Discreteness of interior transmission eigenvalues revisited

Hoai-Minh Nguyen* Quoc-Hung Nguyen†

October 25, 2016

Abstract

This paper is devoted to the discreteness of the transmission eigenvalue problems. It is known that this problem is not self-adjoint and a priori estimates are non-standard and do not hold in general. Two approaches are used. The first one is based on the multiplier technique and the second one is based on the Fourier analysis. The key point of the analysis is to establish the compactness and the uniqueness for Cauchy problems under various conditions. Using these approaches, we are able to rediscover quite a few known discreteness results in the literature and obtain various new results for which only the information near the boundary are required and there might be no contrast of the coefficients on the boundary.

Contents

1 Introduction 1

2 On the case $A_1 \geq A_2$ and $\Sigma_1 \geq \Sigma_2$ in a neighborhood of $\Gamma$ - Proof of Theorem 1 7

3 On the case $A_1 \geq A_2$ and $\Sigma_1$ not greater than $\Sigma_2$ in a neighborhood of $\Gamma$ 18

3.1 On the case $A_1 \geq A_2$ and $\Sigma_1 \leq \Sigma_2$ in a neighborhood of $\Gamma$ - Proof of Theorem 2 18

3.2 On the case $A_1 \geq A_2$ globally in $\Omega$ 25

4 On the case $A_1 = A_2$ in a neighborhood of $\Gamma$ - Proof of Theorem 3 26

5 On the case $A_1$ and $A_2$ satisfying the complementary condition on $\Gamma$ - Proof of Theorem 4 36

1 Introduction

Let $d \geq 2$ and $\Omega$ be a smooth bounded simply connected open subset of $\mathbb{R}^d$ and denote $\Gamma = \partial \Omega$. Let $A_1, A_2$ be two real symmetric matrix-valued functions and $\Sigma_1, \Sigma_2$ be two bounded

*hoai-minh.nguyen@epfl.ch, École Polytechnique Fédérale de Lausanne, SB MATHAA CAMA, Station 8, CH-1015 Lausanne, Switzerland.
†quoc-hung.nguyen@epfl.ch, École Polytechnique Fédérale de Lausanne, SB MATHAA CAMA, Station 8, CH-1015 Lausanne, Switzerland.
positive functions all defined in $\Omega$ such that $A_1$ and $A_2$ are uniformly elliptic and $\Sigma_1, \Sigma_2$ are bounded below by a positive constant in $\Omega$, i.e., for some constant $\Lambda \geq 1$, and for $j = 1, 2$,

$$\Lambda^{-1}|\xi|^2 \leq \langle A_j(x)\xi, \xi \rangle \leq \Lambda|\xi|^2 \quad \text{for all } \xi \in \mathbb{R}^d, \text{ for a.e. } x \in \Omega, \quad (1.1)$$

and

$$\Lambda^{-1} \leq \Sigma_j(x) \leq \Lambda \text{ for a.e. } x \in \Omega. \quad (1.2)$$

Here and in what follows $\langle \cdot, \cdot \rangle$ denotes the Euclidean scalar product. A complex number $\lambda$ is called an eigenvalue of the interior transmission eigenvalue (ITE) problem associated with the pairs $(A_1, \Sigma_1)$ and $(A_2, \Sigma_2)$ in $\Omega$ if there is a non-zero pair of functions $(u_1, u_2) \in [H^1(\Omega)]^2$ satisfying the system

$$\begin{cases}
\text{div}(A_1 \nabla u_1) - \lambda \Sigma_1 u_1 = 0 & \text{in } \Omega, \\
\text{div}(A_2 \nabla u_2) - \lambda \Sigma_2 u_2 = 0 & \text{in } \Omega,
\end{cases} \quad (1.3)$$

and the boundary conditions

$$u_1 = u_2, \quad A_1 \nabla u_1 \cdot \nu = A_2 \nabla u_2 \cdot \nu \text{ on } \Gamma, \quad (1.4)$$

where $\nu$ denotes the outward normal unit vector on $\Gamma$. Such a pair $(u_1, u_2)$ is called an eigenfunction pair of (1.3) and (1.4).

The ITE problem was introduced in the middle of eighties by Kirsch in [18] and Colton and Monk in [10]. One of its interests comes from the connection between the density of the range of the far field operator with the injectivity of ITE problem. The ITE problem is known to be not self-adjoint. More severely, a priori estimates are non-standard and do not hold in general. These create many difficulties for the investigation and a source of interesting problems. Some important directions in the study of the ITE problems are the discreteness of the spectrum [6, 11, 13, 15, 18, 20, 21, 22, 34, 35], the completeness of the generalized eigenfunctions [3, 22], and the Weyl-laws for the spectrum [15, 17, 21, 22]. The reader can find a review on the ITE problem in [9].

This paper is devoted to the discreteness of the eigenvalues for the ITE problem. The discreteness of the spectrum of (1.3)-(1.4) was established by Rynne and Sleeman in [34] and Colton, Kirsch, Paivarinta in [11] in the case $A_1 = A_2 = I$, the identity matrix, $\Sigma_1 - \Sigma_2 = \Sigma_1 - 1 \geq c > 0$ in $\Omega$ for some constant $c$. Using the fact that $\Sigma_1 - \Sigma_2 > c > 0$, they can transform (1.3)-(1.4) into an elliptic equation of fourth order and then derive the discreteness of the spectrum of the ITE problem from the new one. The discreteness of the spectrum of the ITE problem involving only the information near the boundary was established quite recently. Using the T-coercivity method, which is related to the (Banach-Necas-Babuska) inf-sup condition, Bonnet-BenDhia, Chesnel, and Haddar in [1] obtained the discreteness of the spectrum of the ITE problem under the two assumptions i) $(A_1, \Sigma_1) = (I, 1)$ or $(A_2, \Sigma_2) = (I, 1)$ and ii) $A_1 - A_2 \geq cI$ and $\Sigma_1 - \Sigma_2 \geq c$ in a neighborhood of $\Gamma$ for some positive constant $c$. In [35], Sylvester showed that the discreteness takes place if $A_1 = A_2 = I$ in $\Omega$ and $\Sigma_1 - \Sigma_2 = \Sigma_1 - 1 \geq c > 0$ in a neighbourhood of $\Gamma$ for some constant $c$ using the theory of upper triangular compact operators and an apriori estimate for the ITE problem. His method also works in the case $A_1 = A_2$ smooth in $\Omega$. In [20], Lakshtanov and Vainberg obtained the discreteness of ITE problem under
the assumption that $(A_1, \Sigma_1) = (I, I)$, $(A_1, A_2, \Sigma_1, \Sigma_2)$ is smooth in $\Omega$ and satisfies the so-called parameter-elliptic conditions. In a related direction, La kshtanov and Vainberg in [21] obtained the discreteness of the ITE problem for the case $(A_1, A_2, \Sigma_1, \Sigma_2) = (I, I, 1, \Sigma_2)$ in $\Omega$, $\Sigma_2$ is smooth in $\bar{\Omega}$, and $\Sigma_2 = 1$ and $\partial_\nu \Sigma_2 \neq 0$ on $\Gamma$.

In this paper, we follow the standard strategy used in the context of Laplace operator to investigate the discreteness of the eigenvalues for the ITE problem. To this end, we first establish the well-posedness of the following system, for some $\gamma_0 \in \mathbb{C}$,

$$
\begin{aligned}
\text{div}(A_1 \nabla u_1) - \gamma_0 \Sigma_1 u_1 &= g_1 & \text{in } \Omega, \\
\text{div}(A_2 \nabla u_2) - \gamma_0 \Sigma_2 u_2 &= g_2 & \text{in } \Omega, \\
u_1 = u_2, & & A_1 \nabla u_1 \cdot \nu = A_2 \nabla u_2 \cdot \nu & \text{on } \Gamma
\end{aligned}
$$

(1.5)

where $(g_1, g_2)$ is a given pair of functions in an appropriate space. We then define the operator $T(f_1, f_2) = (u_1, u_2)$ where $(u_1, u_2)$ is the unique solution of (1.5) with $(g_1, g_2) = (\Sigma_1 f_1, \Sigma_2 f_2)$.

The discreteness of the eigenvalues for the ITE problem can be now derived from the discreteness of the spectrum of $T$. In this paper, the well-posedness and the compactness of the operator $T$, for some appropriate choice $\lambda_0 \in \mathbb{C}$ are obtained via the multiplier technique and the Fourier analysis. The first approach has roots from the work of H-M. Nguyen in [29] where he studies the well-posedness of the Helmholtz equations with sign changing coefficients. Various ideas in this paper were introduced there. The Fourier analysis is quite natural in this setting and the idea is to compute solutions of Cauchy problems in half space and derive the information from this.

Throughout this paper

$$
d_\Gamma(x) := \text{dist}(x, \Gamma) := \inf_{y \in \Gamma} |x - y| \text{ for } x \in \mathbb{R}^d,
$$

(1.7)

and for two matrices $M_1, M_2$, we denote $M_1 \geq M_2$ if $M_1 - M_2$ is a nonnegative matrix.

Our first result of the discreteness of the eigenvalues for the ITE problem is

**Theorem 1.** Assume that for some $0 \leq \alpha < 2$ and for some positive constant $c$,

$$
A_1 - A_2 \geq cd_\Gamma^\alpha I \quad \text{and} \quad \Sigma_1 - \Sigma_2 \geq 0 \quad \text{in a neighborhood of } \Gamma.
$$

(1.8)

The spectrum of (1.3)-(1.4) is discrete.

No regularity assumption on $(A_1, A_2, \Sigma_1, \Sigma_2)$ is imposed in Theorem 1. The conditions in (1.8) are only required near $\Gamma$ and neither upper bound of $A_1 - A_2$ nor the one of $\Sigma_1 - \Sigma_2$ is required. The proof of Theorem 1 is based on the multiplier technique and given in Section 2. To our knowledge, Theorem 1 is new even in the case $\alpha = 0$. Applying Theorem 1 for $\alpha = 0$, one recovers and extends the discreteness results obtained by Bonnet-BenDhia, Chesnel, and Haddar in [4] mentioned.

When the second inequality in (1.8) is reversed and $\Sigma_2$ is large enough, the discreteness also holds. More precisely, we have
**Theorem 2.** Assume that for some $0 \leq \alpha < 2$ and for some positive constant $c$,

$$A_1 - A_2 \geq cd_1^{\alpha} I \quad \text{in a neighborhood of } \Gamma.$$  \hfill (1.9)

Let $\Lambda_1 \geq 1$ be such that

$$\Lambda_1^{-1} I \leq A_j \leq \Lambda_1 I \quad \text{for } j = 1, 2, \quad \text{and} \quad \Lambda_1^{-1} \Sigma_1 \leq \Sigma_1 \leq \Lambda_1 \quad \text{in } \Omega. \quad (1.10)$$

For every $0 < \Lambda_2 < 1$ there exists $K(\Lambda_2) > 0$ depending only on $\Lambda_1$, $\Lambda_2$, $c$, $\alpha$, and the neighborhood of $\Gamma$ in (1.9) such that if

$$\Lambda_2 K \leq \Sigma_2 \leq K \quad \text{in } \Omega, \quad (1.11)$$

for some $K > K(\Lambda_2)$, then the spectrum of (1.3) - (1.4) is discrete.

As far as we know, Theorem 2 is new even in the case $\alpha = 0$. The proof of Theorem 2 is based on the multiplier technique and presented in Section 3. Another related result of Theorem 2 assuming only a mild condition on $(\Sigma_1, \Sigma_2)$ is given in Proposition 8; however, one requires that $A_1 \geq A_2 \geq cd_1^{\alpha} I$ in the whole domain $\Omega$.

We next deal with the case $A_1 = A_2$ in a neighborhood of $\Gamma$.

**Theorem 3.** Assume that for some $0 \leq \beta < 2$ and for some positive constants $c$,

$$A_1 = A_2 \quad \text{and} \quad \Sigma_1 - \Sigma_2 \geq cd_1^{\beta} \quad \text{in a neighborhood of } \Gamma. \quad (1.12)$$

The spectrum of (1.3) - (1.4) is discrete.

Applying Theorem 3 for the case $A_1 = A_2$ smooth in $\Omega$, and $\beta = 0$, one recovers and extends the results mentioned of Sylvester’s in [35] and Lakshtanov and Vainberg’s in [21]. In comparison with their results, the novelty of Theorem 3 relies on the two facts: i) only the information near the boundary on $(A_1, A_2, \Sigma_1, \Sigma_2)$ is assumed and ii) no regularity on the coefficients is required.

In the case $(A_1, A_2, \Sigma_1, \Sigma_2)$ is continuous in a neighborhood of $\Gamma$, we prove the following result which involves the complementing condition due to Agmon, Douglas, and Nirenberg in [1]:

**Theorem 4.** Assume that $A_1, A_2, \Sigma_1, \Sigma_2$ are continuous in a neighborhood of $\Gamma$ and the following two conditions hold:

i) For all $x \in \Gamma$, $A_1(x), A_2(x)$ satisfy the complementing boundary condition with respect to $\nu(x)$,

ii) For all $x \in \Gamma$, $\langle A_1(x)\nu(x), \nu(x) \rangle \Sigma_1(x) \neq \langle A_2(x)\nu(x), \nu(x) \rangle \Sigma_2(x)$.

The spectrum of (1.3) - (1.4) is discrete.

The complementing condition corresponding to the ITE problem has the following explicit characterization given in [29, Proposition 1] (see e.g. [29, Definition 2] for the definition of the complementing condition for the ITE problem).
Let \( e \) be a unit vector in \( \mathbb{R}^d \) and let \( A_1 \) and \( A_2 \) be two constant positive symmetric matrices. Then \( A_1 \) and \( A_2 \) satisfy the complementing boundary condition with respect to \( e \) if and only if

\[
\langle A_2 e, e \rangle \langle A_2 \xi, \xi \rangle - \langle A_1 e, e \rangle \langle A_1 \xi, \xi \rangle - \langle A_1 e, \xi \rangle^2 \neq 0 \quad \forall \xi \in \mathcal{P} \setminus \{0\},
\]

where \( \mathcal{P} := \{ \xi \in \mathbb{R}^d; \langle \xi, e \rangle = 0 \} \).

We next give some comments on the conditions (i) and (ii) of Theorem 4:

- It is shown in [29, Proposition 1] that if \( A_1 > A_2 \) then \( A_1 \) and \( A_2 \) satisfy the complementing boundary condition for all unit vector \( e \).
- It is clear from (1.13) that if \( d = 2 \) then (1.13) is equivalent to the condition \( \det A_1 \neq \det A_2 \).
- Assume that \( A_1 \) and \( A_2 \) are isotropic on \( \Gamma \), i.e., \( A_1 = a_1 I \) and \( A_2 = a_2 I \) for some (positive) functions \( a_1, a_2 \) on \( \Gamma \). Then (i) and (ii) are equivalent to the conditions \( a_1 \neq a_2 \) and \( a_1 \Sigma_1 \neq a_2 \Sigma_2 \) on \( \Gamma \).

The proof of Theorem 4 is via Fourier analysis and given in Section 5. Applying Theorem 4 and Proposition 1, one recovers and extends the result of Lakhtanov and Vainberg’s in [20].

The discreteness of the eigenvalues of a general elliptic system equipped complementing boundary condition was investigated by Faierman [15] and references therein. The conditions given in [15] are not explicit and require that the coefficients are of class \( C^1(\overline{\Omega}) \). The proof in [15] involved the work of Agmon, Douglis, and Nirenberg in [1].

It is worth noting that ITE problem is invariant after a change of variables. In fact, let \( F \) be a diffeomorphism from \( \Omega \) into itself\(^\dagger\) such that \( F(x) = x \) on \( \Gamma \). Set \( U_1 = u_1 \circ F^{-1} \) in \( \Omega \). By the rule of a change of variables (see e.g. [24, Lemma 2]), \( u_1 \in H^1(\Omega) \) is a solution to the equation

\[
\text{div}(A_1 \nabla u_1) - \lambda \Sigma_1 u_1 = 0 \quad \text{in} \ \Omega
\]

if and only if

\[
\text{div}(F_* A_1 \nabla U_1) - \lambda F_* \Sigma_1 U_1 = 0 \quad \text{in} \ \Omega,
\]

where

\[
F_* A_1(y) = \frac{DF(x)A_1(x)DF^T(x)}{\det DF(x)} \quad \text{and} \quad F_* \Sigma_1(y) = \frac{\Sigma_1(x)}{\det DF(x)} \quad \text{with} \quad x = F^{-1}(y). \tag{1.14}
\]

Moreover,

\[
U_1 = u_1 \quad \text{and} \quad F_* A_1 \nabla U_1 \cdot \nu = A_1 \nabla u_1 \cdot \nu \quad \text{on} \ \Gamma. \tag{1.15}
\]

\(^\dagger\)In this paper, this means that \( F \) is bijective and \( F, F^{-1} \in C^1(\overline{\Omega}) \).
It follows that \((u_1, u_2) \in [H^1(\Omega)]^2\) is a solution of (1.3)-(1.4) if and only if \((U_1, U_2) = (u_1 \circ F^{-1}, u_2) \in [H^1(\Omega)]^2\) is a solution of the system
\[
\begin{align*}
\text{div}(F_* A_1 \nabla U_1) - \lambda F_* \Sigma_1 U_1 &= 0 \quad \text{in } \Omega \\
\text{div}(A_2 \nabla U_2) - \lambda \Sigma_2 U_2 &= 0 \quad \text{in } \Omega \\
U_1 &= U_2, \quad F_* A_1 \nabla U_1 \cdot \nu = A_2 \nabla U_2 \cdot \nu \quad \text{on } \Gamma.
\end{align*}
\tag{1.16}
\]
Moreover, \((u_1, u_2) \neq 0\) if and only if \((U_1, U_2) \neq 0\). Using this observation, we can extend the previous results in which the conditions on \((A_1, \Sigma_1)\) and \((A_2, \Sigma_2)\) are involved to the case where the same conditions hold for \((F_* A_1, F_* \Sigma_1)\) and \((A_2, \Sigma_2)\) for some diffeomorphism \(F\) verifying the condition \(F(x) = x\) on \(\Gamma\). We state here two results following from Theorems 1 and 3 in this direction as an illustration. A variant of Theorem 2 is left to the reader. Theorem 4 is already invariant under a change of variables.

**Proposition 2.** Assume that for some diffeomorphisms \(F\) from \(\Omega\) into itself such that \(F(x) = x\) on \(\Gamma\), and for \(0 \leq \alpha < 2\) and for some positive constant \(c\),
\[
F_* A_1 - A_2 \geq cd_0^2 I \quad \text{and} \quad F_* \Sigma_1 - \Sigma_2 \geq 0 \quad \text{in a neighborhood of } \Gamma.
\tag{1.17}
\]
The spectrum of (1.3)-(1.4) is discrete.

**Proposition 3.** Assume that for some diffeomorphism \(F\) from \(\Omega\) into itself such that \(F(x) = x\) on \(\Gamma\),
\[
F_* A_1 = A_2 \quad \text{and} \quad F_* \Sigma_1 - \Sigma_2 \geq cd_0^{2 \beta} \quad \text{in a neighborhood of } \Gamma,
\tag{1.18}
\]
for some \(0 \leq \beta < 2\) and for some positive constants \(c\). The spectrum of (1.3)-(1.4) is discrete.

Applying Propositions 2 and 3 one is able to obtain the discreteness of the spectrum of the ITE problem even in the case \(A_1 - A_2\) changes the sign in a neighborhood of \(\Gamma\).

