

Existence and regularity for an energy maximization problem in two dimensions

Spyridon Kamvissis

Max Planck Institute for Mathematics in the Sciences, Leipzig, Germany

Evguenii A. Rakhmanov

Department of Mathematics, University of South Florida, Tampa, Florida 33620

(Received 8 February 2005; accepted 6 June 2005; published online 27 July 2005)

We consider the variational problem of maximizing the weighted equilibrium Green's energy of a distribution of charges free to move in a subset of the upper half-plane, under a particular external field. We show that this problem admits a solution and that, under some conditions, this solution is an S-curve (in the sense of Gonchar-Rakhmanov). The above problem appears in the theory of the semiclassical limit of the integrable focusing nonlinear Schrödinger equation. In particular, its solution provides a justification of a crucial step in the asymptotic theory of nonlinear steepest descent for the inverse scattering problem of the associated linear non-self-adjoint Zakharov-Shabat operator and the equivalent Riemann-Hilbert factorization problem. © 2005 American Institute of Physics.

[DOI: 10.1063/1.1985069]

I. INTRODUCTION

Let $\mathbb{H} = \{z : \text{Im } z > 0\}$ be the complex upper-half plane and $\bar{\mathbb{H}} = \{z : \text{Im } z \geq 0\} \cup \{\infty\}$ be the closure of \mathbb{H} . Let also $\mathbb{K} = \{z : \text{Im } z > 0\} \setminus \{z : \text{Re } z = 0, 0 < \text{Im } z \leq A\}$, where A is a positive constant. In the closure of this space, $\bar{\mathbb{K}}$, we consider the points ix_+ and ix_- , where $0 \leq x < A$ as distinct. In other words, we cut a slit in the upper half-plane along the segment $(0, iA)$ and distinguish between the two sides of the slit. The point infinity belongs to $\bar{\mathbb{K}}$, but not \mathbb{K} . We define \mathbb{F} to be the set of all “continua” F in $\bar{\mathbb{K}}$ (i.e., connected compact sets) containing the distinguished points $0_+, 0_-$.

Next, let $\rho^0(z)$ be a given complex-valued function on $\bar{\mathbb{H}}$ satisfying

$\rho^0(z)$ is holomorphic in \mathbb{H} ,

$\rho^0(z)$ is continuous in $\bar{\mathbb{H}}$,

$\text{Re}[\rho^0(z)] = 0$, for $z \in [0, iA]$,

$\text{Im}[\rho^0(z)] > 0$, for $z \in (0, iA] \cup \mathbb{R}$.

Define $G(z; \eta)$ to be the Green's function for the upper half-plane

$$G(z; \eta) = \log \frac{|z - \eta^*|}{|z - \eta|} \quad (2)$$

and let $d\mu^0(\eta)$ be the non-negative measure $-\rho^0(\eta)d\eta$ on the segment $[0, iA]$ oriented from 0 to iA . The star denotes complex conjugation. Let the “external field” ϕ be defined by

$$\phi(z) = - \int G(z; \eta) d\mu^0(\eta) - \operatorname{Re} \left(i\pi J \int_z^{iA} \rho^0(\eta) d\eta + 2iJ(zx + z^2t) \right), \quad (3)$$

where x, t are real parameters with $t \geq 0$ and $J=1$, for $x \geq 0$, while $J=-1$, for $x < 0$. Re denotes the real part.

The particular form of this field is dictated by the particular application to the dynamical system we are interested in. The conditions (1) are natural in view of this application. But many of our results in this paper are valid if the term $zx + z^2t$ is replaced by any polynomial in z . Here x, t are in fact the space and time variables for the associated partial differential equation (PDE) problem [see (9) and (10) below].

Let \mathbb{M} be the set of all positive Borel measures on $\bar{\mathbb{K}}$, such that both the free energy

$$E(\mu) = \int \int G(x, y) d\mu(x) d\mu(y), \quad \mu \in \mathbb{M} \quad (4)$$

and $\int \phi d\mu$ are finite. Also, let

$$V^\mu(z) = \int G(z, x) d\mu(x), \quad \mu \in \mathbb{M} \quad (5)$$

be the Green's potential of the measure μ .

The weighted energy of the field ϕ is

$$E_\phi(\mu) = E(\mu) + 2 \int \phi d\mu, \quad \mu \in \mathbb{M}. \quad (6)$$

Now, given any continuum $F \in \mathbb{F}$, the equilibrium measure λ^F supported in F is defined by

$$E_\phi(\lambda^F) = \min_{\mu \in M(F)} E_\phi(\mu), \quad (7)$$

where $M(F)$ is the set of measures in \mathbb{M} which are supported in F , provided such a measure exists. $E_\phi(\lambda^F)$ is the equilibrium energy of F .

The aim of this paper is to prove the existence of a so-called S-curve (Ref. 1) joining the points 0_+ and 0_- and lying entirely in $\bar{\mathbb{K}}$, at least under some extra assumptions. By S-curve we mean an oriented curve F such that the equilibrium measure λ^F exists, its support consists of a finite union of analytic arcs and at any interior point of $\operatorname{supp} \mu$,

$$\frac{d}{dn_+}(\phi + V^{\lambda^F}) = \frac{d}{dn_-}(\phi + V^{\lambda^F}), \quad (8)$$

where the two derivatives above denote the normal (to $\operatorname{supp} \mu$) derivatives.

To prove the existence of the S-curve we will first need to prove the existence of a continuum F maximizing the equilibrium energy over \mathbb{F} . Then we will show that the maximizer is in fact an S-curve.

It is not always true that an equilibrium measure exists for a given continuum. The Gauss-Frostman theorem (Ref. 2, p.135) guarantees the existence of the equilibrium measure when F does not touch the boundary of the domain \mathbb{H} . This is not the case here. Still, as we show in the next section, in the particular case of our special external field, for any given x, t and for a large class of continua F not containing infinity, the weighted energy is bounded below and λ^F exists. So, in particular, we do know that the supremum of the equilibrium weighted energies over all continua is greater than $-\infty$.

S-curves were first defined in Ref. 1, where the concept first arose in connection with the problem of rational approximation of analytic functions. Our own motivation comes from a seem-

ingly completely different problem, which is the analysis of the so-called semiclassical asymptotics for the focusing nonlinear Schrödinger equation. More precisely, we are interested in studying the behavior of solutions of

$$i\hbar\partial_t\psi + \frac{\hbar^2}{2}\partial_x^2\psi + |\psi|^2\psi = 0, \quad (9)$$

$$\text{under } \psi(x,0) = \psi_0(x),$$

in the so-called semiclassical limit, i.e., as $\hbar \rightarrow 0$. For a concrete discussion, let us here assume that $\psi_0(x)$ is a positive “bell-shaped” function; in other words assume that

$$\psi_0(x) > 0, \quad x \in \mathbb{R},$$

$$\psi_0(-x) = \psi_0(x),$$

$$\psi_0 \text{ has one single local maximum at } 0, \quad \psi_0(0) = A, \quad (10)$$

$$\psi_0''(0) < 0,$$

$$\psi_0 \text{ is Schwartz.}$$

This is a completely integrable PDE and can be solved via the method of inverse scattering. The semiclassical limit is analyzed in the recent research monograph.³ In Chap. 8 of Ref. 3 it is noted that the semiclassical problem is related and can be reduced to a particular “electrostatic” variational problem of maximizing the equilibrium energy of a distribution of charges that are free to move under a given external electrostatic field (assuming that the WKB-approximated density of the eigenvalues admits a holomorphic extension in the upper half-plane). In fact, it is pointed out that the existence and regularity of an S-curve implies the existence of the so-called “g-function” necessary to justify the otherwise rigorous methods employed in Ref. 3.

We would like to point out that the problem of the existence of the “g function” for the semiclassical nonlinear Schrödinger problem is not a mere technicality of isolated interest. Rather, it is an instance of a crucial element in the asymptotic theory of Riemann-Hilbert problem factorizations associated to integrable systems. This asymptotic method has been made rigorous and systematic in Ref. 4 where in fact the term “nonlinear steepest descent method” was first employed to stress the relation with the classical “steepest descent method” initiated by Riemann in the study of exponential integrals with a large phase parameter. Such exponential integrals appear in the solution of Cauchy problems for linear evolution equations, when one employs the method of Fourier transforms. In the case of nonlinear integrable equations, on the other hand, the nonlinear analog of the Fourier transform is the scattering transform and the inverse problem is now a Riemann-Hilbert factorization problem. While in the “linear steepest descent method” the contour of integration must be deformed to a union of contours of “steepest descent” which will make the explicit integration of the integral possible, in the case of the “nonlinear steepest descent method” one deforms the original Riemann-Hilbert factorization contour to appropriate steepest descent contours where the resulting Riemann-Hilbert problems are explicitly solvable.

In the linear case, if the phase and the critical points of the phase are real it may not be necessary to deform the integration contour. One has rather a Laplace integral problem on the contour given. For Riemann-Hilbert problems the analog is the self-adjointness of the underlying Lax operator. In this case the spectrum of the associated linear Lax operator is real and the original Riemann-Hilbert contour is real. The “deformation contour” must then stay near the real line. One novelty of the semiclassical problem for (9) and (10) studied in Ref. 3 however is that, due to the non-self-adjointness of the underlying Lax operator, the “target contour” is very specific (if not

unique) and by no means obvious. It is best characterized via the solution of a maximin energy problem, in fact it is an S-curve. The term “nonlinear steepest descent method” thus acquires full meaning in the non-self-adjoint case.

Given the importance and the recent popularity of the “steepest descent method” and the various different applications to such topics as soliton theory, orthogonal polynomials, solvable models in statistical mechanics, random matrices, combinatorics and representation theory, we believe that the present work offers an important contribution. In particular we expect that the results of this paper may be useful in the treatment of Riemann-Hilbert problems arising in the analysis of general complex or normal random matrices.

On the other hand, we believe that the main results of this paper, Theorems 3, 4, 5, 7, 8 are interesting on their own. This paper can be read without the applications to dynamical systems in mind. It concerns existence and regularity of a solution to an energy variational maximin problem in the complex plane.

The method used to prove the existence of the S-curves arising in the solution of the “max-min” energy problem was first outlined in Ref. 1 and further developed in Ref. 5, at least for logarithmic potentials. But, the concrete particular problem addressed in this paper involves additional technical issues.

The main points of the proof of our results are as follows:

- (i) Appropriate definition of the underlying space of continua (connected compact sets) and its topology. This ensures the compactness of our space of continua which is crucial in proving the existence of an energy maximizing element.
- (ii) Proof of the semicontinuity of the energy functional that takes a continuum to the energy of its associated equilibrium measure (Theorem 3).
- (iii) Proof of existence of an energy maximizing continuum (Theorem 4).
- (iv) A discussion of how some assumptions ensure that the maximizing continuum does not touch the boundary of the underlying space except at a finite number of points. This ensures that variations of continua can be taken.
- (v) Proof of formula (22) involving the support of the equilibrium measure on the maximizing continuum and the external field (Theorem 5).
- (vi) Proof that the support of the equilibrium measure on the maximizing continuum consists of a union of finitely many analytic arcs.
- (vii) Proof that the maximizing continuum is an S-curve (Theorems 7 and 8).

