

COMMUTATION METHODS APPLIED TO THE mKdV-EQUATION

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Dedicated to Raphael Hoegh-Krohn (1938–1988)

ABSTRACT. An explicit construction of solutions of the modified Korteweg-de Vries equation given a solution of the (ordinary) Korteweg-de Vries equation is provided. Our theory is based on commutation methods (i.e., $N = 1$ supersymmetry) underlying Miura's transformation that links solutions of the two evolution equations.

1. INTRODUCTION

In connection with the extensively studied Korteweg-de Vries (KdV-) equation

$$(1.1) \quad \text{KdV}(V) := V_t - 6V V_x + V_{xxx} = 0, \quad (t, x) \in \mathbb{R}^2,$$

its cousins, the modified Korteweg-de Vries (mKdV $_{\pm}$) equations

$$(1.2\pm) \quad \text{mKdV}_{\pm}(\phi) := \phi_t \pm 6\phi^2\phi_x + \phi_{xxx} = 0, \quad (t, x) \in \mathbb{R}^2,$$

have also been investigated. In fact, (1.2+) has been treated in some detail in the literature: See, e.g., [12, 51, 67, 68, 69, 74, 111] for existence and uniqueness questions, [58, 105, 107, 113, 114] for the derivation of the N -soliton solutions and [94, 108] for the general approach to (1.2+) via inverse scattering techniques. The Lax pair for (1.2 \pm) has been given in [107] and finally (1.2 \pm) have been found to be subordinate to the AKNS-ZS-theory [4, 93] (cf. [34]). For more recent work on (1.2+) see, e.g., [11, 13, 70, 118].

Surprisingly enough, apart from existence and uniqueness questions of solutions of (1.2-) in [12, 51, 68, 69, 74, 111], no detailed study of (1.2-) seems to have appeared until 1984 when Grosse [54] (see also [55]) derived the N -soliton solutions of (1.2-).

The present paper is devoted to a detailed investigation of real-valued solutions of the mKdV $_{-}$ -equation (1.2-) (from now on simply denoted by mKdV).

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Our theory is based on Miura's fundamental relation [88]

$$(1.3) \quad \text{KdV}(\phi^2 \pm \phi_x) = [2\phi \pm \partial_x] \text{mKdV}(\phi)$$

that yields solutions $\phi^2 \pm \phi_x$ of the KdV-equation (1.1) given a solution ϕ of the mKdV-equation (1.2-). In view of the fact that mKdV has a higher nonlinearity than KdV and the latter has been studied so thoroughly in the literature, one would clearly like to invert that procedure, i.e., given a solution V of (1.1) one would like to construct all solutions ϕ of (1.2-) that are linked to each other via Miura's transformation

$$(1.4) \quad V = \phi^2 - \phi_x \quad (\text{resp. } V = \phi^2 + \phi_x).$$

For smooth, bounded, and real-valued solutions V, ϕ this has been achieved in [50] from a theoretical point of view (see Theorem 7.9). Our main objective in this paper is to illustrate this method at work, i.e., to explicitly construct new classes of solutions of the mKdV-equation such as soliton-like solutions, periodic solutions, solitons relative to periodic background solutions and to describe some of their properties. (Here soliton-like means solutions $V(\phi)$ of the (m)KdV-equation that approach time-independent (possibly different) asymptotic values $V_{\pm}(\phi_{\pm}) \in \mathbb{R}$ as $x \rightarrow \pm\infty$ and periodic always refers to spatially periodic solutions, i.e., for some $a > 0$:

$$V(t, x + a) = V(t, x)(\phi(t, x + a) = \phi(t, x))$$

for all $(t, x) \in \mathbb{R}^2$.)

Before describing the content of this paper in more detail we would like to explain the main ideas of [50] in an informal manner. Consider the Schrödinger operators

$$(1.5) \quad H_j = -\partial_x^2 + V_j, \quad j = 1, 2,$$

in $L^2(\mathbb{R})$, where the V_j are defined according to Miura's transformation (1.4), i.e.,

$$(1.6) \quad V_j = \phi^2 + (-1)^j \phi_x, \quad j = 1, 2.$$

(In order to simplify matters, V_j, ϕ , and its partial derivatives are assumed to be real-valued, smooth, and bounded.) By inspection one verifies the factorizations

$$(1.7) \quad H_1 = A^* A, \quad H_2 = A A^*,$$

where

$$(1.8) \quad A = \partial_x + \phi.$$

In particular, this forces H_j to be nonnegative, i.e.,

$$(1.9) \quad H_j \geq 0, \quad j = 1, 2.$$

In order to solve the Riccati-equation (1.6) for ϕ if, for instance, V_1 is given, one introduces the linearization

$$(1.10) \quad \phi = -\psi_{0,x}/\psi_0,$$

where ψ_0 satisfies

$$(1.11) \quad H_1 \psi = 0 \quad \text{and hence} \quad A\psi = 0$$

(the latter equation being equivalent to (1.10)) in the distributional sense. Equation (1.10) shows two things: First of all we need $\psi_0 > 0$ in order to get nonsingular ϕ 's. Secondly, ϕ will be uniquely determined by V_1 only if the equation $H_1 \psi = 0$ has a unique (up to multiple of constants) positive solution ψ_0 . Otherwise, we expect a one-parameter family of solutions

$$(1.12) \quad \phi_\sigma = -\partial_x \ln \psi_\sigma,$$

where

$$(1.13) \quad \psi_\sigma = (1 - \sigma)\psi_- + (1 + \sigma)\psi_+, \quad \sigma \in [-1, 1], \quad \psi_\pm > 0, \quad H_1 \psi_\pm = 0,$$

and ψ_+ and ψ_- are linearly independent. Thus we are forced to a close examination of the zero-energy spectral properties of H_1 which will lead us to the concept of (sub)criticality of H_1 . (Roughly speaking, H_1 is (sub)critical if the equation $H_1 \psi = 0$ has 1 (2 linearly independent) positive distributional solution(s).)

So far we have ignored the t -dependence of V_1 and ϕ . Let us assume that V_1 satisfies the KdV-equation (1.1) in what follows. Then the commutation relation for the Lax pair $(H_1(t), B_{V_1}(t))$,

$$(1.14) \quad \partial_t H_1 - [B_{V_1}, H_1] = \text{KdV}(V_1),$$

where

$$(1.15) \quad B_{V_1}(t) = -4\partial_x^3 + 6V(t, \cdot)\partial_x + 3V_x(t, \cdot), \quad t \in \mathbb{R},$$

proves that $H_1(t)$ is unitarily equivalent to $H_1(0)$ for all $t \in \mathbb{R}$ and hence (1.9) holds for all $t \in \mathbb{R}$ if it holds, e.g., at $t = 0$. Moreover, assuming that the ψ_\pm in (1.13) evolve according to B_{V_1} , i.e.,

$$(1.16) \quad \psi_{\pm,t}(t, x) = (B_{V_1}(t)\psi_{\pm}(t))(x), \quad (t, x) \in \mathbb{R}^2,$$

fixes the time-dependence of ψ_\pm in such a way that the Wronskian of ψ_+ and ψ_- becomes t -independent. In particular, the KdV-flow leaves the number of linearly independent, positive solutions of $H_1(t)\psi(t) = 0$ invariant. In other words, given $V_1(t, x)$ with $\text{KdV}(V_1) = 0$, we either get a unique solution $\phi_0(t, x)$ of (1.6) for $j = 1$ or we get a one-parameter family $\phi_\sigma(t, x)$, $\sigma \in [-1, 1]$, of such solutions. What remains to be indicated is why, in either case, ϕ_σ also satisfies the mKdV-equation (1.2-) if σ is time-independent. This

can be made plausible by recalling that the mKdV-equation subordinates to the AKNS-ZS framework and

$$(1.17) \quad Q(t) = \begin{pmatrix} 0 & A(t)^* \\ A(t) & 0 \end{pmatrix}, \quad B(t) = B_{V_1}(t) \oplus B_{V_2}(t),$$

$$A(t) = \partial_x + \phi(t, \cdot), \quad V_j = \phi^2 + (-1)^j \phi_x,$$

represent a Lax pair for the mKdV-equation, i.e.,

$$(1.18) \quad \partial_t Q - [B, Q] = \text{mKdV}(\phi) \begin{pmatrix} 0 & 1 \\ 1 & 0 \end{pmatrix}.$$

But then the identity

$$(1.19) \quad Q(t)^2 = \begin{pmatrix} A(t)^* A(t) & 0 \\ 0 & A(t) A(t)^* \end{pmatrix} = H_1(t) \oplus H_2(t)$$

shows the intimate connection between the Lax pairs (1.14) and (1.18). In particular, this gives a hint as to why the parameter σ in (1.13) should be time-independent, since otherwise ψ_σ in (1.13) would not evolve according to B_{V_1} (see (1.16)). These statements will be justified in the first part of §7, where we shall reprove the main result of [50], viz. that the time independence of σ in (1.13) is not only sufficient but also necessary for ϕ_σ to be a solution of the mKdV-equation.

Summarizing the above, we infer that the connection between the KdV- and mKdV-equations, effected by the Miura transformation

$$(1.20) \quad \text{KdV}(V_j) = 0 \xleftarrow[\substack{\text{Miura} \\ V_j = \phi^2 + (-1)^j \phi_x}]{\text{Miura}} \text{mKdV}(\phi) = 0, \quad j = 1, 2,$$

reflects itself in the connection between the two Lax pairs

$$(1.21) \quad (H_j, B_{V_j}) \longleftrightarrow (Q, B), \quad Q^2 = H_1 \oplus H_2, \quad B = B_{V_1} \oplus B_{V_2}.$$

While the connection between B_{V_1} , B_{V_2} , and B is simple (see (1.21)), the connection between the Schrödinger operators $H_1 = A^* A$, $H_2 = A A^*$ and the Dirac operator $Q = \begin{pmatrix} 0 & A^* \\ A & 0 \end{pmatrix}$ (being of AKNS-type) is more subtle (see §2).

Thus a realization of the program formulated in the paragraph following (1.4) requires the study of different topics such as commutation methods, i.e., the connections between operators of the type $A^* A$, $A A^*$, and $\begin{pmatrix} 0 & A^* \\ A & 0 \end{pmatrix}$, spectral and scattering theory for Schrödinger and Dirac operators, and inverse spectral and scattering theory for Schrödinger operators. In order to make this exposition reasonably self-contained we have included a concise treatment of these subjects in §§2–6: §2 provides all the general facts on commutation methods needed in §§7–10. While the main body of this section summarizes results from various sources in the literature, Theorem 2.3(ii) (although presumably known to some experts) apparently has not appeared in print before. §3 is devoted to a detailed study of spectral and scattering properties of one-dimensional Schrödinger operators. In Theorems 3.1–3.3 we generalize known results in the case where the

potential decays sufficiently fast as $x \rightarrow \pm\infty$ to the case where the potential has a nontrivial spatial asymptotics. This generalization is crucial in transferring soliton-like solutions of the KdV-equation to that of the mKdV-equation (cf. Remark 7.2(iii)). The remaining material in §3 summarizes some elements of Floquet theory for Schrödinger operators. §4, based on §§2 and 3, carries over the results on Schrödinger operators to that of Dirac operators of the type

$$(1.22) \quad \begin{pmatrix} m & A^* \\ A & -m \end{pmatrix}, \quad m \in \mathbb{R}, \quad A = \frac{d}{dx} + \phi.$$

In particular, Theorems 4.2 and 4.3 and Lemma 4.5 appear to be new. §5 gives a short account of relative scattering to be used in §§8 and 10, where soliton solutions of the (m)KdV-equation relative to a trivial (i.e., constant) or periodic background solution of the (m)KdV-equation are shown to be characterized by the property of being relative reflectionless with respect to the background. In §6 we recall the process of adding eigenvalues into the spectral gaps of a background Hamiltonian since this represents a general method of constructing solitons relative to a nontrivial (not necessarily periodic) background. Theorem 6.1, in the first part of §6, extends a known result where the background potential vanishes sufficiently fast as $x \rightarrow \pm\infty$, while the second part is devoted to a study of the corresponding periodic problem.

§7 is the main section and consists of three parts. In the first one we reprove the basic result of [50] (in fact, we provide more details than in [50]). In particular, we relate the existence of nonsingular solutions ϕ of the mKdV-equation and of Miura's transformation (1.6) with $j = 1$, given a KdV-solution V_1 , to the fact whether $H_1(0)$ in (1.5) is nonnegative or not. Moreover, recalling the notion of (sub)criticality of Schrödinger operators, we show in Theorem 7.9 that ϕ is uniquely determined by V_1 iff $H_1(0)$ is critical. The second part of §7 contains the main body of our new results. Given Theorem 7.9 we describe in detail the construction of soliton-like solutions and their basic properties in Theorem 7.14 and Remarks 7.15–7.17. The case of periodic solutions is treated in Theorem 7.18 and Remark 7.19. It turns out that, given a periodic solution V_1 of the KdV-equation with $H_1(0)$ being subcritical, only the special values $\sigma = \pm 1$ in (1.13) lead to periodic solutions $\phi_{\pm 1}(t, x)$ of the mKdV-equation. The final part of §7 contains hints at various possible extensions of our approach to a generalized mKdV-equation of the type

$$(1.23) \quad \phi_t - 6(\phi^2 - \Phi_0)\phi_x + \phi_{xxx} = 0, \quad \Phi_0 \in \mathbb{R}, \quad (t, x) \in \mathbb{R}^2,$$

and to the entire hierarchy of (m)KdV-equations.

In the remaining §§8–11 we transfer particular classes of solutions of the KdV-equation to the mKdV-equation. §8 (based on §§2–7) is devoted to the transfer of the class of KdV-soliton solutions to the mKdV-equation, thereby providing an alternative way to derive all soliton solutions of the mKdV-equation. (Their original derivation in [54] uses inverse scattering techniques for Dirac systems.) In §9 (using §§2–4 and 7) we consider the special example

of a periodic two-zone solution V_1 of the KdV-equation and transfer it to the mKdV-equation. In the course of this we also extend a well-known theorem of Hochstadt's [59] in connection with periodic two-zone Schrödinger operators to the case of Dirac operators (1.22). In §10 (relying on §§2–9) we first recall the construction of solitons relative to the periodic two-zone solution of the KdV-solution V_1 considered in §9 (originally this has been done in [76]) and then transfer this solution to the mKdV-equation. Finally, in §11, we indicate how to extend the framework of §7 to singular solutions of the (m)KdV-equation by dropping the condition of nonnegativity of $H_1(0)$. All mKdV-results in §§9–11 are new. In contrast to §8 which is complete, §§9–11 only present the simplest possible nontrivial solutions of the corresponding classes involved, e.g., a transfer of the whole class of periodic n -zone solutions ($n \in \mathbb{N}$) of the KdV-equation to that of the mKdV-equation based on hyperelliptic function theory (although particularly interesting because of links to algebraic geometry) would necessarily catapult this paper beyond any reasonable length. Thus we confined ourselves in §§9 and 10 to the case of elliptic functions only. A similar remark applies to §11. (We note, however, that a transfer of the full class of periodic n -zone KdV-solutions expressed in terms of Riemann's theta function associated with the underlying hyperelliptic curve of genus $n - 1$ and a transfer of a large class of singular KdV-solutions expressed in terms of Wronski determinants has been completed in the meantime and will be published elsewhere.)

We finally remark that the methods developed in [50] and in the present paper are by no means restricted to the KdV- and mKdV-equations. In fact, they also apply in a $1 + 2$ -dimensional context such as the Kadomtsev-Petviashvili-equation and its modified analog as well as in the context of discrete nonlinear systems such as the Toda lattice and its modified version the Kac-van Moerbeke lattice (corresponding manuscripts are in preparation).

2. BASIC FACTS ON COMMUTATION

In this section we review certain formulas based on commutation and add a few new results of this type in Theorem 2.3(ii).

We introduce the following hypothesis:

(H.2.1) Let \mathfrak{H}_j , $j = 1, 2$, be separable, complex Hilbert spaces and $A : \mathfrak{D}(A) \subseteq \mathfrak{H}_1 \rightarrow \mathfrak{H}_2$ a densely defined, closed linear operator.

We are particularly concerned with abstract Dirac operators in $\mathfrak{H}_1 \oplus \mathfrak{H}_2$ of the type

$$(2.1) \quad Q := \begin{pmatrix} 0 & A^* \\ A & 0 \end{pmatrix}, \quad \mathfrak{D}(Q) = \mathfrak{D}(A) \oplus \mathfrak{D}(A^*),$$

and

$$(2.2) \quad Q_{m_1, m_2} := \begin{pmatrix} m_1 1 & A^* \\ A & -m_2 1 \end{pmatrix}, \\ \mathfrak{D}(Q_{m_1, m_2}) = \mathfrak{D}(Q), \quad m_j \in \mathbb{R}, \quad j = 1, 2.$$

Because of the identity

$$(2.3) \quad Q_{m_1, m_2} = [(m_1 - m_2)/2] 1 + Q_{(m_1+m_2)/2},$$

where

$$(2.4) \quad Q_m := \begin{pmatrix} m1 & A^* \\ A & -m1 \end{pmatrix}, \quad m \in \mathbb{R},$$

we shall only consider the symmetric version (2.4). We note that

$$(2.5) \quad Q_m^2 = \begin{pmatrix} H_{1,m} & 0 \\ 0 & H_{2,m} \end{pmatrix} = Q^2 + m^2 1, \quad m \in \mathbb{R},$$

where

$$(2.6) \quad \begin{aligned} H_{1,m} &:= A^*A + m^2 \quad \text{on } \mathfrak{D}(A^*A) \subseteq \mathfrak{H}_1, \\ H_{2,m} &:= AA^* + m^2 \quad \text{on } \mathfrak{D}(AA^*) \subseteq \mathfrak{H}_2, \quad m \in \mathbb{R}, \\ H_{1,0} &:= H_1 = A^*A, \quad H_{2,0} := H_2 = AA^* \end{aligned}$$

(we delete the symbol 1 in multiples of the identity operator in \mathfrak{H}_j , $j = 1, 2$, and $\mathfrak{H}_1 \oplus \mathfrak{H}_2$ from now on).

Next we collect a few useful commutation formulas.

Theorem 2.1 [32]. *Assume hypothesis (H.2.1) and let $f \in L^\infty(\mathbb{R})$. Then*

(i)

$$(2.7) \quad A = SH_1^{1/2} = H_2^{1/2}S,$$

where $S := \text{sgn}(A)$ denotes the partial isometry with initial set $\overline{\text{Ran}[(A^*A)^{1/2}]}$ and final set $\overline{\text{Ran}(A)}$ in the polar decomposition of A .

(ii)

$$(2.8) \quad \text{sgn}(Q) = \begin{pmatrix} 0 & S^* \\ S & 0 \end{pmatrix}$$

is the corresponding partial isometry in the polar decomposition of Q .

(iii)

$$(2.9) \quad \begin{aligned} \text{Ker}(A) &= \text{Ker}(H_1) = \text{Ran}(A^*)^\perp, \\ \text{Ker}(A^*) &= \text{Ker}(H_2) = \text{Ran}(A)^\perp. \end{aligned}$$

(iv)

$$(2.10) \quad SH_1^{n/2} \subseteq H_2^{n/2}S, \quad n \in \mathbb{N},$$

$$(2.11) \quad Sf(H_1) = f(H_2)S.$$

In particular, H_1 and H_2 are essentially isospectral, i.e.,

$$(2.12) \quad \sigma(H_1) \setminus \{0\} = \sigma(H_2) \setminus \{0\}$$

and

$$(2.13) \quad \begin{aligned} H_1 \psi &= E \psi, & E \neq 0 \Rightarrow H_2(A\psi) &= E(A\psi), & \psi &\in \mathfrak{D}(H_1), \\ H_2 \tilde{\psi} &= \tilde{E} \tilde{\psi}, & \tilde{E} \neq 0 \Rightarrow H_1(A^* \tilde{\psi}) &= \tilde{E}(A^* \tilde{\psi}), & \tilde{\psi} &\in \mathfrak{D}(H_2), \end{aligned}$$

with multiplicities preserved.

(v)

$$(2.14) \quad \begin{aligned} 1 + z(H_2 - z)^{-1} &\supseteq A(H_1 - z)^{-1} A^*, \\ 1 + z(H_1 - z)^{-1} &\supseteq A^*(H_2 - z)^{-1} A, \quad z \in \mathbb{C} \setminus \{\sigma(H_1) \cup \sigma(H_2)\}. \end{aligned}$$

(vi)

$$(2.15) \quad A^* f(H_2) \supseteq f(H_1) A^*, \quad A f(H_1) \supseteq f(H_2) A.$$

Sketch of proof. Let T denote a densely defined, closed linear operator from $\mathfrak{D}(T) \subseteq \mathfrak{H}' \rightarrow \mathfrak{H}''$, where \mathfrak{H}' , \mathfrak{H}'' are separable, complex Hilbert spaces. Then $\text{sgn}(T)$ is defined by

$$(2.16) \quad \text{sgn}(T) := \begin{cases} \overline{T(T^*T)^{-1/2}} & \text{on } \text{Ker}(T)^\perp, \\ 0 & \text{on } \text{Ker}(T). \end{cases}$$

The spectral theorem for selfadjoint operators together with the following obvious equalities then proves (i)–(vi):

$$(2.17) \quad \begin{aligned} Q &= |Q| \text{sgn}(Q) = \text{sgn}(Q) |Q|, \\ \text{Ker}(Q) &= \text{Ker}(Q^2) = \text{Ran}(Q)^\perp, \\ 1 + z(Q^2 - z)^{-1} &= Q^2(Q^2 - z)^{-1} \supseteq Q(Q^2 - z)^{-1} Q, \quad z \in \mathbb{C} \setminus \sigma(Q^2), \\ Q f(Q^2) &\supseteq f(Q^2) Q. \quad \square \end{aligned}$$

Remark 2.2. Theorem 2.1 is due to Deift [32]. The idea of proof, presented above, is due to Nelson (unpublished) and has been summarized in [109].

Theorem 2.3. Assume (H.2.1). Then

(i)

$$(2.18) \quad (Q_m - z)^{-1} = \begin{pmatrix} (H_{1,m} - z^2)^{-1}(z + m) & A^*(H_{2,m} - z^2)^{-1} \\ A(H_{1,m} - z^2)^{-1} & (H_{2,m} - z^2)^{-1}(z - m) \end{pmatrix},$$

$$z^2 \in \mathbb{C} \setminus \{\sigma(H_{1,m}) \cup \sigma(H_{2,m})\}.$$

(ii) Define

$$(2.19) \quad R := i\sigma_3 \text{sgn}(Q) = \begin{pmatrix} 0 & iS^* \\ -iS & 0 \end{pmatrix}, \quad \sigma_3 = \begin{pmatrix} 1 & 0 \\ 0 & -1 \end{pmatrix};$$

then

$$(2.20) \quad \text{Ker}(R) = \text{Ker}(Q) = \text{Ker}(A) \oplus \text{Ker}(A^*),$$

$$(2.21) \quad R^2 = \begin{pmatrix} S^*S & 0 \\ 0 & SS^* \end{pmatrix}, \quad Q_m R = -R Q_m, \quad m \in \mathbb{R}$$

(together with $A = S|A| = |A^*|S$, $A^* = S^*|A^*| = |A|S^*$), shows that Q_m is unitarily equivalent to $-Q_m$ on the subspace $\text{Ker}(Q)^\perp$ of $\mathfrak{H}_1 \oplus \mathfrak{H}_2$, i.e.,

$$(2.22) \quad Q_m = -RQ_mR \quad \text{on } \text{Ker}(Q)^\perp, \quad m \in \mathbb{R}.$$

Moreover,

$$(2.23) \quad Q_m^2 \geq m^2, \quad m \in \mathbb{R},$$

and Q is unitarily equivalent to $-Q$, i.e.,

$$(2.24) \quad \sigma_3 Q \sigma_3 = -Q.$$

Proof. (i) is most easily verified by multiplying with $(Q_m - z)$ from the left and right respectively. (ii) is obvious from the fact that S^*S is the projection onto $\text{Ker}(A)$ and SS^* is the projection onto $\text{Ker}(A^*)$. \square

Remark 2.4. (i) By identity (2.3), these results extend to Q_{m_1, m_2} in a straightforward manner.

(ii) Since by (2.22),

$$(2.25) \quad Q_m^2 = RQ_m^2R \quad \text{on } \text{Ker}(Q)^\perp, \quad m \in \mathbb{R},$$

we infer that $A^*A|_{\text{Ker}(A)^\perp}$ and $AA^*|_{\text{Ker}(A^*)^\perp}$ are unitarily equivalent, i.e.,

$$(2.26) \quad \begin{aligned} A^*A &= S^*AA^*S \quad \text{on } \text{Ker}(A)^\perp, \\ AA^* &= SA^*AS^* \quad \text{on } \text{Ker}(A^*)^\perp. \end{aligned}$$

Thus we recover Theorem 3 of [32].

