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# Application of the Darbous transformation to inverse problems with variable spectral parameters 

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Received 21 January 1991. in final form 11 June 1991


#### Abstract

We show that Darboux transformations can be used to find explicit solutions of deformed equations which have recently been found through the inverse method whth a variable spectral parameter. We construct in particular the soliton solution of the deformed Maxwell-Bloch equations, and analyse its asymptotic behaviour.


## 1. Introduction

The inverse method of solving nonlinear partal differential equations was generalized in [1] to the case of a variable spectral parameter. This extension of the usual inverse method gave rise [1] to a new series of completely integrable nomlinear equations, and to interesting integrable deformations of known nonlnear equationsł̣. Reference [1] did not provide, however, any practical method of integrating the new equations. In this article [4] we therefore propose a simple algebrace method of integration which is applicable to the new deformed equations. The method we use here is a modified version [4] of the Darboux transiomation (DT) method which has been successfully applied to a number of completely integrable equations [2-5, 9, 10]. An alternative method based on a finitely-many gaps infegration has been proposed in [6].

In the inverse method the nonlinear equation to be solved appears as the compatibility condition of an auxiliary linear problem

$$
\begin{equation*}
\psi_{z}=U \psi \quad \psi_{t}=V \psi \tag{11}
\end{equation*}
$$

where $\psi_{x} \equiv(\partial / \partial x) \psi$ etc, and matrices $U$ and $V$ have only simple poles in the complex $\lambda$-plane,

$$
\begin{align*}
& U(\lambda, x, t)=U_{0}(x, t)+\sum_{n=1}^{N_{1}} \frac{U_{n}(x, t)}{\lambda-\nu_{n}} \\
& V(\lambda, \vec{z}, t)=V_{0}(x, t)+\sum_{n=1}^{N_{2}} \frac{V_{n}(x, t)}{\lambda-\mu_{n}} \tag{12}
\end{align*}
$$

[^0]The spectral parameter $\lambda$ can now depend on space and time variables $x, t$, and on a hidden spectral parameter $z: \lambda=\lambda(x, t, z)$. Also the poles of the matrices $U$ and $V$ are dependent on $x$ and $t: \nu_{n}=\nu_{n}(x, t), \mu_{n}=\mu_{n}(x, t)$. The spectral pazameter $\lambda$ is found [1] to satisfy

$$
\begin{equation*}
\lambda_{t}=-\sum_{n=1}^{N_{2}} \frac{b_{n}}{\lambda-\mu_{n}} \quad \lambda_{x}=-\sum_{n=1}^{N_{1}} \frac{c_{n}}{\lambda-\nu_{n}} \tag{1.3}
\end{equation*}
$$

and therefore the poles $\nu_{n}, \mu_{n}$ satisfy [1]

$$
\begin{equation*}
\nu_{n_{t}}=-\sum_{m=1}^{N_{2}} \frac{b_{m}}{\nu_{n}-\mu_{m}} \quad \mu_{n_{s}}=\sum_{m=1}^{N_{1}} \frac{c_{n}}{\mu_{n}-\nu_{m}} \tag{1.4a}
\end{equation*}
$$

where
$c_{n_{t}}=2 c_{n} \sum_{m=1}^{N_{2}} \frac{b_{n n}}{\left(\nu_{n}-\mu_{m}\right)^{2}} b_{n_{x}}=2 b_{n} \sum_{m=1}^{N_{1}} \frac{c_{m}}{\left(\mu_{n}-\nu_{m}\right)^{2}}$.
The compatibility condation of system (1.1) is

$$
\begin{equation*}
U_{z}-V_{w}+[U, V]=0 \tag{1.5}
\end{equation*}
$$

We demand this condition to be satisfied identically for $\lambda$ and the linearly independent factors $\left(\lambda-\nu_{n}\right)^{-1}$ and $\left(\lambda-\mu_{n}\right)^{-1}$. Derivatives with respect to the spectral parameter in (1.5) will result in new terms as compared with those appearing in the usual inverse method.

