ELLiptic AlgEBro-geoMETrIC SOlUTIONS OF THE KdV 
AND AKNS hiERARCHIES – AN ANALYTIC aPPROACH

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A BSTRACT. We provide an overview of elliptic algebro-geometric solutions of 
the KdV and AKNS hierarchies, with special emphasis on Floquet theoretic 
and spectral theoretic methods. Our treatment includes an effective charac-
terization of all stationary elliptic KdV and AKNS solutions based on a theory 
developed by Hermite and Picard.

1. INTRODUCTION.

The story of J. Scott Russell chasing a soliton for a mile or two along the 
Edinburgh-Glasgow channel in 1834 has been told many times. It is the start-
ing point of more than 160 years of an exciting history embracing a variety of deep 
mathematical ideas ranging from applied mathematics to algebraic geometry, Lie 
groups, and differential geometry. In this article we want to tell a certain aspect of 
this story, predominantly from an analytic point of view with special emphasis on 
Floquet theoretic and spectral theoretic methods. Other aspects of this story have 
recently been told by Lax [147], Palais [173], and Terng and Uhlenbeck [210].

We are interested in the class of evolution equations which permit being cast 
into the form $L_t = [P, L]$, where $(P, L)$ is a pair of operators, a so called Lax 
pair, and $[P, L]$ denotes their commutator. Typically, for a fixed operator $L$, there 
is a sequence of operators $P$ such that $L_t = [P, L]$ defines an evolution equation. 
Hence we are actually concerned with hierarchies of evolution equations. Stationary 
solutions of such equations, corresponding to commuting operators $P$ and $L$, are 
related to algebraic curves and are therefore called algebro-geometric solutions. The 
stationary solutions of higher-order equations in the hierarchy play a decisive role 
in the study of the Cauchy problem for the lower-order equations in the hierarchy. 
Therefore, and because of the connection to algebraic geometry, the stationary 
problem has drawn considerable attention.

The central object of this article is the class of elliptic algebro-geometric solutions 
of the KdV and AKNS hierarchies (cf. Section 2.3). In a sense, this class represents 
a natural generalization of the class of soliton solutions and enjoys a much richer 
structure due to its connections to algebraic geometry inherent in its construction. 
However, while frequently the algebraic aspects of this construction dominate the 
stage, we will purposely portray a different analytical view often neglected (and to 
some extent forgotten) in this context. Consequently, we will focus in the following 
on the interplay between spectral (and Floquet) theoretic properties of the Lax
pairs (cf. Section 2.3) defining the integrable evolution equations in question on the one hand, and the construction and properties of the underlying compact Riemann surface associated with algebro-geometric solutions on the other.

In our quest to characterize the class of elliptic algebro-geometric solutions of soliton equations in an effective manner, we rely heavily on a marvelous theory developed by Fuchs, Halphen, Hermite, Mittag-Leffler, and especially, Picard. We have chosen to provide a rather extensive bibliography concerning this classical work since some of the cornerstones of this theory seem to have been forgotten, and at times appear to be independently rediscovered (cf. Subsection 2.9 and the bibliographical remarks at the end of Remark 3.9). In particular, Picard’s theorem (Theorem 2.7), a key in the aforementioned characterization problem, apparently had not been used in the extensive body of literature surrounding elliptic algebro-geometric solutions of various hierarchies of soliton equations. Picard’s theorem suggests one consider the independent variable of the differential equation in question as a complex variable, and study the consequences of assuming the existence of a meromorphic fundamental system of solutions. This led to the discovery of a connection between the existence of such a meromorphic fundamental system of solutions of a linear differential equation \( Ly = zy \) for all values of \( z \) and algebro-geometric properties of \( L \). In particular, it connects the existence of a meromorphic fundamental system of solutions of \( Ly = zy \) to the integrability of the nonlinear equations associated with \( L \) via the Lax pair formalism, as we demonstrated in the KdV and AKNS cases (cf. Theorems 3.12 and 4.8). As it turns out, Picard was describing solutions of differential equations with elliptic coefficients, which represent simultaneous Floquet solutions with respect to all fundamental periods of the underlying period lattice of the torus in question. Hence Floquet theory, and consequently, spectral theory, naturally enters when analyzing Lax pairs for completely integrable evolution equations and their elliptic algebro-geometric solutions.

Section 2 provides an introduction into the KdV hierarchy and its elliptic algebro-geometric solutions, reviews the necessary Floquet and spectral theoretic background for Hill (Lax) operators, presents Picard’s theorem on linear differential equations with elliptic coefficients, and especially, provides a historical perspective of the subject. Section 3 is devoted to the KdV hierarchy and its stationary solutions. The case of rational, periodic, and elliptic stationary KdV solutions is described in detail and the role of meromorphic fundamental systems of solutions for the corresponding Lax operator associated with stationary KdV solutions is underscored. This section culminates in an explicit characterization of all elliptic, simply periodic, and rational stationary KdV solutions (see [97], [98], [234] for the original results). Our final Section 4 describes analogous results for the AKNS hierarchy, in particular, it contains an effective characterization of all elliptic algebro-geometric AKNS solutions (cf. [19] for the original proof). We hope the enormous bibliography at the end (still necessarily incomplete), will lose some of its intimidating character once the reader begins to appreciate the large body of knowledge amassed by some of the giants in the field, starting with Hermite.

Finally we remark that connections between completely integrable systems and Seiberg-Witten theory via elliptic Calogero-Moser-type models, have recently led to a strong resurgence of the field of elliptic solutions of soliton equations. It seems difficult to keep up with all the current activities in the corresponding preprint archives. Hence we refer, for instance, to Babelon and Talon [12], Donagi and Markman [17], Donagi and Witten [18], Itoyama and Morozov [20], Krichever [40],
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Krichever, Babelon, Billey, and Talon [141], Krichever and Phong [143], Krichever, Wiegmann, and Zabrodin [144], Krichever and Zabrodin [144], Kuznetsov, Nijhoff, and Sklyanin [145], Levin and Olshanetsky [149], Marshakov [152], and Vaninsky [229], from which the interested reader can easily find further sources.

2. Background and some History.

2.1. The Korteweg-de Vries Equation. In the first few decades after their discovery, solitary waves were considered, for instance, by Stokes, Boussinesq, and Lord Rayleigh. But the most far reaching, and in a sense, lasting contribution was made in 1895 by Korteweg and de Vries [129], who deduced their celebrated equation, which may be cast in the form

\[ q_t = \frac{1}{4}q_{xxx} + \frac{3}{2}qq_x. \]

(2.1.1)

(However, as recently alluded to by Pego [178], Boussinesq originally noted a system of equations equivalent to (2.1.1) and used it to study solitary waves in the 1870’s.) In 1960 Gardner and Morikawa [79] used the KdV equation to describe collisionless-plasma magnetohydromagnetic waves. Since then the KdV equation has been rederived again and again in different contexts as a model equation describing a considerable variety of physical phenomena, and can now be considered one of the basic equations in mathematical physics. Even more important is the fact that this development triggered the examination of a whole class of other nonlinear evolution equations, some of which are of great physical relevance, such as the nonlinear Schrödinger equation, the sine-Gordon equation, the Toda lattice, the Boussinesq equation, the Kadomtsev-Petviashvili equation, etc.

But the description of wave phenomena alone would not have turned the subject into the kind of industry it is today. In 1955 Fermi, Pasta, and Ulam [67] studied a system of nonlinear oscillators which may be viewed as a discretized version of the KdV equation on a finite interval with periodic boundary conditions. To their surprise they found that energy shows little tendency toward equipartition among the degrees of freedom. Later Kruskal and Zabusky [241] examined the KdV equation numerically and observed the formation of solitary waves which, using their words, “pass through one another without losing their identity”. In order to emphasize this particle-like behavior Kruskal and Zabusky coined the term soliton.

In 1968 Miura [164] introduced the transformation

\[ q(x,t) = z - v(z,x,t)^2 - v_x(z,x,t) \]

(2.1.2)

in which \( z \) is a (generally complex) spectral parameter. This transformation, now known as Miura’s transformation, relates the KdV equation to a variant of the so called modified KdV (mKdV) equation,

\[ v_t = v_{xxx} - 6v^2v_x + 6zv_x. \]

(2.1.3)

Since (2.1.2) is a Riccati equation for \( v \), it may be transformed into the linear equation

\[ (L(t)y)(z,x,t) = y''(z,x,t) + q(x,t)y(z,x,t) = zy(z,x,t) \]

by introducing \( v(z,x,t) = y'(z,x,t)/y(z,x,t) \) (primes denote derivatives with respect to \( x \)). Gardner, Kruskal, and Miura [165] showed that the eigenvalues of the \( L^2(\mathbb{R}) \)-operator associated with \( L(t) \) do not depend on \( t \), that is, they are constants
of the motion under the KdV flow. This was the starting point for another seminal work, this time by Gardner, Greene, Kruskal, and Miura [77], in which they used the inverse scattering method to solve the Cauchy problem for the KdV equation with rapidly decaying initial data. This method, which represents a nonlinear analog of the Fourier transform to solve linear partial differential equations, initially consists of computing the scattering data of $L(0)$, then propagating them in time (which is simple), and finally reconstructing the potential $q(x,t)$ in $L(t)$ via the Gelfand-Levitan, or rather, the Marchenko equation. Miraculously, this function $q(x,t)$ is the desired solution of the Cauchy problem. See, for instance, Ablowitz and Clarkson [1], Ch. 2, Asano and Kato [1], Chs. 5,6, Dodd, Eilbeck, Gibbon, and Morris [46], Ch. 4, Drazin and Johnson [52], Ch. 4, Gardner, Greene, Kruskal, and Miura [78], Iliev, Khristov, and Kirchev [117], Ch. 3, Lax [147], Marchenko [153], Ch. 4, and Palais [173] for more details.

2.2. Hamiltonian Systems. The dynamics of a classical mechanical particle system is described as a flow on a symplectic space, called the phase space. The phase space is the cotangent bundle of a Riemannian manifold, called the configuration space, describing the positions of the particles. The flow is given by Hamilton’s equation $q_t = \text{grad}_S(H)$, where grad$_S$ is a symplectic gradient and $H$ is the so-called Hamiltonian of the system. If there are only finitely many, say $n$, degrees of freedom (as in a system of finitely many interacting particles) then $q = (\xi, \eta)$, where $\xi$ is a vector of local coordinates on the configuration space (called generalized coordinates) and $\eta$ is the vector of the associated momenta (called generalized momenta). $H$ is then a function of $\xi, \eta,$ and $t$, and the symplectic gradient is given by

$$\text{grad}_S = \begin{pmatrix} 0 & I_n \\ -I_n & 0 \end{pmatrix} \text{grad}(\xi,\eta),$$

where $I_n$ is the $n \times n$ identity matrix. A canonical transformation is a change of coordinates in phase space which leaves the form of Hamilton’s equation invariant. Sometimes there exists a canonical transformation such that the transformed Hamiltonian is a function of the transformed generalized momenta alone which are then called action variables while the transformed generalized coordinates are called angle variables. The action variables are then constants of the motion while the angle variables change linearly with time. If this happens the system is called completely integrable.

In 1968 Gardner, Kruskal, and Miura [165] explicitly constructed an infinite sequence of constants of the motion providing the first hint toward a Hamiltonian structure of the KdV equation. Next, in 1971, Faddeev and Zakharov [242], in an attempt to explain the unusual behavior of the KdV equation, showed that it can be viewed as an infinite-dimensional completely integrable Hamiltonian system. Take, for instance, the Schwartz space $S(\mathbb{R})$ as phase space which can be viewed as a symplectic manifold. A symplectic gradient on $S(\mathbb{R})$ is given by

$$\text{grad}_S = \frac{\partial}{\partial x} \text{grad},$$

where grad$(F)$ denotes the Lagrangian (or variational) derivative of $F : S(\mathbb{R}) \to \mathbb{R}$. More precisely, if $F$ is of the form

$$F(q) = \int_{-\infty}^{\infty} dx \bar{F}(q, q', \ldots, q^{(n)}),$$
where $\tilde{F} : \mathbb{R}^{n+1} \to \mathbb{R}$ is a polynomial function without a constant term, then
\[ \text{grad}(F) = \frac{\partial \tilde{F}}{\partial q} - \left( \frac{\partial \tilde{F}}{\partial q} \right)' + \cdots + (-1)^n \left( \frac{\partial \tilde{F}}{\partial q^{(n)}} \right) \]
and hence
\[ \text{grad}_S(F) = \left( \frac{\partial \tilde{F}}{\partial q} \right)' - \left( \frac{\partial \tilde{F}}{\partial q} \right)'' + \cdots + (-1)^n \left( \frac{\partial \tilde{F}}{\partial q^{(n)}} \right)^{(n+1)} \]
using primes to denote derivatives with respect to $x$. The Hamiltonian flow
\[ \dot{q} = \text{grad}_S(F) \]
then becomes the nonlinear evolution equation
\[ q_t = \left( \frac{\partial \tilde{F}}{\partial q} \right)' - \left( \frac{\partial \tilde{F}}{\partial q} \right)'' + \cdots + (-1)^n \left( \frac{\partial \tilde{F}}{\partial q^{(n)}} \right)^{(n+1)}. \]
In the special case of the KdV flow, the Hamiltonian $H$ is given by
\[ H : S(\mathbb{R}) \to \mathbb{R}, \quad q \mapsto \frac{1}{4} \int_{-\infty}^{\infty} dx \ (q(x)^3 - \frac{1}{2} q'(x)^2), \]
and Hamilton’s equation $q_t = \text{grad}_S(H)$ results in the KdV equation (2.1.1). Analogous considerations apply to the periodic case replacing $\int_{-\infty}^{\infty} dx$ by $\int_{0}^{\Omega} dx$, with $\Omega > 0$ the fundamental period of $q$. We refer, for instance, to Palais [173] for more information.

2.3. Lax Pairs and the KdV Hierarchy. Gardner, Greene, Kruskal, and Miura [14] had shown that the $L^2$-spectrum of $L(t) = \partial^2/\partial x^2 + q(x, t)$ is independent of $t$ whenever $q$ is a solution of the KdV equation. This discovery inspired Lax to conjecture that the operators $L(t)$ are all unitarily equivalent to one another, that is, $L(0) = U(t)^{-1} L(t) U(t)$ for some family of unitary operators $U(t)$. This then led to his celebrated commutator representation [146],
\[ q_t = L_t = [P_3, L], \]
where $P_3$ denotes the third-order ordinary differential expression defined by
\[ P_3 = \frac{1}{4} \frac{\partial^3}{\partial x^3} + \frac{3}{2} \frac{\partial}{\partial x} + \frac{3}{4} q_x, \]
which governs the time evolution of $U(t)$ according to $U_t(t) = P_3(t) U(t)$, $U(0) = I$. Indeed, the latter equation implies $(U(t)^{-1})_t = -U(t)^{-1} P_3(t)$ and hence formally,
\[ \frac{d}{dt} (U(t)^{-1} L(t) U(t)) = U(t)^{-1} (L_t(t) - [P_3(t), L(t)]) U(t) = 0 \]
yields
\[ U(t)^{-1} L(t) U(t) = L(0), \quad t \in \mathbb{R}. \]
(Functional analytic arguments guaranteeing self-adjointness of $L(t)$ and unitarity of $U(t)$ in the Hilbert space $L^2(\mathbb{R}; dx)$ can easily be supplied.) The letters $P$ and $L$ were first used by Gelfand and Dickey [83] in honor of Peter Lax. Accordingly, $(P_3, L)$ is called a Lax pair.

In fact, Lax designed his procedure for a far more general setting: it applies to evolution equations for $q$ whenever there is a one-to-one correspondence between $q(\cdot, t)$ and a (formally) symmetric operator $L(t)$. The goal is then to find a (formally) skew-symmetric operator $P(t)$ satisfying $U_t(t) = P(t) U(t)$, where $L(0) =$
In particular, returning to the case where \( L(t) = \partial^2/\partial x^2 + q(x, t) \), Lax showed that one can always find an odd-order differential expression \( P_{2n+1}(t) \) such that \([P_{2n+1}, L(t)]\) is an operator of multiplication. The coefficients of \( P_{2n+1} \) are then differential polynomials in \( q(x, t) \), that is, polynomials of \( q(x, t) \) and its \( x \)-derivatives. The commutator \([P_{2n+1}, L] \) is then also a differential polynomial in \( q \) and the evolution equation

\[
q_t = [P_{2n+1}, L]
\]

defines the \( n \)-th equation of the KdV hierarchy.

