

# Approximate monochromatic reconstruction algorithms for non-linearized two-dimensional multi-channel inverse problems

Roman Novikov, Matteo Santacesaria

#### ▶ To cite this version:

Roman Novikov, Matteo Santacesaria. Approximate monochromatic reconstruction algorithms for non-linearized two-dimensional multi-channel inverse problems. 2011. hal-00594674v1

### HAL Id: hal-00594674 https://hal.science/hal-00594674v1

Preprint submitted on 20 May 2011 (v1), last revised 8 Nov 2011 (v3)

**HAL** is a multi-disciplinary open access archive for the deposit and dissemination of scientific research documents, whether they are published or not. The documents may come from teaching and research institutions in France or abroad, or from public or private research centers.

L'archive ouverte pluridisciplinaire **HAL**, est destinée au dépôt et à la diffusion de documents scientifiques de niveau recherche, publiés ou non, émanant des établissements d'enseignement et de recherche français ou étrangers, des laboratoires publics ou privés.

#### APPROXIMATE MONOCHROMATIC RECONSTRUCTION ALGORITHMS FOR NON-LINEARIZED TWO-DIMENSIONAL MULTI-CHANNEL INVERSE **PROBLEMS**

#### ROMAN G. NOVIKOV AND MATTEO SANTACESARIA

ABSTRACT. We consider two inverse problems for the multi-channel two-dimensional Schrödinger equation at fixed positive energy, i.e. the equation  $-\Delta \psi + V(x)\psi = E\psi$  at fixed positive E, where V is a matrixvalued potential. The first is the Gel'fand-Calderón inverse problem on a bounded domain D and the second is the inverse fixed-energy scattering problem on the whole plane  $\mathbb{R}^2$ . We present in this paper two algorithms which give efficient approximate solutions to these problems: in particular, in both cases we show that the potential V is reconstructed with Lipschitz stability by these algorithms up to  $O(E^{-(m-2)/2})$  in the uniform norm as  $E \to +\infty$ , under the assumptions that V is m-times differentiable in  $L^1$ , for  $m \geq 3$ , and has sufficient boundary decay.

#### 1. Introduction

We consider the equation

$$(1.1) -\Delta \psi + V(x)\psi = E\psi, x \in \mathbb{R}^2, E > 0,$$

where

V is a sufficiently regular  $M_n(\mathbb{C})$ -valued function on  $\mathbb{R}^2$ (1.2)with sufficient decay at infinity,

 $M_n(\mathbb{C})$  is the set of the  $n \times n$  complex matrices. This equation will also be considered on a domain D, where

D is an open bounded domain in  $\mathbb{R}^2$  with a  $C^2$  boundary. (1.3)

Equation (1.1) at fixed E can be considered as rather general multichannel Schrödinger (resp. acoustic) equation on D at a fixed energy (resp. frequency) related to E. It arises, in particular, as a 2D approximation to the following 3D equation

(1.4) 
$$-\Delta_{x,z}\psi + v(x,z)\psi = E\psi, \quad (x,z) \in \Omega = D \times L,$$

where L = [a, b],  $a, b \in \mathbb{R}$ , v is a sufficiently regular complex-valued function on  $\Omega$  and  $\psi|_{D \times \partial L} = 0$  (for example): see [17, Sec. 2]. In this framework, the approximate 2D matrix-valued potential V is given by

$$(1.5) V_{ij}(x) = \lambda_i \delta_{ij} + \int_L \bar{\phi}_i(z) v(x, z) \phi_j(z) dz, x \in D,$$

for  $1 \leq i, j \leq n$ , where  $n \in \mathbb{N}$ ,  $\{\phi_j\}_{j \in \mathbb{N}}$  is the orthonormal basis of  $L^2(L)$  given by the eigenfunctions of  $-\frac{d^2}{dz^2}$  such that  $\phi_j|_{\partial L} = 0$ ,  $-\frac{d^2\phi_j}{dz^2} = \lambda_j\phi_j$ , for  $j \in N$ , and  $\delta_{ij} = 1$  if i = j and 0 otherwise.

In addition, equation (1.1) can be seen as a particular case of the 2D Schrödinger equation in an external Yang-Mills field.

For equation (1.1) on D we consider the Dirichlet-to-Neumann map  $\Phi(E)$  such that

(1.6) 
$$\Phi(E)(\psi|_{\partial D}) = \frac{\partial \psi}{\partial \nu}\Big|_{\partial D}$$

for all sufficiently regular solution  $\psi$  of (1.1) on  $\bar{D} = D \cup \partial D$ , where  $\nu$  is the outer normal of  $\partial D$ . Here we assume also that

(1.7) E is not a Dirichlet eigenvalue for the operator  $-\Delta + V$  in D.

This construction gives rise to the following inverse boundary value problem on D:

#### **Problem 1.** Given $\Phi(E)$ , find V on D.

On the other hand, for equation (1.1) on  $\mathbb{R}^2$ , under assumptions (1.2), we consider the scattering amplitude f defined as follows: we consider the continuous solutions  $\psi^+(x,k)$  of (1.1), where k is a parameter,  $k \in \mathbb{R}^2$ ,  $k^2 = E$ , such that

(1.8) 
$$\psi^{+}(x,k) = e^{ikx}I - i\pi\sqrt{2\pi}e^{-i\frac{\pi}{4}}f\left(k,|k|\frac{x}{|x|}\right)\frac{e^{i|k||x|}}{\sqrt{|k||x|}} + o\left(\frac{1}{\sqrt{|x|}}\right), \quad \text{as } |x| \to \infty,$$

for some a priori unknown  $M_n(\mathbb{C})$ -valued function f, where I is the identity matrix. The function f on  $\mathscr{M}_E = \{(k, l) \in \mathbb{R}^2 \times \mathbb{R}^2 : k^2 = l^2 = E\}$  arising in (1.8) is the scattering amplitude for the potential V in the framework of equation (1.1).

This construction gives rise to the following inverse scattering problem on  $\mathbb{R}^2$ :

**Problem 2.** Given f on  $\mathcal{M}_E$ , find V on  $\mathbb{R}^2$ .

Problems 1 and 2 can be considered as multi-channel fixed-energy analogues in dimension d = 2 of inverse problems formulated in [8] in dimension  $d \geq 2$ . Note that Problems 1 and 2 are not overdetermined, in the sense that we consider the reconstruction of a  $M_n(\mathbb{C})$ -valued function V of two variables from  $M_n(\mathbb{C})$ -valued inverse problem data dependent on two variables. In addition, the history of inverse problems for the two-dimensional Schrödinger equation at fixed energy goes back to [5] (see also [13, 9] and reference therein). Note also that Problem 1 can be considered as a model problem for the monochromatic ocean tomography (e.g. see [1] for similar problems arising in this tomography).

As regards efficient algorithms for solving Problems 1 and 2 for the scalar case, i.e. for n=1, see [12, 13, 14, 15]. In addition, as concerns numerical implementations of these algorithms for Problem 2 for n=1, see [2, 4], and references therein.

Nevertheless, the fixed-energy global uniqueness for Problem 1 (and for Problem 2 with compactly supported V) for n=1 was completely proved only recently in [3]. The reconstruction scheme of [3] is not optimal with respect to its stability properties, and, therefore, is not efficient numerically in comparison with the aforementioned 2D reconstructions of [12, 13, 14, 15], but it is very efficient for proving some global mathematical results. In particular: a related global logarithmic stability estimate for Problem 1 for n=1 was proved in [16]; global uniqueness and reconstruction results for Problem 1 for  $n \geq 2$  were obtained in [17]; a global logarithmic stability estimate for Problem 1 for  $n \geq 2$  was proved in [18]. In addition, Problem 2 with compactly supported V can be reduced, for  $n \geq 2$ , to Problem 1, as in [12] for n=1. This implies, at least, global uniqueness for Problem 2 (in the compactly supported case). On the other hand, the uniqueness for Problem 2 fails already for scalar (n=1) real-valued spherically-symmetric potentials V of the Schwartz class on  $\mathbb{R}^2$  (see [10]).