We shall next formulate the DT method of integration. Let $\left(\psi_{3}, \phi_{3}\right)^{\mathrm{T}}$ be the solution of system (1.1) corresponding to the spectral parameter $\lambda_{\jmath}$ 可 $\lambda\left(x, t, z_{j}\right) ; \jmath=1,2$. If we define the matrices $\Lambda, \Phi, P_{n}$ and $Q_{n}$ such that

$$
\begin{align*}
& A=\left(\begin{array}{cc}
\lambda_{1} & 0 \\
0 & \lambda_{2}
\end{array}\right) \quad \Phi=\left(\begin{array}{ll}
\psi_{1} & \psi_{2} \\
\phi_{1} & \phi_{2}
\end{array}\right) \\
& P_{n}=\left(\begin{array}{cc}
\left(\lambda_{1}-\nu_{n}\right)^{-1} & 0 \\
0 & \left(\lambda_{2}-\nu_{n}\right)^{-1}
\end{array}\right) \quad Q_{n}=\left(\begin{array}{cc}
\left(\lambda_{1}-\mu_{n}\right)^{-1} & 0 \\
0 & \left(\lambda_{2}-\mu_{n}\right)^{-1}
\end{array}\right) \tag{1.6}
\end{align*}
$$

then the system (1.1) can be expressed in the form

$$
\begin{align*}
& \Phi_{z}=U_{0} \Phi+\sum_{n=1}^{N_{1}} U_{n} \Phi P_{n} \\
& \bar{\Phi}_{t}=V_{0} \Phi+\sum_{n=1}^{N_{i}} V_{n} \Phi Q_{n} \tag{1.7}
\end{align*}
$$

The DT method is based on the fact that system (1.7) is covariant with respect to linear transformations $\mathscr{C}_{\boldsymbol{\omega}}^{\boldsymbol{\omega}}=\boldsymbol{C}[\Phi]$. Also the transformed matrices $\tilde{U}_{n}, \tilde{V}_{m}$; $n=0,1, \ldots, N_{1}, m=0,1, \ldots, N_{2}$, will satisfy equation (1.5).

We consider a particularly simple ansatz for the linear transformation $\mathcal{C}$, for which

$$
\begin{equation*}
\overline{6}=S_{0} \Phi \bar{A}+S_{1} \Phi \tag{18}
\end{equation*}
$$

where matrices $S_{0}$ and $S_{1}$ act as free parameters. Under transformation (1.8) equations (1.7) should transform by covanance into

$$
\begin{equation*}
\tilde{\Phi}_{x}=\tilde{U}_{0} \tilde{\Phi}+\sum_{n=1}^{N_{1}} \tilde{U}_{n} \check{\Phi} P_{n} \tilde{\Phi}_{t}=\tilde{V}_{0} \tilde{\Phi}+\sum_{n=1}^{N_{2}} \tilde{V}_{n} \tilde{\Phi} Q_{n} \tag{1.9}
\end{equation*}
$$

By substituting (1.8) into (1.9) and demanding that the resulting equations should be satisfied identically for $\Phi \Lambda^{n}, \Phi P_{n}$ and $\Phi Q_{n}$, we arrive at equations

$$
\begin{align*}
& \sum_{n=1}^{N_{1}} S_{0} U_{n}+S_{1 z}+S_{1} U_{0}=\tilde{U}_{0} S_{1}+\sum_{n=1}^{N_{1}} \tilde{U}_{n} S_{0} \\
& S_{0 z}+S_{0} U_{0}=\tilde{U}_{0} S_{0}  \tag{110}\\
& \left(\nu_{n} S_{0}+S_{1}\right) U_{n}-c_{n} S_{0}=\tilde{U}_{n}\left(\nu_{n} S_{0}+S_{1}\right) \quad n=1, \cdots, N_{1}
\end{align*}
$$

and

$$
\begin{align*}
& \sum_{n=1}^{N_{2}} S_{0} V_{n}+S_{1 t}+S_{1} V_{0}=\tilde{V}_{0} S_{1}+\sum_{n=1}^{N_{2}} \tilde{V}_{n} S_{0} \\
& S_{0 t}+S_{0} V_{0}=\tilde{V}_{0} S_{0}  \tag{1.11}\\
& \left(\mu_{n} S_{0}+S_{1}\right) V_{n}-b_{n} S_{0}=\tilde{V}_{n}\left(\mu_{n} S_{0}+S_{1}\right) \quad n=1, \cdots, N_{2}
\end{align*}
$$