More precisely, if \( L = \partial^2/\partial x^2 + q \), then any ordinary differential expression \( P \), for which \([P, L] = 0\) is the sum of a polynomial \( K \in \mathbb{C}[L] \) and an odd-order differential expression \( P_{2n+1} \), that is,

\[
P = K(L) + P_{2n+1}.
\]

Upon rescaling the time variable, \( P_{2n+1} \) can be assumed to be monic. Moreover, choosing the polynomial \( K \) appropriately, it may be shown that

\[
P_{2n+1} = \sum_{j=0}^{n} \left( -\frac{1}{2} f'_j + f_j \frac{d}{dx} \right) L^{n-j}
\]

for some integer \( n \in \mathbb{N}_0(= \mathbb{N} \cup \{0\}) \), where the functions \( f_j \) satisfy the recursion relation

\[
f_0 = 1,
\]

\[
f'_{j+1} = \frac{1}{4} f''_j + q f'_j + \frac{1}{2} q' f_j, \quad j = 0, \ldots, n.
\]

Then

\[
[P, L] = [P_{2n+1}, L] = \frac{1}{4} f'''_n + q f'_n + \frac{1}{2} q' f_n
\]

and introducing

\[
F_n(z, x, t) = \sum_{j=0}^{n} f_{n-j}(x, t) z^j,
\]

this recursion relation becomes

\[
[P_{2n+1}, L] = \frac{1}{2} F'''_n + 2(q - z) F' + q' F_n.
\]

Conversely, if \( F_n(z, x) \) is a polynomial in \( z \) such that \( F'''_n + 4(q - z) F'_n + 2 q' F_n \) does not depend on \( z \), then its coefficients define a differential expression \( P_{2n+1} \) such that equation \((2.3.4)\) is satisfied. In particular, if \( F'''_n + 4(q - z) F'_n + 2q' F_n = 0 \), then \( q \) is a stationary (i.e., \( t \)-independent) solution of one of the \( n \)-th equations in the KdV hierarchy.

The first few equations of the KdV hierarchy explicitly read,

\[
q_t = q_x,
\]

\[
q_t = \frac{1}{4} q_{xxx} + \frac{3}{2} q q_x + c_1 q_x,
\]

\[
q_t = \frac{1}{16} q_{xxxx} + \frac{5}{8} q q_{xxx} + \frac{5}{4} q_x q_{xx} + \frac{15}{8} q^2 q_x + c_1 \left( \frac{1}{4} q_{xxx} + \frac{3}{2} q_x \right) + c_2 q_x,
\]

etc.,
where the $c_\ell$ denote integration constants encountered in solving the recursion (2.3.3) for $f_{j+1}$. Here we ought to mention that each of these equations is a completely integrable Hamiltonian system and that the sequence of these equations is intimately related to the sequence of conservation laws discovered by Gardner, Kruskal, and Miura [163].

However, the Lax method can be further generalized in a variety of ways: If one considers instead of $L$ a first-order $2 \times 2$-matrix differential expression (i.e., a Dirac-type operator), one arrives at the ZS hierarchy [243] which is associated with the nonlinear Schrödinger equation, and more generally, at the AKNS hierarchy [9] which includes the KdV, the mKdV, and the nonlinear Schrödinger hierarchies as special cases. If one considers instead of the second-order differential expression $L$ a differential expressions of any order, one arrives at the Gelfand-Dickey hierarchy [115]. Moreover, the formalism is not confined to differential expressions but extends to formal pseudo-differential expressions (including rational functions of $d/dx$) needed in connection with the sine-Gordon and (modified) Kadomtsev-Petviashvili hierarchies.

2.4. The Algebra of Commuting Differential Expressions. As seen from (2.3.1), stationary (i.e., time-independent) Lax equations naturally lead to commuting differential expressions. Independently of this fact, the question of commuting differential expressions was raised by Floquet [70] in 1879. Some 25 years later, the question was again considered by Wallenberg [231] and Schur [193]. To this end, Schur developed the algebra of symbolic, that is, formal pseudo-differential expressions. He proves, in particular, the following statement: if $A$ and $B$ are two monic commuting differential expressions then the coefficients of $A$ are differential polynomials in the coefficients of $B$ (for a contemporary approach to these results, see, e.g., Wilson [236]). The decisive step, however, was done by Burchnall and Chaundy in the 1920’s. They proved the following result in [30].

Theorem 2.1. Let $A$ and $B$ be ordinary differential expressions of relatively prime orders $m$ and $n$, respectively. Then $[A, B] = 0$ if and only if there exists a polynomial $f$ of the form

$$f(\alpha, \beta) = \alpha^n - \beta^m + \sum_{j,k \geq 0, m_j + n_k < mn} c_{j,k} \alpha^j \beta^k,$$

(2.4.1)

such that $f(A, B) = 0$.

$f$ is called the Burchnall-Chaundy polynomial of $A$ and $B$. In [30] Burchnall and Chaundy constructed differential expressions $A$ and $B$ satisfying the equation $f(A, B) = 0$ for a given polynomial $f$ of the form (2.4.1) (cf. also Baker [10], Burchnall and Chaundy [32]). More recent treatments of the Burchnall-Chaundy theory can be found, for instance, in Carlson and Goodearl [37], Gatto and Greco [12], Giertz, Kwong, and Zettl [102], Greco and Previant [104], Krichever [133], [134], Mumford [168], and Wilson [237].

Finally, we briefly return to the Lax pairs $(P, L)$ discussed in Section 2.3. Because of the Burchnall-Chaundy relationship we call $L$ (or the set of its coefficients) algebro-geometric if there exists a corresponding $P$ such that $[P, L] = 0$. We will provide more precise definitions in the contexts of the KdV and the AKNS hierarchies in Definitions 3.1 and 4.2, respectively.
2.5. Elliptic Functions in a Nutshell. Elliptic functions provide some of the most important examples of algebro-geometric potentials. We present here a very brief account of Weierstrass’ point of view. For general references see, for instance, Akhiezer [7], Chandrasekharan [38], Markushevich [154], and Whittaker and Watson [233].

A function \( f : \mathbb{C} \to \mathbb{C} \cup \{ \infty \} \) with two periods \( a \) and \( b \), the ratio of which is not real, is called doubly periodic. If all its periods are of the form \( m_1a + m_2b \) where \( m_1 \) and \( m_2 \) are integers then \( a \) and \( b \) are called fundamental periods of \( f \).

A doubly periodic meromorphic function is called elliptic.

It is customary to denote the fundamental periods of an elliptic function by \( 2\omega_1 \) and \( 2\omega_3 \) with \( \text{Im}(\omega_3/\omega_1) > 0 \). We also introduce \( \omega_2 = \omega_1 + \omega_3 \) and \( \omega_4 = 0 \). The numbers \( \omega_1, \ldots, \omega_4 \) are called half-periods. The fundamental period parallelogram (f.p.p.) \( \Delta \) is the half-open region consisting of the line segments \([0, 2\omega_1), [0, 2\omega_3)\) and the interior of the parallelogram with vertices \( 0, 2\omega_1, 2\omega_2, 2\omega_3 \).

The class of elliptic functions with fundamental periods \( 2\omega_1, 2\omega_3 \) is closed under addition, subtraction, multiplication, division by non-zero divisors and differentiation. If \( f \) is an entire elliptic function then it is constant. A non-constant elliptic function \( f \) must have at least one pole in \( \Delta \) and the total number of poles in \( \Delta \) is finite. The total number of poles of an elliptic function \( f \) in \( \Delta \) (counting multiplicities) is called the order of \( f \). The sum of residues of an elliptic function \( f \) at all its poles in \( \Delta \) equals zero. In particular, the order of a non-constant elliptic function \( f \) is at least 2. The total number of points in \( \Delta \) where the non-constant elliptic function \( f \) assumes the value \( A \) (counting multiplicities), denoted by \( n(A) \), is equal to the order of \( f \). In particular, \( n(\infty) \geq 2 \). Furthermore, \( s(A) \), the sum of all the points in \( \Delta \) where the non-constant elliptic function \( f \) assumes the value \( A \), is congruent to \( s(\infty) \), the sum of all the points in \( \Delta \) where \( f \) has a pole, that is, \( s(A) = s(\infty) + 2m_1\omega_1 + 2m_3\omega_3 \), where \( m_1 \) and \( m_3 \) are certain integers.

The function

\[
\wp(z; \omega_1, \omega_3) = \frac{1}{z^2} + \sum_{m, n \in \mathbb{Z} \setminus \{0,0\}} \frac{1}{(z - 2m\omega_1 - 2n\omega_3)^2} - \frac{1}{(2m\omega_1 + 2n\omega_3)^2},
\]

or \( \wp(z) \) for short, was introduced by Weierstrass. It is an even elliptic function of order 2 with fundamental periods \( 2\omega_1 \) and \( 2\omega_3 \). Its derivative \( \wp'(z) \) is an odd elliptic function of order 3 with fundamental periods \( 2\omega_1 \) and \( 2\omega_3 \). Every elliptic function may be written as \( R_1(\wp(z)) + R_2(\wp(z))\wp'(z) \) where \( R_1 \) and \( R_2 \) are rational functions of \( \wp \).

The numbers

\[
g_2 = 60 \sum_{m, n \in \mathbb{Z} \setminus \{0,0\}} \frac{1}{(2m\omega_1 + 2n\omega_3)^4},
\]

\[
g_3 = 140 \sum_{m, n \in \mathbb{Z} \setminus \{0,0\}} \frac{1}{(2m\omega_1 + 2n\omega_3)^6}
\]

are called the invariants of \( \wp \). Since the coefficients of the Laurent expansions of \( \wp(z) \) and \( \wp'(z) \) at \( z = 0 \) are polynomials of \( g_2 \) and \( g_3 \) with rational coefficients, the function \( \wp(z; \omega_1, \omega_3) \) is also uniquely characterized by its invariants \( g_2 \) and \( g_3 \). One frequently also uses the notation \( \wp(z|g_2, g_3) \).
The function $\varphi(z)$ satisfies the first order differential equation
\[(2.5.1) \quad \varphi'(z)^2 = 4\varphi(z)^3 - g_2\varphi(z) - g_3\]
and hence the equations
\[\varphi''(z) = 6\varphi(z)^2 - g_2/2 \quad \text{and} \quad \varphi'''(z) = 12\varphi'(z)\varphi(z)\]
which shows that $-2\varphi$ is a stationary solution of the KdV equation.

The function $\varphi'$, being of order 3, has three zeros in $\Delta$. Since $\varphi'$ is odd and elliptic it is obvious that these zeros are the half-periods $\omega_1, \omega_2 = \omega_1 + \omega_3$ and $\omega_3$. Let $e_j = \varphi(\omega_j)$, $j = 1, 2, 3$. Then (2.5.1) implies that $4e_j^3 - g_2e_j - g_3 = 0$ for $j = 1, 2, 3$. Therefore
\[
0 = e_1 + e_2 + e_3,
\]
\[
g_2 = -4(e_1e_2 + e_1e_3 + e_2e_3) = 2(e_1^2 + e_2^2 + e_3^2),
\]
\[
g_3 = 4e_1e_2e_3 = \frac{4}{3}(e_1^3 + e_2^3 + e_3^3).
\]

Weierstrass also introduced two other functions denoted by $\zeta$ and $\sigma$. The Weierstrass $\zeta$-function is defined by
\[
\frac{d}{dz}\zeta(z) = -\varphi(z), \quad \lim_{z \to 0} (\zeta(z) - \frac{1}{z}) = 0.
\]
It is a meromorphic function with simple poles at $2n\omega_1 + 2n\omega_3, m, n \in \mathbb{Z}$ having residues 1. It is not periodic but quasi-periodic in the sense that
\[
\zeta(z + 2\omega_j) = \zeta(z) + 2\eta_j, \quad j = 1, 2, 3, 4,
\]
where $\eta_j = \zeta(\omega_j)$ for $j = 1, 2, 3$ and $\eta_4 = 0$.

The Weierstrass $\sigma$-function is defined by
\[
\frac{\sigma'(z)}{\sigma(z)} = \zeta(z), \quad \lim_{z \to 0} \frac{\sigma(z)}{z} = 1.
\]
$\sigma$ is an entire function with simple zeros at the points $2m\omega_1 + 2n\omega_3, m, n \in \mathbb{Z}$. Under translation by a period $\sigma$ behaves according to
\[
\sigma(z + 2\omega_j) = -\sigma(z)e^{2\eta_j(z + \omega_j)}, \quad j = 1, 2, 3.
\]

Next we recall the following fundamental theorems.

**Theorem 2.2.** Given an elliptic function $f$ with fundamental periods $2\omega_1$ and $2\omega_3$, let $b_1, \ldots, b_r$ be the distinct poles of $f$ in $\Delta$. Suppose the principal part of the Laurent expansion near $b_k$ is given by
\[
\sum_{j=1}^{\beta_k} \frac{A_{j,k}}{(z - b_k)^j}, \quad k = 1, \ldots, r.
\]
Then
\[
f(z) = C + \sum_{k=1}^{r} \sum_{j=1}^{\beta_k} (-1)^{j-1} \frac{A_{j,k}}{(j-1)!} \zeta^{(j-1)}(z - b_k),
\]
where $C$ is a suitable constant and $\zeta$ is constructed from the fundamental periods $2\omega_1$ and $2\omega_3$. Conversely, every such function is an elliptic function if $\sum_{k=1}^{r} A_{1,k} = 0$. 
Theorem 2.3. Given an elliptic function $f$ of order $n$ with fundamental periods $2\omega_1$ and $2\omega_3$, let $a_1, \ldots, a_n$ and $b_1, \ldots, b_n$ be the zeros and poles of $f$ in $\Delta$ repeated according to their multiplicities. Then

$$f(z) = C \frac{\sigma(z - a_1) \cdots \sigma(z - a_n)}{\sigma(z - b_1) \cdots \sigma(z - b_{n-1}) \sigma(z - b'_n)},$$

where $C$ is a suitable constant, $\sigma$ is constructed from the fundamental periods $2\omega_1$ and $2\omega_3$ and where

$$b'_n - b_n = (a_1 + \ldots + a_n) - (b_1 + \ldots + b_n)$$

is a period of $f$. Conversely, every such function is an elliptic function.

Finally, we turn to elliptic functions of the second kind, the central object in our analysis. A meromorphic function $\psi : \mathbb{C} \to \mathbb{C} \cup \{ \infty \}$ for which there exist two complex constants $\omega_1$ and $\omega_3$ with non-real ratio and two complex constants $\rho_1$ and $\rho_3$ such that for $i = 1, 3$

$$\psi(z + 2\omega_i) = \rho_i \psi(z)$$

is called elliptic of the second kind. We call $2\omega_1$ and $2\omega_3$ the quasi-periods of $\psi$. Together with $2\omega_1$ and $2\omega_3$, $2m_1\omega_1 + 2m_3\omega_3$ are also quasi-periods of $\psi$ if $m_1$ and $m_3$ are integers. If every quasi-period of $\psi$ can be written as an integer linear combination of $2\omega_1$ and $2\omega_3$ then these are called fundamental quasi-periods.

Theorem 2.4. A function $\psi$ which is elliptic of the second kind and has fundamental quasi-periods $2\omega_1$ and $2\omega_3$ can always be put in the form

$$\psi(z) = C \exp(\lambda z) \frac{\sigma(z - a_1) \cdots \sigma(z - a_n)}{\sigma(z - b_1) \cdots \sigma(z - b_n)}$$

for suitable constants $C$, $\lambda$, $a_1, \ldots, a_n$ and $b_1, \ldots, b_n$. Here $\sigma$ is constructed from the fundamental periods $2\omega_1$ and $2\omega_3$. Conversely, every such function is elliptic of the second kind.

Theorem 2.5. Given numbers $\alpha_1, \ldots, \alpha_m$ and $\beta_1, \ldots, \beta_m$ such that $\beta_k \neq \beta_\ell (\text{mod } \Delta)$ for $k \neq \ell$, the following identity holds

$$\prod_{j=1}^{m} \frac{\sigma(x - \alpha_j)}{\sigma(x - \beta_j)} = \sum_{j=1}^{m} \prod_{k=1}^{m} \frac{\sigma(\beta_j - \alpha_k)}{\prod_{\ell=1, \ell \neq j}^{m} \sigma(\beta_j - \beta_\ell)} \frac{\sigma(x - \beta_j + \beta - \alpha)}{\sigma(x - \beta_j) \sigma(\beta - \alpha)},$$

where

$$\alpha = \sum_{j=1}^{m} \alpha_j \text{ and } \beta = \sum_{j=1}^{m} \beta_j$$

and $\sigma$ is constructed from the fundamental periods $2\omega_1$ and $2\omega_3$.

Sketch of proof. Since this result seems less familiar than those above we will briefly sketch its proof. Denote the left and right-hand sides of (2.5.3) by $f$ and $g$, respectively. Both $f$ and $g$ are elliptic functions of the second kind associated with the same Floquet multipliers (with respect to translations by $2\omega_k$, $k = 1, 3$). Their quotient is therefore an elliptic function. Also $f$ and $g$ have the same poles and zeros taking multiplicities into account. The statement about the zeros is a consequence of the identity (cf. [235], p. 451)

$$\sum_{j=1}^{m} \left( \prod_{k=1}^{m} \frac{\sigma(\gamma_j - \delta_k)}{\prod_{\ell=1, \ell \neq j}^{m} \sigma(\gamma_j - \gamma_\ell)} \right) = 0 \text{ if } \sum_{j=1}^{m} \gamma_j = \sum_{j=1}^{m} \delta_j.$$


Therefore \( f/g \) is entire and hence constant. Since all residues are equal to one \( f = g \).

2.6. Hill’s Equation and its Spectral Theory. The study of linear homogeneous differential equations with periodic coefficients predates the late nineteenth century, but the equation \((Ly)(x) = y''(x) + q(x)y(x) = zy(x)\), where \( q(x) \) is a continuous, real-valued, periodic function of a real variable \( x \) and \( z \) is a real parameter, has generally been called Hill’s equation since its appearance in the study of the lunar perigee by Hill \([14]\) in 1877. In addition to its applications in celestial mechanics, this equation has found countless applications in quantum mechanics, where it becomes Schrödinger’s equation and is used, for instance, to model crystal structures of solids.