The main purpose of the present work consists in generalizing the aforementioned reconstruction approach of [14, 15] to the case of Problems 1 and 2 for  $n \geq 2$ . As well as for n = 1 this functional analytic approach gives an efficient non-linear approximation  $V_{appr}(x, E)$  to the unknown V(x) of Problems 1 and 2. The reconstruction of  $V_{appr}(x, E)$  from  $\Phi(E)$  for Problem 1 and from f on  $\mathcal{M}_E$  for Problem 2 is realized with some Lipschitz stability and is based on solving linear integral equations; see Algorithms 1 and 2 of Section 3, Theorems 6.1, 6.2 and Remarks 6.4, 6.5 of Section 6. Among these linear integral equations, the most important ones arise from a non-local Riemann-Hilbert problem. Another important part of these equations is used for transforming  $\Phi(E)$  for Problem 1 and f on  $\mathcal{M}_E$  for Problem 2

into  $M_n(\mathbb{C})$ -valued Faddeev function analogues  $h_{\pm}$  on  $\mathscr{M}_E$ , involved in the formulation of the above-mentioned Riemann-Hilbert problem. In addition,

$$||V_{appr}(\cdot, E) - V|| = \varepsilon(E)$$

rapidly decays as  $E \to +\infty$ , where  $\|\cdot\|$  denotes an appropriate norm. In particular,  $\varepsilon(E) = O(E^{-\infty})$  as  $E \to +\infty$  if  $\|\cdot\|$  is specified as  $\|\cdot\|_{L^{\infty}(D)}$  and  $V \in C^{\infty}(\mathbb{R}^2, M_n(\mathbb{C}))$ , supp  $V \subset D$ , for Problem 1 and if  $\|\cdot\|$  is specified as  $\|\cdot\|_{L^{\infty}(\mathbb{R}^2)}$  and  $V \in \mathcal{S}(\mathbb{R}^2, M_n(\mathbb{C}))$ , for Problem 2, where  $\mathcal{S}$  denotes the Schwartz class. In addition, no reconstruction algorithms for Problems 1 and 2 — comparable, with respect to their stability, with Algorithms 1 and 2 and with an approximation error decaying more rapidly than  $O(E^{-\frac{1}{2}})$  as  $E \to \infty$  — are available in the preceding literature, even for  $V \in C^{\infty}(\mathbb{R}^2, M_n(\mathbb{C}))$ , supp  $V \subset D \subset \mathbb{R}^2$ , when  $n \geq 2$  (in general).

Finally, in spite of the fact that some excellent properties of Algorithms 1 and 2 are proved assuming that V is sufficiently smooth and that E is sufficiently great in comparison with (some norm of) V, we expect that these algorithms will work rather well even for V with discontinuities and for the case when E is not very big in comparison with V. This expectation is based on numerical results for Algorithm 2 for the case n=1; see [4] and references therein. Numerical implementations of Algorithm 1 for  $n \geq 1$  and Algorithm 2 for  $n \geq 2$  are in preparation.

**Acknowledgements.** We thank V. A. Burov, O. D. Rumyantseva, S. N. Sergeev for very useful discussions.

#### 2. Faddeev functions

In this section we recall some preliminary definitions.

Under assumptions (1.2), we consider the Faddeev functions  $G(x,k) = e^{ikx}g(x,k)$ ,  $\psi(x,k)$ , h(k,l) (see [6, 7, 12] for n=1):

(2.1) 
$$g(x,k) = -\left(\frac{1}{2\pi}\right)^2 \int_{\mathbb{R}^2} \frac{e^{i\xi x}}{\xi^2 + 2k\xi} d\xi,$$

(2.2) 
$$\psi(x,k) = e^{ikx}I + \int_{\mathbb{R}^2} G(x-y,k)V(y)\psi(y,k)dy,$$

(2.3) 
$$h(k,l) = \left(\frac{1}{2\pi}\right)^2 \int_{\mathbb{D}^2} e^{-ilx} V(x) \psi(x,k) dx,$$

where  $x = (x_1, x_2) \in \mathbb{R}^2$ ,  $k = (k_1, k_2) \in \mathbb{C}^2 \setminus \mathbb{R}^2$ ,  $l = (l_1, l_2) \in \mathbb{C}^2$ ,  $\text{Im} k = \text{Im} l \neq 0$  and l is the identity matrix. We recall that

(2.4) 
$$(\Delta + k^2)G(x,k) = \delta(x), \quad x \in \mathbb{R}^2, k \in \mathbb{C}^2 \setminus \mathbb{R}^2.$$

In addition, formula (2.2) at fixed k is considered as an equation for

(2.5) 
$$\psi(x,k) = e^{ikx}\mu(x,k),$$

where  $\mu$  is sought in  $L^{\infty}(\mathbb{R}^2, M_n(\mathbb{C}))$ . As a corollary of (2.1), (2.2) and (2.4),  $\psi$  satisfies (1.1) for  $E = k^2 = k_1^2 + k_2^2$ . In addition, h in (2.3) is a generalised scattering amplitude in the complex domain for the potential V.

For  $\gamma \in S^1 = \{ \gamma \in \mathbb{R}^2 : |\gamma| = 1 \}$ , we consider

(2.6) 
$$G_{\gamma}(x,k) = G(x,k+i0\gamma),$$

(2.7) 
$$\psi_{\gamma}(x,k) = e^{ikx}\mu_{\gamma}(x,k), \qquad \mu_{\gamma}(x,k) = \mu(x,k+i0\gamma),$$

$$(2.8) h_{\gamma}(k,l) = h(k+i0\gamma, l+i0\gamma),$$

where  $x \in \mathbb{R}^2$ ,  $k \in \mathbb{R}^2$ ,  $l \in \mathbb{R}^2$ .

In addition, the functions

(2.9) 
$$G^{+}(x,k) = G_{k/|k|}(x,k) = -\frac{i}{4}H_{0}^{1}(|x||k|),$$

(2.10) 
$$\psi^+(x,k) = e^{ikx}\mu^+(x,k), \qquad \mu^+(x,k) = \mu_{k/|k|}(x,k),$$

$$(2.11) f(k,l) = h_{k/|k|}(k,l),$$

for  $k, l \in \mathbb{R}^2$ , |k| = |l|, are functions from the classical scattering theory; in particular, f is the scattering amplitude of (1.8) and  $H_0^1$  is the Hankel function of the first type. We also define

(2.12) 
$$h_{\pm}(k,l) = h_{\pm \hat{k}_{\perp}}(k,l),$$

(2.13) 
$$\mu_{\pm}(x,k) = \mu_{\pm \hat{k}_{\perp}}(x,k), \qquad \psi_{\pm}(x,k) = \psi_{\pm \hat{k}_{\perp}}(x,k),$$

where  $k, l, x \in \mathbb{R}^2$ , |k| = |l|,  $\hat{k}_{\perp} = |k|^{-1}(-k_2, k_1)$  for  $k = (k_1, k_2)$ . Note that  $\mu_+ \neq \mu^+$  and  $\psi_+ \neq \psi^+$ , in general. We shall consider, in particular, the following restriction of the function h:

(2.14) 
$$b(k) = h(k, -\bar{k}), \quad \text{for } k \in \mathbb{C}^2, \ k^2 = E > 0.$$

