Riemann-Hilbert Problems with Lots of Discrete Spectrum: Asymptotics and Applications

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Abstract

I will discuss several situations in which an asymptotic limit of interest leads one to consider the construction of a matrix-valued meromorphic function with principal part data specified at asymptotically many poles. Applications include semiclassical asymptotics of integrable nonlinear wave problems (KdV, NLS, sine-Gordon) as well as statistical combinatorics (discrete orthogonal polynomial ensembles, i.e. discrete analogues of random matrix theory).



Outline

- I. Introduction
- II. The Semiclassical Focusing Nonlinear Schrödinger Equation
- III. Discrete Orthogonal Polynomials
- V. The Semiclassical Sine-Gordon Equation
- VI. Conclusions



Recall the analysis of Lax and Levermore (1983) of the zero-dispersion limit of the Cauchy problem for the Korteweg-de Vries equation. For each $\varepsilon > 0$ there exists a unique global solution of

$$u_t + uu_x + \varepsilon^2 u_{xxx} = 0, \quad x \in \mathbb{R}, \quad t > 0,$$

subject to ε -independent initial data $u(x, 0; \varepsilon) = u_0(x)$. The "zero-dispersion limit" analyzed by L&L refers to the asymptotic analysis of the family of solutions $u(x, t; \varepsilon)$ as $\varepsilon \to 0$.



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L&L used inverse-scattering to solve the Cauchy problem, assuming that $u_0(x) > 0$, rapidly decreasing as $|x| \to \infty$, and having a single critical point (local max). The first step: analyze the stationary Schrödinger equation

$$-6\varepsilon^2\psi_{xx} + V(x)\psi = E\psi$$

where E is the spectral parameter and $V(x) := -u_0(x)$ is a potential well.



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- 1. The reflection coefficient for E > 0 fixed is "as small in ε as V is smooth".
- 2. The number $N(\varepsilon)$ of discrete eigenvalues (all simple) is large, proportional to ε^{-1} . The eigenvalues are approximately located according to the *Bohr-Sommerfeld quantization rule*: $E_n = E_n^0 + O(\varepsilon^2)$ where

$$\Phi(E_n^0) = \pi \varepsilon \left(n + \frac{1}{2} \right) , \quad n = 0, \dots, N(\varepsilon) - 1 , \quad \Phi(E) := \frac{1}{\sqrt{6}} \int_{x_-(E)}^{x_+(E)} \sqrt{E - V(s)} \, ds \, .$$

Here $x_{-}(E) < x_{+}(E)$ are the *turning points* (branches of V^{-1}). The asymptotic number of eigenvalues is

$$N(\boldsymbol{\varepsilon}) = \left[\frac{1}{2} + \frac{1}{\pi \boldsymbol{\varepsilon} \sqrt{6}} \int_{-\infty}^{+\infty} \sqrt{-V(x)} \, dx \right] \, .$$



L&L therefore modified the spectral data associated with $V(x) = -u_0(x)$ by taking the reflection coefficient to be zero and using the approximate eigenvalues $\{E_n^0\}$. Thus, the approximate solution of the KdV Cauchy problem is a pure ensemble of $N(\varepsilon)$ solitons.



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The multisoliton solution of KdV is specified by a collection of $N(\varepsilon)$ discrete eigenvalues and a "norming constant" for each. For such purely discrete data, the inverse-scattering procedure collapses to a problem of linear algebra in dimension $N(\varepsilon)$. Cramer's rule leads to the Kay-Moses determinantal formula

$$u(x,t;\epsilon) = 12\epsilon^2 \frac{\partial^2}{\partial x^2} \log(\tau) \quad \text{where} \quad \tau := \det\left(\delta_{jk} + \frac{F_j F_k}{\kappa_j + \kappa_k}\right)$$

Here $\kappa_n = \sqrt{-E_n}$ and $F_n = e^{(\kappa_n x - 4\kappa_n^3 t + \beta_n)/\epsilon}$ and $\{\beta_n\}$ amount to the norming constants.



