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A Convergent Data Completion Algorithm Using Surface Integral Equations

Yosra Boukari $^{\mathrm{a,b}}$ and Houssem Haddar $^{\mathrm{b}*}$

^a ENIT-LAMSIN, Université de Tunis el-Manar, BP 37, 1002, Tunis, Tunisia

b INRIA Saclay Ile de France/CMAP Ecole Polytechnique, Route de Saclay, 91128, Palaiseau Cedex, France

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Abstract

We propose and analyze a data completion algorithm based on the representation of the solution in terms of surface integral operators to solve the Cauchy problem for the Helmholtz or the Laplace equations. The proposed method is non iterative and intrinsically handle the case of noisy and incompatible data. In order to cope with the ill-posedness of the problem, our formulation is compatible with standard regularization methods associated with linear ill posed inverse problems and leads to convergent scheme. We numerically validate our method with different synthetic examples using a Tikhonov regularization.

1 Introduction

We are concerned with the design of robust and quick numerical algorithms that solve the ill-posed Cauchy problem in the framework of the Helmholtz or the Laplace equations. More specifically, we would like to determine the Cauchy data associated with a solution of the Helmholtz or the Laplace equation in a part of the boundary of a bounded domain from the knowledge of the Cauchy data on the complementary part. We consider the setting of a Lipschitz domain whose boundary can be split in two parts: an exterior boundary and an interior boundary that we denote by Γ_1 and Γ_2 respectively. The inverse problem under concern consists in reconstructing the interior Cauchy data on Γ_1 by knowing the Cauchy data on the exterior boundary Γ_2 .

This problem is known to be severely ill-posed [16, 4] and only weak logarithmic conditional stability can be obtained in general (See for instance [8, 11, 17, 20, 27, 17] and references therein). On the numerical side, a large variety of methods has been proposed in the literature and it is hard to give an extensive overview of them. Let us however indicate some of them in order to put into perspective the method we shall introduce. A first class of method named quasi-reversibility relate missing data to the boundary values of a solution to a stable fourth order problem. It was originally proposed by Lattes and Lions in 1967 [24] and then extensively studied in [22, 7, 9, 14, 11]. This method is not iterative. However, the numerical implementation requires solvers for a fourth order problem, which may be computationally complex and expensive. The extension to Helmholtz type equations has not been studied yet.

In the class of iterative methods, several approaches have been proposed for the Laplace operator by either using a Kozlov [23] alternating scheme type [21, 6, 12, 5] or by minimizing an appropriate functional [2, 1, 15]. The case of Helmholtz equation has been addressed in only a few works [3, 25].

[∗]Corresponding author. Email: haddar@cmap.polytechnique.fr

In this paper, we propose a new method to solve this inverse problem. Using the boundary integral representation of the solution along with the trace of the solution and its normal derivative on Γ_1 and Γ_2 , we derive a system of linear integral equation that is satisfied by the known and unknown Cauchy data. The solution of the inverse problem can be obtained by inverting the linear (matrix) operator that applies to the unknown Cauchy data. The main development here is the analysis of the structure of this operator. We observe for instance that it can be decomposed into an ill-posed component and a projection. This decomposition allows us to prove the injectivity of this operator and characterize the closure of its range. We then show that the system naturally handles the case of noisy and incompatible data, i.e. data that do not belong to the closure of the range of the operator to be inverted. Indeed, in our formulation the available Cauchy data are multiplied by an operator whose range closure is the same as for the operator to be inverted. This makes possible the use of classical regularization techniques for noisy data, such as the Tikhonov regularization. The main important consequence is a convergence result for the exact solution as the noise level goes to zero.

Some numerical results are then provided showing the efficiency of the method for different shapes of the domain, different values of the wave number and different noise levels. Note in addition that the implemented algorithm has the advantage to be fast as it does not rely on an iterative scheme.

Let us finally note that our approach shares similarities with the one proposed in [10] for the Laplace problem. However, for the latter, although also numerically efficient, one cannot ensure convergence towards the exact solution since it does not belong (in general) to the range of the used integral representation. For other conributions on solving the Cauchy problem using the integral representation and iterative methods, we refer the reader to [25, 18, 19].

The outline of the article is as follows. In Section 2 we present the problem and derive the linear equation to be inverted. In Section 3 we analyze the obtained system and associated regularization. In Section 4 we provide numerical experiments that validate the proposed approach.

2 A data completion algorithm based on integral equations

We start this section by introducing the Cauchy problem associated with the Helmholtz equation (or the Laplace equation if the wave number is zero). Let $\Omega_1 \subset \mathbb{R}^m$ be a bounded Lipschitz domain with connected complement and let $\Omega_2 \subset \mathbb{R}^m$ be a Lipschitz subdomain of Ω_1 . We shall assume that Ω_2 is strictly included in Ω_1 . We denote by Γ_1 and Γ_2 respectively the boundary of Ω_1 and Ω_2 and ν is the normal vector on $\Gamma_1 \cup \Gamma_2$ directed to the exterior of Ω_1 and Ω_2 (see Figure 1). Let $u \in H^1(\Omega_1 \setminus \Omega_2)$ solution of the following Helmholtz equation:

$$
\Delta u + k^2 u = 0 \quad \text{in } \Omega_1 \setminus \Omega_2,\tag{1}
$$

where k denotes the wave number. Our data completion problem consists in determining $(u, \partial_{\nu}u)_{|\Gamma_2}$ while knowing $(f, g) := (u, \partial_{\nu}u)_{|\Gamma_1}$.

Our algorithm is based on the inversion of the surface representation of u in terms of the boundary values $(u, \partial_{\nu}u)_{|\Gamma_1}$ and $(u, \partial_{\nu}u)_{|\Gamma_2}$. More precisely, since $u \in H^1(\Omega_1 \setminus \Omega_2)$ satisfies the Helmholtz equation in $\Omega_1 \setminus \Omega_2$, then u has the following representation

$$
u(x) = -DL_{\Gamma_1}(u_{|\Gamma_1}) + SL_{\Gamma_1}(\partial_\nu u_{|\Gamma_1}) + DL_{\Gamma_2}(u_{|\Gamma_2}) - SL_{\Gamma_2}(\partial_\nu u_{|\Gamma_2})
$$
\n(2)

for $x \in \Omega_1 \setminus \overline{\Omega_2}$ where

$$
SL_{\Gamma_i}\psi(x):=\int_{\Gamma_i}\psi(y)\Phi(x,y)ds(y),\quad DL_{\Gamma_i}\varphi(x):=\int_{\Gamma_i}\varphi(y)\frac{\partial\Phi(x,y)}{\partial\nu(y)}ds(y),
$$

Figure 1: Representation of the domain

for $x \in \mathbb{R}^m \setminus \Gamma_i$, $i = 1, 2$, with Φ the Green function defined for $k \neq 0$ by

$$
\Phi(x,y) := \begin{cases} \n\frac{1}{4\pi} \frac{e^{ik|x-y|}}{|x-y|} & m = 3, \\
\frac{i}{4} H_0^{(1)}(k|x-y|) & m = 2,\n\end{cases}
$$
\n(3)

for $x \neq y$ and $H_0^{(1)}$ being the Hankel function of the first kind of order zero and for $k = 0$, $\Phi(x, y)$ is given by

$$
\Phi(x,y) := \begin{cases} \n\frac{1}{4\pi} \frac{1}{|x-y|} & m = 3, \\
\frac{-1}{2\pi} \log(|x-y|) & m = 2.\n\end{cases}
$$
\n(4)