The paper is organized as follows. In the rest of Sec. I, we introduce the appropriate topology for our set of continua that will provide the necessary compactness. In Sec. II, we prove a “Gauss-Frostman” type theorem which shows that the variational problem that we wish to solve is not vacuous. In Sec. III, we present the proof of upper semicontinuity of a particularly defined “energy functional.” In Sec. IV, we present a proof of existence of a solution of the variational problem. Existence is thus derived from the semicontinuity and the compactness results acquired earlier. In Sec. V, we show that, at least under a simplifying assumption, the “max-min” solution of the variational problem does not touch the boundary of the underlying domain, except possibly at some special points. This enables us to eventually take variations and show that the max-min property implies regularity of the support of the solution and the S-property in Secs. VI and VII. By regularity, we mean that the support of the maximizing measure is a finite union of analytic arcs. In Sec. VIII, we conclude by stating the consequence of the above results in regard to the semiclassical limit of the nonlinear Schrödinger equation.

We also include three appendixes. The first one discusses in detail some topological facts regarding the set of closed subsets of a compact space, equipped with the so-called Hausdorff distance. The fact that such a space is compact is vital for proving existence of a solution for the variational problem. The second appendix presents the semiclassical asymptotics for the initial value problem (9) and (10) in terms of theta functions, under the S-curve assumption (as in Ref. 3). It is included so that the connection with the original motivating problem of semiclassical NLS is made more explicit. The third appendix sketches an argument on how to get rid of the simpli-

fying assumption introduced in Sec. V. We feel that the argument leading to Theorem 9 is quite transparent, although an absolutely rigorous proof would require a thorough revisiting of the discussion and results in Ref. 3.

Following Ref. 6 (see Appendix A) we introduce an appropriate topology on \mathbb{F} . We think of the closed upper half-plane $\bar{\mathbb{H}}$ as a compact space in the Riemann sphere. We thus choose to equip $\bar{\mathbb{H}}$ with the “chordal” distance, denoted by d_0 , that is the distance between the images of z and ζ under the stereographic projection. This induces naturally a distance in $\bar{\mathbb{K}}$ [so $d_0(0_+, 0_-) \neq 0$]. We also denote by d_0 the induced distance between compact sets E, F in $\bar{\mathbb{K}}$, $d_0(E, F) = \max_{z \in E} \min_{\zeta \in F} d_0(z, \zeta)$. Then, we define the so-called Hausdorff metric on the set $I(\bar{\mathbb{K}})$ of closed nonempty subsets of $\bar{\mathbb{K}}$ as follows:

$$d_{\mathbb{K}}(A, B) = \sup(d_0(A, B), d_0(B, A)). \quad (11)$$

In Appendix A, we prove the following.

Lemma A.1: The Hausdorff metric defined by (11) is indeed a metric. The set $I(\bar{\mathbb{K}})$ is compact and complete.

Now, it is easy to see that \mathbb{F} is a closed subset of $I(\bar{\mathbb{K}})$. Hence \mathbb{F} is also compact and complete.

Remarks: (1) Because of the particular symmetry $\psi(x) = \psi(-x)$ of the solution to the Cauchy problem (9) and (10) we will restrict ourselves to the case $x \geq 0$ from now on. We then set $J=1$ and the external field is

$$\phi(z) = - \int G(z; \eta) d\mu^0(\eta) - \operatorname{Re} \left(i\pi \int_z^{iA} \rho^0(\eta) d\eta + 2i(zx + z^2t) \right). \quad (3')$$

(2) The function ρ^0 expresses the density of eigenvalues of the Lax operator associated to (9), in the limit as $h \rightarrow 0$. WKB theory can be used to derive an expression for ρ^0 in terms of the initial data $\psi^0(x)$ via an Abel transform (see Ref. 3), from which it follows that

$$\operatorname{Re}[\rho^0(z)] = 0 \quad \text{for } z \in [0, iA],$$

$$\operatorname{Im}[\rho^0(z)] > 0 \quad \text{for } z \in (0, iA].$$

The rest of the conditions (1) are not a necessary consequence of WKB theory. In particular, it is not *a priori* clear what the analyticity properties of ρ^0 are. In this paper, we *assume*, for simplicity, that ρ^0 admits a continuous extension in the closed upper complex plane which is holomorphic in the open upper complex plane. We also assume that $\operatorname{Im} \rho^0$ is positive in the real axis. This will be used later to show that the maximizing continuum does not touch the real line, except at $0_+, 0_-, \infty$. It is a simplifying but not essential assumption. All conditions (1) are satisfied in the simple case where the initial data are given by $\psi(x, 0) = A \operatorname{sech} x$, where A is a positive constant.

(3) It follows that ϕ is a subharmonic function in \mathbb{H} which is actually harmonic in \mathbb{K} ; it also follows that it is upper semicontinuous in \mathbb{H} . It is then subharmonic and upper semicontinuous in $\bar{\mathbb{H}}$ except at infinity.

(4) Even though in the end we wish that the maximum of $E_\phi(\lambda^F)$ over “continua” F is a regular curve, we will begin by studying the variational problem over the set of continua \mathbb{F} and only later (in Sec. VI) we will show that the maximizing continuum is in fact a nice curve. The reason is that the set \mathbb{F} is compact, so once we prove in Sec. III the upper semicontinuity of the energy functional, existence of a maximizing continuum will follow immediately.

II. A GAUSS-FROSTMAN THEOREM

We claim that for any continuum $F \in \mathbb{F}$, not containing the point ∞ and approaching 0_+ , 0_- nontangentially to the real line, the weighted energy is bounded below and the equilibrium measure λ^F exists. This is not true for any external field, but it is true for the field given by (3') because of the particular behavior of the function ρ^0 near zero.

We begin by considering the equilibrium measure on the particular contour F_0 that wraps itself around the straight line segment $[0, iA]$, say λ_0^F . We have

Proposition 1: Consider the contour $F_0 \in \mathbb{F}$ consisting of the straight line segments joining 0_+ to $iA_+ = iA$ and $iA = iA_-$ to 0_- . The equilibrium measure λ_0^F exists. Its support is the imaginary segment $[0, ib_0(x)]$, for some $0 < b_0(x) \leq A$, lying on the right-hand side of the slit $[0, iA]$. It can be written as $\rho(z)dz$ where $\rho(z)$ is a differentiable function in $[0, ib_0(x)]$, such that $\rho - \rho^0$ belongs in the Hölder class with exponent $1/2$.

Proof: See section 6.2.1 of Ref. 3; $\rho(z)$ can be expressed explicitly. The Hölder condition follows from 5.14 of Ref. 3. Note that the field ϕ is independent of time on F_0 , so λ_0^F is also independent of time.

From Proposition 1, it follows that the maximum equilibrium energy over continua is bounded below,

$$\max_{F \in \mathbb{F}} E_\phi(\lambda^F) = \max_{F \in \mathbb{F}} \min_{\mu \in M(F)} E_\phi(\mu) > -\infty. \tag{12}$$

The following formula is easy to verify:

$$E_\phi(\mu) - E_\phi(\lambda^F) = E(\mu - \lambda^F) + 2 \int (V^{\lambda^F} + \phi) d(\mu - \lambda^F), \tag{13}$$

for any μ which is a positive measure on the continuum F . Here

$$V^{\lambda^F}(u) = \int G(u, v) d\lambda^F(v),$$

where again $G(u, v)$ is the Green function for the upper half-plane.

To show that $E_\phi(\mu)$ is bounded below, all we need to show is that the difference $E_\phi(\mu) - E_\phi(\lambda^F)$ is bounded below.

Note that since $V^{\lambda^F} + \phi = 0$, on $\text{supp}(\lambda^F)$, the integral in (13) can be written as $\int (V^{\lambda^F} + \phi) d\mu$. We have

$$\begin{aligned} V^{\lambda^F}(z) + \phi(z) &= \int_0^{b_0(x)} \log \frac{|z + iu|}{|z - iu|} (-\rho)_{t=0} du + \phi \\ &= -\text{Re} \left(\int_0^{b_0(x)} \log \frac{|z + iu|}{|z - iu|} u^{1/2} du \right) + O(|z|) = O(|z|) \text{ near } z = 0. \end{aligned}$$

So we can write $V^{\lambda^F} + \phi \geq c(A, x)|z|$ in a neighborhood of $z=0$, where $c(A, x)$ will be some negative constant independent of z . Note that the dependence on t is not suppressed, but it is of order $O(|z|^2)$.

It is now not hard to see that the $O(|z|)$ decay implies our result.

Write $\mu = M\sigma$, where $M > 0$ is the total mass of μ and σ is a probability measure (on F). Choose ϵ such that for $|u| < \epsilon$ we have $V^{\lambda^F} + \phi \geq c(A, x)|u|$. Then

$$\begin{aligned} E_\phi(\mu) - E_\phi(\lambda^F) &\geq \int G(u, v) d(\mu - \lambda^F)(u) d(\mu - \lambda^F)(v) + 2 \int_{|v| \geq \epsilon} (V^{\lambda^F} + \phi)(v) d\mu(v) \\ &\quad + \int_{|v| < \epsilon} 2c(A, x)|v| d(\mu - \lambda^F)(v). \end{aligned} \tag{14}$$

The first integral of the right-hand side (RHS) can be written as $\int_{|u| \geq \epsilon, |v| \geq \epsilon} + 2 \int_{|u| < \epsilon, |v| \geq \epsilon} + \int_{|u|, |v| < \epsilon}$. The sum of the first integral plus the second term on the RHS of (14) is bounded below, by the standard Gauss-Frostman theorem (Ref. 6, p. 135). It remains to consider

$$\begin{aligned} & \left(\int_F + \int_{|v| \geq \epsilon} \right) \left(\int_{|u| < \epsilon} G(u, v) d(\mu - \lambda^F)(u) \right) d(\mu - \lambda^F)(v) + \int_{|v| < \epsilon} 2c(A, x) |v| d(\mu - \lambda^F)(v) \\ & \geq \left(\int_F + \int_{|v| \geq \epsilon} \right) \left(\int_{|u| < \epsilon} G(u, v) d(M\sigma - \lambda^F)(u) + 2c(A, x) |v| \right) d(M\sigma - \lambda^F)(v). \end{aligned} \quad (15)$$

Now, it is easy to see that for M large, and since F is nontangential to the real line,

$$\int_{|u| < \epsilon} G(u, v) d(M\sigma - \lambda^F)(u) + 2c(A, x) |v| \geq |v \log v| + 2c(A, x) |v|. \quad (16)$$

By adjusting ϵ if necessary, we have $|v \log v| > 2|c(A, x)| |v|$, and hence the integral in (16) is positive. Integrating again with respect to $M d\sigma - d\lambda^F$, again for M large, we see that the integral of (15) is positive.

Since when M is bounded above we have our estimates trivially, we get boundedness below over all positive M .

We have thus proved one part of our (generalized) Gauss-Frostman theorem.