The above assertions may be supplemented by the following results.

Theorem 2.5 [109]. Assume (H.2.1).

(i) There exists a unitary operator U_m on $\text{Ker}(Q)^\perp$ such that

$$(2.27) \quad U_m Q_m U_m^* = \begin{pmatrix} H_{1,m}^{1/2} & 0 \\ 0 & -H_{2,m}^{1/2} \end{pmatrix} \quad \text{on } \text{Ker}(Q)^\perp,$$

$$(2.28) \quad U_m = \begin{pmatrix} a_{1,+} & a_{1,-}S^* \\ -a_{2,-}S & a_{2,+} \end{pmatrix},$$

$$a_{j,\pm} := 2^{-1/2} \{1 \pm m(H_j + m^2)^{-1/2}\}^{1/2}, \quad j = 1, 2, \quad m \in \mathbb{R}.$$

(ii) There exists a unitary operator W_m on $\text{Ker}(Q)^\perp$ such that

$$(2.29) \quad W_m Q_m W_m^* = |Q_m| \text{sgn}(Q) = \begin{pmatrix} 0 & H_{1,m}^{1/2} S^* \\ H_{2,m}^{1/2} S & 0 \end{pmatrix} \quad \text{on } \text{Ker}(Q)^\perp,$$

$$(2.30) \quad W_m := 2^{-1/2} \begin{pmatrix} 1 & -S^* \\ S & 1 \end{pmatrix} U_m, \quad m \in \mathbb{R}.$$

3. ONE-DIMENSIONAL SCHRÖDINGER OPERATORS

In this section we study Schrödinger operators $H_1 = A^*A$, $H_2 = AA^*$ obtained from a concrete realization of A in $L^2(\mathbb{R})$, formally given by $\frac{d}{dx} + \phi(x)$, where ϕ satisfies one of the following hypotheses:

(H.3.1) $\phi, \phi' \in L^\infty(\mathbb{R})$ real-valued,

$$\lim_{x \rightarrow \pm\infty} \phi(x) := \phi_\pm \in \mathbb{R}, \quad \phi_-^2 \leq \phi_+^2,$$

$$\pm \int_0^{\pm\infty} dx(1+x^2)|\phi(x) - \phi_\pm| < \infty, \quad \int_{\mathbb{R}} dx(1+x^2)|\phi'(x)| < \infty.$$

(H.3.2) $\phi, \phi' \in L^\infty(\mathbb{R})$ real-valued, for some $a > 0$: $\phi(x+a) = \phi(x)$, $x \in \mathbb{R}$.

The factorization of 2nd-order Sturm-Liouville operators has a long history: It dates back at least to Jacobi [63], respectively Darboux [29]. Other important contributions are [19, 20, 28, 99] (see also the references therein) and especially [32]. Quite recently these ideas became popular again in connection with supersymmetric quantum mechanics [8, 16, 17, 24, 49, 64, 91, 104, 109] (see also the references cited therein).

In both cases (H.3.1) and (H.3.2) A is defined in $L^2(\mathbb{R})$ by

$$(3.1) \quad A = \frac{d}{dx} + \phi, \quad \mathfrak{D}(A) = H^1(\mathbb{R}),$$

and hence

$$(3.2) \quad H_j = -\frac{d^2}{dx^2} + \phi^2 + (-1)^j \phi' := -\frac{d^2}{dx^2} + V_j,$$

$$\mathfrak{D}(H_j) = H^2(\mathbb{R}), \quad j = 1, 2.$$

Next we briefly summarize spectral and scattering theory for systems governed by (H.3.1). (For more details see, e.g., [16, 23, 27, 31, 33, 49, 96].) Let $f_{j,\pm}(k_\pm, x)$, defined by

$$(3.3) \quad f_{j,\pm}(k_\pm, x) = e^{\pm i k_\pm x} - \int_x^{\pm\infty} dx' k_\pm^{-1} \sin[k_\pm(x-x')] [V_j(x') - V_\pm] f_{j,\pm}(k_\pm, x'),$$

$$(3.4) \quad k_\pm := (z - V_\pm)^{1/2}, \quad \text{Im } k_\pm \geq 0, \quad z \in \mathbb{C}, \quad V_\pm := \phi_\pm^2, \quad j = 1, 2,$$

be the Jost solutions of H_j , i.e.,

$$(3.5) \quad H_j f_{j,\pm}(k_\pm, x) = z f_{j,\pm}(k_\pm, x), \quad z \in \mathbb{C}, \quad j = 1, 2,$$

in the distributional sense. Clearly

$$(3.6) \quad W(f_{j,\mp}(-k_\mp), f_{j,\mp}(k_\mp)) = \mp 2i k_\mp, \quad z > V_\mp, \quad j = 1, 2,$$

where $W(F, G)(x) := F(x)G'(x) - F'(x)G(x)$ denotes the Wronskian of F and G . Spectral properties of H_j are reviewed in

Theorem 3.1. Assume (H.3.1). Then

$$(3.7) \quad \sigma_{\text{ess}}(H_j) = \sigma_{\text{ac}}(H_j) = [\phi_-^2, \infty), \quad \sigma_{\text{sc}}(H_j) = \emptyset, \quad j = 1, 2.$$

Moreover, the H_j , $j = 1, 2$, have simple spectrum in (ϕ_-^2, ϕ_+^2) (iff $\phi_-^2 < \phi_+^2$, otherwise delete this assertion) and spectral multiplicity two in (ϕ_+^2, ∞) . In addition, the H_j , $j = 1, 2$, have finitely many simple eigenvalues $\lambda_{j,l} = \phi_{\pm}^2 - \kappa_{j,\pm,l}^2$ in $[0, \phi_-^2)$ determined by

$$(3.8) \quad W(f_{j,-}(i\kappa_{j,-,l}), f_{j,+}(i\kappa_{j,+,l})) = 0, \quad j = 1, 2.$$

If $\phi_-^2 > 0$, the (simple) eigenvalues of H_1 and H_2 coincide in $(0, \phi_-^2)$. There are no eigenvalues embedded into the essential spectrum and there are no threshold eigenvalues, i.e.,

$$(3.9) \quad \sigma_p(H_j) \cap [\phi_-^2, \infty) = \emptyset, \quad j = 1, 2.$$

If $W(f_{j,-}(0), f_{j,+}(i(\phi_+^2 - \phi_-^2)^{1/2})) = 0$, then H_j has a threshold resonance whose wave function ψ_j (suitably normalized) satisfies

$$(3.10) \quad \begin{aligned} |\psi_j(x) - 1| &\leq c(1 + |x|)^{-1}, & x \leq 0, \\ |\psi_j(x)| &\leq c'e^{-(\phi_+^2 - \phi_-^2)^{1/2}x}, & x \geq 0, \quad j = 1, 2. \end{aligned}$$

Next we describe the unitary on-shell scattering matrix $S_j(\lambda)$ associated with H_j .

Theorem 3.2. Assume (H.3.1). Then if

(i) $\lambda > \phi_+^2$: $S_j(\lambda)$ in \mathbb{C}^2 is given by

$$(3.11) \quad S_j(\lambda) = \begin{pmatrix} T_j(\lambda) & R_j^r(\lambda) \\ R_j^l(\lambda) & T_j(\lambda) \end{pmatrix}, \quad j = 1, 2,$$

where the transmission and reflection coefficients from the left and right are given by

$$(3.12) \quad \begin{aligned} T_j(\lambda) &= 2i(k_+k_-)^{1/2}/W(f_{j,-}(k_-), f_{j,+}(k_+)), \\ R_j^l(\lambda) &= -W(f_{j,-}(-k_-), f_{j,+}(k_+))/W(f_{j,-}(k_-), f_{j,+}(k_+)), \\ R_j^r(\lambda) &= -W(f_{j,-}(k_-), f_{j,+}(-k_+))/W(f_{j,-}(k_-), f_{j,+}(k_+)), \\ k_{\pm} &= (\lambda - \phi_{\pm}^2)^{1/2} > 0, \quad \lambda > \phi_+^2, \quad j = 1, 2. \end{aligned}$$

(ii) $\phi_-^2 < \lambda < \phi_+^2$:

$$(3.13) \quad S_j(\lambda) = -\overline{W(f_{j,-}(k_-), f_{j,+}(k_+))}/W(f_{j,-}(k_-), f_{j,+}(k_+)),$$

$j = 1, 2.$

Unitarity of $S_j(\lambda)$ now trivially follows from (3.6) and the identity

$$(3.14) \quad \begin{aligned} W(F, G)(x)W(f, g)(x) &= W(F, g)(x)W(f, G)(x) \\ &\quad - W(F, f)(x)W(g, G)(x). \end{aligned}$$

In particular,

$$(3.15) \quad \begin{aligned} |T_j(\lambda)|^2 + |R_j^l(\lambda)|^2 &= 1 = |T_j(\lambda)|^2 + |R_j^r(\lambda)|^2, \\ |R_j^l(\lambda)| &= |R_j^r(\lambda)|, \quad \lambda > \phi_+^2, \quad j = 1, 2. \end{aligned}$$

So far we have treated $S_1(\lambda)$ and $S_2(\lambda)$ separately. However, §2 suggests that actually there should be a close connection between them. First of all we infer that

$$(3.16) \quad \begin{cases} f_{1,\pm}(k_\pm, x), \\ f_{2,\pm}(k_\pm, x) = (\pm ik_\pm + \phi_\pm)^{-1} (A f_{1,\pm})(k_\pm, x) \end{cases}$$

are correctly normalized Jost solutions of H_1 and H_2 . Using the elementary identity

$$(3.17) \quad W(Af(z), Ag(z)) = zW(f(z), g(z)), \quad z \in \mathbb{C},$$

where $f(z, x)$, $g(z, x)$ are any distributional solutions of

$$(3.18) \quad (A^* A \psi(z))(x) = (H_1 \psi(z))(x) = z\psi(z, x), \quad z \in \mathbb{C},$$

we obtain

Theorem 3.3. *Assume (H.3.1). Then if*

(i) $\lambda > \phi_+^2$:

$$(3.19) \quad \begin{aligned} T_1(\lambda) &= (ik_- + \phi_-)(ik_+ + \phi_+)^{-1} T_2(\lambda), \\ R_1^l(\lambda) &= (ik_- + \phi_-)(-ik_- + \phi_-)^{-1} R_2^l(\lambda), \\ R_1^r(\lambda) &= (-ik_+ + \phi_+)(ik_+ + \phi_+)^{-1} R_2^r(\lambda). \end{aligned}$$

(ii) $\phi_-^2 < \lambda < \phi_+^2$:

$$(3.20) \quad S_1(\lambda) = (ik_- + \phi_-)(-ik_- + \phi_-)^{-1} S_2(\lambda).$$

Remark 3.4. We also note that the norming constants $c_{j,\pm,l}$, associated with nonzero eigenvalues $\lambda_{j,l} = \phi_\pm^2 - \kappa_{j,\pm,l}^2$ of H_j , $j = 1, 2$, are connected with each other as follows. Let

$$(3.21) \quad |c_{j,\pm,l}| := \|f_{j,\pm,l}\|_2^{-1}, \quad j = 1, 2, \quad l = 1, 2, \dots,$$

where

$$(3.22) \quad \begin{aligned} H_j f_{j,\pm,l} &= \lambda_{j,l} f_{j,\pm,l}, \quad f_{j,\pm,l} \in H^2(\mathbb{R}), \\ f_{j,\pm,l}(ik_{j,\pm,l}, x) &\underset{x \rightarrow \pm\infty}{=} e^{\mp \kappa_{j,\pm,l} x} + o(e^{\mp \kappa_{j,\pm,l} x}), \\ 0 < \kappa_{j,\pm,l} &\leq |\phi_\pm|, \quad j = 1, 2, \\ f_{1,\pm,l}(ik_{1,\pm,l}, x) &= (\phi_\pm \pm \kappa_{1,\pm,l})^{-1} (A^* f_{2,\pm,l})(ik_{1,\pm,l}, x), \\ 0 < \kappa_{1,\pm,l} &< |\phi_\pm|, \quad l = 1, 2, \dots \end{aligned}$$

Then

(3.23)

$$|c_{1,\pm,l}| = \|f_{1,\pm,l}\|_2^{-1} = (\phi_{\pm} \pm \kappa_{1,\pm,l}) \|A^* f_{2,\pm,l}\|_2^{-1} \\ = [(\phi_{\pm} \pm \kappa_{1,\pm,l})(\phi_{\pm} \mp \kappa_{1,\pm,l})^{-1}]^{1/2} |c_{2,\pm,l}|, \quad l = 1, 2, \dots$$

Remark 3.5. As discussed in [49], the hypothesis (H.3.1) in Theorems 3.1–3.3 can be considerably weakened in the sense that for most of the results stated above one only needs the zeroth or first moments of $|\phi(x) - \phi_{\pm}|$, $|\phi'(x)|$ to be finite.

Now we briefly turn to systems governed by (H.3.2). Since now $V_j(x) = \phi^2(x) + (-1)^j \phi'(x)$ are periodic, Floquet theory (see, e.g. [30, 38, 62, 79, 85, 86, 95, 110, 112]) applies. To fix the notation, we briefly introduce normalized Floquet solutions [43, 45]:

$$(3.24) \quad f_{j,\pm}(k, x) := \psi_{j,\pm}(k, x) a^{1/2} / \|\psi_j(k)\|_2,$$

$$(3.25) \quad \|\psi_j(k)\|_2^2 := \int_0^a dx \psi_{j,+}(k, x) \psi_{j,-}(k, x),$$

where the Floquet solutions $\psi_{j,\pm}(k, x)$ are defined as follows. Let

$$(3.26) \quad \begin{aligned} (H_j \theta_j(z))(x) &= z \theta_j(z, x), \\ (H_j \varphi_j(z))(x) &= z \varphi_j(z, x), \quad z \in \mathbb{C}, j = 1, 2, \end{aligned}$$

in the distributional sense with

$$(3.27) \quad \begin{aligned} \theta_j(z, 0) &= 1, \quad \theta'_j(z, 0) = 0, \\ \varphi_j(z, 0) &= 0, \quad \varphi'_j(z, 0) = 1, \quad z \in \mathbb{C}, j = 1, 2, \end{aligned}$$

and define the discriminant F_j by

$$(3.28) \quad F_j(z) := [\theta_j(z, a) + \varphi'_j(z, a)]/2, \quad z \in \mathbb{C}, j = 1, 2.$$

By taking into account that H_1 and H_2 are essentially isospectral, one derives that

$$(3.29) \quad \begin{aligned} \theta_1(z, x), \varphi_1(z, x), \\ \theta_2(z, x) &= A\{z^{-1}\phi(0)\theta_1(z, x) + [1 - z^{-1}\phi(0)^2]\varphi_1(z, x)\}, \\ \varphi_2(z, x) &= A\{-z^{-1}\theta_1(z, x) + z^{-1}\phi(0)\varphi_1(z, x)\}, \quad z \in \mathbb{C} \setminus \{0\}, \end{aligned}$$

satisfy the boundary conditions (3.27). Thus we actually infer

$$(3.30) \quad F_1(z) = F_2(z) := F(z), \quad z \in \mathbb{C}.$$

Now the $\psi_{j,\pm}$ are defined by

$$(3.31) \quad \psi_{j,\pm}(k, x) := \theta_j(z(k), x) + m_{j,\pm}(z(k))\varphi_j(z(k), x), \quad j = 1, 2,$$

where

$$(3.32) \quad m_{j,\pm}(z) := \{\varphi'_j(z, a) - \theta_j(z, a) \mp 2\sqrt{F(z)^2 - 1}\}/2\varphi_j(z, a),$$

$z \in \mathbb{C} \setminus \{\lambda_n \in \mathbb{R} \mid n \in \mathbb{N}\}$, $\sqrt{F(\lambda)^2 - 1} > 0$ for $\lambda < \inf[\sigma(H_j)]$, $j = 1, 2$, with λ_n , $n \in \mathbb{N}$, given by

$$(3.33) \quad F_z(\lambda_n) = 0, \quad n \in \mathbb{N},$$

and the Floquet parameter k defined by

$$(3.34) \quad k(z) := \arcsin[i\sqrt{F(z)^2 - 1}]/a, \quad z \in \mathbb{C}.$$

In addition one has

$$(3.35) \quad f_{j,\pm}(k, x) = e^{\pm ikx} p_{j,\pm}(k, x), \quad p_{j,\pm}(k, x+a) = p_{j,\pm}(k, x), \\ j = 1, 2,$$

$$(3.36) \quad f_{j,\pm}(-k, x) = f_{j,\mp}(k, x) = \overline{f_{j,\pm}(k, x)}, \quad k \in \mathbb{R}, j = 1, 2,$$

$$(3.37) \quad F_z(z) = -\varphi_j(z, a) \parallel \psi_j(k(z)) \parallel_2^2 / 2, \quad z \in \mathbb{C}, j = 1, 2,$$

$$(3.38) \quad W(f_{j,\mp}(-k), f_{j,\mp}(k)) = \pm ai \sin(ka) / F_z(z(k)) \\ = \mp a \sqrt{F(z(k))^2 - 1} / F_z(z(k)), \\ z \in \mathbb{C} \setminus \{\lambda_n \mid n \in \mathbb{N}\}, j = 1, 2.$$

Theorem 3.6. Assume (H.3.2). Then the spectra of H_1 and H_2 coincide. They are purely absolutely continuous and of multiplicity two. Moreover,

$$(3.39) \quad \sigma(H_j) = \sigma_{ac}(H_j) = \{\lambda \in \mathbb{R} \mid |F(\lambda)| \leq 1\} = \bigcup_{n \in \mathbb{N}} \sigma_n, \\ \sigma_n := [E_{2(n-1)}, E_{2n-1}], \quad n \in \mathbb{N}, \\ 0 \leq E_0 < E_1 \leq E_2 < E_3 \leq E_4 < \dots,$$

$$(3.40) \quad \sigma_p(H_j) = \sigma_{sc}(H_j) = \emptyset, \quad j = 1, 2.$$

Remark 3.7. (i) The fact that $F_1(z) = F_2(z)$ is not surprising when taking into account that [95, p. 384]

$$(3.41) \quad \frac{F_j(z) - 2}{F_j(z') - 2} = \det[(h_0 - z')^{-1}(h_j - z)], \\ z \in \mathbb{C}, \quad z' \in \mathbb{C} \setminus \sigma(h_0), \quad j = 1, 2,$$

where

$$(3.42) \quad h_0 = -\frac{d^2}{dx^2}, \quad \mathfrak{D}(h_0) = \{g \in H^2((0, a)) \mid g(0_+) = g(a_-), \quad g'(0_+) = g'(a_-)\}, \\ h_j = -\frac{d^2}{dx^2} + V_j, \quad \mathfrak{D}(h_j) = \mathfrak{D}(h_0), \quad j = 1, 2,$$

are the reduced operators in $L^2((0, a))$ with periodic boundary conditions,

$$(3.43) \quad F_0(z) := 2 \cos(z^{1/2}), \quad z \in \mathbb{C},$$

the corresponding unperturbed discriminant, and $\det(\cdot)$ denotes the Fredholm determinant. Thus

$$(3.44) \quad \frac{F_1(z) - 2}{F_2(z) - 2} = \det[(h_2 - z)^{-1}(h_1 - z)] = 1, \quad z \in \mathbb{C} \setminus \sigma(h_2),$$

by applying commutation to h_1 and h_2 . More generally, applying the Hadamard factorization theorem to $F_1(\cdot)$, one can prove that $F_1(\cdot)$ is an invariant for all periodic potentials (with period $a > 0$) in the isospectral manifold of V_1 (see, e.g., [21]).

(ii) Not only the effective masses (i.e., $d^2k/d\lambda^2|_{\lambda=E_n}$, $n \in \mathbb{N}$) of V_1 and V_2 coincide (cf. [103]) but also the Floquet parameters and the dispersion laws associated with H_1 and H_2 are identical in spite of the fact that $V_1 \neq V_2$ (excluding the trivial case $\phi = \text{const}$).

For further applications of commutation in the periodic case see §9.

We end this section with a simple observation:

Remark 3.8. Let $\phi, \phi' \in L^\infty(\mathbb{R})$ be real-valued, $A = \frac{d}{dx} + \phi$ on $H^1(\mathbb{R})$. Then $H_1 = A^*A$ and $H_2 = AA^*$ cannot simultaneously have zero as an eigenvalue. Indeed, if

$$(3.45) \quad H_1\psi_1 = 0, \quad \psi_1 \in H^2(\mathbb{R}), \quad \psi_1(x) = C \exp \left[- \int_0^x dx' \phi(x') \right],$$

then

$$(3.46) \quad H_2(\psi_1^{-1}) = 0$$

in the distributional sense and the assumption $\psi_1^{-1} \in L^2(\mathbb{R})$ (together with the fact that 0 must be a nondegenerate eigenvalue of H_1) leads to the contradiction

$$(3.47) \quad \infty = \int_{\mathbb{R}} dx \psi_1(x) \psi_1(x)^{-1} \leq \|\psi_1\|_2 \|\psi_1^{-1}\|_2 < \infty.$$

4. ONE-DIMENSIONAL DIRAC OPERATORS

Given the explicit realizations of A in $L^2(\mathbb{R})$ of the foregoing section, we now study the associated Dirac operator

$$(4.1) \quad Q_m = \begin{pmatrix} m & A^* \\ A & -m \end{pmatrix}, \quad m \in \mathbb{R},$$

in $L^2(\mathbb{R}) \otimes \mathbb{C}^2$.

Since Theorems 2.3 and 2.5 directly apply to (4.1), we do not need to repeat them here.

First we consider systems satisfying (H.3.1). Combining Theorem 3.1 and the results of §2 then yields the following theorem.

Theorem 4.1. *Assume (H.3.1). Then*

$$(4.2) \quad \begin{aligned} \sigma_{\text{ess}}(Q_m) &= \sigma_{\text{ac}}(Q_m) = (-\infty, -(\phi_-^2 + m^2)^{1/2}] \cup [(\phi_-^2 + m^2)^{1/2}, \infty), \\ \sigma_{\text{sc}}(Q_m) &= \emptyset, \quad m \in \mathbb{R}, \end{aligned}$$

and Q_m has spectral multiplicity one in $(-(\phi_+^2 + m^2)^{1/2}, -(\phi_-^2 + m^2)^{1/2}) \cup ((\phi_-^2 + m^2)^{1/2}, (\phi_+^2 + m^2)^{1/2})$ (iff $\phi_-^2 < \phi_+^2$, otherwise delete this assertion) and spectral multiplicity two in $(-\infty, -(\phi_+^2 + m^2)^{1/2}) \cup ((\phi_+^2 + m^2)^{1/2}, \infty)$. In addition Q_m has finitely many simple eigenvalues in $(-(\phi_-^2 + m^2)^{1/2}, (\phi_-^2 + m^2)^{1/2})$ (assuming $\phi_- \neq 0$ or $m \neq 0$) symmetrically placed w.r.t. zero with the possible exception at $\pm m$ (if both $\phi_- \neq 0, m \neq 0$). There are no eigenvalues embedded in the essential spectrum and there are no threshold bound states, i.e.,

$$(4.3) \quad \sigma_p(Q_m) \cap \{(-\infty, -(\phi_-^2 + m^2)^{1/2}] \cup [(\phi_-^2 + m^2)^{1/2}, \infty)\} = \emptyset.$$

Next we turn to the on-shell scattering matrix associated with Q_m . In view of Theorem 2.5(i) and (3.5), the Jost solutions of Q_m are given by

$$(4.4) \quad \begin{aligned} \Psi_{1,\pm}(Z, k_{\pm}, x) &= \begin{pmatrix} f_{1,\pm}(k_{\pm}, x) \\ (Z+m)^{-1}(A f_{1,\pm})(k_{\pm}, x) \end{pmatrix}, \\ \Psi_{2,\pm}(Z, k_{\pm}, x) &= \begin{pmatrix} -(Z+m)^{-1}(A^* f_{2,\pm})(k_{\pm}, x) \\ f_{2,\pm}(k_{\pm}, x) \end{pmatrix}, \end{aligned}$$

$$(4.5) \quad \begin{aligned} (Q_m \Psi_{j,\pm})(Z, k_{\pm}, x) &= (-1)^{j+1} Z \Psi_{j,\pm}(Z, k_{\pm}, x), \\ k_{\pm} &= (z - \phi_{\pm})^2, \quad \text{Im } k_{\pm} \geq 0, \quad Z^2 = z + m^2, \quad z \in \mathbb{C} \setminus \{0\}, \quad j = 1, 2. \end{aligned}$$

In analogy to Theorem 3.2, the unitary on-shell scattering matrix $S(E)$ corresponding to Q_m then reads

Theorem 4.2. *Assume (H.3.1). Then if*

(i) $|E| > (\phi_+^2 + m^2)^{1/2}$: $S(E)$ in \mathbb{C}^2 is given by

$$(4.6) \quad S(E) = \begin{pmatrix} T(E) & R'(E) \\ R^l(E) & T(E) \end{pmatrix},$$

where

$$(4.7) \quad \begin{aligned} T(E) &= [2i(k_+ k_-)^{1/2} / (|E| + m)] / W(\Psi_-(E, k_-), \Psi_+(E, k_+)), \\ R^l(E) &= -W(\Psi_-(E, -k_-), \Psi_+(E, k_+)) / W(\Psi_-(E, k_-), \Psi_+(E, k_+)), \\ R'(E) &= -W(\Psi_-(E, k_-), \Psi_+(E, -k_+)) / W(\Psi_-(E, k_-), \Psi_+(E, k_+)), \\ k_{\pm} &= (\lambda - \phi_{\pm}^2)^{1/2}, \quad E^2 = \lambda + m^2 > \phi_+^2 + m^2. \end{aligned}$$

(ii) $(\phi_-^2 + m^2)^{1/2} < |E| < (\phi_+^2 + m^2)^{1/2}$:

$$(4.8) \quad S(E) = -\overline{W(\Psi_-(E, k_-), \Psi_+(E, k_+))} / W(\Psi_-(E, k_-), \Psi_+(E, k_+)).$$

Here

$$(4.9) \quad \Psi_{\pm}(E, k_{\pm}, x) := \begin{cases} \Psi_{1,\pm}(E, k_{\pm}, x), & E > (\phi_-^2 + m^2)^{1/2}, \\ \Psi_{2,\pm}(-E, k_{\pm}, x), & E < -(\phi_-^2 + m^2)^{1/2}, \end{cases}$$

and $W(\Phi, \Psi)$ denotes the determinant of $\Phi, \Psi \in \mathbb{C}^2$.