Periodic differential equations are usually studied by applying Floquet theory. Floquet theory (first developed by Floquet \([71, 72]\) starting in 1880) specifies the general structure of solutions of systems of periodic differential equations. Consider the equation \( \tilde{y}'(x) = Q(x)y(x) \), where \( Q \) is an \( n \times n \) matrix whose entries are continuous (for simplicity) and periodic with period \( \Omega > 0 \) and \( y(x) \) is \( \mathbb{C}^n \)-valued. Let \( Y \) be the space of solutions of \( \tilde{y}'(x) = Q(x)y(x) \) and \( T_{\Omega} \) the restriction of \( y \mapsto y(\cdot + \Omega) \) to \( Y \). Floquet theory then amounts to the study of the operator \( T_{\Omega} \).

Since \( T_{\Omega} \) maps the \( n \)-dimensional vector space \( Y \) to itself, the problem is reduced to a problem in linear algebra. The eigenvalues and eigenfunctions of \( T_{\Omega} \) are called Floquet multipliers and Floquet functions, respectively. For general references on Floquet theory see, for instance, Arscott \([10]\), Coddington and Levinson \([15]\), Ch. 3, Eastham \([58]\), Ince \([119]\), Sect. 10.8, Magnus and Winkler \([151]\), Marchenko \([153]\), Sect. 3.4, McKean and van Moerbeke \([160]\), and Yakubovich and Starzhinskii \([240]\).

In the special case of Hill’s equation \((Ly)(x) = y''(x) + q(x)y(x) = zy(x)\), we will denote the translation operator restricted to the set of solutions \( y(z, x), y'(z, x)' \in \mathcal{Y}(z) \) of the associated first-order system, by \( T_{\Omega}(z) \).

Several differential operators (resp., boundary value problems) are studied in connection with Hill’s equation,

\[(Ly)(x) = y''(x) + q(x)y(x) = zy(x), \quad q(x + \Omega) = q(x).\]

1. The maximally defined operator \( H \) in \( L^2(\mathbb{R}) \) associated with \( L \).
2. Auxiliary operators in \( L^2([x_0, x_0 + \Omega]) \) associated with \( L \) and certain families of boundary conditions, in particular, Dirichlet boundary conditions.
3. The operators in \( L^2([x_0, x_0 + \Omega]) \) associated with \( L \) and cyclic boundary conditions \( y(x_0 + \Omega) = \exp(i\theta)y(x_0), \quad y'(x_0 + \Omega) = \exp(i\theta)y'(x_0), \theta \in [0, 2\pi] \), in particular, periodic (\( \theta = 0 \)) and anti-periodic (\( \theta = \pi \)) boundary conditions.

Floquet theory provides us with a handle on the problem of determining spectral properties of these operators: the periodic eigenvalues are given as the (necessarily real) zeros of \( \text{tr}((T_{\Omega}(z)) - 2 \), while the anti-periodic eigenvalues are given as the (necessarily real) zeros of \( \text{tr}(T_{\Omega}(z)) + 2 \), and the spectrum of \( H \) coincides with the conditional stability set \( \mathcal{S}(q) \), that is, the set of all values \( z \in \mathbb{R} \) such that \((Ly)(x) = zy(x)\) has a nontrivial bounded solution with respect to \( x \in \mathbb{R} \). This was first shown by Wintner \([238, 239]\) in 1947/48. The conditional stability set, in turn, may be characterized as the set of all \( z \in \mathbb{R} \) such that \( Ly = zy \) has a Floquet multiplier of absolute value one. Since

\[
\det(T_{\Omega}(z)) = 1,
\]
the Floquet multipliers, being the eigenvalues of $T_\Omega(z)$, are given as the zeros of $\rho^2 - \rho \text{tr}(T_\Omega(z)) + 1$ and we get

$$S(q) = \{ z \in \mathbb{R} \mid -2 \leq \text{tr}(T_\Omega(z)) \leq 2 \}.$$ 

The conditional stability set consists of countably (possibly finitely) many closed intervals (plus possibly a half-line), whose endpoints coincide with points where only one (linearly independent) Floquet solution exists (which is necessarily (anti-)periodic). This was first shown by Hamel [110] in 1913 (see also Liapunov [150], who corrected a mistake in Hamel's paper). Hence the (anti-)periodic eigenvalues determine the spectrum of $H$. From the periodic and anti-periodic eigenvalues repeated according to their multiplicity and ordered as a decreasing sequence denoted by $E_0, E_1, \ldots$ (observing that for normal operators the algebraic and geometric multiplicities of eigenvalues coincide), one infers

$$\sigma(H) = \bigcup_{j=0}^{\infty} [E_{2j+1}, E_{2j}].$$

Typically, the spectral bands $[E_{2j+1}, E_{2j}]$ are separated by spectral gaps ($E_{2j+2}, E_{2j+1}$). However, since the (anti-)periodic eigenvalues may be twofold degenerate, some gaps (or even all gaps) may close. If only finitely many gaps are present, one calls $q$, the potential coefficient in $L$, a finite-gap, or a finite-band potential.

A trivial example of a finite-band potential is of course the constant potential $q(x) = c$. The first in this century to discuss a nontrivial example was Ince [118] around 1940, who treated in depth Lamé’s potential

$$(2.6.1) \quad q(x) = -s(s+1)\varphi(x + \omega_3), \quad s \in \mathbb{N},$$

with fundamental half periods $\omega_1 \in \mathbb{R}$ and $\omega_3 \in i\mathbb{R}$. Under these conditions $q$ is real-valued and real-analytic in $x$. Ince showed that for all but $2n+1$ values of $z \in \mathbb{R}$, the equation $y''(x) + q(x)y(x) = zy(x)$ has two linearly independent (anti-)periodic eigenvalues, that is, $q$ in (2.6.1) is a finite-band potential (see also Akhiezer [1], Erdelyi [66], Turbiner [228], and Ward [232]). However, a closer look at the classical works of Hermite [113] and Halphen [109] (see also Klein [225] on Lamé’s equation at the end of last century reveals that these results were actually well-known (but not yet put in a Floquet-theoretic language) as can be inferred, for instance, from Whittaker-Watson’s treatise [223], Sect. 23.7.

In 1909 Birkhoff [23] compared eigenvalues of various boundary value problems associated with $L$ on a finite interval. His results show that there is precisely one Dirichlet eigenvalue in each of the intervals $[E_{2j+2}, E_{2j+1}]$ (i.e., in the closure of the gaps). In contrast to the (anti-)periodic eigenvalues, the Dirichlet eigenvalues associated with the interval $[x_0, x_0 + \Omega]$ vary with $x_0$ in $[E_{2j+2}, E_{2j+1}]$ if $E_{2j+2} < E_{2j+1}$. Writing $\mathbb{R} = \bigcup_{n \in \mathbb{Z}} [x_0 + n\Omega, x_0 + (n+1)\Omega]$, a comparison of the Dirichlet problems on $[x_0, x_0 + \Omega]$ and on $(-\infty, x_0) \cup (x_0, \infty)$ assuming $q(x)$ to be $\Omega$-periodic, yields the same Dirichlet spectra in the spectral gaps $(E_{2j+2}, E_{2j+1})$ with $E_{2j+2} < E_{2j+1}$. Explicit formulas for the monotone rate of change of various kinds of eigenvalues (including Dirichlet eigenvalues) with respect to varying $x_0 \in \mathbb{R}$ can be found in Kong and Zettl [126, 127] and the literature therein. In [89] this phenomenon is related to Green’s functions and rank-one perturbations of resolvents.

If $q$ is a locally integrable, real-valued, periodic function with period $\Omega > 0$ we therefore have equivalence of the following statements:
(A) $q$ is a finite-band potential.
(B) For only finitely many values of $z$ the differential equation $y''(x) + q(x)y(x) = zy(x)$ fails to have two linearly independent Floquet solutions.
(C) The Dirichlet boundary value problem on the interval $[x_0, x_0 + \Omega]$ has only a finite number of eigenvalues depending on $x_0$.

Spectral theoretic aspects of complex-valued periodic potentials $q$ have been investigated by Rofe-Beketov [188] and Tkachenko [211] (see also Birnir [24], [25], Kotani [131], McGarvey [157]–[159], Sansuc and Tkachenko [190]–[192], and Tkachenko [213], [213], [213] for recent results). They found that the spectrum of $H$ and the conditional stability set of $L$ still coincide, that is,

$$ S(q) = \{z \in \mathbb{C} \mid -2 \leq \text{tr}(T_0(z)) \leq 2\}.$$

Since the conditional stability set is given as the preimage of $[-2, 2]$ under an entire function, it turns out that the spectrum of $H$ consists of countably many (possibly finitely many) regular analytic arcs. While the term “finite-gap potential” is now rendered meaningless, it still makes perfect sense to call a potential a finite-band potential if the spectrum is a finite union of regular analytic arcs. In contrast to the real-valued case, their spectral arcs can now cross each other, see, for instance, [93], [176] for explicit examples exhibiting this phenomenon.

Returning to the real-valued case, the isospectral set $I(q_0)$ of a given periodic potential $q_0 \in C(\mathbb{R})$ with period $\Omega > 0$ (i.e., the set of all $\Omega$-periodic $q \in C(\mathbb{R})$ whose $2\Omega$-periodic eigenvalues coincide with that of $q_0$) turns out to be a manifold, in fact, a (generally infinite dimensional) torus generated by a product of circles. Each circle is uniquely associated with a spectral gap $(E_{2j+2}, E_{2j+1})$, $E_{2j+2} < E_{2j+1}$ and the periodic motion of the Dirichlet eigenvalue in the gap $(E_{2j+2}, E_{2j+1})$ corresponding to the interval $[x, x + \Omega]$ as a function of $x \in \mathbb{R}$. In the special case where $q$ is real-valued and periodic, and the $L^2$-spectrum of $-d^2/dx^2 + q(x)$ is a half-line $(-\infty, E_0]$, Borg [27] proved the celebrated inverse spectral result $q(x) = E_0$ for all $x \in \mathbb{R}$. A quick proof of this uniqueness result and extensions to real-valued reflectionless potentials $q$ with associated half-line spectrum $(-\infty, E_0]$ follows from the trace formula proved in [33] (observing $\xi(\lambda, x) = 1/2$ for $\lambda < E_0$ in the corresponding trace formula (3.1) for $q(x)$ in [33]). Similarly, if $q$ is real-valued and periodic, and the $L^2$-spectrum of $-d^2/dx^2 + q(x)$ is of the type $(-\infty, E_2] \cup [E_1, E_0]$, $E_2 < E_1$, Hochstadt [116] proved that $q(x) = -2\varphi(x + \omega_3 + \alpha)$, that is, the Lamé potential (2.6.1) for $n = 1$, where $\varphi(x)$ denotes the Weierstrass elliptic function associated with some period lattice $\omega_1 > 0$, $-i\omega_3 > 0$ (depending on $E_0$, $E_1$, and $E_2$) and $\alpha \in \mathbb{R}$. The isospectral torus for real-valued periodic potentials with three or more spectral bands is described, for instance, in Buys and Finkel [34], Finkel, Isaacson, and Trubowitz [68], Gesztesy, Simon, and Teschl [90], Gesztesy and Weikard [92], Iwasaki [24], McKeve van Moerbeke [167], and McKeve and Trubowitz [161].

It should be emphasized in this context that the assumption of real-valuedness of $q(x)$ cannot be dropped as shown by the well-known example $q(x) = \exp(ix)$, with associated $L^2$-spectrum $(-\infty, 0]$ (see the paragraph following Remark 8.3 for more details).

2.7. Periodic KdV Potentials. In 1974, Novikov [169] investigated the Cauchy problem of the KdV equation in the case of periodic initial data. He noted that the right generalization of multi-soliton solutions, which are stationary solutions of
appropriate higher-order KdV equations, are the finite-band potentials, that is, he proved the following result.

**Theorem 2.6.** If $q$ is a real-valued, periodic, stationary solution of an $n$-th order KdV equation, then the $L^2(\mathbb{R})$-spectrum associated with $d^2/dx^2 + q(x)$ has at most $n$ finite bands.

Within a year, Dubrovin [53] and Flaschka [69] also proved the converse, and Dubrovin and Novikov [57] used these results to solve the Cauchy problem of the KdV equation in the case of finite-band initial data. Hence the statement,

(D) $q$ is a stationary solution of a higher-order KdV equation,

is equivalent to any of the statements (A)–(C) in Section 2.6 for real-valued $q$.

Moreover, it became clear from these investigations that generically, finite-band potentials will only be quasi-periodic with respect to $x$ and not periodic in $x$.

About the same time, Its and Matveev [123] derived their celebrated formula for finite-gap solutions $q(x, t)$ in terms of the Riemann theta function associated with the underlying hyperelliptic curve and a fixed homology basis on it (cf. [17], Ch. 3, [87], Ch. 1, [170], Ch. II). Subsequent extensions of this formula to general (matrix-valued) integrable systems were developed by Dubrovin [54]–[56] and Krichever [133], [134], [137]. A new approach to finite-band solutions of the KdV hierarchy in terms of Kleinian functions was recently developed by Buchstaber, Enol’skii, and Leykin [28], [29].

### 2.8. Elliptic KdV Potentials

While the considerations of the preceding subsections pertain to general solutions of the stationary KdV hierarchy, we now concentrate on the additional restriction that $q$ be an elliptic function and hence return to our main subject, elliptic finite-band potentials $q$ for $L = d^2/dx^2 + q(x)$, or, equivalently, elliptic solutions of the stationary KdV hierarchy. The remarkable finite-gap example of the Lamé potentials (2.6.1) due to Hermite [113] and Halphen [109] in the last century, brought back into the limelight by Ince [118] around 1940, remained the only explicit elliptic finite-gap example until the KdV flow

$q_t = \frac{1}{4} q_{xxx} + \frac{3}{2} q_x,$

with initial condition $q(x, 0) = -6\wp(x + \omega_3)$, was explicitly integrated by Dubrovin and Novikov [57] in 1975 (see also Enol’skii [61], [63], Its and Enol’skii [122]) and found to be of the type

$$q(x, t) = -2 \sum_{j=1}^{3} \wp(x - x_j(t))$$

for appropriate $\{x_j(t)\}_{1 \leq j \leq 3}$. Due to the unitary evolution operator $U_n(t)$ constructed with the help of $P_{2n+1}(t)$ via $U_{n,1}(t) = P_{2n+1}(t)U_n(t)$, all potentials $q(x, t)$ in (2.8.1) are isospectral to $q(x, 0) = -6\wp(x + \omega_3)$.

In 1977, Airault, McKean and Moser, in their seminal paper [5], presented the first systematic study of the isospectral torus $I_{\mathbb{R}}(q_0)$ of real-valued smooth potentials $q_0(x)$ of the type

$$q_0(x) = -2 \sum_{j=1}^{M} \wp(x - x_j)$$

with a finite-gap spectrum. Among a variety of results they proved that any element $q$ of $I_{\mathbb{R}}(q_0)$ is an elliptic function of the type (2.8.2) (with different $x_j$), with $M$ constant throughout $I_{\mathbb{R}}(q_0)$ and $\dim I_{\mathbb{R}}(q_0) \leq M$. In particular, if $q_0$ evolves
according to any equation of the KdV hierarchy it remains an elliptic finite-gap potential. The potential (2.8.2) is intimately connected with completely integrable many-body systems of the Calogero-Moser-type [35], [166] (see also Bennequin [22], Birnir [26], Calogero [36], Chudnovsky and Chudnovsky [39], Olshanetsky and Perelomov [171], and Ruijsenaars [189]). This connection with integrable particle systems was subsequently exploited by Krichever [136] in his construction of elliptic algebro-geometric solutions of the Kadomtsev-Petviashvili equation. In the KdV context of (2.8.2), Krichever’s approach relies on the ansatz
\[
\psi_0(z,x) = e^{\kappa(z)x} \sum_{j=1}^{M} A_j(z) \Phi(x - x_j, \rho(z)),
\]
for the Floquet solutions of \( L_0 = d^2/dx^2 + q_0(x) \), where
\[
\Phi(x, \rho) = \frac{\sigma(x - \rho)}{\sigma(x) \sigma(-\rho)} e^{\sigma(\rho)x}
\]
(assuming for simplicity the generic case \( x_j \neq x_k \mod \Delta \) for \( j \neq k \)). Applying \( L_0 \) to (2.8.3) then yields an M-sheeted covering of the torus associated with the fundamental periods \( 2\omega_1, 2\omega_3 \) and hence a description of the underlying algebraic curve. (We will briefly comment on this ansatz in Remark 3.9.)

The next breakthrough occurred in 1988 when Verdier [230] published new explicit examples of elliptic finite-gap potentials. Verdier’s examples spurred a flurry of activities and inspired Belokolos and Enol’skii [19], [20], Smirnov [196], and subsequently Taimanov [208] and Kostov and Enol’skii [130] to find further such examples by combining the reduction process of Abelian integrals to elliptic integrals (see Babich, Bobenko and Matveev [13], [14], Belokolos, Bobenko, Enol’skii, Its, and Matveev [17], Ch. 7, and Belokolos, Bobenko, Mateev, and Enol’skii [18]) with the aforementioned techniques of Krichever [136], [137]. This development finally culminated in a series of recent results of Treibich and Verdier [220]–[225] where it was shown that a general complex-valued potential of the form
\[
q(x) = -\sum_{j=1}^{4} d_j \wp(x - \omega_j)
\]
(\( \omega_2 = \omega_1 + \omega_3, \omega_4 = 0 \)) is a finite-gap potential if and only if \( d_j/2 \) are triangular numbers, that is, if and only if
\[
d_j = s_j(s_j + 1) \text{ for some } s_j \in \mathbb{Z}, 1 \leq j \leq 4.
\]
We shall from now on refer to potentials of the type
\[
q(x) = -\sum_{j=1}^{4} s_j(s_j + 1) \wp(x - \omega_j) \text{, } s_j \in \mathbb{Z}, 1 \leq j \leq 4
\]
as Treibich-Verdier potentials. The methods of Treibich and Verdier (see also Colombo, Pirola, and Previato [14], Previato [184], [185], Previato and Verdier [187]) are based on hyperelliptic tangent covers of the torus \( \mathbb{C}/\Lambda \) (\( \Lambda \) being the period lattice generated by \( 2\omega_1 \) and \( 2\omega_3 \)).