In deriving equations (1.10) and (1.11) we have used equation (1.3) and the simple relationships

$$
\begin{equation*}
P_{n} A=I+\nu_{n} P_{n} \quad Q_{n} A=I+\mu_{n} Q_{n} . \tag{1.12}
\end{equation*}
$$

The matrices $S_{0}$ and $S_{1}$ are fixed by specific cholces for the linear transformation $\mathcal{C}$. We consider here the case $\tilde{\Phi}=0$, which is sathsfied for specific solutions of the auxiliary hnear problem, $\bar{\omega}=\Phi_{0}, \Lambda=\Lambda_{0}$,

$$
\Phi_{0}=\left(\begin{array}{cc}
\psi_{1}^{0} & \psi_{2}^{0}  \tag{113}\\
\phi_{1}^{0} & \phi_{2}^{0}
\end{array}\right) \quad \Lambda_{0}=\left(\begin{array}{cc}
\lambda\left(x, t, z_{1}^{0}\right) & 0 \\
0 & \lambda\left(x, t, z_{2}^{0}\right)
\end{array}\right)
$$

where linearly independent $\left(\psi_{j}{ }^{0}, \phi_{j}{ }^{0}\right)^{\mathrm{T}}, j=1.2$ are to be chosen. The condition $\tilde{\Phi}=0$ means simply that, through (18),

$$
\begin{equation*}
S_{1}=-S_{0} \bar{\omega}_{0} \Lambda_{0} \bar{\omega}_{0}^{-1} \tag{1.14}
\end{equation*}
$$

and $S_{0}$ will be defined by nomalization which is to be chosen for each particular system. For a given seed solution $\left\{U_{n}, V_{n}\right\}$ of (1.5), and the corresponding matrix function $\Phi_{0}$, equations (1.10) and (1.11) together with (1.8) and (1.14) will provide a new solution $\left\{\tilde{U}_{n}, \vec{V}_{n}\right\}$ of $(15)$, and, at the same time, a solution of equation (1.9).

## 2. Deformed Maxwell-Bloch equations

The system of Maxwell-Bloch (MB) equations [7],

$$
\begin{align*}
& E_{\eta}=\rho \quad \rho_{\xi}=N E \\
& \mathcal{N}_{\xi}+\frac{1}{2}(E \bar{\rho}+\bar{E} \rho)=0 \tag{2.1}
\end{align*}
$$

has been used to for example describe self-induced transparency (SIT) of a dispersive medium for coherent pulse propagation. Here $\bar{\rho}, \bar{E}$ denote the complex conjugate of $\rho, E$. In SIT $E$ is the complex electric field amplitude, $\rho$ is the induced (complex) polarization and $\mathcal{N}$ is the normalized population inversion of the medium.

In [1] a number of completely integrable deformations of equation (2.1) were reported; we consider here one of physical interest, i.e.

$$
\begin{align*}
& E_{\eta}=\rho \quad \rho_{\xi}=\mathcal{N} E-\tilde{c} \rho \\
& \mathcal{N}_{\xi}+\frac{1}{2}(E \bar{\rho}+\bar{E} \rho)=4 c-\tilde{c} \mathcal{N} . \tag{2.2}
\end{align*}
$$

Here functions $c$ and $\bar{c}$ are determined by equations

$$
\begin{equation*}
\tilde{c}_{\eta}=0 \quad c_{\xi}+2 c \tilde{c}=0 \tag{23}
\end{equation*}
$$

We may think of the deformed equation (2.2) as describing a coherent pulse propagation in a dispersive medium consisting of two-level atoms, includung relaxation rate $\tilde{c}$ and source of inversion sc The relasation rate and source of inversion are connected in a non-trivial way through equation (2.3).