Taking into account the identity

$$(4.10) \quad W(\Psi_-(E, \sigma k_-), \Psi_+(E, \sigma' k_+)) = W(f_{1,-}(\sigma k_-), f_{1,+}(\sigma' k_+)) / (|E| + m), \\ \sigma, \sigma' = \pm 1, E > (\phi_-^2 + m^2)^{1/2},$$

and similarly for $E < -(\phi_-^2 + m^2)^{1/2}$, we finally obtain

Theorem 4.3. Assume (H.3.1). Then

$$(4.11) \quad S(E) = \begin{cases} S_1(E^2 - m^2), & E > (\phi_-^2 + m^2)^{1/2}, \\ S_2(E^2 - m^2), & E < -(\phi_-^2 + m^2)^{1/2}. \end{cases}$$

Next we briefly discuss systems satisfying (H.3.2). Theorems 2.5(i) and 3.6 imply

Theorem 4.4. Assume (H.3.2). Then the spectrum of Q_m is purely absolutely continuous, symmetric with respect to the origin, and of multiplicity two. Moreover,

$$(4.12) \quad \sigma(Q_m) = \bigcup_{n \in \mathbb{Z} \setminus \{0\}} \Sigma_n, \\ \Sigma_n = [|E_{2(n-1)} + m^2|^{1/2}, |E_{2n-1} + m^2|^{1/2}], \quad \Sigma_{-n} = -\Sigma_n, \quad n \in \mathbb{N},$$

with $E_n, n \in \mathbb{N}_0$, given by (3.39), and

$$(4.13) \quad \sigma_p(Q_m) = \sigma_{sc}(Q_m) = \emptyset, \quad m \in \mathbb{R}.$$

We also state

Lemma 4.5. Assume (H.3.2). Let $F_{Q_m}(Z), Z \in \mathbb{C}$, be the discriminant of Q_m . Then (cf. (3.30))

$$(4.14) \quad F_{Q_m}(Z) = F(Z^2 - m^2), \quad Z \in \mathbb{C}.$$

Proof. Similar to (4.4) and (4.9) we get a fundamental matrix Φ_m for Q_m :

$$(4.15) \quad \Phi_m(E, x) = \begin{pmatrix} \theta_1(\lambda, x) - \phi(0)\varphi_1(\lambda, x) & (E + m)\varphi_1(\lambda, x) \\ (E + m)^{-1}[A(\theta_1 - \phi(0)\varphi_1)](\lambda, x) & (A\varphi_1)(\lambda, x) \end{pmatrix}, \\ E > m, \\ \Phi_m(E, x) = \begin{pmatrix} -(A^*\varphi_2)(\lambda, x) & -(E + m)^{-1}[A^*(\theta_2 + \phi(0)\varphi_2)](\lambda, x) \\ (E + m)\varphi_2(\lambda, x) & \theta_2(\lambda, x) + \phi(0)\varphi_2(\lambda, x) \end{pmatrix}, \\ E < -m, \lambda = E^2 - m^2,$$

with $\theta_j, \varphi_j, j = 1, 2$, defined in (3.27) and

$$(4.16) \quad \Phi_m(E, 0) = 1, \quad E \in \mathbb{R} \setminus [-m, m].$$

Using (3.30), we obtain

$$(4.17) \quad F_{Q_m}(E) := 2^{-1} \operatorname{Tr} [\Phi_m(E, a)] = F(\lambda), \quad E \in \mathbb{R} \setminus [-m, m].$$

Analytic continuation w.r.t. E of both sides in (4.17) then yields (4.14). \square

For further results in the periodic case see Theorem 9.3.

5. RELATIVE SCATTERING

In this section we consider relative scattering, i.e., scattering for the pair $(H_j + W, H_j)$ with $H_j, j = 1, 2$, satisfying (H.3.1) or (H.3.2) and W satisfying hypothesis

$$(H.5.1) \quad W \in L^\infty(\mathbb{R}) \cap L^1(\mathbb{R}; (1+x^2)dx) \text{ real-valued.}$$

By (H.5.1), W is a relatively compact perturbation of H_j and hence

$$(5.1) \quad \tilde{H}_j = H_j + W, \quad \mathfrak{D}(\tilde{H}_j) = \mathfrak{D}(H_j), \quad j = 1, 2,$$

is well defined. Let $\tilde{f}_{j,\pm}(p_\pm, x)$ denote the perturbed Jost solutions of \tilde{H}_j defined by

$$(5.2) \quad \begin{aligned} \tilde{f}_{j,\pm}(p_\pm, x) &= f_{j,\pm}(p_\pm, x) \\ &\quad + \int_x^{\pm\infty} dx' g_j(p_\pm, x, x') W(x') \tilde{f}_{j,\pm}(p_\pm, x'), \quad j = 1, 2, \end{aligned}$$

$$(5.3) \quad \begin{aligned} g_j(p_\pm, x, x') &:= [W(f_{j,-}(p_-), f_{j,+}(p_+))]^{-1} \\ &\quad \cdot [f_{j,+}(p_+, x') f_{j,-}(p_-, x) \\ &\quad - f_{j,+}(p_+, x) f_{j,-}(p_-, x')], \quad j = 1, 2, \end{aligned}$$

$$(5.4) \quad p_\pm := \begin{cases} k_\pm & \text{in case (H.3.1),} \\ k & \text{in case (H.3.2).} \end{cases}$$

We also define

$$(5.5) \quad \begin{aligned} \mp W(p_\mp) &:= W(\tilde{f}_{j,\mp}(-p_\mp), \tilde{f}_{j,\mp}(p_\mp)) \\ &= W(f_{j,\mp}(-p_\mp), f_{j,\mp}(p_\mp)) \\ &= \begin{cases} \mp 2ik_\mp & \text{in case (H.3.1),} \\ \pm i \sin(ka)/F_z(z(k)) & \text{in case (H.3.2).} \end{cases} \end{aligned}$$

Concerning spectral properties of \tilde{H}_j we recall

Theorem 5.1. Assume (H.5.1) and either (H.3.1) or (H.3.2). Then

$$(5.6) \quad \sigma_{\text{ess}}(\tilde{H}_j) = \sigma_{\text{ac}}(\tilde{H}_j) = \sigma_{\text{ess}}(H_j), \quad \sigma_{\text{sc}}(\tilde{H}_j) = \emptyset, \quad j = 1, 2,$$

and the multiplicity of $\sigma_{\text{ess}}(\tilde{H}_j)$ coincides with that of $\sigma_{\text{ess}}(H_j)$, $j = 1, 2$. In addition, the (necessarily simple) eigenvalues of \tilde{H}_j are determined by

$$(5.7) \quad W(\tilde{f}_{j,-}(p_{j,-,l}), \tilde{f}_{j,+}(p_{j,+,l})) = 0, \quad j = 1, 2, \quad l = 1, 2, \dots$$

Moreover, there are no eigenvalues embedded in the essential spectrum and there are no threshold eigenvalues, i.e.,

$$(5.8) \quad \sigma_p(\tilde{H}_j) \cap \sigma_{\text{ess}}(H_j) = \emptyset, \quad j = 1, 2.$$

Next we turn to the relative on-shell scattering matrix $S_{j,\text{rel}}(\lambda)$ w.r.t. the pair (\tilde{H}_j, H_j) , $j = 1, 2$.

Theorem 5.2. Assume (H.3.1) and (H.5.1). Let $\tilde{S}_j(\lambda)$, $\lambda > \phi_-^2$, denote the on-shell scattering matrix associated with \tilde{H}_j . Then $\tilde{S}_j(\lambda)$ is given by (3.11)–(3.13) with $f_{j,\pm}(\sigma k_{\pm})$ replaced by $\tilde{f}_{j,\pm}(\sigma k_{\pm})$, $\sigma = \pm 1$, $j = 1, 2$, and

$$(5.9) \quad S_{j,\text{rel}}(\lambda) = \tilde{S}_j(\lambda) S_j(\lambda)^{-1}, \quad \lambda > \phi_-^2, \quad j = 1, 2.$$

Proof. For simplicity we assume that $\phi_+ = \phi_- = 0$. To fix notations, let A, B, C be selfadjoint operators in a separable, complex Hilbert space \mathfrak{H} and suppose $(B - z_0)^{-1} - (A - z_0)^{-1} \in \mathfrak{B}_1(\mathfrak{H})$, $(C - z_0)^{-1} - (A - z_0)^{-1} \in \mathfrak{B}_1(\mathfrak{H})$ for some $z_0 \in \mathbb{C} \setminus \mathbb{R}$. (Here $\mathfrak{B}_1(\cdot)$ denotes the set of trace class operators.) Then the wave operators

$$(5.10) \quad \Omega_{\pm}(B, A) := \text{s-}\lim_{t \rightarrow \pm\infty} e^{itB} e^{-itA} P_{\text{ac}}(A)$$

exist ($P_{\text{ac}}(A)$ is the projection onto the absolutely continuous subspace of A) and the scattering operator

$$(5.11) \quad S(B, A) := \Omega_+(B, A)^* \Omega_-(B, A)$$

is unitary in $P_{\text{ac}}(A)\mathfrak{H}$. In view of the applications below, assume $L^2((0, \infty); \mathbb{C}^2)$ to be the spectral representation and $U_A^{\pm} : P_{\text{ac}}(A)\mathfrak{H} \rightarrow L^2((0, \infty); \mathbb{C}^2)$ (constructed from the generalized eigenfunctions of A) the spectral transformations of $AP_{\text{ac}}(A)$, i.e.,

$$(5.12) \quad (U_A^{\pm} AP_{\text{ac}}(A)g)(k) = k^2 (U_A^{\pm} g)(k), \quad k \in (0, \infty), \quad g \in \mathfrak{D}(AP_{\text{ac}}(A)).$$

Since $S(B, A)$ commutes with $AP_{\text{ac}}(A)$, we get for a.e. $k \in (0, \infty)$,

$$(5.13) \quad (U_A^{\pm} S(B, A)(U_A^{\pm})^* f)(k) = S(B, A, k) f(k), \quad f(k) \in \mathbb{C}^2,$$

where $S(B, A, k)$, $k > 0$, is unitary in \mathbb{C}^2 . Using the chain rule

$$(5.14) \quad \Omega_{\pm}(C, B) \Omega_{\pm}(B, A) = \Omega_{\pm}(C, A)$$

and

$$(5.15) \quad \Omega_{\pm}(B, A)^* \Omega_{\pm}(B, A) = P_{\text{ac}}(A), \quad \Omega_{\pm}(B, A) = \Omega_{\pm}(A, B)^*,$$

one obtains

$$(5.16) \quad S(C, B) = \Omega_+(B, A)S(C, A)S(B, A)^* \Omega_+(B, A)^*.$$

Finally, taking into account

$$(5.17) \quad \Omega_{\pm}(B, A) = (U_B^{\pm})^* U_A^{\pm}, \quad (U_B^{\pm})^* U_B^{\pm} = P_{\text{ac}}(B), \quad U_B^{\pm} (U_B^{\pm})^* = 1,$$

equation (5.16) implies

$$(5.18) \quad U_B^+ S(C, B) (U_B^+)^* = [U_A^+ S(C, A) (U_A^+)^*] [U_A^+ S(B, A) (U_A^+)^*]^*.$$

Together with (5.13) this yields (5.9) by identifying $C = \tilde{H}_j$, $B = H_j$, $A = -d^2/dx^2$, $\mathfrak{D}(A) = H^2(\mathbb{R})$. The general case $\phi_+ \neq 0$, $\phi_- \neq 0$ can be treated in a similar manner by introducing two comparison dynamics A_{\pm} as $t \rightarrow \pm\infty$ in (5.10) (see, e.g., [96]). \square

Similarly one gets (see also [44, 46])

Theorem 5.3. *Assume (H.3.2) and (H.5.1). Then*

$$(5.19) \quad S_{j, \text{rel}}(\lambda) = \begin{pmatrix} T_{j, \text{rel}}(\lambda) & R_{j, \text{rel}}^r(\lambda) \\ R_{j, \text{rel}}^l(\lambda) & T_{j, \text{rel}}(\lambda) \end{pmatrix}, \quad \lambda \in \sigma(H_j)^0, j = 1, 2,$$

where

$$(5.20) \quad \begin{aligned} T_{j, \text{rel}}(\lambda) &= W(k)/W(\tilde{f}_{j, -}(k), \tilde{f}_{j, +}(k)), \\ R_{j, \text{rel}}^l(\lambda) &= -W(\tilde{f}_{j, -}(-k), \tilde{f}_{j, +}(k))/W(\tilde{f}_{j, -}(k), \tilde{f}_{j, +}(k)), \\ R_{j, \text{rel}}^r(\lambda) &= -W(\tilde{f}_{j, -}(k), \tilde{f}_{j, +}(-k))/W(\tilde{f}_{j, -}(k), \tilde{f}_{j, +}(k)), \\ &\lambda \in \sigma(H_j)^{\circ}, j = 1, 2. \end{aligned}$$

(Here M° denotes the interior of a subset $M \subset \mathbb{R}$.) I.e., formally, $S_j(\lambda) = 1$ in the case of (H.3.2).

Relative scattering theory in connection with the Korteweg-deVries equation has also been studied in [89].

6. ADDING EIGENVALUES

As will become clear in §10, the notion of an N -soliton solution of the KdV-equation relative to a nontrivial background KdV-solution is intimately connected with the process of adding N -eigenvalues, say $\{E_1, \dots, E_N\}$ to the (background) Hamiltonian $H_j = -d^2/dx^2 + V_j$ in (3.2), where $\sigma(H_j) \cap \{E_1, \dots, E_N\} = \emptyset$. Moreover, if \tilde{H}_j denotes the new Schrödinger operator obtained by this process, then one needs the remaining spectral properties of H_j and \tilde{H}_j to coincide in the sense that

$$(6.1) \quad \sigma(\tilde{H}_j) \setminus \{E_1, \dots, E_N\} = \sigma(H_j), \quad \sigma_{\text{ess}}(\tilde{H}_j) = \sigma_{\text{ess}}(H_j)$$

and similarly for the absolutely continuous spectrum and the spectral multiplicity respectively.

In this section we shall consider two cases relevant in §§8 and 10: First the case where $\phi_+^2 = \phi_-^2$ in (H.3.1) and then the periodic case governed by (H.3.2). (Throughout this section j will be fixed to be either 1 or 2.)

Combining §3 with Theorems 3.6 and 3.7 of [33] one gets

Theorem 6.1. Assume (H.3.1) with $0 \leq \phi_-^2 = \phi_+^2 =: \phi_\infty^2$. Let $\lambda_{j,1} < \lambda_{j,2} < \dots < \lambda_{j,N_j}$, $\lambda_{j,l} = \phi_\infty^2 - \kappa_{j,l}^2$, $j = 1, 2$, $l = 1, \dots, N_j$, denote the eigenvalues of H_j and let $\tilde{\kappa}_{j,1} < \tilde{\kappa}_{j,2} < \dots < \tilde{\kappa}_{j,\tilde{N}_j}$, $\alpha_{j,1}, \alpha_{j,2}, \dots, \alpha_{j,\tilde{N}_j}$ be arbitrary positive numbers with $\kappa_{j,l} \neq \tilde{\kappa}_{j,m}$, $l = 1, \dots, N_j$, $m = 1, \dots, \tilde{N}_j$. Then

$$(6.2) \quad \begin{aligned} \tilde{V}_j(x) := & V_j(x) - 2 \frac{d^2}{dx^2} \ln W(f_{j,+}(i\kappa_{j,1}), \dots, f_{j,+}(i\kappa_{j,N_j})) \\ & - 2 \frac{d^2}{dx^2} \ln W(a_{j,1}, \dots, a_{j,\tilde{N}_j}) \end{aligned}$$

(where W denotes the Wronskian determinant [33]), with

$$(6.3) \quad a_{j,m} = \frac{W(f_{j,+}(i\kappa_{j,1}), \dots, f_{j,+}(i\kappa_{j,N_j}), (-1)^{m+1} f_{j,+}(i\tilde{\kappa}_{j,m}) + \alpha_{j,m} (-1)^{N_j} f_{j,-}(i\tilde{\kappa}_{j,m}))}{\prod_{l=1}^{N_j} (\kappa_{j,l} - \tilde{\kappa}_{j,m}) W(f_{j,+}(i\kappa_{j,1}), \dots, f_{j,+}(i\kappa_{j,N_j}))},$$

$$m = 1, \dots, \tilde{N}_j,$$

supports \tilde{N}_j bound states $\tilde{\lambda}_{j,\tilde{N}_j} < \dots < \tilde{\lambda}_{j,1}$, $\tilde{\lambda}_{j,m} = \phi_\infty^2 - \tilde{\kappa}_{j,m}^2$, $m = 1, \dots, \tilde{N}_j$, with norming constants

$$(6.4) \quad \begin{aligned} \tilde{c}_{j,+}^2 = & \frac{2(-1)^{m-1} \tilde{\kappa}_{j,m}}{\alpha_{j,m}} \prod_{n \neq m} \left(\frac{\tilde{\kappa}_{j,n} + \tilde{\kappa}_{j,m}}{\tilde{\kappa}_{j,n} - \tilde{\kappa}_{j,m}} \right) \\ & \cdot \prod_{l=1}^{N_j} \left(\frac{\tilde{\kappa}_{j,m} - \kappa_{j,l}}{\tilde{\kappa}_{j,m} + \kappa_{j,l}} \right) T_j(\tilde{\lambda}_{j,m}), \quad m = 1, \dots, \tilde{N}_j, \end{aligned}$$

transmission coefficient

$$(6.5) \quad \tilde{T}_j(\lambda) = \prod_{l=1}^{N_j} \left(\frac{k_-(\lambda) - i\kappa_{j,l}}{k_-(\lambda) + i\kappa_{j,l}} \right) \prod_{m=1}^{\tilde{N}_j} \left(\frac{k_-(\lambda) + i\tilde{\kappa}_{j,m}}{k_-(\lambda) - i\tilde{\kappa}_{j,m}} \right) T_j(\lambda),$$

and reflection coefficients

$$(6.6) \quad \tilde{R}_j^{l,r}(\lambda) = (-1)^{\tilde{N}_j - N_j} \prod_{l=1}^{N_j} \left(\frac{k_-(\lambda) - i\kappa_{j,l}}{k_-(\lambda) + i\kappa_{j,l}} \right) \prod_{m=1}^{\tilde{N}_j} \left(\frac{k_-(\lambda) + i\tilde{\kappa}_{j,m}}{k_-(\lambda) - i\tilde{\kappa}_{j,m}} \right) R_j^{l,r}(\lambda).$$

Remark 6.2. If $\kappa_{j,L} = \tilde{\kappa}_{j,M}$ for some L, M , $1 \leq L \leq N_j$, $1 \leq M \leq \tilde{N}_j$, (6.3) must be replaced by

$$(6.7) \quad a_{jM} = \frac{W(f_{j,+}(i\kappa_{j,L+1}; -L), \dots, f_{j,+}(i\kappa_{j,N_j}; -L), (-1)^{L+1} f_{j,+}(i\kappa_{j,L}; -L) + \alpha_{j,M} (-1)^{N_j-L} f_{j,-}(i\kappa_{j,L}; -L))}{\prod_{l=L+1}^{N_j} (\kappa_{j,l} - \kappa_{j,L}) W(f_{j,+}(i\kappa_{j,L+1}; -L), \dots, f_{j,+}(i\kappa_{j,N_j}; -L))}.$$

Here $f_{j,\pm}(k_-; -L)$ are the Jost solutions after having removed the lowest lying L bound states from V_j , i.e.,

(6.8)

$$f_{j,\pm}(k_-; -L) = (\pm 1)^L \frac{W(f_{j,+}(i\kappa_{j,1}), \dots, f_{j,+}(i\kappa_{j,L}); f_{j\pm}(k_-))}{\prod_{l=1}^L (ik_- + \kappa_{j,l}) W(f_{j,+}(i\kappa_{j,1}), \dots, f_{j,+}(i\kappa_{j,L}))},$$

$k_- \neq \kappa_{j,L},$

and

(6.9)

$$f_{j,\pm}(i\kappa_{j,L}; -L) = \mp \frac{c_{j,\mp,L}^2(-(L-1))}{T_j(\lambda_{j,L}; -L)} f_{j,\mp}(i\kappa_{j,L}; -(L-1))^{-1} \cdot \int_{\pm\infty}^x f_{j,\mp}^2(i\kappa_{j,L}; -(L-1)) dy,$$

with

(6.10)

$$T_j(\lambda_{j,L}; -L) = \lim_{k_- \rightarrow i\kappa_{j,L}} \prod_{l=1}^L \left(\frac{k_- - i\kappa_{j,l}}{k_- + i\kappa_{j,l}} \right) T_j(\lambda)$$

and

(6.11)

$$c_{j,+,L}^2(-(L-1)) = (-1)^{(L-1)} \prod_{l=1}^{L-1} \left(\frac{\kappa_{j,l} - \kappa_{j,L}}{\kappa_{j,l} + \kappa_{j,L}} \right) c_{j,+,L}^2,$$

(6.12)

$$c_{j,-,L}^2(-(L-1)) = \frac{4\kappa_{j,L}^2 T_j^2(\lambda_{j,L}; -L)}{c_{j,+,L}^2(-(L-1))}.$$

Remark 6.3. If V_j has no bound states, i.e., $N_j = 0$, (6.2) reduces to the formula given by [33]

(6.13)

$$\tilde{V}_j(x) = V_j(x) - 2 \frac{d^2}{dx^2} \ln W(a_{j,1}, \dots, a_{j,\tilde{N}_j})$$

with

(6.14)

$$a_{j,m} = (-1)^{m+1} f_{j,+}(i\tilde{\kappa}_{j,m}) + \alpha_{j,m} f_{j,-}(i\tilde{\kappa}_{j,m}), \quad m = 1, \dots, \tilde{N}_j.$$

For $V_j \equiv 0$ (6.13) is equivalent to the \tilde{N}_j -soliton formula obtained from the Marchenko equation, i.e.,

(6.15)

$$\tilde{V}_j(x) = -2 \frac{d^2}{dx^2} \ln \{ \det[1 + C_{j,\pm}(x)] \}$$

(see [48, 57, 71, 80, 106, 115]), where

(6.16)

$$C_{j,\pm}(x) = [c_{j,\pm,l,m}(x)]_{l,m=1}^{\tilde{N}_j},$$

$$c_{j,\pm,l,m}(x) = (\tilde{\kappa}_{j,l} + \tilde{\kappa}_{j,m})^{-1} \tilde{c}_{j,\pm,l} \tilde{c}_{j,\pm,m} e^{\mp(\tilde{\kappa}_{j,l} + \tilde{\kappa}_{j,m})x},$$

$$l, m = 1, \dots, \tilde{N}_j.$$

The norming constants $\tilde{c}_{j,\pm,m}$ are then related to $\alpha_{j,m}$ via

$$\begin{aligned}
 \alpha_{j,m} &= \frac{2(-1)^{m-1}\tilde{\kappa}_{j,m}}{\tilde{c}_{j,+,m}^2} \prod_{n \neq m}^{\tilde{N}_j} \left(\frac{\tilde{\kappa}_{j,n} + \tilde{\kappa}_{j,m}}{\tilde{\kappa}_{j,n} - \tilde{\kappa}_{j,m}} \right) \\
 (6.17) \quad &= \frac{(-1)^{m-1}\tilde{c}_{j,-,m}^2}{2\tilde{\kappa}_{j,m}} \prod_{n \neq m}^{\tilde{N}_j} \left(\frac{\tilde{\kappa}_{j,n} - \tilde{\kappa}_{j,m}}{\tilde{\kappa}_{j,n} + \tilde{\kappa}_{j,m}} \right), \quad m = 1, \dots, \tilde{N}_j.
 \end{aligned}$$

For the connection between the Marchenko procedure and commutation methods see, e.g., [24].

