} = \{1, 2, 3, \ldots\}$ and $B = \{\beta_1, \beta_2, \beta_3, \ldots\}$, where the numbers $\beta_j$ are as in Lemma 3.1. We shall prove in the Appendix that the $\beta_j$ are transcendental numbers.

An elaborate notation is needed for the exponents in various series that we shall construct, because (a) the non-linearity in (2.1), and the presence of the weighting function $w$ there, cause the exponents in the expansion of $\Psi$ to form a fairly complicated sequence, (b) further exponents enter, and finally depart, during the course of the construction.

In what follows $(m_1, \ldots, m_{r+1})$ is a multi-index: $m_i \in \mathbb{N}_0$ for $i = 1, \ldots, r+1$ and $r$ is any element of $\mathbb{N}$. Define

$$A = \{m_1 \beta_1 + \ldots + m_r \beta_r + m_{r+1} : m_1 + \ldots + m_r + m_{r+1} \geq 1\},$$

since $A$ is a countable set, we may write $A = \{\lambda_1, \lambda_2, \lambda_3, \ldots\}$ with $\lambda_j < \lambda_{j+1}$. The first few numbers $\lambda_j$ are

$\beta_1$, $2\beta_1$, 2, $3\beta_1$, $\beta_1 + 2$, $\beta_2$, $4\beta_1$, $2\beta_1 + 2$, $\beta_1 + \beta_2$, 4, $5\beta_1$, $3\beta_1 + 2$, $2\beta_1 + \beta_2$, etc.

The following subsets of $A$ will be needed.

$$A_B = \{m_1 \beta_1 + \ldots + m_r \beta_r + m_{r+1} : m_1 + \ldots + m_r > 1\},$$

so that

$$A = A_B \cup \{2, 4, 6, \ldots\},$$

and

$$A_3 = \{2\} \cup \{m_1 \beta_1 + \ldots + m_r \beta_r + m_{r+1} : m_1 + \ldots + m_r + m_{r+1} > 2\},$$

so that

$$A = A_3 \cup B.$$

Thus each element of $A_B$ involves at least one of the $\beta_j$; we expect the exponents in the expansion of $\Psi$ to be in $A_B$. The exponents in $A_3$ arise if, for example, we combine three series as follows: form the product of two series with exponents in $A$ and add a series with exponents in $\{2, 4, 6, \ldots\}$.

Truncated subsets will be denoted by
\[ A_\beta(n) = \{ \lambda_j \in \Lambda_\beta : \lambda_j \leq \lambda_n \} = \{ \lambda_j \in \Lambda_\beta : j < n \} , \]
\[ A_3(n) = \{ \lambda_j \in \Lambda_3 : j < n \} , \]
\[ B(n) = \{ \beta_j \in B : \beta_j < \lambda_n \} . \]

4.2. Bootstrapping from \( \gamma \) to \( \rho_\gamma \)

We now prove that a given expansion of \( \gamma \), with exponents in \( A_\beta(n) \), implies an expansion of \( \rho_\gamma \) with exponents in \( A_3(N) \) \{1,3,5,... \}, where \( N > n+1 \). This will allow us to improve the given approximation to \( \gamma \) by the method that we used to pass from Lemma 3.2 to Theorem 3.3. We shall prove in Theorem 4.5 that, if \( A_3(N) \cap B(N) = \emptyset \), then \( \gamma \) has an expansion with exponents in the larger set \( A_\beta(N) \). Hence, if \( A_3 \cap B = \emptyset \) (which is a number-theoretic problem addressed in the Appendix), then \( \gamma \) has an infinite asymptotic expansion with exponents in \( A_\beta \).

We begin by expanding the function \( F\gamma = \sqrt{3} \gamma \) in the definition (3.2) of \( \rho_\gamma \); then we establish two basic properties of the linear integral operator in (3.1) and in the first term of \( \rho_\gamma \). The second term of \( \rho_\gamma \) has a very simple expansion. Combining these results, we obtain the expansion of \( \rho_\gamma \), for a given expansion of \( \gamma \).

**LEMMA 4.1.** Suppose that, for certain constants \( A_{\mu} \),
\[ \gamma(n) = \sum_{\lambda_j \in A_\beta(n)} B_j n^j + O(n^\alpha) \; \text{as} \; n \to \infty , \tag{4.1} \]
where \( \lambda_n < \mu < \lambda_{n+1} \) and \( \mu + \beta_1 \in \Lambda \) \{1,3,5,... \}. Define \( N = N(n,\mu) \) by
\[ \lambda_N = \max(\lambda_j \in \Lambda : \lambda_j < \mu + \beta_1) . \]
(Since \( \lambda_{n+1} < \lambda_n + \beta_1 < \mu + \beta_1 \), we have \( N(n,\mu) > n+1 \).) Then
\[ (F\gamma)(n) = \sqrt{3} \gamma(n) = \sum_{\lambda_j \in A_3(n)} B_j n^j + O(n^{\alpha + \beta_1}) \; \text{as} \; n \to \infty , \tag{4.2} \]
where the coefficients \( B_j \) depend only on the \( A_{\mu} \) in (4.1).