We now describe briefly the ideas of the proof of Theorems 1, 2, 3, and 4. The proof of Theorems 1 and 2 are based on the multiplier technique while the proof of Theorem 3 is based on the Fourier analysis. The key point is to establish the well-posedness and the compactness of \(T\) given by (1.6) for some \(\gamma_0 \in \mathbb{C}\). Concerning Theorem 1 the well-posedness of and the compactness of \(T\) take place for \(\gamma_0 = \lambda_0\) for some large positive \(\lambda_0\). The existence of a solution of the system (1.5) is based on a priori estimates (Lemma 1) for a Cauchy problem with its roots from [29] and via the limiting absorption process (Lemma 7). To require the information only near the boundary \(\Gamma\), \(\lambda_0\) is chosen to be large and a standard exponential decay estimate for elliptic equations is involved (Lemma 2). Since the exponent \(\alpha\) might be positive, some weighted spaces are involved and the solutions are not in \([H^1(\Omega)]^2\) as usual. The proof of the uniqueness of (1.5) faces some issues due to the lack of the regularity of the solutions. To overcome this, we introduce the concept of viscosity solution (Definition 2), the terminology is inspired by from the one of systems of conservation laws. The compactness of \(T\) follows from the condition \(0 \leq \alpha < 2\) in (1.8). Concerning Theorem 2 the well-posedness of and the compactness of \(T\)

---

Footnote: One can use two diffeomorphisms \(F_1, F_2\) and require the corresponding conditions on \((F_1, A_1, F_1, \Sigma_1)\) and \((F_2, A_2, F_2, \Sigma_2)\) to obtain the discreteness of the ITE problem. However, the same conditions hold by using the diffeomorphisms \(F_1 \circ F_2^{-1}, I\).
take place for $\gamma_0 = i\lambda_0$ for some small positive $\lambda_0$. This modification is necessary for the proof of the uniqueness. As in the proof of Theorem 1, a priori estimates hold for (3.1) (Lemma 5). Nevertheless, the uniqueness of (3.1) even for smooth solutions does not follow directly from the a priori estimates as in the proof of Theorem 1. Additional arguments are required in this case (Lemma 10). Beside these points, the proof of Theorem 2 follows similarly as the one of Theorem 1. The proof of Theorem 3 is somehow in the spirit of the one of Theorem 1 but with the following key difference. While the same arguments as used in the proof of Theorem 1 imply the uniqueness of $T$ with $\gamma_0 = \lambda_0$ for large $\lambda_0$, they do not imply the compactness of $T$ since the first condition in (1.8) does not hold. To be able to deal with the situation in which $A_1 = A_2$ in a neighborhood of $\Gamma$, we make some modifications on $T$. The idea is to take into account the fact from the a priori estimates $u_1 - u_2$ is more regular than $u_1$ and $u_2$. After the modification of $T$, the proof of Theorem 3 is in the spirit of the one of Theorem 1 but the functional spaces used in this case is different and somehow more involved and the theory of compact analytic operator is used. The proof of Theorem 4 is via the Fourier analysis. The approach is based on the computation of solutions of the Cauchy problems in half space via Fourier analysis. We then use local charts and the exponential decays of the solutions to deal with the general case. The analysis is also in the spirit of the work Agmon, Douglis, and Nirenberg [1].

The analysis in this paper is devoted to the study of the Cauchy problem (1.3) and (1.4). There is a connection between the study of the Cauchy problem and the Helmholtz equations with sign changing coefficients modelling negative index materials. This connection was pointed out by H-M. Nguyen in [29]. Various ideas in this paper were introduced there. It is worth noting that resonance might appear in various interesting potential applications of negative index materials such as superlensing using complementary media [25], cloaking using complementary media [28, 32], cloaking a source via anomalous localized resonance, see e.g. [26, 27, 30], and cloaking an object via anomalous localized resonance in [31].

The paper is organized as follows. The proof of Theorems 1, 2, 3, and 4 are given in Sections 2, 3, 4, and 5 respectively.

2 On the case $A_1 \geq A_2$ and $\Sigma_1 \geq \Sigma_2$ in a neighborhood of $\Gamma$ - Proof of Theorem 1

The section is devoted to the proof of Theorem 1. To this end, as mentioned in the introduction, we first establish the well-posedness of the following system, for some $\lambda_0 > 0$,

$$
\begin{align*}
\text{div}(A_1 \nabla u_1) - \lambda_0 \Sigma_1 u_1 &= g_1 \quad \text{in } \Omega, \\
\text{div}(A_2 \nabla u_2) - \lambda_0 \Sigma_2 u_2 &= g_2 \quad \text{in } \Omega, \\
u_1 &= u_2, \quad A_1 \nabla u_1 \cdot \nu = A_2 \nabla u_2 \cdot \nu \quad \text{on } \Gamma,
\end{align*}
$$

(2.1)

for a given pair $(g_1, g_2)$ in an appropriate space. We then define the operator

$$
T_1(f_1, f_2) = (u_1, u_2) \quad \text{where } (u_1, u_2) \text{ is the unique solution of (2.1) with } (g_1, g_2) = (\Sigma_1 f_1, \Sigma_2 f_2)
$$

(2.2)

and prove the compactness for it.
The proof of the well-posedness of (2.1) is as follows. The existence of a solution of (2.1) is based on a priori estimates (Lemma 1) for a Cauchy problem with roots from [29] and via a limiting absorption process (Lemma 7). To require only the information near the boundary, $\lambda_0$ is chosen to be large and a standard exponential decay estimate for elliptic equations is involved (Lemma 2). The solutions are not in $[H^1(\Omega)]^2$ as usual and hence it is not clear whether the a priori estimates hold for the difference of two solutions. To overcome this issue, we introduce the concept of viscosity solutions (Definition 2). We now present the details of the proof. We start with the following result (Lemma 1) which plays an important role in the proof of Theorem 1. Let Lemma 1 is a more detailed version of [29, Lemma 9] and its proof follows closely from there.

**Lemma 1.** Let $g = (g_1, g_2) \in [L^2(\Omega)]^2$ and $h \in H^{-1/2}(\Gamma)$, and $\lambda_0 \geq 0$, and let $v = (v_1, v_2) \in [H^1(\Omega)]^2$ satisfy

\[
\begin{align*}
\text{div}(A_1 \nabla v_1) - \lambda_0 \Sigma_1 v_1 &= g_1 & \text{in } \Omega, \\
\text{div}(A_2 \nabla v_2) - \lambda_0 \Sigma_2 v_2 &= g_2 & \text{in } \Omega, \\
v_1 &= v_2, \quad A_1 \nabla v_1 \cdot \nu = A_2 \nabla v_2 \cdot \nu + h & \text{on } \Gamma.
\end{align*}
\]

We have, with $w = v_1 - v_2$,

\[
\int_\Omega \langle A_1 \nabla w, \nabla w \rangle + \lambda_0 \Sigma_1 |w|^2 = - \int_\Omega (g_1 - g_2) \bar{w} + \int_\Omega \lambda_0 (\Sigma_2 - \Sigma_1) v_2 \bar{w} + \int_\Omega \langle [A_2 - A_1] \nabla v_2, \nabla w \rangle
\]

and

\[
\int_\Omega \langle [A_1 - A_2] \nabla v_2, \nabla w \rangle + \langle [A_1 - A_2] \nabla v_2, \nabla v_2 \rangle + \lambda_0 (\Sigma_1 - \Sigma_2) |v_2|^2
\]

\[
= \int_\Omega g_2 \bar{w} - (\bar{g}_1 - \bar{g}_2) v_2 + \lambda_0 (\Sigma_2 - \Sigma_1) v_2 \bar{w} + \int_\Gamma \bar{h} v_2. \quad (2.5)
\]

As a consequence of (2.4) and (2.5), we obtain

\[
\int_\Omega \langle A_1 \nabla w, \nabla w \rangle + \lambda_0 \Sigma_1 |w|^2 \leq 4 \mathcal{M}(v, g, h) + \int_\Omega \langle [A_2 - A_1] \nabla v_2, \nabla w \rangle + \lambda_0 (\Sigma_2 - \Sigma_1) v_2 \bar{w}, \quad (2.6)
\]

and

\[
\int_\Omega \langle [A_1 - A_2] \nabla v_2, \nabla v_2 \rangle + \lambda_0 (\Sigma_1 - \Sigma_2) |v_2|^2
\]

\[
\leq 4 \mathcal{M}(v, g, h) + \int_\Omega |\langle [A_2 - A_1] \nabla v_2 \nabla w \rangle| + \lambda_0 |(\Sigma_2 - \Sigma_1) v_2 \bar{w}|, \quad (2.7)
\]

where

\[
\mathcal{M}(v, g, h) := ||v||_{L^2(\Omega)} ||g||_{L^2(\Omega)} + ||h||_{H^{-1/2}(\Gamma)} ||v||_{H^{1/2}(\Gamma)}.
\]

Here and in what follows in this paper, for a complex number $z$, $\bar{z}$ denotes its conjugate.

**Proof.** We derive from the definition of $w$ that $w = 0$ on $\Gamma$ and

\[
\text{div}(A_1 \nabla w) - \lambda_0 \Sigma_1 w = g_1 - g_2 - \lambda_0 (\Sigma_2 - \Sigma_1) v_2 + \text{div}([A_2 - A_1] \nabla v_2) \quad \text{in } \Omega. \quad (2.9)
\]
Multiplying this equation by \( \bar{w} \) and integrating on \( \Omega \), we have

\[
\int_{\Omega} \langle A_1 \nabla w, \nabla w \rangle + \lambda_0 \Sigma_1 |w|^2 = -\int_{\Omega} (g_1 - g_2) \bar{w} + \int_{\Omega} \lambda_0 (\Sigma_2 - \Sigma_1)v_2 \bar{w} + \int_{\Omega} \langle [A_2 - A_1] \nabla v_2, \nabla w \rangle;
\]

which is (2.4). Multiplying the equation of \( v_2 \) by \( \bar{w} \) and integrating on \( \Omega \), we obtain

\[
\int_{\Omega} -\langle A_2 \nabla v_2, \nabla w \rangle - \lambda_0 \Sigma_2 v_2 \bar{w} = \int_{\Omega} g_2 \bar{w}.
\]

(2.10)

It is clear that

\[
-A_2 \nabla w = (A_1 - A_2) \nabla w - A_1 \nabla w + (A_2 - A_1) \nabla v_2 + (A_1 - A_2) \nabla v_2 \text{ in } \Omega,
\]

and, by (2.3),

\[
\text{div}(A_1 \nabla w) - \text{div}([A_2 - A_1] \nabla v_2) = \lambda_0 \Sigma_1 w + g_1 - g_2 - \lambda_0 (\Sigma_2 - \Sigma_1)v_2 \text{ in } \Omega.
\]

(2.11)

Since

\[
(A_1 \nabla w - (A_2 - A_1) \nabla v_2) \cdot \nu = (A_1 \nabla v_1 - A_2 \nabla v_2) \cdot \nu = h \text{ on } \Gamma,
\]

it follows from (2.10), (2.11), and (2.12) that

\[
\int_{\Omega} \langle [A_1 - A_2] \nabla v_2, \nabla w \rangle + \langle [A_1 - A_2] \nabla v_2, \nabla v_2 \rangle + \lambda_0 (\Sigma_1 - \Sigma_2)|v_2|^2
\]

\[
= \int_{\Omega} g_2 \bar{w} - (g_1 - g_2)v_2 + \lambda_0 (\Sigma_2 - \Sigma_1)v_2 \bar{w} + \int_{\Gamma} h v_2;
\]

which is (2.5).

Assertion (2.6) and (2.7) are direct consequences of (2.4) and (2.5). The proof is complete.

□

In what follows, we denote, for \( s > 0 \),

\[
\Omega_s := \{ x \in \Omega; d_{\Gamma}(x) < s \}.
\]

The following exponential decay property for elliptic equations is useful for our analysis.

**Lemma 2.** Let \( \lambda > 1 \), \( f \in L^2(\Omega) \), \( A \) be a matrix-valued function, and \( \Sigma \) be real function defined in \( \Omega \) such that, for some \( \Lambda > 1 \),

\[
\Lambda^{-1} \leq A \leq \Lambda \quad \text{and} \quad \Lambda^{-1} \leq \Sigma \leq \Lambda \text{ in } \Omega.
\]

(2.13)

Let \( u \in H^1_{\text{loc}}(\Omega) \) be a solution to the equation \( \text{div}(A \nabla u) - \lambda \Sigma u = f \) in \( \Omega \). For all \( s > 0 \), there exist two positive constants \( c_1 \) and \( c_2 \), depending only on \( \Lambda \), \( s \), and \( \Omega \), such that

\[
\|u\|_{H^1(\Omega_s)} \leq c_1 \exp(-c_2 \sqrt{\lambda}) \|u\|_{L^2(\Omega_s)} + c_1 \|f\|_{L^2(\Omega)}.
\]

(2.14)
Proof. Lemma 2 is quite standard and its proof presented here is in the spirit of the one of [16] Theorem 2.2. Let $U \in H^1_0(\Omega)$ be the unique solution of the equation $\text{div}(A\nabla U) - \lambda \Sigma U = f$ in $\Omega$. Then $\|U\|_{H^1(\Omega)} \leq C\|f\|_{L^2(\Omega)}$. Here and in what follows in this proof $C$ denotes a positive constant depending only on $\Lambda$, $s$, and $\Omega$. By considering $u - U$, without loss of generality, one might assume that $f = 0$ in $\Omega$; this is assumed from now. Fix $\varphi \in C^2(\Omega)$ such that $\varphi = cs$ in $\Omega \setminus \Omega_s$ and $\varphi = 0$ in $\Omega_{s/2}$, and $|\nabla \varphi| \leq c$ in $\Omega$ where $c$ is a small positive constant defined later (the smallness of $c$ depends only on $\Lambda$ and $\Omega$, it is independent of $s$). Set $\phi(x) = e^{\sqrt{\lambda} \varphi(x)}$ and $v(x) = u(x)\phi(x)$ for $x \in \Omega$. Since $\text{div}(A\nabla u) - \lambda \Sigma u = 0$ in $\Omega$, it follows that
\[
\text{div}(A\nabla v) - \lambda \Sigma v = \sqrt{\lambda} \text{div}(vA\nabla \varphi) + \sqrt{\lambda} A\nabla v \nabla \varphi - \lambda v A\nabla \varphi \nabla \varphi \text{ in } \Omega.
\] (2.15)

Note that, by Cacciopoli’s inequality,
\[
\|u\|_{H^1(\Omega_{s/2})} \leq C\sqrt{\lambda} \|u\|_{L^2(\Omega_s)}.
\]
Hence there exists $\tau \in (s/4, s/2)$ such that $\|u\|_{H^1(\partial(\Omega \setminus \Omega_{\tau}))} \leq C\sqrt{\lambda} \|u\|_{L^2(\Omega_s)}$. Multiplying (2.15) by $\bar{v}$, integrating by parts in $\Omega \setminus \Omega_{\tau}$, and using (2.13) and the fact that $\varphi = 0$ in $\Omega_{s/2}$, we have
\[
\int_{\Omega \setminus \Omega_{\tau}} |\nabla v|^2 + \lambda |v|^2 \leq C \left( \int_{\Omega \setminus \Omega_{\tau}} \sqrt{\lambda} |v| |\nabla \varphi| + \lambda |v|^2 |\nabla \varphi|^2 \right) + C\lambda \|u\|_{L^2(\Omega_s)}^2.
\] (2.16)

By taking $c$ small enough, one can absorb the first term in the RHS of (2.16) by the LHS, and the conclusion follows. \hfill \Box

We next recall the following compactness result from [29] Lemma 7\footnote{In fact, [29] Lemma 7 is stated for $u_n \in H^1(\Omega)$, however the result also holds for $u_n \in H^1_{\text{loc}}(\Omega)$ and the proof is almost unchanged.}.

Lemma 3. Let $0 \leq \alpha < 2$ and $(u_n) \subset H^1_{\text{loc}}(\Omega)$. Assume that
\[
\sup_n \int_{\Omega} \left( d_1^\alpha |\nabla u_n|^2 + |u_n|^2 \right) \, dx < +\infty.
\] (2.17)
Then $(u_n)$ is relatively compact in $L^2(\Omega)$.

Using Lemma 3 we can prove

Lemma 4. Let $\sigma$ be a bounded real function defined in $\Omega$ such that $\int_{\Omega} \sigma \neq 0$. There exists a positive constant $C$ such that
\[
\int_{\Omega} |u|^2 \leq C \left( \int_{\Omega} |\nabla u|^2 \right) + \int_{\Omega} d_1^\alpha |\nabla u|^2 \quad \forall u \in H^1_{\text{loc}}(\Omega).
\] (2.18)

Proof. Lemma 4 can be derived from Lemma 3 by a contradiction argument as follows. Assume that (2.18) does not hold. There exists a sequence $(u_n) \subset H^1_{\text{loc}}(\Omega)$ such that
\[
\sup_n \left( \int_{\Omega} |\nabla u_n|^2 \right) + \int_{\Omega} d_1^\alpha |\nabla u_n|^2 \leq \int_{\Omega} |u_n|^2 = 1.
\] (2.19)
By Lemma 3, one may assume that \( u_n \to u \) strongly in \( L^2(\Omega) \) and almost everywhere in \( \Omega \), and \( u_n \to u \) weakly in \( H^1_{loc}(\Omega) \). We deduce from (2.19) that

\[
\left| \int_\Omega \sigma |u|^2 \right| + \int_\Omega d_n^2 |\nabla u|^2 = 0 \quad \text{and} \quad \int_\Omega |u|^2 = 1.
\]

The first identity implies that \( u = 0 \) in \( \Omega \) since \( \Omega \) is connected; this contradicts the second identity. Hence (2.18) holds. \( \square \)

Using Lemma 2 we obtain

**Lemma 5.** Let \( 0 \leq \alpha < 2 \) and assume that, for some \( \tau > 0 \) and \( c > 0 \),

\[
A_1 - A_2 \geq cd_1^0 I \quad \text{and} \quad \Sigma_1 - \Sigma_2 \geq 0 \text{ in } \Omega_\tau, \tag{2.20}
\]

There exists \( \Lambda_0 > 1 \) such that if \( \lambda_0 \geq \Lambda_0 \), and \( v = (v_1, v_2) \in [H^1_{loc}(\Omega)]^2 \) satisfy the system

\[
\begin{cases}
\text{div}(A_1 \nabla v_1) - \lambda_0 \Sigma_1 v_1 = g_1 & \text{in } \Omega, \\
\text{div}(A_2 \nabla v_2) - \lambda_0 \Sigma_2 v_2 = g_2 & \text{in } \Omega,
\end{cases} \tag{2.21}
\]

and (2.6) and (2.7) hold for some \( g = (g_1, g_2) \in [L^2(\Omega)]^2 \) and \( h \in H^{-1/2}(\Gamma) \), then, with \( w = v_1 - v_2 \),

\[
\int_{\Omega_\tau} \langle [A_1 - A_2] \nabla v_2, \nabla v_2 \rangle + \int_{\Omega \setminus \Omega_\tau} |\nabla v_2|^2 + \int_{\Omega} |\nabla w|^2 + |v_2|^2 \leq C \left( \mathcal{M}(v, g, h) + \|g\|^2_{L^2(\Omega)} \right), \tag{2.22}
\]

for some positive constant \( C \) independent of \( v, g, \) and \( h \) where \( \mathcal{M}(v, g, h) \) is defined by (2.8). In particular if \( v \in [H^1(\Omega)]^2 \) is a solution of (2.3) then (2.22) holds.

**Proof.** Since \( A_1 - A_2 \geq 0 \) and \( \Sigma_1 - \Sigma_2 \geq 0 \) in \( \Omega_\tau \), we have, for all \( \gamma > 0 \),

\[
\gamma \int_{\Omega_\tau} \langle [A_1 - A_2] \nabla w, \nabla w \rangle + \frac{1}{\gamma} \int_{\Omega_\tau} \langle [A_1 - A_2] \nabla v_2, \nabla v_2 \rangle \geq 2 \int_{\Omega_\tau} \langle [A_1 - A_2] \nabla w, \nabla v_2 \rangle \tag{2.23}
\]

and

\[
\gamma \int_{\Omega_\tau} (\Sigma_1 - \Sigma_2)|w|^2 + \frac{1}{\gamma} \int_{\Omega_\tau} (\Sigma_1 - \Sigma_2)|v_2|^2 \geq 2 \int_{\Omega_\tau} (\Sigma_1 - \Sigma_2)|wv_2|. \tag{2.24}
\]

Taking \( \gamma > 1 \) and close to 1, and adding (2.6) and (2.7), we derive that

\[
\int_{\Omega_\tau} \langle \nabla w, \nabla w \rangle + \lambda_0 |w|^2 + \langle [A_1 - A_2] \nabla v_2, \nabla v_2 \rangle \leq C \mathcal{M}(v, g, h) + C \int_{\Omega \setminus \Omega_\tau} |\nabla v_2||\nabla w| + \lambda_0 |v_2||w| + |\nabla v_2|^2 + \lambda_0 |v_2|^2. \tag{2.25}
\]

Here and in what follows in this proof, \( C \) denotes a positive constant independent of \( v, g, h \) and \( \lambda_0 \).
On the other hand, applying Lemma 6 below, we have
\[
\int_{\Omega} \lambda d_1^{\alpha+2}|v_2|^2 \leq C \int_{\Omega} d_1^{\alpha}|\nabla v_2|^2 + C\|g_2\|_{L^2(\Omega)}^2.
\]
This implies, by applying Lemma 4 with \( \sigma = d_1^{\alpha+2} \),
\[
\int_{\Omega} |v_2|^2 \leq C \int_{\Omega} d_1^{\alpha}|\nabla v_2|^2 + C\|g_2\|_{L^2(\Omega)}^2.
\tag{2.26}
\]
The conclusion now follows from (2.25) and (2.26) by applying Lemma 2 to \( v_2 \) and taking \( \Lambda_0 \) large enough.