At the end of this section we briefly consider reflectionless scattering relative to a finite-zone periodic background potential V_j fulfilling (H.3.2). Thereby we are looking for potentials $W_{j,\pm}$ giving rise to \tilde{N}_j bound states at energies $\tilde{\lambda}_{j,l} \notin \sigma(H_j)$, $l = 1, \dots, \tilde{N}_j$,

$$(6.18) \quad \tilde{\lambda}_{j,1} < \tilde{\lambda}_{j,2} < \dots < \tilde{\lambda}_{j,\tilde{N}_j}, \quad \tilde{N}_j \in \mathbb{N},$$

and zero relative reflection coefficients $R_{j,\text{rel}}^l(\lambda) = R_{j,\text{rel}}^r(\lambda) = 0$, $\lambda \in \sigma(H_j)$. Assuming the unperturbed potential V_j to provide finitely many energy bands σ_k , $k = 1, \dots, N$, the set of bound states $\{\tilde{\lambda}_{j,l} | l = 1, \dots, \tilde{N}_j\}$ decomposes into subsets $\{\tilde{\lambda}_{j,l} | l = \tilde{N}_{j,k-2}+1, \dots, \tilde{N}_{j,k-1}\}$, ($\tilde{N}_{j,-1} = 0$) of $\tilde{N}_{j,k-1}$ bound states respectively, sitting in the $(k-1)$ th gap ($k = 1, \dots, N$). Applying the Marchenko method [80] we have to solve the integral equation

$$\begin{aligned}
 K_{j,\pm}(x, x') + \Omega_{j,\pm}(x, x') \\
 (6.19) \quad \pm \int_x^{\pm\infty} dy K_{j,\pm}(x, y) \Omega_{j,\pm}(y, x') = 0, \quad x \begin{matrix} < \\ > \end{matrix} x',
 \end{aligned}$$

with kernel

$$(6.20) \quad \Omega_{j,\pm}(x, x') = \sum_{l=1}^{\tilde{N}_j} \tilde{c}_{j,\pm,l}^2 f_{j,\pm}(k_{j,l}, x) f_{j,\pm}(k_{j,l}, x'),$$

where $f_{j,\pm}$ are normalized Floquet solutions belonging to H_j (see (3.24)), the quasi-momenta $k_{j,l}$ are related to $\tilde{\lambda}_{j,l}$ via (3.34), and $\tilde{c}_{j,\pm,l}$ are the norming constants corresponding to the l th bound state $\tilde{\lambda}_{j,l}$. For a separable kernel like (6.20) an ansatz of the form

$$(6.21) \quad K_{j,\pm}(x, x') = \sum_{l=1}^{\tilde{N}_j} \tilde{c}_{j,\pm,l} \eta_{j,\pm,l}(x) f_{j,\pm}(k_{j,l}, x')$$

reduces (6.19) to a set of algebraic equations

$$(6.22) \quad \sum_{m=1}^{\tilde{N}_j} [1 + C_{j,\pm}(x)]_{l,m} \eta_{j,\pm,m}(x) = -\tilde{c}_{j,\pm,l} f_{j,\pm}(k_{j,l}, x),$$

with

$$(6.23) \quad [C_{j,\pm}(x)]_{l,m} = \pm \tilde{c}_{j,\pm,l} \tilde{c}_{j,\pm,m} \int_x^{\pm\infty} dy f_{j,\pm}(k_{j,l}, y) f_{j,\pm}(k_{j,m}, y).$$

Since for $(a_1, \dots, a_N) \in \mathbb{R}^N \setminus \{0\}$

$$(6.24) \quad \sum_{l,m=1}^N a_l [C_{j,\pm}(x)]_{lm} a_m = \pm \int_x^{\pm\infty} dy \left[\sum_{l=1}^N a_l \tilde{c}_{j,\pm,l} f_{j,\pm}(k_{j,l}, y) \right]^2 > 0,$$

$C_{j,\pm}(x)$ is positive definite. Therefore (6.22) can be inverted to give

$$(6.25) \quad K_{j,\pm}(x, x') = - \sum_{l,m=1}^{\tilde{N}_j} [1 + C_{j,\pm}(x)]_{l,m}^{-1} \tilde{c}_{j,\pm,l} \tilde{c}_{j,\pm,m} \cdot f_{j,\pm}(k_{j,l}, x) f_{j,\pm}(k_{j,m}, x').$$

The potentials $W_{j,\pm}$ are now recovered from $K_{j,\pm}$ via

$$(6.26) \quad W_{j,\pm}(x) = \mp 2 \frac{d}{dx} K_{j,\pm}(x, x), \quad x \in \mathbb{R}.$$

Inserting (6.25) into (6.26) we finally obtain

$$(6.27) \quad W_{j,\pm}(x) = -2 \frac{d^2}{dx^2} \ln \{ \det[1 + C_{j,\pm}(x)] \}, \quad x \in \mathbb{R},$$

for potentials relative reflectionless w.r.t. a finite-zone periodic background. Note that in general $W_{j,+} \neq W_{j,-}$ (cf. §10).

7. COMMUTATION, MIURA'S TRANSFORMATION AND SOLUTIONS OF THE mKdV-EQUATION

Using commutation methods, we now take a closer look at the Miura transformation which links solutions of the mKdV- and KdV-equations. In particular, we provide an explicit construction of mKdV-solutions given a corresponding KdV-solution.

We introduce hypothesis

$$(H.7.1) \quad V, \phi \in C^\infty(\mathbb{R}^2) \text{ real-valued, } \partial_x^n V, \partial_x^n \phi \in L^\infty(\mathbb{R}^2), \quad n = 0, 1$$

Given (H.7.1), we are interested in studying real-valued solutions of the KdV-equation

$$(7.1) \quad \text{KdV}(V) := V_t - 6V V_x + V_{xxx} = 0$$

and of the mKdV-equation

$$(7.2) \quad \text{mKdV}(\phi) := \phi_t - 6\phi^2 \phi_x + \phi_{xxx} = 0$$

for $(t, x) \in \mathbb{R}^2$.

Remark 7.1. In the following we simply assume the existence of smooth solutions of (7.1) and (7.2). The corresponding initial value problems have been studied, e.g., in [12, 25, 26, 51, 66, 68, 69, 74, 80, 87, 101, 111] (and references therein).

Introducing the Miura transformation [88]

$$(7.3) \quad V_j(t, x) = \phi(t, x)^2 + (-1)^j \phi_x(t, x), \quad (t, x) \in \mathbb{R}^2, \quad j = 1, 2,$$

a simple calculation shows that

$$(7.4) \quad \text{KdV}(V_j) = [2\phi + (-1)^j \partial_x] \text{mKdV}(\phi), \quad j = 1, 2.$$

Remark 7.2.

(i) Note that (up to smoothness properties) $V_j(t, x)$ for fixed $t \in \mathbb{R}$ corresponds to the potential V_j of §3 under hypothesis (H.3.1), resp. (H.3.2), i.e., $V_1(t, x)$ and $V_2(t, x)$ are related to each other by commutation! This fact seems to have been exploited first by Deift [32] and Adler and Moser [6]. In particular, in [6] it is shown that Miura's transformation (7.3) induces a relation analogous to (7.4) between all higher order KdV- and mKdV-equations (see the end of this section). Further references exploiting commutation are, e.g., [7, 12, 13, 42, 65, 75, 82, 83, 102, 118].

(ii) In contrast to KdV, the transformation $\phi \rightarrow -\phi$ yields again a solution of the mKdV-equation (7.2). This transformation corresponds to interchanging V_1 and V_2 (cf. (7.3)) in the KdV-equation (7.1).

(iii) It has been shown in [3] that rapidly decaying solutions $V(t, x)$ as $|x| \rightarrow \infty$ of (7.1) in general cannot be obtained from rapidly decaying solutions ϕ of (7.2) via Miura's transformation (7.3). This explains the necessity of studying solutions ϕ of (7.2) with nontrivial spatial asymptotics (see also [48]).

(iv) It is obvious from (7.4) that, whenever $\phi \not\equiv 0$ satisfies (H.7.1) and (7.2), both V_1 and V_2 satisfy (7.1) and (H.7.1) and at least one of V_1 or V_2 is not identically zero. On the other hand, if V_1 or V_2 solve (7.1), we cannot trivially reverse that procedure and construct a solution ϕ of (7.2) since $2\phi \mp \partial_x$ on $C^\infty(\mathbb{R}^2)$ has the nontrivial kernel

$$(7.5) \quad f(t) \exp \left[\pm 2 \int_{x_0}^x dy \phi(t, y) \right], \quad f \in C^\infty(\mathbb{R}), \quad x_0 \in \mathbb{R}.$$

One of the striking features of $V(t, x)$, resp. $\phi(t, x)$, satisfying (7.1), resp. (7.2), is that they induce isospectral deformations of the potential in the Schrödinger, resp. Dirac, operator [77]:

Theorem 7.3. (i) Assume V satisfies (H.7.1) and $V_t(t, \cdot) \in L^\infty(\mathbb{R})$, $t \in \mathbb{R}$. If in addition $V(t, x)$ solves (7.1), then the Schrödinger operator $H(t)$ in $L^2(\mathbb{R})$,

$$(7.6) \quad H(t) = -\partial_x^2 + V(t, \cdot), \quad \mathfrak{D}(H(t)) = H^2(\mathbb{R}), \quad t \in \mathbb{R},$$

is unitarily equivalent to $H(0)$ for all $t \in \mathbb{R}$, i.e., there exists a family of unitary operators $U(t)$, $t \in \mathbb{R}$, $U(0) = 1$ in $L^2(\mathbb{R})$ such that

$$(7.7) \quad U(t)^{-1} H(t) U(t) = H(0), \quad t \in \mathbb{R}.$$

(ii) Assume ϕ satisfies (H.7.1), $\phi_{xx} \in L^\infty(\mathbb{R}^2)$, and $\phi_t(t, \cdot), \phi_{tx}(t, \cdot) \in L^\infty(\mathbb{R})$, $t \in \mathbb{R}$. If in addition $\phi(t, x)$ solves (7.2), then the Dirac operator $Q_m(t)$ in $L^2(\mathbb{R}) \otimes \mathbb{C}^2$,

$$(7.8) \quad Q_m(t) = \begin{pmatrix} m & A(t)^* \\ A(t) & -m \end{pmatrix},$$

$$\mathfrak{D}(Q_m(t)) = H^1(\mathbb{R}) \otimes \mathbb{C}^2, \quad t \in \mathbb{R}, \quad m \in \mathbb{R},$$

$$(7.9) \quad A(t) = \partial_x + \phi(t, \cdot), \quad \mathfrak{D}(A(t)) = H^1(\mathbb{R}), \quad t \in \mathbb{R},$$

is unitarily equivalent to $Q_m(0)$ for all $t \in \mathbb{R}$, i.e., there exists a family of unitary operators $W_m(t)$, $t \in \mathbb{R}$, $W_m(0) = 1$ in $L^2(\mathbb{R}) \otimes \mathbb{C}^2$ such that

$$(7.10) \quad W_m(t)^{-1} Q_m(t) W_m(t) = Q_m(0), \quad (m, t) \in \mathbb{R}^2.$$

Proof. The Lax pair

$$(7.11) \quad L(t) = H(t),$$

$$B(t) = -4\partial_x^3 + 6V(t, \cdot)\partial_x + 3V_x(t, \cdot), \quad \mathfrak{D}(B(t)) = H^3(\mathbb{R}), \quad t \in \mathbb{R},$$

together with

$$(7.12) \quad \partial_t L(t) = [B(t), L(t)], \quad \partial_t U(t) = B(t)U(t), \quad U(0) = 1$$

is well known to prove assertion (7.7) (cf., e.g., [39, 77, 78, 80, 93], and the references therein). In order to prove (7.10), one could similarly construct a Lax pair (see, e.g., [107]), however, because of (7.4) and Theorem 2.5(i), (7.10) follows from (7.7). \square

Remark 7.2(iv) (based on (7.4)) shows how to obtain solutions of the KdV-equation (7.1) given solutions of the mKdV-equation (7.2). In the following we concentrate on the reversed procedure, i.e., how to construct solutions V_2 , resp. ϕ , of (7.1), resp. (7.2), given, e.g., a solution V_1 of (7.1), such that (H.7.1) and (7.3) are satisfied. After reviewing the results of [50] (see, in particular, Theorem 7.9) we shall apply them in the cases of soliton-like and periodic solutions of the (m)KdV-equation.

Assume V_1 satisfies (H.7.1) and the KdV-equation (7.1) and introduce

$$(7.13) \quad H_1(t) = -\partial_x^2 + V_1(t, \cdot), \quad \mathfrak{D}(H_1(t)) = H^2(\mathbb{R}), \quad t \in \mathbb{R}.$$

By Theorem 7.3,

$$(7.14) \quad \mathfrak{E}_1 := \inf [\sigma(H_1(t))]$$

is t -independent. As will become clear in the course of this section, equation (7.3), being a Riccati-type equation for ϕ given V_j , will have nonsingular solutions ϕ only if the associated Schrödinger operators $H_j(t) = -\partial_x^2 + V_j(t, \cdot)$, $j = 1, 2$, are nonnegative. Thus we assume from now on (cf. also the discussion following (7.76)–(7.82)) that

$$(7.15) \quad \mathfrak{E}_1 \geq 0.$$

An important role in our construction is played by positive, distributional solutions ψ_1 of

$$(7.16) \quad (H_1(t)\psi_1(t))(x) = 0, \quad (t, x) \in \mathbb{R}^2$$

(i.e., $\psi_1(t, \cdot), \psi_{1,x}(t, \cdot) \in AC_{\text{loc}}(\mathbb{R})$, $t \in \mathbb{R}$). These positive solutions may be classified into principal $\tilde{\psi}_{1,\pm}(t, x) > 0$ and nonprincipal $\hat{\psi}_{1,\pm}(t, x) > 0$ solutions of (7.16) at $x = \pm\infty$, being defined by [56]

$$(7.17) \quad \pm \int_0^{\pm\infty} dx \tilde{\psi}_{1,\pm}(t, x)^{-2} = \infty, \quad \pm \int_0^{\pm\infty} dx \hat{\psi}_{1,\pm}(t, x)^{-2} < \infty,$$

where t is held fixed at the moment. If $\hat{\psi}_{1,\pm}(t, x) > 0$ is a nonprincipal solution of (7.16) at $x = \pm\infty$ then

$$(7.18) \quad \tilde{\psi}_{1,\pm}(t, x) := \pm \hat{\psi}_{1,\pm}(t, x) \int_x^{\pm\infty} dx' \hat{\psi}_{1,\pm}(t, x')^{-2}$$

is principal at $x = \pm\infty$.

Remark 7.4. Fix $t \in \mathbb{R}$ and assume $H_1(t) \geq 0$. In the terminology of [100] (cf. also [90]), $H_1(t)$ is critical iff $W(\tilde{\psi}_{1,-}(t), \tilde{\psi}_{1,+}(t)) = 0$ and hence $H_1(t)$ is subcritical iff $W(\tilde{\psi}_{1,-}(t), \tilde{\psi}_{1,+}(t)) \neq 0$. Or, equivalently, $H_1(t)$ is subcritical iff there exist two linearly independent, positive, distributional zero-energy solutions of $H_1(t)$ and hence $H_1(t)$ is critical iff it has a unique (up to multiples of constants) positive, distributional zero-energy solution.

We also recall

Lemma 7.5. Assume V_1 satisfies (H.7.1). Let $0 < \psi_1 \in C^\infty(\mathbb{R}^2)$ be a distributional solution of $H_1(t)\psi_1(t) = 0$, $t \in \mathbb{R}$, and define

$$(7.19) \quad \phi(t, x) = -\psi_{1,x}(t, x)/\psi_1(t, x), \quad (t, x) \in \mathbb{R}^2.$$

Then ϕ satisfies (H.7.1) and $\partial_x^2 \phi \in L^\infty(\mathbb{R}^2)$, and $V_2 = \phi^2 + \phi_x$ satisfies (H.7.1).

Proof. As shown in Corollary XI.6.5 of [56], $V_1 \in L^\infty(\mathbb{R}^2)$ yields

$$(7.20) \quad \phi \in L^\infty(\mathbb{R}^2).$$

Since $V_1 = \phi^2 - \phi_x$ we infer $\phi_x \in L^\infty(\mathbb{R}^2)$. The rest is trivial. \square

Next we study the time dependence of solutions of (7.16) in more detail.

Theorem 7.6. Assume V_1 satisfies (H.7.1) and the KdV-equation (7.1). Let $\psi_{1,\pm} \in C^\infty(\mathbb{R})$ be real-valued, distributional solutions of $H_1(0)\psi_{1,\pm} = 0$. Then $H_1(t)\psi_{1,\pm}(t) = 0$, $t \in \mathbb{R}$, has unique real-valued, distributional solutions $\psi_{1,\pm} \in C^\infty(\mathbb{R}^2)$ that evolve according to $U_1(t)$ in (7.11), (7.12) (with $H(t) \equiv H_1(t)$), i.e.,

$$(7.21a) \quad \begin{aligned} \psi_{1,\pm,t}(t, x) = & -4\psi_{1,\pm,xxx}(t, x) + 6V_1(t, x)\psi_{1,\pm,x}(t, x) \\ & + 3V_{1,x}(t, x)\psi_{1,\pm}(t, x), \end{aligned}$$

or, equivalently, using $\psi_{1,\pm,xx} = V_1\psi_{1,\pm}$,

$$(7.21b) \quad \begin{aligned} \psi_{1,\pm,t}(t, x) &= 2V_1(t, x)\psi_{1,\pm,x}(t, x) \\ &\quad - V_{1,x}(t, x)\psi_{1,\pm}(t, x), \quad (t, x) \in \mathbb{R}^2, \end{aligned}$$

with

$$(7.21c) \quad \psi_{1,\pm}(0, x) = \psi_{1,\pm}(x), \quad x \in \mathbb{R}.$$

In particular, the Wronskian

$$(7.22) \quad W(\psi_{1,-}(t), \psi_{1,+}(t)) = W(\psi_{1,-}, \psi_{1,+})$$

is independent of $(t, x) \in \mathbb{R}^2$.

Proof. Consider the Volterra-integral equation

$$(7.23) \quad \psi_1(t, x) = c(t) + d(t)x + \int_0^x dx'(x - x')V_1(t, x')\psi_1(t, x'),$$

$$c, d \in C^\infty(\mathbb{R}), \quad (t, x) \in \mathbb{R}^2.$$

Then, iterating (7.23), one obtains $\psi_1 \in C^\infty(\mathbb{R}^2)$ and $H_1(t)\psi_1(t) = 0$ in the distributional sense. Moreover,

$$(7.24) \quad \psi_{1,txx} - V_1\psi_{1,t} = V_{1,t}\psi_1.$$

Next, define

$$(7.25) \quad \Psi(t, x) := 2V_1(t, x)\psi_{1,x}(t, x) - V_{1,x}(t, x)\psi_1(t, x), \quad (t, x) \in \mathbb{R}^2.$$

Then (7.1) implies

$$(7.26) \quad \Psi_{xx} - V_1\Psi = V_{1,t}\psi_1,$$

i.e.,

$$(7.27) \quad [\Psi - \psi_{1,t}]_{xx} = V_1[\Psi - \psi_{1,t}].$$

Since

$$(7.28) \quad \begin{aligned} \psi_{1,t}(t, 0) &= \dot{c}(t), \quad \psi_{1,tx}(t, 0) = \dot{d}(t), \\ \Psi(t, 0) &= 2V_1(t, 0)d(t) - V_{1,x}(t, 0)c(t), \\ \Psi_x(t, 0) &= V_{1,x}(t, 0)d(t) + [2V_1(t, 0)^2 - V_{1,xx}(t, 0)]c(t), \quad t \in \mathbb{R}, \end{aligned}$$

we get

$$(7.29) \quad \begin{aligned} \psi_{1,t}(t, x) &= \Psi(t, x) \\ &= 2V_1(t, x)\psi_{1,x}(t, x) - V_{1,x}(t, x)\psi_1(t, x), \quad (t, x) \in \mathbb{R}^2, \end{aligned}$$

iff

$$(7.30) \quad \frac{d}{dt} \begin{pmatrix} c(t) \\ d(t) \end{pmatrix} = \begin{pmatrix} -V_{1,x}(t, 0) & 2V_1(t, 0) \\ 2V_1(t, 0)^2 - V_{1,xx}(t, 0) & V_{1,x}(t, 0) \end{pmatrix} \begin{pmatrix} c(t) \\ d(t) \end{pmatrix}, \quad t \in \mathbb{R}.$$

Finally (7.22) follows by a simple calculation using $H_1(t)\psi_{1,\pm}(t) = 0$ and (7.21). \square

Unless otherwise stated, from now on the time-dependence of distributional solutions $\psi_1(t, x)$ of $H_1(t)\psi_1(t) = 0$ is always chosen according to (7.21) (occasionally we also use the notation $\psi_1(t, x) = (U_1(t)\psi_1(0))(x)$ to stress that fact). Next we recall

Lemma 7.7 [78]. *Assume V_1 satisfies (H.7.1) and the KdV-equation (7.1). Let $\psi_1 \in C^\infty(\mathbb{R}^2)$ be a distributional solution of $H_1(t)\psi_1(t) = 0$, $t \in \mathbb{R}$, evolving in time according to (7.21). If $\psi_1(t, x(t)) = 0$, $t \in \mathbb{R}$, then $x(t)$ solves*

$$(7.31) \quad \dot{x}(t) = -2V_1(t, x), \quad t \in \mathbb{R}.$$

Conversely, if $\psi_1(t_0, x_0) = 0$, solve

$$(7.32) \quad \dot{x}(t) = -2V_1(t, x), \quad t \in \mathbb{R}, x(t_0) = x_0,$$

to get $\psi_1(t, x(t)) = 0$, $t \in \mathbb{R}$. In particular, if $\psi_1(0, x) > 0$, $x \in \mathbb{R}$, then

$$(7.33) \quad \psi_1(t, x) > 0, \quad (t, x) \in \mathbb{R}^2.$$

Proof. Equation (7.31) is obvious from $d\psi = 0$. Conversely, if $x(t)$, $t \in \mathbb{R}$, is the unique solution of (7.32) (which exists by standard considerations if, e.g., $V_1 \in C^1(\mathbb{R}^2)$, $V_1, V_{1,x} \in L^\infty(\mathbb{R}^2)$) then

$$(7.34) \quad d\psi_1 = \psi_{1,t}dt + \psi_{1,x}dx = -V_{1,x}\psi_1dt$$

yields $\psi_1(t, x(t)) = 0$ since $\psi_1(t_0, x_0) = 0$. Finally, if $\psi_1(0, x) > 0$, $x \in \mathbb{R}$, assume that $\psi_1(t_0, x_0) = 0$ for some $(t_0, x_0) \in \mathbb{R}^2$. Then, propagating ψ_1 from t_0 to 0 in time, $\psi_1(0, x(0))$ would be zero as shown above. \square

Given these results we now may strengthen Remark 7.4.