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Introduction

A natural approach is to expand au:

$$\tau = 1 + \sum_{\text{subsets } S \text{ of } \{0, \dots, N(\epsilon) - 1\}} \det \left(\frac{F_{\alpha} F_{\beta}}{\kappa_{\alpha} + \kappa_{\beta}} \bigg|_{\alpha, \beta \in S} \right)$$



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L&L observed: each term in the sum is positive. They proved that as $N(\varepsilon) \to \infty$ the sum is dominated by its largest term. This leads to a discrete variational problem that may be further approximated by a variational problem for an absolutely continuous equilibrium measure. In this way, L&L proved that $u(x, t; \varepsilon)$ has a weak limit $\overline{u}(x, t)$.



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Venakides (1990) obtained strong asymptotics for $u(x, t; \epsilon)$ by "going to higher order", in particular by quantizing the mass of the equilibrium measure. He found strongly nonlinear oscillations of unit amplitude about the mean $\overline{u}(x, t)$ modeled by algebro-geometric multiphase wave solutions of KdV.



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2. M takes continuous boundary values $\mathbf{M}_{\pm}(z)$, $z \in \Sigma$, that are related by $\mathbf{M}_{+}(z) = \mathbf{M}_{-}(z)\mathbf{V}(z)$ (+ means left, - means right).



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Steepest descent in a nutshell: a systematic construction of a "global parametrix" $\hat{\mathbf{M}}(z)$ such that the linear substitution $\mathbf{M}(z) = \mathbf{E}(z)\hat{\mathbf{M}}(z)$ results in a "small-norm" Riemann-Hilbert problem for the *error* $\mathbf{E}(z)$. By definition, a small-norm problem is one for which simple estimates can establish that $\mathbf{E}(z) \approx \mathbb{I}$ in a suitable sense.



In inverse-scattering problems, Σ is identified with the continuous spectrum, and $V \neq I$ only where the reflection coefficient is nonzero. If there is any discrete spectrum, the Riemann-Hilbert problem must be modified: M may have poles at the eigenvalues, and homogeneous conditions involving the norming constants are imposed to relate the principal part of the Laurent expansion at each pole to the regular part.



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$$\operatorname{Res}_{z=z_p} \mathbf{M}(z) = \lim_{z \to z_p} \mathbf{M}(z) \begin{bmatrix} 0 & c \\ 0 & 0 \end{bmatrix}$$



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$$\mathbf{N}(z) := \mathbf{M}(z) \begin{bmatrix} 1 & -c(z-z_p)^{-1} \\ 0 & 1 \end{bmatrix} \text{ for } |z-z_p| < \delta, \text{ and } \mathbf{N}(z) := \mathbf{M}(z) \text{ for } |z-z_p| > \delta.$$

 $\mathbf{N}(z)$ is analytic for $z \in \mathbb{C} \setminus \Sigma \cup \Sigma_p$, where Σ_p is a small, positively oriented circle about z_p of radius δ .



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Then,
$$\mathbf{N}_+(z) = \mathbf{N}_-(z)\mathbf{V}(z)$$
 for $z \in \Sigma$ and $\mathbf{N}_+(z) = \mathbf{N}_-(z) \begin{bmatrix} 1 & -c(z-z_p)^{-1} \\ 0 & 1 \end{bmatrix}$ for $z \in \Sigma_p$.



This "one-at-a-time" procedure for removing the poles from the problem becomes impractical exactly in situations like that faced by L&L: the poles are accumulating as $\varepsilon \rightarrow 0$.

Return to outline.



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It is interesting to consider how the Deift-Zhou steepest descent method can be applied to such problems. Indeed, by working with different initial data $u_0(x)$ for which there is no discrete spectrum but for which WKB analysis predicts an asymptotic formula for the reflection coefficient, Deift, Venakides, and Zhou (1997) showed how the method can reproduce strong asymptotics in a single step (rather than by going to higher order as required in the L&L approach).

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Moreover, there are a number of problems involving "lots of discrete spectrum" to which the L&L method does not apply at all, but for which the Riemann-Hilbert problem offers an alternative approach. . .

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A problem with many apparent similarities:

$$i\varepsilon\psi_t + \frac{\varepsilon^2}{2}\psi_{xx} + |\psi|^2\psi = 0, \quad x \in \mathbb{R}, \quad t > 0,$$

subject to initial data $\psi(x,0;\epsilon) = A(x)e^{iS(x)/\epsilon}$ where $A(\cdot)$ and $S(\cdot)$ are real and independent of ϵ .