Theorem 1: Let ϕ be given by (3'). Let F be a continuum in $\bar{\mathbb{K}} \setminus \infty$ and suppose that F is contained in some sector $\pi < \alpha < \arg(\lambda) < \beta < 0$ as $\lambda \rightarrow 0$. Let $M(F)$ be the set of measures $\in \mathbb{M}$ which are supported in F . [So, in particular their free energy is finite and $\phi \in L_1(\mu)$.] We have

$$\inf_{\mu \in M(F)} E_\phi(\mu) > -\infty. \quad (17)$$

Furthermore the equilibrium measure on F exists, that is there is a measure $\lambda^F \in M(F)$ such that $E_\phi[F] = E_\phi(\lambda^F) = \inf_{\mu \in M(F)} E_\phi(\mu)$.

Proof: The proof that (17) implies the existence of an equilibrium measure is a well-known theorem. For our particular field ϕ given by (3) it is easy to prove. Indeed, the identity

$$E(\mu - \nu) = 2E_\phi(\mu) + 2E_\phi(\nu) - 4E_\phi\left(\frac{\mu + \nu}{2}\right)$$

implies that any sequence μ_n minimizing $E_\phi(\mu)$ is a Cauchy sequence in (unweighted) energy. Since the space of positive measures is complete (see, for example, Ref. 7, Theorem 1.18, p. 90), there is a measure μ_0 such that $E(\mu_n - \mu_0) \rightarrow 0$. We then have $E(\mu_n) \rightarrow E(\mu_0) < +\infty$ and hence $\mu_n \rightarrow \mu_0$ weakly (see, e.g., Ref. 7, pp. 82–88; this is a standard result).

The fact that $\phi \in L_1(\mu_0)$ is trivial for our particular field.

III. SEMICONTINUITY OF THE ENERGY FUNCTIONAL

We consider the functional that takes a given continuum F to the equilibrium energy on this continuum,

$$\mathbb{E}: F \rightarrow E_\psi[F] = E_\psi(\lambda^F) = \inf_{\mu \in M(F)} \left(E(\mu) + 2 \int \psi d\mu \right) \quad (18)$$

and we want to show that it is continuous, if ψ is continuous in $\bar{\mathbb{H}}$. Note that this is not the case for the field ϕ given by (3'), since it has a singularity at ∞ ; that field is only upper semicontinuous. We will see how to circumvent this difficulty later. For the moment, ψ is simply assumed to be a continuous function in $\bar{\mathbb{H}}$.

Theorem 2: If ψ is a continuous function in $\bar{\mathbb{H}} \setminus \infty$ then the energy functional defined by (18) is continuous at any given continuum F not containing the point ∞ .

Proof: Suppose $G \in \mathbb{F}$, also not containing ∞ , with $d_{\mathbb{K}}(F, G) < d$, a small positive constant. Let $\lambda = \lambda_{\psi}^F$ be the equilibrium measure on F and $\mu = \lambda_{\psi}^G$ be the equilibrium measure on G .

We consider the Green's balayage of μ on F , say $\hat{\mu}$. Then $\text{supp } \hat{\mu} \in F$ and

$$V^{\hat{\mu}} = V^{\mu} \quad \text{on } F,$$

$$\int u \, d\hat{\mu} = \int u \, d\mu,$$

for any function u that is harmonic in $\mathbb{H} \setminus F$ and continuous in $\bar{\mathbb{H}}$.

Similarly consider $\hat{\lambda}$, the balayage of λ to G . We trivially have

$$E_{\psi}[\hat{\lambda}] \leq E_{\psi}(\hat{\lambda}),$$

$$E_{\psi}[\hat{\mu}] \leq E_{\psi}(\hat{\mu}).$$

Lemma 1: Suppose $Q \in \mathbb{F}$, μ some positive measure supported in \mathbb{K} and $\hat{\mu}$ is the Green's balayage to Q . Then

$$V^{\hat{\mu}} = V^{\mu} - V_{Q^c}^{\mu},$$

$$E(\hat{\mu}) = E(\mu) - E_{Q^c}(\mu),$$

where $E_{Q^c}(\mu)$ is the unweighted Green energy with respect to $Q^c = \mathbb{K} \setminus Q$. In particular, since unweighted energies are non-negative,

$$E(\hat{\mu}) \leq E(\mu).$$

Proof: The first identity follows from the fact that $V^{\hat{\mu}} - V^{\mu}$ vanishes on Q and the real line, and is harmonic in Q^c and superharmonic in \mathbb{K} .

Integrating $E(\hat{\mu}) = \int V^{\hat{\mu}} \, d\hat{\mu} = \int V^{\mu} \, d\hat{\mu} - \int V_{Q^c}^{\mu} \, d\hat{\mu} = \int (V^{\mu} - V_{Q^c}^{\mu}) \, d\mu = E(\mu) - E_{Q^c}(\mu)$. The proof of the Lemma follows.

So, let u_{ψ} be a function harmonic in $\mathbb{H} \setminus F$ such that $u_{\psi} = \psi$ on F and $u_{\psi} = 0$ on $\partial\mathbb{H} \setminus F$. By the definition of balayage one has $\int \psi \, d\hat{\mu} = \int u_{\psi} \, d\mu$.

We have $E_{\psi}[\hat{\mu}] \leq E_{\psi}(\hat{\mu}) = E(\hat{\mu}) + 2 \int \psi \, d\hat{\mu} \leq E(\mu) + 2 \int \psi \, d\mu + 2 \int (u_{\psi} - \psi) \, d\mu = E_{\psi}(\mu) + 2 \int (u_{\psi} - \psi) \, d\mu = E_{\psi}[G] + 2 \int (u_{\psi} - \psi) \, d\mu$.

In a small neighborhood of F , $\bar{F}_d = \{z : d(z, F) \leq d\}$, we have

$$\left| 2 \int (u_{\psi} - \psi)(y) \, d\mu(y) \right| \leq C \max_{y \in \bar{F}_d} |u_{\psi}(y) - \psi(y)|. \tag{19}$$

We assumed here that the equilibrium measures on continua near F are bounded above. This is easy to see. Suppose, first, that the point ∞ is not in F . Indeed, on the support of the equilibrium measure λ , we have

$$V^{\lambda} + \psi = 0.$$

If the equilibrium measures on continua near F were unbounded, then so would be the potentials V^{λ} . But ψ is definitely bounded near F . This contradicts the above equality.

Now given $y \in \bar{F}_d$, choose $z \in F$ such that $|z - y| = d$. The above expression (19) is less or equal than

$$C \max_{y \in \bar{F}_d} |u_{\psi}(y) - \psi(y) - u_{\psi}(z) + \psi(z)| \leq o(1) + C \max_{y \in \bar{F}_d} |u_{\psi}(y) - u_{\psi}(z)|.$$

It remains to bound $|u_{\psi}(y) - u_{\psi}(z)|$ by an $o(1)$ quantity.

The next Lemma is due to Milloux and can be found in Ref. 8.

Lemma 2: Suppose D is an open disc of radius R , with center z_0 ; let y be a point in D , F a continuum in \mathbb{C} , containing z_0 , and Ω be the connected component of $D \setminus F$ containing y . Let $w(z)$ be a function harmonic in Ω such that

$$w(z) = 0, \quad z \in F \cap \partial\Omega,$$

$$w(z) = 1, \quad z \in \partial\Omega \setminus F.$$

Then $w(y) \leq C(|y - z_0|/R)^{1/2}$.

Proof: See Ref. 8, p. 347.

Now, select a disc of radius $d^{1/2}$, centered on z . We have $|u_\psi(y) - u_\psi(z)| = o(1)$ on the part of F lying in the disc, while $|u_\psi(y) - u_\psi(z)|$ is bounded by some positive constant M on the disc boundary.

Lemma 3: Let Ω be a domain, $\partial\Omega = F_1 \cup F_2$ and

$$w_1 = 0, \quad z \in F_1 = 1, \quad z \in F_2,$$

$$w_2 = 1, \quad z \in F_1 = 0, \quad z \in F_2.$$

Suppose u is harmonic in Ω and

$$u(z) \leq \epsilon, \quad z \in F_1,$$

$$u(z) \leq M, \quad z \in F_2.$$

Then $u(z) \leq \epsilon w_2(z) + M w_1(z)$.

Proof: Maximum principle.

Now, using Milloux's lemma, we get $|u_\psi(y) - u_\psi(z)| \leq o(1)w_2(z) + Mw_1(z) \leq o(1) + MC(|y - z|/d)^{1/2} \leq o(1) + MCd^{1/2}$. This concludes the proof of Theorem 2.

We now recall that the energy continuity proof was based on the continuity of ψ . In our case, ϕ is upper semicontinuous and discontinuous at ∞ . Still we can prove that the energy is upper semicontinuous and that will be enough.

Theorem 3: For the external field given by (3'), the energy functional defined in (18) is upper semicontinuous.

Proof: We first note that if the external field ϕ' is upper semicontinuous away from infinity then so is the energy functional that takes a given continuum F to the equilibrium energy of F . Indeed, if ϕ' is upper semicontinuous away from infinity, then there exists a sequence of continuous functions (away from infinity) such that $\phi_n \downarrow \phi'$. Each functional $E_{\phi_n}[F]$ is continuous, away from infinity, and $E_{\phi_n}[F] \downarrow E_{\phi'}[F]$. So, $E_{\phi'}[F]$ is upper semicontinuous, away from infinity.

Now consider the field ϕ given by (3'). Let F be a continuum. If ∞ is not in F , then we are done. If $\infty \in F$, let $\lambda = \lambda^F$ be the equilibrium measure. We can assume that on the equilibrium measure ϕ is bounded by 0. Indeed, on the support of the equilibrium measure λ , we have

$$V^\lambda + \phi = 0.$$

But $V^\lambda \geq 0$, so $\phi \leq 0$.

This means that we can change ϕ to $\phi' = \min(\phi, 0)$, which is an upper semicontinuous function. Theorem 3 is proved.

Remark: If we naively consider the functional taking a measure to its weighted energy we will see that it is not continuous even if the external field is continuous. It is essential that the energy functional is defined on equilibrium measures.

IV. PROOF OF EXISTENCE OF A MAXIMIZING CONTINUUM

Theorem 4: For the external field given by (3'), there exists a continuum $F \in \mathbb{F}$ such that the equilibrium measure λ^F exists and

$$E_\phi[F](=E_\phi(\lambda^F)) = \max_{F \in \mathbb{F}} \min_{\mu \in M(F)} E_\phi(\mu).$$

Proof: We know (see, for example, Sec. II) that there is at least one continuum F for which the equilibrium measure exists and $E_\phi(\lambda^F) > -\infty$, for all time. On the other hand, clearly $E_\phi(\lambda^F) \leq 0$ for any F . Hence the supremum over continua in \mathbb{F} is finite (and trivially nonpositive), since \mathbb{F} is compact. Call it L .

We can now take a sequence F_n such that $E_\phi[F_n] \rightarrow L$. Choose a convergent subsequence of continua $F_n \rightarrow F$, say. By upper semicontinuity of the weighted energy functional,

$$\limsup E_\phi[F_n] \leq E_\phi[F] \leq L = \lim E_\phi[F_n].$$

So $L = E_\phi[F]$. The theorem is proved.