The linear system associated with equation (2.2) is [1]

$$
\begin{equation*}
\Phi_{\xi}=U_{1} \Phi \Lambda+U_{0} \bar{\Phi} \quad \Phi_{\eta}=\hat{\rho} \Phi P \tag{2.4}
\end{equation*}
$$

where

$$
\begin{align*}
& U_{1}=\sigma_{3} \quad U_{0}=\frac{1}{2}\left(\begin{array}{cc}
0 & E \\
-\bar{E} & 0
\end{array}\right) \\
& \Lambda=\left(\begin{array}{cc}
\lambda\left(\xi, \eta, z_{1}\right) & 0 \\
0 & \lambda\left(\xi, \eta, z_{2}\right)
\end{array}\right) \\
& \mathcal{N}=\frac{1}{2}(n+\tilde{n}) \tag{2.5}
\end{align*}
$$

( $\sigma_{3}$ is Pauli's matrix).
The matrices $A$ and $P$ satisfy [1] equations

$$
\begin{equation*}
A_{\eta}=4 c P \quad P_{\xi}=\tilde{c} P \tag{2.6}
\end{equation*}
$$

We can now apply to system (2 2)-(2 6) the DT method described in section 1. In this way we find

$$
\begin{align*}
& S_{0} U_{1}=U_{1} S_{0} \\
& \tilde{U}_{0}=S_{1} U_{0} S_{1}^{-1}+S_{1 \xi} S_{1}^{-1}  \tag{2.7a}\\
& \tilde{U}_{0}=S_{0} U_{0} S_{0}^{-1}+\left[S_{1}, U_{1}\right] S_{0}^{-1}+S_{0 \xi} S_{0}^{-1}-\tilde{c} I  \tag{2.7b}\\
& S_{0 \eta}=0 \\
& \tilde{\hat{\rho}}=S_{1} \hat{\rho} S_{1}^{-1}+4 c S_{0} S_{1}^{-1}  \tag{28a}\\
& \tilde{\hat{\rho}}=S_{0} \hat{\rho} S_{0}^{-1}+4 S_{1 刃} S_{0}^{-1} \tag{2.8b}
\end{align*}
$$

This set of equations corresponds to equations (1.10)-(1.11) If we consider the case when $\bar{c}$ is constant, then equation (27b) gives ummedately the solution for $S_{0}$,

$$
\begin{equation*}
S_{0}=S_{0}(\xi)=\mathrm{I} \mathrm{e}^{\bar{\sigma} \xi} \tag{2.9}
\end{equation*}
$$

In order to proceed further we use the fact that, acting on the set of solutions of the linear system (2.4), there is a group of automorphisms which has the generating transformation

$$
\begin{equation*}
\binom{\psi(\lambda)}{\phi(\lambda)} \rightarrow\binom{-\bar{\phi}(-\bar{\lambda})}{\bar{\psi}(-\bar{\lambda})} . \tag{2.10}
\end{equation*}
$$

Transformation (2.10) allows us to take the seed solution (1.13) in the form

$$
\bar{\Phi}_{0}=\left(\begin{array}{cc}
\psi_{1}\left(\lambda_{1}\right) & -\bar{\phi}_{1}\left(-\bar{\lambda}_{1}\right)  \tag{2.11}\\
\phi_{1}\left(\lambda_{1}\right) & \bar{\psi}_{1}\left(-\bar{\lambda}_{1}\right)
\end{array}\right)
$$

It can be easily seen by differentiating $\operatorname{det} \bar{\omega}_{0}$ with respect to $\xi$ and $\eta$ that this determinant can be chosen to be non-zero and, therefore, the matrix $\Phi_{0}$ is invertible for $\lambda=\lambda_{1}$. The corresponding matrix spectral parameter $\Lambda$ is now

$$
\Lambda_{0}=\left(\begin{array}{cc}
\lambda_{1} & 0  \tag{212}\\
0 & -\bar{\lambda}_{1}
\end{array}\right)
$$

where $\lambda_{1} \equiv \lambda\left(\xi, \eta, z_{1}\right)$, and $z_{1}$ is the value of the hidden spectral parameter which specifies the seed solution.