Considering the trace and the normal trace of u on Γ_1 and using the trace properties of double and single layer potentials, we obtain the following identities

$$
f = -K_{\Gamma_1}(f) + \frac{f}{2} + S_{\Gamma_1}(g) + DL_{\Gamma_2 \to \Gamma_1}(u_{|\Gamma_2}) - SL_{\Gamma_2 \to \Gamma_1}(\partial_\nu u_{|\Gamma_2}),
$$
(5)

$$
g = -T_{\Gamma_1}(f) + K'_{\Gamma_1}(g) + \frac{g}{2} + TL_{\Gamma_2 \to \Gamma_1}(u_{|\Gamma_2}) - DL'_{\Gamma_2 \to \Gamma_1}(\partial_\nu u_{|\Gamma_2}),
$$
(6)

where the boundary integral operators

$$
S_{\Gamma_i}: H^{-1/2}(\Gamma_i) \to H^{1/2}(\Gamma_i), \qquad K_{\Gamma_i}: H^{1/2}(\Gamma_i) \to H^{1/2}(\Gamma_i),
$$

$$
K'_{\Gamma_i}: H^{-1/2}(\Gamma_i) \to H^{-1/2}(\Gamma_i), \qquad T_{\Gamma_i}: H^{1/2}(\Gamma_i) \to H^{-1/2}(\Gamma_i),
$$

are defined for regular densities ψ and φ by

$$
S_{\Gamma_i}\psi(x) := \int_{\Gamma_i} \psi(y)\Phi(x,y)ds(y), \qquad \qquad K_{\Gamma_i}\varphi(x) := \int_{\Gamma_i} \varphi(y)\frac{\partial \Phi(x,y)}{\partial \nu(y)}ds(y),
$$

$$
K'_{\Gamma_i}\psi(x) := \int_{\Gamma_i} \psi(y)\frac{\partial \Phi(x,y)}{\partial \nu(x)}ds(y), \quad T_{\Gamma_i}\varphi(x) := \lim_{\epsilon \longrightarrow 0} \frac{\partial}{\partial \nu(x)} \int_{\Gamma_i, |x-y| > \epsilon} \varphi(y)\frac{\partial \Phi(x,y)}{\partial \nu(y)}ds(y),
$$

for $x \in \Gamma_i$. The boundary integral operators for $i, j = 1, 2$ and $i \neq j$

$$
SL_{\Gamma_i \to \Gamma_j}: H^{-1/2}(\Gamma_i) \longrightarrow H^{1/2}(\Gamma_j), \qquad DL_{\Gamma_i \to \Gamma_j}: H^{1/2}(\Gamma_i) \longrightarrow H^{1/2}(\Gamma_j),
$$

$$
DL'_{\Gamma_i \to \Gamma_j}: H^{-1/2}(\Gamma_i) \longrightarrow H^{-1/2}(\Gamma_j), \qquad TL_{\Gamma_i \to \Gamma_j}: H^{1/2}(\Gamma_i) \longrightarrow H^{-1/2}(\Gamma_j),
$$

are defined for regular densities ψ and φ by

$$
SL_{\Gamma_i \to \Gamma_j} \psi(x) := \int_{\Gamma_i} \psi(y) \Phi(x, y) ds(y), \qquad DL_{\Gamma_i \to \Gamma_j} \varphi(x) := \int_{\Gamma_i} \varphi(y) \frac{\partial \Phi(x, y)}{\partial \nu(y)} ds(y),
$$

$$
DL'_{\Gamma_i \to \Gamma_j} \psi(x) := \int_{\Gamma_i} \psi(y) \frac{\partial \Phi(x, y)}{\partial \nu(x)} ds(y), \quad TL_{\Gamma_i \to \Gamma_j} \varphi(x) := \frac{\partial}{\partial \nu(x)} \int_{\Gamma_i} \varphi(y) \frac{\partial \Phi(x, y)}{\partial \nu(y)} ds(y)
$$

for $x \in \Gamma_j$. We observe that, since $\Gamma_i \neq \Gamma_j$, the kernels of $DL_{\Gamma_i \to \Gamma_j}$, $SL_{\Gamma_i \to \Gamma_j}$, $DL'_{\Gamma_i \to \Gamma_j}$, $TL_{\Gamma_i \to \Gamma_j}$ are smooth and therefore these operators are compact. Equations (5) and (6) can be written in a matrix form as

$$
-\left[\begin{array}{cc} DL_{\Gamma_2 \to \Gamma_1} & -SL_{\Gamma_2 \to \Gamma_1} \\ TL_{\Gamma_2 \to \Gamma_1} & -DL'_{\Gamma_2 \to \Gamma_1} \end{array}\right] \left[\begin{array}{c} u_{|\Gamma_2} \\ \partial_{\nu} u_{|\Gamma_2} \end{array}\right] = \left[\begin{array}{cc} -K_{\Gamma_1} & S_{\Gamma_1} \\ -T_{\Gamma_1} & K'_{\Gamma_1} \end{array}\right] \left[\begin{array}{c} f \\ g \end{array}\right] - \frac{1}{2} \left[\begin{array}{c} f \\ g \end{array}\right].
$$

Using the Calderon projector $P_{\Gamma_1}: H^{1/2}(\Gamma_1) \times H^{-1/2}(\Gamma_1) \longrightarrow H^{1/2}(\Gamma_1) \times H^{-1/2}(\Gamma_1)$ defined by

$$
P_{\Gamma_1}\left[\begin{array}{c}f\\g\end{array}\right] := \left[\begin{array}{cc}-K_{\Gamma_1} & S_{\Gamma_1}\\-T_{\Gamma_1} & K'_{\Gamma_1}\end{array}\right] \left[\begin{array}{c}f\\g\end{array}\right] + \frac{1}{2}\left[\begin{array}{c}f\\g\end{array}\right],\tag{7}
$$

(we recall that $P_{\Gamma_1}^2 = P_{\Gamma_1}$, [26]) and the operator $T_{\Gamma_2 \to \Gamma_1} : H^{1/2}(\Gamma_2) \times H^{-1/2}(\Gamma_2) \longrightarrow H^{1/2}(\Gamma_1) \times$ $H^{-1/2}(\Gamma_1)$ defined by

$$
T_{\Gamma_2 \to \Gamma_1} \left[\begin{array}{c} \varphi \\ \psi \end{array} \right] := - \left[\begin{array}{cc} DL_{\Gamma_2 \to \Gamma_1} & -SL_{\Gamma_2 \to \Gamma_1} \\ TL_{\Gamma_2 \to \Gamma_1} & -DL'_{\Gamma_2 \to \Gamma_1} \end{array} \right] \left[\begin{array}{c} \varphi \\ \psi \end{array} \right],
$$

we obtain the following first equation

$$
T_{\Gamma_2 \to \Gamma_1} \left[\begin{array}{c} u_{|\Gamma_2} \\ \partial_\nu u_{|\Gamma_2} \end{array} \right] = (P_{\Gamma_1} - I_{\Gamma_1}) \left[\begin{array}{c} f \\ g \end{array} \right]
$$
 (8)