**Remark 1.** Using (2.26), one can weaken the assumption \( \Sigma_1 - \Sigma_2 \geq 0 \) in (2.20) by \( \Sigma_1 - \Sigma_2 \geq -\hat{c} \) in a neighborhood of \( \Gamma \) for some small positive constant \( \hat{c} \).

In the proof of Lemma 6, we used the following result.

**Lemma 6.** Let \( \lambda > 1 \), \( \alpha \geq 0 \), \( f \in L^2(\Omega) \), \( A \) be a matrix-valued function, and \( \Sigma \) be real function defined in \( \Omega \) such that
\[
\Lambda^{-1} \leq A \leq \Lambda \quad \text{and} \quad \Lambda^{-1} \leq \Sigma \leq \Lambda \quad \text{in} \quad \Omega.
\]
Let \( u \in H^1_{loc}(\Omega) \) be a solution to the equation \( \text{div}(A\nabla u) - \lambda \Sigma u = f \) in \( \Omega \). We have
\[
\int_{\Omega} \lambda d_1^{\alpha+2}|u|^2 \leq C \int_{\Omega} d_1^{\alpha}|\nabla u|^2 + C\|f\|_{L^2(\Omega)}^2 \tag{2.27}
\]
and
\[
\int_{\Omega} d_1^{\alpha+2}|\nabla u|^2 \leq C \int_{\Omega} \lambda d_1^{\alpha}|u|^2 + C\|f\|_{L^2(\Omega)}^2 \tag{2.28}
\]
for some positive constant \( C \) independent of \( \lambda \), \( f \), and \( u \).

**Proof.** We only prove (2.27); the proof of (2.28) follows similarly. Set \( \Gamma_s = \partial(\Omega \setminus \Omega_s) \) and \( \varphi_s(x) = \text{dist}(x, \Gamma_s)^{\alpha/2+1} \). Multiplying the equation of \( u \) by \( \tilde{u}\varphi_s^2 \) and integrating in \( \Omega \setminus \Omega_s \), we have
\[
\int_{\Omega \setminus \Omega_s} \lambda |\tilde{u}\varphi_s|^2 - \int_{\Omega \setminus \Omega_s} |\nabla u|^2 \varphi_s^2 \leq C \int_{\Omega \setminus \Omega_s} |\nabla u| \|\nabla \varphi_s\| |\tilde{u}\varphi_s| + |f||\tilde{u}\varphi_s^2|.
\]
Since \( |\nabla u| \|\nabla \varphi_s\| |\tilde{u}\varphi_s| \leq |\nabla \varphi_s|^2 |\nabla u|^2 + |\varphi_s|^2 |\tilde{u}|^2 \leq C \text{dist}^\alpha(x, \Gamma_s) |\nabla u|^2 + |\varphi_s|^2 |\tilde{u}|^2 \), it follows that
\[
\int_{\Omega \setminus \Omega_s} \lambda |\tilde{u}\varphi_s|^2 \leq C \int_{\Omega \setminus \Omega_s} \text{dist}(x, \Gamma_s)^\alpha |\nabla u|^2 + |f|^2.
\]
Letting \( s \to 0 \), applying the Fatou lemma to the LHS and the dominated convergence theorem to the RHS, one obtains (2.26).

**Remark 2.** The same conclusion of Lemma 6 holds if instead of imposing that \( \text{div}(A\nabla u) - \lambda \Sigma u = f \) in \( \Omega \), one assumes \( \text{div}(A\nabla u) - i\lambda \Sigma u = f \) in \( \Omega \).
The next basic result involving the limiting absorption process is

**Lemma 7.** Let \( \lambda_0 \geq 0, 0 < \delta < 1, g = (g_1, g_2) \in [L^2(\Omega)]^2 \). There exists a unique solution \( v_\delta = (v_{1,\delta}, v_{2,\delta}) \in [H^1(\Omega)]^2 \) of the system

\[
\begin{cases}
\text{div} \left( (1 + i\delta)A_1 \nabla v_{1,\delta} \right) - \lambda_0 \Sigma_1 v_{1,\delta} - i\delta v_{1,\delta} = g_1 & \text{in } \Omega, \\
\text{div} \left( (1 - i\delta)A_2 \nabla v_{2,\delta} \right) - \lambda_0 \Sigma_2 v_{2,\delta} + i\delta v_{2,\delta} = g_2 & \text{in } \Omega, \\
v_{1,\delta} = v_{2,\delta}, (1 + i\delta)A_1 \nabla v_{1,\delta} \cdot \nu = (1 - i\delta)A_2 \nabla v_{2,\delta} \cdot \nu & \text{on } \Gamma.
\end{cases}
\]

Moreover,

\[
\|v_\delta\|_{H^1(\Omega)}^2 \leq \frac{C}{\delta} \|g\|_{L^2(\Omega)} \|v_\delta\|_{L^2(\Omega)},
\]

for some positive constant \( C \) independent of \( g \) and \( \delta \). Consequently,

\[
\|v_\delta\|_{H^1(\Omega)} \leq \frac{C}{\delta} \|g\|_{L^2(\Omega)}.
\]

**Proof.** Set

\[
X := \{ \phi = (\phi_1, \phi_2) \in [H^1(\Omega)]^2 : \phi_1 - \phi_2 \in H^1_0(\Omega) \}.
\]

Then \( X \) is a Hilbert space equipped the scalar product:

\[
\langle (\varphi_1, \varphi_2), (\phi_1, \phi_2) \rangle_X := \langle \varphi_1, \phi_1 \rangle_{H^1(\Omega)} + \langle \varphi_2, \phi_2 \rangle_{H^1(\Omega)}.
\]

Let \( a \) be the bilinear form on \( X \times X \) defined as follows, for \( \varphi = (\varphi_1, \varphi_1), \phi = (\phi_1, \phi_2) \in X, \)

\[
a(\varphi, \phi) = \int_\Omega (1 + i\delta) \langle A_1 \nabla \varphi_1, \nabla \phi_1 \rangle + \int_\Omega \lambda_0 \Sigma_1 \varphi_1 \phi_1 + \int_\Omega i\delta \varphi_1 \phi_1
\]

\[
- \left[ \int_\Omega (1 - i\delta) \langle A_2 \nabla \varphi_2, \nabla \phi_2 \rangle + \lambda_0 \int_\Omega \Sigma_2 \varphi_2 \phi_2 - \int_\Omega i\delta \varphi_2 \phi_2 \right],
\]

and let \( b \) be the linear form on \( X \) given by

\[
b(\phi) = \int_\Omega -g_1 \phi_1 + g_2 \phi_2.
\]

It is easy to check that \( a \) is a continuous on \( X \times X \) and \( b \) is continuous on \( X \); moreover,

\[
|b(\phi)| \leq \|(g_1, g_2)\|_{L^2(\Omega)} \|(\phi_1, \phi_2)\|_{L^2(\Omega)}.
\]

On the other hand, by considering the imaginary part of \( a \), we have

\[
\Im(a(\phi, \phi)) \geq C\delta \|\phi\|_X^2;
\]

which implies the coercivity of \( a \). By Lax-Milgram’s theorem, there exists a unique \( \varphi \in X \) such that

\[
a(\varphi, \phi) = b(\phi) \text{ for all } \phi \in X.
\]
One can check that \((v_{1,\delta}, v_{2,\delta}) = \varphi\) is a solution of \((2.29)\); moreover, if \((v_{1,\delta}, v_{2,\delta}) \in X\) is a solution of \((2.29)\) then \(\varphi = (v_{1,\delta}, v_{2,\delta})\) is a solution of \((2.30)\). The proof is complete. \(\square\)

Let \(\tau > 0\) be such that \((1.8)\) holds in \(\Omega_\tau\). Define

\[
\mathcal{H}(\Omega) := \left\{ (u_1, u_2) \in [H^1_0(\Omega) \cap L^2(\Omega)]^2; \ u_1 - u_2 \in H^1_0(\Omega), \ \|(u_1, u_2)\|_{\mathcal{H}(\Omega)} < +\infty \right\}, \tag{2.37}
\]

where, for \(u = (u_1, u_2)\) and \(v = (v_1, v_2)\) in \([H^1_0(\Omega)]^2\),

\[
\langle u, v \rangle_{\mathcal{H}(\Omega)} = \int_{\Omega} \nabla(u_1 - u_2) \nabla(v_1 - v_2) + \int_{\Omega \setminus \Omega_\tau} \nabla u \nabla v + \int_{\Omega_\tau} \langle [A_1 - A_2] \nabla u, \nabla v \rangle + \int_{\Omega} uv. \tag{2.38}
\]

One can check that \(\mathcal{H}(\Omega)\) equipped the scalar product in \((2.38)\) is a Hilbert space.

We have

**Definition 1.** Let \((g_1, g_2) \in [L^2(\Omega)]^2\). A pair of functions \((v_1, v_2) \in \mathcal{H}(\Omega)\) is called a weak solution of \((2.1)\) if

\[
\begin{align*}
\text{div}(A_1 \nabla v_1) - \lambda_0 \Sigma_1 v_1 &= g_1 \quad \text{in } \Omega, \\
\text{div}(A_2 \nabla v_2) - \lambda_0 \Sigma_2 v_2 &= g_2 \quad \text{in } \Omega, \\
(A_1 \nabla v_1 - A_2 \nabla v_2) \cdot \nu &= 0 \quad \text{on } \Gamma
\end{align*}
\tag{2.39}
\]

**Remark 3.** Since \(\text{div}(A_1 \nabla v_1 - A_2 \nabla v_2) \in L^2(\Omega)\) and \(A_1 \nabla v_1 - A_2 \nabla v_2 = A_1 \nabla (v_1 - v_2) + (A_1 - A_2) \nabla v_2 \in L^2(\Omega)\) by \((2.38)\). The last identity on \(\Gamma\) in \((2.39)\) makes sense.

Using Lemmas 4 and 7 one can construct a weak solution \((v_1, v_2) \in \mathcal{H}(\Omega)\) of \((2.1)\). More precisely, we have

**Proposition 4.** Assume \((1.8)\) and let \(\delta \in (0, 1)\), \(g = (g_1, g_2) \in [L^2(\Omega)]^2\) and \(v_\delta = (v_{1,\delta}, v_{2,\delta}) \in [H^1(\Omega)]^2\) be the unique solution of \((2.29)\). There exists \(\Lambda_0 > 1\) such that if \(\lambda_0 > \Lambda_0\) then

\[
\|v_\delta\|_{\mathcal{H}(\Omega)} \leq C\|g\|_{L^2(\Omega)}, \tag{2.40}
\]

for some positive constant \(C\) independent of \(\delta\) and \(g\). As a consequence, there exists a weak solution \(v = (v_1, v_2) \in \mathcal{H}(\Omega)\) of \((2.1)\) such that

\[
\|v\|_{\mathcal{H}(\Omega)} \leq C\|g\|_{L^2(\Omega)}. \tag{2.41}
\]

**Proof.** Applying Lemma 7 we have

\[
\|v_\delta\|_{H^1(\Omega)} \leq \frac{C}{\delta}\|g\|_{L^2(\Omega)}. \tag{2.42}
\]

Here and in what follows in the proof, \(C\) denotes a positive constant independent of \(\delta\) and \(g\). Rewriting the system of \(v_\delta\), we get

\[
\begin{align*}
\text{div} \left( A_1 \nabla v_{1,\delta} \right) - \lambda_0 \Sigma_1 v_{1,\delta} &= g_{1,\delta} \quad \text{in } \Omega, \\
\text{div} \left( A_2 \nabla v_{2,\delta} \right) - \lambda_0 \Sigma_2 v_{2,\delta} &= g_{2,\delta} \quad \text{in } \Omega, \\
A_1 \nabla v_{1,\delta} \cdot \nu - A_2 \nabla v_{2,\delta} \cdot \nu &= h_\delta \quad \text{on } \Gamma
\end{align*}
\tag{2.43}
\]
where
\[ g_{1,\delta} = \frac{1}{1 + i\delta} \left( g_1 + i\delta v_{1,\delta} - i\delta\lambda_0 \Sigma_1 v_{1,\delta} \right), \quad g_{2,\delta} = \frac{1}{1 - i\delta} \left( g_2 - i\delta v_{2,\delta} + i\delta\lambda_0 \Sigma_2 v_{2,\delta} \right), \] (2.44)
and
\[ h_{\delta} = -i\delta(A_1 \nabla v_{1,\delta} + A_2 \nabla v_{2,\delta}) \cdot \nu. \]
Denote \( g_{\delta} = (g_{1,\delta}, g_{2,\delta}). \) For \( \lambda_0 \geq \Lambda_0 \) and for large \( \lambda_0, \) we have, by Lemmas 1 and 5 and (2.42),
\[ \|v_{\delta}\|_{H(\Omega)} \leq C \left( \mathcal{M}(v_{\delta}, g_{\delta}, h_{\delta}) + \|g_{\delta}\|_{L^2(\Omega)}^2 \right) \leq C\|g\|_{L^2(\Omega)}^2 + C\|g\|_{L^2(\Omega)}\|v\|_{L^2(\Omega)}. \] (2.45)
Here we used the fact that, by (2.42),
\[ \|g_{\delta}\|_{L^2(\Omega)} \leq C\|g\|_{L^2(\Omega)} \]
and in addition by the trace theory and (2.30),
\[ \|h_{\delta}\|_{H^{-1/2}(\Gamma)} \|v_{\delta}\|_{H^{1/2}(\Gamma)} \leq C\delta \left( \|g_{\delta}\|_{L^2(\Omega)} + \|v_{\delta}\|_{H^1(\Omega)} \right) \|v_{\delta}\|_{H^1(\Omega)} \]
\[ \leq C \left( \|g\|_{L^2(\Omega)}^2 + \|g\|_{L^2(\Omega)}\|v_{\delta}\|_{L^2(\Omega)} \right). \]

We derive from (2.45) that
\[ \|v_{\delta}\|_{H(\Omega)} \leq C\|g\|_{L^2(\Omega)}; \] (2.46)
which is (2.40). Assertion (2.41) follows from (2.40) by a standard compactness argument. The details are left to the reader. \( \square \)

We next discuss the uniqueness of the weak solutions. To obtain the uniqueness, one needs to show that \((v_1, v_2) = (0, 0)\) if \((v_1, v_2)\) is a weak solution of (2.39) with \(g_1 = g_2 = 0.\) A natural way for this is to obtain (2.6) and (2.7) with \(g_1 = g_2 = 0\) and \(h = 0\) for \((v_1, v_2)\) and then to apply Lemma 5. Note that estimate (2.6) follows from (2.4) and estimate (2.7) follows from (2.5). It is not clear that (2.4) and estimate (2.7) hold for \((v_1, v_2) \in H(\Omega)\) which satisfies (2.39). One way to remedy the situation is to consider only solutions of (2.39) which come from the limiting absorption principle process given in Lemma 7. Even making this restriction, the proof of the uniqueness is still not straightforward since (2.4) and (2.5) are non-linear with respect to \((v_1, v_2).\) To overcome this difficulty, we use an kind of relaxation argument which we now introduce. Consider \((\hat{v}_1, \hat{v}_2) \in H(\Omega)\) and \((\hat{g}_1, \hat{g}_2) \in [L^2(\Omega)]^2\) such that
\[
\left\{ \begin{array}{l}
\text{div}(A_1 \nabla \hat{v}_1) - \lambda_0 \Sigma_1 \hat{v}_1 = \hat{g}_1 \text{ in } \Omega, \\
\text{div}(A_2 \nabla \hat{v}_2) - \lambda_0 \Sigma_2 \hat{v}_2 = \hat{g}_2 \text{ in } \Omega, \\
(A_1 \nabla \hat{v}_1 - A_2 \nabla \hat{v}_2) \cdot \nu = 0 \text{ on } \Gamma.
\end{array} \right. \] (2.47)
Let \((g_1, g_2) \in [L^2(\Omega)]^2\) and \((v_1, v_2) \in [H^1(\Omega)]^2\) be a solution of the system
\[
\left\{ \begin{array}{l}
\text{div}(A_1 \nabla v_1) - \lambda_0 \Sigma_1 v_1 = g_1 \text{ in } \Omega, \\
\text{div}(A_2 \nabla v_2) - \lambda_0 \Sigma_2 v_2 = g_2 \text{ in } \Omega, \\
v_1 = v_2 \text{ on } \Gamma.
\end{array} \right. \]
Set \( w = v_1 - v_2 \) and \( \tilde{w} = \tilde{v}_1 - \tilde{v}_2 \). We claim that
\[
\int_{\Omega} \langle A_1 \nabla w, \nabla \tilde{w} \rangle + \lambda_0 \Sigma_1 w \tilde{w} = - \int_{\Omega} (g_1 - g_2) \tilde{w} + \int_{\Omega} \lambda_0 (\Sigma_2 - \Sigma_1) v_2 \tilde{w} + \int_{\Omega} \langle [A_2 - A_1] \nabla v_2, \nabla \tilde{w} \rangle
\] (2.48)
and
\[
\int_{\Omega} \langle [A_1 - A_2] \nabla v_2, \nabla \tilde{w} \rangle + \langle [A_1 - A_2] \nabla \tilde{v}_2 \rangle + \lambda_0 (\Sigma_1 - \Sigma_2) v_2 \tilde{v}_2
\]
\[
= \int_{\Omega} g_2 \tilde{w} - (\tilde{g}_1 - \tilde{g}_2) v_2 + \lambda_0 (\Sigma_2 - \Sigma_1) v_2 \tilde{w}. \quad (2.49)
\]
Indeed, we have
\[
\text{div}(A_1 \nabla w) - \lambda_0 \Sigma_1 w = g_1 - g_2 - \lambda_0 (\Sigma_2 - \Sigma_1) v_2 + \text{div}([A_2 - A_1] \nabla v_2) \text{ in } \Omega.
\]
Multiplying this equation by \( \tilde{w} \) and integrating in \( \Omega \), we obtain (2.48). On the other hand, multiplying the equation of \( v_2 \) by \( \tilde{w} \) and integrating in \( \Omega \), we derive that
\[
\int_{\Omega} - \langle A_2 \nabla v_2, \nabla \tilde{w} \rangle - \lambda_0 \Sigma_2 v_2 \tilde{w} = \int_{\Omega} g_2 \tilde{w}. \quad (2.50)
\]
We have
\[
-A_2 \nabla \tilde{w} = (A_1 - A_2) \nabla \tilde{w} - A_1 \nabla \tilde{w} + (A_2 - A_1) \nabla \tilde{v}_2 + (A_1 - A_2) \nabla \tilde{v}_2 \text{ in } \Omega
\]
and, as in (2.12),
\[
\text{div}(A_1 \nabla \tilde{w}) - \text{div}([A_2 - A_1] \nabla \tilde{v}_2) = \lambda_0 \Sigma_1 \tilde{w} + \tilde{g}_1 - \tilde{g}_2 - \lambda_0 (\Sigma_2 - \Sigma_1) \tilde{v}_2 \text{ in } \Omega.
\]
Since
\[
(A_1 \nabla \tilde{w} - (A_2 - A_1) \nabla \tilde{v}_2) \cdot \nu = (A_1 \nabla \tilde{v}_1 - A_2 \nabla \tilde{v}_2) \cdot \nu = 0 \text{ on } \Gamma,
\]
it follows from (2.50) that
\[
\int_{\Omega} \langle [A_1 - A_2] \nabla v_2, \nabla \tilde{w} \rangle + \langle [A_1 - A_2] \nabla \tilde{v}_2 \rangle + \lambda_0 (\Sigma_1 - \Sigma_2) v_2 \tilde{v}_2
\]
\[
= \int_{\Omega} g_2 \tilde{w} - (\tilde{g}_1 - \tilde{g}_2) v_2 + \lambda_0 (\Sigma_2 - \Sigma_1) v_2 \tilde{w}, \quad (2.49)
\]
which is (2.49).

We are ready to introduce the notion of viscosity solutions of (2.1).

**Definition 2.** Let \((g_1, g_2) \in [L^2(\Omega)]^2\). A weak solution \((v_1, v_2) \in H(\Omega)\) of (2.39) is called a viscosity solution if (2.48) and (2.49) hold for any \((\tilde{v}_1, \tilde{v}_2) \in H(\Omega)\) and \((\tilde{g}_1, \tilde{g}_2) \in [L^2(\Omega)]^2\) satisfying (2.47).

