



Existence Results of Two Mixed Boundary Value Elliptic PDEs in \mathbb{R}^{N}

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Abstract

We study the existence of a solution to the mixed boundary value problem for Helmholtz and Poisson type equations in a bounded Lipschitz domain $\Omega \subset \mathbb{R}^N$ and in $\mathbb{R}^N \setminus \Omega$ for $N \ge 3$. The boundary $\partial \Omega$ of Ω is the decomposition of $\Gamma_1, \Gamma_2 \subset \partial \Omega$ such that $\partial \Omega = \Gamma = \overline{\Gamma}_1 \cup \Gamma_2 = \Gamma_1 \cup \overline{\Gamma}_2$ and $\Gamma_1 \cap \Gamma_2 = \emptyset$. We have shown that if the Neumann data $f_2 \in H^{-\frac{1}{2}}(\Gamma_2)$ and the Dirichlet data $f_1 \in H^{\frac{1}{2}}(\Gamma_1)$ then the Helmholtz problem with mixed boundary data admits a unique solution. We have also shown the existence of a weak solution to a mixed boundary value problem governed by the Poisson equation with a measure data and the Dirichlet, Neumann data belongs to $f_1 \in H^{\frac{1}{2}}(\Gamma_1), f_2 \in H^{-\frac{1}{2}}(\Gamma_2)$, respectively.

Keywords Mixed boundary value problem \cdot Sobolev space \cdot Newton potential \cdot Boundary integral operator \cdot Layer potentials \cdot Radon measure

Mathematics Subject Classification 35J25 · 31B10 · 35J20

1 Introduction

The Poisson problem with mixed Dirichlet–Neumann boundary conditions deals with conductivity, heat transfer, metallurgical melting, wave phenomena, elasticity and electrostatics in mathematical physics and engineering. The detailed applications can be found in Dauge (1992), Fabrikant (1991), Jochmann (1999), Lagnese (1983), Maz'ya and Rossmann (2006), Rempel and Schulze (1989), Simanca (1988), Sneddon (1966), Wendland et al. (1979) and the references therein. A common problem of interest found in the literature is the following mixed boundary value (MBVP).

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$$Lu = h \text{ in } \Omega,$$

$$u = f \text{ on } \Gamma_1,$$

$$Mu = g \text{ on } \Gamma_2,$$

(1.1)

where, Ω is a bounded Lipschitz domain in \mathbb{R}^N for $N \ge 3$. The boundary of Ω , which will be denoted by Γ , is the disjoint union of Γ_1 and Γ_2 which are subsets of Γ such that $\overline{\Gamma}_1 \cup \Gamma_2 = \Gamma_1 \cup \overline{\Gamma}_2 = \Gamma$ and $\Gamma_1 \cap \Gamma_2 = \emptyset$. Further, *L* is a second order elliptic operator, *M* is a general first order oblique differential operator on Γ_2 .

Lieberman (1986, 1989) considered the problem (1.1) and proved the existence and Hölder continuity of classical solutions with smooth data. The techniques used in the corresponding Dirichlet problem ($\Gamma_2 = \emptyset$) and oblique derivative problem ($\Gamma = \Gamma_2$) of (1.1) are helpful to show the existence of solutions to the mixed boundary value problem. It is worth to mention the work due to Azzam and Kreyszig (1982), as they have provided the regularity result for MBVP in a plane domain with corners, where the Dirichlet data belongs to $C^{2,\alpha}(\Gamma \setminus \{0\})$ and the remaining boundary data is in $C^{1,\alpha}(\Gamma \setminus \{0\})$. The work due to Sykes and Brown (2001) deals with the boundary regularity of problem (1.1) with Dirichlet and Neumann boundary conditions where L is the Laplacian operator, h = 0 in Ω , $g \in L^p(\Gamma_2)$, $f \in W^{1,p}(\Gamma_1)$ for 1 $and the angle between <math>\Gamma_1$, Γ_2 should be strictly less than π in the interface. Sykes and Brown (2001) drew motivation from Brown (1994), who considered the two boundary data as $f \in H^1(\Gamma_1)$ and $g \in L^2(\Gamma_2)$.

Not much of literature is found for MBVP involving a measure data, although Liang and Rodrigues (1996) considered a problem involving measure data both on the domain and on the boundary Γ_2 . Some work has been done by Gallouët and Sire (2011), where the nonlinearity lies on the boundary with measure supported on the domain Ω and on the boundary Γ_2 . The MBVP in Gallouët and Sire (2011) possesses a weak solution *u* in $W^{1,q}(\Omega)$, for all $1 < q < \frac{N}{N-1}$ and the trace of *u* on Γ lies in $W^{1-\frac{1}{q},q}(\Gamma)$, for all $1 \le q < \frac{N}{N-1}$.

In this article we have considered the following two mixed boundary value problems. The first problem (P1) is

$$-\Delta u - \lambda^2 u = h \text{ in } \Omega,$$

$$u = f_1 \text{ on } \Gamma_1,$$

$$\frac{\partial u}{\partial \hat{n}} = f_2 \text{ on } \Gamma_2$$
(1.2)

and

$$-\Delta u - \lambda^2 u = 0 \text{ in } \mathbb{R}^N \setminus \overline{\Omega},$$

$$u = f_1 \text{ on } \Gamma_1,$$

$$\frac{\partial u}{\partial \hat{n}} = f_2 \text{ on } \Gamma_2,$$
(1.3)

where *u* satisfies the following conditions at infinity, i.e. $|x| \to \infty$.