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Note: formal zero-dispersion limit of KdV is inviscid Burgers' equation: $u_t + uu_x = 0$ (hyperbolic). Fact: the well-posed Cauchy problem for this equation governs the early stages of the zero-dispersion limit for KdV.



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$$\rho_t + u\rho_x + \rho u_x = 0 \quad \text{and} \quad u_t - \rho_x + uu_x = \epsilon^2 F[\rho],$$

subject to $\rho(x, 0; \epsilon) = A(x)^2$ and $u(x, 0; \epsilon) = S'(x)$. The formal limit (neglecting $\epsilon^2 F[\rho]$) is a Cauchy problem for an elliptic system. This is an ill-posed (formal) limit problem!



$$\boldsymbol{\varepsilon} \mathbf{u}_x = \begin{bmatrix} -i\lambda & A(x)e^{iS(x)/\boldsymbol{\varepsilon}} \\ -A(x)e^{-iS(x)/\boldsymbol{\varepsilon}} & i\lambda \end{bmatrix} \mathbf{u}, \text{ where } \lambda \in \mathbb{C} \text{ is the spectral parameter.}$$



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A remarkable fact: Klaus and Shaw (2002) showed that if $S(x) \equiv 0$ and A(x) is a positive $L^1(\mathbb{R})$ function with a single critical point (a local max) then all eigenvalues are purely imaginary numbers.



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WKB calculations for Klaus-Shaw potentials yield results analogous to those for the Schrödinger operator:

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- 1. For $\lambda \in \mathbb{R}$, $\lambda \neq 0$, the reflection coefficient is as small in ϵ as $A(\cdot)$ is smooth.
- 2. Bohr-Sommerfeld quantization rule for eigenvalues $\lambda_n \in [0, i \max A(x)]$: $\lambda_n \approx \lambda_n^0$ where

$$\Psi(\lambda_n^0) = \pi \varepsilon \left(n + \frac{1}{2} \right), \quad n = 0, \dots, N(\varepsilon) - 1, \quad \Psi(\lambda) := \int_{x_-(\lambda)}^{x_+(\lambda)} \sqrt{A(s) + \lambda^2} \, ds \, .$$

Here $x_{-}(\lambda) < x_{+}(\lambda)$ are turning points. The asymptotic number of positive imaginary eigenvalues is

$$N(\epsilon) = \left\lfloor \frac{1}{2} + \frac{1}{\pi\epsilon} \int_{-\infty}^{+\infty} A(x) \, dx \right\rfloor \, .$$



This suggests an approach completely parallel to that applied by L&L to KdV: neglect the reflection coefficient and consider the asymptotic behavior of the "semiclassical soliton ensemble" given by the reflectionless potential associated with the WKB approximations to the eigenvalues.



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Here $\{\lambda_k\}$ are the eigenvalues in the upper half-plane and $E_k := e^{i(\lambda_k x + \lambda_k^2 t + \beta_k)/\epsilon}$ where $\{\beta_k\}$ are like the norming constants.


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The Lax-Levermore method fails because the principal minors expansion of τ consists of both positive and negative terms!



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$$\boldsymbol{\epsilon} = \boldsymbol{\epsilon}_N := \frac{1}{\pi N} \int_{-\infty}^{+\infty} A(x) \, dx \,, \quad N = 1, 2, 3, \dots$$

This makes the reflection coefficient uniformly small. Note $N(\epsilon) = N$.



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2. A(-x) = A(x). This allows a simple derivation via WKB of the proportionality constants (related to norming constants): $\gamma_k^0 := (-1)^{k+1}$.



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- 2. A(-x) = A(x). This allows a simple derivation via WKB of the proportionality constants (related to norming constants): $\gamma_k^0 := (-1)^{k+1}$.
- 3. $A(\cdot)$ is analytic. This comes in by making $\Psi(\lambda)$ analytic, *not* by allowing solution of the limiting (ill-posed) Cauchy problem.