We are now in the position to state and prove the crucial result of this section.
Proposition 5. Assume \([1.8]\). There exists \(\Lambda_0 > 1\) such that for \(\lambda_0 \geq \Lambda_0\) and for \(g = (g_1, g_2) \in [L^2(\Omega)]^2\), there exists a unique viscosity solution \(v = (v_1, v_2) \in \mathcal{H}(\Omega)\) of \((2.1)\). Moreover,

\[
\|v\|_{\mathcal{H}(\Omega)} \leq C\|g\|_{L^2(\Omega)},
\]

for some positive constant \(C\) independent of \(g\).

Proof. Let \(v_\delta = (v_1, v_2) \in [H^1(\Omega)]^2\) be the unique solution of \((2.29)\) and set \(w_\delta = v_1 - v_2\). By Proposition 4, we have, for large \(\Lambda\),

\[
\|v\|_{\mathcal{H}(\Omega)} \leq C\|g\|_{L^2(\Omega)},
\]

for some positive constant \(C\) independent of \(g\).

Proof. Let \(v_\delta = (v_1, v_2) \in [H^1(\Omega)]^2\) be the unique solution of \((2.29)\) and set \(w_\delta = v_1 - v_2\). By Proposition 4, we have, for large \(\Lambda_0\),

\[
\|v_\delta\|_{\mathcal{H}(\Omega)} \leq C\|g\|_{L^2(\Omega)}.
\]

Let \(v\) be the weak limit of \((v_\delta_n)\) in \(\mathcal{H}(\Omega)\) for some sequence \((\delta_n) \to 0\) such that \(w_\delta_n \rightharpoonup w = v_1 - v_2\) weakly in \(H^1(\Omega)\). Then \(v\) is a weak solution of \((2.1)\). We prove that \(v\) is a viscosity solution. Define \(g_1, g_2, \delta_2\) via \((2.44)\). Then, by \((2.48)\) and \((2.49)\),

\[
\int_\Omega (A_1 \nabla w_\delta, \nabla \tilde{w}) + \lambda_0 \Sigma_1 w_\delta \tilde{w} = -\int_\Omega (g_1, g_2) \tilde{w} + \int_\Omega \lambda_0 (\Sigma_2 - \Sigma_1) v_2 \tilde{w} + \int_\Omega \langle [A_2 - A_1] \nabla v_2, \nabla \tilde{w}\rangle
\]

and

\[
\int_\Omega \langle [A_1 - A_2] \nabla v_2, \nabla \tilde{w}\rangle + \langle [A_1 - A_2] \nabla v_2, \nabla \tilde{v}_2\rangle + \lambda_0 (\Sigma_2 - \Sigma_1) v_2 \tilde{v}_2
\]

\[
= \int_\Omega [g_2, \tilde{w} - \tilde{v}_2] g_2 + \lambda_0 (\Sigma_2 - \Sigma_1) v_2 \tilde{w}.
\]

By \((2.52)\), it follows from \((2.30)\) of Lemma 7 that

\[
\|v_\delta\|_{H^1(\Omega)} \leq \frac{C}{\sqrt{\delta}}\|g\|_{L^2(\Omega)}.
\]

By letting \(\delta \to 0\) in \((2.53)\) and \((2.54)\), and using \((2.55)\), one obtains \((2.48)\) and \((2.49)\) for \(v\). Hence \(v\) is a viscosity solution. The uniqueness of \(v\) follows from the definition of viscosity solutions as follows. Assume that \((v_1, v_2)\) and \((\tilde{v}_1, \tilde{v}_2)\) are two viscosity solutions. Set

\[(V_1, V_2) = (v_1 - \tilde{v}_1, v_2 - \tilde{v}_2) \text{ in } \Omega.
\]

Then \((V_1, V_2)\) is a viscosity solution corresponding to the pair of data \((0, 0)\). By choosing \((\hat{v}_1, \hat{v}_2) = (V_1, V_2)\), one obtains \((2.41)\) and \((2.5)\) where \((v_1, v_2) = (V_1, V_2)\) and \((g_1, g_2, h) = (0, 0, 0)\). It follows from Lemma 11 that \((2.0)\) and \((2.7)\) holds with \((g_1, g_2, h) = (0, 0, 0)\). This implies \((V_1, V_2) = (0, 0)\) by Lemma 11 for large \(\lambda_0\). The proof is complete.

\[
\Box
\]

Fix \(\lambda_0 > \Lambda_0\) where \(\Lambda_0\) is the constant in Proposition 5. Define

\[
T_1 : [L^2(\Omega)]^2 \to [L^2(\Omega)]^2
\]

\[
(f_1, f_2) \mapsto (u_1, u_2),
\]

where \((u_1, u_2) \in \mathcal{H}(\Omega)\) is the unique viscosity solution of \((2.2)\) with \((g_1, g_2) = (\Sigma_1 f_1, \Sigma_2 f_2)\). Using \((2.51)\) and applying Lemma 3 we derive that \(T_1\) is compact.
We are ready to give the

**Proof of Theorem 1** By the theory of compact operator see, e.g., [5], the spectrum of \( T_1 \) is discrete. It is clear that if \((u_1, u_2)\) is an eigenfunction pair of the ITE problem corresponding to the eigenvalue \( \lambda \) then \((u_1, u_2)\) is eigenfunction pair of \( T_1 \) corresponding to the eigenvalue \( \frac{1}{\lambda - \lambda_0} \). Hence, the spectrum of the ITE problem is discrete. □

**Remark 4.** The condition \( 0 \leq \alpha < 2 \) is necessary to guarantee the compactness of \( T_1 \).

**Remark 5.** Theorem 1 also holds if the condition \( \Sigma_1 - \Sigma_2 \geq 0 \) in (1.8) is replaced by \( \Sigma_1 - \Sigma_2 \geq -\hat{c} \) for some small positive constant \( \hat{c} \) since Lemma 5 holds in this case (see Remark 1).

### 3 On the case \( A_1 \geq A_2 \) and \( \Sigma_1 \) not greater than \( \Sigma_2 \) in a neighborhood of \( \Gamma \)

This section contains two sections. In the first one, we deal with the case \( A_1 \geq A_2 \) and \( \Sigma_1 \leq \Sigma_2 \) in a neighborhood of \( \Gamma \) and give the proof of Theorem 2. In the second one, we deal with the case \( A_1 \geq A_2 \) globally in \( \Omega \). The main result in this section is Proposition 8.

#### 3.1 On the case \( A_1 \geq A_2 \) and \( \Sigma_1 \leq \Sigma_2 \) in a neighborhood of \( \Gamma \)-Proof of Theorem 2

The section is devoted to the proof of Theorem 2. As mentioned in the introduction, we first establish the well-posedness of the following system, for some \( \lambda_0 > 0 \) (small),

\[
\begin{cases}
\text{div}(A_1 \nabla u_1) - i\lambda_0 \Sigma_1 u_1 = g_1 & \text{in } \Omega, \\
\text{div}(A_2 \nabla u_2) - i\lambda_0 \Sigma_2 u_2 = g_2 & \text{in } \Omega, \\
u_1 = u_2, & A_1 \nabla u_1 \cdot \nu = A_2 \nabla u_2 \cdot \nu & \text{on } \Gamma,
\end{cases}
\tag{3.1}
\]

for a given pair \((g_1, g_2)\) in an appropriate space. We then define the operator \( T_2(f_1, f_2) = (u_1, u_2) \) where \((u_1, u_2)\) is the unique solution of \((3.1)\) with \((g_1, g_2) = (\Sigma_1 f_1, \Sigma_2 f_2)\) \((3.2)\) and prove the compactness for \( T_2 \). System \((3.1)\) is slightly different from the one \((2.1)\) where the constant \( i \) appears in front of \( \lambda_0 \). This modification is necessary for the proof of the uniqueness. As in the proof of Theorem 1, similar a priori estimates hold for \((3.1)\) (Lemma 8). Nevertheless, the uniqueness of \((3.1)\) even for smooth solutions does not follow directly from the a priori estimates as in the proof of Theorem 1. Additional arguments are required and the condition on the largeness of \( \Sigma_2 \) is involved (Lemma 10). Beside this point, the proof of Theorem 2 is in the spirit of the one of Theorem 1.

We now process the proof of Theorem 2. We first establish a variant of Lemma 11.
Lemma 8. Let $\lambda_0 > 0$, $g = (g_1, g_2) \in [L^2(\Omega)]^2$, and $h \in H^{-1/2}(\Gamma)$. Assume that $v = (v_1, v_2) \in [H^1(\Omega)]^2$ satisfy the system

\[
\begin{cases}
\text{div}(A_1 \nabla v_1) - i\lambda_0 \Sigma_1 v_1 = g_1 & \text{in } \Omega, \\
\text{div}(A_2 \nabla v_2) - i\lambda_0 \Sigma_2 v_2 = g_2 & \text{in } \Omega, \\
v_1 = v_2, & A_1 \nabla v_1 \cdot \nu = A_2 \nabla v_2 \cdot \nu + h \text{ on } \Gamma.
\end{cases}
\] (3.3)

We have, with $w = v_2 - v_1$,

\[
\int_{\Omega} \langle A_1 \nabla w, \nabla w \rangle + i\lambda_0 \Sigma_1 |w|^2 = - \int_{\Omega} (g_1 - g_2) \bar{w} + \int_{\Omega} i\lambda_0 (\Sigma_2 - \Sigma_1) v_2 \bar{w} + \int_{\Omega} \langle [A_2 - A_1] \nabla v_2, \nabla w \rangle
\] (3.4)

and

\[
\int_{\Omega} \langle [A_1 - A_2] \nabla v_2, \nabla w \rangle + \langle [A_1 - A_2] \nabla v_2, \nabla v_2 \rangle + i\lambda_0 (\Sigma_2 - \Sigma_1) |v_2|^2
\]
\[
= \int_{\Omega} g_2 \bar{w} - (\bar{g}_1 - \bar{g}_2) v_2 + i\lambda_0 (\Sigma_1 + \Sigma_2) v_2 \bar{w} + \int_{\Gamma} \bar{h} v_2. \] (3.5)

As a consequence of (3.4) and (3.5), we obtain

\[
\lambda_0 \left| \int_{\Omega} (\Sigma_2 - \Sigma_1) |v_2|^2 \right| \leq 4 \mathcal{M}(v, g, h) + \lambda_0 \int_{\Omega} \Sigma_1 |w|^2 + 2\lambda_0 \int_{\Omega} \Sigma_1 |v_2||w| \] (3.6)

and

\[
\int_{\Omega} \langle A_1 \nabla w, \nabla w \rangle + \int_{\Omega} \langle [A_1 - A_2] \nabla v_2, \nabla v_2 \rangle
\]
\[
\leq 4 \mathcal{M}(v, g, h) + 2 \int_{\Omega} \left| \langle [A_2 - A_1] \nabla v_2, \nabla w \rangle \right| + 2\lambda_0 \int_{\Omega} \Sigma_2 |v_2||w|, \) (3.7)

where $\mathcal{M}(v, g, h)$ is defined in (2.8).

**Proof.** The proof is in the spirit of the one of Lemma 1; nevertheless, different coefficients appear in the conclusion due to the effect of the constant $i$ in the front of $\lambda_0$. For the convenience of the reader, we present the proof. From the definition of $w$, we derive that $w = 0$ on $\Gamma$ and

\[
\text{div}(A_1 \nabla w) - i\lambda_0 \Sigma_1 w = g_1 - g_2 - i\lambda_0 (\Sigma_2 - \Sigma_1) v_2 + \text{div}([A_2 - A_1] \nabla v_2) \text{ in } \Omega. \] (3.8)

Multiplying this equation by $\bar{w}$ and integrating on $\Omega$, we have

\[
\int_{\Omega} \langle A_1 \nabla w, \nabla w \rangle + i\lambda_0 \Sigma_1 |w|^2 = - \int_{\Omega} (g_1 - g_2) \bar{w} + \int_{\Omega} i\lambda_0 (\Sigma_2 - \Sigma_1) v_2 \bar{w} + \int_{\Omega} \langle [A_2 - A_1] \nabla v_2, \nabla w \rangle;
\] (3.9)

which is (3.4). Multiplying the equation of $v_2$ by $\bar{w}$ and integrating on $\Omega$, we obtain

\[
\int_{\Omega} -\langle A_2 \nabla v_2, \nabla w \rangle - i\lambda_0 \Sigma_2 v_2 \bar{w} = \int_{\Omega} g_2 \bar{w}. \] (3.10)
It is clear that
\[ -A_2 \nabla w = (A_1 - A_2) \nabla w - A_1 \nabla w + (A_2 - A_1) \nabla v_2 + (A_1 - A_2) \nabla v_2 \text{ in } \Omega, \quad (3.11) \]
and, by (3.8),
\[ \text{div}(A_1 \nabla w) - \text{div}((A_2 - A_1) \nabla v_2) = i\lambda_0 \Sigma_1 w + g_1 - g_2 - i\lambda_0 (\Sigma_2 - \Sigma_1)v_2 \text{ in } \Omega. \quad (3.12) \]
Since
\[ (A_1 \nabla w - (A_2 - A_1) \nabla v_2) \cdot \nu = (A_1 \nabla v_1 - A_2 \nabla v_2) \cdot \nu = h \text{ on } \Gamma, \]
it follows from (3.10), (3.11), and (3.12) that
\[ \hat{\Omega} \langle [A_1 - A_2] \nabla v_2, \nabla w \rangle + \langle [A_1 - A_2] \nabla v_2, \nabla v_2 \rangle + i\lambda_0 (\Sigma_2 - \Sigma_1)|v_2|^2 \]
\[ = \int_{\Omega} g_2 \overline{w} - (g_1 - \overline{g_2})v_2 + i\lambda_0 (\Sigma_1 + \Sigma_2)v_2 \overline{w} + \int_{\Gamma} h v_2; \quad (3.13) \]
which is (3.5). Subtracting (3.4) from (3.13) and considering the imaginary part yields
\[ \lambda_0 \left| \int_{\Omega} (\Sigma_2 - \Sigma_1)|v_2|^2 \right| \leq 4M(v, g, h) + \lambda_0 \int_{\Omega} |w|^2 + 2\lambda_0 \int_{\Omega} |w||v_2|; \quad (3.14) \]
which is (3.6). Adding (3.4) and (3.13) and considering the real part implies
\[ \int_{\Omega} \langle A_1 \nabla w, \nabla w \rangle + \int_{\Omega} \langle [A_1 - A_2] \nabla v_2, \nabla v_2 \rangle \]
\[ \leq 4M(v, g, h) + 2\lambda_0 \int_{\Omega} \Sigma_2 |v_2||w| + 2 \int_{\Omega} \langle [A_2 - A_1] \nabla v_2, \nabla w \rangle. \quad (3.15) \]
which is (3.7). The proof is complete. □

Here is a variant of Lemma 2.

**Lemma 9.** Let \( \lambda > 1, f \in L^2(\Omega), A \) be a matrix-valued function, and \( \Sigma \) be real function defined in \( \Omega \) such that
\[ \Lambda^{-1} \leq A \leq \Lambda \quad \text{and} \quad \Lambda^{-1} \leq \Sigma \leq \Lambda \text{ in } \Omega. \]
Let \( u \in H^1_{\text{loc}}(\Omega) \) be a solution to the equation \( \text{div}(A \nabla u) - i\lambda \Sigma u = f \) in \( \Omega \). For all \( s > 0 \), there exist two positive constants \( c_1 \) and \( c_2 \), depending only on \( \Lambda, s, \) and \( \Omega \), such that
\[ \|u\|_{H^1(\Omega; \Omega_s)} \leq c_1 \exp(-c_2 \sqrt{s})\|u\|_{L^2(\Omega_s)} + c_1 \|f\|_{L^2(\Omega)}. \]

**Proof.** The proof of Lemma 9 is similar to the one of Lemma 2. The details are omitted. □

The following result is a variant of Lemma 5 and plays an important role in the proof of the uniqueness in Theorem 2.

20
Lemma 10. Assume that for some $c, \tau > 0$ and for some $0 \leq \alpha < 2$,
\[
A_1 - A_2 \geq c d_0^\alpha I \text{ in } \Omega. \tag{3.16}
\]
Let $\Lambda_1 \geq 1$ be such that,
\[
\Lambda_1^{-1} I \leq A_j \leq \Lambda_1 I \text{ for } j = 1, 2, \text{ and } \Lambda_1^{-1} \leq \Sigma_1 \leq \Lambda_1 \text{ in } \Omega. \tag{3.17}
\]
For every $0 < \Lambda_2 < 1$ there exists $K_1 > 1$ depending only on $\Lambda_1, \Lambda_2, c, \alpha$, and $\tau$ such that if
\[
\Lambda_2 K \leq \Sigma_2 \leq K \text{ in } \Omega, \tag{3.18}
\]
for some $K > K_1$, then there exists $\lambda_0 > 0$, depending on $K$, such that if $g = (g_1, g_2) \in [L^2(\Omega)]^2$, $h \in H^{-1/2}(\Gamma)$, and $v = (v_1, v_2) \in H(\Omega)$ verify
\[
\begin{align*}
\text{div}(A_1 \nabla v_1) - i \lambda_0 \Sigma_1 v_1 &= g_1 \text{ in } \Omega, \\
\text{div}(A_2 \nabla v_2) - i \lambda_0 \Sigma_2 v_2 &= g_2 \text{ in } \Omega,
\end{align*} \tag{3.19}
\]
and $(3.6)$ and $(3.7)$, we have
\[
\|v\|^2_{H(\Omega)} \leq C \left( M(v, g, h) + \|g\|^2_{L^2(\Omega)} \right), \tag{3.20}
\]
for some positive constant $C$ independent of $v, g, and h$ where $M = M(v, g, h)$ is defined in $(2.8)$. In particular if $v \in [H^1(\Omega)]^2$ is a solution of $(3.3)$ then $(3.20)$ holds.

Proof. Fix $0 < \varepsilon < 1$ a small constant which is defined later and set $\lambda_0 = \varepsilon K^{-1/2}$. Note that $\lambda_0 K = \varepsilon K^{1/2}$ is large if $K$ is large; this fact is assumed from now. We derive from $(3.6)$ and Lemma 9 that, for large $K$,
\[
\int_\Omega \lambda_0 K |v_2|^2 \leq C \int_\Omega \lambda_0 |v_2||w| + \lambda_0 |w|^2 + CM(v, g, h).
\]
Here and in what follows in the proof, $C$ denotes a positive constant depending only on $\Lambda_1, \Lambda_2, \Omega, and \tau$; it is independent of $K$ and $\varepsilon$. This implies, for large $K$,
\[
\int_\Omega \lambda_0 K |v_2|^2 \leq C \int_\Omega \lambda_0 |w|^2 + CM(v, g, h). \tag{3.21}
\]
Using the fact, for $\gamma > 0$,
\[
\gamma \int_\Omega_{r} \langle [A_1 - A_2] \nabla w, \nabla w \rangle + \frac{1}{\gamma} \int_\Omega_{r} \langle [A_1 - A_2] \nabla v_2, \nabla v_2 \rangle \geq 2 \int_\Omega_{r} \left| \langle [A_1 - A_2] \nabla w, \nabla v_2 \rangle \right|
\]
and taking $\gamma > 1$ and close to 1, we derive from $(3.7)$ that
\[
\int_\Omega \langle A_1 \nabla w, \nabla w \rangle + \int_\Omega_{r} \langle [A_1 - A_2] \nabla v_2, \nabla v_2 \rangle \leq C \int_\Omega \lambda_0 K |v_2||w| + C \int_{\Omega_{r}} |\nabla v_2|^2 + CM(v, g, h). \tag{3.22}
\]
Since

\[ 2\lambda_0K|v_2||w| \leq \lambda_0K^{3/2}|v_2|^2 + \lambda_0K^{1/2}|w|^2, \]

it follows from (3.21) and (3.22) that

\[ \int_\Omega \langle A_1 \nabla w, \nabla w \rangle + \int_{\Omega} \langle [A_1 - A_2] \nabla v_2, \nabla v_2 \rangle \]

\[ \leq C \int_\Omega \lambda_0K^{1/2}|w|^2 + C \int_{\Omega} |\nabla v_2|^2 + CK^{1/2}M(v, g, h). \] (3.23)

Choosing \( \varepsilon \) small enough, one can absorb the first term in the RHS of (3.23) by the first term of the LHS of (3.23) (recall that \( \lambda_0K^{1/2} = \varepsilon \)) and obtains

\[ \int_\Omega \langle A_1 \nabla w, \nabla w \rangle + \int_{\Omega} \langle [A_1 - A_2] \nabla v_2, \nabla v_2 \rangle \leq C \int_\Omega |\nabla v_2|^2 + CK^{1/2}M(v, g, h). \] (3.24)

Similar to (2.26), we have

\[ \int_\Omega |v_2|^2 \leq C \int_\Omega d_1^2 |\nabla v_2|^2 + C\|g_2\|_{L^2(\Omega)}^2. \] (3.25)

Estimate (3.20) now follows from (3.24) and (3.25) by applying Lemma 9 with \( v = v_2, A = A_2, \Sigma = \Sigma_2/K, \) and \( \lambda = \lambda_0K = \varepsilon K^{1/2}. \)

The following result is a variant of Lemma 7 and its proof uses essentially Lemma 10.