V. ACCEPTABILITY OF THE CONTINUUM

We have thus shown that a solution of the maximum-minimum problem exists. We do not know yet that the maximizing continuum is a contour. Clearly the pieces of the continuum lying in the region where the external field is positive do not support the equilibrium measure and by the continuity of the external field they can be perturbed to a finite union of analytic arcs. The real problem is to show that the support of the equilibrium measure is a finite union of analytic arcs. This will follow from the analyticity properties of the external field.

Note that the maximizing continuum cannot be unique, since the subset where the equilibrium measure is zero can be perturbed without changing the energy. A more interesting question is whether the support of the equilibrium measure of the maximizing contour is unique. We do not know the answer to this question but it is not important as far as the application to the semiclassical limit of NLS is concerned. (See Appendix B.)

It is important however, that the maximizing continuum does not touch the boundary of the space \mathbb{K} except of course at the points 0^- , 0^+ , and perhaps at ∞ . This is to guarantee that variations with respect to the maximizing contour can be properly taken.

The proof of the acceptability of the continuum requires two things.

- (i) The continuum does not touch the real negative axis.
- (ii) The continuum does not touch the real positive axis.

We will also make the following assumption.

Assumption (A): The continuum maximizing the equilibrium energy does not touch the linear segment $(0, iA]$.

Remark: Assumption (A) is not satisfied at $t=0$, where in fact the continuum is a contour F_0 wrapping around the linear segment $[0, iA]$. However, the case $t=0$ is well understood. The equilibrium measure for F_0 exists and its support is connected. On the other hand, assumption (A) is satisfied for small $t > 0$. (See Chap. 6 of Ref. 3.)

Remark: It is conceivable that at some positive t_0 there is an x for which assumption (A) is not satisfied. It can in fact be dropped but the analysis of the semiclassical limit of the nonlinear Schrödinger equation will get more tedious; see Appendix C.

Proposition 2: The continuum maximizing the equilibrium energy does not touch the real axis except at the points zero and possibly infinity.

Proof: (i) If $z < 0$, then $\phi(z) = \pi \int_z^0 |\rho^0(\eta)| d\eta > 0$.

This follows from an easy calculation, using the conditions defining ρ^0 . But we can always delete the strictly positive measure lying in a region where the field is positive and make the energy smaller. So even the solution of the "inner" minimizing problem must lie away from the real negative axis.

(ii) If $z > 0$, then again a short calculation shows that $d\phi/d \operatorname{Im} z > 0$, for $t > 0$.

It is crucial here that if $u \in \mathbb{R}$ then $G(u, v) = 0$, while if both u, v are off the real line $G(u, v) > 0$. Hence, for any configuration that involves a continuum including points on the real line, we can find a configuration with no points on the real line, by pushing measures up away from the real axis, which has greater (unweighted *and* weighted) energy. So, suppose the maximizing continuum touches the axis. We can always push the measures up away from the real axis and end up with a continuum that has greater minimal energy, thus arriving at a contradiction.

The proposition is now proved.

Remark: It is also important to consider the point at infinity. We cannot prove that the continuum does not hit this point. [In fact, our numerics (Ref. 3, Chap. 6) show that it may well do so.] In connection with the semiclassical problem (9) and (10) as analyzed in Ref. 3, it might seem at first that the maximizing continuum should not pass through infinity. Indeed, the transformations (2.17) and (4.1) of Ref. 3 implicitly assume that the continuum C lies in \mathbb{C} . Otherwise, one would lose the appropriate normalization for M at infinity. However, one must simply notice that infinity is just an arbitrary choice of normalizing point, once we view our Riemann-Hilbert problems in the compact Riemann sphere. The important observation is that the composition of transformations (2.17) and (4.1) (which are purely formal, i.e., no estimates are required and no approximation is needed) does not introduce any bad (essential) singularities. In the end, the asymptotic behavior of \tilde{N}^σ is still the identity as $z \rightarrow \infty$ in the lower half-plane and nonsingular as $z \rightarrow \infty$ in the upper half-plane. So, in the end it is acceptable for a continuum to go through the point infinity.

VI. TAKING SMALL VARIATIONS

We now complexify the external field and extend it to a function in the whole complex plane, by turning a Green's potential to a logarithmic potential. We will thus be able to make direct use of the results of Ref. 5.

We let, for any complex z ,

$$V(z) = - \int_{-iA}^{iA} \log(z - \eta) \rho^0(\eta) d\eta - \left(2ixz + 2itz^2 + i\pi \int_z^{iA} \rho^0(\eta) d\eta \right) \quad (20)$$

and $V_R = \operatorname{Re} V$ be the real part of V . In the lower half-plane the function ρ^0 is extended simply by

$$\rho^0(\eta^*) = (\rho^0(\eta))^* \quad (20')$$

Note right away that the field ϕ defined in (3') is the restriction of $V_R = \operatorname{Re} V$ to the closed upper half-plane.

The actual contour of the logarithmic integral is chosen to be the linear segment χ joining the points $-iA, 0, iA$. The branch of the logarithm function $\log(z - \eta)$ is then defined to agree with the principal branch as $z \rightarrow \infty$, and with jump across the very contour χ .

The unweighted Green's energy (4) can be written as

$$E_{V_R}(\mu) = \int_{\operatorname{supp} \mu} \int_{\operatorname{supp} \mu} \log \frac{1}{|u - v|} d\mu(u) d\mu(v) + 2 \int_{\operatorname{supp} \mu} V_R(u) d\mu(u), \quad (21)$$

where the measures μ are extended to the lower half complex plane by

$$\mu(z^*) = -\mu(z). \quad (21')$$

(So they are "signed" measures.)

Having established in Sec. V that the contour solving the variational problem does not touch the real line we can take small variations of measures and contours, never intersecting the real line, and keeping the points $0_-, 0_+$ fixed. In view of (21') we can think of them as variations of measures symmetric under (21') in the full complex plane, never intersecting the real line, and

keeping the points $0_-, 0_+$ fixed. The perturbed measures do not change sign. The fact that ∞ can belong to the contour is not a problem. Our variations will keep it automatically fixed.

The first step is to show that the solution of the variational problem satisfies a crucial relation.

Remark: It is not hard to see that the variational problem of Theorem 4 is actually *equivalent* to the variational problem of maximizing equilibrium measures on continua in the whole complex plane, under the symmetry (21') and the condition that measures are to be positive in the upper half-plane and negative in the lower half-plane.

Theorem 5: Let F be the maximizing continuum of Theorem 4 and λ^F be the equilibrium measure minimizing the weighted logarithmic energy (6) under the external field $V_R = \text{Re } V$ where V is given by (20). Let μ be the extension of λ^F to the lower complex plane via $\mu(z^*) = -\mu(z)$. Then

$$\begin{aligned} \text{Re} \left(\int \frac{d\mu(u)}{u-z} + V'(z) \right)^2 &= \text{Re}(V'(z))^2 - 2 \text{Re} \int \frac{V'(z) - V'(u)}{z-u} d\mu(u) \\ &\quad + \text{Re} \left(\frac{1}{z^2} \int 2(u+z)V'(u) d\mu(u) \right). \end{aligned} \tag{22}$$

Proof: We first need to prove the following.

Theorem 6: Let Γ be a critical point of the functional taking a continuum $\Gamma \in \mathbb{F}$ to $E_{V_R}(\lambda^\Gamma)$, and assume that Γ is not tangent to \mathbb{R} . Also assume that Γ does not touch the segment $[0, iA]$ except at zero. Let μ be the extension of λ^Γ via $\mu(z^*) = -\mu(z)$, O_Γ be an open set containing the interior of $\Gamma \cup \Gamma^*$ and $h \in C^1(O_\Gamma)$ such that $h(0) = 0$. We have

$$\text{Re} \left(\int \int \frac{h(u) - h(v)}{u-v} d\mu(u) d\mu(v) \right) = 2 \text{Re} \left(\int V'(u) h(u) d\mu(u) \right). \tag{23}$$

Proof: Consider the family of (signed) measures $\{\mu^\tau, \tau \in \mathbb{C}, |\tau| < \tau_0\}$ defined by $d\mu^\tau(z^\tau) = d\mu(z)$ where $z^\tau = z + \tau h(z)$, or equivalently, $\int f(z) d\mu^\tau(z) = \int f(z) d\mu(z)$, $f \in L_1(O_\Gamma)$. Assume that τ is small enough so that the support of the deformed continuum does not hit the linear segment $(0, iA]$ or a nonzero point in the real line.

With $\hat{h} = \hat{h}(u, v) = [h(u) - h(v)] / (u - v)$, we have $(u^\tau - v^\tau) / (u - v) = 1 + \tau \hat{h}$, so that $\log(1/|u^\tau - v^\tau|) - \log(1/|u - v|) = -\log|1 + \tau \hat{h}| = -\text{Re}(\tau \hat{h}) + O(\tau^2)$.

Integrating with respect to $d\mu(u) d\mu(v)$ we arrive at

$$E(\mu^\tau) - E(\mu) = -\text{Re} \left(\tau \int \hat{h} d\mu(u) d\mu(v) \right) + O(\tau^2), \tag{24}$$

where $E(\mu)$ denotes the free logarithmic energy of the measure μ . Also,

$$\int V_R d\mu^\tau - \int V_R d\mu = 2 \int (V_R(u^\tau) - V_R(u)) d\mu(u) = 2 \text{Re} \left(\tau \int V'(u) h(u) d\mu(u) \right) + O(\tau^2).$$

Combining with the above,

$$E_{V_R}(\mu^\tau) - E_{V_R}(\mu) = \text{Re} \left(-\tau \int \hat{h} d\mu(u) d\mu(v) + 2\tau \int V' h d\mu \right) + O(\tau^2). \tag{25}$$

So, if μ is (the symmetric extension of) a critical point of the map $\mu \rightarrow E_{V_R}(\mu)$ the linear part of the increment is zero. In other words given a C^1 function h and a measure μ the function $E_{V_R}(\mu^\tau)$ of τ is differentiable at $\tau = 0$ and the derivative is

$$\text{Re}(-H(\mu)), \quad \text{where } H(\mu) = \int \int \hat{h} d\mu^2 - 2 \int V' h d\mu. \tag{26}$$

But what we really want is the derivative of the energy as a function of the equilibrium measure. This function can be shown to be differentiable and its derivative can be set to zero at a critical continuum.

Indeed, we need to show the following.

Lemma 4:

$$\frac{d}{d\tau} E_{V_R}((\lambda^\Gamma)^\tau)|_{\tau=0} = \frac{d}{d\tau} E_{V_R}(\lambda^{\Gamma^\tau})|_{\tau=0} = 0.$$

In the relation above $\Gamma_\tau = \text{supp}(\lambda^{\Gamma^\tau})$. The first derivative is of a function of general measures. The second derivative is of a function of equilibrium measures.

Proof: Define the measure σ_τ with support Γ and such that $(\sigma_\tau)^\tau = \lambda^{\Gamma^\tau}$.

LEMMA 5: With H defined by (26), we have

$$\text{Re } H(\sigma_\tau) \rightarrow \text{Re } H(\lambda^\Gamma),$$

as $\tau \rightarrow 0$.