Equations (2.3) and (2.6) are easily solved for constant $\tilde{c}$ to give

$$
\begin{align*}
& c=c_{0} \mathrm{e}^{-2 \varepsilon \xi} \\
& \lambda_{1} \equiv \lambda\left(\xi, n, z_{1}\right)=\left(2 c_{0} \eta+z_{1}\right)^{1 / 2} \mathrm{e}^{-\tau \xi} \tag{2.13}
\end{align*}
$$

Equations (2.7b) and (2.8b) can now be solved in terms of (2.11) and (212) In this way we find the Darboux dressing expressions for the electric field amplitude $E$, polarization $\rho$ and inversion $\mathcal{N}$,

$$
\begin{align*}
& \tilde{E}=E+4\left(\lambda_{1}+\bar{\lambda}_{1}\right) \frac{\psi_{1} \bar{\phi}_{1}}{\left|\psi_{1}\right|^{2}+\left|\phi_{1}\right|^{2}}  \tag{2.14a}\\
& \tilde{\mathcal{N}}=\mathcal{N}-2 \frac{\partial}{\partial \eta}\left[\left(\lambda_{1}+\bar{\lambda}_{1}\right) \frac{\left|\psi_{1}\right|^{2}-\left|\phi_{1}\right|^{2}}{\left|\psi_{1}\right|^{2}+\left|\phi_{1}\right|^{2}}\right]  \tag{2.14b}\\
& \tilde{\rho}=\rho+4 \frac{\partial}{\partial \eta} \frac{\left(\lambda_{1}+\bar{\lambda}_{1}\right) \psi_{1} \bar{\phi}_{1}}{|\psi|^{2}+\left|\phi_{1}\right|^{2}} \tag{214c}
\end{align*}
$$

where $E, \mathcal{N}, \rho$ comprise the seed solution to be found from equation (2.2). The solution for the electric field amplitude (2.14a) gives immediately the energy density of the electric field,
$|\tilde{E}|^{2}=|E|^{2}+4 \tilde{\tilde{C}}\left(\lambda_{1}+\bar{\lambda}_{1}\right) \frac{\partial}{\partial \xi} \ln \left(\left|\psi_{1}\right|^{2}+\left|\phi_{1}\right|^{2}\right)+4 \frac{\partial^{2}}{\partial \tilde{\xi}^{2}} \ln \left(\left|\psi_{1}\right|^{2}+\left|\phi_{1}\right|^{2}\right)$.

It is straightforward to work out the Darboux dressing equations for an $N$-step iteration. We give these equations in the appendix.

Thanks to its algebraic nature the DT method allows wide classes of solutions to be built up, generally in an arbitrary background, which is advantageous in the explicit computations. As far as the MB system is concerned, some examples of successful application of the DT formalism can be found in [4, 5, 8] where soliton solutions in pentedic and constant backgrounds are considered. In a forthcoming publication we will consider the question of an optical pulse with a background of a finitely-many gaps poiential by Darboux dressing of the corresponding $\Psi$-function [6]. The question of the interaction of an optical pulse with a spontaneous radiation field is investigated by means of the DT method in [8].

The simplest seed solution of equation (2.2) is that for zero electric field,

$$
\begin{align*}
& E=0 \quad \rho=0 \\
& \mathcal{N}(\xi, \eta)=\mathcal{N}_{0} \mathrm{e}^{-\bar{c} \xi}+\frac{4 c_{0}}{\tilde{c}}\left(\mathrm{e}^{-\bar{c} \xi}-\mathrm{e}^{-2 \tilde{z} \xi}\right) \tag{2.16}
\end{align*}
$$

When this seed solution is used in the auxiliary linear problem (2.4), we find that

$$
\begin{equation*}
\psi_{1}=c_{1} \mathrm{e}^{\theta(\xi, \eta)} \quad \phi_{1}=c_{2} \mathrm{e}^{-\theta(\xi, \eta)} \tag{217}
\end{equation*}
$$

where $c_{1,2}$ are constants and

$$
\begin{equation*}
\theta(\xi, \eta)=\frac{\mathcal{N}_{0}}{4 c_{0}}\left[\left(2 c_{0} \eta+z_{1}\right)^{1 / 2}-z_{1}^{1 / 2}\right]+\tilde{c}^{-1}\left(2 c_{0} \eta+z_{1}\right)^{1 / 2}\left(1-\mathrm{e}^{-\tilde{c} \xi}\right) \tag{2.18}
\end{equation*}
$$