**Lemma 11.** Assume that for some \( c, \tau > 0 \) and for some \( 0 \leq \alpha < 2, \)

\[ A_1 - A_2 \geq cd_1^2I \text{ in } \Omega. \] (3.26)

Let \( \Lambda_1 \geq 1 \) be such that

\[ \Lambda_1^{-1}I \leq A_j \leq \Lambda_1 I \text{ for } j = 1, 2, \quad \text{and} \quad \Lambda_1^{-1} \leq \Sigma_1 \leq \Lambda_1 \text{ in } \Omega. \] (3.27)

For every \( 0 < \Lambda_2 < 1 \) there exists \( K_1 > 1 \) depending only on \( \Lambda_1, \Lambda_2, c, \alpha, \) and \( \tau \) such that if

\[ \Lambda_2 K \leq \Sigma_2 \leq K \text{ in } \Omega, \] (3.28)

for some \( K > K_1, \) then there exists \( \lambda_0 > 0, \) depending on \( K, \) such that for \( \delta \in (0, 1), g = (g_1, g_2) \in [L^2(\Omega)]^2, \) there exists a unique solution \( v_\delta = (v_{1,\delta}, v_{2,\delta}) \in [H^1(\Omega)]^2 \) of the system

\[
\begin{cases}
\text{div} \left( (1 + \delta)A_1 \nabla v_{1,\delta} \right) - i\lambda_0 \Omega v_{1,\delta} = g_1 & \text{in } \Omega, \\
\text{div} \left( A_2 \nabla v_{2,\delta} \right) - i\lambda_0 \Omega_2 v_{2,\delta} = g_2 & \text{in } \Omega, \\
v_{1,\delta} = v_{2,\delta}, & (1 + \delta)A_1 \nabla v_{1,\delta} \cdot \nu = A_2 \nabla v_{2,\delta} \cdot \nu \text{ on } \Gamma.
\end{cases}
\] (3.29)

Moreover,

\[ \|v_\delta\|_{H^1(\Omega)}^2 \leq C \left( \|g\|_{L^2(\Omega)}^2 \|v_\delta\|_{L^2(\Omega)} + \|g\|_{L^2(\Omega)}^2 \right), \] (3.30)

for some positive constant \( C \) independent of \( g \) and \( \delta. \) Consequently,

\[ \|v_\delta\|_{H^1(\Omega)} \leq C \|g\|_{L^2(\Omega)}. \]
Proof. Applying Lemma 10 with $A_1 = (1 + \delta)A_1$ and $A_2 = A_2$, there exists $K_1 > 1$ depending only on $A_1, A_2, c, \alpha$, and \( \tau \) such that if

$$A_2 K \leq \Sigma_2 \leq K$$

in \( \Omega \),

(3.31)

for some $K > K_1$, there exists $\lambda_0 > 0$, depending on $K$, such that if $v_\delta = (v_1, v_2, \delta) \in [H^1(\Omega)]^2$ $(0 < \delta < 1)$ is a solution of (3.29) then

$$\|v_\delta\|_{H^1(\Omega)}^2 + \|v_1, v_2, \partial_\delta\|_{H^1(\Omega)}^2 + \|v_\delta\|_{L^2(\Omega)}^2 + \int_{\Omega^\tau} \left\langle \left[(1 + \delta)A_1 - A_2\right] \nabla v_\delta, \nabla v_\delta \right\rangle$$

$$\leq CM(v_\delta, g, 0) + C\|g\|_{L^2(\Omega)}^2,$$

(3.32)

for some positive constant $C$ independent of $g, v_\delta$, and $\delta$. This implies the uniqueness of $v_\delta$ and estimate (3.30). The existence of $v_\delta$ follows from Fredholm’s theory and can be proceeded as follows. Define

$$T_{1, \varepsilon} : [L^2(\Omega)]^2 \rightarrow [L^2(\Omega)]^2$$

$$
\begin{align*}
(g_1, g_2) &\mapsto (v_1, v_2),
\end{align*}
$$

(3.33)

where $(v_1, v_2) \in [H^1(\Omega)]^2$ is the unique solution of (2.1) where $(A_1, A_2)$ and $(\Sigma_1, \Sigma_2)$ are replaced by $((1 + \delta)A_1, A_2)$ and $(2\varepsilon^{-1}, \varepsilon^{-1})$ for small $\varepsilon$. $T_{1, \varepsilon}$ is well-defined since $T_1$ given in (2.56) is well-defined $(\alpha = 0$ in this case). The existence of $v_\delta$ follows from the uniqueness of $v_\delta$ by applying the Fredholm theory for the operator $I - T_{1, \varepsilon}B : [L^2(\Omega)]^2 \rightarrow [L^2(\Omega)]^2$ where $B : [L^2(\Omega)]^2 \rightarrow [L^2(\Omega)]^2$ is defined by

$$B(v_1, v_2) = (i\lambda_0 \Sigma_1 v_1 - 2\varepsilon^{-1}v_1, i\lambda_0 \Sigma_2 v_2 - \varepsilon^{-1}v_2).$$

Clearly $B$ is invertible if $\varepsilon$ is small enough. Note that $I - T_{1, \varepsilon}B$ is injective by the uniqueness of (3.29). By the Fredholm theory it is bijective. Hence for any $g \in [L^2(\Omega)]^2$ there exists $u \in [L^2(\Omega)]^2$ such that $u - T_{1, \varepsilon}B(u) = B^{-1}g$. Set $v = u - B^{-1}g$. Then $T_{1, \varepsilon}B(v + B^{-1}g) = v$. In other words, $v$ is a solution of (3.29). The proof is complete.

We are now in the position to give the definition of a weak solution of (2.1). Let $\tau > 0$ be such that (1.9) holds in $\Omega^\tau$. Define $\mathcal{H}(\Omega)$ as in (2.37) with the scalar product given in (2.38). The notion of weak solution is similar to the one in Definition 1 and is given in

Definition 3. Let $(g_1, g_2) \in [L^2(\Omega)]^2$. A pair of functions $(v_1, v_2) \in \mathcal{H}(\Omega)$ is called a weak solution of (2.1) if

$$\begin{cases}
\text{div}(A_1 \nabla v_1) - i\lambda_0 \Sigma_1 v_1 = g_1 & \text{in } \Omega, \\
\text{div}(A_2 \nabla v_2) - i\lambda_0 \Sigma_2 v_2 = g_2 & \text{in } \Omega, \\
(A_1 \nabla v_1 - A_2 \nabla v_2) \cdot \nu = 0 & \text{on } \Gamma.
\end{cases}$$

(3.34)

Using Lemmas 10 and 11, one can construct a weak solution $(v_1, v_2) \in \mathcal{H}(\Omega)$ of (3.31). More precisely, we have

Proposition 6. Let $\delta \in (0, 1)$, $g = (g_1, g_2) \in [L^2(\Omega)]^2$, and $v_\delta = (v_1, v_2, \delta) \in [H^1(\Omega)]^2$ be the unique solution of (3.29). Assume (3.26), (3.27), and (3.28). Let $K > K_1 > 1$ and $\lambda_0 > 0$ be as in Lemma 77. We have

$$\|v_\delta\|_{\mathcal{H}(\Omega)} \leq C\|g\|_{L^2(\Omega)},$$

(3.35)
for some positive constant $C$ independent of $\delta$ and $(g_1, g_2)$. As a consequence, there exists a weak solution $v = (v_1, v_2) \in \mathcal{H}(\Omega)$ of (3.1) such that

$$
\|v\|_{\mathcal{H}(\Omega)} \leq C\|g\|_{L^2(\Omega)}. \quad (3.36)
$$

**Proof.** The proof is similar to the one of Proposition 4. However, instead of using Lemmas 1 and 7, one applies Lemmas 8, 10, and 11. The details are left to the reader. \(\square\)

As in the proof of Theorem 1 in Section 2 we introduce the concept of viscosity solutions using a relaxation argument. Let $(\hat{v}_1, \hat{v}_2) \in \mathcal{H}(\Omega)$ and $(\hat{g}_1, \hat{g}_2) \in [L^2(\Omega)]^2$ be such that

$$
\begin{align*}
\text{div}(A_1 \nabla \hat{v}_1) - i\lambda_0 \Sigma_1 \hat{v}_1 &= \hat{g}_1 \quad \text{in } \Omega, \\
\text{div}(A_2 \nabla \hat{v}_2) - i\lambda_0 \Sigma_2 \hat{v}_2 &= \hat{g}_2 \quad \text{in } \Omega, \\
(A_1 \nabla \hat{v}_1 - A_2 \nabla \hat{v}_2) \cdot \nu &= \text{on } \Gamma.
\end{align*}
$$

(3.37)

Let $(g_1, g_2) \in [L^2(\Omega)]^2$ and $(v_1, v_2) \in [H^1(\Omega)]^2$ be a weak solution of the system

$$
\begin{align*}
\text{div}(A_1 \nabla v_1) - \lambda_0 \Sigma_1 v_1 &= g_1 \quad \text{in } \Omega, \\
\text{div}(A_2 \nabla v_2) - \lambda_0 \Sigma_2 v_2 &= g_2 \quad \text{in } \Omega, \\
v_1 &= v_2 \quad \text{on } \Gamma.
\end{align*}
$$

(3.38)

Set $w = v_1 - v_2$ and $\hat{w} = \hat{v}_1 - \hat{v}_2$. Invoking the same arguments used to derive (2.48) and (2.49), we have

$$
\int_\Omega (A_1 \nabla w, \nabla \hat{w}) + i\lambda_0 \Sigma_1 w \hat{w} = -\int_\Omega (g_1 - g_2) \hat{w} + \int_\Omega i\lambda_0 (\Sigma_2 - \Sigma_1) v_2 \hat{w} + \int_\Omega \langle [A_2 - A_1] \nabla v_2, \nabla \hat{w} \rangle
$$

(3.39)

and

$$
\int_\Omega \langle (A_1 - A_2) \nabla v_2, \nabla \hat{w} \rangle + \langle (A_1 - A_2) \nabla v_2, \nabla \hat{v}_2 \rangle + i\lambda_0 (\Sigma_2 - \Sigma_1) v_2 \hat{v}_2 = \int_\Omega g_2 \hat{w} - (\hat{g}_1 - \hat{g}_2)v_2 + i\lambda_0 (\Sigma_1 + \Sigma_2) v_2 \hat{w}. \quad (3.39)
$$

We are ready now to define the notion of viscosity solution of (3.1).

**Definition 4.** Let $(g_1, g_2) \in [L^2(\Omega)]^2$. A weak solution $(v_1, v_2) \in \mathcal{H}(\Omega)$ of (3.1) is called a viscosity solution if (3.38) and (3.39) hold for any $(v_1, \hat{v}_2) \in \mathcal{H}(\Omega)$ and $(\hat{g}_1, \hat{g}_2) \in [L^2(\Omega)]^2$ which satisfy (3.37).

Applying Lemma 8 and using Lemma 11, we can prove

**Proposition 7.** Assume (3.26), (3.27), and (3.28). Let $K_1 > 1$ and $\lambda_0 > 1$ be as in Lemma 11. For $g = (g_1, g_2) \in [L^2(\Omega)]^2$, there exists a unique viscosity solution $v = (v_1, v_2) \in \mathcal{H}(\Omega)$ of (3.1). Moreover,

$$
\|v\|_{\mathcal{H}(\Omega)} \leq C\|g\|_{L^2(\Omega)}. \quad (3.40)
$$

for some positive constant $C$ independent of $(g_1, g_2)$.
Proof. The proof of Proposition 7 is similar to the one of Proposition 5. The details are left to the reader. □

Fix $\lambda_0$ in Proposition 7. Define

$$T_2 : [L^2(\Omega)]^2 \rightarrow [L^2(\Omega)]^2$$

$$(f_1, f_2) \mapsto (u_1, u_2),$$

where $(u_1, u_2) \in H(\Omega)$ is the unique viscosity solution of (3.1). The compactness of $T_2$ is a consequence of Lemma 3.

We are ready to give the Proof of Theorem 2. By the theory of compact operator, the spectrum of $T_2$ is discrete. It is clear that if $(u_1, u_2)$ is an eigenfunction pair of the ITE problem corresponding to $\lambda$ then $(u_1, u_2)$ is eigenfunction pair of $T$ corresponding to the eigenvalue $\frac{\lambda}{\lambda - i \lambda_0}$. Hence, the spectrum of the ITE problem is discrete. □

Remark 6. As in the proof of Theorem 1, the condition $0 \leq \alpha < 2$ in Theorem 2 is necessary to guarantee the compactness of $T_2$.

3.2 On the case $A_1 \geq A_2$ globally in $\Omega$

The main result of this section is Proposition 8. Assume that

$$A_1 - A_2 \geq cd_1^2 I \text{ in } \Omega \text{ and } \int_{\Omega} (\Sigma_1 - \Sigma_2) \neq 0,$$

for some $0 \leq \alpha < 2$ and for some $c > 0$. The spectrum of (3.3) - (3.4) is discrete.

In Proposition 8 the condition $A_1 - A_2 \geq cd_1^2 I$ is required in the whole domain $\Omega$ not in a neighborhood of $\Gamma$; this is different from the context of Theorem 2. Applying Proposition 8 with $\alpha = 0$ and $(A_1, \Sigma_1) = (I, 1)$ in $\Omega$, one rediscovers the discreteness result given in [4, Theorem 4.4].

The proof of Proposition 8 is based on the following lemma

Lemma 12. Assume that for some $c > 0$ and for some $0 \leq \alpha < 2$,

$$A_1 - A_2 \geq cd_1^2 I \text{ in } \Omega \text{ and } \int_{\Omega} (\Sigma_1 - \Sigma_2) \neq 0.$$

There exists $\Lambda_0 > 0$ such that if $0 < \lambda_0 < \Lambda_0$, and $v = (v_1, v_2) \in [H^1_{\text{loc}}(\Omega)]^2$ verifies

$$\begin{cases}
\text{div}(A_1 \nabla v_1) - i\lambda_0 \Sigma_1 v_1 = g_1 \text{ in } \Omega, \\
\text{div}(A_2 \nabla v_2) - i\lambda_0 \Sigma_2 v_2 = g_2 \text{ in } \Omega,
\end{cases}$$

and (3.6) and (3.7) for some $g = (g_1, g_2) \in [L^2(\Omega)]^2$ and $h \in H^{-1/2}(\Gamma)$, then

$$\|v\|_{H^1(\Omega)}^2 \leq C\mathcal{M}(v, g, h),$$

for some positive constant $C$ independent of $v$, $g$, and $h$ where $\mathcal{M}(v, g, h)$ is defined in (2.8),
Proof. The proof is in the spirit of the one of Lemma 11. From (3.6), we have, for all $0 < \gamma < 1$,
\[
\lambda_0 \left| \int_{\Omega} (\Sigma_2 - \Sigma_1) |v_2|^2 \right| \leq \int_{\Omega} \gamma \lambda_0 |v_2|^2 + C \gamma^{-1} \lambda_0 |w|^2 + M(v, g, h). \tag{3.45}
\]
Here and in what follows in the proof, $C$ denotes a positive constant depending only on $\Lambda_1, \Lambda_2$, $\Omega$; it is independent of $\lambda_0$ and $\gamma$. As in (3.23), we deduce from (3.7) that
\[
\int_{\Omega} (A_1 \nabla w, \nabla w) + \int_{\Omega} \langle [A_1 - A_2] \nabla v_2, \nabla v_2 \rangle \leq \int_{\Omega} \gamma \lambda_0 |v_2|^2 + C \gamma^{-1} \lambda_0 |w|^2 + C M(v, g, h). \tag{3.46}
\]
Note that $\int_{\Omega} |w|^2 \leq C \int_{\Omega} |\nabla w|^2$ by Poincaré’s inequality, and
\[
\int_{\Omega} |v_2|^2 \leq C \left| \int_{\Omega} (\Sigma_2 - \Sigma_1) |v_2|^2 \right| + C \int_{\Omega} d_1^2 |\nabla v_2|^2, \tag{3.47}
\]
by Lemma 4. Fixing $\gamma$ small enough in (3.45) and (3.46), adding (3.45) and (3.46), and using (3.47), we derive that if $\lambda_0$ is small enough,
\[
\int_{\Omega} (A_1 \nabla w, \nabla w) + \int_{\Omega} \langle [A_1 - A_2] \nabla v_2, \nabla v_2 \rangle + \lambda_0 \int_{\Omega} |v_2|^2 \leq C M(v, g, h). \tag{3.48}
\]
The conclusion now follows from (3.48) by noting that $w = v_1 - v_2$. \qed

We are ready to present

Proof of Proposition 8. The proof follows as in the one of Theorem 2 by considering the transformation
\[
T_2 : [L^2(\Omega)]^2 \rightarrow [L^2(\Omega)]^2
\]
\[
(f_1, f_2) \mapsto (u_1, u_2), \tag{3.49}
\]
where $(u_1, u_2) \in \mathcal{H}(\Omega)$ is the unique viscosity solution of (3.1). The details are left to the reader. \qed

4 On the case $A_1 = A_2$ in a neighborhood of $\Gamma$ - Proof of Theorem 3

The condition $A_1 - A_2 \geq cd_1^\alpha I$ in a neighborhood of $\Gamma$ with $0 \leq \alpha < 2$ plays a crucial role in establishing the compactness of $T$, more precisely $T_1$ and $T_2$ (see Remarks 4 and 5). To be able to deal with case $A_1 = A_2$ in a neighborhood of $\Gamma$, we make some modifications on $T$. The idea is to take into account the fact that $u_1 - u_2 \in H_0^1(\Omega)$ which is more regular $u_1$ and $u_2$ which are in general not in $H^1(\Omega)$. A modification on $T$ was also used in the work of Sylvester’s [35]. Nevertheless, the modification in [35] requires the condition $A_1 = A_2$ globally in $\Omega$ and does not fit for our situation.
The motivation for reformulating the problem is as follows. Let \( \lambda \in \mathbb{C} \) be an eigenvalue of the ITE problem and let \((u_1, u_2) \in [H^1(\Omega)]^2\) be a corresponding pair of eigenfunctions. Then

\[
\begin{aligned}
\text{div}(A_1 \nabla u_1) - \lambda \Sigma_1 u_1 &= 0 \quad \text{in } \Omega, \\
\text{div}(A_2 \nabla u_2) - \lambda \Sigma_2 u_2 &= 0 \quad \text{in } \Omega, \\
u_1 = u_2, \quad A_1 \nabla u_1 \cdot \nu &= A_2 \nabla u_2 \cdot \nu \quad \text{on } \Gamma.
\end{aligned}
\]  

(4.1)

Fix \( \lambda_0 \neq 0 \) and set \( w = u_1 - u_2 \) in \( \Omega \). From (4.1), we have

\[
\text{div}(A_1 \nabla w) - \lambda_0 \Sigma_1 w = (\lambda - \lambda_0)\Sigma_1 u_1 - \text{div}(A_1 \nabla u_2) + \lambda_0 \Sigma_1 u_2
\]

\[
= (\lambda - \lambda_0)\Sigma_1 (u_1 - u_2) - \lambda \Sigma_2 u_2 - \text{div}([A_1 - A_2] \nabla u_2) + \lambda \Sigma_1 u_2 \quad \text{in } \Omega.
\]

Define

\[
\hat{w} = w + \lambda u_2/\lambda_0 \quad \text{in } \Omega.
\]

Then

\[
\text{div}(A_1 \nabla \hat{w}) - \lambda_0 \Sigma_1 \hat{w} = (\lambda - \lambda_0)\Sigma_1 (u_1 - u_2) + \frac{\lambda}{\lambda_0} (\lambda - \lambda_0) \Sigma_2 u_2 + \frac{\lambda - \lambda_0}{\lambda_0} \text{div}([A_1 - A_2] \nabla u_2) \quad \text{in } \Omega.
\]

(4.2)

It is clear that

\[
\text{div}(A_2 \nabla u_2) - \lambda_0 \Sigma_2 u_2 = (\lambda - \lambda_0)\Sigma_2 u_2 \quad \text{in } \Omega.
\]

(4.3)

Set

\[
U_1 = \hat{w} = (u_1 - u_2) + \lambda u_2/\lambda_0 \quad \text{and} \quad U_2 = \lambda u_2/\lambda_0 \quad \text{in } \Omega.
\]