Neglecting the reflection coefficient and taking the WKB eigenvalues and proportionality constants as exact spectral data, let

$$c_k(x,t) := \frac{1}{\gamma_k} \operatorname{Res}_{\lambda=\lambda_k} W(\lambda), \quad W(\lambda) := e^{2i(\lambda x + \lambda^2 t)/\varepsilon} \prod_{n=0}^{N-1} \frac{\lambda - \lambda_n^*}{\lambda - \lambda_n}.$$

The Riemann-Hilbert problem of inverse scattering is to find a 2×2 matrix $\mathbf{m}(\lambda)$, $\lambda \in \mathbb{C}$, with the following properties:

1. $\mathbf{m}(\lambda)$ is a rational function of λ with simple poles confined to $\{\lambda_n, \lambda_n^*\}$ such that for $k = 0, \ldots, N-1$:

$$\operatorname{Res}_{\lambda=\lambda_k} \mathbf{m}(\lambda) = \lim_{\lambda \to \lambda_k} \mathbf{m}(\lambda) \begin{bmatrix} 0 & 0 \\ c_k(x,t) & 0 \end{bmatrix}, \qquad \operatorname{Res}_{\lambda=\lambda_k^*} \mathbf{m}(\lambda) = \lim_{\lambda \to \lambda_k^*} \mathbf{m}(\lambda) \begin{bmatrix} 0 & -c_k(x,t)^* \\ 0 & 0 \end{bmatrix}.$$



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2. $\lim_{\lambda \to \infty} \mathbf{m}(\lambda) = \mathbb{I}.$

The semiclassical soliton ensemble itself is given by $\psi(x, t; \epsilon) = 2i \lim_{\lambda \to \infty} \lambda m_{12}(\lambda)$.



This is not a traditional RHP: no jumps, just poles. Indeed it seems to be simple, solvable by partial fractions:

$$\mathbf{m}(\lambda) = \mathbb{I} + \sum_{k=0}^{N-1} \frac{\mathbf{a}_k}{\lambda - \lambda_k} + \sum_{k=0}^{N-1} \frac{\mathbf{b}_k}{\lambda - \lambda_k^*}$$

where \mathbf{a}_k and \mathbf{b}_k are matrix coefficients determined by the residue conditions, leading to a linear algebra problem for the coefficients.



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However, for N reasonably large, it can be solved numerically on a grid of independent (x, t)-values. The resulting plots reveal marvelous structures (here $A(x) = 2 \operatorname{sech}(x)$ and $\epsilon = 2/N$ with N = 40):





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Kamvissis, McLaughlin, and M (2003) studied the discrete RHP by removing all the poles "at once": the key observation is that it is only necessary to find an analytic function interpolating the proportionality constants $\{\gamma_k\}$ at the corresponding eigenvalues $\{\lambda_k\}$. Such a function may be constructed from the WKB phase integral; indeed

$$\gamma_k = -ie^{-i\Psi(\lambda_k)/\epsilon}, \quad k = 0, 1, 2, \dots, N-1.$$



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The Deift-Zhou technique may now be applied, the key step of which is to stabilize the problem with a "g-function". Supposing that $g(\lambda)$ is a function analytic in $\mathbb{C} \setminus \Sigma \cup \Sigma^*$ with $g \to 0$ as $\lambda \to \infty$ and $g(\lambda) + g(\lambda^*)^* = 0$ define a new unknown: $\mathbf{N}(\lambda) := \mathbf{M}(\lambda)e^{-g(\lambda)\sigma_3/\varepsilon}$.



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Note that for λ bounded away from the eigenvalues a continuum limit of Riemann sums shows that

$$W(\lambda) := \exp\left(\frac{1}{\varepsilon} \left[2i\lambda x + 2i\lambda^2 t + \sum_{n=0}^{N-1} \varepsilon \log(\lambda - \lambda_n^*) - \sum_{n=0}^{N-1} \varepsilon \log(\lambda - \lambda_n)\right]\right)$$



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Thus, the jump condition satisfied by $\mathbf{N}(\lambda)$ on Σ is $\mathbf{N}_{+}(\lambda) = \mathbf{N}_{-}(\lambda) \begin{bmatrix} e^{i\theta(\lambda)/\epsilon} & 0\\ iS(\lambda)e^{\phi(\lambda)/\epsilon} & e^{-i\theta(\lambda)/\epsilon} \end{bmatrix}$, where $\theta(\lambda) := i(g_{+}(\lambda) - g_{-}(\lambda))$ and $\phi(\lambda) := F(\lambda) - g_{+}(\lambda) - g_{-}(\lambda)$.