We now introduce the notations

(2.15) 
$$z = x_1 + ix_2, \qquad \bar{z} = x_1 - ix_2,$$

$$\frac{\partial}{\partial z} = \frac{1}{2} \left( \frac{\partial}{\partial x_1} - i \frac{\partial}{\partial x_2} \right), \qquad \frac{\partial}{\partial \bar{z}} = \frac{1}{2} \left( \frac{\partial}{\partial x_1} + i \frac{\partial}{\partial x_2} \right),$$

$$\lambda = E^{-1/2} (k_1 + ik_2), \qquad \lambda' = E^{-1/2} (l_1 + il_2),$$

where  $x = (x_1, x_2) \in \mathbb{R}^2$ ,  $k = (k_1, k_2)$ ,  $l = (l_1, l_2) \in \mathbb{C}^2$ ,  $k^2 = l^2 = E \in \mathbb{R}_+$ . In the new notations

(2.16a) 
$$k_1 = \frac{1}{2}E^{1/2}(\lambda + \lambda^{-1}), \qquad k_2 = \frac{i}{2}E^{1/2}(\lambda^{-1} - \lambda),$$

(2.16b) 
$$l_1 = \frac{1}{2}E^{1/2}(\lambda' + {\lambda'}^{-1}), \qquad l_2 = \frac{i}{2}E^{1/2}(\lambda'^{-1} - {\lambda'}),$$

(2.16c) 
$$\exp(ikx) = \exp[\frac{i}{2}E^{1/2}(\lambda \bar{z} + \lambda^{-1}z)],$$

where  $\lambda, \lambda' \in \mathbb{C} \setminus \{0\}$ ,  $z \in \mathbb{C}$  and the Schrödinger equation (1.1) takes the form

(2.17) 
$$-4\frac{\partial^2}{\partial z\partial\bar{z}}\psi + V(z)\psi = E\psi, \qquad z \in \mathbb{C}.$$

In addition, the functions f from (1.8) and (2.11),  $h_{\pm}$  from (2.12),  $\mu^{+}, \psi^{+}$  from (2.10),  $\mu_{\pm}, \psi_{\pm}$  from (2.13),  $\psi$  from (2.2),  $\mu$  from (2.5) and b from (2.14) take the form

(2.18) 
$$f = f(\lambda, \lambda', E), \qquad h_{\pm} = h_{\pm}(\lambda, \lambda', E), \mu^{+} = \mu^{+}(z, \lambda, E), \qquad \psi^{+} = \psi^{+}(z, \lambda, E), \mu_{\pm} = \mu_{\pm}(z, \lambda, E), \qquad \psi_{\pm} = \psi_{\pm}(z, \lambda, E),$$

where  $\lambda, \lambda' \in T$ ,  $z \in \mathbb{C}$ ,  $E \in \mathbb{R}_+$ ,

(2.19) 
$$\mu = \mu(z, \lambda, E), \qquad \psi = \psi(z, \lambda, E), \qquad b = b(\lambda, E),$$

where  $\lambda \in \mathbb{C} \setminus T$ ,  $z \in \mathbb{C}$ ,  $E \in \mathbb{R}_+$ . Here

$$(2.20) T = \{ \zeta : \zeta \in \mathbb{C}, |\zeta| = 1 \}.$$

Under assumption (1.2), for E sufficiently large the function  $\mu(z, \lambda, E)$  has the following properties (see [14, 15] for n = 1 and Section 4 for  $n \ge 2$ ):

- (2.21)  $\mu(z, \lambda, E)$  is continuous in  $\lambda \in \mathbb{C} \setminus T$ ;
- (2.22)  $\mu(z, \lambda(1 \mp 0), E) = \mu_{+}(z, \lambda, E)$  for  $\lambda \in T$ ;
- (2.23)  $\mu_{+}(z,\lambda,E) = \mu^{+}(z,\lambda,E)$

$$+ \pi i \int_{T} \mu^{+}(z, \lambda'', E) \chi_{+} \left( \pm i \left( \frac{\lambda}{\lambda''} - \frac{\lambda''}{\lambda} \right) \right) h_{\pm}(\lambda, \lambda'', z, E) |d\lambda''|,$$

for  $\lambda \in T$ , where

$$(2.24) \chi_{+}(s) = 0 \text{ for } s < 0, \quad \chi_{+}(s) = 1 \text{ for } s > 0.$$

$$(2.25) \quad h_{\pm}(\lambda, \lambda', z, E) = \exp\left[-\frac{i}{2}E^{1/2}\left(\lambda \bar{z} + \frac{z}{\lambda} - \lambda' \bar{z} - \frac{z}{\lambda'}\right)\right] h_{\pm}(\lambda, \lambda', E);$$

(2.26) 
$$\frac{\partial}{\partial \bar{\lambda}} \mu(z, \lambda, E) = \mu \left( z, -\frac{1}{\bar{\lambda}}, E \right) r(\lambda, z, E),$$

for  $\lambda \in \mathbb{C} \setminus T$ , where

(2.27)

$$r(\lambda, z, E) = \exp\left[-\frac{i}{2}E^{1/2}\left(\lambda\bar{z} + \frac{z}{\lambda} + \bar{\lambda}z + \frac{\bar{z}}{\bar{\lambda}}\right)\right] \frac{\pi}{\bar{\lambda}}\operatorname{sign}(\lambda\bar{\lambda} - 1)b(\lambda, E),$$

where b is defined by means of (2.14) and (2.16a);

(2.28) 
$$\mu(z,\lambda,E) = I + \frac{\mu_{-1}(z,E)}{\lambda} + o\left(\frac{1}{\lambda}\right), \qquad \lambda \to \infty,$$

(2.29) 
$$V(z) = 2iE^{1/2} \frac{\partial}{\partial z} \mu_{-1}(z, E).$$

The following formula is valid (see [15] for n = 1 and Section 4 for  $n \ge 2$ ):

$$(2.30) V(z) = 2iE^{1/2} \frac{\partial}{\partial z} \left( \frac{1}{\pi} \int_{D_{-}} \mu(z, -\frac{1}{\zeta}, E) r(\zeta, z, E) d\text{Re}\zeta d\text{Im}\zeta \right) + \frac{1}{2\pi i} \int_{T} \mu_{-}(z, \zeta, E) i\zeta |d\zeta| ,$$

for  $z \in \mathbb{C}$ , E sufficiently large and  $D_{-} = \{\zeta : \zeta \in \mathbb{C}, |\zeta| > 1\}.$ 

#### 3. Reconstruction algorithms

We present here Algorithms 1 and 2, which yield approximate but sufficiently stable solutions to Problems 1 and 2, respectively. These algorithms have a final common part: the reconstruction of the approximate potential  $V_{appr}$  starting from  $h_{\pm}$  of (2.12). Thus, for the sake of clarity, we first give the different initial parts of the algorithms—that is, the reconstruction of  $h_{\pm}$  starting from  $\Phi(E)$  for Algorithm 1 and from f for Algorithm 2—and then the final common part.

Note that in both algorithms we consider in particular the functions  $\psi_{\pm}, h_{\pm}, \mu_{-}$  of (2.12), (2.13) and  $\mu^{+}$  of (2.10). In addition, in Algorithm 1, in the definitions of these functions we assume that  $V \equiv 0$  on  $\mathbb{R}^{2} \setminus D$ .