The result (2.14a) now means that a non-zero electric field amplitude satisfying equations (2.2) is

$$
\begin{equation*}
\tilde{E}(\xi, \eta)=2\left(\lambda_{1}+\bar{\lambda}_{1}\right) \frac{\exp \left(\theta(\xi, \eta)-\bar{\theta}(\xi, \eta)+\mathrm{i} \phi_{0}\right)}{\cosh \left(\theta(\xi, \eta)+\bar{\theta}(\xi, \eta)+\psi_{0}\right)} \tag{2.19}
\end{equation*}
$$

Corresponding expressions for the polarization and unversion are easily found from (2.14b) and (2.14c)

$$
\begin{align*}
& \tilde{\rho}=\frac{\lambda_{1}+\bar{\lambda}_{1}}{\left|\lambda_{1}\right|^{2}} \frac{\exp \left(\theta-\bar{\theta}+1 \phi_{0}\right)}{\cosh \left(\theta+\bar{\theta}+\psi_{0}\right)}\left\{4 c-\left(\lambda_{1}-\bar{\lambda}_{1}\right) \mathcal{N}-\left(\lambda_{1}+\bar{\lambda}_{1}\right) \frac{\mathcal{N}}{\cosh \left(\theta+\bar{\theta}+\psi_{0}\right)}\right\}  \tag{2.19a}\\
& \tilde{\mathcal{N}}=\mathcal{N}-2 \frac{\lambda_{1}+\bar{\lambda}_{1}}{\left|\lambda_{1}\right|^{2}}\left\{c \tanh \left(\theta+\bar{\theta}+\psi_{0}\right)+\frac{\mathcal{N}\left(\lambda_{1}+\bar{\lambda}_{1}\right)}{4 \cosh ^{2}\left(\theta+\bar{\theta}+\psi_{0}\right)}\right\} \tag{2.19b}
\end{align*}
$$

Here $\mathcal{N}$ has the value (2.16).
It is easy to show that for $c_{0}=\tilde{c}=0$ solution (2.19) is the usual one-soliton solution of SIT in the case of an nfinitely narrow line width [7]. Indeed, taking $z_{1}{ }^{1 / 2}=\alpha+i \beta$, we get
$\tilde{E}(\xi, \eta)=4 \alpha \exp -\mathrm{i}\left(\frac{\beta \mathcal{N}_{0}}{2\left(\alpha^{2}+\beta^{2}\right)} \eta-\phi_{0}\right)\left[\cosh 2 \alpha\left(\xi-\xi_{0}+\frac{\mathcal{N}_{0}}{4\left(\alpha^{2}+\beta^{2}\right)} \eta\right)\right]^{-1}$.

In order to analyse the solution (219) in more detail we shall make the transformation to the real space variables $x, t$ :

$$
\begin{equation*}
\xi=\Omega\left(t-\frac{x}{c}\right) \quad \eta=\frac{\Omega}{c} x \tag{220}
\end{equation*}
$$

where $\Omega^{2}=2 \pi \hbar^{-1} n_{0} \omega_{0} \mu^{2}$, and $c$ is the velocity of light. Here $n_{0}$ is the density of the two-level atoms, $\hbar \omega_{0}$ is the energy difference between the levels and $\mu$ is the dipole momentum. We then find from (2.19) that, asymptotically when $t \rightarrow \infty$, the group velocity $v$ of the electromagnetic pulse be baves as

$$
\begin{equation*}
v=c+\frac{c A \Omega^{-3 / 2}}{4 \sqrt{2 c_{0}}\left(1+\left(\tilde{c} N_{0} / 4 c_{0}\right)\right)} t^{-3 / 2}+\mathrm{o}\left(t^{-3 / 2}\right) \tag{2.21}
\end{equation*}
$$