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One then aims to choose the contour Σ and the function $g(\lambda)$ so that as $\epsilon \to 0$ the jump of \mathbb{N} takes two alternative forms. This *determines* both g and several arcs of $\Sigma \cup \Sigma^*$. The local dynamics of $\psi(x, t; \epsilon)$ are given in terms of Θ for the double cover of \mathbb{C} with cuts on these arcs.



Sometimes this procedure fails, because the contour arcs it predicts intersect the imaginary intervals where the eigenvalues are accumulating as $\varepsilon \to 0$. Thus, Σ does not completely encircle all of the eigenvalues and some poles have never been removed!



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For example, when t = 0, the theory predicts that one arc of Σ should be contained within the imaginary interval $[0, i \max A(x)]$. In M (2002) it is shown how to nonetheless make use of the g-function at t = 0. The key is to use two different analytic interpolants of the proportionality constants $\{\gamma_k\}$ at the eigenvalues $\{\lambda_k\}$:



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The two interpolants used correspond to two choices of $j \in \mathbb{Z}$ in the more general formula

$$\gamma_k = -i(-1)^j e^{-i(2j+1)\Psi(\lambda_k)/\varepsilon}, \quad k = 0, 1, 2, \dots, N-1.$$

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Return to outline.

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It turns out that the g-function causes headaches not only when extra symmetry is present as at t = 0, but for other (x, t) as well. Lyng and M (2007) showed that another collision of Σ with the interval $[0, i \max A(x)]$ occurs somewhere between the primary and secondary caustic curves:

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The situation is repaired using three interpolants: j = -1, j = 0, and j = 1. Return to outline.



Consider random tilings of a hexagon with rhombus-shaped tiles:



We suppose that all tilings are equally likely.



Discrete Orthogonal Polynomials

Here is a typical tiling (thanks to Jim Propp):


Using an interpretation of random tilings as a problem of nonintersecting random walks, Johansson (2000) proved the following about the positions of the blue tiles in any vertical slice of the figure:

Theorem 1. [Johansson] Let $\xi_1 < \cdots < \xi_L$ denote the integer-valued positions (counted from the bottom) of the blue tiles in the m^{th} vertical slice (from the left) of the hexagon with side lengths $\mathfrak{a} \ge \mathfrak{b}$ and \mathfrak{c} . This slice contains

$$N := \mathfrak{c} + \frac{\mathfrak{a} - |m - \mathfrak{a}|}{2} + \frac{\mathfrak{b} - |m - \mathfrak{b}|}{2}$$

possible positions of tiles, and $L = N - \mathfrak{c}$. Then, the probability $P_m(\xi_1, \ldots, \xi_L)$ of finding this configuration of blue tiles is

$$P_m(\xi_1, \dots, \xi_L) = \frac{1}{Z} \prod_{1 \le j < k \le L} (\xi_j - \xi_k)^2 \prod_{j=1}^L w(\xi_j),$$

where

$$w(\xi) := \frac{(\xi + |m - \mathfrak{a}|)!(N - \xi - 1 + |m - \mathfrak{b}|)!}{\xi!(N - \xi - 1)!} \,.$$



K. Johansson, Comm. Math. Phys., 209, 437-476, 2000.

Note: $w(\xi)$ is the weight function for the family of "discrete orthogonal" Hahn polynomials $\{p_{N,k}(z)\}_{k=0}^{N-1}$ defined by the orthogonality conditions

$$\sum_{j=0}^{N-1} p_{N,k}(x_{N,j}) p_{N,l}(x_{N,j})^* w(x_{N,j}) = \delta_{kl}, \quad x_{N,j} := j.$$



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By standard arguments of random matrix theory, all multipoint correlations and gap probabilities are therefore encoded in the Hahn polynomials, and asymptotics of these polynomials leads to facts about the scaling limit in which $\mathfrak{a}, \mathfrak{b}, \mathfrak{c} \to \infty$.