It follows from (4.2) and (4.3) that, in \( \Omega \),

\[
\begin{aligned}
\text{div}(A_1 \nabla U_1) - \lambda_0 \Sigma_1 U_1 &= (\lambda - \lambda_0)\Sigma_1 (U_1 - U_2) + (\lambda - \lambda_0) \Sigma_2 U_2 + \frac{\lambda - \lambda_0}{\lambda} \text{div}([A_1 - A_2] \nabla U_2) \\
&= (\lambda - \lambda_0)\Sigma_1 (U_1 - U_2) + (\lambda - \lambda_0) \Sigma_2 U_2 + \frac{\lambda - \lambda_0}{\lambda} \text{div}([A_1 - A_2] \nabla U_2)
\end{aligned}
\]  

(4.4)

and

\[
\begin{aligned}
\text{div}(A_2 \nabla U_2) - \lambda_0 \Sigma_2 U_2 &= (\lambda - \lambda_0)\Sigma_2 U_2.
\end{aligned}
\]  

(4.5)

Since \( A_1 = A_2 \) near \( \Gamma \), we also have

\[
U_1 = U_2, \quad A_1 \nabla U_1 \cdot \nu = A_2 \nabla U_2 \cdot \nu \quad \text{on } \Gamma.
\]

(4.6)

One can rewrite (4.4), (4.5), and (4.6) under the form

\[
T_{3,\lambda}(U_1, U_2) = \frac{1}{\lambda - \lambda_0}(U_1, U_2).
\]

Here

\[
T_{3,\lambda}(f_1, f_2) := (v_1, v_2),
\]

\footnote{The goal is to eliminate \( \Sigma_1 u_2 \) from the equation of \( \hat{w} \).}
where \((v_1, v_2)\) satisfies the system
\[
\begin{aligned}
\text{div}(A_1 \nabla v_1) - \lambda_0 \Sigma_1 v_1 &= \Sigma_1 (f_1 - f_2) + \Sigma_2 f_2 + \lambda^{-1} \text{div} ([A_1 - A_2] \nabla f_2) & \text{in } \Omega, \\
\text{div}(A_2 \nabla v_2) - \lambda_0 \Sigma_2 v_2 &= \Sigma_2 f_2 & \text{in } \Omega, \\
v_1 = v_2, & A_1 \nabla v_1 \cdot \nu = A_2 \nabla v_2 \cdot \nu & \text{on } \Gamma.
\end{aligned}
\]

The discreteness of \(T\) can now be derived from the discreteness of \(T_{3,\lambda}\). To this end, we first establish the well-posedness of the system
\[
\begin{aligned}
\text{div}(A_1 \nabla v_1) - \lambda_0 \Sigma_1 v_1 &= g_1 + \text{div}(G_1) & \text{in } \Omega, \\
\text{div}(A_2 \nabla v_2) - \lambda_0 \Sigma_2 v_2 &= g_2 & \text{in } \Omega, \\
v_1 = v_2, & A_1 \nabla v_1 \cdot \nu = A_2 \nabla v_2 \cdot \nu & \text{on } \Gamma,
\end{aligned}
\]
for appropriate functions \(g_1, g_2\) and vector field \(G_1\). We follow the strategy used in the proof of Theorem \[\text{Lemma 13.}\] We first introduce some notations. Given \(s \in \mathbb{R}\), we denote
\[
L^2(\Omega, d_1^s) := \left\{ \psi \in L^2_{\text{loc}}(\Omega); \|\psi\|_{L^2(\Omega, d_1^s)} < +\infty \right\},
\]
where
\[
\|\psi\|_{L^2(\Omega, d_1^s)}^2 := \int_\Omega d_1^s |\psi|^2.
\]
We also define, for a given \(0 \leq \beta_1 < 2\),
\[
\mathcal{H}_0(\Omega) = \left\{ (\psi_1, \psi_2) \in L^2_{\text{loc}}(\Omega); \psi_1, \psi_2 \in L^2(\Omega, d_1^{\beta_1}) \text{ and } \psi_1 - \psi_2 \in L^2(\Omega, d_1^{-\beta_1}) \right\}
\]
and
\[
\| (\psi_1, \psi_2) \|_{\mathcal{H}_0(\Omega)} := \| (\psi_1, \psi_2) \|_{L^2(\Omega)} + \| \psi_1 - \psi_2 \|_{L^2(\Omega, d_1^{-\beta_1})}.
\]

Here is a variant of Lemma \[\text{Lemma 13.}\]

**Lemma 13.** Let \(0 \leq \beta_1 < 2\), \(\lambda_0 > 0\), \(g = (g_1, g_2) \in \mathcal{H}_0(\Omega)\), and \(G_1 \in [L^2(\Omega)]^d\) with \(\text{supp } G_1 \subset \subset \Omega\). Assume that \(v = (v_1, v_2) \in [H^1(\Omega)]^2\) satisfy the system
\[
\begin{aligned}
\text{div}(A_1 \nabla v_1) - \lambda_0 \Sigma_1 v_1 &= g_1 + \text{div}(G_1) & \text{in } \Omega, \\
\text{div}(A_2 \nabla v_2) - \lambda_0 \Sigma_2 v_2 &= g_2 & \text{in } \Omega, \\
v_1 = v_2, & A_1 \nabla v_1 \cdot \nu = A_2 \nabla v_2 \cdot \nu & \text{on } \Gamma.
\end{aligned}
\]

We have, with \(w = v_1 - v_2\),
\[
\int_\Omega \langle A_1 \nabla w, \nabla \bar{w} \rangle + \lambda_0 \Sigma_1 |w|^2
\]
\[
= - \int_\Omega (g_1 - g_2) \bar{w} + \int_\Omega \langle G_1, \nabla \bar{w} \rangle + \int_\Omega \lambda_0 (\Sigma_2 - \Sigma_1) v_2 \bar{w} + \int_\Omega \langle [A_2 - A_1] \nabla v_2, \nabla \bar{w} \rangle
\]
and
\[
\int_\Omega \langle [A_1 - A_2] \nabla v_2, \nabla w \rangle + \langle [A_1 - A_2] \nabla v_2, \nabla v_2 \rangle + \lambda_0 (\Sigma_1 - \Sigma_2)|v_2|^2 = \int_\Omega g_2 \bar{w} - (\bar{g}_1 - \bar{g}_2)v_2 + \lambda_0 (\Sigma_2 - \Sigma_1)v_2 \bar{w} + \int_\Omega \langle \nabla v_2, G_1 \rangle. \tag{4.11}
\]

As a consequence of (4.10) and (4.11), we obtain
\[
\int_\Omega \langle A_1 \nabla w, \nabla w \rangle + \lambda_0 \Sigma_1 |w|^2 \leq 4\mathcal{N}(v, g, G_1) + \int_\Omega |\langle [A_2 - A_1] \nabla v_2, \nabla w \rangle| + \lambda_0 |(\Sigma_2 - \Sigma_1)v_2 \bar{w}|, \tag{4.12}
\]
\[
\int_\Omega \langle [A_1 - A_2] \nabla v_2, \nabla v_2 \rangle + \lambda_0 (\Sigma_1 - \Sigma_2)|v_2|^2 \leq 4\mathcal{N}(v, g, G_1) + \int_\Omega |\langle [A_2 - A_1] \nabla v_2, \nabla w \rangle| + \lambda_0 |(\Sigma_2 - \Sigma_1)v_2 \bar{w}|, \tag{4.13}
\]
where
\[
\mathcal{N}(v, g, G_1) = \int_\Omega |g||w| + |g_1 - g_2||v| + |G_1||\nabla v_2|. \tag{4.14}
\]

**Proof.** The proof of Lemma 13 follows closely to the one of Lemma 1 by noting that one still has (2.4) and (2.5) in which \(g_1\) is replaced by \(g_1 + \text{div}(G_1)\) and \(h = 0\), i.e., (4.10) and (4.11) holds. The details are left to the reader. \(\square\)

We next establish a variant of Lemma 5.

**Lemma 14.** Let \(0 \leq \beta < \beta_1 < 2\) and assume that, for some \(c, \tau > 0\),
\[
A_1 \geq A_2 \quad \text{and} \quad \Sigma_1 - \Sigma_2 \geq cd_1^3 \text{ in } \Omega_\tau. \tag{4.15}
\]
There exists \(\Lambda_0 > 1\) such that if \(\lambda_0 \geq \Lambda_0\), and \(v = (v_1, v_2) \in [H^1_0(\Omega)]^2\) satisfies the system
\[
\begin{cases}
\text{div}(A_1 \nabla v_1) - \lambda_0 \Sigma_1 v_1 = g_1 + \text{div}(G_1) & \text{in } \Omega, \\
\text{div}(A_2 \nabla v_2) - \lambda_0 \Sigma_2 v_2 = g_2 & \text{in } \Omega,
\end{cases} \tag{4.16}
\]
and (4.12) and (4.13) for some \(g = (g_1, g_2) \in H_0(\Omega)\), and \(G_1 \in [L^2(\Omega)]^d\) with \(\text{supp} G_1 \subset \subset \Omega\), then, for \(w = v_1 - v_2\) in \(\Omega\),
\[
\int_{\Omega_\tau} \langle [A_1 - A_2] \nabla v_2, \nabla v_2 \rangle + \int_{\Omega_\tau} (\Sigma_1 - \Sigma_2)|v_2|^2 + \int_{\Omega_\tau} (|\nabla v_2|^2 + |v_2|^2) + \int_{\Omega} |\nabla w|^2 \leq C \left( \mathcal{N}(v, g, G_1) + \|g_2\|_{L^2(\Omega_\tau)}^2 \right), \tag{4.17}
\]
for some positive constant \(C\) independent of \(v, g,\) and \(G_1\) where \(\mathcal{N}(v, g, G_1)\) is defined by (4.14). In particular, if \(v \in [H^1(\Omega)]^2\) is a solution of (4.9) then (4.17) holds.
Proof. The proof is in the spirit of the one of Lemma 5 and even simpler. Adding (4.12) and (4.13), using the fact, for $\gamma > 0$,

$$
\gamma \int_{\Omega_r} \langle [A_1 - A_2] \nabla w, \nabla w \rangle + \frac{1}{\gamma} \int_{\Omega_r} \langle [A_1 - A_2] \nabla v_2, \nabla v_2 \rangle \geq 2 \int_{\Omega_r} \langle [A_1 - A_2] \nabla w, \nabla v_2 \rangle
$$

and

$$
\gamma \int_{\Omega_r} (\Sigma_1 - \Sigma_2) |w|^2 + \frac{1}{\gamma} \int_{\Omega_r} (\Sigma_1 - \Sigma_2) |v_2|^2 \geq 2 \int_{\Omega_r} (\Sigma_1 - \Sigma_2) |w v_2|,
$$

and taking $\gamma > 1$ and close to 1 in the previous two inequalities, we obtain

$$
\int_{\Omega} |\nabla w|^2 + \lambda_0 |w|^2 + \lambda_0 (\Sigma_1 - \Sigma_2) |v_2|^2 + \int_{\Omega_r} \langle [A_1 - A_2] \nabla v_2, \nabla v_2 \rangle
$$

$$
\leq CN(v, g, G_1) + C \int_{\Omega \setminus \Omega_r} |\nabla v_2|^2 + \lambda_0 |v_2w|, \quad (4.18)
$$

for some positive constant $C$ independent of $v$, $g$, $G_1$, and $\lambda_0$. Since, for $\gamma > 0$,

$$
2 |v_2w| \leq \gamma |v_2|^2 + \gamma^{-1} |w|^2, \quad (4.19)
$$

by applying Lemma 2 to $v_2 \in \Omega \setminus \Omega_{r/2}$ after fixing $\gamma > 1$ large enough in (4.19), we have, for large $\lambda_0$,

$$
\int_{\Omega_r} \langle [A_1 - A_2] \nabla v_2, \nabla v_2 \rangle + \int_{\Omega_r} (\Sigma_1 - \Sigma_2) |v_2|^2 + \int_{\Omega \setminus \Omega_r} (|\nabla v_2|^2 + |v_2|^2) + \int_{\Omega} |\nabla w|^2
$$

$$
\leq C \left( N(v, g, G_1) + \|g_2\|_{L^2(\Omega \setminus \Omega_{r/2})} \right),
$$

for some positive constant $C$ independent of $v$, $g$, and $G_1$; which is the conclusion. \hfill \Box

We now introduce functional spaces. In what follows in this section, we assume that (1.12) holds in $\Omega_\tau$ for some $\tau > 0$ and set $\beta_1 = (2 + \beta)/2$. Define

$$
\mathcal{H}_1(\Omega) := \left\{ (v_1, v_2) \in [H^1_{\text{loc}}(\Omega)]^2; \quad v_1, v_2 \in H^1_0(\Omega), \; \|(v_1, v_2)\|_{\mathcal{H}_1(\Omega)} < +\infty \right\}, \quad (4.20)
$$

where the norm is generated from the following scalar product, with $u = (u_1, u_2)$ and $v = (v_1, v_2)$,

$$
\langle u, v \rangle_{\mathcal{H}_1(\Omega)} = \int_{\Omega_r} \nabla (u_1 - v_2) \nabla (\bar{v}_1 - \bar{v}_2) + \int_{\Omega} d_0^{\beta_1 + 2} \nabla u_2 \nabla \bar{v}_2 + \int_{\Omega_\tau} u \bar{v} + \int_{\Omega_r} (\Sigma_1 - \Sigma_2) u \bar{v}. \quad (4.21)
$$

One can verify that $\mathcal{H}_1(\Omega)$ is a Hilbert space.

We are ready to give the

Definition 5. Let $(g_1, g_2) \in \mathcal{H}_0(\Omega)$ and $G_1 \in [L^2(\Omega)]^d$ such that $\text{supp} G_1 \subset \Omega \setminus \Omega_{\tau}$. A pair of functions $(v_1, v_2) \in \mathcal{H}_1(\Omega)$ is called a weak solution of (1.3) if

$$
\begin{cases}
\text{div}(A_1 \nabla v_1) - \lambda_0 \Sigma_1 v_1 = g_1 + \text{div}(G_1) \quad &\text{in } \Omega, \\
\text{div}(A_2 \nabla v_2) - \lambda_0 \Sigma_2 v_2 = g_2 \quad &\text{in } \Omega, \\
(A_1 \nabla v_1 - A_2 \nabla v_2) \cdot \nu = 0 \quad &\text{on } \Gamma.
\end{cases} \quad (4.22)
$$
Remark 7. Since \( \text{div}(A_1 \nabla v_1 - A_2 \nabla v_2) \in L^2(\Omega \setminus \Omega_\tau) \) and \( A_1 \nabla v_1 - A_2 \nabla v_2 = A_1 \nabla (v_1 - v_2) + (A_1 - A_2) \nabla v_2 \in L^2(\Omega) \) by (4.21). The last identity on \( \Gamma \) in (4.22) makes sense.

Applying Lemma\[14\] we can prove

**Lemma 15.** Let \( g = (g_1, g_2) \in \mathcal{H}_0(\Omega) \) and \( G_1 \in L^2(\Omega) \) with \( \text{supp} G_1 \subset \Omega \setminus \Omega_\tau \). There exists \( \Lambda_0 > 1 \) such that if \( \Lambda_0 \geq \Lambda_0 \), then for all \( 0 < \delta < 1 \), there exists a unique solution \( v_\delta = (v_{1,\delta}, v_{2,\delta}) \in [H^1(\Omega)]^2 \) satisfy the system

\[
\begin{align*}
\text{div} \left( (1 + \delta) A_1 \nabla v_{1,\delta} \right) - \lambda_0 (1 + \delta) \Sigma_1 v_{1,\delta} &= g_1 + \text{div}(G_1) \quad \text{in} \ \Omega, \\
\text{div} \left( A_2 \nabla v_{2,\delta} \right) - \lambda_0 \Sigma_2 v_{2,\delta} &= g_2 \quad \text{in} \ \Omega, \\
v_{1,\delta} &= v_{2,\delta}, \quad (1 + \delta) A_1 \nabla v_{1,\delta} \cdot \nu = A_2 \nabla v_{2,\delta} \cdot \nu \quad \text{on} \ \Gamma.
\end{align*}
\]

(4.23)

Moreover,

\[
\int_\Omega \delta (|\nabla v_{2,\delta}|^2 + |v_{2,\delta}|^2) + \|v_\delta\|^2_{\tilde{H}_1(\Omega)} \leq C \left( \|g\|^2_{\mathcal{H}_0(\Omega)} + \|G_1\|^2_{L^2(\Omega)} \right),
\]

(4.24)

for some positive constant \( C \) independent of \( \delta, g, \) and \( G_1 \). As a consequence, there exists a weak solution \( v = (v_1, v_2) \in \mathcal{H}_1(\Omega) \) of (4.23) such that

\[
\|v\|^2_{\tilde{H}_1(\Omega)} \leq C \left( \|g\|^2_{\mathcal{H}_0(\Omega)} + \|G_1\|^2_{L^2(\Omega)} \right).
\]

(4.25)

**Proof.** Assume that \( v_\delta = (v_{1,\delta}, v_{2,\delta}) \in [H^1(\Omega)]^2 \) is a solution of (4.23). Set \( w_\delta = v_{1,\delta} - v_{2,\delta} \). Applying Lemma\[14\] to \((A_1, A_2) = ((1 + \delta) A_1, A_2)\) and \((\Sigma_1, \Sigma_2) = ((1 + \delta) \Sigma_1, \Sigma_2)\), one obtains, for large \( \lambda_0 \),

\[
\begin{align*}
\int_{\Omega_\tau} \langle (1 + \delta) A_1 - A_2 \rangle \nabla v_{2,\delta} , \nabla v_{2,\delta} \rangle + \int_{\Omega_\tau} (1 + \delta) \Sigma_1 - \Sigma_2) |v_{2,\delta}|^2 \\
+ \int_{\Omega \setminus \Omega_\tau} (|\nabla v_{2,\delta}|^2 + |v_{2,\delta}|^2) + \int_\Omega |\nabla w_\delta|^2 \leq C \left( \mathcal{N}(v_\delta, g, G_1) + \|g_2\|^2_{L^2(\Omega \setminus \Omega_\tau)} \right).
\end{align*}
\]

(4.26)

We have

\[
\begin{align*}
\int_\Omega |g| |w_\delta| &\leq \left( \int_\Omega d_\Gamma^{\beta_1} |g|^2 \right)^{1/2} \left( \int_\Omega d_\Gamma^{-\beta_1} |w_\delta|^2 \right)^{1/2}, \\
\int_\Omega |g_1 - g_2| |v_\delta| &\leq \left( \int_\Omega d_\Gamma^{-\beta_1} |g_1 - g_2|^2 \right)^{1/2} \left( \int_\Omega d_\Gamma^{\beta_1} |v_\delta|^2 \right)^{1/2}, \\
\int_\Omega |G_1| |\nabla v_{2,\delta}| &\leq \left( \int_{\Omega \setminus \Omega_\tau} |G_1|^2 \right)^{1/2} \left( \int_{\Omega \setminus \Omega_\tau} |\nabla v_{2,\delta}|^2 \right)^{1/2}.
\end{align*}
\]

Noting that \( w_\delta \in H^1_0(\Omega) \) and hence, by the Hardy inequality, we obtain

\[
\int_\Omega d_\Gamma^{-2} |w_\delta|^2 \leq C_\Omega \int_\Omega |\nabla w_\delta|^2.
\]

Since \( 0 \leq \beta < \beta_1 < 2 \), it follows from the definition of \( \mathcal{N} \) that

\[
\mathcal{N}(v_\delta, g, G_1) \leq C \left( \|g\|^2_{\mathcal{H}_0(\Omega)} + \|G_1\|^2_{L^2(\Omega \setminus \Omega_\tau)} \right) \|v_\delta\|_{\tilde{H}_1(\Omega)}.
\]

(4.27)

31
By Lemma 6, we have
\[ \int_\Omega d_\tau^{\beta_2 + 2} |\nabla v_2|^2 \leq C \int_\Omega d_\tau^{\beta_2} |v_2|^2 + d_\tau^{\beta_2} |g_2|^2. \] (4.28)

Since \( \beta_1 > \beta \), we derive from (4.28) that
\[
\int_\Omega ((1 + \delta)A_1 - A_2) \nabla v_{2,\delta}, \nabla v_{2,\delta} + \int_\Omega ((1 + \delta)\Sigma_1 - \Sigma_2) |v_{2,\delta}|^2 + \int_{\Omega \setminus \Omega_r} (|\nabla v_{2,\delta}|^2 + |v_{2,\delta}|^2)
\]
\[+ \int_\Omega |\nabla w|^2 + \|g\|_{H_0(\Omega)}^2 \geq C \left( \|v_{\delta}\|_{H^1(\Omega)}^2 + \int_\Omega \delta (|\nabla v_{2,\delta}|^2 + |v_{2,\delta}|^2) \right). \] (4.29)

Combining (4.26), (4.27), and (4.29) yields (4.24). This in turn implies the uniqueness of \( v_\delta \).