**Proof.** Again we abbreviate \( (F\gamma)(n) \) to \( F\gamma(n) \), and similarly for other functions of the same kind. We have
\[ E_\gamma(n) = 2w(n) \sin \left( \frac{\pi}{6} + \gamma(n) \right), \]

where
\[ w(n) = \left(1+n^2\right)^{-1/2} \left(1+2n^2\right)^{-1/2} = 1 + \sum_{k=1}^{\infty} C_k n^{2k}, \text{ say, for } n < 1. \]

Define two new operators \( H \) and \( J \) by
\[
H_\gamma(n) = E_\gamma(n) - 1 - \sqrt{3} \gamma(n)
\]
\[ = w(n)\{\cos \gamma(n) - 1 + \sqrt{3}(\sin \gamma(n) - \gamma(n))\} + (w(n) - 1)\{1 + \sqrt{3} \gamma(n)\}, \tag{4.3} \]
and
\[ J_\gamma(n) = \frac{1}{n} \int_0^n (E_\gamma - 1). \]

Let us now restrict attention to those \( n \), say \( n < n_0 \), for which \( |J_\gamma(n)| < \frac{1}{2} \). Then
\[
F_\gamma(n) = \frac{d}{dn} \left[ n \log (1 + J_\gamma(n)) \right]
\]
\[ = \frac{d}{dn} \left[ n \left( J_\gamma - \frac{1}{2} J_\gamma^2 + \frac{1}{3} J_\gamma^3 - \ldots \right) \right], \quad J_\gamma \equiv J_\gamma(n), \]
\[ = E_\gamma(n) - 1 + \frac{d}{dn} \left( -\frac{1}{2} n J_\gamma^2 + \frac{1}{3} n J_\gamma^3 - \ldots \right). \]

whence
\[ F_\gamma(n) - \sqrt{3} \gamma(n) = H_\gamma(n) + \frac{d}{dn} \left( -\frac{1}{2} n J_\gamma^2 + \frac{1}{3} n J_\gamma^3 - \ldots \right). \tag{4.4} \]

Although this may seem a very involved formula, it provides a good way of expanding
\[ F_\gamma(n) - \sqrt{3} \gamma(n), \] given an expansion of \( \gamma(n) \).

Equation (4.3) yields
\[
H_\gamma(n) = (1 + \sum_{k=1}^{\infty} C_k n^{2k})\left( -\frac{\gamma^2}{21} + \frac{\gamma^4}{41} - \ldots + \sqrt{3} \left( -\frac{\gamma^3}{31} + \frac{\gamma^5}{51} - \ldots \right) \right)
\]
\[ + \sum_{k=1}^{\infty} C_k n^{2k}\{1 + \sqrt{3} \gamma(n)\}, \tag{4.5} \]
\[ = \sum_{j=1}^{\infty} D_j A_{3j}(n) D_j n^4 + O(n^{4+\delta}), \]
where the constants \( D_j \) depend only on the \( A_n \) in (4.1) (and on the \( C_k \)). Consider now the remaining terms on the right of (4.4). Since
\[ J_\gamma(n) = \frac{1}{n} \int_0^n (\sqrt{3} \gamma + H_\gamma), \]

-21-
equation (4.5) ensures that the expansion of \( J \) has exponents in \( \Lambda \), and hence that the expansion of

\[
\frac{d}{dn} (n^m J^m) = (1-m)J^m + mJ^m \left[ (\sqrt{3} J + H) \right], \quad m \geq 2,
\]

has exponents in \( \Lambda \). Since \( \sqrt{3} J + H \) and hence \( J \), are known with error \( O(n^u) \), it follows that \( J^m \) and \( J^m \left[ (\sqrt{3} J + H) \right] \) are known with error \( O(n^{u-1}) \) for \( m \geq 2 \). We apply those two conclusions to the terms in (4.4) that follow \( H \), and the lemma is proved.