Lemma 7.8. *Assume V_1 satisfies (H.7.1) and the KdV-equation (7.1). Then $H_1(t)$ is subcritical (resp. critical) for all $t \in \mathbb{R}$ iff $H_1(0)$ is subcritical (resp. critical).*

Proof. $H_1(0)$ is subcritical iff there exist solutions

$$(7.35) \quad \begin{aligned} \psi_{1,\pm}(0, x) &> 0, \quad x \in \mathbb{R}, \\ H_1(0)\psi_{1,\pm}(0) &= 0, \quad W(\psi_{1,-}(0), \psi_{1,+}(0)) \neq 0. \end{aligned}$$

By Theorem 7.6 and Lemma 7.7,

$$(7.36) \quad \begin{aligned} \psi_{1,\pm}(t, x) &:= (U_1(t)\psi_{1,\pm}(0))(x) > 0, \\ (t, x) &\in \mathbb{R}^2, H_1(t)\psi_{1,\pm}(t) = 0, t \in \mathbb{R}, \end{aligned}$$

and

$$(7.37) \quad W(\psi_{1,-}(t), \psi_{1,+}(t)) = W(\psi_{1,-}(0), \psi_{1,+}(0)) \neq 0, \quad t \in \mathbb{R},$$

i.e., $H_1(t)$ is subcritical for all $t \in \mathbb{R}$. \square

Our main result finally reads [50]

Theorem 7.9. *Assume V_1 satisfies (H.7.1) and the KdV-equation (7.1). Let $0 < \psi_{1,\pm} \in C^\infty(\mathbb{R}^2)$ be positive, distributional solutions of $H_1(t)\psi_1(t) = 0$, $t \in \mathbb{R}$, evolving in time according to (7.21). Define*

$$(7.38) \quad \begin{aligned} \psi_{1,\sigma}(t, x) := & 2^{-1} [1 - \sigma(t)] \psi_{1,-}(t, x) \\ & + 2^{-1} [1 + \sigma(t)] \psi_{1,+}(t, x), \quad (t, x) \in \mathbb{R}^2, \end{aligned}$$

$$(7.39) \quad \phi_\sigma(t, x) := -\psi_{1,\sigma,x}(t, x)/\psi_{1,\sigma}(t, x), \quad (t, x) \in \mathbb{R}^2,$$

$$(7.40) \quad V_{2,\sigma}(t, x) := \phi_\sigma(t, x)^2 + \phi_{\sigma,x}(t, x), \quad (t, x) \in \mathbb{R}^2,$$

where $\sigma: \mathbb{R} \rightarrow [-1, 1]$, $\sigma \in C^\infty(\mathbb{R})$. Then ϕ_σ and $V_{2,\sigma}$ satisfy (H.7.1). In addition,

$$(7.41) \quad \text{mKdV}(\phi_\sigma) = 0, \quad \text{KdV}(V_{2,\sigma}) = 0 \text{ iff } \dot{\sigma} = 0 \text{ or } W(\psi_{1,-}, \psi_{1,+}) = 0.$$

Proof. A lengthy, though straightforward computation yields

$$(7.42) \quad \begin{aligned} \text{mKdV}(\phi) = \psi_1^{-2} \{ & \psi_{1,t} \psi_{1,x} - \psi_1 \psi_{1,xt} - 6V_1 \psi_{1,x}^2 + 3\psi_{1,xx}^2 \\ & + 4\psi_{1,x} \psi_{1,xxx} - \psi_1 \psi_{1,xxxx} \}, \end{aligned}$$

where $\phi := -\psi_{1,x}/\psi_1$ and $\psi_1 \in C^\infty(\mathbb{R}^2)$ is any positive, distributional solution of $H_1(t)\psi_1(t) = 0$ not necessarily evolving in time according to (7.21). Specializing to $\psi_1 = \psi_{1,\sigma}$ and taking into account (7.21) for $\psi_{1,\pm}$ finally leads to

$$(7.43) \quad \text{mKdV}(\phi_\sigma) = -\psi_{1,\sigma}^{-2} \dot{\sigma} W(\psi_{1,-}, \psi_{1,+})/2,$$

$$\text{KdV}(V_{2,\sigma}) = [2\phi_\sigma + \partial_x] \text{mKdV}(\phi_\sigma) = 2\psi_{1,\sigma}^{-3} \psi_{1,\sigma,x} \dot{\sigma} W(\psi_{1,-}, \psi_{1,+}). \quad \square$$

Remark 7.10. If $H_1(0)$ is critical, we necessarily have $W(\psi_{1,-}, \psi_{1,+}) = 0$ by Remark 7.4 and hence ϕ_σ is actually independent of σ . If $H_1(0)$ is subcritical, then Theorem 7.9 guarantees a one-parameter family of solutions ϕ_σ , $\sigma \in [-1, 1]$, of the mKdV-equation (7.2) (choose $\psi_{1,\pm}$ such that $W(\psi_{1,-}, \psi_{1,+}) \neq 0$). Moreover, the explicit construction (7.39) yields all smooth solutions ϕ of $\text{mKdV}(\phi) = 0$ related to V_1 via $V_1 = \phi^2 - \phi_x$. (This is obvious in the critical case since the positive solution of $H_1(t)\psi_1(t) = 0$ is unique up to multiples of constants. In the subcritical case choose again $\psi_{1,\pm}$ such that $W(\psi_{1,-}, \psi_{1,+}) \neq 0$ and observe that $\phi_\sigma = -\psi_{1,\sigma,x}/\psi_{1,\sigma}$, $\sigma \in [-1, 1]$, is the general solution of the Riccati-equation $\phi_x = \phi^2 - V_1$ on \mathbb{R} .) Finally we remark that in the case where H_1 is critical, $\psi_1(t)$ need not necessarily satisfy (7.21) since another (time-dependent) multiple of it (which drops out in the definition of ϕ) will satisfy (7.21) by Theorem 7.6.

Remark 7.11. The “if” part in Theorem 7.9 is known and follows, e.g., from prolongation methods developed in [117] (see also [47, 98, 119] and the references therein). In general, these techniques do not distinguish between singular

and nonsingular solutions ϕ of the mKdV-equation. It is the “only if” part in Theorem 7.9 that yields a classification into singular and nonsingular solutions ϕ (depending on whether $H_1(0)$ is nonnegative or not) and a uniqueness or nonuniqueness result for ϕ in Miura’s transformation (7.3) depending on whether $H_1(0)$ is critical or subcritical. A different approach to this problem can be found in §38 of [12]. There the authors use the fact that the scattering data for an n th order ordinary differential operator can also be viewed as scattering data of a certain first order system. When specializing to Schrödinger operators (i.e., $n = 2$) and solving an associated matrix factorization problem they are able to treat the Cauchy problem for solutions of the KdV- and mKdV-equations related to each other via Miura’s transformation. Their approach assumes rapidly decreasing solutions of the KdV-equation at $\pm\infty$.

Remark 7.12. (i) Given ϕ_σ according to (7.39), $H_1(t)$ is recovered from ϕ_σ via

$$(7.44) \quad \begin{aligned} H_1(t) &= -\partial_x^2 + V_1(t, \cdot) = A_\sigma(t)^* A_\sigma(t), & \mathfrak{D}(H_1(t)) &= H^2(\mathbb{R}), \\ A_\sigma(t) &= \partial_x + \phi_\sigma(t, \cdot), & \mathfrak{D}(A_\sigma(t)) &= H^1(\mathbb{R}), \\ V_1(t, x) &= \phi_\sigma(t, x)^2 - \phi_{\sigma, x}(t, x), & (t, x) &\in \mathbb{R}^2. \end{aligned}$$

Similarly, $H_{2, \sigma}(t)$ is recovered from ϕ_σ via

$$(7.45) \quad \begin{aligned} H_{2, \sigma}(t) &= A_\sigma(t) A_\sigma(t)^* = -\partial_x^2 + V_{2, \sigma}(t, \cdot), & \mathfrak{D}(H_{2, \sigma}(t)) &= H^2(\mathbb{R}), \\ V_{2, \sigma}(t) &= \phi_\sigma(t, x)^2 + \phi_{\sigma, x}(t, x), & (t, x) &\in \mathbb{R}^2. \end{aligned}$$

(ii) If one starts with $H_2(t)$ instead of $H_1(t)$, then ϕ_σ is defined by

$$(7.46) \quad \phi_\sigma(t, x) := \psi_{2, \sigma, x}(t, x) / \psi_{2, \sigma}(t, x), \quad (t, x) \in \mathbb{R}^2,$$

where

$$(7.47) \quad \begin{aligned} \psi_{2, \sigma}(t, x) &:= 2^{-1} [1 - \sigma(t)] \psi_{2, -}(t, x) \\ &\quad + 2^{-1} [1 + \sigma(t)] \psi_{2, +}(t, x), \quad (t, x) \in \mathbb{R}^2, \end{aligned}$$

and $0 < \psi_{2, \pm} \in C^\infty(\mathbb{R}^2)$ are positive, distributional solutions of $H_2(t)\psi_2(t) = 0$ evolving in time according to (7.21) (with V_1 replaced by V_2).

In the case where both V_1 and V_2 are given and linked to ϕ via Miura’s transformation (7.3), one reconstructs ϕ as follows:

Lemma 7.13. Assume that V_j , $j = 1, 2$, satisfy (H.7.1) and Miura’s transformation (7.3) for some $\phi \in C^\infty(\mathbb{R}^2)$. Then ϕ satisfies (H.7.1) and

$$(7.48) \quad \phi(t, x) = \{[V_{1, x}(t, x) + V_{2, x}(t, x)] / 2[V_2(t, x) - V_1(t, x)]\}, \quad (t, x) \in \mathbb{R}^2.$$

Moreover, if V_1 and V_2 satisfy the KdV-equation (7.1), then ϕ satisfies the mKdV-equation (7.2).

Proof. Equation (7.48) follows from

$$(7.49) \quad 2\phi_x = V_2 - V_1, \quad 2\phi^2 = V_1 + V_2, \quad 4\phi\phi_x = V_{1, x} + V_{2, x}.$$

The rest is obvious from (7.4). \square

Next we illustrate the content of Theorem 7.9 in two important special cases, viz., soliton-like and periodic solutions V (resp. ϕ) of the (m)KdV-equation. For the soliton-like case we introduce hypothesis

(H.7.2)

(i) $V \in C^\infty(\mathbb{R}^2)$ real-valued, $\partial_x^n V \in L^\infty(\mathbb{R}^2)$, $n = 0, 1$, $\lim_{x \rightarrow \pm\infty} V(t, x) = V_\pm \in \mathbb{R}$, $0 < V_- \leq V_+$, $\pm \int_0^{\pm\infty} dx(1+x^2)|V(t, x) - V_\pm| < \infty$, $t \in \mathbb{R}$.

(ii) $\phi \in C^\infty(\mathbb{R}^2)$ real-valued, $\partial_x^n \phi \in L^\infty(\mathbb{R}^2)$, $n = 0, 1$, $\lim_{x \rightarrow \pm\infty} \phi(t, x) = \phi_\pm \in \mathbb{R}$, $0 < \phi_-^2 \leq \phi_+^2$, $\pm \int_0^{\pm\infty} dx(1+x^2)|\phi(t, x) - \phi_\pm| < \infty$, $\int_{\mathbb{R}} dx(1+x^2)|\phi_x(t, x)| < \infty$, $t \in \mathbb{R}$.

(Here V_\pm , ϕ_\pm are t -independent constants.)

For the periodic case we introduce

(H.7.3) $V, \phi \in C^\infty(\mathbb{R}^2)$ real-valued,

$$\text{for some } a > 0: V(t, x+a) = V(t, x), \quad (t, x) \in \mathbb{R}^2,$$

$$\phi(t, x+a) = \phi(t, x), \quad (t, x) \in \mathbb{R}^2.$$

We first study systems governed by (H.7.2). Assume (cf. (7.15))

$$\mathfrak{E}_1 = 0 \quad (\text{i.e., } H_1(t), t \in \mathbb{R}, \text{ is critical}).$$

In this case (since $V_\pm > 0$)

$$(7.50) \quad H_1(t)\psi_{1,0}(t) = 0, \quad 0 < \psi_{1,0}(t, \cdot) \in H^2(\mathbb{R}), \quad t \in \mathbb{R},$$

has a unique solution $\psi_{1,0} \in C^\infty(\mathbb{R}^2)$ up to multiples of constants (this follows, e.g., from nonoscillation theorems [37]) which, without loss of generality (cf. the end of Remark 7.10), may be assumed to evolve according to (7.21). Then one has

$$(7.51) \quad \psi_{1,0}(t, x) \underset{x \rightarrow \pm\infty}{=} e^{\mp V_\pm^{1/2}(x+2V_\pm t)} + o(e^{\mp V_\pm^{1/2}x}), \quad t \in \mathbb{R},$$

and thus, according to (7.39), defines

$$(7.52) \quad \phi_0(t, x) := -\psi_{1,0,x}(t, x)/\psi_{1,0}(t, x), \quad (t, x) \in \mathbb{R}^2.$$

Next suppose that

$$\mathfrak{E}_1 > 0 \quad (\text{i.e., } H_1(t), t \in \mathbb{R}, \text{ is subcritical}).$$

In this case we consider

$$(7.53) \quad H_1(t)\psi_{1,\pm}(t) = 0, \quad 0 < \psi_{1,\pm}(t, \cdot) \in L_{\text{loc}}^\infty(\mathbb{R}), \quad t \in \mathbb{R},$$

where again $\psi_{1,\pm}$ evolve in time according to (7.21) and hence

$$(7.54) \quad \begin{aligned} \psi_{1,+}(t, x) &\underset{x \rightarrow \pm\infty}{=} e^{-V_+^{1/2}(x+2V_+ t)} + o(e^{-V_+^{1/2}x}), \\ \psi_{1,-}(t, x) &\underset{x \rightarrow \pm\infty}{=} e^{V_-^{1/2}(x+2V_- t)} + o(e^{V_-^{1/2}x}), \end{aligned} \quad t \in \mathbb{R},$$

and define $\psi_{1,\sigma}$ and ϕ_σ as in (7.38) and (7.39), i.e.,

$$(7.55) \quad \phi_\sigma(t, x) := -\psi_{1,\sigma,x}(t, x)/\psi_{1,\sigma}(t, x), \quad (t, x) \in \mathbb{R}^2.$$

Then we have

Theorem 7.14. *Assume that V_1 satisfies (H.7.2(i)). Then ϕ_σ , defined in (7.52), resp. (7.55), satisfies (H.7.2(ii)) and $V_{2,\sigma} = \phi_\sigma^2 + \phi_{\sigma,x}$, $\sigma \in [-1, 1]$, satisfies (H.7.2(i)). Moreover,*

$$(7.56) \quad \phi_{\sigma,\pm} := \lim_{x \rightarrow \pm\infty} \phi_\sigma(t, x) = \begin{cases} \pm V_\pm^{1/2}, & \sigma = 0, \mathfrak{E}_1 = 0, \\ V_\pm^{1/2}, & \sigma = +1, \mathfrak{E}_1 > 0, \\ -V_\pm^{1/2}, & \sigma = -1, \mathfrak{E}_1 > 0, \\ \mp V_\pm^{1/2}, & \sigma \in (-1, 1), \mathfrak{E}_1 > 0, \end{cases}$$

and

$$(7.57) \quad \begin{aligned} 0 &\in \sigma_d(H_1(t)), \quad 0 \notin \sigma_d(H_{2,\sigma}(t)), \quad \phi_{\sigma,-} < 0 < \phi_{\sigma,+}, \\ 0 &\notin \sigma_d(H_1(t)), \quad 0 \in \sigma_d(H_{2,\sigma}(t)), \quad \phi_{\sigma,+} < 0 < \phi_{\sigma,-}, \\ 0 &\notin \{\sigma_d(H_1(t)) \cup \sigma_d(H_{2,\sigma}(t))\}, \quad \text{sgn}(\phi_{\sigma,-}) = \text{sgn}(\phi_{\sigma,+}), \\ &t \in \mathbb{R}, \sigma \in [-1, 1]. \end{aligned}$$

Proof. Let $\sigma = 0$, $\mathfrak{E}_1 = 0$. Then standard Volterra-integral equation techniques yield (cf., e.g., [49])

$$(7.58) \quad \begin{pmatrix} \psi_{1,0}(t, x) \\ \psi_{1,0,x}(t, x) \end{pmatrix}_{x \rightarrow \pm\infty} = \begin{pmatrix} 1 \\ \mp V_\pm^{1/2} \end{pmatrix} e^{\mp V_\pm^{1/2}(x+2V_\pm t)} + o(e^{\mp V_\pm^{1/2}x}),$$

$t \in \mathbb{R},$

and

$$(7.59) \quad \phi_0(t, x)_{x \rightarrow \pm\infty} = \pm V_\pm^{1/2} + O(|x|^{-2})g_\pm(t, x), \quad t \in \mathbb{R},$$

where $g_\pm(t, \cdot) \in L^1((0, \pm\infty))$, $t \in \mathbb{R}$. Thus

$$(7.60) \quad \begin{aligned} &\pm \int_0^{\pm\infty} dx(1 + |x|^2)|\phi_{0,x}(t, x)| \\ &= \pm \int_0^{\pm\infty} dx(1 + |x|^2)|V_1(t, x) - \phi_0(t, x)|^2 \\ &\leq \pm \int_0^{\pm\infty} dx(1 + |x|^2)|V_1(t, x) - V_\pm| \\ &\pm \int_0^{\pm\infty} dx(1 + |x|^2)|\phi_0(t, x)^2 - V_\pm| < \infty, \quad t \in \mathbb{R}, \end{aligned}$$

by (H.7.2(i)) and (7.59). Consequently, also

$$(7.61) \quad \pm \int_0^{\pm\infty} dx(1 + |x|^2)|\phi_0(t, x) \mp V_\pm^{1/2}| < \infty, \quad t \in \mathbb{R}.$$

Similarly,

$$\begin{aligned}
 & \pm \int_0^{\pm\infty} dx(1+|x|^2)|V_{2,0}(t, x) - V_{\pm}| \\
 (7.62) \quad & \leq \pm \int_0^{\pm\infty} dx(1+|x|^2)|\phi_0(t, x)^2 - V_{\pm}| \\
 & \pm \int_0^{\pm\infty} dx(1+|x|^2)|\phi_{0,x}(t, x)| < \infty, \quad t \in \mathbb{R}.
 \end{aligned}$$

The case $\mathfrak{E}_1 > 0$ follows similarly. (7.57) is discussed, e.g., in [16, 49]. \square

Remark 7.15. A comparison of (7.56) and (7.57) together with Remark 3.8 shows that, whenever $\text{sgn}(\phi_{\sigma,-}) \neq \text{sgn}(\phi_{\sigma,+})$, then $\mathfrak{E}_1 = 0$ or $\mathfrak{E}_2 = 0$ and hence either H_1 is critical and $H_{2,\sigma}$ is subcritical or vice versa. In particular, $H_{2,\sigma}$ is critical if $\mathfrak{E}_1 > 0$, $\sigma \in (-1, 1)$. Only in the case where $\text{sgn}(\phi_{\sigma,-}) = \text{sgn}(\phi_{\sigma,+})$ one has $\mathfrak{E}_1 = \mathfrak{E}_2 > 0$ and thus both H_1 and $H_{2,\sigma}$ are subcritical iff $\sigma = \pm 1$.

Remark 7.16. In order to avoid too many case distinctions in the definition of ϕ_{σ} we have restricted ourselves to the case $0 < V_- \leq V_+$. Clearly $V_- = 0$ (i.e., $\sigma_p(H_1(t)) = \emptyset$, $\sigma(H_1(t)) = \sigma_{\text{ac}}(H_1(t)) = [0, \infty)$, $t \in \mathbb{R}$) could be discussed along exactly the same lines by considering zero-energy (i.e., threshold) wave functions of $H_1(t)$ behaving asymptotically like $c + dx$ as $|x| \rightarrow \infty$ (see, e.g., [14, 15, 72]).

Remark 7.17. If V_j , $j = 1, 2$, in Lemma 7.13 actually satisfy (H.7.2(i)) then ϕ satisfies (H.7.2(ii)) and we get

$$\begin{aligned}
 (7.63) \quad & \phi_+ - \phi_- = \frac{1}{2} \int_{\mathbb{R}} dy [V_2(t, y) - V_1(t, y)], \quad t \in \mathbb{R}, \\
 & \phi(t, x) = \phi_{\pm} - \frac{1}{2} \int_x^{\pm\infty} dy [V_2(t, y) - V_1(t, y)], \quad (t, x) \in \mathbb{R}^2,
 \end{aligned}$$

in addition to (7.48).

Next we turn to systems satisfying (H.7.3). Again we first treat the case

$$\mathfrak{E}_1 = 0 \quad (\text{i.e., } H_1(t), \quad t \in \mathbb{R}, \text{ is critical}).$$

We consider

$$(7.64) \quad H_1(t)\psi_{1,0}(t) = 0, \quad 0 < \psi_{1,0}(t, \cdot) \in L^{\infty}(\mathbb{R}), \quad t \in \mathbb{R},$$

and suppose that $\psi_{1,0}$ evolves in time according to (7.21). By Floquet theory (take $k = 0$ in (3.35)), $\psi_{1,0}$ is periodic in x , i.e.,

$$(7.65) \quad \psi_{1,0}(t, x+a) = \psi_{1,0}(t, x), \quad (t, x) \in \mathbb{R}^2,$$

and we define according to (7.39)

$$(7.66) \quad \phi_0(t, x) := -\psi_{1,0,x}(t, x)/\psi_{1,0}(t, x), \quad (t, x) \in \mathbb{R}^2.$$

In the case

$$\mathfrak{E}_1 > 0 \quad (\text{i.e., } H_1(t), \quad t \in \mathbb{R}, \text{ is subcritical})$$

we consider

(7.67)

$$H_1(t)\psi_{1,\pm}(t) = 0 \quad (\text{distributional sense}),$$

$$0 < \psi_{1,\pm}(t, \cdot) \in L_{\text{loc}}^\infty(\mathbb{R}), \quad t \in \mathbb{R},$$

where

$$(7.68) \quad \psi_{1,\pm}(t, x) \underset{x \rightarrow \pm\infty}{=} e^{\mp\kappa x} + o(e^{\mp\kappa x}), \quad t \in \mathbb{R},$$

$$(7.69) \quad \kappa = -ik(0) > 0,$$

(cf. (3.34)) and $\psi_{1,\pm}$ satisfy (7.21). In accordance with (7.39) we then define

$$(7.70) \quad \phi_{\pm 1}(t, x) := -\psi_{1,\pm,x}(t, x)/\psi_{1,\pm}(t, x), \quad (t, x) \in \mathbb{R}^2.$$

We have

Theorem 7.18. Assume that V_1 satisfies (H.7.3). Then ϕ_σ , defined in (7.66), resp. (7.70), and $V_{2,\sigma} = \phi_\sigma^2 + \phi_{\sigma,x}$, $\sigma = 0, \pm 1$, both satisfy (H.7.3). In particular,

$$(7.71) \quad \begin{aligned} \sigma(H_1(t)) &= \sigma(H_{2,\sigma}(t)) = \sigma(H_1(0)), \\ \sigma_p(H_{j,(\sigma)}(t)) &= \emptyset, \quad j = 1, 2, \quad t \in \mathbb{R}, \quad \sigma = 0, \pm 1. \end{aligned}$$

Proof. Because of (7.64), (7.65), (7.67), and (7.68), we have (cf. (3.35))

$$(7.72) \quad \psi_{1,\sigma}(t, x) \geq \varepsilon(t) > 0, \quad t \in \mathbb{R}, \quad x \in [0, a], \quad \sigma = 0, \pm 1.$$

Using again standard Volterra-integral equations for $\psi_{1,\sigma}$ one establishes the corresponding differentiability properties of $\psi_{1,\sigma}$, $\psi_{1,\sigma,x}$ and hence that of ϕ_σ . \square

Remark 7.19. Clearly $\phi_\sigma(t, x)$, defined according to (7.39), is not periodic in x for $\sigma \in (-1, 1)$ in the subcritical case. Thus we confined ourselves to $\sigma = \pm 1$ in (7.70).

Remark 7.20. In the case where $\mathfrak{E}_1 = \mathfrak{E}_2 > 0$ we have

$$(7.73) \quad \begin{aligned} H_1 &= A(t)^* A(t) = \hat{H}_1(t) + \mathfrak{E}_1, \quad \hat{H}_1(t) = \hat{A}(t)^* \hat{A}(t), \\ \hat{A}(t) &= \partial_x + \hat{\phi}_\sigma(t, \cdot), \quad \mathfrak{D}(\hat{A}(t)) = H^1(\mathbb{R}), \quad \sigma = \pm 1, \quad t \in \mathbb{R}, \end{aligned}$$

where $\hat{\phi}_\sigma$ satisfies (H.7.2(ii)) or (H.7.3) and

$$(7.74) \quad \begin{aligned} \phi_\sigma(t, x)^2 - \phi_{\sigma,x}(t, x) &= \hat{\phi}_\sigma(t, x)^2 - \hat{\phi}_{\sigma,x}(t, x) + \mathfrak{E}_1, \\ \sigma &= \pm 1, \quad (t, x) \in \mathbb{R}^2. \end{aligned}$$

In particular, if V_j satisfy (H.7.2(i)), then $0 \in \sigma_p(\hat{H}_1(t))$ and

$$(7.75) \quad \hat{\phi}_{\sigma,\pm} = \lim_{x \rightarrow \pm\infty} \hat{\phi}_\sigma(t, x) = \pm(V_\pm - \mathfrak{E}_1)^{1/2}.$$

(Similarly, $\hat{\phi}_\sigma \rightarrow -\hat{\phi}_\sigma$ if $H_1(t) \rightarrow H_2(t)$.)