where I_{Γ_i} is the identity operator on $H^{1/2}(\Gamma_i) \times H^{-1/2}(\Gamma_i)$, $i = 1, 2$. Similarly, considering the trace and the normal trace of u on Γ_2 , we obtain a second relationship between (f, g) and $(u, \partial_{\nu}u)_{\Gamma_2}$, namely

$$
u_{|\Gamma_2} = -DL_{\Gamma_1 \to \Gamma_2}(f) + SL_{\Gamma_1 \to \Gamma_2}(g) + K_{\Gamma_2}(u_{|\Gamma_2}) + \frac{u_{|\Gamma_2}}{2} - S_{\Gamma_2}(\partial_\nu u_{|\Gamma_2})
$$

$$
\partial_\nu u_{|\Gamma_2} = -TL_{\Gamma_1 \to \Gamma_2}(f) + DL'_{\Gamma_1 \to \Gamma_2}(g) + T_{\Gamma_2}(u_{|\Gamma_2}) - K'_{\Gamma_2}(\partial_\nu u_{|\Gamma_2}) + \frac{\partial_\nu u_{|\Gamma_2}}{2}
$$

Therefore, we also have

$$
P_{\Gamma_2} \left[\begin{array}{c} u_{\Gamma_2} \\ \partial_\nu u_{\Gamma_2} \end{array} \right] = T_{\Gamma_2 \to \Gamma_1} \left[\begin{array}{c} f \\ g \end{array} \right]. \tag{9}
$$

with P_{Γ_2} and $T_{\Gamma_1 \to \Gamma_2}$ defined similarly to P_{Γ_1} and $T_{\Gamma_2 \to \Gamma_1}$. By grouping (8) and (9) we end up with the following system

$$
\begin{bmatrix} T_{\Gamma_2 \to \Gamma_1} \\ P_{\Gamma_2} \end{bmatrix} \begin{bmatrix} u_{|\Gamma_2} \\ \partial_\nu u_{|\Gamma_2} \end{bmatrix} = \begin{bmatrix} (P_{\Gamma_1} - I_{\Gamma_1}) \\ T_{\Gamma_1 \to \Gamma_2} \end{bmatrix} \begin{bmatrix} f \\ g \end{bmatrix}.
$$
 (10)

In order to simplify the notation, we introduce $A : H^{1/2}(\Gamma_2) \times H^{-1/2}(\Gamma_2) \longrightarrow H^{1/2}(\Gamma_1) \times$ $H^{-1/2}(\Gamma_1) \times H^{1/2}(\Gamma_2) \times H^{-1/2}(\Gamma_2)$ defined as

$$
A\begin{bmatrix} \varphi \\ \psi \end{bmatrix} = \begin{bmatrix} T_{\Gamma_2 \to \Gamma_1} \\ P_{\Gamma_2} \end{bmatrix} \begin{bmatrix} \varphi \\ \psi \end{bmatrix}
$$
 (11)

and $B: H^{1/2}(\Gamma_1) \times H^{-1/2}(\Gamma_1) \longrightarrow H^{1/2}(\Gamma_1) \times H^{-1/2}(\Gamma_1) \times H^{1/2}(\Gamma_2) \times H^{-1/2}(\Gamma_2)$ defined as

$$
B\left[\begin{array}{c} \varphi \\ \psi \end{array}\right] = \left[\begin{array}{c} (P_{\Gamma_1} - I_{\Gamma_1}) \\ T_{\Gamma_1 \to \Gamma_2} \end{array}\right] \left[\begin{array}{c} \varphi \\ \psi \end{array}\right]. \tag{12}
$$

Then our inverse problem can be simply written as seeking for $(\varphi, \psi) \in H^{1/2}(\Gamma_2) \times H^{-1/2}(\Gamma_2)$ satisfying

$$
A\left[\begin{array}{c} \varphi \\ \psi \end{array}\right] = B\left[\begin{array}{c} f \\ g \end{array}\right].
$$
 (13)

The ill-posedness of this problem can be seen in the compactness of the operator $T_{\Gamma_2 \to \Gamma_1}$ (all of its components are integral operators with smooth kernels). In order to apply a regularization strategy we should first investigate the injectivity of A and its adjoint A^* . As we shall see, this property fails for the adjoint but is replaced with $Range(A) = Range(B)$. The latter also ensures that the right hand side of (13) is always in the closure of the range of A for noisy data.

Remark 1. Let us notice that conceptually our method can be extended to the case where Γ_1 and Γ_2 are parts of the same boundary. However, the functional space setting has to modified to ensure proper continuity properties at the surface boundaries (see for instance [26] for trace spaces on open surfaces).

Remark 2. The method (as any boundary integral method) is not suited for inhomogeneous media. However, it can handle the case of piecewise constant backgrounds with known interfaces by incorporating the interfaces as parts of Γ_2 . Another alternative would be to (numerically) compute the Green function of the background, but this would be very costly.

3 Analysis of operators A and B

We first need to characterize the range and the kernel of $T_{\Gamma_i \to \Gamma_j}$, $i \neq j$. To simplify the notation we introduce the following spaces $(i = 1, 2)$:

$$
\mathcal{H}(\Omega_i) := \{ u \in H^1(\Omega_i), \Delta u + k^2 u = 0 \text{ in } \Omega_i \}
$$

$$
\tilde{\mathcal{H}}(\Omega_i) := \{ u \in H^1_{loc}(\mathbf{R}^m \setminus \Omega_1), \Delta u + k^2 u = 0 \text{ in } \mathbf{R}^m \setminus \Omega_i \text{ and } u \text{ satisfies the } R.C.\}
$$

where R.C. refers to the Sommerfeld radiation condition which is defined for $k \neq 0$ as follows

$$
\lim_{r=|x|\to+\infty} r^{\frac{m-1}{2}} (\partial_r u - iku) = 0,
$$
\n(14)

uniformly in all directions $\hat{x} = \frac{x}{|x|}$. In the Laplace case, i.e. $k = 0$, the R.C. simply refers here to $u(x)$ vanishing at $|x| = \infty$. The following classical notation will be used

$$
v^{\pm}_{|\Gamma_i}(x) := \lim_{h \to 0^+} v(x \pm h\nu) \quad \text{and} \quad \partial_{\nu} v^{\pm}_{|\Gamma_i}(x) := \lim_{h \to 0^+} \nu \cdot \nabla v(x \pm h\nu) \quad \text{for} \quad x \in \Gamma_i;
$$

$$
[v]_{|\Gamma_i} := v^{\pm}_{|\Gamma_i} - v^-_{|\Gamma_i} \quad \text{and} \quad [\partial_{\nu} v]_{|\Gamma_i} := \partial_{\nu} v^{\pm}_{|\Gamma_i} - \partial_{\nu} v^-_{|\Gamma_i}.
$$

Lemma 3. The following identities hold.