**Algorithm 1**  $(\Phi(E) \longrightarrow h_{\pm})$ . Given  $\Phi(E)$ , for E sufficiently large, we first reconstruct  $\psi_{\pm}(x,k)|_{\partial D}, k \in \mathbb{R}^2, k^2 = E$ , with the help of the following Fredholm linear integral equation (see [12] for n = 1 and Section 4 for  $n \geq 2$ ):

(3.1) 
$$\psi_{\pm}(x,k)|_{\partial D} = e^{ikx}I + \int_{\partial D} A_{\pm}(x,y,k)\psi_{\pm}(y,k)dy, \quad k \in \mathbb{R}^2, \ k^2 = E,$$

where

(3.2) 
$$A_{\pm}(x,y,k) = \int_{\partial D} G_{\pm}(x-\xi,k) (\Phi - \Phi_0) (\xi,y,E) d\xi, \quad x,y \in \partial D,$$

(3.3) 
$$G_{\pm}(x,k) = G^{+}(x,k) - \frac{1}{4\pi i} \int_{S^{1}} e^{i|k|\theta x} \chi_{+}(\pm \theta k_{\perp}) d\theta,$$

I is the identity matrix,  $(\Phi - \Phi_0)(x, y, E)$  is the Schwartz kernel of the operator  $\Phi(E) - \Phi_0(E)$ ,  $\Phi_0(E)$  is the Dirichlet-to-Neumann operator associated to the zero potential in D at fixed energy E,  $G^+(x, k)$  is defined in (2.9),  $k_{\perp} = (-k_2, k_1)$  for  $k = (k_1, k_2)$ , dy,  $d\xi$  denote the standard Euclidean measure on the boundary  $\partial D$  and  $d\theta$  denotes the standard Euclidean measure on  $S^1$ .

Then, in order to obtain  $h_{\pm}$ , it is sufficient to use the following formula (see [12] for n=1 and Section 4 for  $n \geq 2$ ):

$$h_{\pm}(k,l) = \frac{1}{(2\pi)^2} \int_{\partial D} \int_{\partial D} e^{-ilx} (\Phi - \Phi_0)(x,y,E) \psi_{\pm}(y,k) dy dx, \quad (k,l) \in \mathscr{M}_E.$$

**Algorithm 2**  $(f \longrightarrow h_{\pm})$ . Starting from f on  $\mathcal{M}_E$  (for E sufficiently large), one directly recovers  $h_{\pm}$  solving the following integral equation (see [13, 15] for n = 1 and Section 4 for  $n \geq 2$ ):

(3.5)
$$h_{\pm}(\lambda, \lambda', E) - \pi i \int_{T} f(\lambda'', \lambda', E) \chi_{+} \left( \pm i \left( \frac{\lambda}{\lambda''} - \frac{\lambda''}{\lambda} \right) \right) h_{\pm}(\lambda, \lambda'', E) |d\lambda''|$$

$$= f(\lambda, \lambda', E), \qquad (\lambda, \lambda') \in T \times T.$$

Algorithms 1 and 2  $(h_{\pm} \longrightarrow V_{appr})$ . We begin with the construction of  $\tilde{\mu}^+$ , an approximation to  $\mu^+$ ; this is done by solving the following integral equation arising from the non-local Riemann-Hilbert problem (2.21)-(2.28) for  $\mu$  in the approximation that  $b \equiv 0$  at fixed E (see [15] for n = 1 and Section 4 for  $n \geq 2$ ):

$$(3.6) \quad \tilde{\mu}^+(z,\lambda,E) + \int_T \tilde{\mu}^+(z,\lambda',E) B(\lambda,\lambda',z,E) |d\lambda'| = I, \ \lambda \in T, \ z \in \mathbb{C},$$

where E is sufficiently large and

$$(3.7)$$

$$B(\lambda, \lambda', z, E) = \frac{1}{2} \int_{T} h_{-}(\zeta, \lambda', z, E) \chi_{+} \left( -i \left( \frac{\zeta}{\lambda'} - \frac{\lambda'}{\zeta} \right) \right) \frac{d\zeta}{\zeta - \lambda(1 - 0)}$$

$$- \frac{1}{2} \int_{T} h_{+}(\zeta, \lambda', z, E) \chi_{+} \left( i \left( \frac{\zeta}{\lambda'} - \frac{\lambda'}{\zeta} \right) \right) \frac{d\zeta}{\zeta - \lambda(1 + 0)},$$

where  $\chi_+$ ,  $h_{\pm}$  are defined in (2.24), (2.25). Then one can obtain an approximation  $\tilde{\mu}_-$  to  $\mu_-$  via (2.23), used as follows:

(3.8) 
$$\tilde{\mu}_{-}(z,\lambda,E) = \tilde{\mu}^{+}(z,\lambda,E)$$
  
  $+ \pi i \int_{T} \tilde{\mu}^{+}(z,\lambda'',E) \chi_{+} \left(-i\left(\frac{\lambda}{\lambda''} - \frac{\lambda''}{\lambda}\right)\right) h_{-}(\lambda,\lambda'',z,E) |d\lambda''|,$ 

for  $\lambda \in T$ ,  $z \in \mathbb{C}$ . Finally, the approximate potential  $V_{appr}(\cdot, E)$  can be obtained using the following formula (see [14, 15] for n = 1 and Section 4 for  $n \geq 2$ ):

$$(3.9) V_{appr}(z,E) = 2iE^{1/2} \frac{\partial}{\partial z} \left( \frac{1}{2\pi i} \int_{T} \tilde{\mu}_{-}(z,\zeta,E) i\zeta |d\zeta| \right).$$

The approximate potential  $V_{appr}$  depends in a non-linear way on  $\Phi(E)$  in Algorithm 1 and on f on  $\mathcal{M}_E$  in Algorithm 2, in spite of the fact that both algorithms are based on solving linear integral equation. In the linear approximation near zero potential, the following formulas hold:

(3.10) 
$$h_{\pm}(k,l) \approx \frac{1}{(2\pi)^2} \int_{\partial D} \int_{\partial D} e^{i(-lx+ky)} (\Phi - \Phi_0)(x,y,E) dx dy, \quad (k,l) \in \mathcal{M}_E,$$

for linearised Algorithm 1;

$$(3.11) h_{+}(\lambda, \lambda', E) \approx f(\lambda, \lambda', E), \lambda, \lambda' \in T,$$

for linearised Algorithm 2;

(3.12) 
$$V_{appr}(z, E) \approx \frac{1}{\pi} E^{1/2} \int_{T} w(z, \lambda, E) i\lambda |d\lambda|,$$

where  $z \in D$  for linearised Algorithm 1 and  $z \in \mathbb{C}$  for linearised Algorithm 2, and

(3.13) 
$$w(z,\lambda,E) = \frac{\partial}{\partial z} \left( \pi i \int_{T} \exp \left[ -\frac{i}{2} E^{1/2} \left( \lambda \bar{z} + \frac{z}{\lambda} - \lambda' \bar{z} - \frac{z}{\lambda'} \right) \right] \times \operatorname{sign} \left( -i \left( \frac{\lambda}{\lambda'} - \frac{\lambda'}{\lambda} \right) \right) h_{\pm}(\lambda,\lambda',E) |d\lambda'| \right),$$

for  $z \in \mathbb{C}$ ,  $\lambda \in T$ , E > 0.