Next we recall a few useful formulas. Define

$$(7.76) \quad u_j(t, x) := \phi(t, x)^2 + (-1)^j \phi_x(t, x) - \Phi_0, \\ (t, x) \in \mathbb{R}^2, \Phi_0 \in \mathbb{R}, j = 1, 2;$$

then

$$(7.77) \quad \text{KdV}(u_j) = [2\phi + (-1)^j \partial_x][\phi_t - 6(\phi^2 - \Phi_0)\phi_x + \phi_{xxx}], \quad j = 1, 2.$$

Next let

$$(7.78) \quad y = x - 6\Phi_0 t, \quad u(t, x) := w(t, y) - \Phi_0, \quad (t, x) \in \mathbb{R}^2;$$

then

$$(7.79) \quad u_t - 6uu_x + u_{xxx} = w_t - 6ww_y + w_{yyy}.$$

(This simple fact reproduces the results in [10].) Similarly, let

$$(7.80) \quad \phi(t, x) := \chi(t, y), \quad (t, x) \in \mathbb{R}^2;$$

then

$$(7.81) \quad \phi_t - 6(\phi^2 - \Phi_0)\phi_x + \phi_{xxx} = \chi_t - 6\chi^2\chi_y + \chi_{yyy}.$$

Equations (7.76)–(7.81) show that using the methods of this section one can construct solutions to the generalized mKdV-equation

$$(7.82) \quad \phi_t - 6(\phi^2 - \Phi_0)\phi_x + \phi_{xxx} = 0, \quad \Phi_0 \in \mathbb{R}.$$

This has been studied in [119] (some special solutions of (7.82) have been considered in [47]). Alternatively one could use the generalized Miura transformation (7.76) together with the Galilei transformation $(t, x) \rightarrow (t, y = x - 6\Phi_0 t)$ in V_j in order to treat general Hamiltonians $H_j(t)$ bounded from below by $\Phi_0 \in \mathbb{R}$, $H_j(t) = -\partial_x^2 + V_j(t, \cdot) \geq \Phi_0$, $j = 1, 2$.

Finally, we briefly indicate the generalization to higher-order (m)KdV-equations. These equations are recursively defined by

$$(7.83) \quad \begin{aligned} \text{KdV}_n(V) &:= V_t - 2X_{n,x} = 0, & n \in \mathbb{N}_0, \\ \text{mKdV}_n(\phi) &:= \phi_t - Y_{n,x} = 0, & n \in \mathbb{N}_0, \end{aligned}$$

where

$$(7.84) \quad \begin{aligned} X_0 &= 4, \\ X_{n+1,x} &= -\frac{1}{4}X_{n,xxx} + VX_{n,x} + \frac{1}{2}V_xX_n, & n \in \mathbb{N}_0, \\ Y_0 &= 4, \\ Y_{n+1,x} &= -\frac{1}{4}Y_{n,xxx} + \phi^2Y_{n,x} \\ &\quad + \phi_x \left[\int^x dx' (\phi Y_{n,x'}) + c_n \right], & n \in \mathbb{N}_0. \end{aligned}$$

Here, c_n are chosen to be the integration constants in X_n , $n \in \mathbb{N}_0$. Explicitly we have, e.g.,

$$\begin{aligned}
 (7.85) \quad & X_0 = c_0 = 4, \\
 & X_1 = 2V + c_1, \\
 & X_2 = -\frac{1}{2}V_{xx} + \frac{3}{2}V^2 + (c_1/2)V + c_2, \\
 & X_3 = \frac{1}{8}V_{xxxx} - \frac{5}{4}V V_{xx} - \frac{5}{8}V_x^2 + \frac{5}{4}V^3 \\
 & \quad - (c_1/8)V_{xx} + (3c_1/8)V^2 + (c_2/2)V + c_3, \\
 & Y_0 = d_0 = 4, \\
 & Y_1 = 4\phi + d_1, \\
 & Y_2 = -\phi_{xx} + 2\phi^3 + c_1\phi + d_2, \\
 & Y_3 = \frac{1}{4}\phi_{xxxx} - \frac{5}{2}\phi^2\phi_{xx} - \frac{5}{2}\phi\phi_x^2 + \frac{3}{2}\phi^5 \\
 & \quad - (c_1/4)\phi_{xx} + (c_1/2)\phi^3 + c_2\phi + d_3, \quad c_j, d_j \in \mathbb{R}, \quad j \in \mathbb{N},
 \end{aligned}$$

etc. Thus,

$$\begin{aligned}
 (7.86) \quad & \text{KdV}_0(V) = V_t = 0, \\
 & \text{KdV}_1(V) = V_t - 4V_x = 0, \\
 & \text{KdV}_2(V) \equiv \text{KdV}(V) = V_t + V_{xxx} - 6VV_x - c_1V_x = 0, \\
 & \text{KdV}_3(V) = V_t - \frac{1}{4}V_{xxxx} + \frac{5}{2}VV_{xx} + 5V_xV_{xx} - \frac{15}{2}V^2V_x \\
 & \quad + (c_1/4)V_{xxx} - (3c_1/2)VV_x - c_2V_x = 0, \\
 & \text{mKdV}_0(\phi) = \phi_t = 0, \\
 & \text{mKdV}_1(\phi) = \phi_t - 4\phi_x = 0, \\
 & \text{mKdV}_2(\phi) \equiv \text{mKdV}(\phi) = \phi_t + \phi_{xxx} - 6\phi^2\phi_x - c_1\phi_x = 0, \\
 & \text{mKdV}_3(\phi) = \phi_t - \frac{1}{4}\phi_{xxxx} + \frac{5}{2}\phi^2\phi_{xx} + 10\phi\phi_x\phi_{xx} + \frac{5}{2}\phi_x^3 \\
 & \quad - \frac{15}{2}\phi^4\phi_x + (c_1/4)\phi_{xxx} - (3c_1/2)\phi^2\phi_x - c_2\phi_x = 0,
 \end{aligned}$$

etc. (usually one chooses $c_j = d_j = 0$, $j \in \mathbb{N}$). We end up with

Remark 7.21. As proven in [6], (7.4) extends to the whole (m)KdV $_n$ -hierarchy introduced above, i.e.,

$$(7.87) \quad \text{KdV}_n(V_j) = [2\phi + (-1)^j \partial_x] \text{mKdV}_n(\phi), \quad j = 1, 2, \quad n \in \mathbb{N}_0,$$

where V_j , $j = 1, 2$, and ϕ are related by Miura's transformation (7.3). Checking the proof of our main Theorem 7.9 shows that it works as well in the present situation, i.e., Theorem 7.9 extends to the entire (m)KdV $_n$ -hierarchy, $n \in \mathbb{N}_0$ (in particular, (7.43) holds with mKdV replaced by mKdV $_n$, $n \in \mathbb{N}$) assuming $\partial_x^m V_1 \in L^\infty(\mathbb{R}^2)$, $0 \leq m \leq 2n - 3$, $n \geq 2$.

8. SOLITONS OF THE mKdV-EQUATION

In this section we rederive the N -soliton solutions of the mKdV-equation (originally obtained by Grosse [54, 55]) given the corresponding soliton solutions of the KdV-equation. In contrast to [54, 55] we do not use inverse scattering techniques for the Dirac operator.

Assume that $V_1(t, x)$ is an N -soliton solution of the KdV-equation, i.e., that according to Remark 6.3 (and (7.78), (7.79))

$$(8.1) \quad V_1(t, x) = V_\infty - 2\partial_x^2 \ln\{\det[1 + C_{1,\pm}(t, x)]\}, \quad (t, x) \in \mathbb{R}^2,$$

$$C_{1,\pm}(t, x) = [c_{1,\pm,l,m}(t, x)]_{l,m=1}^N, \quad N \in \mathbb{N},$$

$$(8.2) \quad c_{1,\pm,l,m}(t, x) = (\kappa_l + \kappa_m)^{-1} c_{1,\pm,l}(0) c_{1,\pm,m}(0) e^{\mp(\kappa_l + \kappa_m)x} e^{\pm 4(\kappa_l^3 + \kappa_m^3)t} \cdot e^{\mp(\kappa_l + \kappa_m)6V_\infty t}, \quad l, m = 1, \dots, N, \quad (t, x) \in \mathbb{R}^2,$$

$$(8.3) \quad 0 < \kappa_N < \kappa_{N-1} < \dots < \kappa_1, \quad \kappa_1^2 \leq V_\infty,$$

$$(8.4) \quad c_{1,\pm,l}(t) = c_{1,\pm,l}(0) e^{\pm 4\kappa_l^3 t}, \quad t \in \mathbb{R}, \quad l = 1, \dots, N,$$

$$(8.5) \quad H_1(t) = -\partial_x^2 + V_1(t, \cdot), \quad \mathfrak{D}(H_1(t)) = H^2(\mathbb{R}), \quad t \in \mathbb{R},$$

$$(8.6) \quad \sigma_p(H_1(t)) = \{\lambda_{1,l} := V_\infty - \kappa_l^2 \mid l = 1, \dots, N\},$$

$$(8.7) \quad \begin{aligned} H_1(0)f_{1,\pm,l}(0) &= (V_\infty - \kappa_l^2)f_{1,\pm,l}(0), \quad f_{1,\pm,l}(0) \in H^2(\mathbb{R}), \\ f_{1,\pm,l}(0, x) &\underset{x \rightarrow \pm\infty}{=} e^{\mp\kappa_l x} + o(e^{\mp\kappa_l x}), \end{aligned}$$

$$(8.8) \quad |c_{1,\pm,l}(0)| = \|f_{1,\pm,l}(0)\|_2^{-1}, \quad l = 1, \dots, N.$$

The assumption that $H_1(0)$ is critical, i.e.,

$$(8.9) \quad \mathfrak{E}_1 = \inf[\sigma(H_1(0))] = 0,$$

then yields

$$(8.10) \quad V_\infty = \kappa_1^2, \quad \lambda_{1,1} = 0.$$

Next we turn to the construction of $H_2(t)$. First we note that

$$(8.11) \quad \begin{aligned} H_1(t) &= A_0(t)^* A_0(t), \quad t \in \mathbb{R}, \\ A_0(t) &= \partial_x + \phi_0(t, \cdot), \quad \mathfrak{D}(A_0(t)) = H^1(\mathbb{R}), \quad t \in \mathbb{R}, \end{aligned}$$

and (cf. Lemma 7.14), since $\mathfrak{E}_1 = 0$,

$$(8.12) \quad \phi_{0,\pm} = \lim_{x \rightarrow \pm\infty} \phi_0(t, x) = \pm V_\infty^{1/2} = \pm \kappa_1.$$

Since $\kappa_1 > 0$ by hypothesis,

$$(8.13) \quad 0 \in \sigma_p(A_0(t)), \quad 0 \notin \sigma_p(A_0(t)^*), \quad t \in \mathbb{R}.$$

Thus commutation implies for $H_{2,0}(t) = A_0(t)A_0(t)^*$,

$$(8.14) \quad \begin{aligned} \sigma_p(H_{2,0}(t)) &= \{\lambda_{2,l} = \kappa_1^2 - \kappa_l^2 \mid l = 2, \dots, N\} \\ &= \sigma_p(H_1(t)) \setminus \{0\}, \quad t \in \mathbb{R}. \end{aligned}$$

We also recall that

$$(8.15) \quad \sigma_{\text{ess}}(H_{j,(0)}(t)) = \sigma_{\text{ac}}(H_{j,(0)}(t)) = [V_\infty, \infty), \quad t \in \mathbb{R}, j = 1, 2.$$

Moreover, since all eigenvalues of $H_{j,(0)}(t)$, $j = 1, 2$, are simple, we obtain from (3.23),

$$(8.16) \quad c_{2,\pm,0,l}(0) = [(\kappa_1 - \kappa_l)/(\kappa_1 + \kappa_l)]^{1/2} c_{1,\pm,l}(0), \quad l = 2, \dots, N.$$

Thus $V_{2,0}(t, x)$ is an $(N - 1)$ -soliton solution of the KdV-equation,

$$(8.17) \quad V_{2,0}(t, x) = \kappa_1^2 - 2\partial_x^2 \ln\{\det[1 + C_{2,\pm,0}(t, x)]\}, \quad (t, x) \in \mathbb{R}^2,$$

where

$$(8.18) \quad \begin{aligned} C_{2,\pm,0}(t, x) &= [c_{2,\pm,0,l,m}(t, x)]_{l,m=2}^N, \quad N \geq 2, \\ c_{2,\pm,0,l,m}(t, x) &= \left[\frac{(\kappa_{1\mp}\kappa_l)(\kappa_{1\mp}\kappa_m)}{(\kappa_{1\pm}\kappa_l)(\kappa_{1\pm}\kappa_m)} \right]^{1/2} c_{1,\pm,l,m}(t, x), \\ &\quad l, m = 2, \dots, N, \quad (t, x) \in \mathbb{R}^2. \end{aligned}$$

It remains to compute $\phi_0(t, x)$.

Theorem 8.1. Assume $\mathfrak{E}_1 = 0$ and let

$$(8.19) \quad V_{j,(0)}(t, x) = \kappa_1^2 - 2\partial_x^2 \ln\{\det[1 + C_{j,\pm,(0)}(t, x)]\}, \quad (t, x) \in \mathbb{R}^2, \quad j = 1, 2$$

(cf. (8.1) and (8.17)), be an N -soliton, resp. $(N - 1)$ -soliton, solution for the KdV-equation (7.1), $N \in \mathbb{N}$. Let V_1 , $V_{2,0}$, and ϕ_0 be connected by the Miura transformation (7.3). Then

$$(8.20) \quad \phi_0(t, x) = \pm\kappa_1 + \partial_x \ln\{\det[1 + C_{1,\pm}(t, x)] / \det[1 + C_{2,\pm,(0)}(t, x)]\}, \quad (t, x) \in \mathbb{R}^2,$$

$$\phi_{0,\pm} = \pm\kappa_1$$

is an $(2N - 1)$ -soliton solution for the mKdV-equation (7.2). Moreover, up to an overall sign, the solutions (8.20) represent all reflectionless potentials of the associated Dirac operator $Q(t)$, $t \in \mathbb{R}$, under the assumption that $Q(0)$ has a zero-eigenvalue, i.e., $0 \in \sigma_d(Q(0))$.

Proof. Clearly (7.56), (7.63), and (8.10) imply (8.20). The fact that the $\phi_0(t, x)$ in (8.20) are reflectionless potentials of $Q_m(t)$, $t \in \mathbb{R}$, is a consequence of (4.11) and §6. Since $\phi_{0,\pm} = \pm\kappa_1 \gtrless 0$,

$$(8.21) \quad 0 \in \sigma_p(A_0(t)), \quad 0 \notin \sigma_p(A_0(t)^*), \quad t \in \mathbb{R},$$

implying (cf. (7.57))

$$(8.22) \quad 0 \in \sigma_d(Q(t)), \quad t \in \mathbb{R}.$$

Conversely, if $0 \in \sigma_p(Q(0))$ and, e.g., $\phi_{0,-} < 0 < \phi_{0,+}$, then (8.22) holds and by (4.11), $V_{j,(0)}(t, x)$ must be reflectionless potentials for $H_{j,(0)}(t)$, $j = 1, 2$, of the type (8.19). \square

Remark 8.2. The symmetry $\phi(t, x) \rightarrow -\phi(t, x)$ of the mKdV-equation (7.2) is connected with the interchange $A(t) \rightarrow A(t)^*$, or, equivalently, with $H_1(t) \rightarrow H_2(t)$. This explains the open overall sign in the last part of Theorem 8.1.

It remains to treat the case where $0 \notin \sigma(Q(t))$, $t \in \mathbb{R}$, i.e., instead of assumption (8.9) we now assume that $H_1(0)$ is subcritical, and, hence,

$$(8.23) \quad \mathfrak{E}_1 > 0 \quad (\text{i.e., } \mathfrak{E}_1 = \mathfrak{E}_2 > 0).$$

Then (8.1)–(8.8) and (8.11), (8.12) still hold whereas (8.10) turns into

$$(8.24) \quad V_\infty = \lambda_{1,1} + \kappa_1^2 = \mathfrak{E}_1 + \kappa_1^2.$$

Combining Remarks 7.15 and 8.2, we may restrict ourselves to $\sigma = \pm 1$ in (7.55) and thus (cf. (7.56))

$$(8.25) \quad \phi_{\sigma,\pm} = \lim_{x \rightarrow \pm\infty} \phi_\sigma(t, x) = \begin{cases} V_\infty^{1/2}, & \sigma = +1, \\ -V_\infty^{1/2}, & \sigma = -1, \end{cases}$$

implies

$$(8.26) \quad \sigma(H_1(t)) = \sigma(H_{2,\sigma}(t)) = \sigma(H_1(0)), \quad t \in \mathbb{R},$$

and also the validity of (8.15). It remains to compute the norming constants. We get

$$(8.27) \quad c_{2,\pm,\sigma,l}(0) = \left(\frac{\sigma V_\infty^{1/2} \mp \kappa_l}{\sigma V_\infty^{1/2} \pm \kappa_l} \right)^{1/2} c_{1,\pm,l}(0), \quad l = 1, \dots, N, \sigma = \pm 1,$$

and thus $V_{2,\sigma}(t, x)$ is the N -soliton potential

$$(8.28) \quad V_{2,\sigma}(t, x) = V_\infty - 2\partial_x^2 \ln\{\det[1 + C_{2,\pm,\sigma}(t, x)]\},$$

$$\sigma = \pm 1, \quad (t, x) \in \mathbb{R}^2,$$

where

$$(8.29) \quad C_{2,\pm,\sigma}(t, x) = [c_{2,\pm,\sigma,l,m}(t, x)]_{l,m=1}^N, \quad N \in \mathbb{N},$$

$$c_{2,\pm,\sigma,l,m}(t, x) = \left[\frac{(\sigma V_\infty^{1/2} \mp \kappa_l)(\sigma V_\infty^{1/2} \mp \kappa_m)}{(\sigma V_\infty^{1/2} \pm \kappa_l)(\sigma V_\infty^{1/2} \pm \kappa_m)} \right]^{1/2} c_{1,\pm,l,m}(t, x),$$

$$\sigma = \pm 1, \quad l, m = 1, \dots, N, \quad (t, x) \in \mathbb{R}^2.$$

Finally we compute ϕ_σ again.

Theorem 8.3. Assume $\mathfrak{E}_1 > 0$ and let

$$(8.30) \quad V_{j,(\sigma)}(t, x) = V_\infty - 2\partial_x^2 \ln\{\det[1 + C_{j,\pm,(\sigma)}(t, x)]\},$$

$(t, x) \in \mathbb{R}^2$, $j = 1, 2$, $\sigma = \pm 1$ (cf. (8.1) and (8.28)), be N -soliton solutions of the KdV-equation (7.1), $N \in \mathbb{N}$. Let V_1 , $V_{2,\sigma}$, and ϕ_σ , $\sigma = \pm 1$, be connected by the Miura transformation (7.3). Then

$$(8.31) \quad \begin{aligned} \phi_\sigma(t, x) &= \sigma V_\infty^{1/2} + \partial_x \ln\{\det[1 + C_{1,\pm}(t, x)] / \det[1 + C_{2,\pm,\sigma}(t, x)]\}, \\ \sigma &= \pm 1, (t, x) \in \mathbb{R}^2, \\ \phi_{\sigma,\pm} &= \sigma V_\infty^{1/2}, \quad \sigma = \pm 1 \end{aligned}$$

is a $2N$ -soliton solution of the mKdV-equation (7.2). Moreover, up to an overall sign, the solutions (8.31) represent all reflectionless potentials of the associated Dirac operator $Q(t)$, $t \in \mathbb{R}$, under the assumption that $Q(0)$ has no zero-eigenvalue, i.e., $0 \notin \sigma_d(Q(0))$.

Proof. Similar to that of Theorem 8.1. \square

Example 8.4. $N = 1$.

(i) $\mathfrak{E}_1 = \inf[\sigma(H_1(0))] = 0$.

$$(8.32) \quad \begin{aligned} V_1(t, x) &= \kappa_1^2 - 2\kappa_1^2 \cosh^{-2}(\kappa_1 x + 2\kappa_1^3 t), \\ V_\infty &= \kappa_1^2, \quad c_{1,\pm,1}(0)^2 = 2\kappa_1, \\ V_{2,0}(t, x) &= \kappa_1^2, \\ \phi_0(t, x) &= \kappa_1 \tanh(\kappa_1 x + 2\kappa_1^3 t), \quad \kappa_1 > 0, (t, x) \in \mathbb{R}^2. \end{aligned}$$

(ii) $\mathfrak{E}_1 = \inf[\sigma(H_1(0))] > 0$.

$$(8.33) \quad \begin{aligned} V_1(t, x) &= (\mathfrak{E}_1 + \kappa_1^2) - 2\kappa_1^2 \cosh^{-2}(\kappa_1 x + (2\kappa_1^3 + 6\kappa_1 \mathfrak{E}_1)t), \\ V_\infty &= \mathfrak{E}_1 + \kappa_1^2, \quad c_{1,\pm,1}(0)^2 = 2\kappa_1, \\ V_{2,\sigma}(t, x) &= (\mathfrak{E}_1 + \kappa_1^2) - 2\kappa_1^2 \cosh^{-2}(\kappa_1 x + (2\kappa_1^3 + 6\kappa_1 \mathfrak{E}_1)t + \sigma \tfrac{1}{2} \ln \gamma_1), \\ \gamma_1 &:= [(\mathfrak{E}_1 + \kappa_1^2)^{1/2} + \kappa_1] / [(\mathfrak{E}_1 + \kappa_1^2)^{1/2} - \kappa_1], \quad \sigma = \pm 1, \\ \phi_\sigma(t, x) &= \sigma(\mathfrak{E}_1 + \kappa_1^2)^{1/2} - \sigma \kappa_1^2 \cosh^{-2}(\kappa_1 x + (2\kappa_1^3 + 6\kappa_1 \mathfrak{E}_1)t) \\ &\quad \cdot [(\mathfrak{E}_1 + \kappa_1^2)^{1/2} + \sigma \kappa_1 \tanh(\kappa_1 x + (2\kappa_1^3 + 6\kappa_1 \mathfrak{E}_1)t)]^{-1}, \\ &\quad \kappa_1 > 0, (t, x) \in \mathbb{R}^2, \sigma = \pm 1. \end{aligned}$$

9. THE TWO-ZONE MODEL

In this section we investigate the simplest (nontrivial) finite zone model and illustrate the results of §7. As a by-product we obtain a generalization of a

well-known result of Hochstadt's [59] in Theorem 9.3. Let

$$(9.1) \quad V_1(t, x) = 2 \wp(x + 6 \wp(\omega)t + \omega') + \wp(\omega), \quad (t, x) \in \mathbb{R}^2,$$

where $\wp(z) := \wp(z; \omega; \omega')$ denotes the Weierstrass \wp -function [5] with real half-period $\omega > 0$ and purely imaginary half-period ω' ($\text{Im } \omega' > 0$). As is well known [21, 35, 36, 40, 61, 78, 84, 92, 93], V_1 satisfies the KdV-equation (7.1), i.e.,

$$(9.2) \quad \text{KdV}(V_1) = 0.$$

The corresponding periodic Schrödinger operator

$$(9.3) \quad H_1(t) = -\partial_x^2 + V_1(t, \cdot), \quad \mathfrak{D}(H_1(t)) = H^2(\mathbb{R}), \quad t \in \mathbb{R}$$

(with period $a = 2\omega$), has the absolutely continuous band spectrum [76, 93]

$$(9.4) \quad \begin{aligned} \sigma(H_1(t)) &= [0, E_1] \cup [E_2, \infty), \quad t \in \mathbb{R}, \\ E_1 &= \wp(\omega) - \wp(\omega + \omega'), \quad E_2 = \wp(\omega) - \wp(\omega'), \\ \sigma_p(H_1(t)) &= \sigma_{\text{sc}}(H_1(t)) = \emptyset, \quad t \in \mathbb{R}. \end{aligned}$$

The associated normalized Floquet solutions read [76]

$$(9.5) \quad \begin{aligned} f_{1,\pm}(k, t, x) &= \pm [\omega^{-1} \zeta(\omega) + \wp(b)]^{-1/2} \frac{\sigma(x + 6 \wp(\omega)t + \omega' \pm b)}{\sigma(x + 6 \wp(\omega)t + \omega') \sigma(\pm b)} \\ &\quad \cdot e^{\mp \{ [x + 6 \wp(\omega)t] \zeta(b) + \zeta(\omega') b \}} \\ &\quad \text{for } \lambda \in [0, E_1], \quad -i(b - \omega) \in [0, -i\omega'], \\ f_{1,\pm}(k, t, x) &= [-\omega^{-1} \zeta(\omega) - \wp(b)]^{-1/2} \frac{\sigma(x + 6 \wp(\omega)t + \omega' \pm b)}{\sigma(x + 6 \wp(\omega)t + \omega') \sigma(\pm b)} \\ &\quad \cdot e^{\mp \{ [x + 6 \wp(\omega)t] \zeta(b) + \zeta(\omega') b \}} \\ &\quad \text{for } \lambda \in [E_2, \infty), \quad -ib \in [-i\omega', 0), \end{aligned}$$

where the Floquet parameter k and the energy λ in terms of the parameter b are given by

$$(9.6) \quad \begin{aligned} k(b) &= i \zeta(b) - i \omega^{-1} \zeta(\omega) b, \\ \lambda(b) &= \wp(\omega) - \wp(b). \end{aligned}$$

Here b runs counterclockwise through the perimeter of the fundamental quarter rectangle with vertices $0, \omega, \omega + \omega', \omega'$, and $\zeta(z) := \zeta(z; \omega, \omega')$ and $\sigma(z) := \sigma(z; \omega, \omega')$ are the Weierstrass ζ - and σ -functions respectively [5].