a)
$$
\operatorname{Ker}(T_{\Gamma_2 \to \Gamma_1}) = \{(u, \partial_\nu u)_{|\Gamma_2}; u \in \mathcal{H}(\Omega_2)\} = \operatorname{Ker}(P_{\Gamma_2} - I_{\Gamma_2}).
$$

b) $\operatorname{Ker}(T_{\Gamma_1 \to \Gamma_2}) = \{(u, \partial_\nu u)_{|\Gamma_1}; u \in \tilde{\mathcal{H}}(\Omega_1)\} = \operatorname{Ker}(P_{\Gamma_1}).$

Proof. We start by proving a). Let $(\varphi, \psi) \in H^{1/2}(\Gamma_2) \times H^{-1/2}(\Gamma_2)$ such that $T_{\Gamma_2 \to \Gamma_1}(\varphi, \psi) = 0$. We denote by v the solution of the Helmholtz (or Laplace) equation in $\mathbb{R}^m \setminus \Gamma_2$ given by:

$$
v = -DL_{\Gamma_2}(\varphi) + SL_{\Gamma_2}(\psi), \qquad \text{in } \mathbf{R}^m \setminus \Gamma_2. \tag{15}
$$

Knowing that $T_{\Gamma_2 \to \Gamma_1}(\varphi, \psi) = 0$, it is clear that $(v, \partial_\nu v)$ vanishes on Γ_1 . Using the unique continuation principle for radiating solutions, we get that $v = 0$ in $\mathbb{R}^m \setminus \Omega_2$. Consequently, we obtain

$$
(v^+, \partial_\nu v^+)|_{\Gamma_2} = (0,0)
$$
 on Γ_2 .

Since $[v]_{|\Gamma_2} = -\varphi$ and $[\partial_\nu v]_{|\Gamma_2} = -\psi$, we get $(\varphi, \psi) = (v^-, \partial_\nu v^-)_{|\Gamma_2}$. This proves that $\text{Ker}(T_{\Gamma_2 \to \Gamma_1}) \subset$ $\{(u, \partial_{\nu}u)_{|\Gamma_2}; u \in \mathcal{H}(\Omega_2)\}.$

For the reverse inclusion, consider $(\varphi, \psi) = (u, \partial_{\nu}u)_{|\Gamma_2}$ where $u \in \mathcal{H}(\Omega_2)$. For $x \in \mathbb{R}^m \setminus \overline{\Omega_2}$ we have, using the Green formula in Ω_2 ,

$$
\int_{\Gamma_2} \partial_{\nu} u(y) \Phi(x, y) ds(y) - \int_{\Gamma_2} u(y) \partial_{\nu} \Phi(x, y) ds(y) = 0.
$$

Therefore, the function v be defined by (15) satisfies $v = 0$ in $\mathbb{R}^m \setminus \overline{\Omega_2}$. Hence

$$
(v, \partial_\nu v)_{|\Gamma_1} = 0
$$
 and $T_{\Gamma_2 \to \Gamma_1}(\varphi, \psi) = (0, 0).$

Consequently $\text{Ker}(T_{\Gamma_2 \to \Gamma_1}) = \{(u, \partial_\nu u)_{|\Gamma_2}; u \in \mathcal{H}(\Omega_2)\}\.$

We now prove the second equality in a), i.e. $\text{Ker}(T_{\Gamma_2 \to \Gamma_1}) = \text{Ker}(P_{\Gamma_2} - I_{\Gamma_2})$. In fact we proved in the first part that if $(\varphi, \psi) \in \text{Ker}(T_{\Gamma_2 \to \Gamma_1})$ then $(\varphi, \psi) = (v, \partial_\nu v)_{|\Gamma_2}$ where v satisfies the Helmholtz equation inside Ω_2 and can be represented as

$$
v = -DL_{\Gamma_2}(v_{|\Gamma_2}) + SL_{\Gamma_2}(\partial_\nu v_{|\Gamma_2}) \quad \text{in} \quad \Omega_2.
$$

Using the trace and the normal derivative trace of v on Γ_2 we obtain

$$
\begin{array}{rcl} v_{|\Gamma_2} & = & -K_{\Gamma_2}(v_{|\Gamma_2}) + \frac{v_{|\Gamma_2}}{2} + S_{\Gamma_2}(\partial_\nu v_{|\Gamma_2}), \\[2mm] \partial_\nu v_{|\Gamma_2} & = & -T_{\Gamma_2}(v_{|\Gamma_2}) + K'_{\Gamma_2}(\partial_\nu v_{|\Gamma_2}) + \frac{\partial_\nu v_{|\Gamma_2}}{2}, \end{array}
$$

which proves that $(v, \partial_\nu v)_{|\Gamma_2} \in \text{Ker}(P_{\Gamma_2} - I_{\Gamma_2})$. On the other hand let $(\varphi, \psi) \in \text{Ker}(P_{\Gamma_2} - I_{\Gamma_2})$ and let the function v defined by (15) . Using again trace formulas for double and single layer potentials, we obtain

$$
\left(\begin{array}{c}v_{|\Gamma_2}^+\\ \partial_\nu v_{|\Gamma_2}^+\end{array}\right)=(P_{\Gamma_2}-I_{\Gamma_2})\left(\begin{array}{c}\psi\\ \varphi\end{array}\right)=\left(\begin{array}{c}0\\ 0\end{array}\right).
$$

Therefore the function v replaced by 0 in Ω_2 is locally H^1 and is solution to the Helmholtz equation in \mathbb{R}^m . Unique continuation principle implies that v is zero outside Ω_2 , which proves that $(v, \partial_\nu v)_{|\Gamma_1}$ is also zero. Consequently $(\varphi, \psi) \in \text{Ker}(T_{\Gamma_2 \to \Gamma_1})$.

The proof of b) uses similar arguments as the proof of a) and is based on considering v defined by (15) where Γ_2 is replaced with Γ_1 , namely

$$
v = -DL_{\Gamma_1}(\varphi) + SL_{\Gamma_1}(\psi), \qquad \text{in } \mathbf{R}^m \setminus \Gamma_1 \tag{16}
$$

for $(\varphi, \psi) \in H^{1/2}(\Gamma_1) \times H^{-1/2}(\Gamma_1)$. Assume that $T_{\Gamma_1 \to \Gamma_2}(\varphi, \psi) = 0$. Then $(v, \partial_\nu v)$ vanishes on Γ₂. Using the unique continuation principle for solutions of the Helmholtz equation inside $Ω_1$ we get $v = 0$ in Ω_1 . Consequently, we obtain

$$
(v^-,\partial_\nu v^-)_{|\Gamma_1} = (0,0) \quad \text{on } \Gamma_1.
$$

Since $[v]_{\Gamma_1} = -\varphi$ and $[\partial_\nu v]_{\Gamma_1} = -\psi$, we get $(\varphi, \psi) = -(v^+, \partial_\nu v^+)_{\Gamma_1}$. This proves that $\text{Ker}(T_{\Gamma_1 \to \Gamma_2}) \subset \Big\{ (u, \partial_\nu u)_{|\Gamma_1}; \ u \in \tilde{\mathcal{H}}(\Omega_1) \Big\}.$

For the reverse inclusion, consider $(\varphi, \psi) = (u, \partial_{\nu}u)_{|\Gamma_1}$ where $u \in \tilde{\mathcal{H}}(\Omega_1)$. For $x \in \Omega_1$, using the fact that u satisfies the radiation condition we get

$$
\int_{\Gamma_1} \partial_{\nu} u(y) \Phi(x, y) ds(y) - \int_{\Gamma_1} u(y) \partial_{\nu} \Phi(x, y) ds(y) = 0.
$$