The state of the art of elliptic finite-gap solutions up to 1993 was recently reviewed in a special issue of Acta Applicandae Math., see, for instance, Belokolos and Enol’skii [21], Enol’skii and Kostov [19], Krichever [139], Smirnov [198], Taimanov [209], and Treibich [221]. For more recent results see Eilbeck and Enol’skii [59],
Since a complete characterization of all elliptic finite-gap solutions of the stationary KdV hierarchy was still open at that time, we developed a new approach to this characterization problem to be described in Sections 3 and 4. As alluded to at the end of our introduction, Calogero-Moser-type models are again an intensive object of study.

Since we will also discuss stationary rational solutions of the KdV hierarchy in Section 3 we should mention the case where $\phi(x)$ in (2.8.2) degenerates into $x^{-2}$ as discussed, for instance by Airault, McKean, and Moser [10], Krichever [133], [138], Moser [160], [167], Pelinovsky [176], and Shiota [195].

2.9. Linear Differential Equations in the Complex Domain. While all the developments described in previous subsections were in place around 1993, one final point, the connection of this subject to the classical area of differential equations in the complex domain, was made only around 1994 when we started to work on [108]. In the following we will remind the reader about this fundamental, but thus far missing piece, which plays a decisive role in the remainder of this review.

In the late 1830’s, Lamé studied Laplace’s equation in confocal coordinates. After some appropriate changes of variables this led to the differential equation (2.6.1), it’s Weierstrass form. At the end of the last century, Lamé’s equation (especially in Jacobi’s form) was studied intensively by Hermite [112], who obtained the general solution of the equation for integer values of $n$ and any value of $z$. Picard proceeded to consider first general second-order equations [180] (see also Floquet [71], [72], Mittag-Leffler [163], and Halphen [107], [108], [109]), and finally first-order systems [182] whose coefficients are elliptic functions. Consider the differential equation

$$y'(x) = Q(x)y(x),$$

(2.9.1)

where $y(x)$ is $C^n$-valued and $Q$ is a $n \times n$ matrix whose entries are elliptic functions with a common period lattice spanned by the fundamental periods $2\omega_1$ and $2\omega_3$. Let $\mathcal{Y}$ be the space of solutions of (2.9.1) and denote the restriction of the translation operator $y \mapsto y(\cdot + 2\omega_j)$ to $\mathcal{Y}$ by $T_j$. Using this notation, Picard’s theorem reads as follows (see also Akhiezer [6], Sects. 58, 59, Burkhardt [33], Ch. 15, Forsyth [76], Ch. IX, Gray [103], Sect. 6.1, Halphen [107], [108], [109], Ch. XIII, Ince [119], Sect. 15.6, Krause [132], Vol. 2, Ch. 3, and Picard [183], Sect. III.V).

**Theorem 2.7.** Assume that the first-order system (2.9.1) has a meromorphic fundamental system of solutions. Then there exists at least one solution $\psi$ which is elliptic of the second kind, that is, the components of $\psi$ are meromorphic and $T_j \psi = \rho_j \psi$ for $j = 1, 3$ for suitable constants $\rho_1, \rho_3 \in \mathbb{C}\backslash\{0\}$. If in addition, one of the operators $T_1$ and $T_3$ has distinct eigenvalues, then there exists a fundamental system of solutions of (2.9.1) which are elliptic of the second kind.

The explicit Floquet-type structure of solutions of (2.9.1) in terms of a doubly periodic vector, powers of $x$, powers of the Weierstrass zeta-function, and an exponential contribution, has recently been determined in [11].

About the time Floquet, Fuchs, Hermite, Mittag-Leffler, and Picard (cf. the historical discussions in Gray [103], Ch. VI) developed the theory of differential equations with elliptic coefficients, Floquet [71], [72] also established his celebrated results for linear, homogeneous differential equations with simply-periodic coefficients. From a historical perspective it is perhaps interesting to note that Floquet
assumed the coefficients, as well as the general solution of the equation, to be meromorphic in order to arrive at the existence of periodic solutions of the second kind, that is, he obtained the precise analog of Picard’s result. Only later was it realized that his theorem extends to continuous periodic coefficients on \( \mathbb{R} \) without any reference to meromorphic fundamental systems. The solutions Floquet called “periodic of the second kind”, are today generally called Floquet solutions.

Next, we mention another theorem with a similar flavor that concerns differential equations with rational coefficients and meromorphic fundamental systems of solutions and hence is applicable to the study of rational algebro-geometric KdV potentials. In 1885, Halphen \cite{halphen} published the following result.

**Theorem 2.8.** Assume that the differential equation

\[
y^{(n)}(x) + q_1(x)y^{(n-1)}(x) + \cdots + q_n(x)y(x) = 0
\]

has rational coefficients bounded at infinity and a meromorphic fundamental system of solutions. Then the general solution is of the form

\[
R_1(x) \exp(\lambda_1 x) + \cdots + R_n(x) \exp(\lambda_n x),
\]

where \( R_1, \ldots, R_n \) are rational functions.

It should be emphasized that the principal hypothesis in Theorems 2.7 and 2.8, the existence of a meromorphic fundamental system of solutions, can be verified in a straightforward manner by applying the Frobenius method (see, e.g., Coddington and Levinson \cite{coddington}, Sect. 4.8, Forsyth \cite{forsyth}, Ch. III, Hille \cite{hille}, Ch. 9, Ince \cite{ince}, Ch. XVI, and Whittaker and Watson \cite{whittaker}, Ch. X) to each pole of \( q \) in the fundamental period parallelogram.

Finally we mention an observation made by Appell \cite{appell}. Let \( y_1(x) \) and \( y_2(x) \) be linearly independent solutions of

\[
y''(x) + u(x)y(x) = 0.
\]

Then \( y_1(x)^2, y_1(x)y_2(x), \) and \( y_2(x)^2 \) are linearly independent solutions of

\[
w'''(x) + 4u(x)w'(x) + 2u'(x)w(x) = 0.
\]

This equation is easily integrated and yields

\[
g'(x)^2 - 2g(x)g''(x) - 4u(x)g(x) = W(y_1, y_2)^2,
\]

where \( g(x) \) is the product of any two solutions \( y_1(x) \) and \( y_2(x) \) of \( y''(x) + u(x)y(x) = 0 \) and \( W(y_1, y_2) \) denotes their (x-independent) Wronskian. This innocent looking fact will be of great importance in our analysis later. In fact, a comparison of equations (2.3.4) and (2.9.2) reveals another connection between the KdV hierarchy and the Schrödinger-type equation

\[
y''(x) + u(x)y(x) = 0.
\]

Moreover, since the formal Green’s function \( G(z, x, x') \) of \( \frac{d^2}{dx^2} + q(x) \) on the diagonal \( x = x' \) is of the type \( y_1(z, x)y_2(z, x)/W(y_1(z, x), y_2(z)) \), \( \ref{2.9.3} \), with \( u(x) = q(x) - z \), is equivalent to the well-known universal nonlinear second-order differential equation satisfied by \( G(z, x, x) \) (see, also Gelfand and Dickey \cite{gelfand_dickey}).

It should be noted that Drach \cite{drach}, \cite{drach2} (see also \cite{drach3}) used \( \ref{2.9.2} \) to derive a class of completely integrable systems now known as the stationary KdV hierarchy as early as 1918/19. It appears he was the first to make the explicit connection between completely integrable systems and spectral theory. More than 55 years later, Gelfand and Dickey \cite{gelfand_dickey}, \cite{gelfand_dickey2} also based some of their celebrated work on the KdV hierarchy on \( \ref{2.9.2} \).

### 3. Algebro-Geometric, and Especially, Elliptic KdV Potentials.

**Definition 3.1.** Suppose \( q \) is meromorphic and let \( L \) be the differential expression \( L = \frac{d^2}{dx^2} + q(x) \). Then \( q \) is called an algebro-geometric KdV potential (or
simply \textit{algebro-geometric}) if \( q \) is a solution of some equation of the stationary KdV hierarchy.

Equivalently, we could define the meromorphic function \( q \) as algebro-geometric if any one of the following three conditions is satisfied.

1. There exists an odd-order differential expression \( P \) such that \([P, L] = 0\) (according to the result of Burchnall and Chaundy).
2. There exists an ordinary differential expression \( P \) of odd order and a polynomial \( R \) such that \( P^2 = R(L) \).
3. There exists a function \( F : \mathbb{C}^2 \to \mathbb{C}_\infty \), which is a polynomial in the first variable, meromorphic in the second, and which satisfies \( F'''(z, x) + 4q(x) - z)F''(z, x) + 2q'(x)F(z, x) = 0 \) (cf. equation (2.3.4)).

It can be shown (see Theorem 6.10 by Segal and Wilson \cite{194}) that any solution of any of the stationary KdV equations is necessarily meromorphic. Hence the assumption that \( q \) is meromorphic is actually redundant in Definition 3.1.

3.1. Periodic KdV Potentials. If \( q \) is real-valued, locally integrable, and periodic we obtain the equivalence of the following statements from the works of Birkhoff and Hamel described in Section 2.6.

1. \( q \) is algebro-geometric with \( P^2 = \prod_{j=0}^{2n}(L - \lambda_j) \).
2. \( q \) is finite-band with spectrum the union of \((-\infty, \lambda_{2n}] \) and \( n \) compact bands \([\lambda_{2j+1}, \lambda_{2j}], j = 0, \ldots, n - 1 \).
3. \( Ly = zy \) has two linearly independent Floquet solutions for all \( z \in \mathbb{C} \) with the exception of the \( 2n + 1 \) values \( \lambda_0, \ldots, \lambda_{2n} \).
4. \( q \) has \( n \) movable Dirichlet eigenvalues, precisely one in each of the closures of the spectral gaps \( (\lambda_{2j}, \lambda_{2j-1}), j = 1, \ldots, n \).

As mentioned previously, the classical example for finite-band potentials are the Lamé potentials \( q(x) = -s(s + 1)\varphi(x + c; \omega_1, \omega_3), n \in \mathbb{N}, \) with \( c = \omega_3 \) purely imaginary and \( \omega_1 \) real. However, Lamé potentials are algebro-geometric for general choices of the half-periods as well as for general choices of \( c \in \mathbb{C} \). This suggests the study of complex-valued potentials with inverse square singularities as in \cite{18} and \cite{233}, and we will subsequently report on some of these results.

Let \( q \) be a complex-valued, periodic function with period \( \Omega > 0 \), which is locally integrable on \( \mathbb{R} \setminus \Sigma \), where \( \Sigma \subset \mathbb{R} \) is a discrete set (i.e., a set without finite accumulation points). Moreover, \( q \) is assumed to be meromorphic near each \( \xi \in \Sigma \) with principal part \(-s(s+1)/(x-\xi)^2\), where \( s = s(\xi) \in \mathbb{N} \). Then one can define unique solutions of initial value problems of the differential equation \( y''(x) + q(x)y(x) = zy(x) \) on \( \mathbb{R} \setminus \Sigma \) (with initial conditions at \( x_0 \in \mathbb{R} \setminus \Sigma \)) by analytic continuation around the singularities. Even though the potential is no longer continuous, Floquet theory (see Section 2.6) remains essentially unchanged.

To apply Floquet theory we first introduce a basis in \( \mathcal{Y}(z) \), the space of solutions of \( y''(x) + q(x)y(x) = zy(x) \). Let \( c(z, \cdot, x_0), s(z, \cdot, x_0) \in \mathcal{Y}(z) \) be defined by the initial conditions \( c(z, x_0, x_0) = s'(z, x_0, x_0) = 1 \) and \( c'(z, x_0, x_0) = s(z, x_0, x_0) = 0 \) (prime denoting the derivative with respect to the second variable). Let \( \rho_k \) denote the Floquet multipliers of \( y''(x) + q(x)y(x) = zy(x) \), that is, the eigenvalues of the translation operator \( T_{\Omega}(z) \). An important role is played by \( \text{tr}(T_{\Omega}(z)) \), which is sometimes called the Floquet discriminant and which, in our basis, is given by

\[
\text{tr}(T_{\Omega}(z)) = c(z, x_0 + \Omega, x_0) + s'(z, x_0 + \Omega, x_0).
\]
Since the trace is of course independent of the chosen basis in \( Y(z) \), the dependence of the right-hand side on \( x_0 \) is only apparent. The Floquet solutions, may be expressed as

\[
f_{\pm}(z, x_0, x) = s(z, x_0 + \Omega, x_0) c(z, x_0, x) + (\rho_{\pm} - c(z, x_0 + \Omega, x_0)) s(z, x_0, x),
\]

and their Wronskian is given by

\[
W(f_+(z, \cdot, x_0), f_-(z, \cdot, x_0)) = -s(z, x_0 + \Omega, x_0) \sqrt{(\text{tr}(T\Omega(z)))^2 - 4}.
\]

Next, consider the function

\[
g(z, x) = \frac{f_+(z, x, x_0) f_-(z, x, x_0)}{W(f_+(z, \cdot, x_0), f_-(z, \cdot, x_0))}.
\]

As our notation for \( g(z, x) \) suggests, the dependence of this function on \( x_0 \) is only apparent since \( f_{\pm}(z, \cdot, x_1) \) are just multiples of \( f_{\pm}(z, \cdot, x_0) \) and the right-hand side of (3.1.1) is independent of normalization. In particular, we may replace \( x_0 \) by \( x \) in (3.1.1) to obtain

\[
g(z, x) = \frac{-s(z, x + \Omega, x)}{\sqrt{(\text{tr}(T\Omega(z)))^2 - 4}}.
\]

The function \( s(\cdot, x + \Omega, x) \) is an entire function of order of growth 1/2. The zeros of \( s(\cdot, x + \Omega, x) \) are the Dirichlet eigenvalues of \( d^2/dx^2 + q \) on the interval \([x, x + \Omega]\) and their order, which we denote by \( d(z, x) \), is the algebraic multiplicity of the corresponding Dirichlet eigenvalue (we set \( d(z, x) = 0 \) if \( s(z, x + \Omega, x) \neq 0 \)). We also introduce \( d_i(z) = \min\{d(z, x) \mid x \in \mathbb{R} \setminus \Sigma\} \) and \( d_m(z, x) = d(z, x) - d_i(z) \), the immovable and movable parts of \( d(z, x) \), respectively. The quantity \( \sum_{z \in \mathbb{C}} d_m(z, x) \), which is independent of \( x \), is called the number of movable Dirichlet eigenvalues. Using Hadamard’s factorization theorem we write \( g(z, x) = F(z, x) D(z) \) collecting in \( F \) the factors depending on \( x \) and in \( D \) the factors independent of \( x \). Then the multiplicity of a zero \( z \) of \( D \) is just \( d_i(z) \), while the multiplicity of a zero \( z \) of \( F(\cdot, x) \) is \( d_m(z, x) \). One then obtains from Appell’s equation (2.9.3) that

\[
F'(z, x)^2 - 2F(z, x)F''(z, x) - 4(q - z)F(z, x)^2 = ((\text{tr}(T\Omega(z)))^2 - 4)/D(z)^2.
\]

Recall that a zero \( z \) of \( (\text{tr}(T\Omega(z)))^2 - 4 \) is an (anti-)periodic eigenvalue whose multiplicity we denote by \( p(z) \). Since the left-hand side of equation (3.1.2) is entire as a function of \( z \) we obtain the following result.

**Theorem 3.2.** There exists an entire function \( R \) such that \( (\text{tr}(T\Omega(z)))^2 - 4 = R(z) D(z)^2 \). In particular, \( p(z) - 2d_i(z) \geq 0 \) for every \( z \in \mathbb{C} \).

There are, at most, countably many points where \( p(z) > 0 \) since these points are isolated. Therefore, there are at most countably many points where \( p(z) - 2d_i(z) > 0 \). These include all algebraically simple (anti-)periodic eigenvalues (where \( p(z) = 1 \)) but may well include other points too.