**Lemma 4.2.** For \( \xi \in (0,1) \) and \( p > -1 \),

\[
\int_0^1 r(\xi, n) \xi^p dn = \begin{cases} \tan (\pi s/2) \xi^p + \frac{1}{\pi} \sum_{k=1}^{\infty} \frac{\xi^{2k-1}}{k(1+p-2k)} & \text{if } p \notin \{1, 3, 5, \ldots\} \\ \frac{2}{\pi(1+p)} \xi^p \log \frac{1}{\xi} + \frac{\xi^{2k-1}}{k(1+p-2k)} & \text{if } p \notin \{1, 3, 5, \ldots\} \end{cases}
\]

**Proof.** Denote the integral by \( I_p(\xi) \), and define

\[
\phi_p(\xi) = \begin{cases} \xi^p, & 0 < \xi < 1, \\
p, & 0 < \xi \leq 1, \\
0, & 1 < \xi,
\end{cases}
\]

so that \( \phi_p(s) = \frac{1}{s+p} \) for \( s > -p \).

Then

\[
I_p(\xi) = \int_0^1 r(\xi, n) \phi_p(n) dn \quad (0 < \xi < 1)
\]

\[
= \int_0^1 r(1, u) \phi_p(\xi u) du
\]

\[
= R(1-s) \phi_p(s) \quad \text{for} \quad -\min(p, 1) < s < 1,
\]

where we have used the product rule, the notation \( r(1, u) = R(u) \), and the statement in (3.11b) that \( R(1-s) \) exists and is analytic for \( -1 < s < 1 \). By the inversion formula and (3.11b),

\[
I_p(\xi) = \frac{1}{2\pi i} \int_{c-i \infty}^{c+i \infty} \frac{\tan(\pi s/2)}{(s-1)(s+p)} ds, \quad -\min(p, 1) < c < 1.
\]
Consider the integral (of this integrand) around the (positively directed) rectangular contour with corners at \( c \pm 12n \) and \(-2n \pm 12n\); here \( n \in \mathbb{N} \). On the parts \( |t| < 2n \) and \(-2n < t < c, \ t = \pm 2n \) of this path, \( |\tan(\pi/2)| \) is bounded; hence the integrand is \( O(|s|^{-2}) \) there, and the contribution of these three sides tends to zero as \( n \to \infty \). Accordingly, \( \int P(\xi) \) equals the sum of the residues at poles with \( s < c \). If \( p \notin \{1,3,5,\ldots\} \), these are simple poles at \( s = -p \) and \( s = -1,-3,-5,\ldots \); if \( p \in \{1,3,5,\ldots\} \), there is a double pole at \( s = -p \), and simple poles at \( s = -(2k-1) \), \( k \in \mathbb{N} \setminus \{1+p\} \). Evaluating the residues, we obtain the result of the lemma.

**Lemma 4.3.** If \( p \in (2m+1, 2m+3) \) for some \( m \in \mathbb{N}_0 \) and \( |\xi(n)| < \text{const.} \ n^p \), then, for \( \xi \in (0,1) \),

$$\int_0^1 r(\xi,n)\xi(n)\,dn = \frac{1}{\pi} \sum_{k=1}^{m+1} \frac{2k+1}{k} \int_0^1 \xi(n)\xi(n)^{-2k}dn + O(\xi^p).$$

**Proof.** We have

$$\int_0^1 r(\xi,n)\xi(n)\,dn = \int_0^\xi r(\xi,n)\xi(n)\,dn + \frac{1}{\pi} \int_0^\xi \log(1 - \frac{\xi^2}{n^2})\xi(n)\,dn$$

$$= \int_0^1 r(1,u)\xi(u)\,du + \frac{1}{\pi} \sum_{k=1}^{m+1} \frac{2k-1}{k} \{ \int_0^1 - \int_0^\xi \xi(n)\xi(n)^{-2k}dn \}$$

$$+ \frac{1}{\pi} \sum_{k=m+2}^{\infty} \frac{2k-1}{k} \int_0^\xi \xi(n)\xi(n)^{-2k}dn.$$

The assumption that \( |\xi(n)| < \text{const.} \ n^p \), in which \( p \) is not an odd integer, ensures that the first, third and fourth integrals in this expression are respectively \( O(\xi^p) \), \( O(\xi^{p-2k+1}) \) with \( p > 2k-1 \) and \( O(\xi^{p-2k+1}) \) with \( p < 2k-1 \); the fourth is \( O(\xi^{p-2k+1}/k) \) as \( \xi \to 0 \) and \( k \to \infty \). This proves the lemma.