The existence of \( v_\delta \) can be established via an approximation argument by first assuming that \( g \in [L^2(\Omega)]^2 \) and \( G_1 \in [L^2(\Omega)]^d \) and then considering the general case; the existence in first case follows from Proposition 4. The details of this fact are left to the reader. The existence and estimate of a weak solution \( v \in H_1(\Omega) \) of (4.28) follows directly from the existence and the estimate of \( v_\delta \). The details are omitted.

As in the proof of Theorem 1 in Section 2, we introduce the concept of viscosity solutions.

**Definition 6.** Let \((g_1, g_2) \in H_0(\Omega)\) and \(G_1 \in L^2(\Omega)\) with \(\text{supp} G_1 \subset \Omega \setminus \Omega_r\). A weak solution \((v_1, v_2) \in H_1(\Omega)\) of (4.28) is called a viscosity solution if for any \((\hat{v}_1, \hat{v}_2) \in H_1(\Omega)\), \((\hat{g}_1, \hat{g}_2) \in H_0(\Omega)\), and \(\hat{G}_1 \in [L^2(\Omega)]^d\) with \(\text{supp} \hat{G}_1 \subset \Omega \setminus \Omega_r\) such that
\[
\begin{cases}
\text{div}(A_1 \nabla \hat{v}_1) - \lambda_0 \Sigma_1 \hat{v}_1 = \hat{g}_1 + \text{div}(\hat{G}_1) & \text{in } \Omega, \\
\text{div}(A_2 \nabla \hat{v}_2) - \lambda_0 \Sigma_2 \hat{v}_2 = \hat{g}_2 & \text{in } \Omega, \\
A_1 \nabla \hat{v}_1 \cdot \nu - A_2 \nabla \hat{v}_2 \cdot \nu = 0 & \text{on } \Gamma.
\end{cases}
\] (4.30)

we have, with \(w = v_1 - v_2\) and \(\hat{w} = \hat{v}_1 - \hat{v}_2\),
\[
\int_\Omega \langle A_1 \nabla w, \nabla \hat{w} \rangle + \lambda_0 \Sigma_1 w \hat{w} = -\int_\Omega (g_1 - g_2) \hat{w} + \int_\Omega \lambda_0 (\Sigma_2 - \Sigma_1) v_2 \hat{w}
\]
\[+ \int_\Omega \langle [A_2 - A_1] \nabla v_2, \nabla \hat{w} \rangle + \int_\Omega G_1 \nabla \hat{w}. \] (4.31)

and
\[
\int_\Omega \langle [A_1 - A_2] \nabla v_2, \nabla \hat{w} \rangle + \langle [A_1 - A_2] \nabla v_2, \nabla \hat{v}_2 \rangle + \lambda_0 (\Sigma_1 - \Sigma_2) v_2 \hat{v}_2 \]
\[= \int_\Omega g_2 \hat{w} - (\hat{g}_1 - \hat{g}_2) v_2 + \hat{G}_1 \nabla v_2 + \lambda_0 (\Sigma_2 - \Sigma_1) v_2 \hat{w}. \] (4.32)

We are now in a position to state and prove the key result of this section.

\[ \text{In Proposition 4, } G_1 = 0; \text{ nevertheless, the same proof gives the same result in the case } G_1 \in L^2(\Omega) \text{ with } \text{supp} G_1 \subset \subset \Omega. \]
Proposition 5. There exists $\Lambda_0 > 1$ such that for $\lambda_0 > \Lambda_0$ and for $g = (g_1, g_2) \in H_0(\Omega)$ and $G_1 \in [L^2(\Omega)]^d$ with $\text{supp} G_1 \subset \Omega \setminus \Omega_\tau$, there exists a unique viscosity solution $v = (v_1, v_2) \in H_1(\Omega)$ of \eqref{1.8}. Moreover,

$$\|v\|_{H_1(\Omega)} \leq C \left( \|g\|_{H_0(\Omega)} + \|G_1\|_{L^2(\Omega)} \right),$$

for some positive constant $C$ independent of $g$ and $G_1$.

Proof. We first prove that viscosity solutions exist. By Lemma [15] there exists a unique solution $v_\delta = (v_{1,\delta}, v_{2,\delta}) \in [H^1(\Omega)]^2$ ($0 < \delta < 1$) of the system

$$\begin{cases}
\text{div}((1 + \delta) A_1 \nabla v_{1,\delta}) - \lambda_0 (1 + \delta) \Sigma_1 v_{1,\delta} = g_1 + \text{div}(G_1) & \text{in } \Omega, \\
\text{div}(A_2 \nabla v_{2,\delta}) - \lambda_0 \Sigma_2 v_{2,\delta} = g_2 & \text{in } \Omega,
\end{cases}$$

Moreover,

$$\int_{\Omega} \delta \left( \|\nabla v_{2,\delta}\|^2 + |v_{2,\delta}|^2 \right) + \|v_\delta\|^2_{H_1(\Omega)} \leq C \left( \|g\|^2_{H_0(\Omega)} + \|G_1\|^2_{L^2(\Omega)} \right),$$

for some positive constant $C$ independent of $\delta$, $g$, and $G_1$.

Set $w_\delta = v_{1,\delta} - v_{2,\delta}$. Let $v$ be the weak limit of $(v_{\delta_n})$ in $H_1^1(\Omega)$ for some sequence $(\delta_n) \to 0$ such that $w_{\delta_n} \rightharpoonup w := v_1 - v_2$ weakly in $H_0^1(\Omega)$. Then $v$ is a weak solution of \eqref{1.8}. We prove that $v$ is a viscosity solution. Multiplying the equation of $w_\delta$,

$$\text{div} \left( (1 + \delta) A_1 \nabla w_\delta \right) - \lambda_0 (1 + \delta) \Sigma_1 w_\delta = g_1 + \text{div}(G_1) - g_2 - \lambda_0 (\Sigma_2 - (1 + \delta) \Sigma_1) v_{2,\delta} + \text{div}([A_2 - (1 + \delta) A_1] \nabla v_{2,\delta})$$

by $\tilde{w}$ and integrating in $\Omega$, we obtain

$$\int_{\Omega} \langle (1 + \delta) A_1 \nabla w_\delta, \nabla \tilde{w} \rangle + \lambda_0 (1 + \delta) \Sigma_1 w_\delta \tilde{w} = - \int_{\Omega} (g_1 - g_2) \tilde{w}$$

$$+ \int_{\Omega} \lambda_0 (\Sigma_2 - (1 + \delta) \Sigma_1) v_{2,\delta} \tilde{w} + \int_{\Omega} \langle [A_2 - (1 + \delta) A_1] \nabla v_{2,\delta}, \nabla \tilde{w} \rangle + \int_{\Omega} G_1 \nabla \tilde{w}. \tag{4.35}$$

Similar to \eqref{2.49}, we have

$$\int_{\Omega} \langle [A_1 - A_2] \nabla v_{2,\delta}, \nabla \tilde{w} \rangle + \langle [A_1 - A_2] \nabla \tilde{v}_2, \nabla \tilde{w} \rangle + \lambda_0 (\Sigma_1 - \Sigma_2) v_{2,\delta} \tilde{v}_2$$

$$= \int_{\Omega} g_2 \tilde{w} - (\tilde{g}_1 - \tilde{g}_2) v_{2,\delta} + \tilde{G}_1 \nabla v_2 + \lambda_0 (\Sigma_2 - \Sigma_1) v_{2,\delta} \tilde{w}. \tag{4.36}$$

There is no term $(1 + \delta)$ in front of $A_1$ and $\Sigma_1$ in \eqref{4.36} since we only use the equation of $v_2$ and the system of $\tilde{v}$ here. Choosing $\delta = \delta_n$, letting $n \to +\infty$ in \eqref{4.35} and \eqref{4.36} and using \eqref{4.34}, we obtain \eqref{1.31} and \eqref{1.32}. Therefore, $v$ is a viscosity solution.

The uniqueness of viscosity solutions is now standard from its definition as in the proof of Proposition 5. \hfill \square
Set
\[ \mathcal{H}_0(\Omega) := \left\{ (v_1, v_2) \in [L^2_{loc}(\Omega)]^2; \| (v_1, v_2) \|_{\mathcal{H}_0(\Omega)} < +\infty \right\}, \tag{4.37} \]
where
\[ \langle u, v \rangle_{\mathcal{H}_0(\Omega)} = \int_{\Omega} \nabla u_2 \nabla v_2 + \int_{\Omega} d_\Gamma^{-\beta_1} (u_1 - u_2)(\bar{v}_1 - \bar{v}_2) + \int_{\Omega} d_\Gamma^{\beta_1} u \bar{v}. \tag{4.38} \]
Fix \( \lambda > \Lambda_0 \) where \( \Lambda_0 \) is the constant in Proposition 9. Define
\[ T_{3, \lambda}: \mathcal{H}_0(\Omega) \to \mathcal{H}_0(\Omega), \quad (f_1, f_2) \mapsto (u_1, u_2), \tag{4.39} \]
where \((u_1, u_2) \in \mathcal{H}_1(\Omega)\) is the unique viscosity solution of (4.8) with
\[ g_1 = \Sigma_1 (f_1 - f_2) + \Sigma_2 f_2, \quad G_1 = \lambda^{-1} (A_1 - A_2) \nabla f_2, \quad \text{and} \quad g_2 = \Sigma_2 f_2. \tag{4.40} \]
We have

**Lemma 16.** There exists \( \Lambda_0 > 1 \) such that if \( \lambda_0 > \Lambda_0 \) then \( T_{3, \lambda} \) is a compact operator.

**Proof.** Let \((f_{1,n}, f_{2,n})\) be an arbitrary bounded sequence in \( \mathcal{H}_0(\Omega) \) and set
\[ (u_{1,n}, u_{2,n}) = T_{3, \lambda}(f_{1,n}, f_{2,n}). \]
Denote \( g_{1,n}, g_{2,n}, G_1 \) as in (4.40) where \((f_1, f_2)\) is replaced by \((f_{1,n}, f_{2,n})\). It follows from (4.33) that \((u_{1,n}, u_{2,n})\) is bounded sequence in \( \mathcal{H}_1(\Omega) \). Using the equations of \( u_{1,n} \) and \( u_{2,n} \), we derive that \((u_{1,n}, u_{2,n})\) is bounded in \( H^1_{loc}(\Omega) \). Since \( \text{div}(A_2 \nabla u_{2,n}) - \lambda_0 \Sigma_2 u_{2,n} = g_{2,n} \), \((u_{2,n})\) is bounded in \( H^1_{loc}(\Omega) \), and \((g_{2,n})\) is bounded in \( L^2_{loc}(\Omega) \), one might assume that
\[ (u_{2,n}) \text{ converges strongly in } H^1_{loc}(\Omega). \tag{4.41} \]
By Hardy’s inequality, we have, with \( w_n = u_{1,n} - u_{2,n} \),
\[ \int_{\Omega} d_\Gamma^{-2} |w_n|^2 \leq C_{\Omega} \int_{\Omega} |\nabla w_n|^2. \]
Since the embedding \( H^1 \subset L^2 \) is compact, applying Lemma 17 below to \( \psi_n = w_n \), \( s = -2 \), and \( t = -\beta_1 \), one might assume in addition that
\[ (w_n) \text{ converges in } L^2(\Omega, d_\Gamma^{-\beta_1}). \tag{4.42} \]
Since the embedding \( H^1 \subset L^2 \) is compact, applying Lemma 17 below to \( \psi_n = u_n \), \( s = \beta \), and \( t = \beta_1 > \beta \), one might assume in addition that
\[ (u_n) \text{ converges in } [L^2(\Omega, d_\Gamma^{\beta_1})]^2. \tag{4.43} \]
By Lemma 6 we have
\[ \int_{\Omega} d_\Gamma^{\beta + 2} |\nabla u_{2,n}|^2 \leq C \int_{\Omega} d_\Gamma^{\beta} |u_{2,n}|^2 + d_\Gamma^{\beta} |g_{2,n}|^2 \leq C. \]
Using \( (4.41) \) and applying Lemma 17 to \( \psi_n = \nabla u_{2,n} \), \( s = \beta + 2 \), and \( t = \beta_1 + 2 \), one can assume in addition that
\[ \nabla u_{2,n} \text{ converges in } [L^2(\Omega, d_\Gamma^{\beta_1 + 2})]^d. \tag{4.44} \]
The convergence of \((u_n)\) in \( \mathcal{H}_0(\Omega) \) now follows from \( (4.41), (4.42), (4.43), \) and \( (4.44) \). \( \square \)
Remark 8. The embedding \( \mathcal{H}_1(\Omega) \) into \( \mathcal{H}_0(\Omega) \) is not compact. Nevertheless, \( T_{3,\lambda} \) is compact as shown in Lemma 16.

In the proof of Lemma 16, we used the following compactness result.

Lemma 17. Let \( m \geq 1 \) and \( (\psi_n) \subset [L^2_{loc}(\Omega)]^m \) be such that for every \( K \subset\subset \Omega \), \( (\psi_n|_K) \) is relatively compact in \( [L^2(K)]^m \). Assume that, for some \( s \in \mathbb{R} \),

\[
\sup_n \|\psi_n\|_{[L^2(\Omega,d^\tau_T)]^m} < +\infty.
\]

Then \( (\psi_n) \) is relatively compact in \( [L^2(\Omega,d^t_T)]^m \) for every \( t > s \).

Proof. The proof is almost in the spirit of the one [29, Lemma 7]. For the sake of completeness, we present the proof. Set

\[
C = \sup_n \|\psi_n\|_{[L^2(\Omega,d^t_T)]^m}.
\]

We have, for \( \tau > 0 \),

\[
\int_{\Omega_T} d^\tau_T |\psi_n|^2 \leq \int_{\Omega_{\tau}} d^\tau_T |\psi_n|^2 \leq C \tau^{-2}. \quad (4.45)
\]

Fix \( \varepsilon > 0 \) arbitrary. Fix \( \tau > 0 \) small enough such that

\[
\|\psi_n\|_{L^2(\Omega_\tau)} \leq \varepsilon/2 \quad \forall n \in \mathbb{N}. \quad (4.46)
\]

Such a constant \( \tau \) exists by (4.45). Since \( (\psi_n) \) is relatively compact in \( L^2(\Omega \setminus \Omega_\tau) \), there exist \( \psi_{n_1}, \cdots, \psi_{n_k} \) such that

\[
\left\{ \psi_n \in L^2(\Omega \setminus \Omega_\tau); \ n \in \mathbb{N} \right\} \subset \bigcup_{j=1}^k \left\{ \psi \in L^2(\Omega \setminus \Omega_\tau); \|\psi - \psi_{n_j}\|_{L^2(\Omega_\tau,d^t_T)} \leq \varepsilon/2 \right\}. \quad (4.47)
\]

A combination of (4.46) and (4.47) yields

\[
\left\{ \psi_n \in L^2(\Omega,d^t_T); \ n \in \mathbb{N} \right\} \subset \bigcup_{j=1}^k \left\{ \psi \in L^2(\Omega); \|\psi - \psi_{n_j}\|_{L^2(\Omega)} \leq 3\varepsilon/2 \right\}.
\]

Therefore, \( (\psi_n) \) is relatively compact in \( L^2(\Omega,d^t_T) \).

We are ready to give

Proof of Theorem 31. By the spectral theory of a family of compact analytic operators, the spectrum of \( T_{3,\lambda} \) is discrete (see, e.g., [23, Theorem 8.92]). The conclusion follows by the definition of \( T_{3,\lambda} \).
5 On the case $A_1$ and $A_2$ satisfying the complementary condition on $\Gamma$ - Proof of Theorem 4

This section is devoted to the proof of Theorem 4. We first establish the well-posedness in $[H^1(\Omega)]^2$ of the following system (Proposition 10), for $\lambda_0 > 0$ (large),

\[
\begin{align*}
\text{div}(A_1 \nabla u_1) - i\lambda_0 \Sigma_1 u_1 &= g_1 & \text{in } \Omega, \\
\text{div}(A_2 \nabla u_2) - i\lambda_0 \Sigma_2 u_2 &= g_2 & \text{in } \Omega, \\
u_1 &= u_2, & A_1 \nabla u_1 \cdot \nu = A_2 \nabla u_2 \cdot \nu & \text{on } \Gamma,
\end{align*}
\]

(5.1)

for a given pair $(g_1, g_2) \in [L^2(\Omega)]^2$ under the following three assumptions, which are the assumptions of Theorem 4

$C1)$ $A_1$, $A_2$, $\Sigma_1$, $\Sigma_2$ are continuous in a neighbourhood of $\Gamma$.

$C2)$ For all $x \in \Gamma$, $A_1(x)$ and $A_2(x)$ satisfy the following condition, with $\nu = \nu(x)$,

\[
\langle A_2(x) e, e \rangle \langle A_2(x) \xi, \xi \rangle - \langle A_2(x) e, \xi \rangle^2 \neq \langle A_1(x) e, e \rangle \langle A_1(x) \xi, \xi \rangle - \langle A_1(x) e, \xi \rangle^2
\]

for all $\xi \in \mathcal{P}(x) \setminus \{0\}$, where

\[
\mathcal{P}(x) := \{ \xi \in \mathbb{R}^d ; \langle \xi, e \rangle = 0 \}.
\]

$C3)$ For all $x \in \Gamma$, $A_1(x)$, $A_2(x)$, $\Sigma_1(x)$, and $\Sigma_2(x)$ verify

\[
\langle A_1(x) \nu(x), \nu(x) \rangle \Sigma_1(x) \neq \langle A_2(x) \nu(x), \nu(x) \rangle \Sigma_2(x).
\]

(5.2)

We then define the operator

\[
T_4(f_1, f_2) = (u_1, u_2) \text{ where } (u_1, u_2) \text{ is the unique solution of (5.1) with } (g_1, g_2) = (\Sigma_1 f_1, \Sigma_2 f_2).
\]

(5.3)

The basic ingredient in the proof of Proposition 10 is the following lemma:

Lemma 18. Let $\lambda \geq 1$, $A_1, A_2$ be two constant positive symmetric matrices and let $\Sigma_1, \Sigma_2$ be two positive constants. Assume that

\[
\langle A_2 e_d, e_d \rangle \langle A_2 \xi, \xi \rangle - \langle A_2 e_d, \xi \rangle^2 \neq \langle A_1 e_d, e_d \rangle \langle A_1 \xi, \xi \rangle - \langle A_1 e_d, \xi \rangle^2
\]

for all $\xi \in \mathcal{P} \setminus \{0\}$, where

\[
\mathcal{P} := \{ \xi \in \mathbb{R}^d ; \langle \xi, e_d \rangle = 0 \},
\]

and

\[
\langle A_1 e_d, e_d \rangle \Sigma_1 \neq \langle A_2 e_d, e_d \rangle \Sigma_2.
\]

(5.4)

(5.5)

Let $g = (g_1, g_2) \in [L^2(\mathbb{R}^d_+)]^2$, $G = (G_1, G_2) \in ([L^2(\mathbb{R}^d_+)]^d)^2$, $\varphi \in H^{1/2}(\mathbb{R}^d_0)$, $\phi \in H^{-1/2}(\mathbb{R}^d_0)$. Assume that $v = (v_1, v_2) \in [H^1(\mathbb{R}^d_+)]^2$ is a solution of the system

\[
\begin{align*}
\text{div}(A_1 \nabla v_1) - i\lambda \Sigma_1 v_1 &= g_1 + \text{div}(G_1) & \text{in } \mathbb{R}^d_+, \\
\text{div}(A_2 \nabla v_2) - i\lambda \Sigma_2 v_2 &= g_2 + \text{div}(G_2) & \text{in } \mathbb{R}^d_+, \\
v_1 - v_2 &= \varphi, & (A_1 \nabla v_1 - G_1) \cdot e_d - (A_2 \nabla v_2 - G_2) \cdot e_d &= \phi & \text{on } \mathbb{R}^d_0.
\end{align*}
\]

(5.6)
We have
\[ \|v\|_{H^1(\mathbb{R}_+^d)} + \lambda^{1/2}\|v\|_{L^2(\mathbb{R}_+^d)} \]
\[ \leq C \left( \lambda^{-1/2}\|g\|_{L^2(\mathbb{R}_+^d)} + \|G\|_{L^2(\mathbb{R}_+^d)} + \|\varphi\|_{H^{1/2}(\mathbb{R}_+^d)} + \lambda^{1/4}\|\varphi\|_{L^2(\mathbb{R}_+^d)} + \|\phi\|_{H^{-1/2}(\mathbb{R}_+^d)} \right), \]  
(5.7)
where \( C \) is a positive constant depending only on \( d, A_1, A_2, \Sigma_1, \) and \( \Sigma_2 \).