We call a Floquet solution \( \psi(z_0, \cdot) \) of \( y''(x) + q(x)y(x) = z_0 y(x) \) regular if there exist Floquet solutions \( \psi(z, \cdot) \) of \( y''(x) + q(x)y(x) = z y(x) \), which converge pointwise to \( \psi(z_0, \cdot) \) as \( z \) tends to \( z_0 \). It was shown in [23] that the set of regular Floquet solutions forms a line bundle on the topological space \( M_F \) obtained from the curve \( \rho^2 - (\text{tr}(T\Omega(z)))^2 \rho + 1 = 0 \) by desingularization at all points where \( p(z) = 2d_i(z) > 0 \). The space \( M_F \) can be viewed as a double cover of the complex plane branched precisely at all points \( z \) where \( p(z) - 2d_i(z) > 0 \). In particular, the equation \( y''(x) +
$q(x)y(x) = zy(x)$ has two linearly independent regular Floquet solutions if and only if $p(z) - 2d_i(z) = 0$. In other words, $p(z) - 2d_i(z) > 0$ indicates a defect in the structure of regular Floquet solutions. The number $\text{def}(q) = \sum_{z \in \mathbb{C}} (p(z) - 2d_i(z)) = \text{deg}(R)$, which is a positive integer or infinity, will therefore subsequently be called the Floquet defect of $q$. When $q$ is real-valued and nonsingular, then $p(z) - 2d_i(z)$ is always zero, except when $p(z) = 1$, which forces $d_i(z) = 0$. In this case $\text{def}(q)$ counts the number of points $z \in \mathbb{C}$ where only one linearly independent Floquet solution exists. In general, however, $y''(x) + q(x)y(x) = zy(x)$ may have two linearly independent Floquet solutions with only one being regular.

**Theorem 3.3.** The following statements are equivalent:

1. $\text{def}(q) = 2n + 1$.
2. There are $n$ movable Dirichlet eigenvalues.
3. There exists a differential expression $P_{2n+1}$ of order $2n+1$ but none of smaller odd order commuting with $L$. This differential expression satisfies $P_{2n+1}^2 = R_{2n+1}(z) = \prod_{z \in \mathbb{C}} (L - z)^{p(z) - 2d_i(z)}$ and hence $q$ is algebro-geometric.

**Sketch of proof.** Generally $p(z) \leq 2$ and $d(z, x_0) \leq 1$ for all suitably large $z$. Moreover, asymptotically, any Dirichlet eigenvalue is close to two (anti-)periodic eigenvalues. Therefore, $\text{def}(q) = 2n + 1$ implies $d_m(z, x_0)$ is different from zero for only finitely many $z \in \mathbb{C}$ and $\sum_{z \in \mathbb{C}} d_m(z, x)$, which equals the degree of $F('', x_0)$, must be finite. Equation (3.1.2) then yields $\text{deg} F('', x_0) = n$. Hence (1) implies (2). The converse of this follows immediately from equation (3.1.3).

Differentiation of equation (3.1.2) shows that the third criterion after Definition 3.1 is satisfied when $F('', x_0)$ is a polynomial. Hence (2) implies (3). We remark here that $R_{2n+1}$ is a constant multiple of $R$. To prove that (3) implies (2) one has to show that the zeros of the function $F_n$ defined by $P_{2n+1}$ are precisely the movable Dirichlet eigenvalues of $L$. This follows from applying $P_{2n+1}$ given by (3.3.3), successively to the generalized Dirichlet eigenfunctions (i.e., the eigenfunctions corresponding to the algebraic eigenspace) associated with the movable Dirichlet eigenvalues of $L$ (cf. [233] for more details).

**Remark 3.4.** The quantity $p(z) - 2d_i(z)$ proves to be of utmost importance. Determining the points where it takes on positive values, this quantity then governs the structure of the line bundle of regular Floquet solutions and determines the entire function $R(z) = ((\text{tr}(T(z)))^2 - 4)/D(z)^2$, which defines the algebraic curve associated with $q$ if $q$ is algebro-geometric and $R$ is a polynomial. When $q$ is real-valued and locally integrable on $\mathbb{R}$, then geometric and algebraic multiplicities of Dirichlet and (anti-)periodic boundary value problems coincide. Therefore, $0 \leq p(z) \leq 2$ and $0 \leq d(z, x_0) \leq 1$. Thus, $d_i(z) = 1$ enforces $p(z) = 2$. In addition, if $p(z) = 2$, then all solutions are (anti-)periodic, implying $d_i(z) = 1$. In this case the questions, “When is $p(z) - 2d_i(z) > 0$?” and “When is $p(z) = 1$?”, are equivalent. Hence, in determining the edges of the spectral gaps, the role played by the Dirichlet eigenvalues and, in particular, the distinction between movable and immovable Dirichlet eigenvalues, is secondary in the case of real-valued locally integrable potentials.

**Remark 3.5.** One may also show that every algebro-geometric potential is a finite-band potential, that is, the conditional stability set (which coincides with the spectrum of $H$ when no singularities are present) consists of finitely many regular analytic arcs.
When \( q \) is real and has no singularities, the converse is also true (Dubrovin [13]). However, in general this is not the case as the following example shows. Let \( q(x) = e^{i2x} \). Then a fundamental system of solutions of \( y''(x) + q(x)y(x) = -\nu^2 y(x) \) is given in terms of Bessel functions (cf. [2], Ch. 9) by \( y_1(x) = J_\nu(\iota \nu x), y_2(x) = Y_\nu(\iota \nu x) \). Note that \( y_1 \) is always a Floquet solution with multiplier \( e^{\iota \nu x} \). Hence \( z = -\nu^2 \) is in the conditional stability set if and only if \( \nu \in \mathbb{R} \). Consequently, \( S(L) = (-\infty, 0] \).

However, \( s(\nu^2, x_0 + \pi, x_0) = \pi J_\nu(\iota \nu x_0)J_{-\nu}(\iota \nu x_0) \), which is entire as a function of \( x_0 \). Hence \( d_i(-\nu^2) = 0 \) for all \( \nu \in \mathbb{C} \), that is, every Dirichlet eigenvalue is movable. Thus, \( \text{def}(q) = \infty \), and hence \( q \) is not algebro-geometric. More general examples of this type have been studied systematically by Gasymov [80], [81], Guillemin and Uribe [106], and Pastur and Tkachenko [174], [175].

3.2. Picard-KdV Potentials.

**Definition 3.6.** Let \( q \) be an elliptic function. Then \( q \) is called a Picard-KdV potential (or simply a Picard potential) if the equation \( y''(x) + q(x)y(x) = zy(x) \) has a meromorphic fundamental system of solutions with respect to \( x \) for all values of the spectral parameter \( z \in \mathbb{C} \).

**Theorem 3.7.** If \( q \) is a Picard potential then it may be represented as

\[
q(x) = C - \sum_{j=1}^{m} s_j(s_j + 1)\varphi(x - b_j)
\]

for suitable integers \( s_1, \ldots, s_m \) and complex numbers \( C, b_1, \ldots, b_m \), where the \( b_j \) are pairwise distinct mod \( \Delta \).

**Sketch of proof.** Every singularity of \( y''(x) + q(x)y(x) = zy(x) \) must be a regular singular point with integer indices. From the partial fraction expansion for elliptic functions (Theorem 2.2) one obtains

\[
q(x) = C - \sum_{j=1}^{m} (s_j(s_j + 1)\varphi(x - b_j) + B_j \zeta(x - b_j)).
\]

In a vicinity of \( b_j \) there is a solution of the form

\[
\psi(x) = (x - b_j)^{-s_j} \sum_{k=0}^{\infty} \beta_k (x - b_j)^k,
\]

where \( \beta_0 = 1 \). Next we use the Frobenius method to show that \( B_j = 0 \). Let

\[
q(x) = \frac{-s_j(s_j + 1)}{(x - b_j)^2} + \frac{B_j}{x - b_j} + \sum_{k=0}^{\infty} C_{j,k}(x - b_j)^k
\]

and insert \( \psi \) into the differential equation \( y'' + (q - z)y = 0 \) to get

\[
0 = f(-s_j)(x - b_j)^{-s_j} + \{ f(1 - s_j)\beta_1 + G_1 \}(x - b_j)^{1 - s_j}
\]

\[
+ \cdots + \{ f(k - s_j)\beta_k + G_k \}(z - b_j)^{k - s_j} + \cdots,
\]

where \( f(\ell) = (\ell + s_j)(\ell - s_j - 1) \) and \( G_k = B_j\beta_{k-1} + (C_{j,0} - z)\beta_{k-2} + C_{j,1}\beta_{k-3} + \cdots + C_{j,k-2}\beta_0 \).

Now \( f(-s_j) = 0 \) and \( \beta_1, \ldots, \beta_{s_j} \) may be determined recursively so that the coefficients of \( (x - b_j)^{1-s_j}, \ldots, (x - b_j)^{s_j} \) vanish. But since \( f(s_j + 1) = 0 \), the coefficient of \( (x - b_j)^{s_j+1} \) is just \( G_{2s_j+1} \), which therefore must vanish for all \( z \in \mathbb{C} \). On the
other hand, if \( B_j \neq 0 \), one can show by induction that \( G_{2s_j+1} \) is a polynomial in \( z \) of degree \( s_j \). This contradiction completes the proof.

If \( q \) is a Picard potential then, by Picard’s theorem (Theorem 2.7), the equation \( y''(x) + q(x)y(x) = zy(x) \) has at least one solution which is elliptic of the second kind. Using Theorem 2.4 and the special structure of the Picard potential (3.2.1) this solution may be represented as

\[
(3.2.2) \quad \psi_a(x) = \frac{\prod_{j=1}^{s} \sigma(x - a_j(z))}{\prod_{j=1}^{b} \sigma(x - b_j)^{s_j}} \exp(\lambda_a(x)z),
\]

where \( s = \sum_{j=1}^{m} s_j \) and \( a(z) = (a_1(z), \ldots, a_s(z)) \). At \( b_j \) the function \( \psi_a \) has a pole of order \( s_j \) or a zero of order \( s_j + 1 \). For later notational purposes we allow for \( s_j = 0 \), in which case \( \psi_a \) has either no pole and no zero or a simple zero. For subsequent use we define

\[
M_1 = \{ j \in \{1, \ldots, m\} \mid \text{ord}_{b_j}(\psi_a) = -s_j \},
\]

\[
M_2 = \{ j \in \{1, \ldots, m\} \mid \text{ord}_{b_j}(\psi_a) = s_j + 1 \}.
\]

The function \( \psi_a \) is a solution of \( y''(x) + q(x)y(x) = zy(x) \) if and only if

\[
(3.2.3) \quad \lambda_a + \sum_{j=1}^{s} \zeta(b_r - a_j) - \sum_{j=1, j \neq r}^{m} s_j \zeta(b_r - b_j) = 0,
\]

\[
(3.2.4) \quad z = C - \sum_{j=1}^{s} (1 - 2s_r) \varphi(b_r - a_j) - \sum_{j=1, j \neq r}^{m} s_j (s_j + 2s_r) \varphi(b_r - b_j),
\]

where \( r \) is chosen such that \( s_r \neq 0 \) and \( \text{ord}_{b_r}(\psi_a) = -s_r \) which is always possible. In the case of a Lamé potential these conditions are recorded, for instance, by Burkhardt [33]. The subscript \( a \) in \( \lambda_a \) expresses the dependence of \( \lambda_a \) on \( a = (a_1, \ldots, a_s) \), which in turn depends on \( z \).

In [14] we have developed a method for even Picard potentials (i.e., potentials \( q \) satisfying \( q(x_0 + x) = q(x_0 - x) \) for some \( x_0 \in \mathbb{C} \)) to determine all points where two regular Floquet solutions fail to exist. (For simplicity we will assume \( x_0 = 0 \) from now on.) First note that an even Picard potential is of the form

\[
(3.2.5) \quad q(x) = C - \sum_{k=1}^{4} s_k (s_k + 1) \varphi(x - \omega_k) - \sum_{k=1}^{\tilde{m}} r_k (r_k + 1) [\varphi(x - b_k) + \varphi(x + b_k)],
\]

where the \( b_k \) are pairwise distinct and different from half-periods, the \( s_k \) are non-negative and the \( r_k \) are positive. If \( \tilde{m} = 0 \), that is, if

\[
q(x) = C - \sum_{k=1}^{4} s_k (s_k + 1) \varphi(x - \omega_k),
\]

then \( q \) is called a Treibich-Verdier potential following the work of Verdier [230] and Treibich and Verdier [218, 227]. If in addition, only one of the numbers \( s_k \) is different from zero, then \( q \) is a Lamé potential. In order to make use of previous results we will adopt the following notation: \( m = 2\tilde{m} + 4 \), \( b_k + \omega_k = -b_k \), \( s_k + \tilde{m} = \tilde{m} \) for \( k = 1, \ldots, \tilde{m} \), and \( b_{k+2\tilde{m}} = \omega_k \), \( s_{k+2\tilde{m}} = s_k \) for \( k = 1, \ldots, 4 \). Let \( \delta \) be the number of \( a_j \)’s which do not appear in \( \{b_1, \ldots, b_m\} \).
If \( q \) is an even Picard potential of the type (3.2.5), and if \( \psi_a \) given by (3.2.2) is a solution of the differential equation \( y''(x) + q(x)y(x) = zy(x) \), then so is the function \( \psi_{-a} \), which is obtained by replacing every \( a_j \) with \(-a_j \) in (3.2.2) and (3.2.3), since \( \psi_{-a}(x) = (-1)^{s_1 + s_2 + s_3} \psi_a(-x) \).

Next we compute the Wronskian of the two solutions \( \psi_a \) and \( \psi_{-a} \). One obtains an expression which involves \( x \), but since the Wronskian does not depend on \( x \), one may evaluate it anywhere. Choosing \( x = a_\ell \) for any \( a_\ell \) which does not appear in \( \{ b_1, ..., b_m \} \), one finds that the Wronskian is a nonzero multiple of

\[
\frac{\sigma(2a_\ell)}{\sigma(a_\ell - b_r)\sigma(a_\ell + b_r)} \prod_{j=1, j \neq \ell}^{s} \sigma(a_\ell - a_j)\sigma(a_\ell + a_j).
\]

Since \( a_\ell \) is different from all the \( b_j \) and, in particular, different from the half-periods, one infers \( \sigma(2a_\ell) \neq 0 \). Also \( a_\ell \neq a_k \) if \( k \neq \ell \). Therefore we find that the Wronskian is zero if and only if \( \sigma(a_\ell + a_j) = 0 \) for some \( j \in \{ 1, ..., \ell - 1, \ell + 1, ..., s \} \) and hence \( a_j = -a_\ell \) (mod \( \Delta \)). In particular, we find that the number \( \hat{s} \) is even and define \( d = \hat{s}/2 \).

Choosing now \( x = b_\ell \) for any \( \ell \in M_2 \) the Wronskian can be written as a nonzero multiple of

\[
\frac{\sigma(2b_\ell)2^{2s_\ell + 1}}{\sigma(b_\ell - b_r)\sigma(b_\ell + b_r)} \prod_{j=1, j \neq \ell}^{m} \sigma(b_\ell - b_j)2^{2s_\ell + 1},
\]

which is zero if and only if \( b_\ell \) is a half-period or if there is a \( j \in M_2 \) such that \( b_j = -b_\ell \) (mod \( \Delta \)).

In summary we have found the following: if \( \psi_a \) and \( \psi_{-a} \) are linearly dependent solutions of \( \psi'' + q\psi = z\psi \), then some of the numbers \( a_1, ..., a_s \) may be half-periods while all others appear in pairs \( (a_j, a_{\ell j}) \) with \( a_{\ell j} = -a_j \). Moreover, if \( a_j \) is equal to a half-period \( \omega_k \), which is a pole of \( q \) of the form \(-r_k/(z - \omega_k)^2 \), then there are exactly 2\( r_k + 1 \) of the \( a_j \) which are equal to this half-period. If \( a_j \) equals a pole \( b_\ell \) of the form \(-s_\ell/(z - b_\ell)^2 \), where \( b_\ell \) is not a half-period, then there are exactly 2\( s_\ell + 1 \) of the \( a_m \) which are equal to this pole and exactly 2\( s_\ell + 1 \) other \( a_m \)’s which are equal to the pole \(-b_\ell \).

This information is now being used to rewrite the solution \( \psi_a \) of \( \psi'' + q\psi = z\psi \) for those values of the spectral parameter \( z \) where \( W(\psi_a, \psi_{-a}) = 0 \) as a product of two functions. The first function is fixed, depending only on the poles of the potential \( q \), on the half-periods, and the exponents associated with these. The second function is a polynomial in \( \psi(x) \), whose coefficients depend on those \( a_j \) which are neither half-periods nor poles of \( q \) and which are as yet undetermined. According to the above argument there must be an even number, 2\( d \), of those, and half of them are just negatives of the other half.

Hence we define \( t_\ell = \mathrm{ord}_{b_\ell}(\psi_a) \) for \( \ell = 1, ..., 2\hat{m} + 4 \) and obtain \( \psi_a(x) = f(x)Q(\varphi(x)) \), where

\[
f(x) = e^{\lambda_\varphi x} \left( \prod_{k=1}^{4} \sigma(x - \omega_k)^{t_k + 2\hat{m}} \right) \left( \prod_{k=1}^{2\hat{m}} \sigma(x - b_\ell)^{t_\ell} \right)^{2d},
\]

\[
Q(\varphi(x)) = \prod_{j=1}^{d} (\varphi(x) - \varphi(a_j)) = \sum_{j=0}^{d} c_j \varphi(x)^j.
\]
Here we used the fact that $\sigma(z - aj)\sigma(z + aj) = -\sigma(z)^2\sigma(aj)^2(\varphi(z) - \varphi(aj))$. Moreover, we dropped the non-zero constant factor $(-1)^d \prod_{j=1}^{d} \sigma(aj)^2$. This yields

$$\psi_a''(x) + q(x)\psi_a(x) = f(x) \left\{ \sum_{k=0}^{d} \sum_{j=0}^{d} S_{k+1,j+1} c_j \psi(x)^k + \sum_{k=1}^{m} \sum_{j=0}^{d} T_{k,j+1} c_j \psi(x)^{k-1} \right\} \prod_{j=1}^{m} (\varphi(z) - \varphi(b_j))$$

for suitable constants $S_{j,k}$ and $T_{j,k}$ depending on $q$ and the numbers $t_k$. Let $S$ and $T$ be the matrices with entries $S_{j,k}$ and $T_{j,k}$.