4. Derivation of some formulas and equations of Section 2 and 3 for the matrix case

The following formula and equations will be useful:

(4.1) 
$$\psi_{\gamma}(x,k) = \psi^{+}(x,k) + 2\pi i \int_{\mathbb{R}^{2}} \psi^{+}(x,l) \delta(l^{2} - k^{2}) \chi_{+}((l-k)\gamma) h_{\gamma}(k,l) dl,$$
(4.2) 
$$h_{\gamma}(k,l) = f(k,l) + 2\pi i \int_{\mathbb{R}^{2}} f(m,l) \delta(m^{2} - k^{2}) \chi_{+}((m-k)\gamma) h_{\gamma}(k,m) dm,$$

for  $\gamma \in S^1$ ,  $x, k, l \in \mathbb{R}^2$ ,  $k^2 = E \in \mathbb{R}_+$  sufficiently large,

(4.3) 
$$\frac{\partial \mu}{\partial \bar{k}_j}(x,k) = -2\pi \int_{\mathbb{R}^2} \xi_j e^{i\xi x} \mu(x,k+\xi) H(k,-\xi) \delta(\xi^2 + 2k\xi) d\xi,$$

(4.4) 
$$\frac{\partial H}{\partial \bar{k}_j}(k,p) = -2\pi \int_{\mathbb{R}^2} \xi_j H(k+\xi,p+\xi) H(k,-\xi) \delta(\xi^2 + 2k\xi) d\xi,$$

for  $j = 1, 2, k \in \mathbb{C}^2 \setminus \mathbb{R}^2$ ,  $k^2 = E \in \mathbb{R}_+$  sufficiently large,  $x, p \in \mathbb{R}^2$ , where H(k, p) = h(k, k - p),  $\delta$  is the Dirac delta and the other functions were already defined in Section 2. Formula (4.1) and equations (4.2)-(4.4) are proved in [7, 11] for the scalar case: the proof can be straightforwardly generalized to the matrix case, where one only has to pay attention to the order of the terms (which is indeed different from the formulation given in the quoted papers, but coherent with similar results obtained in [19]).

Now formula (2.23) follows directly from (4.1), (2.5), (2.25) using notations (2.15); equation (2.26) follows from (4.3) taking into account (2.14), (2.27) and notations (2.15). In addition, equation (3.5) is a direct consequence of (4.2) (with  $\gamma = \pm \hat{k}_{\perp}$ ) using notations (2.15).

Formula (2.30) follows from (2.29), (2.28), (2.26), (2.22), (2.21) (these can be proved exactly as in the scalar case) and the Cauchy–Pompeiu formula

$$(4.5) u(\lambda) = \frac{1}{2\pi i} \int_{\partial \mathcal{D}} u(\zeta) \frac{d\zeta}{\zeta - \lambda} - \frac{1}{\pi} \int_{\mathcal{D}} \frac{\partial u(\zeta)}{\partial \bar{\zeta}} \frac{d\text{Re}\zeta \, d\text{Im}\zeta}{\zeta - \lambda}, \quad \lambda \in \mathcal{D},$$

for any sufficiently regular  $M_n(\mathbb{C})$ -valued u in  $\mathcal{D}$ , where  $\partial \mathcal{D}$  is sufficiently regular. In addition, formula (3.9) is just formula (2.30) without the first term in the sum.

Equation (3.6) is an approximation of the following (exact) equation for  $u^+$ :

$$(4.6) \qquad \mu^{+}(z,\lambda,E) + \int_{T} \mu^{+}(z,\lambda',E)B(\lambda,\lambda',z,E)|d\lambda'| = I + \varphi(z,\lambda,E),$$

for  $\lambda \in T$ ,  $z \in \mathbb{C}$ , where

(4.7) 
$$\varphi(z,\lambda,E) = -\frac{1}{\pi} \int_{\mathbb{C}} \mu\left(z, -\frac{1}{\overline{\zeta}}, E\right) r(\zeta, z, E) \frac{d\text{Re}\zeta \, d\text{Im}\zeta}{\zeta - \lambda}.$$

The derivation of (4.6) can be found in [15] for the scalar case and its generalisation to the matrix case is straightforward (paying attention to the order of the terms).

Formula (3.3) is a result of [7], while formulas (3.1), (3.2), (3.4) are results of [12] for the scalar case and can be proved for the matrix case following the scheme of [17], where similar formulas appear.

#### 5. Function spaces and some estimates

We introduce some function spaces, which will be useful to prove the high-energy convergence of our algorithms. For  $m \in \mathbb{N}$ ,  $\varepsilon > 0$  we consider

$$W^{m,1}(\mathbb{R}^2, M_n(\mathbb{C})) = \{u : \partial^k u \in L^1(\mathbb{R}^2, M_n(\mathbb{C})) \text{ for } |k| \leq m\},$$

$$W^{m,1}_{\varepsilon}(\mathbb{R}^2, M_n(\mathbb{C})) = \{u : \varkappa^{\varepsilon} \partial^k u \in L^1(\mathbb{R}^2, M_n(\mathbb{C})) \text{ for } |k| \leq m\},$$

$$(\varkappa^{\varepsilon} u)(x) = (1 + |x|^2)^{\varepsilon/2} u(x), \qquad k \in (\mathbb{N} \cup 0)^2, \qquad |k| = k_1 + k_2,$$

$$\partial^k = \partial_1^{k_1} \partial_2^{k_2}, \qquad \partial_j = \frac{\partial}{\partial x_j};$$

for  $\alpha \in ]0,1]$ ,  $s \in \mathbb{R}$  we consider

$$C^{\alpha,s}(\mathbb{R}^2, M_n(\mathbb{C})) = \{ u : ||u||_{\alpha,s} < \infty \},$$

where

$$||u||_{\alpha,s} = ||\varkappa^{s}u||_{\alpha}$$

$$||w||_{\alpha} = \sup_{p,\xi \in \mathbb{R}^{2}, |\xi| \le 1} \left( |w(p)| + \frac{|w(p+\xi) - w(p)|}{|\xi|^{\alpha}} \right),$$

$$(\varkappa^{s}u)(p) = (1+|p|^{2})^{s/2}u(p), \qquad |u(p)| = \max_{1 \le i,j \le n} |u_{ij}(p)|;$$

in addition we consider  $\mathcal{H}_{\alpha,s}(\mathbb{R}^2, M_n(\mathbb{C}))$ , defined as the closure of  $C_0^{\infty}(\mathbb{R}^2, M_n(\mathbb{C}))$  (the space of infinitely smooth functions with compact support) in  $\|\cdot\|_{\alpha,s}$ .

Let

(5.1) 
$$\widehat{V}(p) = \frac{1}{(2\pi)^2} \int_{\mathbb{R}^2} e^{ipx} V(x) dx, \qquad p \in \mathbb{R}^2.$$

If a matrix-valued potential V satisfies

(5.2) 
$$V \in W_{\varepsilon}^{m,1}(\mathbb{R}^2, M_n(\mathbb{C})) \text{ for some } \varepsilon > 0, \ m \in \mathbb{N},$$

then

(5.3) 
$$\widehat{V} \in \mathcal{H}_{\alpha,s}(\mathbb{R}^2, M_n(\mathbb{C})), \quad \alpha \in ]0,1], \quad s \in \mathbb{R}_+,$$

where  $\alpha = \min(1, \varepsilon), s = m$ . Let

(5.4) 
$$\Sigma(r) = (1 - r)^{-1}r.$$

We have the following results:

**Proposition 5.1.** Let the condition (5.3) be valid. Then

$$(5.5a) |f(k,l) - \widehat{V}(k-l)| \le \Sigma(r) ||\widehat{V}||_{\alpha,s} (1 + |k-l|^2)^{-s/2},$$

(5.5b) 
$$|H_{\gamma}(k,p) - \widehat{V}(p)| \le \Sigma(r) ||\widehat{V}||_{\alpha,s} (1+p^2)^{-s/2},$$

for 
$$r = |k|^{-\sigma} c_1(\alpha, s, \sigma, n) \|\widehat{V}\|_{\alpha, s} < 1, \ k, l, p \in \mathbb{R}^2, \ \gamma \in S^1, \ k^2 \ge 1,$$

(5.5c) 
$$|H(k,p) - \widehat{V}(p)| \le \Sigma(r) \|\widehat{V}\|_{\alpha,s} (1+p^2)^{-s/2},$$

for  $r = |\operatorname{Re} k|^{-\sigma} c_1(\alpha, s, \sigma, n) \|\widehat{V}\|_{\alpha, s} < 1, \ k \in \mathbb{C}^2 \setminus \mathbb{R}^2, \ p \in \mathbb{R}^2, \ \mathbb{R} \ni k^2 \ge 1.$ In particular