Proof: By (25) and (26), we have

$$E_{V_R}((\lambda^\Gamma)^\tau) - E_{V_R}(\lambda^\Gamma) = -\text{Re}(\tau H(\lambda^\Gamma) + O(\tau^2)),$$

$$E_{V_R}(\lambda^{\Gamma^\tau}) - E_{V_R}(\sigma_\tau) = -\text{Re}(\tau H(\sigma_\tau) + O(\tau^2)).$$

On the other hand, $E_{V_R}(\sigma_\tau) \geq E_{V_R}(\lambda^\Gamma)$, and $E_{V_R}((\lambda^\Gamma)^\tau) \geq E_{V_R}(\lambda^{\Gamma^\tau})$. It follows that

$$E_{V_R}(\sigma_\tau) - \text{Re}(\tau H(\sigma_\tau)) + O(\tau^2) = E_{V_R}(\lambda^{\Gamma^\tau}) \leq E_{V_R}((\lambda^\Gamma)^\tau) = E_{V_R}(\lambda^\Gamma) - \text{Re}(\tau H(\lambda^\Gamma)) + O(\tau^2).$$

Hence $E_{V_R}(\sigma_\tau) \rightarrow E_{V_R}(\lambda^\Gamma)$.

As in the proof of Theorem 1, it follows that $\sigma_\tau \rightarrow \lambda^\Gamma$ weakly; see Ref. 7, pp. 82–88. It then follows immediately that $H(\sigma_\tau) \rightarrow H(\lambda^\Gamma)$. This proves Lemma 5.

To complete the proof of Lemma 4, we note that $0 \geq E_{V_R}((\lambda^\Gamma)^\tau) - E_{V_R}(\lambda^\Gamma) \geq E_{V_R}(\lambda^{\Gamma^\tau}) - E_{V_R}(\sigma_\tau) = -\text{Re}(\tau H(\sigma_\tau)) + O(\tau^2)$. Hence the derivative of $E_{V_R}((\lambda^\Gamma)^\tau)$ at $\tau=0$ is equal to the derivative of $E_{V_R}(\lambda^{\Gamma^\tau})$ at $\tau=0$ which is equal to $\text{Re } H(\lambda^\Gamma)$. This proves Lemma 4 and Theorem 6.

Proof of Theorem 5: Consider the Schiffer variation, i.e., take $h(u) = u^2/(u-z)$ where z is some fixed point not in Γ . Note that $h(0) = 0$ so that the deformation $z^\tau = z + \tau h(z)$ keeps the points $0_+, 0_-$ fixed. Also assume that τ is small enough so that the support of the deformed continuum does not hit the linear segment $(0, iA]$ or a nonzero point in the real line. We have

$$\hat{h} = \hat{h}(u, v) = \frac{h(u) - h(v)}{u - v} = 1 - \frac{z^2}{(u-z)(v-z)},$$

and therefore

$$\text{Re} \left(\int \int \hat{h}(u, v) d\mu(u) d\mu(v) \right) = \text{Re} \left[\int \int du(u) d\mu(v) - z^2 \left(\int_{\text{supp } \mu} \frac{d\mu(u)}{u-z} \right)^2 \right].$$

Next

$$\begin{aligned} \operatorname{Re}\left(\int 2V'(u)h(u)d\mu(u)\right) &= 2 \operatorname{Re}\left(\int (u+z)V'(u)d\mu(u) + z^2 \int \frac{V'(u)d\mu(u)}{u-z}\right) \\ &= \operatorname{Re}\left(\int 2V'(u)(u+z)d\mu(u) + 2z^2 \int \frac{V'(u)-V'(z)}{u-z}d\mu(u)\right. \\ &\quad \left.+ 2z^2V'(z) \int \frac{d\mu(u)}{u-z}\right). \end{aligned}$$

Theorem 5 now follows from Theorem 6.

Remark: If our continuum is allowed to touch the point iA [so we slightly weaken assumption (A)] then we may need to keep points $\pm iA$ fixed under a small variation. We can then choose the Schiffer variation $h(u)=u^2(u^2+A^2)/(u-z)$. We will arrive at a similar and equally useful formula.

In general if one wants to keep points a_1, \dots, a_s fixed, the appropriate Schiffer variation is $h(u)=\prod_{i=1}^s(u-a_i)/(u-z)$.

Proposition 3: The support of the equilibrium measure consists of a finite number of analytic arcs.

Proof: Theorem 5 above implies that the support of μ is the level set of the real part of a function that is analytic except at countably many branch points. In fact, $\operatorname{supp} \mu$ is characterized by $\int \log(1/|u-z|)d\mu(u) + V_R(z) = 0$. From Theorem 5 we get

$$\operatorname{Re}\left(\int \frac{d\mu(u)}{u-z} + V'(z)\right) = \operatorname{Re}[(R_\mu(z))^{1/2}], \quad (27)$$

where

$$R_\mu(z) = (V'(z))^2 - 2 \int_{\operatorname{supp} \mu} \frac{V'(z)-V'(u)}{z-u}d\mu(u) + \frac{1}{z^2} \left(\int_{\operatorname{supp} \mu} 2(u+z)V'(u)d\mu(u) \right). \quad (28)$$

This is a function analytic in K , with possibly a pole at zero. By integrating, we have that $\operatorname{supp}(\mu)$ is characterized by

$$\operatorname{Re} \int^z (R_\mu)^{1/2} dz = 0. \quad (29)$$

The locus defined by (29) is a union of arcs with endpoints at zeros of R_μ . To see this, consider (29) as an equation of two real variables $f(u, v) = 0$ and try to solve for u as an analytic expression of v . One can only do this if the derivative of f is nonzero, which means (via the Cauchy-Riemann equations) that $R_\mu \neq 0$. The points where $R_\mu = 0$ are exactly the points where the analytic arcs cannot be continued.

Note that

$$R_\mu(z) \sim -[16t^2z^2 + \pi^2(\rho^0(z))^2], \quad \text{as } z \rightarrow \infty, \quad (30)$$

$$R_\mu(z) \sim \frac{1}{z^2} \int 2uV'(u)d\mu(u), \quad \text{as } z \rightarrow 0.$$

By conditions (1) for ρ^0 , R_μ is blowing up at the point ∞ (at least for $t > 0$; but the case $t = 0$ is well understood: the equilibrium measure consists of a single analytic arc; see Sec. V). Hence it can only have finitely many zeros near infinity, otherwise they would have to accumulate near ∞ and then R_μ would be 0 there. On the other hand, R_μ cannot have an accumulation point of zeros at $z=0$, because even if the pole at 0 were removed (the coefficients of $1/z^2$, $1/z$ being zero), R_μ would be holomorphically extended across $z=0$. So, R_μ can only have a finite number of zeros in \bar{K} . It follows that the support of the maximizing equilibrium measure consists of only finitely many arcs.

REMARK: Of course, conditions (1) can be weakened. We could allow ρ^0 to have a pole at infinity of order other than two. But our aim here is not to prove the most general theorem possible, but instead illustrate a method that can be applied in the most general settings under appropriate amendments.

Remark: The assumption that ρ^0 is continuous and hence bounded at infinity is only needed to prove the finiteness of the components of the support of the equilibrium measure of the maximizing continuum. If it is dropped then we may have an infinite number of components for isolated values of x, t . This will result in infinite genus representations of the semiclassical asymptotics. Of course infinite genus solutions of the focusing NLS equation are known and well understood. So the analysis of Ref. 3 is expected to also apply in that case, although it will be more tedious.

For a justification of the “finite gap ansatz,” concerning the semiclassical limit of focusing NLS, it only remains to verify the “S-property.”

VII. THE S-PROPERTY

Theorem 7: (The S-property.)

Let C be the contour maximizing the equilibrium energy, for the field given by (3') with conditions (1). Let μ be the extension of its equilibrium measure to the full complex plane via (21'). Assume for simplicity that assumption (A) holds. Let $X(z) = \int_{\text{supp } \mu} \log[1/(u-z)] d\mu(u)$, $X_R(z) = \text{Re } X(z) = \int_{\text{supp } \mu} \log(1/|u-z|) d\mu(u)$, $W_\mu = X'$. Then, at any interior point of $\text{supp } \mu$ other than zero,

$$\frac{d}{dn_+}(V_R + X_R) = \frac{d}{dn_-}(V_R + X_R), \quad (8')$$

where the two derivatives above denote the normal derivatives, on the + and - sides, respectively.

Proof: From Theorem 5, we have

$$|\text{Re}(W_\mu(z) + V'(z))| = |\text{Re}(R_\mu)|^{1/2}.$$

Using the definition for X , the above relation becomes

$$\left| \frac{d}{dz} \text{Re}(X + V) \right| = |\text{Re}(R_\mu)|^{1/2}.$$

Now, $\text{Re}(X+V)=0$ on the support of the equilibrium measure. So, in particular $\text{Re}(X+V)$ is constant along the equilibrium measure. Hence $|(d/dz)\text{Re}(X+V)|$ must be equal to the modulus of *each* normal derivative across the equilibrium measure. So,

$$\left| \frac{d}{dn_\pm}(V_R + \text{Re } X) \right| = \left| \frac{d}{dz} \text{Re}(X + V) \right| = |\text{Re}(R_\mu)|^{1/2}.$$

Hence,

$$\left| \frac{d}{dn_+}(V_R + \text{Re } X) \right| = \left| \frac{d}{dn_-}(V_R + \text{Re } X) \right|.$$

But it is easy to see that both LHS and RHS quantities inside the modulus sign are negative. This is because $V_R + \text{Re } X = 0$ on $\text{supp } \mu$ and negative on each side of $\text{supp } \mu$. Hence result.

Remark: Once Theorem 7 is proved it follows by the Cauchy-Riemann equations that $(V_I + \text{Im } X)_+ + (V_I + \text{Im } X)_-$ is constant on each connected component of $\text{supp } \mu$, which means that $\text{Im } \tilde{\phi}$ is constant on connected components of the contour, where $\tilde{\phi}$ is as defined in formula (4.13) of Ref. 3. This proves the existence of the appropriate “g functions” in Ref. 3.

We recapitulate our results in the following theorem, set in the upper complex half-plane. Note that (8') is the “doubled up” version of (8).

Theorem 8: Let ϕ be given by (3'), where ρ^0 satisfies conditions (1). Under assumption (A), there is a piecewise smooth contour $C \in \mathbb{F}$, containing points $0_+, 0_-$ and otherwise lying in the cut upper half-plane \mathbb{K} , with equilibrium measure λ^C , such that $\text{supp}(\lambda^C)$ consists of a union of finitely many analytic arcs and

$$E_\phi(\lambda^C) = \max_{C' \in \mathbb{F}} E_\phi(\lambda^{C'}) = \max_{C' \in \mathbb{F}} [\inf_{\mu \in M(\mathbb{F})} E_\phi(\mu)].$$

On each interior point of $\text{supp}(\lambda^C)$ we have

$$\frac{d}{dn_+}(\phi + V^{\lambda^C}) = \frac{d}{dn_-}(\phi + V^{\lambda^C}), \tag{8''}$$

where V^{λ^C} is the Green's potential of the equilibrium measure λ^C [see (5)] and the two derivatives above are the normal derivatives.

A curve satisfying (8) such that the support of its equilibrium measure consists of a union of finitely many analytic arcs is called an S-curve.