Hence we obtain a solution $\psi_a$ of the equation $y''(x) + q(x)y(x) = zy(x)$ satisfying $W(\psi_a, \psi_{-a}) = 0$ if and only if

$$S\gamma = z\gamma \text{ and } T\gamma = 0,$$

where $\gamma = (c_0, ..., c_d)^T$.

For any given even Picard potential there are several (but finitely many) choices to distribute some (or all) of the parameters $a_1, ..., a_s$ among the half-periods and/or poles of $q$. Accordingly, there are several (but finitely many) of the above described constraint eigenvalue problems to solve in order to find all the values of the spectral parameter $z \in \mathbb{C}$ where $W(\psi_a, \psi_{-a}) = 0$. In each case there are only finitely many eigenvalues of the associated matrix $S$, some (or possibly all) of which may be in contradiction to the constraint $T\gamma = 0$. Thus, we proved the following result.

**Theorem 3.8.** Let $q$ be an even Picard potential. Then, for every complex number $z$, the differential equation $y''(x) + q(x)y(x) = zy(x)$ has a solution $\psi_a$ of the form (3.2.2), that is, a solution which is elliptic of the second kind. Similarly, the function $\psi_{-a}$ is a solution of the same equation (for the same value of $z$) and also elliptic of the second kind. With respect to any period of $q$, the functions $\psi_a$ and $\psi_{-a}$ are regular Floquet solutions. For all but a finite number of values of $z \in \mathbb{C}$ these two solutions are linearly independent and therefore $p(z) - 2d_t(z) = 0$. Hence $\text{def}(q) < \infty$, that is, $q$ is algebro-geometric.

In the case of a Treibich-Verdier potential

$$q(x) = -\sum_{j=1}^{4} s_j(s_j + 1)\varphi(x - \omega_j), \quad s_j \in \mathbb{N} \cup \{0\}, \quad j = 1, \ldots, 4,$$

the matrix $T$ is absent and therefore one only has to find the eigenvalues of $S$ for all possible choices of the numbers $t_k$. This was performed in [14] for Lamé potentials and in [93] for Treibich-Verdier potentials. The arithmetic genus of the curve $P^2 = R(L)$ associated with $q$ is given in the following table, where $s = s_1 + s_2 + s_3 + s_4$ (and, without loss of generality, $s_1 \geq s_2 \geq s_3 \geq s_4 \geq 0$).

| $s$       | # of finite branch points | genus   |
|-----------|---------------------------|---------|
| even $s_2 + s_3 \leq s_1 + s_4$ | $2s_1 + 1$ | $s_1$ |
| even $s_2 + s_3 \geq s_1 + s_4$ | $s_1 + s_2 + s_3 - s_4 + 1$ | $\frac{1}{2} - s_4$ |
| odd $s_2 + s_3 + s_4 < s_1$ | $2s_1 + 1$ | $s_1$ |
| odd $s_2 + s_3 + s_4 \geq s_1$ | $s_1 + s_2 + s_3 + s_4 + 2$ | $\frac{5}{2} - s_4$ |

Shortly after this table appeared in our paper [13], Armando Treibich kindly sent us the following formula for the arithmetic genus of the Treibich-Verdier curve. It
is given by
\[ \max\{2s_1, s + 1 - (1 + (-1)^s)(s_4 + (1/2))\}/2 \]
in agreement with our column on the right.

**Remark 3.9.** We briefly return to the Picard potential \( q(x) \) in (3.2.1). Assuming first the special case where \( s_j = 1, j = 1, \ldots, m \), that is,
\[ q(x) = -2 \sum_{j=1}^{m} \wp(x - b_j), \quad b_k \neq b_\ell \text{ (mod } \Delta) \text{ for } k \neq \ell, \]
Theorem 2.5 yields Krichever’s ansatz (2.8.3) in [136], that is,
\[ \psi_\alpha(z)(x) = \exp(\lambda_\alpha(z)x) \prod_{j=1}^{m} \frac{\sigma(x - a_j(z))}{\sigma(x - b_j)} = e^{\kappa(z)x} \sum_{j=1}^{m} A_j(z) \Phi(x - b_j, \rho(z)), \]
and where \( \kappa, A_j \) and \( \rho_j \) are suitably chosen.

In the general case, where
\[ q(x) = -\sum_{j=1}^{m} s_j(s_j + 1) \wp(x - b_j), \quad b_k \neq b_\ell \text{ (mod } \Delta) \text{ for } k \neq \ell, \]
one can use an extension of Theorem 2.5 (using, e.g., l’Hospital’s rule as several \( b_j \)'s are confluenting into a \( b_{j_0} \) and converting \( b \)-derivatives into an \( x \)-derivative) to arrive at the corresponding analog
\[ \psi_\alpha(z)(x) = \exp(\lambda_\alpha(z)x) \prod_{j=1}^{m} \frac{\sigma(x - a_j(z))}{\sigma(x - b_j)} = e^{\kappa(z)x} \sum_{j=1}^{m} s_j \sum_{k=0}^{s_j - 1} A_{j,k}(z) \Phi^{(k)}(x - b_j, \rho(z)), \quad s = \sum_{j=1}^{m} s_j \]
of (3.2.7). The extended ansatz (3.2.9) was recently used by Enol’skii and Eilbeck [59], [64], Enol’skii and Kostov [65], [130], in the context of Treibich-Verdier potentials. The special case \( m = 1, s_1 = 2 \) can be found in Hermite [113], p. 374–377 (see also Forsyth [76], p. 475–476), the general case (3.2.9) is discussed in Krause’s monograph [132], Vol. 1, p. 292–296, Vol. 2, p. 183, 259–264.

We emphasize that Picard’s theorem, Theorem 2.7, yields the \( \sigma \)-function representations of \( \psi_\alpha(z)(x) \) in (3.2.7) and (3.2.9) and hence Krichever’s ansatz (2.8.3) and its extension (3.2.9). In particular, an alternative characterization of all stationary elliptic KdV and AKNS potentials to the one provided in Theorem 3.12 and Theorem 4.8, respectively, can be based on these observations (cf. [101]). In the remainder of this review we do not pursue this avenue but stress Floquet-theoretic and Green’s function methods instead.

### 3.3. Necessary Conditions for a KdV Potential to be Algebro-Geometric.

**Theorem 3.10.** Suppose \( q \) is an algebro-geometric KdV potential. Then \((Ly)(x) = y'''(x) + q(x)y(x) = zy(x)\) has a meromorphic fundamental system of solutions with respect to \( x \) for all values of the spectral parameter \( z \in \mathbb{C} \).
Sketch of proof. First one shows that every pole of $q$ is of second order. This follows since only in this case there is a balance between the growth of the order of the poles of $f_j''$ and $4qf_j'' + 2q_j' f_j$, considering the recursion relation (2.3.3). Such a balance is necessary since eventually all $f_j$ vanish. Next one shows by a similar argument that the leading term of the Laurent expansion of $q$ about a pole $x_0$ must be $-s(s+1)/(x-x_0)^2$ for some integer $s$. This guarantees that the exponents of the singularity $x_0$ of the differential equation $y''(x) + q(x)y(x) = zy(x)$ are integers. Finally one has to show the absence of logarithmic terms in the solutions of $y''(x) + q(x)y(x) = zy(x)$ which follows from a careful analysis using the Frobenius method. More details are provided in [233].

We note that $\tau$-function results in Segal and Wilson [194] imply Theorem 3.10. (In the case of nonsingular curves $K_n$ this simply follows from the standard theta-function representation of the Baker-Akhiezer function.) While Segal and Wilson rely on loop group techniques (and study the Gelfand-Dickey hierarchy), the above proof represents a completely elementary alternative for the KdV hierarchy.

3.4. Sufficient Conditions for a KdV Potential to be Algebrao-Geometric – Characterizations of Elliptic, Simply Periodic, and Rational Stationary KdV Solutions.

**Theorem 3.11.** $q$ is an algebrao-geometric potential if, for each $z \in \mathbb{C}$, $y''(x) + q(x)y(x) = zy(x)$ has a meromorphic fundamental system of solutions with respect to $x$ and if one of the following three conditions is satisfied.

(i) $q$ is rational and bounded near infinity.

(ii) $q$ is simply periodic with period $\Omega$ and there exists a positive number $R$ such that $q$ is bounded in $\{x \in \mathbb{C} \mid |\text{Im}(x/\Omega)| \geq R\}$.

(iii) $q$ is elliptic.

**Sketch of proof.** We will first sketch the proof in [85] which treats the case where $q$ is elliptic and prominently uses the double periodicity of $q$.

Let $q$ be an elliptic function with fundamental periods $2\omega_1$ and $2\omega_3$. Assume, without loss of generality, that $\text{Im}(\omega_3/\omega_1) > 0$. Introduce $t_j = \omega_j/\omega_1$ and define

$q_j(x) = t_j^2 q(t_j x).

The transformation $\xi = t_j x$, $\psi(\xi) = w(x)$ transforms $\psi'' + q \psi = z \psi$ into $w'' + q_jw = t_j^2 zw$.

Note that $q_j$ has the period $2|\omega_j|$. With the aid of Rouché’s theorem, one may determine the asymptotic distribution of the (anti-)periodic eigenvalues of $d^2/dx^2 + q_j$ and verify that they all stay close to the real axis while their real parts tend to $-\infty$.

(Never mind that $q_j$ may have inverse square singularities, see [233].) Equivalently, all $2\omega_j$-(anti-)periodic eigenvalues of $d^2/dx^2 + q$ lie in the half-strip $\Sigma_j$ given by

$$\Sigma_j = \{ z \in \mathbb{C} \mid |\text{Im}(t_j^2 z)| \leq C, \text{Re}(t_j^2 z) \leq C \}, \quad j = 1, 3$$

for some constant $C > 0$. Note that the angle between the strips $\Sigma_1$ and $\Sigma_3$ is positive and less than $2\pi$ and therefore, they intersect only in a finite part of the complex plane. Hence, for every point $z$ outside a sufficiently large compact disk, at least one of the translation operators $T_1$ or $T_3$ has distinct eigenvalues. Assuming, according to our hypothesis, that all solutions of $y''(x) + q(x)y(x) = zy(x)$ are meromorphic, we may invoke Picard’s theorem to conclude that there are two linearly independent solutions which are elliptic of the second kind for
any $z$ outside that disk. These solutions are Floquet solutions with respect to any period of $q$ and since the disk is compact, we have shown that there are at most finitely many points which lack two linearly independent Floquet solutions. Applying Theorem 3.3 then proves that $q$ is algebro-geometric.

Next we sketch the proof for the case when $q$ is rational. That proof can easily be extended for the simply periodic and elliptic potentials. In the latter case one uses the algebraic properties of elliptic functions rather than their double periodicity. For more details see [234].

Suppose that $q$ is rational and bounded at infinity. Let $z_0 = \lim_{x \to \infty} q(x)$. From Halphen's theorem (Theorem 2.8) we obtain for $z \neq z_0$ the existence of linearly independent solutions

$$y_{\pm}(z, x) = R_{\pm}(z, x) \exp(\pm \sqrt{z - z_0} x),$$

where $R_{\pm}(z, \cdot)$ are rational functions. The poles of $R_{\pm}(z, \cdot)$ are determined as the singular points of the differential equation and their orders as the exponents of the corresponding singularities. Next define the function $g(z, x) = y_+(z, x)y_-(z, x)$. Then there exists a polynomial $v$ in $x$ such that

$$v(x)^2 g(z, x) = \sum_{j=0}^{d} c_j(z) x^j.$$

Next note that the functions $v^2 g(z, x)$, $v^3 g'(z, x)$, $v^4 g''(z, x)$, and $v^5 g'''(z, x)$ are polynomials in $x$ whose coefficients are homogeneous polynomials of degree one in $c_0, ..., c_d$. Also $v^2 q$ and $v^3 q'$ are polynomials. Hence $v^5 (g''' + 4(q - z)g' + 2q'g)$ is also a polynomial in $x$, whose coefficients are homogeneous polynomials of degree one in $c_0, ..., c_d$. The coefficients of the polynomials $c_\ell$ in this last expression are polynomials in $z$ of degree at most one, that is,

$$(3.4.1)\quad v^5 (g''' + 4(q - z)g' + 2q'g) = \sum_{j=0}^{N} \sum_{\ell=0}^{d} (\alpha_{j, \ell} + \beta_{j, \ell} z) c_\ell x^j$$

for suitable numbers $N$, $\alpha_{j, \ell}$, and $\beta_{j, \ell}$ which only depend on $q$. From Appell's equation (2.9.2) it follows upon differentiation that the expression (3.4.1) vanishes identically. This gives rise to a homogeneous system of $N + 1$ linear equations which has a nontrivial solution. Solving the system shows now that the coefficients $c_\ell$ are rational functions of $z$. Therefore

$$g(z, x) = \frac{F(z, x)}{\gamma(z)},$$

where $F(\cdot, x)$ and $\gamma$ are polynomials and $F(z, \cdot)$ is a rational function. Hence $q$ is algebro-geometric by item (3) following Definition 3.1.

Combining the results of this subsection and its preceding one yields the following explicit characterization of all elliptic algebro-geometric KdV potentials (a problem posed, for instance, by Novikov, Manakov, Pitaevskii, and Zakharov [170], p. 152), originally proven in [98] (see also [97], [100]).

**Theorem 3.12.** Let $q$ be an elliptic function. Then $q$ is an elliptic algebro-geometric KdV potential if and only if it is a Picard potential.

Similarly, these results characterize the stationary rational KdV potentials vanishing at infinity studied, for instance, by Adler and Moser [4], Airault, McKean,
and Moser [3], Calogero [33, 36], Chudnovsky and Chudnovsky [39, 41], Grinevich
[105], Krichever [135], Matveev [155], and Sokolov [207].

4. Algebra-Geometric and Especially, Elliptic AKNS Potentials.

Most of the results in this section are taken from [99]. Since we take here a slightly different point of departure, the notation differs a bit from that in [99].

4.1. The AKNS Hierarchy. Let \( L = Jd/dx + Q(x) \), where

(4.1.1)

\[
J = \begin{pmatrix} i & 0 \\ 0 & -i \end{pmatrix} \quad \text{and} \quad Q(x) = \begin{pmatrix} 0 & -i(q(x)) \\ ip(x) & 0 \end{pmatrix}.
\]

Note that \( J^2 = -I \), the \( 2 \times 2 \) identity matrix, and that \( JQ + QJ = 0 \). As mentioned in Section 2.3 one may consider the Lax pair \((P_{n+1}, L)\), where \( P_{n+1} \) is a \( 2 \times 2 \)-matrix-valued differential expression of order \( n + 1 \) such that \([P_{n+1}, L]\) is an operator of multiplication. Writing

\[
P_{n+1} = \sum_{\ell=0}^{n+1} C_{n+1-\ell}(x)L^\ell,
\]

and utilizing that the commutator \([C_j, L]\) may be written as \( D_j(x)L + E_j(x) \), the condition that \([P_{n+1}, L]\) is an operator of multiplication yields \( D_{j+1}(x) = -E_j(x) \)
for \( j = 0, \ldots, n \), \( D_0 = 0 \), and \([P_{n+1}, L] = E_{n+1}(x)\). Note that \( C_j(x) \) can be expressed as

\[
C_j(x) = k_j(x)I + \nu_j(x)J + W_j(x), \quad j = 0, \ldots, n + 1,
\]

where the \( k_j(x) \) and \( \nu_j(x) \) are scalar-valued and the matrices \( W_j(x) \) have vanishing diagonal elements. Then \((\ell' = d/dx) D_j = 2W_j \) and \( E_j = -W_jQ - QW_j + \nu_j'x - k_j'J + 2\nu_jJQ - JW_j' \).