**Lemma 4.4.** Let the hypotheses in Lemma 4.1 hold, let \( \{\text{odd}\} \) denote the set \( \{1,3,5,\ldots\} \) of odd positive integers, and let \( 2p-1 \) be the largest odd integer less than \( u + 8_1 \). Then, as \( \xi \to 0 \),
\[ p_Y(\xi) = \sum_{k=1}^{P} \kappa_k e^{2\kappa-1} + \sum_{\lambda_j \in A_3(N) \setminus \mathbb{N}} \sum_{\lambda_j \in (N) \setminus \{\text{odd}\}} G_j e^{\lambda_j} \log \frac{1}{\xi} + \frac{\mu + \theta_1}{\lambda_j} + O(\xi) \]  

where

\[ F_j = \frac{\tan(\frac{\lambda_j \pi}{2})}{3(1+\lambda_j)} B_j, \quad G_j = \frac{2}{3\pi(1+\lambda_j)} B_j, \]  

but the coefficients \( \kappa_k \) are defined in terms of the global behavior of \( Y \) (and not merely in terms of the previous \( \lambda_j \) and \( B_j \)).

Remark. In (4.6), the set \( \lambda_j(N) \cap \{\text{odd}\} \) is probably empty for all \( N \), but, if it is not, the exponents in it cause no complication (in contrast to any exponents that there may be in the set \( \lambda_j(N) \cap \mathbb{N}(N) \) to be considered presently), because the logarithmic terms in (4.6) will be cancelled when we compute the corresponding expansion of \( Y \).

Proof. Recall from (3.2) that

\[ p_Y(\xi) = \frac{1}{3} \int_0^1 x(\xi)(\gamma(\eta)) - \sqrt{3} \gamma(\eta) d\eta + \frac{1}{3} \int_1^\infty x(\xi)(\gamma(\eta)) d\eta, \]  

and observe that, for \( \xi \in (0,1) \),

\[ \int_1^\infty x(\xi)(\gamma(\eta)) d\eta = -\frac{1}{\pi} \int_1^\infty \log(1 - \frac{2}{\eta}) x(\xi)(\gamma(\eta)) d\eta \]

\[ = \frac{1}{\pi} \sum_{k=1}^{\infty} \frac{2k-1}{k} \int_1^\infty (\gamma(\eta))^{-2k} d\eta, \]  

where the last integral is \( O(k^{-1}) \) as \( k \to +\infty \) because \( (\gamma(\eta)) = O(\log \eta) \) as \( \eta \to +\infty \).

For the first integral in (4.8), we apply Lemma 4.2 to the terms \( \lambda_j \) in (4.2), and Lemma 4.3 to the O-term there, recalling that \( k + \beta_1 \) is not an odd integer. The terms of form \( a_k e^{2k-1} \) (where \( k \in \mathbb{N} \) and \( a_k = O(k^{-2}) \) as \( k \to +\infty \)) that result from the \( B_j \) in (4.2), from the O-term there and from (4.9) are all collected in the first sum of (4.6) if \( 2k-1 < k + \beta_1 \) (so that \( k < P \)) or in the O-term of (4.6) if \( 2k-1 > k + \beta_1 \). In the second sum of (4.6) we have replaced \( \lambda_j(N) \setminus \{\text{odd}\} \) by \( \lambda_j(N) \setminus \mathbb{N} \) because (4.7) shows that \( F_j = 0 \) when \( \lambda_j \) is an even integer.
4.3. The expansion of $\gamma$

Lemma 4.4 provides a generalization of step (i) in the proof of Theorem 3.3, and we are now ready to generalize step (ii). To avoid a conceivable but highly improbable complication, we make

**ASSUMPTION A(N).** The sets $A_3(N)$ and $B(N)$ are disjoint.

The reason for this will be explained after the proof of Theorem 4.5. We have verified Assumption A(N) numerically for $N < 100$, and conjecture that it is true for all $N$.

**THEOREM 4.5.** If the hypotheses of Lemma 4.1 and Assumption A(N) hold, then

$$\gamma(\xi) = \sum_{j} \xi j \in A_3(N) \text{ and } 0(\xi^\nu) \text{ as } \xi \rightarrow 0,$$

for some $\nu \in (\nu_N^\infty).$ Here is a new coefficient $\lambda_j$, with $j$ such that $\lambda_j \not\in A_3(N) \setminus B(N)$, is determined by the previous coefficients $\lambda_i$, with $i$ such that $\lambda_i \in A_3(n)$, if and only if $\lambda_j \not\in B(N)$.