The group velocity of the pulse approaches thus the velocity of light In the same limit $t \rightarrow \infty$, the width in space $\Delta x$ of the pulse is asymptotically

$$
\begin{equation*}
\Delta x=\frac{c \chi^{(1)} \Omega^{-3 / 2}}{\sqrt{2 c_{0}}\left(1+\left(\tilde{c} N_{0} / 4 c_{0}\right)\right)} t^{-1 / 2}+o\left(t^{-1 / 2}\right) \tag{2.22}
\end{equation*}
$$

where $\chi^{(1)}$ is the fixed phase $\chi^{(1)}=\theta+\bar{\theta}+\psi_{0}$. Similarly the point of maximum amplitude of the pulse ( 219 ) moves in the $2 t$ plane along the asymptotic trajectory
$x-c t=\frac{c}{\tilde{c} \Omega}\left(1+\frac{\tilde{c} N_{0}}{4 c_{0}}\right)-\frac{c A \Omega^{-3 / 2}}{2 \sqrt{2 \varepsilon_{0}}\left(1+\left(\tilde{c} N_{0} / 4 c_{0}\right)\right)} t^{-1 / 2}+o\left(t^{-1 / 2}\right)$.
Finally it should be noted that, in the limit $t \rightarrow \infty$, the exponential factor $\exp (-\tilde{c} \xi)=\exp (-\tilde{c} \Omega(t-(x / c)))$ tends to one and, therefore, the amplitude of the pulse (2.19) diverges as $t^{1 / 2}$

## Acknowledgments

The author would like to express his gratitude to VB Matveev for continuous encouragement and to J Tmonen for useful discussions and help in completing this work.

## Appendiz

In this appendix we derive the Darboux dressing equations for an $N$-step iteration.
To this end we have to generalize ansatz (1.8) such that we assume
$\Phi[N]=\sum_{j=0}^{N} S_{j} \Phi \mathrm{~A}^{N-3} S_{3}=\left(\begin{array}{ll}a_{j} & b_{j} \\ c_{j} & d_{j}\end{array}\right)$.
One can now proceed in the same way as in section 2 for the one-step iteration. Inserting ansatz (A1) into equation (2.4) we find that the followng relationships must be satisfied:

$$
\begin{align*}
& S_{0} U_{1}=U_{1} S_{0} \\
& U_{0}[N]=S_{N} U_{0} S_{N}^{-1}+S_{N \xi} S_{N}^{-1} \quad 0 \leqslant j \leqslant N-1  \tag{A2}\\
& U_{0}[N]=S_{j} U_{0} S_{j}^{-1}+S_{j \xi} S_{j}^{-1}-\tilde{c}(N-j) I+\left[S_{j+1}, U_{1} 1 S_{j}^{-1}\right.
\end{align*}
$$

and
$S_{0 \eta}=0$
$\hat{\rho}[N]=S_{0} \hat{\rho} S_{0}^{-1}+4 S_{1 \eta} S_{0}^{-1}$
$\hat{\rho}[N]=S_{N} \hat{\rho} S_{N}^{-1}+4 c S_{N-1} S_{N}^{-1}$
$\hat{\rho}[N]=S_{j-1} \hat{\rho} S_{j-1}^{-1}+4 S_{j \eta} S_{j-1}^{-1}+4 c(N-\jmath+2) S_{j-2} S_{j-1}^{-1} \quad 2 \leqslant j \leqslant N$.
These equations will determine the dressed variables in terms of matrices $S$, which we still have to define.