Here and in what follows in this section,
\[ \mathbb{R}^d_+ := \{ x = (x', x_d) \in \mathbb{R}^{d-1} \times \mathbb{R}; \ x_d > 0 \} \quad \text{and} \quad \mathbb{R}^d_0 := \{ x = (x', x_d) \in \mathbb{R}^{d-1} \times \mathbb{R}; \ x_d = 0 \}. \]

**Proof.** Let \( u_j \in H^1(\mathbb{R}_+^d) \) \((j = 1, 2)\) be the unique solution of the system
\[ \text{div}(A_j \nabla u_j) - i\lambda \Sigma_j u_j = g_j + \text{div}(G_j) \quad \text{in} \ \mathbb{R}_+^d, \]  
and
\[ (A_j \nabla u_j - G_j) \cdot e_d = \phi_j, \]  
with \( \phi_1 = \phi \) and \( \phi_2 = 0 \). Multiplying the equation of \( u_j \) by \( \bar{u}_j \) and integrating on \( \mathbb{R}_+^d \), we have
\[ \int_{\mathbb{R}_+^d} \langle A_j \nabla u_j, \nabla u_j \rangle + i\lambda \Sigma_j |u_j|^2 = \int_{\mathbb{R}_+^d} G_j \nabla \bar{u}_j - g_j \bar{u}_j + \int_{\mathbb{R}_0^d} \phi_j \bar{u}_j; \]
which implies
\[ \|\nabla u_j\|_{L^2(\mathbb{R}_+^d)} + \lambda^{1/2}\|u_j\|_{L^2(\mathbb{R}_+^d)} \leq C\left( \lambda^{-1/2}\|g_j\|_{L^2(\mathbb{R}_+^d)} + \|G_j\|_{L^2(\mathbb{R}_+^d)} + \|\phi_j\|_{H^{1/2}(\mathbb{R}_+^d)} \right). \]  
(5.8)

Since, for \( j = 1, 2, \)
\[ \int_{\mathbb{R}^{d-1}} |u_j(x')|^2 \, dx' \leq 2 \int_{\mathbb{R}_+^d} |\partial_{x_d} u_j||u_j|, \]
it follows that, for \( j = 1, 2, \)
\[ \|u_j\|_{L^2(\mathbb{R}_+^d)}^2 \leq 2\|u_j\|_{L^2(\mathbb{R}_+^d)}\|\nabla u_j\|_{L^2(\mathbb{R}_+^d)}. \]  
(5.9)

A combination of \( (5.8) \) and \( (5.9) \) yields, for \( j = 1, 2, \)
\[ \|u_j\|_{H^1(\mathbb{R}_+^d)} + \lambda^{1/2}\|u_j\|_{L^2(\mathbb{R}_+^d)} + \|u_j\|_{H^{1/2}(\mathbb{R}_+^d)} + \lambda^{1/4}\|u_j\|_{L^2(\mathbb{R}_+^d)} \]
\[ \leq C\left( \lambda^{-1/2}\|g_j\|_{L^2(\mathbb{R}_+^d)} + \|G_j\|_{L^2(\mathbb{R}_+^d)} + \|\phi_j\|_{H^{-1/2}} \right). \]  
(5.10)

By considering the system of \((v_1 - u_1, v_2 - u_2)\), one might assume that \(g_1 = g_2 = 0, G_1 = G_2 = 0,\) and \( \phi = 0 \). In what follows, we make this assumption.

Let \( \hat{v}_j(\xi', t) \) for \( j = 1, 2 \) and \( \hat{\varphi}(\xi', t) \) be the Fourier transform of \( v_j \) and \( \varphi \) with respect to \( x' \in \mathbb{R}^{d-1}, \) i.e., for \((\xi', t) \in \mathbb{R}^{d-1} \times (0, +\infty), \)
\[ \hat{v}_j(\xi', t) = \int_{\mathbb{R}^{d-1}} v_j(x', t)e^{-ix' \cdot \xi'} \, dx' \quad \text{for} \ j = 1, 2, \quad \text{and} \quad \hat{\varphi}(\xi', t) = \int_{\mathbb{R}^{d-1}} \varphi(x')e^{-ix' \cdot \xi'} \, dx'. \]
Since
\[ \text{div}(A_j \nabla v_j) - i \lambda \Sigma_j v_j = 0 \text{ in } \mathbb{R}^d, \]
it follows that
\[ a_j \ddot{v}_j(t) + 2ib_j \dot{v}_j(t) - (c_j + i \lambda \Sigma_j) \dot{v}_j(t) = 0 \text{ for } t > 0, \quad (5.11) \]
where
\[ a_j = (A_j)_{d,d}, \quad b_j = \sum_{k=1}^{d-1} (A_j)_{d,k} \xi_k, \quad \text{and} \quad c_j = \sum_{k=1}^{d-1} \sum_{l=1}^{d-1} (A_j)_{k,l} \xi_k \xi_l. \quad (5.12) \]
Here \( (A_j)_{k,l} \) denotes the \((k, l)\) component of \( A_j \) for \( j = 1, 2 \) and the symmetry of \( A_j \) is used.

Define, for \( j = 1, 2 \),
\[ \Delta_j = -b_j^2 + a_j (c_j + i \lambda \Sigma_j). \quad (5.13) \]
Denote \( \xi = (\xi', 0) \). Since \( A_j \) is symmetric and positive, it is clear that, for \( j = 1, 2 \),
\[ a_j = (A_j e_d, e_d) > 0, \quad b_j = (A_j \xi, e_d), \quad \text{and} \quad a_j c_j - b_j^2 = (A_j e_d, e_d) \langle A_j \xi, \xi \rangle - (A_j e_d, \xi)^2 > 0. \quad (5.14) \]

For \( j = 1, 2 \), let \( \sqrt{\Delta_j} \) denote the square root of \( \Delta_j \) with positive real part and set
\[ \eta_j = (-ib_j - \sqrt{\Delta_j})/a_j. \]
Since \( \dot{v}_j(\xi', t) \in L^2(\mathbb{R}^d_0) \), we derive from (5.11) that, for \( j = 1, 2 \),
\[ \dot{v}_j(\xi', t) = \alpha_j(\xi') e^{\eta_j t}, \quad (5.15) \]
for some \( \alpha_j \in L^2(\mathbb{R}^d_0) \). Using the fact that \( v_1 - v_2 = \varphi \) and \( A_1 \nabla v_1 \cdot e_d - A_2 \nabla v_2 \cdot e_d = 0 \) on \( \mathbb{R}^d_0 \), we derive that
\[ \alpha_1(\xi') - \alpha_2(\xi') = \dot{\varphi}(\xi') \quad \text{and} \quad \alpha_1(\xi') (iA_1 \xi + \eta_1 A_1 e_d, e_d) - \alpha_2(\xi') (iA_2 \xi + \eta_2 A_2 e_d, e_d) = 0. \quad (5.16) \]

Since, by (5.14),
\[ (A_j \xi, e_d) - (A_j e_d, e_d) b_j/a_j = 0 \text{ for } j = 1, 2, \]
the last identity of (5.16) implies
\[ \alpha_1(\xi') \sqrt{\Delta_1} = \alpha_2(\xi') \sqrt{\Delta_2}. \]
Combining this identity and the first one of (5.16) yields
\[ \alpha_1(\xi') = \frac{\dot{\varphi}(\xi') \sqrt{\Delta_2}}{\sqrt{\Delta_2} - \sqrt{\Delta_1}}. \quad (5.17) \]
Note that, by (5.4), (5.5), (5.13), and (5.14),
\[ |\Delta_2 - \Delta_1|^2 = \left( (A_2 e_d, e_d) (A_2 \xi, \xi) - (A_2 e_d, \xi)^2 - (A_1 e_d, e_d) (A_1 \xi, \xi) + (A_1 e_d, \xi)^2 \right)^2 \]
\[ + \lambda^2 \left( (A_1 e_d, e_d) \Sigma_1 - (A_2 e_d, e_d) \Sigma_2 \right)^2 \]
\[ \geq C (|\xi|^4 + \lambda^2) \]
38
and

\[ |\Delta_j| \leq C(|\xi|^2 + \lambda). \]

This implies

\[ \frac{\sqrt{\Delta_2}}{|\sqrt{\Delta_2} - \sqrt{\Delta_1}|} = \frac{\sqrt{\Delta_2}(\sqrt{\Delta_2} + \sqrt{\Delta_1})}{|\Delta_2 - \Delta_1|} \leq C. \]  \hspace{1cm} (5.18)

Since, by Parseval’s theorem,

\[ \int_{\mathbb{R}^d_+} |v_1|^2 = C_d \int_0^\infty \int_{\mathbb{R}^{d-1}} |\hat{v}_1(\xi', t)|^2 \, d\xi' \, dt, \]

it follows from (5.15), (5.17), and (5.18) that

\[ \int_{\mathbb{R}^d_+} |v_1|^2 \leq C \int_{\mathbb{R}^{d-1}} \int_0^\infty |\hat{\varphi}(\xi')|^2 |e^{2\eta t}| \, dt \, d\xi'. \]  \hspace{1cm} (5.19)

Using the fact

\[ \int_0^\infty |e^{2\eta t}| \, dt \leq \frac{1}{2|\Re(\eta)|} \quad \text{and} \quad \frac{1}{|\Re(\eta)|} \leq \frac{C}{|\xi'| + \sqrt{\lambda}}, \]

we deduce from (5.19) that

\[ \int_{\mathbb{R}^d_+} |v_1|^2 \leq \int_{\mathbb{R}^{d-1}} \frac{C|\hat{\varphi}(\xi')|^2}{|\xi'| + \sqrt{\lambda}} \, d\xi' \leq C\lambda^{-1/2}||\varphi||^2_{L^2(\mathbb{R}^d_+)} \]

which yields the corresponding estimate for \( ||v_1||_{L^2(\mathbb{R}^d_+)} \). The estimates follows similarly. The details are left to the reader. \( \square \)

**Remark 9.** The computation of \( \hat{v}_j \) in Lemma 18 has roots from the proof of Proposition 1.

**Proposition 10.** Assume C1), C2), and C3). There exists a constant \( \Lambda_0 > 0 \) depending only on \( A_1, A_2, \Sigma_1, \Sigma_2, \) and \( \Omega \), such that for \( \lambda > \Lambda_0 \) and \( g = (g_1, g_2) \in [L^2(\Omega)]^2 \), there exists a unique solution \( \nu = (v_1, v_2) \in [H^1(\Omega)]^2 \) satisfy

\[
\begin{aligned}
\operatorname{div}(A_1 \nabla v_1) - i\lambda \Sigma_1 v_1 &= g_1 & \text{in } \Omega, \\
\operatorname{div}(A_2 \nabla v_2) - i\lambda \Sigma_2 v_2 &= g_2 & \text{in } \Omega, \\
v_1 - v_2 &= 0, & A_1 \nabla v_1 \cdot \nu - A_2 \nabla v_2 \cdot \nu &= 0 & \text{on } \Gamma.
\end{aligned}
\]

Moreover,

\[ ||v||_{H^1(\Omega)} + \lambda^{1/2}||v||_{L^2(\Omega)} \leq C\lambda^{-1/2}||g||_{L^2(\Omega)}, \]  \hspace{1cm} (5.21)

for some positive constant \( C \) independent of \( g \) and \( \lambda \).

**Proof.** We first assume the existence of \( v \in [H^1(\Omega)]^2 \) and derive the estimate for \( v \). Assume that \( (A_1, A_2, \Sigma_1, \Sigma_2) \in C(\hat{\Omega}_{\tau_0}) \) for some \( \tau_0 > 0 \). Using local charts, applying Lemma 18 with \( \varphi = 0 \) and \( \phi = 0 \), and involving the standard freezing coefficient technique, we have

\[ ||v||_{H^1(\Omega_{\tau})} + \lambda^{1/2}||v||_{L^2(\Omega_{\tau})} \leq C \left( ||v||_{H^1(\Omega_{2\tau} \setminus \Omega_{\tau})} + \lambda^{-1/2}||(g_1, g_2)||_{L^2(\Omega)} \right), \]  \hspace{1cm} (5.22)
if $\tau > 0$ is small enough and $\lambda > 1$. Here and in what follows $C$ denotes a positive constant independent of $v$, $g$, and $\lambda$; $C$ depends on $(A_1, A_2, \Sigma_1, \Sigma_2)$, $\Omega$, and $\tau$. Fix such a positive constant $\tau$. Applying Lemma 9 to $v$, we have

$$\|v\|_{H^1(\Omega, \Omega)} \leq c_1 \exp(-c_2 \sqrt{\lambda}) \|v\|_{L^2(\Omega)} + c_1 \|g\|_{L^2(\Omega)},$$

(5.23)

for some positive $c_1$ and $c_2$ independent of $\lambda$, $v$, and $g$. Combining (5.22) and (5.23) yields, for large $\lambda$,

$$\|v\|_{H^1(\Omega)} + \lambda^{1/2} \|v\|_{L^2(\Omega)} \leq C \|g\|_{L^2(\Omega)};$$

which is (5.21). The existence of $v \in [H^1(\Omega)]^2$ follows from the uniqueness via the Fredholm theory by noting the well-posedness in $[H^1(\Omega)]^2$ of the system

\[
\begin{aligned}
\text{div}(A_1 \nabla v_1) - i\lambda \Sigma_1 v_1 &= g_1 & \text{in } \Omega, \\
\text{div}(A_2 \nabla v_2) + i\lambda \Sigma_2 v_2 &= g_2 & \text{in } \Omega, \\
v_1 - v_2 &= 0, & A_1 \nabla v_1 \cdot \nu - A_2 \nabla v_2 \cdot \nu &= 0 & \text{on } \Gamma,
\end{aligned}
\]

for $g \in [L^2(\Omega)]^2$ by Lax-Milgram’s theory. The details are left to the reader. □

We are ready to present

**Proof of Theorem 4.** Fix $\lambda > \Lambda_0$ where $\Lambda_0$ is the positive constant in Proposition 10. Define

$$T_4 : [L^2(\Omega)]^2 \rightarrow [L^2(\Omega)]^2$$

(5.24)

where $(u_1, u_2) \in [H^1(\Omega)]^2$ is the unique solution of (5.1) with $(g_1, g_2) = (\Sigma_1 f_1, \Sigma_2 f_2)$. Since $H^1(\Omega) \subset L^2(\Omega)$ is compact, so $T_4$ is compact. The conclusion of Theorem 4 follows. □

**References**

[1] S. Agmon, A. Douglis, and L. Nirenberg, *Estimates near the boundary for solutions of elliptic partial differential equations satisfying general boundary conditions. II.*, Comm. Pure Appl. Math. 17 (1964), 35–92.

[2] S. Agmon, *Lectures on elliptic boundary value problems*, second ed., American Mathematical Soc., vol 369, 2010.

[3] E. Blåsten, L. Paivarinta *Completeness of generalized transmission eigenstates*, Inverse Problem 29 (2013), 104002.

[4] A. S. Bonnet-BenDhia, L. Chesnel, H. Haddar, *On the use of T-coercivity to study the interior transmission eigenvalue problem*, C. R. Math. Acad. Sci. Paris, 349 (2011), 647–651.

[5] H. Brezis, *Functional Analysis, Sobolev Spaces and Partial Differential Equations*, Universitext, Springer, 2010.
[6] F. Cakoni, D. Colton, H. Haddar, *The linear sampling method for anisotropic media*, J. Comput. Appl. Math. **146** (2002), 285–299.

[7] F. Cakoni, D. Colton, H. Haddar. *The interior transmission problem for regions with cavities* SIAM J. Math. Analysis **42** (2010), 145–162.

[8] F. Cakoni, H. Haddar. *On the existence of transmission eigenvalues in an inhomogeneous medium*, Applicable Analysis **88** (2009), 475–493.

[9] F. Cakoni, H. Haddar. *Transmission eigenvalues in inverse scattering theory*. Inside out II MSRI Publications **60** (2012), 529–580.

[10] D. Colton and P. Monk, *The inverse scattering problem for acoustic waves in an inhomogeneous medium*, Quart. J. Mech. Appl. Math. **41** (1988), 97–125.

[11] D. Colton, A. Kirsch, L. Paivarinta. *Far-Field patterns for acoustic waves in an inhomogeneous medium*. SIAM J. Math. Anal. **20** (1989) 1472–1483.

[12] D. Colton and R. Kress, *Inverse acoustic and electromagnetic scattering theory*, second ed., Applied Mathematical Sciences, vol. 98, Springer-Verlag, Berlin, 1998.

[13] D. Colton, L. Paivarinta, *Transmission eigenvalues and a problem of Hans Lewy*, J. Comput. Appl. Math. **117** (2000), 91–104.

[14] M. Costabel and E. Stephan, *A direct boundary integral equation method for transmission problems*, J. Math. Anal. Appl. **106** (1985), 367–413.

[15] M. Faierman, *Eigenvalue asymptotics for a boundary problem involving an elliptic system*, Math. Nachr. **279** (2006), 1159-1184.

[16] H. Haddar, P. Joly, and H-M. Nguyen. *Generalized impedance boundary conditions for scattering by strongly absorbing obstacles: the scalar case*. Math. Models Methods Appl. Sci., **15** (2005) 1273–1300.

[17] M. Hitrik, K. Krupchyk, P. Ola, L. Paivarinta, *The interior transmission problem and bounds on transmission eigenvalues*, Math. Res. Lett., **18** (2011), 279–293.

[18] A. Kirsch, *The denseness of the far field patterns for the transmission problem*, IMA J.Appl. Math. **37** (1986), 213–225.

[19] A. Kirsch, *On the existence of transmission eigenvalues*, Inverse Probl. Imaging **3** (2009), 155–172.

[20] E. Lakshtanov and B. Vainberg, *Ellipticity in the interior transmission problem in anisotropic media*, SIAM J. Math. Anal. **44** (2012), 1165–1174.

[21] E. Lakshtanov and B. Vainberg, *Applications of elliptic operator theory to the isotropic interior transmission eigenvalue problem*, Inverse Problems **29** (2013), 104003.
[22] L. Robbiano, *Spectral analysis on interior transmission eigenvalues*, Inverse Problems 29 (2013), 104001.

[23] Y. B. Lopatinskii, *On a method of reducing boundary problems for a system of differential equations of elliptic type to regular integral equations*, Ukrain. Mat. Z 5 (1953), 123–151.

[24] H-M. Nguyen, *Asymptotic behavior of solutions to the Helmholtz equations with sign changing coefficients*, Trans. Amer. Math. Soc. 367 (2015), 6581-6595.

[25] H-M. Nguyen, *Superlensing using complementary media*, Ann. Inst. H. Poincaré Anal. Non Linéaire 32 (2015), 471–484.

[26] H-M. Nguyen, *Cloaking via anomalous localized resonance. A connection between the localized resonance and the blow up of the power for doubly complementary media*, C. R. Math. Acad. Sci. Paris 353 (2015), 41–46.

[27] H-M. Nguyen, *Cloaking via anomalous localized resonance for doubly complementary media in the quasi static regime*, J. Eur. Math. Soc. (JEMS) 17 (2015), 1327-1365.

[28] H-M. Nguyen, *Cloaking using complementary media in the quasistatic regime*, Ann. Inst. H. Poincaré Anal. Non Linéaire, to appear, doi:10.1016/j.anihpc.2015.06.004.

[29] H-M. Nguyen, *Limiting absorption principle and well-posedness for the Helmholtz equation with sign changing coefficients*, J. Math. Pures Appl. 106 (2016), 342–374.

[30] H-M. Nguyen, *Cloaking via anomalous localized resonance for doubly complementary media in the finite frequency regime*, https://arxiv.org/abs/1511.08053.

[31] H-M. Nguyen, *Cloaking an arbitrary object via anomalous localized resonance: the cloak is independent of the object: the acoustic case*, https://arxiv.org/abs/1607.06492.

[32] H-M. Nguyen and H. L. Nguyen, *Cloaking using complementary media for the Helmholtz equation and a three spheres inequality for second order elliptic equations*, Trans. Amer. Math. Soc. Ser. B, 2 (2015), 93-112.

[33] M. Renardy and R.C. Rogers, *An introduction to partial differential equations*, Texts in Applied Mathematics vol. 13, Second ed., New York: Springer-Verlag, 2004.

[34] B. P. Rynne and B. D. Sleeman, *The interior transmission problem and inverse scattering from inhomogeneous media*, SIAM J. Math. Anal. 22 (1991), 1755–1762.

[35] J. Sylvester, *Discreteness of transmission eigenvalues via upper triangular compact operators*, SIAM J. Math. Anal. 44 (2012), 341–354.