**Proof.** Again we use the notation in the proof of Lemma 3.2. In view of (4.6), we define

$$h_1(\xi) = \begin{cases} \sum_{k=1}^{\rho} \frac{\xi^{2k-1}}{k!} + \sum_{j} \xi j \in A_3(N) \setminus B(N) \text{ and } 0(\xi^\nu) \text{ as } \xi \rightarrow 0, \\ 0 \text{ if } 1 < \xi, \end{cases}$$

and $h_2 = h - h_1$. Accordingly, $h_2(\xi) = \mu(\xi) - h_1(\xi)$ if $\xi < 1$, and $h_2(\xi) = O(\xi^{-1} \log \xi)$ as $\xi \rightarrow \infty$, by (3.9). The smallest exponent in $h_1(\xi)$ is $1$, hence, for $\sigma > -1$,

$$h_1(s) = \sum_{k=1}^{\rho} \frac{\xi^{s+2k-1}}{k!} + \sum_{j} \xi j \in A_3(N) \setminus B(N) \text{ and } 0(\xi^\nu) \text{ as } \xi \rightarrow 0,$$

while $h_2(s)$ exists and is analytic for $-\beta_1 < \sigma < 1$, and $h_2(\sigma + i \epsilon) \in L_p^p(\mathbb{R}^\nu)$ for all $\sigma - (-\beta_1, 1)$ and all $p' > 2$. The obvious analytic (more precisely, meromorphic) continuations of $h_1$ and $h$ into the strip of convergence of $h_2$ will also be denoted by $\hat{h}_1$ and $\hat{h}$, respectively.
Let $\xi \in (0,1)$. We proceed from (3.14) and move the path of integration to $\sigma = -v$ for any $v \in (\lambda_1, u+\beta_1)$; the remarks, in the proof of Lemma 3.2, regarding the contribution of paths $t = \frac{\sigma}{v}$ remain valid when the $t$ there is replaced by $v$, and $e^{-\sigma Q(s)}h(s)$ has no singularities for $\sigma \in (-u-B, -\lambda_1)$. There are four sets of poles between the original path of integration ($\sigma = c$) and the new one ($\sigma = -v$); we consider one set at a time.

(a) The points $s = 1 - 2k$, $k \in \{1, 2, \ldots, p\}$, are simple poles of $h_1$, and, since $Q(1 - 2k) = -1$,
\[
\sum_{k=1}^{p} \text{res} = - \sum_{k=1}^{p} \frac{2k-1}{\xi^k}.
\]

(b) The points $s = \lambda_j$, where $\lambda_j \in \Lambda_j \cap \mathbb{N}$, are simple poles of $h_1$ and are not poles of $Q$ by Assumption $\text{A}(N)$, so that
\[
\sum_{\lambda_j} \text{res} = \sum_{\lambda_j \in \Lambda_j \cap \mathbb{N}} \frac{\lambda_j}{\xi} \log \frac{1}{\xi} + \frac{\sqrt{2\pi}}{2} (1 + \lambda_j) \xi^j.
\]

(c) The points $s = -\lambda_j$, where $\lambda_j \in \Lambda_j \cap \mathbb{N}$ and odd, are double poles of $h_1$, and
\[
\sum_{\lambda_j} \text{res} = - \sum_{\lambda_j \in \Lambda_j \cap \mathbb{N}} \frac{\lambda_j}{\xi} \log \frac{1}{\xi} + \frac{\sqrt{2\pi}}{2} (1 + \lambda_j) \xi^j.
\]

(d) The points $s = -\lambda_j$, where $\lambda_j \in \mathbb{B}(N)$, are simple poles of $Q$ and distinct from those of $h_1$ by Assumption $\text{A}(N)$, so that
\[
\sum_{\lambda_j} \text{res} = \sum_{\lambda_j \in \mathbb{B}(N)} \frac{\lambda_j}{\xi} \log \frac{1}{\xi} + \frac{\sqrt{2\pi}}{2} (1 + \lambda_j) \xi^j.
\]

where
\[
H_j = \text{res} Q(s)\big|_{s = -\lambda_j} h(-\lambda_j).
\]

Since $\xi \in (0,1)$, we have $g(\xi) = \gamma(\xi)$ and $h(\xi) = \rho_\gamma(\xi)$; equation (3.14) yields
\[
\gamma(\xi) = \rho_\gamma(\xi) + \sum \text{res} + \frac{1}{2\pi i} \int_{-\mu - i\infty}^{-\mu + i\infty} \xi^{-\sigma} Q(s) h(s) ds.
\]

We substitute for $\rho_\gamma(\xi)$ from (4.6) and for the residues, noting that the terms $\frac{\sqrt{2\pi}}{2} \frac{\lambda_j}{\xi}$, and those containing $\log(1/\xi)$, in (4.6) are cancelled by the residues, there results
\[
\gamma(\xi) = \sum_{\lambda_j \in \Lambda_j \cap \mathbb{N}} \frac{\lambda_j}{\xi} \log \frac{1}{\xi} + \frac{\sqrt{2\pi}}{2} \sum_{\lambda_j \in \Lambda_j \cap \mathbb{N}} (1 + \lambda_j) \xi^j + \frac{\xi}{2\pi} \int_{-\mu}^{-\mu + i\infty} \xi^{-\sigma} h(-\eta + it) dt + O(\xi^{u+\beta_1}).
\]
where \( Q(-v+i.) \in L_{3/2}(-\infty,\infty) \), because \( Q(-v+it) = o(t^{-1}) \) as \( t \to \infty \), and 
\[ h(-v+i.) \in L_{3}(-\infty,\infty) \), so that the integral (without the factor \( \xi' /2\pi \)) is bounded, by Hölder's inequality.