First we choose a set $\left\{\bar{\Phi}_{m}, \Lambda_{m}\right\}_{m=1}^{N}$,

$$
\bar{\Phi}_{m}=\left(\begin{array}{ll}
\psi_{2 m-1} & \psi_{2 m}  \tag{A4}\\
\phi_{2 m-1} & \phi_{2 m}
\end{array}\right) \quad \Lambda_{m}=\left(\begin{array}{cc}
\lambda_{2 m-1} & 0 \\
0 & \lambda_{2 m}
\end{array}\right)
$$

of solutions of the auxiliary linear problem. As before $A_{m}$ is the matrix spectral parameter related to $\omega_{m}$. We then set

$$
\begin{equation*}
\left.\tilde{\Phi}[N]\right|_{A=A_{m}, \bar{W}=\bar{\Phi}_{m}}=0 \quad m=1, \ldots, N . \tag{A5}
\end{equation*}
$$

In terms of the solutions (A4) equation (A5) can be expressed in the form

$$
\begin{align*}
& \sum_{j=1}^{N} \lambda_{n}^{N-3}\left(a_{j} \psi_{n}+b_{j} \phi_{n}\right)=-\lambda_{n}^{N} \psi_{n} \\
& \sum_{j=1}^{N} \lambda_{n}^{N-j}\left(c_{j} \psi_{n}+d_{j} \phi_{n}\right)=-\lambda_{n}^{N} \phi_{n} . \tag{A6}
\end{align*} \quad \text { for } \quad n=1,2, \ldots, 2 N
$$

In deriving equation (A6) we have used the result

$$
\begin{equation*}
S_{0}=\operatorname{diag}\left(a_{0}, d_{0}\right) e^{\varepsilon N \xi} \tag{A7}
\end{equation*}
$$

which can be easily found from equation (A2). Because $S_{0}$ is only a normalizing factor we can choose it to be

$$
\begin{equation*}
S_{0}=\mathrm{I} \cdot \mathrm{e}^{\varepsilon N \xi} . \tag{A8}
\end{equation*}
$$

Equation (A6) can be solved by, for example, the Cramer's rule and we find

$$
\begin{array}{ll}
c_{j}=-\frac{\operatorname{det}\left(v_{n l}^{(2 j-1)}\right)}{\operatorname{det}\left(U_{n l}\right)} & b_{j}=-\frac{\operatorname{det}\left(v_{n j}^{(2 j)}\right)}{\operatorname{det}\left(U_{n l}\right)} \\
c_{3}=-\frac{\operatorname{det}\left(w_{n l}^{(2,-1)}\right)}{\operatorname{det}\left(U_{n l}\right)} & d_{1}=-\frac{\operatorname{det}\left(w_{n l}^{(2 j)}\right)}{\operatorname{det}\left(U_{n l}\right)} \tag{A9}
\end{array} \quad j=1,2, \ldots, N,
$$

where

$$
\left.\begin{array}{rl}
U_{n l} & =\left\{\begin{array}{ll}
\lambda_{n}^{N-k} \psi_{n} & l=2 k-1 \\
\lambda_{n}^{N-k} \phi_{n} & l=2 k
\end{array} \quad \text { for } \quad k=1,2, \ldots, N\right. \\
v_{n l}^{(i)} & =\left\{\begin{array}{ll}
U_{n l} & i \neq i \\
\lambda_{n}^{N} \psi_{n} & l=i
\end{array} \quad \text { for } \quad i=1,2, \ldots, 2 N\right.
\end{array}\right\} \begin{array}{ll}
(\mathrm{e}) & = \begin{cases}U_{n l} & l \neq i \\
\lambda_{n}^{N} \phi_{n} & l=i .\end{cases} \tag{A10}
\end{array}
$$

In order to save the involution property of the system (2.4) together with (25), we should choose, as in expression (2.19), the matrix $\omega_{m}$ of the fixed solutions of the auxiliary problem in the form

$$
\bar{\Phi}_{m}=\left(\begin{array}{cc}
\psi_{2 m-1} & -\bar{\phi}_{2 m-1}  \tag{A11}\\
\phi_{2 m-1} & \bar{\psi}_{2 m-1}
\end{array}\right)
$$

and the associated matrix spectral parameter in the form

$$
\begin{equation*}
\Lambda_{m}=\operatorname{diag}\left(\lambda_{2 m-1} \quad-\bar{\lambda}_{2 m-1}\right) \tag{A12}
\end{equation*}
$$

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    $\ddagger$ It is worth mentioning that the generalization of the traditional inverse problem method for the non-isospectral scattering problems have been studied in [11] and [12]. In this article we consider the case of both space and time dependence of the spectral parameter.