Proof: The fact that the maximizing continuum C is actually a contour is proved as follows. If this were not the case, then we could choose a subset of C , say F , which is a contour, starting at 0_+ and ending at 0_- , and going around the point iA . Clearly, by definition, the equilibrium energy of C is less than the equilibrium energy of F , i.e., $E_\phi(\lambda^C) \leq E_\phi(\lambda^F)$. On the other hand, since C maximizes the equilibrium energy, we have $E_\phi(\lambda^F) \leq E_\phi(\lambda^C)$. So $E_\phi(\lambda^F) = E_\phi(\lambda^C)$.

VIII. CONCLUSION

In view of the interpretation of the variational problem in terms of the semiclassical NLS problem, we have the following result.

Consider the semiclassical limit ($\hbar \rightarrow 0$) of the solution of (9) and (10) with bell-shaped initial data. Replace the initial data by the so-called soliton ensembles data (as introduced in Ref. 3) defined by replacing the scattering data for $\psi(x, 0) = \psi_0(x)$ by their WKB approximation, so that the spectral density of eigenvalues is

$$d\mu_0^{\text{WKB}}(\eta) := \rho^0(\eta)\chi_{[0, iA]}(\eta)d\eta + \rho^0(\eta^*)\chi_{[-iA, 0]}(\eta)d\eta,$$

$$\text{with } \rho^0(\eta) := \frac{\eta}{\pi} \int_{x_-(\eta)}^{x_+(\eta)} \frac{dx}{\sqrt{A(x)^2 + \eta^2}} = \frac{1}{\pi} \frac{d}{d\eta} \int_{x_-(\eta)}^{x_+(\eta)} \sqrt{A(x)^2 + \eta^2} d\mu,$$

for $\eta \in (0, iA)$, where $x_-(\eta) < x_+(\eta)$ are the two real turning points, i.e., $(A(x_\pm))^2 + \eta^2 = 0$, the square root is positive and the imaginary segments $(-iA, 0)$ and $(0, iA)$ are both considered to be oriented from bottom to top to define the differential $d\eta$.

Assume that ρ^0 satisfies conditions (1). Then, under assumption (A), asymptotically as $\hbar \rightarrow 0$, the solution $\psi(x, t)$ admits a "finite genus description." (For a more precise explanation, see Appendix B.)

The proof of this is the main result of Ref. 3, assuming that the variational problem of Sec. I has an S-curve as a solution. But this is now guaranteed by Theorem 8.

Remark: For conditions weaker than the above, the particular spectral density ρ^0 arising in the semiclassical NLS problem can conceivably admit branch singularities in the upper complex plane and condition (1) will not be satisfied. We claim that even in such a case the finite gap genus can be justified, at least generically. The proof of this fact will require setting the variational problem on a Riemann surface with moduli at the branch singularities of ρ^0 .

Remark: Consider the semiclassical problem (9) and (10) in the case of initial data $\psi_0(x) = A \text{sech } x$, where $A > 0$. Then the WKB density is given by $\rho^0 = i$ [see (3.1) and (3.2) of Ref. 3; note that condition (1) is satisfied]. So the finite genus ansatz holds for any x, t , as long as the assumption (A) of Sec. V holds. But then assumption (A) can be eventually dropped; see Appendix C.

Remark: The behavior of a solution of (9) in general depends not only on the eigenvalues of the Lax operator, but also on the associated norming constants and the reflection coefficient. In the special case of the soliton ensembles data the norming constants alternate between -1 and 1 while the reflection coefficient is by definition zero. More generally, for real analytic data decaying at infinity the reflection coefficient is exponentially small everywhere except at zero and can be neglected (although the rigorous proof of this is not trivial).

ACKNOWLEDGMENTS

The first author acknowledges the kind support of the General Secretariat of Research and Technology, Greece, in particular Grant No. 97EL16. He is also grateful to the Department of Mathematics of the University of South Florida for its hospitality during a visit on May 2001 and to the Max Planck Society for support since 2002. Both authors acknowledge the invaluable contribution of collaborators Ken McLaughlin and Peter Miller through stimulating discussions, important comments, corrections, and constructive criticism.

APPENDIX A: COMPACTNESS OF THE SET OF CONTINUA

In this section we prove that the sets $I(\bar{\mathbb{K}})$ and hence \mathbb{F} defined in Sec. I are compact and complete.

As stated in Sec. I, the space we must work with is the upper half-plane: $\mathbb{H} = \{z : \text{Im } z > 0\}$. The closure of this space is $\bar{\mathbb{H}} = \{z : \text{Im } z \geq 0\} \cup \{\infty\}$. Also $\mathbb{K} = \{z : \text{Im } z > 0\} \setminus \{z : \text{Re } z = 0, 0 < \text{Im } z \leq A\}$. In the closure of this space, $\bar{\mathbb{K}}$, we consider the points ix_+ and ix_- , where $0 \leq x < A$ as distinct.

Even though we eventually wish to consider only smooth contours, we are forced to *a priori* work with general closed sets. The reason is that the set of contours is not compact in any reasonable way, so it seems impossible to prove any existence theorem for a variational problem defined only on contours. Instead, we define \mathbb{F} to be the set of all “continua” F in $\bar{\mathbb{K}}$ (i.e., connected compact sets, containing the points $0_+, 0_-$).

Furthermore, we need to introduce an appropriate topology on \mathbb{F} , that will make it a compact set. In this we follow the discussion of Dieudonné (Ref. 6, Chap. III.16).

We think of the closed upper half-plane $\bar{\mathbb{H}}$ as a compact space in the Riemann sphere. We thus choose to equip $\bar{\mathbb{H}}$ with the “chordal” distance, denoted by $d_0(z, \zeta)$, that is the distance between the images of z and ζ under the stereographic projection. This induces naturally a distance in $\bar{\mathbb{K}}$ [so, for example, $d_0(0_+, 0_-) \neq 0$]. We also denote by d_0 the induced distance between compact sets E, F in $\bar{\mathbb{K}}$: $d_0(E, F) = \max_{z \in E} \min_{\zeta \in F} d_0(z, \zeta)$. Then, we define the so-called Hausdorff metric on the set $I(\bar{\mathbb{K}})$ of closed nonempty subsets of $\bar{\mathbb{K}}$ as follows:

$$d_{\mathbb{K}}(A, B) = \sup(d_0(A, B), d_0(B, A)). \quad (\text{A1})$$

Lemma A.1: The Hausdorff metric defined by (A1) is indeed a metric. The set $I(\bar{\mathbb{K}})$ is compact and complete.

Proof: It is clear that $d_{\mathbb{K}}(A, B)$ is non-negative and symmetric by definition. Also if $d_{\mathbb{K}}(A, B) = 0$, then $d_0(A, B) = 0$, hence $\max_{z \in A} \min_{\zeta \in B} d_0(z, \zeta) = 0$ and thus for all $z \in A$, we have $\min_{\zeta \in B} d_0(z, \zeta) = 0$. In other words, $z \in B$. By symmetry, $A = B$.

The triangle inequality follows from the triangle inequality for d_0 . Indeed, suppose $A, B, C \in I(\bar{\mathbb{K}})$. Then $d_{\mathbb{K}}(A, B) = \sup(d_0(A, B), d_0(B, A)) = d_0(A, B)$, without loss of generality. Now,

$$d_0(A, B) = \max_{z \in A} \min_{\zeta \in B} d_0(z, \zeta) \leq \max_{z \in A} \min_{\zeta \in B} \min_{\zeta_0 \in C} (d_0(z, \zeta_0) + d_0(\zeta_0, \zeta)),$$

by the triangle inequality for d_0 . Let $z = z_0 \in A$ be the value of z that maximizes $\min_{\zeta \in B} \min_{\zeta_0 \in C} (d_0(z, \zeta_0) + d_0(\zeta_0, \zeta))$. This is then

$$\begin{aligned} \min_{\zeta \in B} \min_{\zeta_0 \in C} (d_0(z_0, \zeta_0) + d_0(\zeta_0, \zeta)) &\leq \min_{\zeta_0 \in C} d_0(z_0, \zeta_0) + \min_{\zeta \in B} \min_{\zeta_0 \in C} d_0(\zeta_0, \zeta) \\ &\leq \max_{z \in A} \min_{\zeta_0 \in C} d_0(z, \zeta_0) + \max_{\zeta \in B} \min_{\zeta_0 \in C} d_0(\zeta_0, \zeta) \\ &\leq d_0(A, C) + d_0(B, C) \leq d_{\mathbb{K}}(A, C) + d_{\mathbb{K}}(B, C). \end{aligned}$$

The result follows from symmetry.

We will next show that $I(\overline{\mathbb{K}})$ is complete and precompact. Since a precompact, complete metric space is compact [Ref. 6, proposition (3.16.1)] the proof of Lemma A.1 follows.

Lemma A.2: If the metric space \mathbb{E} equipped with a distance d_0 is complete, then so is $I(\mathbb{E})$, the set of closed nonempty subsets of \mathbb{E} , equipped with the Hausdorff distance

$$d_{\mathbb{E}}(A, B) = \sup(d_0(A, B), d_0(B, A)),$$

for any closed nonempty subsets A, B , where $d_0(A, B) = \max_{a \in A} \min_{b \in B} d_0(a, b)$.

Furthermore, if \mathbb{E} is precompact, then so is $I(\mathbb{E})$.

Proof: Suppose \mathbb{E} is complete. Let X_n be a Cauchy sequence in $I(\mathbb{E})$. We will show that X_n converges to $X = \bigcap_{n \geq 0} \overline{\bigcup_{p \geq 0} X_{n+p}}$. (Overbar denotes closure.)

Indeed, given any $\epsilon > 0$,

$$d_0(X_n, X) = \max_{x \in X_n} \min_{y \in X} d_0(x, y) \leq \max_{x \in X_n} \max_{y \in \overline{\bigcup_{p \geq 0} X_{n+p}}} d_0(x, y) < \epsilon,$$

for large n , by the completeness of \mathbb{E} . Similarly,

$$d_0(X, X_n) = \max_{x \in X} \min_{y \in X_n} d_0(x, y) \leq \max_{x \in \overline{\bigcup_{p \geq 0} X_{n+p}}} \min_{y \in X_n} d_0(x, y) < \epsilon.$$

Next, suppose \mathbb{E} is precompact. Then, by definition, given any $\epsilon > 0$, there is a finite set, say $S = \{s_1, s_2, \dots, s_n\}$, where n is a finite integer, such that any point x of \mathbb{E} is at a distance d_0 less than ϵ to the set S . Now, consider the set of subsets of S , which is of course finite. Clearly every closed set is at a distance less than ϵ to a member of that set,

$$d_0(A, S) = \max_{a \in A} \min_{s \in S} d_0(a, s) < \epsilon,$$

$$d_0(S, A) = \max_{s \in S} \min_{a \in A} d_0(a, s) < \epsilon,$$

for any closed nonempty set A . Hence $d_{\mathbb{E}}(A, S) < \epsilon$.

So, any closed nonempty set A is at a distance less than ϵ to the finite power set of S . So $I(\mathbb{E})$ is precompact.