We claim that all the exponents \( \lambda_j \) in (4.11) belong to \( \Lambda_B(N) \). Indeed, if 
\( \lambda_j \in \Lambda \setminus \Lambda_B \), then \( \lambda_j \in \{2, 4, 6, \ldots \} \), and this is impossible in (4.11) because there
\( \lambda_j \notin \mathbb{N} \) or \( \lambda_j \notin \{\text{odd} \} \) or \( \lambda_j \notin \mathbb{B} \). Since in (4.11) each \( \lambda_j < \lambda_B \), it follows that \( \lambda_j \in \Lambda_B(N) \).

Comparison with (4.1) shows that in (4.11) the coefficient of \( \xi^j \) must equal \( \lambda_j \)
whenever \( \lambda_j \in \Lambda_B(n) \). Consider the coefficient of \( \xi^j \) when \( \lambda_j \notin \Lambda_B(N) \setminus \Lambda_B(n) \); if
\( \lambda_j \notin \mathbb{B}(N) \), then (4.7) and Lemma 4.1 show that the coefficient is determined by the
previous \( \lambda_i \); if \( \lambda_j \notin \mathbb{B}(N) \), then (4.10) shows that the coefficient is not so determined.
Of course, we relabel this new coefficient \( \lambda_j \) in either case.

Remark. Assume that the hypotheses of Lemma 4.1 hold but that Assumption \( A(n) \) is
false, for some particular \( n \) and \( M(n,\mu) \). Then there exists an exponent
\( \lambda_m \notin \Lambda_B(N) \cap \mathbb{B}(N) \), and \( \hat{Q}(s)h(s) \) has a double pole at \( s = -\lambda_m \). This causes the term
with \( j = m \) in (4.11) to be replaced by
\[ \int_{\lambda_m}^{\lambda_m} e^{(a+b \log \frac{1}{\xi})} \],
where \( a, b \) are constants and \( b \neq 0 \). At the \( N^{th} \) stage (and perhaps earlier), the hypo-
thesis (4.1), with \( n \) replaced by \( M(n,\mu) \), must be modified to include the logarithmic
term; at still later stages, higher and higher powers of logarithms accrue in the process
of expanding \( FY - \sqrt{3y} \).

COROLLARY 4.6. If \( \psi \) is a solution of (2.1) and \( \lambda_3 \cap \mathbb{B} = \emptyset \), then
\[ \psi(\xi) \sim \frac{\pi}{6} + \sum_{\lambda_j \in \Lambda_B} \frac{\lambda_j}{\lambda_j} \xi^j \quad \text{as} \quad \xi \to 0 \quad (4.12) \]
in the sense of an asymptotic expansion:
\[ \psi(\xi) - \frac{\pi}{6} - \sum_{\lambda_j \in \Lambda_B(n)} \frac{\lambda_j}{\lambda_j} \xi^j = o(\xi^n) \quad \text{as} \quad \xi \to 0 \],
for all \( n \in \mathbb{N} \).
COROLLARY 4.7. If \( A_3 \cap B = \emptyset \), then the complex co-ordinate \( z = x + iy \) is related to the complex potential \( x = \phi + i\psi \) by

\[
z(\chi) = \int_{\phi}^{\chi} \left( -g - \frac{1}{x} \right)^{1/3} \left[ 1 + \sum_{j=1}^{\lambda} K_j \left( i\chi \right)^{-\lambda} \right] \, d\chi \quad \text{as} \quad \chi \to 0,
\]

(4.13)

where the constant \( g \) is the gravitational acceleration, \( \arg(i\chi) \in [-\pi/2, \pi/2] \), and the coefficients \( K_j \) are real. If the set \( \{i\} \cup B \) is linearly independent over the rationals, then the sum in (4.13) is only over exponents in

\[
A_{\beta,0} = \{m \beta_1 + \ldots + m_r \beta_r : m_1 + \ldots + m_r > 1\};
\]

in other words, \( K_j = 0 \) whenever \( \lambda_j \notin A_{\beta,0} \).