(5.6a) 
$$|f(k,l)| \le 2||\widehat{V}||_{\alpha,s} (1+|k-l|^2)^{-s/2}, \quad k,l \in \mathbb{R}^2,$$

(5.6b) 
$$|H_{\gamma}(k,p)| \le 2||\widehat{V}||_{\alpha,s}(1+p^2)^{-s/2}, \quad k, p \in \mathbb{R}^2, \quad \gamma \in S^1,$$

(5.6c) 
$$|H(k,p)| \le 2\|\widehat{V}\|_{\alpha,s} (1+p^2)^{-s/2}, \qquad k \in \mathbb{C}^2 \setminus \mathbb{R}^2, \qquad p \in \mathbb{R}^2$$

for 
$$k^2 \ge E_1 = \max(1, (2c_1(\alpha, s, \sigma, n) \| \widehat{V} \|_{\alpha, s})^{2/\sigma})$$
, where  $H_{\gamma}(k, l) = h_{\gamma}(k, k - l)$ ,  $0 < \alpha < 1$ ,  $s > 0$ ,  $0 < \sigma < \min(1, s)$ .

**Lemma 5.2.** Under condition (5.3), we have the following estimates:

(5.7a)

$$\left|\mu_{\gamma}(x,k) - I\right| + \left|\frac{\partial \mu_{\gamma}(x,k)}{\partial x_{1}}\right| + \left|\frac{\partial \mu_{\gamma}(x,k)}{\partial x_{2}}\right| \leq |k|^{-\sigma} c_{2}(\alpha,s,\sigma,n) \|\widehat{V}\|_{\alpha,s},$$

for  $x = (x_1, x_2) \in \mathbb{R}^2$ ,  $k \in \mathbb{R}^2$ ,  $\gamma \in S^1$ ,

(5.7b)

$$|\mu(x,k) - I| + \left| \frac{\partial \mu(x,k)}{\partial x_1} \right| + \left| \frac{\partial \mu(x,k)}{\partial x_2} \right| \le |\operatorname{Re} k|^{-\sigma} c_2(\alpha, s, \sigma, n) \|\widehat{V}\|_{\alpha,s},$$

for 
$$x = (x_1, x_2) \in \mathbb{R}^2$$
,  $k \in \mathbb{C}^2 \setminus \mathbb{R}^2$ ,  $k^2 \geq E_1(\alpha, s, \sigma, n, ||\widehat{V}||_{\alpha, s})$ , where  $0 < \alpha < 1, s > 1, 0 < \sigma < \min(1, s - 1)$ .

Proposition 5.1 and Lemma 5.2 for the scalar case (n = 1) were given in [15] and their generalisation to the matrix case  $(n \ge 2)$  is straightforward.

## 6. Lipschitz stability and rapid convergence of Algorithms 1 and 2 for $E \to +\infty$

We present here the main results concerning stability and convergence of our algorithms.

**Theorem 6.1** (Stability and convergence of Algorithm 1). Let  $V \in W^{m,1}(\mathbb{R}^2, M_n(\mathbb{C}))$ ,  $m \geq 3$ , supp  $V \subset D$  and let  $\Phi(E)$  be the Dirichlet-to-Neumann operator of (1.6) at fixed energy E, where  $E \geq E_2(\alpha, s, \sigma, n, \|\widehat{V}\|_{\alpha,s})$ ,  $0 < \alpha \leq 1$ , s = m,  $0 < \sigma < 1$  and E is not a Dirichlet eigenvalue of  $-\Delta + V$  and  $-\Delta$  in D. Then V is reconstructed from  $\Phi(E)$  with Lipschitz stability via Algorithm 1 up to  $O(E^{-(m-2)/2})$  in the uniform norm as  $E \to +\infty$ .

**Theorem 6.2** (Stability and convergence of Algorithm 2). Let V satisfy (5.2), for  $m \geq 3$ , and let f be the scattering amplitude of (1.8) at fixed energy  $E \geq E_2(\alpha, s, \sigma, n, \|\widehat{V}\|_{\alpha,s})$ , where  $\alpha = \min(1, \varepsilon)$ , s = m and  $0 < \sigma < 1$ . Then V is reconstructed from f on  $\mathcal{M}_E$  with Lipschitz stability via Algorithms 2 up to  $O(E^{-(m-2)/2})$  in the uniform norm as  $E \to +\infty$ .

The constant  $E_2$  of Theorems 6.1 and 6.2 is precisely stated in Remark 6.3. The Lipschitz stability of Theorems 6.1 and 6.2 is specified in the proofs of these theorems and is summarized in Remarks 6.4 and 6.5. The error term  $O(E^{-(m-2)/2})$  of Theorems 6.1 and 6.2 is made explicit in formula (6.9).

Similarly with the presentation of Algorithms 1 and 2 in section 3, we separate the proofs of Theorems 6.1 and 6.2 in several steps.

Proof of Theorem 6.1 ( $\Phi(E) \longrightarrow h_{\pm}$ ). We have that equation (3.1) is a Fredholm linear integral equation of second kind for  $\psi_{\pm}|_{\partial D} \in L^2(\partial D)$ , which is uniquely solvable with precise data  $\Phi - \Phi_0$  (the proof of the latter fact is the same as in the scalar case; see [12]). Therefore the reconstruction of  $\psi_{\pm}$  via (3.1) is Lipschitz stable, with respect to small errors in  $\Phi - \Phi_0$  (in the  $L^2$  norm of the Schwartz kernel),

As a corollary, the reconstruction of  $h_{\pm}$  in  $L^2(\mathcal{M}_E)$  from  $\Phi(E) - \Phi_0(E)$  via equation (3.1) and formula (3.4) is also Lipschitz stable.

Proof of Theorem 6.2  $(f \longrightarrow h_{\pm})$ . Estimates (5.6) and notations (2.15) give

(6.1a) 
$$|f(\lambda, \lambda', E)| \le 2||\hat{V}||_{\alpha, s} (1 + E|\lambda - \lambda'|^2)^{-s/2}, \quad \lambda, \lambda' \in T,$$

(6.1b) 
$$||f||_{L^2(T\times T)} \le c_3 n ||\widehat{V}||_{\alpha,s} E^{-1/4},$$

for  $E \geq E_1$ ,  $\alpha = \min(1, \varepsilon)$ , s = m. Now, under the assumptions of Theorem 6.2, integral equation (3.5) is uniquely solvable for  $h_{\pm}(\lambda, \cdot, E) \in L^2(T)$  for  $\lambda \in T$ ,  $E \geq E_1$  (this is a consequence of the unique solvability of

integral equation (2.2) for  $E \geq E_1$ ). In addition, by estimate (6.1b), for  $E \geq \max(E_1, (\pi c_3 n \| \hat{V} \|_{\alpha,s})^4)$ , equation (3.5) is uniquely solvable for  $h_{\pm}(\lambda, \cdot, E) \in L^2(T)$ ,  $\lambda \in T$ , and for  $h_{\pm}(\cdot, \cdot, E) \in L^2(T \times T)$  by the method of successive approximations. This implies the Lipschitz stability of the reconstruction of  $h_{\pm}$  on  $T \times T$  from f on  $T \times T$ , with respect to small errors in the  $L^2$  norm.

Proof of Theorems 6.1 and 6.2  $(h_{\pm} \longrightarrow V_{appr})$ . The proof follows as in the scalar case (that was treated in [15]), except for the order of the terms in formulas and integral equations.