Therefore, the function v defined by (16) satisfies $v = 0$ in Ω_1 . Hence

$$
(v, \partial_\nu v)_{|\Gamma_2} = 0
$$
 and $T_{\Gamma_1 \to \Gamma_2}(\varphi, \psi) = (0, 0).$

Consequently $\text{Ker}(T_{\Gamma_1 \to \Gamma_2}) = \{(u, \partial_\nu u)_{|\Gamma_1}; u \in \tilde{\mathcal{H}}(\Omega_1)\}\.$

We now prove the second equality in b), i.e. $\text{Ker}(T_{\Gamma_1 \to \Gamma_2}) = \text{Ker}P_{\Gamma_1}$. The first part of the proof of assertion b) shows that if $(\varphi, \psi) \in \text{Ker}(T_{\Gamma_2 \to \Gamma_1})$ then $(\varphi, \psi) = -(v, \partial_\nu v)_{|\Gamma_2}$ where v satisfies the Helmholtz equation outside Ω_1 and can be represented as

$$
v = DL_{\Gamma_2}(v_{|\Gamma_1}) - SL_{\Gamma_1}(\partial_\nu v_{|\Gamma_1}) \quad \text{in} \quad \mathbf{R}^m \setminus \overline{\Omega_1}.
$$

Using the trace and the normal derivative trace of v on Γ_1 we obtain

$$
\begin{array}{rcl} v_{|\Gamma_1} & = & K_{\Gamma_1}(v_{|\Gamma_1}) + \dfrac{v_{|\Gamma_1}}{2} - S_{\Gamma_1}(\partial_\nu v_{|\Gamma_1}), \\[2mm] \partial_\nu v_{|\Gamma_1} & = & T_{\Gamma_1}(v_{|\Gamma_1}) - K'_{\Gamma_1}(\partial_\nu v_{|\Gamma_1}) + \dfrac{\partial_\nu v_{|\Gamma_1}}{2}, \end{array}
$$

which proves that $(v, \partial_\nu v)_{|\Gamma_1} \in \text{Ker}(P_{\Gamma_1})$. On the other hand let $(\varphi, \psi) \in \text{Ker}(P_{\Gamma_1})$ and let the function v defined by (16) . Using again trace formulas for double and single layer potentials, we obtain

$$
\left(\begin{array}{c}v_{|\Gamma_1}^-\\ \partial_\nu v_{|\Gamma_1}^-\end{array}\right)=P_{\Gamma_1}\left(\begin{array}{c}\psi\\ \varphi\end{array}\right)=\left(\begin{array}{c}0\\ 0\end{array}\right).
$$

Therefore the function v replaced by 0 outside Ω_1 is a H^1 solution to the Helmholtz equation in \mathbf{R}^m . Unique continuation principle implies that v is zero inside Ω_1 , which proves that $(v, \partial_\nu v)_{|\Gamma_2}$ is also zero. Consequently $(\varphi, \psi) \in \text{Ker}(T_{\Gamma_1 \to \Gamma_2})$. is also zero. Consequently $(\varphi, \psi) \in \text{Ker}(T_{\Gamma_1 \to \Gamma_2})$.

We now characterize the range of the operators $T_{\Gamma_2 \to \Gamma_1}$ and $T_{\Gamma_1 \to \Gamma_2}$. This is consequence of the previous lemma using a duality argument. Let us define a duality pairing \langle, \rangle_i between $H^{1/2}(\Gamma_i) \times$ $H^{-1/2}(\Gamma_i)$ and its dual $H^{-1/2}(\Gamma_i) \times H^{1/2}(\Gamma_i)$ such that

$$
\langle (\varphi,\psi),(\varphi',\psi')\rangle_i:=\int_{\Gamma_i} \varphi\, \varphi'\, ds - \int_{\Gamma_i} \psi\, \psi'\, ds
$$

for regular functions (φ, ψ) and (φ', ψ') . Then for $(\varphi_2, \psi_2) \in H^{1/2}(\Gamma_2) \times H^{-1/2}(\Gamma_2)$ and $(\varphi_1, \psi_1) \in$ $H^{1/2}(\Gamma_1) \times H^{-1/2}(\Gamma_1),$

$$
-\langle T_{\Gamma_2 \to \Gamma_1}(\varphi_2, \psi_2), (\psi_1, \varphi_1) \rangle_1
$$

=
$$
\int_{\Gamma_1} \int_{\Gamma_2} \frac{\partial \Phi(x, y)}{\partial \nu(y)} \varphi_2(y) \psi_1(x) - \int_{\Gamma_1} \int_{\Gamma_2} \Phi(x, y) \psi_2 \psi_1(x)
$$

$$
- \int_{\Gamma_1} \int_{\Gamma_2} \frac{\partial^2 \Phi(x, y)}{\partial \nu(x) \partial \nu(y)} \varphi_2(y) \varphi_1(x) + \int_{\Gamma_1} \int_{\Gamma_2} \frac{\partial \Phi(x, y)}{\partial \nu(x)} \psi_2(y) \varphi_1(x),
$$

where the integrals should be understood as duality products. Observing the symmetry between the indexes 1 and 2 in this identity we conclude

$$
\langle T_{\Gamma_2 \to \Gamma_1}(\varphi_2, \psi_2), (\psi_1, \varphi_1) \rangle_1 = \langle T_{\Gamma_1 \to \Gamma_2}(\varphi_1, \psi_1), (\psi_2, \varphi_2) \rangle_2. \tag{17}
$$

Similarly, for $(\varphi, \psi) \in H^{1/2}(\Gamma_i) \times H^{-1/2}(\Gamma_i)$ and $(\tilde{\varphi}, \tilde{\psi}) \in H^{1/2}(\Gamma_i) \times H^{-1/2}(\Gamma_i)$

$$
\label{eq:3.1} \begin{array}{ll} \displaystyle \Big\langle P_{\Gamma_i}(\varphi,\psi), (\tilde{\psi},\tilde{\varphi}) \Big\rangle_i = \int_{\Gamma_i} -K_{\Gamma_i} \varphi \, \tilde{\psi} + \frac{1}{2} \varphi \tilde{\psi} + S_{\Gamma_i} \psi \, \tilde{\psi} dx \\[2mm] \displaystyle & \qquad \qquad - \int_{\Gamma_i} -T_{\Gamma_i} \varphi \, \tilde{\varphi} + K'_{\Gamma_i} \psi \, \tilde{\varphi} + \frac{1}{2} \psi \, \tilde{\varphi} dx. \end{array}
$$

We deduce that

$$
\left\langle P_{\Gamma_i}(\varphi,\psi),(\tilde{\psi},\tilde{\varphi})\right\rangle_i = \left\langle (P_{\Gamma_i} - I_{\Gamma_i})(\tilde{\varphi},\tilde{\psi}),(\psi,\varphi)\right\rangle_i \quad i = 1,2. \tag{18}
$$

These duality identities allows us to deduce the following corollary.