At zero-energy we have

$$(9.7) \quad \begin{aligned} H_1(t) \psi_{1,0}(t) &= 0, \\ 0 < \psi_{1,0}(t, x) &= \frac{\sigma(x + 6 \wp(\omega)t + \omega' + \omega)}{\sigma(x + 6 \wp(\omega)t + \omega') \sigma(\omega)} \\ &\quad \cdot e^{-\{ [x + 6 \wp(\omega)t] \zeta(\omega) + \zeta(\omega') \omega \}}, \quad (t, x) \in \mathbb{R}^2, \end{aligned}$$

and thus (cf. (7.66) and Theorem 7.9) we get

Theorem 9.1. *Let*

$$(9.8) \quad V_1(t, x) = 2\mathfrak{P}(x + 6\mathfrak{P}(\omega)t + \omega') + \mathfrak{P}(\omega), \quad (t, x) \in \mathbb{R}^2,$$

be a two-zone solution of the KdV-equation (7.1) with $\mathfrak{E}_1 = 0$. Let $V_1, V_{2,0}$, and ϕ_0 be connected by Miura's transformation (7.3). Then

$$(9.9) \quad \begin{aligned} \phi_0(t, x) &= -\psi_{1,0,x}(t, x)/\psi_{1,0}(t, x) \\ &= -\frac{1}{2} \frac{\mathfrak{P}'(x + 6\mathfrak{P}(\omega)t + \omega')}{\mathfrak{P}(x + 6\mathfrak{P}(\omega)t + \omega') - \mathfrak{P}(\omega)} \\ &= [\zeta(x + 6\mathfrak{P}(\omega)t + \omega') + \zeta(\omega) - \zeta(x + 6\mathfrak{P}(\omega)t + \omega + \omega')], \\ &\quad (t, x) \in \mathbb{R}^2, \end{aligned}$$

satisfies the mKdV-equation (7.2), i.e.,

$$(9.10) \quad \text{mKdV}(\phi_0) = 0$$

and

$$(9.11) \quad \begin{aligned} V_{2,0}(t, x) &= \phi_0(t, x)^2 + \phi_{0,x}(t, x) \\ &= 2\mathfrak{P}(x + 6\mathfrak{P}(\omega)t + \omega + \omega') + \mathfrak{P}(\omega), \quad (t, x) \in \mathbb{R}^2, \end{aligned}$$

satisfies the KdV-equation (7.1), i.e.,

$$(9.12) \quad \text{KdV}(V_{2,0}) = 0.$$

The fact that $V_{2,0}$ is just a translate of V_1 , i.e.,

$$(9.13) \quad V_{2,0}(t, x) = V_1(t, x + \omega), \quad (t, x) \in \mathbb{R}^2,$$

is no accident (see [116]). In fact one has

Theorem 9.2 [40, 52, 59]. *Let $V \in L^1_{\text{loc}}(\mathbb{R})$ be real-valued and periodic with period $2\omega > 0$ and define the form sum $H = -d^2/dx^2 + V$ in $L^2(\mathbb{R})$. Then*

(i) *H has spectrum*

$$(9.14) \quad \sigma(H) = [E_0, \infty), \quad E_0 \in \mathbb{R}, \quad \text{iff } V = E_0 \text{ a.e.}$$

(ii) *H has spectrum*

$$(9.15) \quad \sigma(H) = [E_0, E_1] \cup [E_2, \infty), \quad E_0 < E_1 < E_2,$$

iff

$$V = 2\mathfrak{P}(\cdot + \omega' + \alpha) + C \text{ a.e.}, \quad C \in \mathbb{R},$$

where

$$(9.16) \quad \alpha \in \mathbb{R}, \quad E_0 = C - \mathfrak{P}(\omega), \quad E_1 = C - \mathfrak{P}(\omega + \omega'), \quad E_2 = C - \mathfrak{P}(\omega').$$

By commutation (cf. Theorem 4.4) this immediately yields

Theorem 9.3. *Let $\phi \in AC_{\text{loc}}(\mathbb{R})$ be real-valued and periodic with period $2\omega > 0$ and define in $L^2(\mathbb{R})$, $A = d/dx + \phi$, $\mathfrak{D}(A) = H^1(\mathbb{R})$. Let $Q_m = \begin{pmatrix} m & A^* \\ A & -m \end{pmatrix}$, $m \in \mathbb{R}$, be the corresponding Dirac operator in $L^2(\mathbb{R})^2$. Then*

(i) Q_m has spectrum

$$(9.17) \quad \sigma(Q_m) = (-\infty, -|m^2 + \phi_0^2|^{1/2}] \cup [|m^2 + \phi_0^2|^{1/2}, \infty), \quad \phi_0 \in \mathbb{R}, \\ \text{iff } \phi = \pm \phi_0.$$

(ii) Q_m has spectrum

$$(9.18) \quad \sigma(Q_m) = (-\infty, -|m^2 + E_2|^{1/2}] \cup [-|m^2 + E_1|^{1/2}, -|m^2 + E_0|^{1/2}] \\ \cup [|m^2 + E_0|^{1/2}, |m^2 + E_1|^{1/2}] \cup [|m^2 + E_2|^{1/2}, \infty), \\ 0 \leq E_0 < E_1 < E_2,$$

iff

$$\phi = \pm \frac{1}{2} \frac{\mathfrak{P}'(\cdot + \omega' + \alpha) - \varepsilon \mathfrak{P}'(b_0)}{\mathfrak{P}(\cdot + \omega' + \alpha) - \mathfrak{P}(b_0)},$$

where

$$(9.19) \quad \alpha \in \mathbb{R}, \quad \varepsilon = \pm 1, \quad E_0 = \mathfrak{P}(b_0) - \mathfrak{P}(\omega), \\ E_1 = \mathfrak{P}(b_0) - \mathfrak{P}(\omega + \omega'), \quad E_2 = \mathfrak{P}(b_0) - \mathfrak{P}(\omega').$$

Interestingly enough, the potential ϕ in (9.18) appears in a recently discovered new class of integrable systems [97] (see also [18]) being generalizations of the Calogero-Moser systems. Periodic, finite-zone Dirac operators have also been studied in [60].

Remark 9.4. Quite generally, one can ask which differential equations do the travelling wave solutions of the KdV-equation, respectively mKdV-equation, $V_j(x - \nu t)$, $j = 1, 2$, resp. $\phi(x - \nu t)$, $(t, x, \nu) \in \mathbb{R}^3$, connected via Miura's transformation (7.3), solve. The answer is given by

$$(9.20) \quad V_j'^2 = 2V_j^3 + \nu V_j^2 + aV_j + b, \quad j = 1, 2,$$

$$(9.21) \quad \phi'^2 = \phi^4 + \nu \phi^2 + c\phi + d,$$

where

$$(9.22) \quad c^2 = 4b - 2\nu a, \quad d = -a/2.$$

The sign ambiguity of c in (9.22) has to do with the symmetry $\phi \rightarrow -\phi$ of the mKdV-equation and a, b, ν need to satisfy additional constraints (e.g., $4b - 2\nu a \geq 0$ to guarantee real-valued solutions). The explicit example (9.8), (9.9), (9.12) leads to

$$(9.23) \quad \nu = -2(E_1 + E_2) = -6\mathfrak{P}(\omega), \\ a = -2(E_2 - E_1)^2 = -2[\mathfrak{P}(\omega + \omega') - \mathfrak{P}(\omega')]^2, \\ b = 2(E_1 + E_2)(E_2 - E_1)^2 = 6\mathfrak{P}(\omega)[\mathfrak{P}(\omega + \omega') - \mathfrak{P}(\omega')]^2, \\ c = 0, \quad d = (E_2 - E_1)^2 = [\mathfrak{P}(\omega + \omega') - \mathfrak{P}(\omega')]^2.$$

10. SOLITONS ON THE TWO-ZONE PERIODIC BACKGROUND

Here we continue the study of the two-zone model of §9 and construct solitons relative to this background. These KdV-solutions are then transferred to the mKdV-case. Let

$$(10.1) \quad V_1(t, x) = 2 \mathfrak{P}(x - \nu t + \omega') - (\nu/6), \quad (t, x, \nu) \in \mathbb{R}^3,$$

$$(10.2) \quad H_1(t) = -\partial_x^2 + V_1(t, \cdot), \quad \mathcal{D}(H_1(t)) = H^2(\mathbb{R}), \quad t \in \mathbb{R}.$$

Then [76, 93]

$$(10.3) \quad \text{KdV}(V_1) = 0$$

and

$$(10.4) \quad \begin{aligned} \sigma(H_1(t)) &= [E_0, E_1] \cup [E_2, \infty), \quad t \in \mathbb{R}, \\ E_0 &= -\mathfrak{P}(\omega) - (\nu/6), \quad E_1 = -\mathfrak{P}(\omega + \omega') - (\nu/6), \\ E_2 &= -\mathfrak{P}(\omega') - (\nu/6), \quad \sigma_p(H_1(t)) = \sigma_{\text{sc}}(H_1(t)) = \emptyset, \quad t \in \mathbb{R}. \end{aligned}$$

We aim to construct reflectionless potentials $W_{1,+}$, resp. $W_{1,-}$ (i.e., solitons), relative to the background V_1 that vanish as $x \rightarrow +\infty$, resp. as $x \rightarrow -\infty$. Applying the Marchenko method [80] this is accomplished as follows (see §6): Assuming

$$(10.5) \quad \nu \leq -6 \mathfrak{P}(\omega)$$

in order to guarantee

$$(10.6) \quad H_1(t) \geq 0,$$

we suppose $W_{1,+}$, resp. $W_{1,-}$, to support N bound states at energies

$$(10.7) \quad 0 \leq \tilde{\lambda}_{1,1} < \tilde{\lambda}_{1,2} < \dots < \tilde{\lambda}_{1,N}, \quad N \in \mathbb{N},$$

where $\tilde{\lambda}_{1,l} := \lambda(b_l)$, $l = 1, \dots, N$, and

$$(10.8) \quad \begin{aligned} 0 \leq \tilde{\lambda}_{1,m} < E_0 \quad (\text{i.e. } \mathfrak{P}(\omega) < \mathfrak{P}(b_m) \leq -\nu/6), \\ m = 1, \dots, N_1, \quad N_1 \leq N \text{ if } \nu < -6\mathfrak{P}(\omega) \end{aligned}$$

(otherwise (10.8) is empty),

$$(10.9) \quad \begin{aligned} E_1 < \tilde{\lambda}_{1,m} < E_2 \quad (\text{i.e. } \mathfrak{P}(\omega') < \mathfrak{P}(b_m) < \mathfrak{P}(\omega + \omega')), \\ m = N_1 + 1, \dots, N. \end{aligned}$$

We also introduce $b_0 \in (0, \omega)$ such that

$$(10.10) \quad \lambda(b_0) = -\mathfrak{P}(b_0) - (\nu/6) = 0.$$

Clearly

$$(10.11) \quad \mathfrak{P}(b_l) \leq \mathfrak{P}(b_0), \quad l = 1, \dots, N.$$

In order to compute $W_{1,\pm}$ one then solves the Marchenko equation

$$(10.12) \quad K_{1,\pm}(t, x, x') + \Omega_{1,\pm}(t, x, x') \pm \int_x^{\pm\infty} dy K_{1,\pm}(t, x, y) \Omega_{1,\pm}(t, y, x') = 0, \quad t \in \mathbb{R}, \quad x \stackrel{<}{>} x',$$

with the ansatz

$$(10.13) \quad \Omega_{1,\pm}(t, x, x') = \sum_{l=1}^N \tilde{c}_{1,\pm,l}(t)^2 f_{1,\pm}(t, k_l, x) f_{1,\pm}(t, k_l, x'),$$

$$(t, x, x') \in \mathbb{R}^3,$$

$$(10.14) \quad f_{1,\pm}(t, k, x) = \frac{\sigma(x - \nu t + \omega' \pm b)}{\sigma(x - \nu t + \omega') \sigma(\pm b)} e^{\mp[\zeta(b)(x - \nu t) + \zeta(\omega')b]}$$

$$\cdot \begin{cases} \pm[\omega^{-1}\zeta(\omega) + \mathfrak{P}(b)]^{-1/2}, & \lambda \in (-\infty, \lambda(\bar{b})), \\ [-\omega^{-1}\zeta(\omega) - \mathfrak{P}(b)]^{-1/2}, & \lambda \in (\lambda(\bar{b}), \infty), \end{cases} \quad (t, x) \in \mathbb{R}^2$$

(cf. Remark 10.2 for a discussion of \bar{b}),

$$(10.15) \quad k(b) = i\zeta(b) - i\omega^{-1}\zeta(\omega)b,$$

$$k_l := k(b_l), \quad \lambda(b) = -\mathfrak{P}(b) - (\nu/6),$$

$$(10.16) \quad \tilde{c}_{1,\pm,l}(t)^2 = \tilde{c}_{1,\pm,l}(0)^2 e^{\mp 4\mathfrak{P}'(b_l)t}, \quad t \in \mathbb{R}, \quad l = 1, \dots, N.$$

This yields

$$(10.17) \quad K_{1,\pm}(t, x, x') = - \sum_{l=1}^N \tilde{c}_{1,\pm,l}(t)^2 \tilde{f}_{1,\pm}(t, k_l, x) f_{1,\pm}(t, k_l, x'),$$

where

$$(10.18) \quad \tilde{f}_{1,\pm}(t, k, x)$$

$$= f_{1,\pm}(t, k, x) \pm \int_x^{\pm\infty} dx' K_{1,\pm}(t, x, x') f_{1,\pm}(t, k, x'),$$

$$(t, x) \in \mathbb{R}^2,$$

$$(10.19) \quad \tilde{H}_{1,\pm}(t) \tilde{f}_{1,\pm}(t, k, x) = \lambda \tilde{f}_{1,\pm}(t, k, x), \quad t \in \mathbb{R} \quad (\text{distributional sense}),$$

$$\tilde{H}_{1,\pm}(t) \tilde{f}_{1,\pm,l}(t) = \tilde{\lambda}_{1,l} \tilde{f}_{1,\pm,l}(t),$$

$$\tilde{f}_{1,\pm,l}(t, \cdot) := \tilde{f}_{1,\pm}(t, k_l, \cdot) \in H^2(\mathbb{R}),$$

$$|\tilde{c}_{1,\pm,l}(0)| = \|\tilde{f}_{1,\pm,l}(0)\|_2^{-1}, \quad l = 1, \dots, N.$$

Here

$$(10.20) \quad \tilde{H}_{1,\pm}(t) = H_1(t) + W_{1,\pm}(t, \cdot), \quad \mathfrak{D}(\tilde{H}_{1,\pm}(t)) = H^2(\mathbb{R}), \quad t \in \mathbb{R},$$

with

$$(10.21) \quad \begin{aligned} W_{1,\pm}(t, x) &= \mp 2\partial_x K_{1,\pm}(t, x, x) \\ &= -2\partial_x^2 \ln\{\det[1 + \tilde{C}_{1,\pm}(t, x)]\}, \quad (t, x) \in \mathbb{R}^2, \end{aligned}$$

$$(10.22) \quad \begin{aligned} \tilde{C}_{1,\pm}(t, x) &= [\tilde{c}_{1,\pm,l,m}(t, x)]_{l,m=1}^N, \quad N \in \mathbb{N}, \\ \tilde{c}_{1,\pm,l,m}(t, x) &= \pm \tilde{c}_{1,\pm,l}(t) \tilde{c}_{1,\pm,m}(t) \\ &\quad \cdot \int_x^{\pm\infty} dx' f_{1,\pm}(t, k_l, x') f_{1,\pm}(t, k_m, x') \\ &= \tilde{c}_{1,\pm,l}(t) \tilde{c}_{1,\pm,m}(t) f_{1,\pm}(t, k_l, x) f_{1,\pm}(t, k_m, x) \\ &\quad \cdot \frac{\sigma(x - \nu t + \omega' \pm b_l \pm b_m) \sigma(x - \nu t + \omega) \sigma(b_l) \sigma(b_m)}{\sigma(x - \nu t + \omega' \pm b_l) \sigma(x - \nu t + \omega' \pm b_m) \sigma(b_l + b_m)}, \\ &\quad l, m = 1, \dots, N, \quad (t, x) \in \mathbb{R}^2. \end{aligned}$$

Finally, we note that

$$(10.23) \quad \text{KdV}(\tilde{V}_{1,\pm}) = 0, \quad \tilde{V}_{1,\pm} = V_1 + W_{1,\pm}.$$

In order to shorten the discussion in the following, we only discuss the special case

$$(10.24) \quad N = 1, \quad 0 \leq \tilde{\lambda}_{1,1} < E_0, \quad \nu < -6 \mathfrak{P}(\omega)$$

in some detail. Clearly (10.17) and (10.18) imply

$$(10.25) \quad \begin{aligned} \tilde{f}_{1,\pm}(t, k, x) &= \left[1 \pm \tilde{c}_{1,\pm,1}(t)^2 \int_x^{\pm\infty} dx' f_{1,\pm}(t, k_1, x')^2 \right]^{-1} \\ &\quad \cdot \left\{ f_{1,\pm}(t, k, x) \pm \tilde{c}_{1,\pm,1}(t)^2 \right. \\ &\quad \cdot \int_x^{\pm\infty} dx' [f_{1,\pm}(t, k, x) f_{1,\pm}(t, k_1, x') \\ &\quad \left. - f_{1,\pm}(t, k_1, x) f_{1,\pm}(t, k, x')] f_{1,\pm}(t, k_1, x') \right\}, \quad (t, x) \in \mathbb{R}^2. \end{aligned}$$

In particular,

$$(10.26) \quad \begin{aligned} \tilde{f}_{1,\pm}(t, k_1, x) &= \left[1 \pm \tilde{c}_{1,\pm,1}(t)^2 \int_x^{\pm\infty} dx' f_{1,\pm}(t, k_1, x')^2 \right]^{-1} \\ &\quad \cdot f_{1,\pm}(t, k_1, x), \quad (t, x) \in \mathbb{R}^2, \end{aligned}$$

and thus

$$(10.27) \quad \tilde{f}_{1,\pm}(t, k_1, x) \Big|_{|x| \rightarrow \infty} = O(e^{-\text{Im } k_1 |x|}), \quad t \in \mathbb{R}.$$

Taking into account

$$\begin{aligned}
 & \pm \int_x^{\pm\infty} dx' f_{1,\pm}(t, k, x')^2 \\
 &= f_{1,\pm}(t, k, x)^2 \frac{\sigma(x - \nu t + \omega' \pm 2b) \sigma(x - \nu t + \omega') \sigma(b)^2}{\sigma(x - \nu t + \omega' \pm b)^2 \sigma(2b)}, \\
 (10.28) \quad & \pm \int_x^{\pm\infty} dx' f_{1,\pm}(t, k, x') f_{1,\pm}(t, k_1, x') \\
 &= f_{1,\pm}(t, k, x) f_{1,\pm}(t, k_1, x) \\
 &\quad \cdot \frac{\sigma(x - \nu t + \omega' \pm b \pm b_1) \sigma(x - \nu t + \omega') \sigma(b) \sigma(b_1)}{\sigma(x - \nu t + \omega' \pm b) \sigma(x - \nu t + \omega' \pm b_1) \sigma(b + b_1)},
 \end{aligned}$$

we get in general

$$\begin{aligned}
 & \tilde{f}_{1,\pm}(t, k, x) \underset{x \rightarrow \pm\infty}{=} O(f_{1,\pm}(t, k, x)) \underset{x \rightarrow \pm\infty}{=} O(e^{\pm i k x}), \\
 (10.29) \quad & \tilde{f}_{1,\pm}(t, k, x) \underset{x \rightarrow \mp\infty}{=} 0(f_{1,\pm}(t, k, x)) \underset{x \rightarrow \mp\infty}{=} O(e^{\pm i k x}), \\
 & k \neq k_1, \quad t \in \mathbb{R}.
 \end{aligned}$$

Similarly,

$$\begin{aligned}
 & W_{1,\pm}(t, x) \underset{x \rightarrow \pm\infty}{=} O(f_{1,\pm}(t, k_1, x)^2) \underset{x \rightarrow \pm\infty}{=} O(e^{\mp 2 \operatorname{Im} k_1 x}), \\
 (10.30) \quad & W_{1,\pm}(t, x) \underset{x \rightarrow \mp\infty}{=} 2 \left[\Re(x - \nu t + \omega' \pm 2b_1) - \Re(x - \nu t + \omega') \right] \\
 & + O(e^{\pm 2 \operatorname{Im} k_1 x}), \quad t \in \mathbb{R}.
 \end{aligned}$$

Remark 10.1. We emphasize the fact that $W_{1,+} \neq W_{1,-}$ and that [76]

$$(10.31) \quad W_{1,\pm}(t, x) \xrightarrow{x \rightarrow \pm\infty} 0, \quad t \in \mathbb{R},$$

but

$$(10.32) \quad W_{1,\pm}(t, x) \xrightarrow{x \rightarrow \mp\infty} V_1(t, x \pm 2b_1) - V_1(t, x) \neq 0, \quad t \in \mathbb{R}.$$

Thus $W_{1,\pm}(t, x)$ are not decaying as $x \rightarrow \mp\infty$ (as has also been observed in [41]) and hence $W_{1,\pm}(t)$ is no relatively compact perturbation of $H_1(t)$, $t \in \mathbb{R}$. However, since, up to exponentially decreasing terms, $V_1(t, x) + W_{1,\pm}(t, x)$ is just a translation ($x \rightarrow x \pm 2b_1$) of $V_1(t, x)$ as $x \rightarrow \mp\infty$, the band spectrum of $H_1(t)$ stays invariant under addition of $W_{1,\pm}(t)$, i.e.,

$$\sigma_{\text{ess}}(\tilde{H}_{1,\pm}(t)) = \sigma(H_1(t)) = [E_0, E_1] \cup [E_2, \infty), \quad t \in \mathbb{R}.$$

Remark 10.2. The special case where $b_{l_0} = \bar{b}$ for some $l_0 \in \{1, \dots, N\}$ in (10.21) with $\mathfrak{P}(\bar{b}) = (2\omega)^{-1} \int_0^{2\omega} dx \mathfrak{P}(x + \omega')$ (i.e., $(\bar{b} - \omega') \in (0, \omega)$) needs a straightforward limiting argument since

$$(10.33) \quad F_z(\lambda_1) = 0, \quad \lambda_1 := -\mathfrak{P}(\bar{b}) - (\nu/6) \in (E_1, E_2)$$

(cf. (3.33)), and hence $f_{1,\pm}(t, k(b), x)$ blows up at $b = \bar{b}$.