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This leads to the general question of how to correctly formulate and compute asymptotics for families of discrete orthogonal polynomials, a question we will now discuss.



Consider the discrete measure

$$u(x) = \sum_{n=0}^{N-1} w_{N,n} \delta(x - x_{N,n})$$

where the nodes of orthogonality $\{x_{N,n}\}$ are prescribed for each N by giving a fixed real-analytic function $\rho^0(x) > 0$ on an interval [a, b] with

$$\int_a^b
ho^0(x)\,dx = 1\,,$$

and then imposing the conditions (resembling a Bohr-Sommerfeld quantization rule)

$$\int_{a}^{x_{N,n}} \rho^{0}(x) \, dx = \frac{2n+1}{2N} \, , \ n = 0, \dots, N-1 \, .$$

The weights $\{w_{N,n}\}$ are assumed to be of the form

$$w_{N,n} = e^{-NV(x_{N,n})} \prod_{\substack{m=0\\m \neq n}}^{N-1} |x_{N,n} - x_{N,m}|^{-1}$$

for some real analytic function V(x) (only weakly dependent on N, if at all).



Let $\mathbb{Z}_N := \{0, \dots, N-1\}, \Delta \subset \mathbb{Z}_N$, and $\nabla := \mathbb{Z}_N \setminus \Delta$. Set $W(z) := e^{-NV(z)} \frac{\prod_{n \in \Delta} (z - x_{N,n})}{\prod_{n \in \nabla} (z - x_{N,n})}$.



A. Borodin, Internat. Math. Res. Not., 9, 467-494, 2000.

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The DOPs satisfy a (fully discrete) RHP. Seek a 2 imes 2 matrix $\mathbf{Q}(z)$, $z \in \mathbb{C}$ such that:

1. $\mathbf{Q}(z)$ is a rational function of z with simple poles confined to the nodes $\{x_{N,n}\}_{n=0}^{N-1}$ such that

$$\operatorname{Res}_{z=x_{N,n}} \mathbf{Q}(z) = \lim_{z \to x_{N,n}} \mathbf{Q}(z) \begin{bmatrix} 0 & (-1)^{N-1-n} \operatorname{Res}_{\zeta=x_{N,n}} W(\zeta) \\ 0 & 0 \end{bmatrix}, \ n \in \nabla$$
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$$\begin{split} \underset{z \to \infty}{\operatorname{Res}} \mathbf{Q}(z) &= \lim_{z \to x_{N,n}} \mathbf{Q}(z) \begin{bmatrix} 0 & (-1)^{N-1-n} \underset{\zeta = x_{N,n}}{\operatorname{Res}} W(\zeta) \\ 0 & 0 \end{bmatrix}, \ n \in \nabla \\ \underset{z = x_{N,n}}{\operatorname{Res}} \mathbf{Q}(z) &= \lim_{z \to x_{N,n}} \mathbf{Q}(z) \begin{bmatrix} 0 & (-1)^{N-1-n} \underset{\zeta = x_{N,n}}{\operatorname{Res}} W(\zeta)^{-1} & 0 \\ (-1)^{N-1-n} \underset{\zeta = x_{N,n}}{\operatorname{Res}} W(\zeta)^{-1} & 0 \end{bmatrix}, \ n \in \Delta \,. \\ \underset{z \to \infty}{\operatorname{lim}} \mathbf{Q}(z) z^{(\#\Delta - k)\sigma_3} = \mathbb{I}. \end{split}$$



2.

A. Borodin, Internat. Math. Res. Not., 9, 467-494, 2000.

Let $\mathbb{Z}_N := \{0, \dots, N-1\}, \Delta \subset \mathbb{Z}_N$, and $\nabla := \mathbb{Z}_N \setminus \Delta$. Set $W(z) := e^{-NV(z)} \frac{\prod_{n \in \Delta} (z - x_{N,n})}{\prod_{n \in \nabla} (z - x_{N,n})}$.