Since \( W_jQ + QW_j \) is a multiple of the identity matrix \( I \), \( D_{j+1} + E_j = 0 \) shows that \( k_j' = 0 \) for all \( j = 0, \ldots, n + 1 \), and that the following recursion relation holds,

(4.1.2)

\[
W_0 = 0, \quad \nu_j' = W_jQ + QW_j, \quad W_{j+1} = \frac{1}{2} J(W_j - 2\nu_jQ), \quad j = 0, \ldots, n + 1.
\]

One concludes

\[
[P_{n+1}, L] = 2v_{n+1}JQ - JW_{n+1}'.
\]

Next, let

\[
K_{n+1}(z) = \sum_{j=0}^{n+1} k_{n+1-j}z^j,
\]

\[
V_{n+1}(z, x) = \sum_{j=0}^{n+1} v_{n+1-j}(x)z^j,
\]

\[
W_{n+1}(z, x) = \sum_{j=1}^{n+1} W_{n+1-j}(x)z^j,
\]

so that the recursion relation (4.1.2) becomes

\[
V_{n+1}(z, x)I = W_{n+1}(z, x)Q(x) + Q(x)W_{n+1}(z, x),
\]

\[
W_{n+1}'(z, x) = 2V_{n+1}(z, x)Q(x) - 2zJW_{n+1}(z, x) + J[P_{n+1}, L].
\]
Hence,
\[(4.1.3)\quad [P_{n+1}, L] = 2z\mathcal{W}_{n+1}(z, x) + 2V_{n+1}(z, x)JQ(x) - JW'_{n+1}(z, x).\]

If \(P_{n+1}\) is of order \(n + 1\) we define the \(n^{th}\) order AKNS equations by
\[\text{AKNS}_n(Q) = Q_t - [P_{n+1}, L] = 0.\]
The first few of these equations are
\[Q_t = -v_0Q' + 2c_1 JQ,\]
\[Q_t = -\frac{v_0}{2}J(Q'' - 2Q^3) - c_1Q' + 2c_2 JQ,\]
\[Q_t = \frac{v_0}{4}(Q'' - 6Q^2Q') - c_1J(Q'' - 2Q^3) - c_2Q' + 2c_3 JQ,\]
\[\text{etc.},\]
where \(v_0\) and \(c_1, c_2, \ldots\) are arbitrary integration constants. Upon rescaling the \(t\) variable one may choose \(v_0 = 1\). In terms of the AKNS pair \((p, q)\) the homogeneous versions (that is, \(c_1 = c_2 = \cdots = 0\)) of these equations read
\[
\begin{align*}
\begin{pmatrix} p_t \\ q_t \end{pmatrix} &= -v_0 \begin{pmatrix} p_x \\ q_x \end{pmatrix}, \\
\begin{pmatrix} p_t \\ q_t \end{pmatrix} &= \frac{iv_0}{2} \begin{pmatrix} p_{xx} - 2p^2q \\ -q_{xx} + 2pq^2 \end{pmatrix}, \\
\begin{pmatrix} p_t \\ q_t \end{pmatrix} &= \frac{iv_0}{4} \begin{pmatrix} p_{xxx} - 6pqpx \\ -q_{xxx} + 6ppqx \end{pmatrix},
\end{align*}
\]
\[\text{etc.}\]

We also mention an interesting scale invariance of the AKNS equations. Suppose \(Q\) satisfies one of the AKNS equations, that is, \(\text{AKNS}_n(Q) = 0\). Suppose \(a \neq 0\) and
\[A = \begin{pmatrix} a & 0 \\ 0 & a^{-1} \end{pmatrix}.\]
Then \(\tilde{Q} = AQA\) also satisfies \(\text{AKNS}_n(\tilde{Q}) = 0\). We omit the straightforward proof which can be found, for instance, in [88].

In the particular case of the nonlinear Schrödinger (NS) hierarchy, where \(p(x, t) = \pm q(x, t)\) the matrix \(A\) is unimodular, that is, \(|a| = 1\).

Note that the KdV hierarchy as well as the modified Korteweg-de Vries (mKdV) hierarchy are contained in the AKNS hierarchy. In fact, setting all integration constants \(c_{2\ell+1}\) equal to zero, the \(n^{th}\) KdV equation is obtained from the \((2n)^{th}\) AKNS system by the constraint \(p(x, t) = 1\), while the \(n^{th}\) mKdV equation is obtained from the \((2n)^{th}\) AKNS system by the constraint \(p(x, t) = \pm q(x, t)\).

Just as in the KdV case one makes the following observation: Suppose a polynomial \(V_{n+1}\) whose coefficients are scalar functions and a polynomial \(\mathcal{W}_{n+1}\) whose coefficients are \(2 \times 2\) matrix-valued functions with zero diagonal to be given. Furthermore, assume that \(V'_{n+1}(z, x)I = \mathcal{W}_{n+1}(z, x)Q(x) + Q(x)\mathcal{W}'_{n+1}(z, x)\) and that \(2z\mathcal{W}_{n+1}(z, x) + 2V_{n+1}(z, x)JQ(x) - JW'_{n+1}(z, x)\) is independent of \(z\). Then the coefficients of \(V_{n+1}(z, x)\) and \(\mathcal{W}_{n+1}(z, x)\) define a differential expression \(P_{n+1}\) which satisfies \((4.1.3)\).

Next, suppose that \([P_{n+1}, L] = 0\) and, without loss of generality, \(v_0 \neq 0\). Then,
\[\begin{align*}
(4.1.4)\quad (P_{n+1} - K_{n+1}(L))^2 &= (JV_{n+1}(L, x) + \mathcal{W}_{n+1}(L, x))^2
\end{align*}\]
\[ \begin{align*}
&= W_{n+1}(L,x)^2 - V_{n+1}(L,x)^2 I \\
&= \sum_{m=0}^{2n+2} a_m(x)L^m,
\end{align*} \]

where

\[ a_m(x)I = \sum_{\ell+k=m} (W_{n+1-\ell}(x)W_{n+1-k}(x) - v_{n+1-\ell}(x)v_{n+1-k}(x))I \]

is a multiple of the identity matrix. Moreover, differentiating \((JV_{n+1}(z,x) + W_{n+1}(z,x))^2\) with respect to \(x\) yields

\[-2V_{n+1}(z,x)V'_{n+1}(z,x)I + W_{n+1}(z,x)W'_{n+1}(z,x)\]

using

\[ V'_{n+1}(z,x)I = W_{n+1}(z,x)Q(x) + Q(x)W_{n+1}(z,x), \]

\[ W'_{n+1}(z,x) = 2V_{n+1}(z,x)Q - 2zJW_{n+1}(z,x), \]

and

\[ 0 = JW_{n+1}(z,x) + W_{n+1}(z,x)J. \]

Hence the coefficients \(a_m(x)\) in \([1.1.4]\) may be interpreted as constant scalars. Since \(a_{2n+2} = -v_0^2\) one infers \((P_{n+1} - R_{n+1}(L))^2 + R_{2n+2}(L) = 0\), where \(R_{2n+2}\) is a polynomial of degree \(2n+2\) with complex coefficients. Hence, if \([P_{n+1}, L] = 0\), then the pair \((P_{n+1}, L)\) is associated with a hyperelliptic curve of (arithmetic) genus \(n\).

Next assume that \(F_n : \mathbb{C}^2 \to \mathbb{C}_\infty\) is a polynomial of degree \(n\) in its first variable with scalar meromorphic coefficients. Denote the leading coefficient by \(-iq(x)\) and let \(p(x)\) be another nonzero meromorphic function. Defining

\[ V_{n+1}(z,x) = \frac{-1}{2q(x)}(F_n'(z,x) + 2izF_n(z,x)) \]

and

\[ W_{n+1}(z,x) = \frac{i}{q(x)} \begin{pmatrix} 0 & q(x)F_n(z,x) \\ q(x)F_n(z,x) & 0 \end{pmatrix}, \]

this implies \(V'_{n+1}I = W_{n+1}Q + QW_{n+1}\) with \(Q\) given as in \([1.1.1]\). Moreover, \(V_{n+1}^2 - W_{n+1}^2 = R_{2n+2}(z,x)I\), where the scalar \(R_{2n+2}(z,x)\) is given by

\[ R_{2n+2}(z,x) = \frac{1}{4q(x)^2}(F_n'(z,x)^2 - 2F_n(z,x)F''_n(z,x) + 4(p(x)q(x) - z^2)F_n(z,x)^2) \]

\[ + \frac{q'(x)}{4q(x)}(2F_n(z,x)F'_n(z,x) + 4izF_n(z,x)^2). \]

If this is constant then differentiation with respect to \(x\) yields

\[ W_{n+1}(z,x)(W_{n+1}(z,x) - 2V_{n+1}(z,x)Q(x)) \]

\[ + (W_{n+1}(z,x) - 2V_{n+1}(z,x)Q(x))W_{n+1}(z,x) = 0. \]

Since \(W'_{n+1}(z,x) - 2V_{n+1}(z,x)Q(x)\) has zero diagonal elements, this may be considered, for each fixed \(x\), a linear homogeneous equation for the off-diagonal elements of \(W'_{n+1}(z,x) - 2V_{n+1}(z,x)Q(x)\). This equation has a one-dimensional space of solutions, in fact, \(W'_{n+1}(z,x) - 2V_{n+1}(z,x)Q(x) = r(x)JW_{n+1}(z,x)\). Comparing the leading coefficients yields \(r(x) = -2z\) and hence we have shown that \(W'_{n+1}(z,x) = -2zJW_{n+1}(z,x) + 2V_{n+1}(z,x)Q(x)\). Summarizing, the following theorem holds.
Theorem 4.1. If \( F_n : C^2 \rightarrow C_\infty \) is a polynomial of degree \( n \) in the first variable and meromorphic in the second, and if the expression \( R_{2n+2}(z,x) \) in (4.1.5) is independent of \( x \), then there exists a \( 2 \times 2 \) matrix-valued differential expression \( P_{n+1} \) of order \( n+1 \) with leading coefficient \( J^{n+2} \) such that \( [P_{n+1}, L] = 0 \).

We are now ready to define the term algebro-geometric in the context of the AKNS hierarchy.

Definition 4.2. Suppose \( Q \) is a \( 2 \times 2 \) matrix-valued meromorphic function and let \( L \) be the differential expression \( L = J d/dx + Q \). Then \( Q \) is called an algebro-geometric AKNS potential (or simply algebro-geometric) if \( Q \) is a solution of some equation of the stationary AKNS hierarchy.

Equivalently, \( Q \) is algebro-geometric if any one of the following two conditions is satisfied.

1. There exists a differential expression \( P_{n+1} \) of order \( n+1 \) with leading coefficient \( J^{n+2} \) such that \( [P_{n+1}, L] = 0 \).
2. There exists a function \( F_n : C^2 \rightarrow C_\infty \) which is a polynomial in the first variable, meromorphic in the second, such that the expression \( R_{2n+2}(z,x) \) in (4.1.5) does not depend on \( x \).

4.2. Periodic AKNS Potentials. Throughout this section we will assume that \( p, q \in C^2(\mathbb{R}) \) are complex-valued periodic functions of period \( \Omega > 0 \). Matrix-valued solutions \( Y \) of the initial value problem \((LY)(x) = zY(x) \) and \( Y(x_0) = Y_0 \) are denoted by \( \Phi(z, \cdot, x_0, Y_0) \), which is also the unique solution of the integral equation

\[
(4.2.1) \quad Y(x) = \exp(z(x - x_0)J)Y_0 + \int_{x_0}^{x} e^{z(x-x')}JQ(x')Y(x') \, dx'.
\]

In contrast to the Sturm-Liouville case, the Volterra integral equation (4.2.1) is not suitable to determine the asymptotic behavior of solutions as \( z \) tends to infinity. To circumvent this difficulty one can follow Marchenko’s approach in [153], Section 1.4, to obtain the asymptotic expansion

\[
\Phi(z, x_0 + \Omega, x_0, I) = \left( \begin{array}{cc} e^{-iz\Omega} & 0 \\ 0 & e^{iz\Omega} \end{array} \right) + \frac{1}{2i\Omega} \left( \begin{array}{cc} \beta e^{-iz\Omega} & 2q(x_0) \sin(z\Omega) \\ 2p(x_0) \sin(z\Omega) & -\beta e^{iz\Omega} \end{array} \right) + O(e^{\text{Im}(z)\Omega z^{-2}}),
\]

(4.2.2)

where

\[
\beta = \int_{x_0}^{x_0 + \Omega} p(t)q(t) \, dt,
\]

provided \( p, q \in C^2(\mathbb{R}) \). From this result we infer that the entries of \( \Phi(\cdot, x_0 + \Omega, x_0, I) \), which are entire functions, have order of growth equal to one whenever \( q(x_0) \) and \( p(x_0) \) are nonzero.

As in the scalar case we introduce \( T_\Omega(z) \), the restriction of the translation operator \( Y \mapsto Y(\cdot + \Omega) \) to the two-dimensional space of solutions \( Y(z) \) of \((LY)(x) = zY(x) \). The Floquet multipliers are then determined as solutions of

\[
\rho^2 - \rho \text{tr}(T_\Omega(z)) + 1 = 0.
\]

They are degenerate if and only if \( \rho(z)^2 = 1 \). The points \( z \) where this happens are the (anti-)periodic eigenvalues which we denote by \( E_j, j \in \mathbb{Z} \). Their algebraic multiplicities are given by \( \text{ord}_{E_j}(T_\Omega^2 - 4) \) and are denoted by \( p(E_j) \). The asymptotic
behavior of the (anti-)periodic eigenvalues may be determined from \((1.2.3)\). It is given by

\[ E_{2j}, E_{2j-1} = \frac{j\pi}{\Omega} + O(|j|^{-1}), \]

if the eigenvalues are labeled in accordance with their algebraic multiplicity.

The conditional stability set \(S(L)\) of \(L\) is again

\[ S(L) = \{ z \in \mathbb{C} | -2 \leq \text{tr}(T_{11}(z)) \leq 2 \}. \]

The spectrum of the maximal operator in \(L^2(\mathbb{R})^2\) associated with \(L\) coincides with the conditional stability set \(S(L)\) of \(L\). One can also prove that the conditional stability set \(S(L)\) consists of a countable number of regular analytic arcs, the spectral bands. At most two spectral bands extend to infinity and at most finitely many spectral bands are closed arcs. The point \(z\) is a band edge, that is, an endpoint of a spectral band, if and only if \((\text{tr}(T_{11}(z)))^2 - 4\) has a zero of odd order. For additional results on nonself-adjoint periodic Dirac-type operators, see Tkachenko [214], [216], [217].

The boundary value problem defined by the differential expression \(L\) and the requirement that the first component of a solution vanishes at \(x_0\) and \(x_0 + \Omega\) will (somewhat artificially but in close analogy to the KdV case) be called the Dirichlet problem for the interval \([x_0, x_0 + \Omega]\). The Dirichlet eigenvalues and their algebraic multiplicities are given as the zeros and their multiplicities of the function

\[ g(z, x_0) = (1, 0)\Phi(z, x_0 + \Omega, x_0, I)(0, 1)^t, \]

that is, the entry in the upper right corner of \(\Phi(z, x_0 + \Omega, x_0, I)\). The algebraic multiplicity of \(z\) as a Dirichlet eigenvalue \(\mu(x)\) is denoted by \(d(z, x)\). The quantities

\[ d_i(z) = \min\{d(z, x) | x \in \mathbb{R}\} \quad \text{and} \quad d_m(z, x) = d(z, x) - d_i(z) \]

will be called the immovable part and the movable part of the algebraic multiplicity \(d(z, x)\), respectively. The sum \(\sum_{z \in \mathbb{C}} d_m(z, x)\), which is independent of \(x\), is called the number of movable Dirichlet eigenvalues. Asymptotically, the Dirichlet eigenvalues are distributed according to

\[ \mu_j(x_0) = \frac{j \pi}{\Omega} + O(|j|^{-1}), \quad j \in \mathbb{Z}. \]

If \(q(x) \neq 0\), the function \(g(\cdot, x)\) is an entire function with order of growth equal to one. Hadamard’s factorization theorem then implies \(g(z, x) = F(z, x)D(z)\), where \(F\) comprises the factors depending on \(x\), while \(D\) contains the factors independent of \(x\). Of course, \(\text{ord}_z(F(\cdot, x)) = d_m(z, x)\) and \(\text{ord}_z(D) = d_i(z)\).

We now turn to the \(x\)-dependence of the function \(g\). The following is the analog of equation \((3.1.2)\) and is obtained by a straightforward computation. The function \(F(z, \cdot)\) satisfies the differential equation

\[ \begin{align*}
q(x)(F'(z, x)^2 - 2F(z, x)F''(z, x) + 4(p(x)q(x) - z^2)F(z, x)^2) \\
+ q'(x)(2F(z, x)F'(z, x) + 4izF(z, x)^2) = q(x)^3((\text{tr}(T_{11}(z)))^2 - 4)/D(z)^2.
\end{align*} \tag{4.2.3} \]

Just as in the KdV case we get the following theorem.

**Theorem 4.3.** There exists an entire function \(R\) such that \((\text{tr}(T_{11}(z)))^2 - 4 = R(z)D(z)^2\). In particular, \(p(z) - 2d_i(z) \geq 0\) for every \(z \in \mathbb{C}\).

We define the Floquet deficiency of \(Q\) (or \(L\)) as \(\text{def}(Q) = \deg(R) \in \mathbb{N}_0 \cup \{\infty\}\). Equation \((4.2.3)\), Theorem \([1.3]\) and the asymptotic distribution of (anti-)periodic and Dirichlet eigenvalues allow one to prove the following central theorem.
Theorem 4.4. The following statements are equivalent:
(1) $\text{def}(Q) = 2n + 2$.
(2) There are $n$ movable Dirichlet eigenvalues.
(3) There exists a $2 \times 2$ matrix-valued differential expression $P_{n+1}$ of order $n + 1$
and leading coefficient $J^{n+2}$ which commutes with $L$. The number $n + 1$ is the
smallest integer with that property. The differential expression $P_{n+1}$ satisfies
$P_{n+1}^2 = R_{2n+2}(L) = \prod_{z \in \mathbb{C}}(L - z)^{p(z) - 2d(z)}$ and hence $Q$ is algebro-geometric.

We remark that in (3) $R_{2n+2}$ is a constant multiple of $R$.

4.3. Necessary Conditions for an AKNS Potential to be Algebro-Geometric. In this section we will prove the following result.

Theorem 4.5. Suppose $Q$ is an algebro-geometric AKNS potential. Then $LY = JY' + QY = zy$ has a meromorphic fundamental system of solutions with respect to $x$ for all values of the spectral parameter $z \in \mathbb{C}$.

A proof of this theorem along the lines of the proof of Theorem 4.10 fails since one can only show that the exponents of a singularity are half integers rather than integers. However, such reasoning can be used to prove the following result.