Corollary 4. The following identities hold.

$$
\frac{\text{Range}(T_{\Gamma_2 \to \Gamma_1})}{\text{Range}(T_{\Gamma_1 \to \Gamma_2})} = \text{Range}(P_{\Gamma_1} - I_{\Gamma_1}) = \text{Ker}(P_{\Gamma_1}),
$$

$$
\frac{\text{Range}(T_{\Gamma_1 \to \Gamma_2})}{\text{Range}(T_{\Gamma_2})} = \text{Range}(P_{\Gamma_2}) = \text{Ker}(P_{\Gamma_2} - I_{\Gamma_2}).
$$

Proof. Identity (17) shows that

$$
\overline{\text{Range}(T_{\Gamma_2 \to \Gamma_1})} = (\text{Ker}(T_{\Gamma_1 \to \Gamma_2}))^{\perp} \text{ and } \overline{\text{Range}(T_{\Gamma_1 \to \Gamma_2})} = (\text{Ker}(T_{\Gamma_2 \to \Gamma_1}))^{\perp}
$$

where the orthogonality is with respect to the duality products $\langle ., . \rangle_i$ defined. We also have using (18)

$$
Ker(P_{\Gamma_i}^{\perp}) = \text{Range}(P_{\Gamma_i} - I_{\Gamma_i}) = \text{Ker}(P_{\Gamma_i}).
$$

 \Box The announced identities are then be deduced as direct consequences of those in Lemma 3. We are now in position to state the main result. We define

$$
X(\Gamma_i) := H^{1/2}(\Gamma_i) \times H^{-1/2}(\Gamma_i)
$$

and the subspace of $X(\Gamma_1) \times X(\Gamma_2)$

$$
Y(\Gamma_1 \times \Gamma_2) := \text{Ker}(P_{\Gamma_1}) \times \text{Ker}(P_{\Gamma_2} - I_{\Gamma_2}).
$$

Theorem 5. The operator $A: X(\Gamma_2) \to Y(\Gamma_1 \times \Gamma_2)$ defined by expression (11) is one to one and has a dense range. Moreover $\text{Range}(A) = \text{Range}(B)$ where B is defined by (12).

Proof. We start by proving the injectivity of the operator A. Since P_{Γ_2} is a projector then

for all
$$
u \in X(\Gamma_2)
$$
, $u = u_1 + u_2$

with $u_1 \in \text{Ker}(P_{\Gamma_2})$ and $u_2 \in \text{Range}(P_{\Gamma_2})$. Assume that $A(u) = 0$. Then

$$
T_{\Gamma_2 \to \Gamma_1}(u) = 0
$$
 and $P_{\Gamma_2}(u) = 0$.

Since $\text{Range}(P_{\Gamma_2}) = \text{Ker}(P_{\Gamma_2} - I_{\Gamma_2})$, we get from Lemma 3 that $u_2 \in \text{Ker}(T_{\Gamma_2 \to \Gamma_1})$. Consequently $T_{\Gamma_2 \to \Gamma_1}(u_1) = 0$ and therefore $u_1 \in \text{Range}(P_{\Gamma_2})$. This implies $u_1 = 0$.

On the other hand, $A(u) = 0$ also implies $P_{\Gamma_2}(u) = 0$ and therefore $P_{\Gamma_2}(u_2) = 0$. This proves that $u_2 = 0$ since we already have $u_2 \in \text{Range}(P_{\Gamma_2})$. The injectivity of A is then proved.

For the characterization of the range of the operator A we use Corollary 4 in which we have proved that $\text{Range}(T_{\Gamma_2 \to \Gamma_1}) = \text{Ker}(P_{\Gamma_1})$. Consequently

$$
\overline{\text{Range}(A)} = \text{Ker}(P_{\Gamma_1}) \times \text{Range}(P_{\Gamma_2}) = \text{Ker}(P_{\Gamma_1}) \times \text{Ker}(P_{\Gamma_2} - I_{\Gamma_2})
$$

Using Corollary 4 we also deduce that

$$
\overline{\text{Range}(B)} = \text{Ker}(P_{\Gamma_1}) \times \overline{\text{Range}(T_{\Gamma_1 \to \Gamma_2})} = \text{Ker}(P_{\Gamma_1}) \times \text{Ker}(P_{\Gamma_2} - I_{\Gamma_2}) = \overline{\text{Range}(A)}.
$$

A straightforward consequence of this result is that any regular regularization strategy can be applied to equation 13 in the case of noisy data. We hereafter state the result for the case of Tikhonov regularization.

Corollary 6. Let $(f, g) \in X(\Gamma_1)$ and $(u_{|\Gamma_2}, \partial_\nu u_{|\Gamma_2}) \in X(\Gamma_2)$ such that

$$
A\left(\begin{array}{c}u_{\mid \Gamma_2} \\ \partial_\nu u_{\mid \Gamma_2}\end{array}\right) = B\left(\begin{array}{c}f \\ g\end{array}\right) \tag{19}
$$

and let $(f^{\delta}, g^{\delta}) \in X(\Gamma_1)$ such that

$$
\left\| \left(\begin{array}{c} f^{\delta} \\ g^{\delta} \end{array} \right) - \left(\begin{array}{c} f \\ g \end{array} \right) \right\|_{X(\Gamma_1)} < \delta.
$$

Consider the Tikhonov solution $(u^{\delta}, \partial_{\nu}u^{\delta})$ satisfying

$$
(\alpha(\delta)I_{\Gamma_2} + A^*A) \begin{pmatrix} u^{\delta} \\ \partial_{\nu}u^{\delta} \end{pmatrix} = A^*B \begin{pmatrix} f^{\delta} \\ g^{\delta} \end{pmatrix}
$$
 (20)

where $\alpha(\delta)$ is determined (for sufficiently small δ) using the Morozov discrepancy principle, i.e. $\alpha(\delta)$ is the unique solution of

$$
\left\| A \left(\begin{array}{c} u^{\delta} \\ \partial_{\nu} u^{\delta} \end{array} \right) - B \left(\begin{array}{c} f^{\delta} \\ g^{\delta} \end{array} \right) \right\|_{Y(\Gamma_1 \times \Gamma_2)} = \delta \| B \|_{X(\Gamma_1) \to Y(\Gamma_1 \times \Gamma_2)}.
$$
(21)

Then

$$
\left\| \left(\begin{array}{c} u^{\delta} \\ \partial_{\nu} u^{\delta} \end{array} \right) - \left(\begin{array}{c} u \\ \partial_{\nu} u \end{array} \right) \right\|_{X(\Gamma_2)} \longrightarrow 0 \qquad \text{as} \qquad \delta \longrightarrow 0. \tag{22}
$$

Proof. The proof of this Corollary can be deduced directly from Theorem 5 using the Tikhonov regularization theory and the Morozov discrepancy principle (See for instance [13, Theorem 4.16]) \Box