The DOPs satisfy a (fully discrete) RHP. Seek a 2 imes 2 matrix $\mathbf{Q}(z)$, $z \in \mathbb{C}$ such that:

1. $\mathbf{Q}(z)$ is a rational function of z with simple poles confined to the nodes $\{x_{N,n}\}_{n=0}^{N-1}$ such that

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The monic scaling of $p_{N,k}(z)$ is given by $\pi_{N,k}(z) = z^k + \cdots = Q_{11}(z) \prod_{n \in \Delta} (z - x_{N,n}).$



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To analyze this RHP, we know based on experience that we should remove the poles. The key is to find an analytic function that interpolates the signs $\{(-1)^{N-1-n}\}$ at the nodes $\{x_{N,n}\}$. Setting

$$\theta^{0}(z) := 2\pi \int_{z}^{b} \rho^{0}(s) \, ds \,, \text{ we have the identities } i e^{-iN\theta^{0}(x_{N,n})/2} = -i e^{iN\theta^{0}(x_{N,n})/2} = (-1)^{N-1-n}$$

Return to outline.



J. Baik, T. Kriecherbauer, K. T.-R. McLaughlin, and P. D. Miller, *Discrete Orthogonal Polynomials: Asymptotics and Applications*, Vol. 164, Annals of Mathematics Studies Series, Princeton University Press, Princeton, 2007.

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This problem is similar to the t = 0 situation for the focusing NLS. There is a g-function for this problem and its support lies in the interval [a, b]. Therefore we need both of these interpolants.

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$$\begin{split} \mathbf{R}(z) &:= \mathbf{Q}(z) \begin{bmatrix} 1 & \mp i e^{\mp i N \theta^0(z)/2} W(z) \\ 0 & 1 \end{bmatrix}, z \in \Omega_{\pm}^{\nabla} \\ \mathbf{R}(z) &:= \mathbf{Q}(z) \begin{bmatrix} 1 & 0 \\ \mp i e^{\mp i N \theta^0(z)/2} W(z)^{-1} & 1 \end{bmatrix}, z \in \Omega_{\pm}^{\Delta}, \end{split}$$

Ω^{Δ}_{+}	$\Omega^{ abla}_+$	Ω^{Δ}_+	Ω^{∇}_{+}
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and for all other $z \in \mathbb{C}$ set $\mathbf{R}(z) := \mathbf{Q}(z)$. We obtain a traditional RHP! Return to outline.



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Going back to nonlinear waves, another interesting Cauchy problem is:

$$\epsilon^2 u_{tt} - \epsilon^2 u_{xx} + \sin(u) = 0, \quad x \in \mathbb{R}, \quad t > 0$$

subject to initial data of the form $u(x, 0; \epsilon) = f(x)$ and $\epsilon u_t(x; 0; \epsilon) = g(x)$, independent of ϵ .



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- 3. Thus, $\epsilon := \ell_J / \ell_0 \approx 0.0005$.



At a more elementary level, it is the equation of motion of a simple mechanical system: a chain of coupled pendula all attached to the same rod:





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This configuration has "topological charge" equal to zero. Small *e* means (i) small mass of each pendulum bob and (ii) weak coupling between neighborhing pendula.



A configuration with unit topological charge is a "kink":





As a first step in analyzing this problem, Buckingham and M (2007) considered special initial data:

$$\sin\left(\frac{f}{2}\right) := \operatorname{sech}(x), \quad \cos\left(\frac{f}{2}\right) := \tanh(x), \quad g := 2\mu \operatorname{sech}(x)$$

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- 2. Taking $\mu = 0$ corresponds to the pendulum bobs being stationary at t = 0.
- 3. Taking $\mu \neq 0$ correponds to giving the pendula near x = 0 a (large) impulse at t = 0.



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Among other things, this allows us to numerically handle the inverse-scattering. Here is a plot of cos(u) over -2.5 < x < 2.5 and 0 < t < 5. Here $\mu = 0$ and N = 0 ($\varepsilon = 1$).





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And here is a plot of $\cos(u)$ over -2.5 < x < 2.5 and 0 < t < 5, now for $\mu = 1$ and N = 0 ($\varepsilon = \sqrt{2}$), bringing some kinks into the mix.





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An observation: the points of phase transition at t = 0 correspond to those values of x for which the initial data is crossing the separatrix in the phase portrait of the simple pendulum:



Return to outline.



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We are now in a position to analyze the discrete RHP associated with these reflectionless potentials and generalizations thereof (work in progress).

Return to outline.



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Thank You!