Estimates (5.6), formula (2.11) and notations (2.15) give

(6.2a) 
$$|h_{\pm}(\lambda, \lambda', E)| \le 2\|\widehat{V}\|_{\alpha, s} (1 + E|\lambda - \lambda'|^2)^{-s/2}, \quad \lambda, \lambda' \in T,$$

(6.2b) 
$$||h_{\pm}||_{L^2(T\times T)} \le c_3 n ||\widehat{V}||_{\alpha,s} E^{-1/4},$$

for  $E \geq E_1$ , s = m and

(6.3)  $0 < \alpha \le 1$  for Theorem 6.1 and  $\alpha = \min(1, \varepsilon)$  for Theorem 6.2.

We define the integral operator B(z, E) as

(6.4) 
$$(B(z, E)u)(\lambda) = \int_{T} u(\lambda')B(\lambda, \lambda', z, E)|d\lambda'|,$$

for  $\lambda \in T$ , where  $B(\lambda, \lambda', z, E)$  is defined in (3.7) and u is a test matrix function. The following decomposition holds

(6.5) 
$$B(z,E) = C_{+}Q_{-}(z,E) - C_{-}Q_{+}(z,E),$$

where

(6.6) 
$$(C_{\pm}u)(\lambda) = \frac{1}{2\pi i} \int_{T} \frac{u(\zeta)}{\zeta - \lambda(1 \mp 0)} d\zeta,$$

(6.7) 
$$(Q_{\pm}u)(\lambda) = \pi i \int_{T} u(\lambda') \chi_{+} \left( \pm i \left( \frac{\lambda}{\lambda'} - \frac{\lambda'}{\lambda} \right) \right) h_{\pm}(\lambda, \lambda', z, E) |d\lambda'|,$$

 $z \in \mathbb{C}$ ,  $\lambda \in T$ ,  $\chi_+$ ,  $h_{\pm}$  are defined in (2.24) and (2.25) and u is a test matrix function. Thanks to (6.2), (6.5) and properties of the Cauchy projectors  $C_{\pm}$  (see [15] for more details), B(z, E) satisfies the estimates

(6.8a) 
$$||B(z,E)u||_{L^{2}(T)} \le c_{4}n||\widehat{V}||_{\alpha,s}E^{-1/4}||u||_{L^{2}(T)},$$

(6.8b) 
$$\left\| \frac{\partial}{\partial z} B(z, E) u \right\|_{L^2(T)} \le c_4 n \|\widehat{V}\|_{\alpha, s} E^{-1/4} \|u\|_{L^2(T)},$$

for  $z \in \mathbb{C}$ ,  $E \geq E_1$ , s = m,  $\alpha$  as in (6.3).

Now by estimate (6.8a), for  $E \ge \max(E_1, (c_4 n \|\hat{V}\|_{\alpha,s})^4)$ , integral equation (3.6) is uniquely solvable for  $\tilde{\mu}^+(z,\cdot,E) \in L^2(T)$ , at fixed  $z \in \mathbb{C}$ , by the

method of successive approximations. This implies the Lipschitz stability of the reconstruction of  $\tilde{\mu}^+(z,\cdot,E)$  on T, at fixed  $z \in \mathbb{C}$ , from  $h_{\pm}$  on  $T \times T$  with respect to small errors in the  $L^2$  norm.

Proof of Theorems 6.1 and 6.2 ( $V_{appr} \longrightarrow V$ ). Our high-energy convergence estimate is as follows:

(6.9) 
$$|V(z) - V_{appr}(z, E)| \le c_5 n \|\widehat{V}\|_{\alpha, s} E^{-(s-2)/2},$$

where  $z \in \mathbb{C}$ ,  $E \geq E_2(\alpha, s, \sigma, n, \|\widehat{V}\|_{\alpha,s})$ ,  $\alpha$  as in (6.3), s = m,  $0 < \sigma < 1$  (see [15] for complete details). This estimate follows from (2.23), (2.30), (3.6)–(3.9), (6.2b) and the following estimates (whose proofs for n = 1 can be found in [15]):

(6.10) 
$$\left| 2iE^{1/2} \frac{\partial}{\partial z} \left( \int_{D_{-}} \mu(z, -\frac{1}{\overline{\zeta}}, E) r(\zeta, z, E) d\operatorname{Re} \zeta d\operatorname{Im} \zeta \right) \right| \leq c_6 n \|\widehat{V}\|_{\alpha, s} E^{-(s-2)/2},$$

(6.11) 
$$\|\mu^{+}(z,\cdot,E) - \tilde{\mu}^{+}(z,\cdot,E)\|_{L^{2}(T,M_{n}(\mathbb{C}))} \le c_{7}n\|\hat{V}\|_{\alpha,s}E^{-s/2},$$

(6.12) 
$$\left\| \frac{\partial \mu^+}{\partial z}(z,\cdot,E) - \frac{\partial \tilde{\mu}^+}{\partial z}(z,\cdot,E) \right\|_{L^2(T,M_n(\mathbb{C}))} \leq c_7 n \|\widehat{V}\|_{\alpha,s} E^{-(s-1)/2},$$

for 
$$z \in \mathbb{C}$$
,  $s = m \ge 3$ ,  $E \ge E_2(\alpha, s, \sigma, n, \|\widehat{V}\|_{\alpha, s})$ ,  $\alpha$  as in (6.3).

**Remark 6.3.** The constant  $E_2$  of Theorems 6.1 and 6.2 can be fixed as some constant such that  $E \geq E_2$  implies that

$$E \ge E_1, \qquad |\mu(z,\lambda,E)| \le 2, \qquad \left|\frac{\partial}{\partial z}\mu(z,\lambda,E)\right| \le 1,$$
  
$$\|B(z,E)\|_{L^2(T)}^{op} \le \frac{1}{2}, \qquad \left\|\frac{\partial}{\partial z}B(z,E)\right\|_{L^2(T)}^{op} \le \frac{1}{2},$$

for  $z \in \mathbb{C}$ ,  $\lambda \in \mathbb{C}$ , where  $\mu$  and B are estimated in (5.7) and (6.8).

Now, let  $\Phi_{V,0}(x,y,E)$ ,  $x,y \in \partial D$ , denote the Schwartz kernel of the operator  $\Phi(E) - \Phi_0(E)$  considered as precise data for Problem 1. Let  $\Phi'_{V,0}$  denote  $\Phi_{V,0}$  with some small errors (for the case of Problem 1) and f' denote f with some small errors (for the case of Problem 2). Let  $V'_{appr}$  denote  $V_{appr}$  reconstructed from  $\Phi'_{V,0}$  via Algorithm 1 (for Problem 1) and from f' via Algorithm 2 (for Problem 2).

The Lipschitz stability of Theorems 6.1 and 6.2 is summarized in the following remarks:

Remark 6.4. Let the assumptions of Theorem 6.1 hold and let

(6.13) 
$$\delta = \|\Phi'_{V,0}(\cdot,\cdot,E) - \Phi_{V,0}(\cdot,\cdot,E)\|_{L^2(\partial D \times \partial D)} \le \delta_1(V,E,D,n).$$

Then

(6.14) 
$$\varepsilon = \|V'_{appr} - V_{appr}\|_{L^{\infty}(D)} \le \eta_1(V, E, D, n)\delta.$$

Here  $\delta_1$  and  $\eta_1$  are some positive constants summarizing the Lipschitz stability of Algorithm 1. In particular,

$$\delta_1(V, E, D, n) \ge \delta_1^0,$$

(6.16) 
$$\eta_1(V, E, D, n) \le \eta_1^0 E,$$

as  $\|\Phi_{V,0}(\cdot,\cdot,E)\|_{L^2(\partial D\times\partial D)}\to 0$ , for some positive (sufficiently small)  $\delta_1^0$  and (sufficiently big)  $\eta_1^0$ , where  $\delta_1^0$  and  $\eta_1^0$  are independent of V and E for fixed D and n.