In the previous corollary, one assumes that $f^{\delta} \in H^{1/2}(\Gamma_1)$. A more reasonable assumption would be that the noisy data is only in $L^2(\Gamma_1)$. In that case a regularization of the data is needed. We hereafter explain a possible method based again on Tikhonov regularization (any other standard method can be indeed used). Let $\{e_n \in L^2(\Gamma_1), n = 1, 2 \ldots\}$ denote an orthonormal basis of $L^2(\Gamma_1)$ formed by the eigenvectors of the Laplace-Beltrami operator on Γ_1 and denote by σ_n the eigenvalue associated with e_n . For $f \in H^s(\Gamma_1)$ we recall that an equivalent definition of the $H^s(\Gamma_1)$ norm is

$$
\|f\|^2_{H^s(\Gamma_1)}:=\sum_{j=1}^\infty (1+\sigma_n^{1/2})^{2s}|(f,e_n)_{L^2(\Gamma_1)}|^2.
$$

Set $\lambda_n := 1/\sqrt{1 + \sigma_n^{1/2}}$ and set for $h \in L^2(\Gamma_1)$

$$
R_{\beta}h := \sum_{n=1}^{\infty} \frac{\lambda_n}{\sqrt{\beta + \lambda_n^2}} (h, e_n)_{L^2(\Gamma_1)} e_n.
$$

The operator R_β is nothing but the Tikhonov regularized inverse of the compact injection $H^{1/2}(\Gamma_1) \to$ $L^2(\Gamma_1)$. In particular, one easily verifies that for $f \in H^{1/2}(\Gamma_1)$

$$
||R_{\beta}f - f||_{H^{1/2}(\Gamma_1)} \to 0 \text{ as } \beta \to 0
$$

and if $f \in H^1(\Gamma_1)$ then

$$
||R_{\beta}f - f||_{H^{1/2}(\Gamma_1)} \le \sqrt{\beta} ||f||_{H^1(\Gamma_1)}.
$$

Morover, for $h \in L^2(\Gamma_1)$

$$
||R_{\beta}h||_{H^{1/2}(\Gamma_1)} \leq \frac{1}{\sqrt{\beta}}||h||_{L^2(\Gamma_1)}.
$$

It is then obvious that if $\beta(\delta)$ is such that $\beta(\delta) \to 0$ and $\delta^2/\beta(\delta) \to 0$ as $\delta \to 0$, then for $f^{\delta} \in L^2(\Gamma_1)$ such that

$$
||f^{\delta} - f||_{L^{2}(\Gamma_1)} \leq \delta
$$

one has $R_{\beta}f^{\delta} \in H^{1/2}(\Gamma_1)$ and

$$
||R_{\beta}f^{\delta}-f||_{H^{1/2}(\Gamma_1)} \leq \delta/\sqrt{\beta(\delta)} + ||R_{\beta(\delta)}f-f||_{H^{1/2}(\Gamma_1)} \to 0 \text{ as } \delta \to 0.
$$

If $f \in H^1(\Gamma_1)$ and $\beta(\delta) = \delta$ then one has the more explicit estimate

$$
||R_{\delta}f^{\delta} - f||_{H^{1/2}(\Gamma_1)} \leq \sqrt{\delta}(1 + ||f||_{H^1(\Gamma_1)}).
$$

One then can state the straightforward following convergence result for L^2 noisy data.

Corollary 7. Let $(f,g) \in X(\Gamma_1)$ and $(u_{|\Gamma_2}, \partial_\nu u_{|\Gamma_2}) \in X(\Gamma_2)$ such that (19) is satisfied and further assume that $f \in H^1(\Gamma_1)$. Let $(f^{\delta}, g^{\delta}) \in X(\Gamma_1)$ such that

$$
\left\| \left(\begin{array}{c} f^{\delta} \\ g^{\delta} \end{array} \right) - \left(\begin{array}{c} f \\ g \end{array} \right) \right\|_{L^{2}(\Gamma_{1}) \times L^{2}(\Gamma_{1})} < \delta
$$

and consider the Tikhonov solution $(u^{\delta}, \partial_{\nu}u^{\delta})$ of (20) where f^{δ} is replaced with $R_{\delta}f^{\delta}$ and $\alpha(\delta)$ is determined (for sufficiently small δ) using the Morozov discrepancy principle, i.e $\alpha(\delta)$ is the is aetermined (for sufficiently small *o*) using the Morozov aiscrepancy principle, *i.e* α (*o*) is the unique solution of (21) with δ in the right hand side replaced by $\sqrt{\delta}(1 + ||f||_{H^1(\Gamma_1)})$ and again f^{δ} is replaced with $R_{\delta}f^{\delta}$. Then

$$
\left\| \left(\begin{array}{c} u^{\delta} \\ \partial_{\nu} u^{\delta} \end{array} \right) - \left(\begin{array}{c} u \\ \partial_{\nu} u \end{array} \right) \right\|_{X(\Gamma_2)} \longrightarrow 0 \qquad \text{as} \qquad \delta \longrightarrow 0. \tag{23}
$$

4 Numerical validation

The numerical experiments are conducted in a 2D setting of the problem. We recall that our algorithm is based on solving equation (20). We assume that Γ_i is parametrized as

$$
\Gamma_i = \{ (x_i(\theta), y_i(\theta)); \theta \in [0, 2\pi] \}
$$

with regular functions x_i and y_i , $i = 1, 2$. Then matrix approximations $A_M \in \mathbb{C}^{4M \times 2M}$ and $B_M \in \mathbb{C}^{4M \times 2M}$ of the operators A and B are constructed using M equally distant points θ_ℓ , $\ell = 1, \ldots M$ with a step $= 2\pi/M$. We employ a collocation method and a Nyström method (with weighted trigonometric interpolation) (see [13]) to approximate the singular parts of the potentials while a trapezoidal rule is used to approximate the nonsingular component of these operators. For a given (noisy) discrete data $F_1 \in \mathbb{C}^{2M}$ on the exterior boundary Γ_1 (that reflect pointwise values of (f, g) at $(x_1(\theta_\ell), y_1(\theta_\ell))$, we construct an approximation of the unknown Cauchy data $F_2 \simeq (u, \partial_\nu u)|_{\Gamma_2}$ using the following expression

$$
F_2(\ell) = \sum_{k=1}^{2M} \frac{\sigma_k}{\alpha + \sigma_k^2} (B_M F_1, t_k) p_k(\ell), \qquad 1 \le \ell \le 2M
$$
 (24)

where $\{\sigma_k, p_k, t_k\}_{1\leq k\leq 2M}$ are the 2M first components of the singular value decomposition of A_M that satisfy $A_M p_k = \sigma_k t_k$, with $\sigma_k > 0$, $p_k \in \mathbb{C}^{2M}$ and $t_k \in \mathbb{C}^{4M}$ and α is the regularization parameter chosen by the Morozov discrepancy principle by solving (21). In the numerical examples below we simply used (obvious) piecewise linear interpolation of data corrupted with random noise (without regularization) which already give a $H^1(\Gamma_1)$ function.

In order to construct a synthetic data that does not coincide with Helmholtz (or Laplace) solutions in all of Ω_1 we consider an inclusion Ω_0 (with boundary Γ_0) strictly included in Ω_2 and construct the data from solutions to the Helmholtz (or Laplace) equation in $\Omega_1 \setminus \Omega_0$ with Neumann boundary

conditions on Γ_0 and a prescribed Dirichlet data on Γ_1 . A numerical solution is obtained using a boundary integral equation formulation of the problem approximated with a Nyström method (and weighted trigonometric interpolation quadratures on an equidistant mesh, as explained in [13, Chapter 3]).