Theorem 4.6. Suppose $Q$ is a meromorphic potential coefficient of $L$. Then the equation $(LY)(x) = zY(x)$ has a fundamental system of solutions meromorphic with respect to $x$ for all values of the spectral parameter $z \in \mathbb{C}$ whenever this holds for infinitely many values of $z$.

The rest of this section is devoted to a proof of Theorem 4.6. Suppose $Q$ is an algebro-geometric AKNS potential. According to the characterizations following Definition 4.3 there is a function $F_n : \mathbb{C}^2 \to \mathbb{C}_\infty$ which is a polynomial of degree $n$ in the first variable, meromorphic in the second, such that the expression $R_{2n+2}(z, x)$ in (4.1.5) does not depend on $x$. With the aid of $F_n$ we define (as above)

$V_{n+1}(z, x) = \frac{-1}{2q(x)}(q(x)F_n'(z, x) + 2izF_n(z, x)),$

$W_{n+1}(z, x) = \frac{i}{q(x)}\begin{pmatrix} 0 & q(x)F_n(z, x) \\ -F_n'(z, x) + p(x)F_n(z, x) & 0 \end{pmatrix},$

and

$P_{n+1} = JV_{n+1}(L, x) + W_{n+1}(L, x).$

The pair $(P_{n+1}, L)$ is associated with the hyperelliptic curve

$K_n = \{(z, w) : w^2 + R_{2n+2}(z) = 0\}$

of arithmetic genus $n$, where

$R_{2n+2}(z)I = V_{n+1}(z, x)^2I - W_{n+1}(z, x)^2 = \prod_{m=1}^{2n+2} (z - E_m)I.$

We emphasize that the points $E_m$ are not necessarily distinct.

If $\Psi = (\psi_1, \psi_2)^t$ is a common solution of $L\Psi = z\Psi$ and $P_{n+1}\Psi = w\Psi$, then, considering the first component of $P_{n+1}\Psi$, we find

$w\psi_1 = (P_{n+1}\Psi)_1 = (JV_{n+1}(z, x) + W_{n+1}(z, x))\Psi_1 = iV_{n+1}\psi_1 + W_{n+1, 1, 2}\psi_2$

$= (iV_{n+1} + W_{n+1, 1, 2}\phi)\psi_1,$
defining $\phi = \psi_2/\psi_1$. We now revert this process and define the meromorphic function $\phi$ on $K_n \times \mathbb{C}$ by

$$
\phi((z, w), x) = \frac{w - iV_{n+1}(z, x)}{W_{n+1,1,2}(z, x)} = \frac{W_{n+1,2,1}}{w + iV_{n+1}(z, x)}.
$$

We remark that $\phi$ can be extended to a meromorphic function on the compactification (projective closure) of the affine curve $K_n$. This compactification is obtained by joining two points (the points at infinity) to $K_n$.

Next we define

$$
\psi_1((z, w), x, x_0) = \exp\left(\int_{x_0}^{x} (-iz + q(x')\phi((z, w), x'))dx'\right),
$$

$$
\psi_2((z, w), x, x_0) = \phi((z, w), x)\psi_1((z, w), x, x_0)
$$

where the simple Jordan arc from $x_0$ to $x$ in $[3.3]$ avoids poles of $q$ and $\phi$. One verifies with the help of $V'_{n+1} = W_{n+1}Q + QW_{n+1}$ and $W'_{n+1} = 2V_{n+1}Q - 2zJW_{n+1}$ that

$$
\phi'((z, w), x) = p(x) - q(x)\phi((z, w), x)^2 + 2iz\phi((z, w), x).
$$

From this one obtains that

$$
\Psi((z, w), x, x_0) = \begin{pmatrix} \psi_1((z, w), x, x_0) \\ \psi_2((z, w), x, x_0) \end{pmatrix}
$$

is a common solution of $L\Psi = z\Psi$ and $P_{n+1}\Psi = w\Psi$.

Let $\mu_1(x_0), ..., \mu_n(x_0)$ denote the zeros of $F_n(\cdot, x_0)$. One then observes that the two branches $\Psi_{\pm}(z, \cdot, x_0) = \Psi((z, \pm w, \cdot, x_0)$ of the function $\Psi((z, w), \cdot, x_0)$ represent a fundamental system of solutions of $Ly = zy$ for all complex numbers $z$ different from $E_1, ..., E_{2n+2}, \mu_1(x_0), ..., \mu_n(x_0)$, since

$$
W(\Psi_-(z, x, x_0), \Psi_+(z, x, x_0)) = \frac{2w}{W_{n+1,1,2}(z, x_0)} = \frac{-2iw}{F_n(z, x_0)},
$$

where $W(f, g)$ denotes the determinant of the two columns $f$ and $g$.

In the special case where $K_n$ is nonsingular, that is, when the points $E_m$ are pairwise distinct, the explicit representation of $\Psi((z, w), x, x_0)$ in terms of the Riemann theta function associated with $K_n$ immediately proves that $\Psi_{\pm}(z, x, x_0)$ are meromorphic with respect to $x \in \mathbb{C}$ for all $z \in \mathbb{C}\{E_1, ..., E_{2n+2}, \mu_1(x_0), ..., \mu_n(x_0)\}$.

A detailed account of this theta function representation can be found, for instance, in Theorem 3.5 of [28]. In the following we demonstrate how to use gauge transformations to reduce the case of singular curves $K_n$ to nonsingular ones.

Let $Q$ be meromorphic on $\mathbb{C}$ and introduce $A(z, x) = (Q(x) + zI)J$ which turns $LY = zY$ into $Y' + AY = 0$. Then consider the gauge transformation

$$
\tilde{\Psi}(z, x) = \Gamma(z, x)\Psi(z, x).
$$

If

$$
\tilde{A}(z, x) = \Gamma(z, x)A(z, x)\Gamma(z, x)^{-1} - \Gamma'(z, x)\Gamma(z, x)^{-1},
$$

then $\tilde{A}(z, x) = (Q(x) + zI)J$, where $\tilde{Q}$ has zero diagonal elements. Moreover, $\Psi' + \tilde{A}\Psi = 0$, that is, $(Jd/dx + \tilde{Q})\Psi = z\Psi$. Next we make the following explicit choice suggested by Konopelchenko [28]. Let $\tilde{z} \in \mathbb{C}$ be fixed, $\Psi^{(0)}(\tilde{z}, \cdot) = (\psi_1^{(0)}(\tilde{z}, \cdot), \psi_2^{(0)}(\tilde{z}, \cdot))^t$ be any solution of $Ly = zy$, and introduce

$$
\phi^{(0)}(\tilde{z}, x) = \psi_2^{(0)}(\tilde{z}, x)/\psi_1^{(0)}(\tilde{z}, x).
$$
Then one defines
\[ \Gamma(z, x) = \frac{1}{2} \begin{pmatrix} 2(z - \bar{z}) - iq(x)\phi^{(0)}(\bar{z}, x) \\ iq^{(0)}(\bar{z}, x) \\ -i \end{pmatrix} \]
for \( z \in \mathbb{C} \setminus \{\bar{z}\} \).

The upper right entry \( G_{1,2}(z, x, x') \) of the Green’s matrix of \( L \) is given by
\[ G_{1,2}(z, x, x') = \frac{i\Psi_{+1}(z, x, x_0)\psi_{-1}(z, x', x_0)}{W(\Psi_{-}(z, x, x_0), \Psi_{+}(z, x', x_0))}, \quad x \geq x'. \]

We want to evaluate \( G_{1,2}(z, x, x') \) on its diagonal (i.e., where \( x = x' \)). Since \( \psi_{+1}(z, x_0)\psi_{-1}(z, x_0) = 1 \) we obtain from (4.3.1) and (4.3.3)
\[ G_{1,2}(z, x, x) = i \frac{W_{n+1,1,2}(z, x)}{2w} = -F_{n}(z, x). \]

Next note that
\[ W(\Psi_{-}(z, \cdot), \Psi_{+}(z, \cdot)) = \det(\Gamma(z, x))W(\Psi_{-}(z, \cdot), \Psi_{+}(z, \cdot)) \]
and that \( \det(\Gamma(z, x)) = -i(z - \bar{z})/2 \). With the help of this fact and some computations one finds that the upper right entry \( \tilde{G}_{1,2}(z, x, x) \) of the Green’s matrix of \( Jd/dx + \tilde{Q} \) is
\[ \tilde{G}_{1,2}(z, x, x) = i \frac{\tilde{\psi}_{+1}(z, x)\tilde{\psi}_{-1}(z, x)}{W(\Psi_{-}(z, \cdot), \Psi_{+}(z, \cdot))} = -\tilde{F}_{n+1}(z, x) \frac{2(z - \bar{z})w}{2w}, \]
where \( \tilde{F}_{n+1}(. , x) \) is a polynomial of degree \( n + 1 \) with leading coefficient \(-i\tilde{q}(x)\).

Moreover, \( \tilde{F}_{n+1}(z, x) \) satisfies
\[ \tilde{q}(2\tilde{F}_{n+1}\tilde{\bar{F}}_{n+1}'' - \tilde{F}_{n+1}'' + 4(z^2 - \tilde{q}\tilde{p})\bar{F}_{n+1}^2) - \tilde{q}^2(2\tilde{F}_{n+1}\tilde{F}_{n+1}' + 4iz\bar{F}_{n+1}^2) = -4\tilde{q}^3R_{2n+2}(z). \]

Hence \( \tilde{Q} \) is an algebro-geometric AKNS potential.

Now suppose that \( (\bar{z}, 0) \) is a singular point of \( K_n \), that is, that \( \bar{z} \) is a zero of \( R_{2n+2} \) of order \( r \geq 2 \). Choose
\[ \phi^{(0)}(\bar{z}, x) = -i V_{n+1}(\bar{z}, x) \]
\[ \frac{W_{n+1,1,2}(\bar{z}, x)}{W(\Psi_{-}(\bar{z}, \cdot), \Psi_{+}(\bar{z}, \cdot))}. \]

Then one may show that \( \tilde{F}_{n+1}(., x) \) has a zero of order at least 2 at \( \bar{z} \), that is,
\[ \tilde{F}_{n+1}(z, x) = (z - \bar{z})^s \tilde{F}_{n}(z, x) \]
for some \( s \geq 2 \) and some polynomial \( \tilde{F}_{n}(., x) \) of degree \( \tilde{n} = n + 1 - s \). From (4.3.4) one obtains
\[ \tilde{q}(2\tilde{F}_{n}\tilde{F}_{n}'' - \tilde{F}_{n}'' + 4(z^2 - \tilde{p}\tilde{q})\tilde{F}_{n}^2) - \tilde{q}^2(2\tilde{F}_{n}\tilde{F}_{n}' + 4iz\tilde{F}_{n}^2) = -4\tilde{q}^3\tilde{R}_{2n-2s+4}(z), \]
where
\[ \tilde{R}_{2n-2s+4}(z) = (z - \bar{z})^{2-2s}R_{2n+2}(z) \]
is a polynomial in \( z \) of degree \( 2n - 2s + 4 \in (0, 2n + 2) \). This proves that \( \tilde{Q} \) is associated with the curve
\[ K_{\tilde{n}} = \{(z, w) : w^2 + (z - \bar{z})^{2-2s}R_{2n+2}(z) = 0\}. \]

Our choice of \( \phi^{(0)}(\bar{z}, x) \) led to a curve \( K_{\tilde{n}} \) which is less singular at \( (\bar{z}, 0) \) than \( K_n \), without changing the local structure of the curve elsewhere. By iterating this procedure one ends up with a curve which is nonsingular at \((\bar{z}, 0)\). Repeating
the procedure for each singular point of \( K_n \) then results in a nonsingular curve and a corresponding Baker-Akhiezer function \( \Psi((z, w), x, x_0) \) meromorphic with respect to \( x \in \mathbb{C} \) using its standard theta function representation (see, e.g., [88]). Since \( \phi^{(0)} = -iV_{n+1}/W_{n+1,1,2} \) is meromorphic, the gauge transformations and their inverses map meromorphic functions to meromorphic functions. Combining these results proves Theorem 4.3.

A systematic account of the effect of Darboux-type transformations on hyperelliptic curves in connection with the KdV, AKNS, and Toda hierarchies will be presented in [86].

4.4. Sufficient Conditions for an AKNS Potential to be Algebro-Geometric – A Characterization of Elliptic Stationary AKNS Solutions.

Definition 4.7. Let \( p, q \) be elliptic functions with a common period lattice. Then \( Q = \begin{pmatrix} 0 & -iq \\ ip & 0 \end{pmatrix} \) is called a Picard-AKNS potential if the equation \( J\psi'(x) + Q(x)\psi(x) = z\psi(x) \) has a meromorphic fundamental system of solutions with respect to \( x \) for all values of the spectral parameter \( z \in \mathbb{C} \).

We note that according to Theorem 4.3 it is sufficient to show the existence of a meromorphic fundamental system for infinitely many values of \( z \) in order to prove that \( Q \) is a Picard-AKNS potential.

Just as in the KdV case the following theorem characterizes all elliptic algebro-geometric AKNS potentials.

Theorem 4.8. Let \( Q = \begin{pmatrix} 0 & -iq \\ ip & 0 \end{pmatrix} \) with \( p, q \) elliptic functions with a common period lattice. Then \( Q \) is an elliptic algebro-geometric AKNS potential if and only if it is a Picard-AKNS potential.

Sketch of proof. The necessity of the criterion is the content of Theorem 4.5. The sufficiency follows from Picard’s Theorem 2.7 and Theorem 4.4 in the same way as in the KdV case.

The transformation

\[
\tilde{\Psi}(x) = \begin{pmatrix} e^{ax+b} & 0 \\ 0 & e^{-(ax+b)} \end{pmatrix} \Psi(x)
\]

allows one to prove the following corollary, which slightly extends the class of algebro-geometric AKNS potentials \( Q \) considered thus far. Such cases have recently been considered by Smirnov [200].

Corollary 4.9. Suppose

\[
Q(x) = \begin{pmatrix} 0 & -iq(x)e^{-2(ax+b)} \\ ip(x)e^{2(ax+b)} & 0 \end{pmatrix},
\]

where \( a, b \in \mathbb{C} \) and \( p, q \) are elliptic functions with a common period lattice. Then \( Q \) is an algebro-geometric AKNS potential if \( J\tilde{\Psi}'(x) + Q(x)\tilde{\Psi}(x) = z\tilde{\Psi}(x) \) has a meromorphic fundamental system of solutions with respect to \( x \) for all values of the spectral parameter \( z \in \mathbb{C} \).
Examples. With the exception of the studies by Christiansen, Eilbeck, Enol’skii, and Kostov \[40\] and Smirnov \[200\], \[203\], \[204\], not too many examples of elliptic solutions \((p,q)\) of the AKNS hierarchy associated with higher (arithmetic) genus curves of the type \(w^2 + R_{2n+2}(z) = 0\) have been worked out in detail. The genus \(n = 1\) case has been considered, for example, by Its \[121\] and Pavlov \[177\]. Moreover, examples for low genus \(n\) for special cases such as the nonlinear Schrödinger and mKdV equation are considered, for instance, in \[8\], \[15\], \[148\], \[162\], \[156\], \[172\], \[201\]. Examples related to equations of the sine-Gordon-type are discussed in \[197\], \[205\], \[206\]. The following examples in (4.5.1) – (4.5.4) are algebro-geometric AKNS potentials as can be proved using the Frobenius method. For details on this procedure see \[99\].

\[
Q(x) = i(n(\zeta(x) - \zeta(x - \omega_2) - \zeta(\omega_2))) \begin{pmatrix} 0 & -1 \\ 1 & 0 \end{pmatrix}, \quad n \in \mathbb{N},
\]

(4.5.2)

\[
Q(x) = \begin{pmatrix} 0 & -i(n+1)\phi(x) \\ 1 & 0 \end{pmatrix}, \quad n \in \mathbb{N},
\]

(4.5.3)

\[
Q(x) = \begin{pmatrix} 0 & i\psi'(x - \omega_2)/(2e_1) \\ -3i\psi'(x)/(2e_1) & 0 \end{pmatrix}, \quad e_2 = 0,
\]

(4.5.4)

\[
Q(x) = \begin{pmatrix} 0 & i\psi(x - \omega_2)/e_1^2 \\ 2i(\psi''(x) - c_1^2)/3 & 0 \end{pmatrix}, \quad e_2 = 0.
\]

Incidentally, if \(p = 1\), as in Example (4.5.2) then \(J\Psi + Q\Psi = z\Psi\) is equivalent to the scalar equation \(\psi''_2 - q\psi_2 = -z^2\psi_2\), where \(\Psi = (\psi_1, \psi_2)^t\) and \(\psi_1 = \psi_2 - iz\psi_2\). Therefore, if \(-q\) is an elliptic algebro-geometric potential of the KdV hierarchy then, by Theorem 3.12, \(\psi_2\) is meromorphic for all values of \(z\). Hence \(\Psi\) is meromorphic for all values of \(z\) and therefore \(Q\) is a Picard-AKNS and hence an algebro-geometric AKNS potential. Conversely, if \(Q\) is an algebro-geometric AKNS potential with \(p = 1\) then \(-q\) is an algebro-geometric potential of the KdV hierarchy. In particular, \(q(x) = n(n+1)\phi(x)\) is again the class of Lamé potentials associated with the KdV hierarchy and hence a special case of the material discussed in Section 3.2.

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