Proof. On the image \( Y = 0 \) of the free boundary, we have

\[
\chi = -a\theta + \sum_{k=1}^{\infty} c_k \theta^{2k}, \quad a > 0,
\]

(4.14)

where \( a \) and \( c_k \) are constants, and the series converges for sufficiently small values of \( |\theta| \). Define \( \theta(\theta) = \Psi(\Psi(\theta)) \); it is a basic hypothesis in the derivation of (1.1) that \( x \) and \( \theta \) are odd functions of \( \theta \), while \( y \) is even. Accordingly, (4.12) and (4.14) imply that

\[
\theta(\theta) \sim -\text{sgn} \left( \frac{x}{6} + \sum_{j=1}^{\lambda} B_j \theta^{1/3} \right) \quad \text{as} \quad \theta \to 0,
\]

(4.15)

for certain coefficients \( B_j \) (which are not those in (4.2)). Moreover, it follows from the boundary condition of the basic problem (or, equivalently, from (1.1), from the fact that \( \theta(\theta) \) is the boundary value of \( \text{Im} \log(dx/d\chi) \), and from suitable choice of the additive constant in \( \text{Re} \log(dx/d\chi) \)) that

\[
\frac{\partial \chi}{\partial \theta} = \left( -g \theta \sin \theta \right)^{-1/3} e^{i\theta(\theta)} \quad \theta \in (-\pi, 0) \setminus \{0\}, \quad Y = 0,
\]

(4.16)

for some constant \( c > 0 \). Using (4.15) in this formula, we obtain (4.13) for \( x(\theta) \) (that is, for \( \arg(i\chi) = \pi/2 \) or \( -\pi/2 \)); that the \( K_j \) are real follows from the symmetry.

To prove that the asymptotic series for \( z(\chi) \) in (4.13) is the only appropriate extension (into the half-plane \( Y < 0 \)) of the series for \( x(\theta) \), we let \( h \) denote the difference between two such extensions to finitely many terms. It then suffices to prove the following: if \( h(x) \) is analytic and bounded in \( D = \{x : 0 < |x| < c, \quad -\pi/2 < \arg(i\chi) < \pi/2\} \), and if, for some constant \( v > 0 \), we have \( h(\phi) = O(|\phi|^v) \) as \( \phi \to 0 \) (with \( \arg(i\theta) = -\pi/2 \) or \( \pi/2 \)), then \( h(x) = O(|x|^v) \) as \( x \to 0 \) in \( D \). Now this follows
from application of the Phragmén-Lindelöf theorem ([11], p. 176) to the function \( f \) defined by \( f(x) = x^{-1/2} \). A suitable auxiliary function is \( w(x) = \exp\left(-\left(1x\right)^{1/2}\right) \).

It remains to prove that, if the set \([1] \cup B\) is linearly independent over the rationals, then \( X_j = 0 \) whenever \( \lambda_j \in A_B \setminus A_B, 0 \). Suppose that

\[
\psi(\theta) = \left(3 \right)^{1/2} \left(1 + b_1(1\theta)^{1/2} + b_n(1\theta)^{n/2} + k(1\theta)^{\lambda} + O(|\theta|^\mu)\right),
\]

where \( 0 < \alpha_1 < \ldots < \alpha_n < \lambda < \mu \), each \( \alpha_j \in A_B, 0 \) and \( \lambda \in A_B \setminus A_B, 0 \). We know, from our construction by way of (4.16), that this approximation may be differentiated term by term, and that the exact function \( z(\theta) \) satisfies the boundary condition of the basic problem. Then a slight variant of a calculation by Norman ([9], p. 262) shows that, under the foregoing hypothesis, the coefficient \( k = 0 \). Repeated application of this argument proves the result.
Appendix. On the numbers $\beta_j$.

We recall that $\beta_j$ denotes the $j$-th positive root of $\sqrt[3]{1+b} = \tan(\beta\pi/2)$; its simplest properties were noted in Lemma 3.1. We now probe a little deeper.

**Lemma A.1.** Each number $\beta_j$ is transcendental.

*Proof.* We abbreviate $\beta_j$ to $\beta$ for any fixed $j \in \mathbb{N}$. It was shown by Grant ([4], p. 260), albeit somewhat tersely, that $\beta$ is irrational. To prove it transcendental, assume the contrary: that $\beta$ is an algebraic irrational number. The equation of which $\beta$ is a root may be written

$$\beta = \frac{1+i\sqrt{3}(1+b)}{1-i\sqrt{3}(1+b)}$$  \hspace{1cm} (A.1)

Here the right-hand side is an algebraic number (because $\beta$ is one, and the algebraic numbers form a field), so, therefore, is the left-hand side. In other words, $(-1)^\beta$ is an algebraic number. But this contradicts the Gelfond-Schneider theorem ([10], p. 76), which states that $a^b$ is transcendental whenever $a$ is algebraic (and neither $0$ nor $1$) and $b$ is algebraic and irrational.