APPENDIX B: THE DESCRIPTION OF THE SEMICLASSICAL LIMIT OF THE FOCUSING NLS EQUATION UNDER THE FINITE GENUS ANSATZ

We present one of the main results of Ref. 3 on the semiclassical asymptotics for problem (9) and (10), in view of the fact that the finite genus ansatz holds. In particular, we fix x, t and use the result that the support of the maximizing measure of Theorems 4 and 8 consists of a finite union of analytic arcs.

First, we define the so-called g function. Let C be the maximizing contour of Theorem 4. *A priori* we seek a function satisfying

$$g(\lambda) \text{ is independent of } \hbar,$$

$$g(\lambda) \text{ is analytic for } \lambda \in \mathbb{C} \setminus (C \cup C^*),$$

$$g(\lambda) \rightarrow 0 \text{ as } \lambda \rightarrow \infty,$$

$$g(\lambda) \text{ assumes continuous boundary values from both sides of } C \cup C^*,$$

denoted by $g_+(g_-)$ on the left (right) of $C \cup C^*$,

$$g(\lambda^*) + g(\lambda)^* = 0 \quad \text{for all } \lambda \in \mathbb{C} \setminus (C \cup C^*).$$

The assumptions above are satisfied if we write g in terms of the maximizing equilibrium measure of Theorem 8, $d\mu = d\lambda^C = \rho(\eta)d\eta$, doubled up according to (21'). Indeed,

$$g(\lambda) = \int_{C \cup C^*} \log(\lambda - \eta)\rho(\eta)d\eta,$$

for an appropriate definition of the logarithm branch (see Ref. 3).

For $\lambda \in C$, define the functions

$$\theta(\lambda) := i(g_+(\lambda) - g_-(\lambda)),$$

$$\begin{aligned} \Phi(\lambda) := & \int_0^{iA} \log(\lambda - \eta)\rho^0(\eta)d\eta + \int_{-iA}^0 \log(\lambda - \eta)\rho^0(\eta)^*d\eta \\ & + 2i\lambda x + 2i\lambda^2 t + i\pi \int_{\lambda}^{iA} \rho^0(\eta)d\eta - g_+(\lambda) - g_-(\lambda), \end{aligned}$$

where $\rho^0(\eta)$ is the holomorphic function (WKB density of eigenvalues) introduced in Sec. I [see conditions (1)].

The finite genus ansatz implies that for each x, t there is a finite positive integer G such that the contour C can be divided into “bands” [the support of $\rho(\eta)d\eta$] and “gaps” (where $\rho=0$). We denote these bands by I_j . More precisely, we define the analytic arcs $I_j, I_j^*, j=1, \dots, G/2$ as follows (they come in conjugate pairs). Let the points $\lambda_j, j=0, \dots, G$, in the open upper half-plane be the branch points of the function g . All such points lie on the contour C and we order them as $\lambda_0, \lambda_1, \dots, \lambda_G$, according to the direction given to C . The points $\lambda_0^*, \lambda_1^*, \dots, \lambda_G^*$ are their complex conjugates. Then let $I_0=[0, \lambda_0]$ be the subarc of C joining points 0 and λ_0 . Similarly, $I_j=[\lambda_{2j-1}, \lambda_{2j}]$, $j=1, \dots, G/2$. The connected components of the set $\mathbb{C} \setminus \cup_j (I_j \cup I_j^*)$ are the so-called “gaps,” for example, the gap Γ_1 joins λ_0 to λ_1 , etc.

It actually follows from the properties of g, ρ that the function $\theta(\lambda)$ defined on C is constant on each of the gaps Γ_j , taking a value which we will denote by θ_j , while the function Φ is constant on each of the bands, taking the value denoted by α_j on the band I_j .

The finite genus ansatz for the given fixed x, t implies that the asymptotics of the solution of (9) and (10) as $\hbar \rightarrow 0$ can be given by the next theorem.

Theorem A.1: Let x_0, t_0 be given. The solution $\psi(x, t)$ of (9) and (10) is asymptotically described (locally) as a slowly modulated $G+1$ phase wavetrain. Setting $x=x_0+\hbar\hat{x}$ and $t=t_0+\hbar\hat{t}$, so that x_0, t_0 are “slow” variables while \hat{x}, \hat{t} are “fast” variables, there exist parameters

$a, U=(U_0, U_1, \dots, U_G)^T, k=(k_0, k_1, \dots, k_G)^T, w=(w_0, w_1, \dots, w_G)^T, Y=(Y_0, Y_1, \dots, Y_G)^T, Z=(Z_0, Z_1, \dots, Z_G)^T$ [depending on the slow variables x_0 and t_0 (but not \hat{x}, \hat{t}) such that

$$\begin{aligned} \psi(x, t) = & \psi(x_0 + \hbar\hat{x}, t_0 + \hbar\hat{t}) \sim a(x_0, t_0) e^{iU_0(x_0, t_0)/\hbar} e^{i(k_0(x_0, t_0)\hat{x} - w_0(x_0, t_0)t)} \\ & \times \frac{\Theta(Y(x_0, t_0) + iU(x_0, t_0)/\hbar + i(k(x_0, t_0)\hat{x} - w(x_0, t_0)t))}{\Theta(Z(x_0, t_0) + iU(x_0, t_0)/\hbar + i(k(x_0, t_0)\hat{x} - w(x_0, t_0)t))}. \end{aligned} \tag{B1}$$

All parameters can be defined in terms of an underlying Riemann surface X . The moduli of X are given by $\lambda_j, j=0, \dots, G$ and their complex conjugates $\lambda_j^*, j=0, \dots, G$. The genus of X is G . The moduli of X vary slowly with x, t , i.e., they depend on x_0, t_0 but not \hat{x}, \hat{t} . For the exact formulas for the parameters as well as the definition of the theta functions we present the following construction.

The Riemann surface X is constructed by cutting two copies of the complex sphere along the slits $I_0 \cup I_0^*, I_j, I_j^*, j=1, \dots, G$, and pasting the “top” copy to the “bottom” copy along these very slits.

We define the homology cycles $a_j, b_j, j=1, \dots, G$ as follows. Cycle a_1 goes around the slit $I_0 \cup I_0^*$ joining λ_0 to λ_0^* , remaining on the top sheet, oriented counterclockwise, a_2 goes through the slits I_{-1} and I_1 starting from the top sheet, also oriented counterclockwise, a_3 goes around the slits $I_{-1}, I_0 \cup I_0^*, I_1$ remaining on the top sheet, oriented counterclockwise, etc. Cycle b_1 goes through I_0 and I_1 oriented counterclockwise, cycle b_2 goes through I_{-1} and I_1 , also oriented counterclockwise, cycle b_3 goes through I_{-1} and I_2 , and around the slits $I_{-1}, I_0 \cup I_0^*, I_1$, oriented counterclockwise, etc.

On X there is a complex G -dimensional linear space of holomorphic differentials, with basis elements $\nu_k(P)$ for $k=1, \dots, G$ that can be written in the form

$$\nu_k(P) = \frac{\sum_{j=0}^{G-1} c_{kj} \lambda(P)^j}{R_X(P)} d\lambda(P),$$

where $R_X(P)$ is a “lifting” of the function $R(\lambda)$ from the cut plane to X , if P is on the first sheet of X then $R_X(P) = R(\lambda(P))$ and if P is on the second sheet of X then $R_X(P) = -R(\lambda(P))$. The coefficients c_{kj} are uniquely determined by the constraint that the differentials satisfy the normalization conditions

$$\oint_{a_j} \nu_k(P) = 2\pi i \delta_{jk}.$$

From the normalized differentials, one defines a $G \times G$ matrix H (the period matrix) by the formula

$$H_{jk} = \oint_{b_j} \nu_k(P).$$

It is a consequence of the standard theory of Riemann surfaces that H is a symmetric matrix whose real part is negative definite.

In particular, we can define the theta function

$$\Theta(w) := \sum_{n \in \mathbb{Z}^G} \exp\left(\frac{1}{2} n^T H n + n^T w\right),$$

where H is the period matrix associated to X . Since the real part of H is negative definite, the series converges.

We arbitrarily fix a base point P_0 on X . The Abel map $A: X \rightarrow \text{Jac}(X)$ is then defined componentwise as follows:

$$A_k(P; P_0) := \int_{P_0}^P \nu_k(P'), \quad k=1, \dots, G,$$

where P' is an integration variable.

A particularly important element of the Jacobian is the Riemann constant vector K which is defined, modulo the lattice Λ , componentwise by

$$K_k := \pi i + \frac{H_{kk}}{2} - \frac{1}{2\pi i} \sum_{\substack{j=1 \\ j \neq k}}^G \oint_{a_j} \left(\nu_j(P) \int_{P_0}^P \nu_k(P') \right),$$

where the index k varies between 1 and G .

Next, we will need to define a certain meromorphic differential on X . Let $\Omega(P)$ be holomorphic away from the points ∞_1 and ∞_2 , where it has the behavior

$$\Omega(P) = dp(\lambda(P)) + \left(\frac{d\lambda(P)}{\lambda(P)^2} \right), \quad P \rightarrow \infty_1,$$

$$\Omega(P) = -dp(\lambda(P)) + O\left(\frac{d\lambda(P)}{\lambda(P)^2} \right), \quad P \rightarrow \infty_2,$$

and made unique by the normalization conditions

$$\oint_{a_j} \Omega(P) = 0, \quad j = 1, \dots, G.$$

Here p is a polynomial, defined as follows.

First, let us introduce the function $R(\lambda)$ defined by

$$R(\lambda)^2 = \prod_{k=0}^G (\lambda - \lambda_k)(\lambda - \lambda_k^*),$$

choosing the particular branch that is cut along the bands I_k^+ and I_k^- and satisfies

$$\lim_{\lambda \rightarrow \infty} \frac{R(\lambda)}{\lambda^{G+1}} = -1.$$

This defines a real function, i.e., one that satisfies $R(\lambda^*) = R(\lambda)^*$. At the bands, we have $R_+(\lambda) = -R_-(\lambda)$, while $R(\lambda)$ is analytic in the gaps. Next, let us introduce the function $k(\lambda)$ defined by

$$k(\lambda) = \frac{1}{2\pi i} \sum_{n=1}^{G/2} \theta_n \int_{\Gamma_n^+ \cup \Gamma_n^-} \frac{d\eta}{(\lambda - \eta)R(\eta)} + \frac{1}{2\pi i} \sum_{n=0}^{G/2} \int_{I_n^+ \cup I_n^-} \frac{\alpha_n d\eta}{(\lambda - \eta)R_+(\eta)}.$$

Next let

$$H(\lambda) = k(\lambda)R(\lambda).$$

The function k satisfies the jump relations

$$k_+(\lambda) - k_-(\lambda) = -\frac{\theta_n}{R(\lambda)}, \quad \lambda \in \Gamma_n^+ \cup \Gamma_n^-,$$

$$k_+(\lambda) - k_-(\lambda) = -\frac{\alpha_n}{R_+(\lambda)}, \quad \lambda \in I_n^+ \cup I_n^-,$$

and is otherwise analytic. It blows up like $(\lambda - \lambda_n)^{-1/2}$ near each endpoint, has continuous boundary values in between the endpoints, and vanishes like $1/\lambda$ for large λ . It is the only such solution of the jump relations. The factor of $R(\lambda)$ renormalizes the singularities at the endpoints, so that, as desired, the boundary values of $H(\lambda)$ are bounded continuous functions. Near infinity, there is the asymptotic expansion

$$H(\lambda) = H_G \lambda^G + H_{G-1} \lambda^{G-1} + \dots + H_1 \lambda + H_0 + O(\lambda^{-1}) = p(\lambda) + O(\lambda^{-1}), \quad (\text{B2})$$

where all coefficients H_j of the polynomial $p(\lambda)$ can be found explicitly by expanding $R(\lambda)$ and the Cauchy integral $k(\lambda)$ for large λ . It is easy to see from the reality of θ_j and α_j that $p(\lambda)$ is a polynomial with real coefficients.