The results discussed so far are due to [76] (see also [73]). Next we transfer them to the mKdV-case. For simplicity we restrict ourselves to a treatment of $\tilde{H}_{1,+}(0)$ only. We first discuss the case where $\tilde{H}_{1,+}(0)$ is critical, i.e., where

$$(10.34) \quad \tilde{\mathfrak{E}}_1 := \inf[\sigma(\tilde{H}_{1,+}(0))] = 0 \quad (\text{i.e. } \tilde{\lambda}_{1,1} = 0).$$

Then we define

$$(10.35) \quad \begin{aligned} \tilde{\phi}_0(t, x) &:= -\tilde{f}_{1,+,x}(t, k_0, x)/\tilde{f}_{1,+}(t, k_0, x) \\ &= \phi_+(t, x) - \left[1 + \tilde{c}_{1,+,1}(t)^2 \int_x^{+\infty} dx' f_{1,+}(t, k_0, x')^2 \right]^{-1} \\ &\quad \cdot \tilde{c}_{1,+,1}(t)^2 f_{1,+}(t, k_0, x)^2, \quad k_0 := k(b_0), \quad (t, x) \in \mathbb{R}^2 \end{aligned}$$

(cf. (10.28)), with $\phi_+(t, x)$ given by (cf. (9.9))

$$(10.36) \quad \begin{aligned} \phi_+(t, x) &= -f_{1,+,x}(t, k_0, x)/f_{1,+}(t, k_0, x) \\ &= -\frac{1}{2} \frac{\mathfrak{P}'(x - \nu t + \omega') - \mathfrak{P}'(b_0)}{\mathfrak{P}(x - \nu t + \omega') - \mathfrak{P}(b_0)} \\ &= \left[\zeta(x - \nu t + \omega') + \zeta(b_0) - \zeta(x - \nu t + \omega' + b_0) \right], \quad (t, x) \in \mathbb{R}^2. \end{aligned}$$

Thus

$$(10.37) \quad \int_{x_0}^x dx' \tilde{\phi}_0(t, x') \underset{x \rightarrow \pm\infty}{=} \pm \operatorname{Im} k_0 x + O(1), \quad t \in \mathbb{R},$$

since $k_0 = k_1$ by assumption. Similarly, in the case where $\tilde{H}_{1,+}(0)$ is subcritical, i.e., where

$$(10.38) \quad \tilde{\mathfrak{E}}_1 > 0 \quad (\text{i.e. } 0 < \tilde{\lambda}_{1,1} < E_0),$$

we define

(10.39)

$$\begin{aligned}
 \tilde{\phi}_{\pm}(t, x) &:= -\tilde{f}_{1, \pm, x}(t, k_0, x)/\tilde{f}_{1, \pm}(t, k_0, x) \\
 &= -\left\{f_{1, \pm}(t, k_0, x) \pm \tilde{c}_{1, \pm, 1}(t)^2 \right. \\
 &\quad \cdot \int_x^{\pm\infty} dx' \left[f_{1, \pm}(t, k_0, x) f_{1, \pm}(t, k_1, x') - f_{1, \pm}(t, k_1, x) f_{1, \pm}(t, k_0, x') \right] \\
 &\quad \cdot f_{1, \pm}(t, k_1, x') \Big\}^{-1} \\
 &\quad \cdot \left\{ f_{1, \pm, x}(t, k_0, x) \pm \tilde{c}_{1, \pm, 1}(t)^2 \int_x^{\pm\infty} dx' [f_{1, \pm, x}(t, k_0, x) f_{1, \pm}(t, k_1, x') \right. \\
 &\quad \left. - f_{1, \pm, x}(t, k_1, x) f_{1, \pm}(t, k_0, x')] f_{1, \pm}(t, k_1, x') \right\} \\
 &\quad \mp \left[1 \pm \tilde{c}_{1, \pm, 1}(t)^2 \int_x^{\pm\infty} dx' f_{1, \pm}(t, k_1, x')^2 \right]^{-1} \\
 &\quad \cdot \tilde{c}_{1, \pm, 1}(t)^2 f_{1, \pm}(t, k_1, x)^2, \quad (t, x) \in \mathbb{R}^2
 \end{aligned}$$

(cf. (10.28)). Thus

$$(10.40) \quad \int_{x_0}^x dx' \tilde{\phi}_{\pm 1}(t, x') \Big|_{x \rightarrow \pm\infty} = \pm \operatorname{Im} k_0 x + O(1), \quad t \in \mathbb{R}.$$

Theorem 7.9 now implies

Theorem 10.3. *Let $\tilde{\phi}_{\sigma}$, $\sigma = 0, \pm 1$, be defined by (10.35), resp. (10.39). Then*

$$(10.41) \quad \operatorname{KdV}(\tilde{V}_{2, \sigma}) = 0, \quad \operatorname{mKdV}(\tilde{\phi}_{\sigma}) = 0, \quad \sigma = 0, \pm 1,$$

where

$$(10.42) \quad \tilde{V}_{2, \sigma}(t, x) := \tilde{\phi}_{\sigma}(t, x)^2 + \tilde{\phi}_{\sigma, x}(t, x), \quad \sigma = 0, \pm 1, \quad (t, x) \in \mathbb{R}^2.$$

Remark 10.4. For $\sigma = 0$ one can actually show that

$$(10.43) \quad \tilde{V}_{2, 0}(t, x) = V_1(t, x + b_0), \quad (t, x) \in \mathbb{R}^2$$

(cf. also (9.13)).

11. SINGULAR SOLUTIONS OF THE KdV- AND mKdV EQUATIONS

In this section we indicate how to transfer singular solutions of the KdV-equation to singular solutions of the mKdV-equation using the methods developed in §7. We also derive rational solutions of the KdV- and mKdV-equations starting from a pure soliton solution and performing a long-wavelength limit.

We first introduce hypothesis

(H.11.1) Let $x_0 \in C^{\infty}(\mathbb{R})$ and $X_0 := \{(t, x_0(t)) \in \mathbb{R}^2 \mid t \in \mathbb{R}\}$.

(i) $V \in C^{\infty}(\mathbb{R}^2 \setminus X_0)$ real-valued, $\lim_{x \rightarrow \pm\infty} V(t, x) = V_{\infty}$, $0 < V_{\infty}$ (independent of t). For each $t \in \mathbb{R}$

$$\int_{|x - x_0(t)| \geq \varepsilon} dx (1 + x^2) |V(t, x) - V_{\infty}| < \infty$$

for all sufficiently small $\varepsilon > 0$. Furthermore, there exists an open superset U of X_0 in \mathbb{R}^2 and a real-valued function $v \in C^\infty(U)$, $\partial_x^n v \in L^\infty(U)$, $n = 0, 1$, such that $V(t, x) = 2[x - x_0(t)]^{-2} + v(t, x)$, $(t, x) \in U \setminus X_0$, $\partial_x v(t, x_0(t)) = 0$, $t \in \mathbb{R}$, and $\partial_x^n V \in L^\infty(\mathbb{R}^2 \setminus U)$, $n = 0, 1$.

(ii) $\phi \in C^\infty(\mathbb{R}^2 \setminus X_0)$ real-valued, $\lim_{x \rightarrow \pm\infty} \phi(t, x) = \phi_\pm$, $0 < \phi_-^2 = \phi_+^2$ (ϕ_\pm independent of t). For each, $t \in \mathbb{R}$ $\int_{|x-x_0(t)| \geq \varepsilon} dx(1+x^2)|\phi(t, x) - \phi_\pm| < \infty$ for all sufficiently small $\varepsilon > 0$. Furthermore, there is an open superset U of X_0 in \mathbb{R}^2 and a real-valued function $\varphi \in C^\infty(U)$, $\partial_x^n \varphi \in L^\infty(U)$, $n = 0, 1$, such that $\phi(t, x) = \pm[x - x_0(t)]^{-1} + \varphi(t, x)$, $(t, x) \in U \setminus X_0$, $\varphi(t, x_0(t)) = 0$, $t \in \mathbb{R}$, and $\partial_x^n \phi \in L^\infty(\mathbb{R}^2 \setminus U)$, $n = 0, 1$.

Remark 11.1. The behavior of V (resp. ϕ) near the singularity $x_0(t)$ results from the requirement that the Laurent expansion of V (resp. ϕ) is compatible with the time evolution described by the KdV-equation (resp. mKdV-equation).

Theorem 11.2. Let V_1 denote a solution of KdV which satisfies (H.7.2(i)) with $V_{1+} = V_{1-} = V_\infty > 0$. Assume $\mathfrak{E}_1 = \inf[\sigma(H_1)] > 0$ and define

$$(11.1) \quad g_\gamma(t, x) := f_{1,+}(t, x, i\beta) - \gamma e^{-4\beta t(2\beta^2 - 3V_\infty)} f_{1,-}(t, x, i\beta),$$

$$(t, x) \in \mathbb{R}^2, \gamma \in \mathbb{R},$$

where $f_{1,\pm}(t, x, k_-)$ are Jost solutions corresponding to V_1 and $\beta = \sqrt{V_\infty}$. If, for some $t_0 \in \mathbb{R}$, $g_\gamma(t_0, x)$ provides precisely one simple zero at $x = x_0(t_0)$, then

$$(11.2) \quad \phi_0(t, x) := -g_{\gamma,x}(t, x)/g_\gamma(t, x), \quad (t, x) \in \mathbb{R}^2 \setminus X_0,$$

is a singular solution of the mKdV-equation satisfying (H.11.1(ii)), and

$$(11.3) \quad V_2(t, x) := \phi_0(t, x)^2 + \phi_{0,x}(t, x), \quad (t, x) \in \mathbb{R}^2 \setminus X_0,$$

is a singular solution of the KdV-equation satisfying (H.11.1(i)). In both cases the path of the singularity $x_0(t)$ is given by the equation

$$(11.4) \quad g_\gamma(t, x_0(t)) = 0, \quad t \in \mathbb{R}.$$

Proof. First we note that

$$(11.5) \quad f_{1,\pm,t}(t, x, i\beta) = \{B_1(t) \mp 2\beta(2\beta^2 - 3V_\infty)\} f_{1,\pm}(t, x, i\beta), \quad (t, x) \in \mathbb{R}^2$$

(cf. (7.11) for the definition of $B_1(t)$). Therefore

$$(11.6) \quad \psi_1(t, x) := e^{2\beta t(2\beta^2 - 3V_\infty)} g_\gamma(t, x), \quad (t, x) \in \mathbb{R}^2,$$

is a distributional solution of $H_1(t)\psi_1(t) = 0$ evolving in time according to (7.21). From the assumption that $g_\gamma(t_0, x)$ provides precisely one zero at

$x = x_0(t_0)$ we conclude, using Lemma 7.7, that $g_\gamma(t, x)$, $t \in \mathbb{R}$, provides precisely one zero at $x = x_0(t)$ satisfying (7.32). Due to this fact we may write

$$(11.7) \quad \psi_1(t, x) = (x - x_0(t))\chi(t, (x - x_0(t))), \quad (t, x) \in \mathbb{R}^2.$$

Inserting (11.7) into (7.21b) and using (7.32) we derive

$$(11.8) \quad \chi_t(t, (x - x_0(t))) = -V_{1,x}(t, x)\chi(t, (x - x_0(t))),$$

and thus

$$(11.9) \quad \chi(t, (x - x_0(t))) = \chi(t_0, (x - x_0(t_0))) \cdot \exp \left[- \int_{t_0}^t V_{1,x}(t', x) dt' \right], \quad (t, x) \in \mathbb{R}^2.$$

Since $\chi(t, (x - x_0(t))) \neq 0$, $(t, x) \in \mathbb{R}^2$, the zero remains simple for arbitrary t . Using $\psi_{1,xx} = V_1\psi_1$, it is easily shown that V_2 (resp. ϕ_0) satisfies (H.11.1(i)) (resp. (H.11.1(ii))). Either by direct calculation, using (11.5), or by a generalization of Theorem 7.9, one finally proves that $\text{KdV}(V_2) = 0$ and $\text{mKdV}(\phi_0) = 0$. \square

In the following lemma we will exhibit a rich class of potentials V_1 with $g_\gamma(t_0, \cdot)$ providing precisely one simple zero:

Lemma 11.3. *Assume V_1 satisfies (H.7.2(i)) with $V_+ = V_- = V_\infty > 0$ and either $V_1(t_0, x) \geq V_\infty$, $x \in \mathbb{R}$, or V_1 supports precisely N discrete eigenvalues $\lambda_N < \lambda_{N-1} < \dots < \lambda_1$, $N \in \mathbb{N}$. Then for $0 < \lambda_N < V_\infty$ the function $g_\gamma(t_0, \cdot)$ defined in (11.1) with $\beta = \sqrt{V_\infty}$, $\gamma > 0$, has precisely one simple zero.*

Proof. In what follows we will neglect the argument t_0 . From the integral equation

$$(11.10) \quad f_{1,\pm}(x, i\beta) = e^{\mp\beta x} + \frac{1}{\beta} \int_x^{\pm\infty} dy \sinh(\beta(y-x)) [V_1(y) - V_\infty] f_{1,\pm}(y, i\beta),$$

$x \in \mathbb{R},$

we infer in case $V_1(x) \geq V_\infty$, $x \in \mathbb{R}$, that $f_{1,+}(x, i\beta)(f_{1,-}(x, i\beta))$ are positive, strictly monotonically decreasing (increasing) functions, continuously differentiable w.r.t. $x \in \mathbb{R}$. In case V_1 supports N bound states $\lambda_1, \dots, \lambda_N$, $N \geq 1$, let us start with $V_1(x; -1)$ (cf. Theorem 6.5) providing $(N-1)$ bound states $\lambda_1, \dots, \lambda_{N-1}$. Let $f_{1,\pm}(x, k_-; -1)$ denote the corresponding Jost solutions. Adding the N th bound state $V_\infty - \kappa_N^2 = \lambda_N$ to obtain V_1 the Jost solutions corresponding to V_1 may be written as

$$(11.11) \quad f_{1,\pm}(x, k_-) = \pm \frac{1}{ik_- - \kappa_N} \frac{1}{\tilde{g}_\alpha(x, i\kappa_N)} W(\tilde{g}_\alpha(x, i\kappa_N), f_{1,\pm}(x, k_-; -1)),$$

$x \in \mathbb{R},$

with

$$(11.12) \quad \tilde{g}_\alpha(x, k_-) = f_{1,+}(x, k_-; -1) + \alpha f_{1,-}(x, k_-; -1), \quad x \in \mathbb{R},$$

and $\alpha > 0$ suitably chosen [33]. For $\beta > \kappa_N$ and $\gamma > 0$ we thus obtain

(11.13)

$$\begin{aligned} g_\gamma(x) &= f_{1,+}(x, i\beta) - \gamma f_{1,-}(x, i\beta) = -\frac{1}{\kappa_n + \beta} \frac{1}{\tilde{g}_\alpha(x, i\kappa_n)} \\ &\quad \cdot W(\tilde{g}_\alpha(x, i\kappa_n), f_{1,+}(x, i\beta; -1) + \gamma f_{1,-}(x, i\beta; -1)) \\ &= (\beta - \kappa_N) \frac{1}{\tilde{g}_\alpha(x, i\kappa_N)} \left[\int_x^\infty \tilde{g}_\alpha(y, i\kappa_N) f_{1,+}(y, i\beta; -1) dy \right. \\ &\quad \left. - \gamma \int_{-\infty}^x \tilde{g}_\alpha(y, i\kappa_N) f_{1,-}(y, i\beta; -1) dy \right], \\ &\quad x \in \mathbb{R}. \end{aligned}$$

Again we observe that the right-hand side is the difference of a positive, strictly monotonically decreasing and a positive, strictly monotonically increasing function, which in addition is continuously differentiable w.r.t. $x \in \mathbb{R}$. Simplicity of the zero of g_γ is easily verified in both cases. \square

Lemma 11.4. Assume that V_1 is an N -soliton solution of the KdV-equation, i.e., that, according to (6.29), (6.30),

(11.14)

$$V_1(t, x) = V_\infty - 2\partial_x \left[\frac{\partial_x [W(a_1(t, x), \dots, a_N(t, x))]}{W(a_1(t, x), \dots, a_N(t, x))} \right], \quad (t, x) \in \mathbb{R}^2,$$

with

$$(11.15) \quad a_j(t, x) = (-1)^{j+1} e^{-\kappa_j x} + \alpha_j e^{-4\kappa_j t(2\kappa_j^2 - 3V_\infty)} e^{\kappa_j x},$$

$$j = 1, \dots, N, (t, x) \in \mathbb{R}^2,$$

and $\kappa_j = \sqrt{V_\infty - \lambda_j}$, α_j , $j = 1, \dots, N$, positive constants. Then the singular solutions of the KdV- and mKdV-equations obtained in Theorem 11.1 are given by

$$(11.16) \quad V_2(t, x) = V_\infty - 2\partial_x \left[\frac{\partial_x [W(a_1(t, x), \dots, a_N(t, x), a_{N+1}(t, x))]}{W(a_1(t, x), \dots, a_N(t, x), a_{N+1}(t, x))} \right],$$

$$(t, x) \in \mathbb{R}^2 \setminus X_0,$$

and

(11.17)

$$\begin{aligned} \phi_0(t, x) &= \frac{\partial_x [W(a_1(t, x), \dots, a_N(t, x))]}{W(a_1(t, x), \dots, a_N(t, x))} \\ &\quad - \frac{\partial_x [W(a_1(t, x), \dots, a_N(t, x), a_{N+1}(t, x))]}{W(a_1(t, x), \dots, a_N(t, x), a_{N+1}(t, x))}, \quad (t, x) \in \mathbb{R}^2 \setminus X_0, \end{aligned}$$

respectively, where a_{N+1} is defined by (11.15) with $\kappa_{N+1} = \beta = \sqrt{V_\infty}$ and $\alpha_{N+1} = -\gamma$. V_2 is an N -soliton plus one pole solution of the KdV-equation, ϕ_0 is a $2N$ -soliton plus one pole solution of the mKdV-equation.

Proof. The lemma is proved by using (11.2), (11.3), and (cf. [33])

$$(11.18) \quad f_{1,\pm}(t, x, i\beta) = (\mp 1)^N \frac{W(a_1, \dots, a_N, e^{\mp \beta x})}{\prod_{j=1}^N (\beta + \kappa_j) W(a_1, \dots, a_N)},$$

$$(t, x) \in \mathbb{R}^2. \quad \square$$

Remark 11.5. We note that the only difference between the N -soliton +1-pole solution (11.16) and an $(N+1)$ -soliton solution (cf. (11.14)) consists in the fact that α_{N+1} is negative in the first case, whereas it is positive in the second. By means of relation (6.33) the constants α_i are connected with the norming constants $c_{\pm,j}^2$. Thus, formally, adding a pole to a solution of KdV corresponds to introducing a “bound state” with negative norming constant $c_{\pm,N+1}^2$ in the associated isospectral Schrödinger operator. A more detailed account on the isospectral problem in connection with singular solutions of the KdV-equation may be found in [9].

Example 11.6. (i) 1-pole solution of the KdV- and mKdV-equations, $\beta > 0$, $\alpha_1 = -1$:

$$(11.19) \quad \begin{aligned} V_1(t, x) &= \beta^2, & (t, x) &\in \mathbb{R}^2, \\ V_2(t, x) &= \beta^2 + 2\beta^2 \sinh^{-2}(\beta x + 2\beta^3 t), & (t, x) &\in \mathbb{R}^2 \setminus X_0, \\ \phi_0(t, x) &= -\beta \coth(\beta x + 2\beta^3 t), & (t, x) &\in \mathbb{R}^2 \setminus X_0, \\ x_0(t) &= -2\beta^2 t, & t &\in \mathbb{R}. \end{aligned}$$

(ii) 1-soliton plus 1-pole solution of the KdV-equation $\beta > \kappa > 0$, $\alpha_1 = -\alpha_2 = 1$:

$$(11.20) \quad \begin{aligned} V_1(t, x) &= \beta^2 - 2\kappa^2 \cosh^{-2}(\kappa x - 2\kappa t(2\kappa^2 - 3\beta^2)), & (t, x) &\in \mathbb{R}^2, \\ V_2(t, x) &= \beta^2 - 4(\beta^2 - \kappa^2)[(\beta^2 - \kappa^2) + \beta^2 \cosh(2\kappa x - 4\kappa t(2\kappa^2 - 3\beta^2)) \\ &\quad - x^2 \cosh(2\beta x + 4\beta^3 t)] \\ &\quad \cdot [(\beta - \kappa) \sinh((\beta + \kappa)x - 4(\beta^3 + \kappa^3)t + 6\beta^2(\beta + \kappa)t) \\ &\quad + (\beta + \kappa) \sinh((\beta - \kappa)x - 4(\beta^3 - \kappa^3)t + 6\beta^2(\beta - \kappa)t)]^{-2}, \\ &\quad (t, x) \in \mathbb{R}^2 \setminus X_0, \\ \kappa \tanh(\kappa x_0(t) - 4\kappa^3 t + 6\kappa \beta^2 t) &= \beta \tanh(\beta x_0(t) + 2\beta^3 t), & t &\in \mathbb{R}. \end{aligned}$$

(iii) 2-soliton plus 1-pole solution of the mKdV-equation, $\beta > \kappa > 0$, $\alpha_1 = -\alpha_2 = 1$:
(11.21)

$$\begin{aligned} \phi_0(t, x) = & \kappa \tanh(\kappa x - 2\kappa t(2\kappa^2 - 3\beta^2)) \\ & - 2(\beta^2 - \kappa^2)[\cosh(\beta x + 2\beta^3 t) \cosh(\kappa x - 2\kappa t(2\kappa^2 - 3\beta^2))] \\ & \cdot [(\beta - \kappa) \sinh((\beta + \kappa)x - 4(\beta^3 + \kappa^3)t + 6\beta^2(\beta + \kappa)t) \\ & + (\beta + \kappa) \sinh((\beta - \kappa)x - 4(\beta^3 - \kappa^3)t + 6\beta^2(\beta - \kappa)t)]^{-1}, \\ & (t, x) \in \mathbb{R}^2 \setminus X_0, \quad x_0(\cdot) \text{ as given in (11.20)}. \end{aligned}$$

Remark 11.7. The singular solutions in (11.19) formally may be recovered from the one-soliton solutions of (8.32) by choosing the phase $\eta_1 = \frac{1}{2} \ln(c_{1,+}^2/2\kappa_1)$ to be $\eta_1 = i\pi/2$ instead of $\eta_1 = 0$. Similarly the singular solutions (11.20) and (11.21) formally could have been obtained from two- or three-soliton solutions, respectively, by appropriately choosing the phase constants. In general, the fact that one can recover singular solutions from regular ones relies on the observation that the norming constants $c_{j,\pm}^2$ merely enter as parameters, which essentially determine the phase of the j th soliton (see, e.g., [2]).

Exploiting this observation, we finally sketch a derivation of rational solutions of the KdV- and mKdV-equations. The simplest one is easily obtained from (11.19) by performing the “long-wave” limit $\beta \rightarrow 0$.

Example 11.8. Rational (stationary) solutions of the KdV- and mKdV-equations:

$$\begin{aligned} (11.22) \quad & V_1(t, x) = 0, \quad (t, x) \in \mathbb{R}^2, \\ & V_2(t, x) = \frac{2}{x^2}, \quad (t, x) \in \mathbb{R}^2 \setminus X_0, \\ & \phi_0(t, x) = -\frac{1}{x}, \quad (t, x) \in \mathbb{R}^2 \setminus X_0, \quad X_0 = \mathbb{R} \times \{0\}. \end{aligned}$$

In order to demonstrate how more general rational solutions of the KdV- and mKdV-equations may be obtained from the corresponding soliton solutions it suffices here to consider the two-soliton solution of the KdV-equation which we write in the form

$$(11.23) \quad V_1(t, x) = \kappa_2^2 - 2\partial_x^2 \{\ln[W(a_1(t, x), a_2(t, x))]\}, \quad (t, x) \in \mathbb{R}^2,$$

with $\kappa_2 > \kappa_1 > 0$. Here a_j , $j = 1, 2$, are defined via (11.15) with $V_\infty = \kappa_2^2$ and α_j are related to the norming constants $c_{j,+}$, $j = 1, 2$, via (6.33). The potential $V_{2,0}$ then is given by (see Theorem 8.1)

$$(11.24) \quad V_{2,0}(t, x) = \kappa_2^2 - 2\partial_x^2 \{\ln[a_1(t, x)]\}, \quad (t, x) \in \mathbb{R}^2,$$

and furthermore (cf. (7.63))

$$(11.25) \quad \phi_0(t, x) = \partial_x \{\ln[W(a_1(t, x), a_2(t, x))/a_1(t, x)]\}, \quad (t, x) \in \mathbb{R}^2.$$

Now taking $\alpha_j = (-1)^j$, $j = 1, 2$, and performing the “long-wave” limit $\kappa_1, \kappa_2 \rightarrow 0$ one obtains

Example 11.9.

$$\begin{aligned}
 V_1(t, x) &= -2\partial_x^2 \ln(x^3 + 12t) \\
 &= \frac{18x^4}{(x^3 + 12t)^2} - \frac{12x}{(x^3 + 12t)}, \quad (t, x) \in \mathbb{R}^2 \setminus X_0, \\
 V_2(t, x) &= -2\partial_x^2 \ln(x) = \frac{2}{x^2}, \quad (t, x) \in \mathbb{R}^2 \setminus \{\mathbb{R} \times \{0\}\}, \\
 \phi_0(t, x) &= \partial_x [\ln(x^3 + 12t) - \ln(x)] \\
 &= \frac{2x^3 - 12t}{x(x^3 + 12t)}, \quad (t, x) \in \mathbb{R}^2 \setminus X_0, \\
 x_0(t) &= -\operatorname{sgn}(t)|12t|^{1/3}, \quad t \in \mathbb{R} \ (\operatorname{sgn}(0) := 0).
 \end{aligned}
 \tag{11.26}$$

Whereas these rational solutions of the KdV-equation are well known [1, 2, 6, 22, 42, 53, 81] the corresponding solutions of the mKdV-equation seem to be new (although one could easily derive them from the results in [6]).

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