**Remark.** In Corollary 4.6 we assumed that the sets

$$A_3 = \{2\} \cup \{m \beta_1 + \ldots + m \beta_r + \sum_{r+1}^m 2 : m_1 + \ldots + m_r + m_{r+1} > 2\}$$

and

$$B = \{\beta_1, \beta_2, \beta_3, \ldots\}$$

are disjoint. It is clear that the truth of this assumption would be implied by the truth of the following, very natural

**Conjecture.** The set $\{1\} \cup B$ is linearly independent over the rationals.

(That is, if $x_1, \ldots, x_k$ are distinct elements of $\{1\} \cup B$, and $a_1, \ldots, a_k$ are rational numbers, then $\sum_{n=1}^k a_n x_n = 0$ implies that $a_1 = \ldots = a_k = 0$.) Unfortunately, our only result in this direction is the following.

**Theorem A.2.** If $j \neq k$, then the set $\{1, \beta_j, \beta_k\}$ is linearly independent over the rationals.

*Proof.* Assume the contrary: then there exist integers $p, q$ and $r$, not all zero, such that $p\beta_j + q\beta_k = r$. Neither $p$ nor $q$ can be zero; we may suppose that $r > 0$.
then at least one of $p$ and $q$ must be positive, say $p > 0$. Now

$$(-1)^r \cdot e^{i\pi r} = \left(e^{i\pi \beta_k}ight)^p \cdot \left(e^{i\pi \beta_k}ight)^q = \frac{1 + i\sqrt{3}(1 + \beta_j)}{1 - i\sqrt{3}(1 + \beta_j)} \cdot \frac{p + i\sqrt{3}(1 + \beta_j)}{q}$$

by (A.1), since $\beta_k = (r-p\beta_j)/q$, we have

$$(-1)^r \left(1 + i\sqrt{3}(1 + \beta_j)\right)^p \left(1 + i\sqrt{3}(1 + \beta_j)\right)^q = \left(1 + i\sqrt{3}(1 + \beta_j)\right)^p \left(1 + i\sqrt{3}(1 + \beta_j)\right)^q.$$  \hspace{1cm} (A.2)

Suppose that $q > 0$, expanding both sides of (A.2), and re-arranging the result, we obtain

$$\sqrt{3} P_1(\beta_j) = P_2(\beta_j),$$

where $P_1$ and $P_2$ are polynomials with (real) rational coefficients, and $\beta_j$ is real. Hence each polynomial is zero, which makes $\beta_j$ an algebraic number and thus contradicts Lemma A.1.

If $q < 0$, we re-write (A.2) to have positive exponents $p$ and $q$ on both sides, and argue as before.
FIGURE 1. Images by various conformal maps of one period of the flow domain, say 
$-\frac{1}{2} \lambda < x < \frac{1}{2} \lambda, \ y < Y(x)$, for the case $b = 1$ of periodic waves in liquid of infinite 
depth. (a) The physical plane, that of $z = x + iy$; (b) plane of the complex potential 
$\chi = \phi + iy$; (c) plane of $\zeta = e^{is}$, showing the variable $s$ in (1.1); (d) plane of 
$\omega = \sigma + i \xi$, showing the variable $\xi$ in (2.1).
References


On the Behavior Near the Crest of Waves of Extreme Form

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The angle $\phi$ which the free boundary of an extreme wave makes with the horizontal is the solution of a singular, non-linear integral equation that does not fit (as far as we know) into the theory of compact operators on Banach spaces. It has been proved only recently that solutions exist and that (as Stokes suggested in 1880) these solutions represent waves with sharp crests of included angle $2\pi/3$. In this paper we use the integral equation, known properties of solutions and the technique of the Mellin transform to obtain the asymptotic expansion.
ABSTRACT (continued)

\[
\phi(s) = \frac{\pi}{6} + \sum_{n=1}^{k} a_n s^n + o(s^k) \text{ as } s \to 0 ,
\]

(\*)

to arbitrary order; the co-ordinate \( s \) is related to distance from the crest as measured by the velocity potential rather than by length. The first few (and probably all) of the exponents \( \mu_n \) are transcendental numbers. We are unable to evaluate the coefficients \( a_n \) explicitly, but define some in terms of global properties of \( \phi \), and the others in terms of earlier coefficients. The derivation of (\*) includes an assumption about a question in number theory; if that assumption should be false, logarithmic terms \( a_{n,j} s^n (\log s)^{m_{n,j}} \), with \( m_{n,j} \) a positive integer, would enter the series at very large values of \( n \).

Our results confirm the heuristic calculations of Grant [4] and Norman [9], and, in effect, remove a tentative element from certain statements in those papers.
END

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