In the first test, we consider as prescribed Dirichlet data on Γ_1 , $f(x) = \Phi(x, x_0)$ where x_0 is a point source fixed in the exterior of Ω_1 . The parametrization of Γ_1 , Γ_2 and Γ_0 are given by 2.5

$$
\Gamma_1 = \{ (1.2 \cos(\theta), 1.5 \sin(\theta) + 0.05 \sin(3\theta)) \}
$$

\n
$$
\Gamma_2 = \{ (0.8 \cos(\theta), \sin(\theta) + 0.1 \sin(3\theta)) \}
$$

\n
$$
\Gamma_0 = \{ (0.5 \cos(\theta) + 0.1 \cos(2\theta), 0.5 \sin(\theta)) \}
$$

The point source is at $(5, 5)$.

Figure 2: Sketch of the geometry for the first example.

The series of Figures 3 - 4 - 5 show the obtained reconstructions for different values of the wave number $k = 0, 2, 4$. We observe that a good accuracy is achieved for noise-free and noisy data. The accuracy-dependence on the frequency (corresponding with a not too small wavelength) seems to be robust.

Figure 3: Reconstruction of u (left) and $\partial_{\nu}u$ (right) for the test problem described in Figure 2 and for a wave number $k = 0$.

Figure 4: Reconstruction of the modulus of u (left) and the modulus of $\partial_\nu u$ (right) for the test problem described in Figure 2 and a wave number $k = 2$.

Figure 5: Reconstruction of the modulus ofu (left) and $\partial_{\nu}u$ (right) for the test problem described in Figure 2 and a wave number $k = 4$.

Figure 6: Sketch of the geometry for the second example.

Figure 7: Reconstruction of u (left) and $\partial_{\nu}u$ (right) for the test problem described in Figure 6 and a wave number $k = 0$

Figure 8: Reconstruction of the norm of the trace of u (at left) and the norm of the trace of $\partial_{\nu}u$ (at right) for the test problem described in Figure 6 on Γ_2 for $k = 2$

In the third test (see Figures 10-11) we consider the same setting of the second test except that the exterior boundary $Γ₂$ is slightly moved far from $Γ₁$. The parametrization of $\Gamma_1,$ Γ_2 and Γ_0 are given by $\Gamma_1 = \{ (2 \cos(\theta) + 0.2 \cos(2\theta), 2.5 \sin(\theta) + 0.2 \sin(2\theta)) \}$ $\Gamma_2 = \{ (0.8 \cos(\theta), 0.4 \sin(\theta) + 0.5 \sin(2\theta) \}$ $\Gamma_0 = \{ (0.4 \cos(\theta) + 0.1 \cos(2\theta), 0.6 \sin(\theta)) \}$ −2.5 −2 −1.5 −1 −0.5 0 0.5 1 1.5 2 2.5 −2.5 −2 −1.5 −1 −0.5 ٥ŀ 0.5 1 F 1.5 2 F 2.5 Γ_{1} Γ_{α} 2 Γ \mathbf{a}

Figure 9: Sketch of the geometry for the third example.

Figure 10: Reconstruction of u (left) and $\partial_{\nu}u$ (right) for the test problem described in Figure 9 and a wave number $k = 0$

Figure 11: Reconstruction of the modulus of u (left) and $\partial_{\nu}u$ (right) for the test problem described in Figure 9 and a wave number $k = 2$

The series of Figures 7-8 and Figures 10-11 show that we have a good reconstruction achieved for different levels of noise and for different shapes of the domains. These figures show also the stability of the proposed algorithm since even by moving the exterior boundary Γ_1 far from Γ_2 the reconstruction is not very sensitive to the distance between the boundary of the given data and the boundary where data is reconstructed. In the follwing Figures 12-14 we consider the same setting as in the third test but we have changed the Dirichlet boundary conditions by imposing $f(x) = e^{ikdx}$ with $d = (0, 1)$. We obtain even better reconstructions for this type data.

Figure 12: Reconstruction of u (left) and $\partial_{\nu}u$ (right) for the test problem described in Figure 9, a wave number $k = 0$ and $f(x) = e^{ikdx}$ with $d = (0, 1)$.

Figure 13: Reconstruction of the modulus of u (left) and $\partial_{\nu}u$ (right) for the test problem described in Figure 9, a wave number $k = 2$ and $f(x) = e^{ikdx}$ with $d = (0, 1)$.

Figure 14: Reconstruction of the modulus of u (left) and $\partial_{\nu}u$ (right) for the test problem described in Figure 9, a wave number $k = 4$ and $f(x) = e^{ikdx}$ with $d = (0, 1)$.

We now give an example of a smooth geometry where the exact solution of the Helmholtz or the Laplace equation can be analytically computed. In fact, in this test we compare the reconstructed solution with an exact solution that is known. Moroever, in the case of the Helmholz equation we compare the real and the imaginary parts of the trace of u and $\partial_{\nu}u$ to show that our algorithm is also efficient to reconstruct both the modulus and the phase. Note that the conclusion also holds for previous examples where we have opted to show only the modulus for the purpose of simplifying the presentation and to provide more examples.

In the fourth test (see Figures 16-18) we will consider a smooth geometry with two boundary Γ_1 and Γ_2 defined _{0.5} by:

$$
\Gamma_1 = \{(\cos(\theta), 1.5\sin(\theta))\}\n\Gamma_2 = \{(0.5\cos(\theta), 0.75\sin(\theta))\}\n\right\}
$$

1 F 1.5

−1.5 −1 −0.5 0 0.5 1 1.5

 Γ_{1}

Γ 2

Figure 15: Sketch of the geometry for the fourth example.

Figure 16: Reconstruction of u (left) and $\partial_{\nu}u$ (right) for the test problem described in Figure 15, a wave number $k = 0$ and $u(x, y) = xy$.

Figure 17: Reconstruction of the real part (left) and the imaginary part (right) of u for the test Figure 17: Reconstruction of the real part (left) and the imaginary part (right) of u for
problem described in Figure 15, a wave number $k = 4$ and $f(x) = e^{ikdx}$ with $d = (\sqrt{2}/2,$ √ $2/2$).

Figure 18: Reconstruction of the real part (left) and the imaginary part (right) of $\partial_{\nu}u$ for the test Figure 18: Reconstruction of the real part (left) and the imaginary part (right) of $\partial_{\nu} u$ for problem described in Figure 15, a wave number $k = 4$ and $f(x) = e^{ikdx}$ with $d = (\sqrt{2}/2,$ √ $2/2$).

We end this paper with a test where we keep the same setting of the geometry as in the last test and we suppose that we have a solution of the Helmholtz equation satisfying the Neumann boundary condition on Γ_2 and a Dirichlet boundary condition on Γ_1 with $f = \Phi(., x_0)$ where $x_0 = 6 + 5i.$

Figure 19: Reconstruction of u (left) and $\partial_{\nu}u$ (right) for the test problem described in Figure 15 and a wave number $k = 0$

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Figure 20: Reconstruction of the modulus of u (left) and $\partial_{\nu}u$ (right) for the test problem described in Figure 15 and a wave number $k = 2$

Figure 21: Reconstruction of the modulus of u (left) and $\partial_{\nu}u$ (right) for the test problem described in Figure 15 and a wave number $k = 4$

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