Remark 6.5. Let the assumptions of Theorem 6.2 hold and let

(6.17) 
$$\delta = \|f - f'\|_{L^2(\mathscr{M}_E)} \le \delta_2(V, E, n).$$

Then

(6.18) 
$$\varepsilon = \|V_{appr} - V'_{appr}\|_{L^{\infty}(\mathbb{R}^2)} \le \eta_2(V, E, n)\delta.$$

Here  $\delta_2$  and  $\eta_2$  are suitable constants summarizing the Lipschitz stability of Algorithms 2. In particular,

$$\delta_2(V, E, n) \ge \delta_2^0,$$

(6.20) 
$$\eta_2(V, E, n) \le \eta_2^0 E,$$

as  $||f||_{L^2(\mathcal{M}_E)} \to 0$ , for some positive (sufficiently small)  $\delta_2^0$  and (sufficiently big)  $\eta_2^0$ , where  $\delta_2^0$  and  $\eta_2^0$  are independent of V and E for fixed n.

Note that in Remark 6.5, the norm  $\|\cdot\|_{L^2(\mathscr{M}_E)}$  is identified with  $\|\cdot\|_{L^2(T\times T)}$ . The property that  $\|f\|_{L^2(\mathscr{M}_E)}\to 0$ , mentioned in Remark 6.5, is fulfilled, in particular, for  $E\to +\infty$ , as a consequence of estimate (6.1b). On the contrary, the property that  $\|\Phi_{V,0}(\cdot,\cdot,E)\|_{L^2(\partial D\times \partial D)}\to 0$ , mentioned in Remark 6.4, is not fulfilled for  $E=E_j, j\to\infty$ , for any sequence  $\{E_j\}_{j\in\mathbb{N}}$  of positive real numbers such that  $E_j\to +\infty$  as  $j\to\infty$ , if  $V\not\equiv 0$ . In this connection, our high-energy conjecture is that

(6.21) 
$$\sup_{j \in \mathbb{N}} \|\Phi_{V,0}(\cdot,\cdot,E_j)\|_{L^2(\partial D \times \partial D)} < +\infty,$$

for some  $\{E_j\}_{j\in\mathbb{N}}$  dependent on V, where  $E_j \to +\infty$  as  $j \to \infty$ .

Note that the E factor in the right side of (6.16) and of (6.20) is related with the choice of the  $L^2$  norm for estimates of the inverse problem data.

For example, for Algorithm 2, at least in the linear approximation (3.11)-(3.13), this factor disappear if  $\|\cdot\|_{L^2(\mathscr{M}_E)}$  is replaced by  $\|\cdot\|_{L^\infty_s(\mathscr{M}_E)}$ , s=m, where

(6.22) 
$$||u||_{L_s^{\infty}(\mathcal{M}_E)} = \sup_{(\lambda,\lambda')\in T\times T} (1+E|\lambda-\lambda'|^2)^{s/2} |u(\lambda,\lambda')|.$$

#### References

- [1] Baykov, S.V., Burov, V.A., Sergeev, S.N., *Mode Tomography of Moving Ocean*, Proc. of the 3rd European Conference on Underwater Acoustics, 1996, 845–850.
- [2] Bogatyrev, A.V., Burov, V.A., Morozov, S.A., Rumyantseva, O.D., Sukhov, E.G. Numerical realization of algorithm for exact solution of two-dimensional monochromatic inverse problem of acoustical scattering, Acoust. Imaging 25, 2000, 65–70.
- [3] Bukhgeim, A. L., Recovering a potential from Cauchy data in the two-dimensional case, J. Inverse Ill-Posed Probl. 16, 2008, no. 1, 19–33.
- [4] Burov, V. A., Alekseenko, N. V., Rumyantseva, O. D., Multifrequency generalization of the Novikov algorithm for the two-dimensional inverse scattering problem, Acoustical Physics 55, 2009, no. 6, 843–856.
- [5] Dubrovin, B. A., Krichever, I. M., Novikov, S. P., The Schrödinger equation in a periodic field and Riemann surfaces, Dokl. Akad. Nauk SSSR 229, 1976, no. 1, 15– 18
- [6] Faddeev, L. D., Growing solutions of the Schrödinger equation, Dokl. Akad. Nauk SSSR 165, 1965, no. 3, 514–517.
- [7] Faddeev, L. D., The inverse problem in the quantum theory of scattering. II, Journal of Mathematical Sciences 5, 1976, no. 3, 334–396.
- [8] Gel'fand, I.M., Some problems of functional analysis and algebra, Proc. Int. Congr. Math., Amsterdam, 1954, 253–276.
- [9] Grinevich, P. G., The scattering transform for the two-dimensional Schrödinger operator with a potential that decreases at infinity at fixed nonzero energy, (Russian) Uspekhi Mat. Nauk 55, 2000, no. 6(336), 3–70; translation in Russian Math. Surveys 55, 2000, no. 6, 1015–1083.
- [10] Grinevich, P. G., Novikov, R. G., Transparent potentials at fixed energy in dimension two. Fixed-energy dispersion relations for the fast decaying potentials Comm. Math. Phys. 174, 1995, no. 2, 409–446.
- [11] Henkin, G. M., Novikov, R. G., The  $\bar{\partial}$ -equation in the multidimensional inverse scattering problem, Russian Mathematical Surveys 42, 1987, no. 3, 109–180.
- [12] Novikov, R. G., Multidimensional inverse spectral problem for the equation  $-\Delta \psi + (v(x) Eu(x))\psi = 0$ , Funkt. Anal. i Pril. **22**, 1988, no. 4, 11–22 (in Russian); English Transl.: Funct. Anal. and Appl. **22**, 1988, 263–272.
- [13] Novikov, R. G., The inverse scattering problem on a fixed energy level for the twodimensional Schrödinger operator, J. Funct. Anal. 103, 1992, no. 2, 409–463.
- [14] Novikov, R. G., Rapidly converging approximation in inverse quantum scattering in dimension 2, Phys. Lett. A 238, 1998, no. 2-3, 73–78.
- [15] Novikov, R. G., Approximate solution of the inverse problem of quantum scattering theory with fixed energy in dimension 2, (Russian) Tr. Mat. Inst. Steklova 225, 1999, Solitony Geom. Topol. na Perekrest., 301–318; translation in Proc. Steklov Inst. Math. 225, 1999, no. 2, 285–302.

- [16] Novikov, R. G., Santacesaria, M., A global stability estimate for the Gel'fand-Calderón inverse problem in two dimensions, J. Inverse Ill-Posed Probl. 18, 2010, 765–785.
- [17] Novikov, R. G., Santacesaria, M., Global uniqueness and reconstruction for the multichannel Gel'fand-Calderón inverse problem in two dimensions, Bulletin des Sciences Mathematiques, 2011, doi:10.1016/j.bulsci.2011.04.007.
- [18] Santacesaria, M., Global stability for the multi-channel Gel'fand-Calderón inverse problem in two dimensions, e-print arXiv:1102.5175.
- [19] Xiaosheng, L., Inverse scattering problem for the Schrödinger operator with external Yang-Mills potentials in two dimensions at fixed energy, Comm. Part. Diff. Eq. **30**, 2005, no.4-6, 451–482.
- (R. G. Novikov and M. Santacesaria) CNRS (UMR 7641), CENTRE DE MATHÉMATIQUES APPLIQUÉES, ÉCOLE POLYTECHNIQUE, 91128, PALAISEAU, FRANCE

 $\textit{E-mail address:} \verb| novikov@cmap.polytechnique.fr|, \verb| santacesaria@cmap.polytechnique.fr| \\$