Thus the polynomial $p(\lambda)$ is defined and hence the meromorphic differential $\Omega(P)$ is defined.

Let the vector $U \in \mathbb{C}^G$ be defined componentwise by

$$U_j := \oint_{b_j} \Omega(P).$$

Note that $\Omega(P)$ has no residues.

Let the vectors V_1, V_2 be defined componentwise by

$$V_{1,k} = (A_k(\lambda_{1+}^*) + A_k(\lambda_{2+}) + A_k(\lambda_{3+}^*) + \cdots + A_k(\lambda_{G+})) + A_k(\infty) + \pi i + \frac{H_{kk}}{2},$$

$$V_{2,k} = (A_k(\lambda_{1+}^*) + A_k(\lambda_{2+}) + A_k(\lambda_{3+}^*) + \cdots + A_k(\lambda_{G+})) - A_k(\infty) + \pi i + \frac{H_{kk}}{2},$$

where $k=1, \dots, G$, and the $+$ index means that the integral for A is to be taken on the first sheet of X , with base point λ_+^0 .

Finally, let

$$a = \frac{\Theta(Z)}{\Theta(Y)} \sum_{k=0}^G (-1)^k \mathfrak{J}(\lambda_k),$$

$$k_n = \partial_x U_n, \quad w_n = -\partial_t U_n, \quad n = 0, \dots, G,$$

where

$$Y = -A(\infty) - V_1, \quad Z = A(\infty) - V_1,$$

and $U_0 = -(\theta_1 + \alpha_0)$ where θ_1 is the (constant in λ) value of the function θ in the gap Γ_1 and α_0 is the (constant) value of the function ϕ in the band I_0 .

Now, the parameters appearing in formula (B1) are completely described.

We simply note here that the U_i and hence the k_i and w_i are real. We also note that the denominator in (B1) never vanishes (for any $x_0, t_0, \hat{x}, \hat{t}$).

Remark: Theorem A.1 presents pointwise asymptotics in x, t . In Ref. 3, these are extended to uniform asymptotics in certain compact sets covering the x, t plane. Error estimates are also given in Ref. 3.

Remark: As mentioned above, we do not know if the support of the equilibrium measure of the maximizing continuum is unique. But the asymptotic formula (B1) depends only on the endpoints λ_j of the analytic subarcs of the support. Since the asymptotic expression (B1) must be unique, it is easy to see that the endpoints also must be unique. Different Riemann surfaces give different formulas (except of course in degenerate cases, a degenerate genus 2 surface can be a pinched genus 0 surface and so on).

APPENDIX C: DROPPING ASSUMPTION (A) OF SEC. V

In Sec. V, we have assumed that the solution of the problem of the maximization of the equilibrium energy is a continuum, say F , which does not intersect the linear segment $[0, iA]$ except of course at $0_+, 0_-$. We also prove that F does not touch the real line, except of course at 0 and possibly ∞ . This enables us to take variations in Sec. VI, keeping fixed a finite number of points, and thus arrive at the identity of Theorem 5, from which we derive the regularity of F and the fact that F is, after all, an S-curve.

In general, it is conceivable that F intersects the linear segment $[0, iA]$ at points other than $0_+, 0_-$. If the set of such points is finite, there is no problem, since we can always consider variations keeping fixed a finite number of points, and arrive at the same result (see the remark after the proof of Theorem 5).

If, on the other hand, this is not the case, we have a different kind of problem, because the function V introduced in Sec. VI (the complexification of the field) is not analytic across the segment $[-iA, iA]$.

What is true, however, is that V is analytic in a Riemann surface consisting of infinitely many sheets, cut along the line segment $[-iA, iA]$. So, the appropriate, underlying space for the (doubled up) variational problem should now be a noncompact Riemann surface, say L .

Compactness is crucial in the proof of a maximizing continuum. But we can compactify the Riemann surface L by compactifying the complex plane. Let the map $C \rightarrow L$ be defined by

$$y = \log(z - iA) - \log(z + iA).$$

The point $z=iA$ corresponds to infinitely many y points, i.e., $y=-\infty+i\theta$, $\theta \in \mathbb{R}$, which will be identified. Similarly, the point $z=-iA$ corresponds to infinitely many points $y=+\infty+i\theta$, $\theta \in \mathbb{R}$, which will also be identified. The point $0 \in C$ corresponds to the points $k\pi i$, k odd.

By compactifying the plane we then compactify the Riemann surface L . The distance between two points in the Riemann surface L is defined to be the corresponding stereographic distance between the images of these points in the compactified C .

With these changes, the proof of the existence of the maximizing continuum in Secs. I, III, and IV goes through virtually unaltered. In Sec. VI, we would have to consider the complex field V as a function defined in the Riemann surface L and all proofs go through. The corresponding result of Sec. VII will give us an S-curve C in the Riemann surface L . We then have the following facts.

Consider the image D of the closed upper half-plane under

$$y = \log(z - iA) - \log(z + iA).$$

Consider continua in D containing the points $y=\pi i$ and $y=-\pi i$. Define the Green's potential and Green's energy of a Borel measure by (4)–(6) and the equilibrium measure by (7). Then there exists a continuum C maximizing the equilibrium energy, for the field given by (3) with conditions (1). C does not touch ∂D except at a finite number of points. By taking variations as in Sec. VI, one sees that C is a S-curve. In particular, the support of the equilibrium measure on C is a union of analytic arcs and at any interior point of $\text{supp } \mu$,

$$\frac{d}{dn_+}(\phi + V^{\lambda^C}) = \frac{d}{dn_-}(\phi + V^{\lambda^C}), \quad (8''')$$

where the two derivatives above denote the normal derivatives.

As far as the consequences of the above facts regarding the semiclassical limit of NLS, some more work on the details is necessary. Indeed, the analysis of the dynamics of the semiclassical problem (9) and (10) presented in Ref. 3 assumes that the S-curve C lies entirely in $\mathbb{K} \cup \{0_+, 0_-\}$. However, the explicit computation of the equilibrium measure on the S-curve and the derivation of the equations implicitly defining the endpoints of the components of the support of the equilibrium measure also make sense when extended to the Riemann surface L . As a result the statement of Theorem A.1 is correct, when interpreted correctly, i.e., allowing for all contours involved in the definition of the line integrals appearing in formula (B1) as well as the bands $I_0 \cup I_0^*, I_j, I_j^*, j = 1, \dots, G$ and the gaps Γ_j to lie in L .

We plan to describe the details of the evolution of the S-curve in the Riemann surface L in a later publication. In particular, one must check that all the deformations described in Ref. 3 are valid; this is indeed true, as one can see by the analysis of Ref. 3 transferred to the Riemann surface L .

Remark: The “discrete-to-continuous” passage.

The following caveat has to be addressed here. One of the problems dealt with in Ref. 3 was to transform the so-called “discrete” Riemann-Hilbert problem (problem 4.1.1 of Ref. 3) to a “continuous” Riemann-Hilbert problem (problem 4.2.1 of Ref. 3). The initial (discrete) problem posed is to construct a (matrix) meromorphic function from the information on its poles. These poles indeed accumulate (as \hbar goes to 0) along the segment $(-iA, iA)$. Then one transforms the

problem to a properly speaking (continuous) holomorphic Riemann-Hilbert problem with jumps, by constructing two contours, one encircling the poles with positive imaginary part and the other encircling the poles with negative imaginary part, and redefining the matrix functions inside the contours so that the singularities are removed. Naturally the information concerning the poles appears in the arising jump along the two loops.

One then (formally at first) approximates the jump on the loops in an obvious way (it is a Riemann sum approximated by an integral). This jump involves a logarithmic function with a cut along the linear segment $[-iA, iA]$.

Now, to rigorously show that the approximation is valid one needs a very careful analysis of what happens near 0, because the ‘‘Riemann sum’’ approximation breaks down there. Indeed this is done in Sec. 4.4.3 of Ref. 3, where, however, it is assumed that the loop (the jump contour) is located exactly at the S-curve (which solves the variational problem), and that this S-curve emanates from $0+$ at an acute, nonright, angle to the horizontal axis.

Even though in Ref. 3 the S-curve is assumed not to touch the linear segment $[-iA, iA]$, this condition is actually only necessary *locally* near 0. There is no problem in allowing the S-curve to intersect the spike, as long as it emanates from 0_+ at an acute, nonright, angle to the horizontal axis. Now it is *proved* in Ref. 3 that the S-curve solving the variational problem must emanate from $0+$ at an acute, nonright, angle to the horizontal axis. [This in fact follows immediately from the measure reality condition (4.18) applied at the origin.] So the proof of Theorem A.1 will go through even if the S-curve intersects the linear segment $[-iA, iA]$.

We then have the following.

Theorem 9: Consider the semiclassical limit ($\hbar \rightarrow 0$) of the solution of (9) and (10) with bell-shaped initial data. Replace the initial data by the so-called soliton ensembles data (as introduced in Ref. 3) defined by replacing the scattering data for $\psi(x, 0) = \psi_0(x)$ by their WKB approximation. Assume, for simplicity, that the spectral density of eigenvalues satisfies conditions (1).

Then, asymptotically as $\hbar \rightarrow 0$, the solution $\psi(x, t)$ admits a ‘‘finite genus description,’’ in the sense of Theorem A.1.

¹A. A. Gonchar and E. A. Rakhmanov, *Math. USSR. Sb.* **62**, 305 (1989).

²E. B. Saff and V. Totik, *Logarithmic Potentials with External Fields* (Springer, Verlag, 1997).

³S. Kamvissis, K. T.-R. McLaughlin, and P. D. Miller, *Semiclassical Soliton Ensembles for the Focusing Nonlinear Schrödinger Equation*, *Annals of Mathematics Studies*, Vol. 154 (Princeton University Press, Princeton, NJ, 2003).

⁴P. Deift and X. Zhou, *Ann. Math.* **137**, 295 (1993).

⁵E. A. Perevozhnikova and E. A. Rakhmanov (unpublished).

⁶J. Dieudonné, *Foundations of Modern Analysis* (Academic, New York, 1969).

⁷N. S. Landkof, *Foundations of Modern Potential Theory* (Springer-Verlag, Berlin, 1972).

⁸G. M. Goluzin, *Geometric Theory of Functions of a Complex Variable*, *Translations of Mathematical Monographs*, Vol. 26 (AMS, Providence, RI, 1969).