For $\lambda = 0$

$$u(x) = O(|x|^{2-N}).$$
(1.4)

For $\lambda \neq 0$ (Sommerfeld's radiation condition)

$$u(x) = O(|x|^{\frac{1-N}{2}}),$$

$$\frac{\partial u(x)}{\partial |x|} - i\lambda u(x) = o(|x|^{\frac{1-N}{2}}).$$
(1.5)

The second problem (P2) is

$$-\Delta u = \mu \text{ in } \Omega,$$

$$u = f_1 \text{ on } \Gamma_1,$$

$$\frac{\partial u}{\partial \hat{n}} = f_2 \text{ on } \Gamma_2$$
(1.6)

and

$$-\Delta u = 0 \text{ in } \mathbb{R}^N \setminus \overline{\Omega},$$

$$u = f_1 \text{ on } \Gamma_1,$$

$$\frac{\partial u}{\partial \hat{n}} = f_2 \text{ on } \Gamma_2,$$
(1.7)

where u satisfies

$$|u(x)| + |x| |\nabla u(x)| = O(|x|^{2-N}), \text{ as } |x| \to \infty.$$
 (1.8)

Throughout the article $\frac{\partial u}{\partial \hat{n}}$ will represent the normal derivative with respect to the outward unit normal \hat{n} to the boundary, $f_1 \in H^{\frac{1}{2}}(\Gamma_1)$, $f_2 \in H^{-\frac{1}{2}}(\Gamma_2)$ are boundary data, $h \in \tilde{H}^{-1}(\Omega)$, μ will denote a bounded Radon measure, $\lambda \in \mathbb{C}$ with $Im(\lambda) \geq 0$ and λ^2 will be different from the eigenvalues of the Laplacian $(-\Delta)$. We will, at some places, refer problem (1.2), (1.6) as interior problems (IP1), (IP2) respectively and (1.3), (1.7) as exterior problems (EP1), (EP2) respectively. This work is motivated from the work of Chang and Choe (2008) and Ernst (1987) where the authors have used the method of layer potentials to show the uniqueness of solution to the homogeneous mixed boundary value problem in both interior and exterior domains. Chang and Choe (2008) has shown that for h = 0 and $\lambda = 0$ the solution u belongs to $H^1(\Omega)$ for the interior problem and belongs to $H^{1}_{loc}(\mathbb{R}^N \setminus \overline{\Omega})$ for the exterior problem. This u also satisfies the following inequality.

$$\|u\|_{H^{1/2}(\Gamma_1)} + \left\|\frac{\partial u}{\partial \hat{n}}\right\|_{H^{-1/2}(\Gamma_2)} \le C\left\{\|f_1\|_{H^{1/2}(\Gamma_1)} + \|f_2\|_{H^{-1/2}(\Gamma_2)}\right\}$$

where C is independent of f_1 , f_2 and h.

The novelty of our work is the consideration of two nonhomogenous mixed boundary value problems and a Radon measure μ as a nonhomogeneous term in (P2), for which the solution space becomes weaker than the Sobolev space $H^1(\Omega)$. We have shown the existence of a unique solution to the problem (P1) by boundary layer method in Sect. 3 using the representation formula of solution and the invertibility of boundary layer operators of Helmholtz equation given in Sects. 2.1 and 2.2, respectively. Due to insufficient given boundary data, we have also found out the cauchy data, i.e. $u|_{\Gamma}$ and $\frac{\partial u}{\partial \hat{n}}|_{\Gamma}$ for the problem (P1). On applying this method, we get unique solutions to the approximating problems to (P2) and then guarantee the existence of solution to the problem (P2). The Helmholtz operator $(-\Delta - \lambda^2)$ is not a coercive operator, so we discuss the problem (P2) with the Laplace operator $(-\Delta)$. Since the continuous embeddings in Lebesgue spaces and Marcinkiewicz spaces fail in unbounded domains, we have not considered measures in the problem (1.7) and in the boundary conditions.

2 Preliminary Definitions and Properties of Boundary Layer Potentials

We will denote several constants by *C* which can only depend on Ω , *N* and independent of the indices of the sequences. The value of *C* can be different from line to line and sometimes, on the same line.

For $1 \le p \le \infty$ and a nonnegative integer k, the Sobolev space $W^{k,p}(\Omega)$ Evans (2010) is defined as

$$W^{k,p}(\Omega) = \{ u \in L^p(\Omega) : D^{\gamma} u \in L^p(\Omega), \text{ for } |\gamma| \le k \}$$

and the norm on vectors in $W^{k,p}(\Omega)$ is defined as

$$\|u\|_{W^{k,p}(\Omega)} = \sum_{|\gamma| \le k} \left\| D^{\gamma} u \right\|_{L^{p}(\Omega)}$$

where Ω is a domain in \mathbb{R}^N . We denote $W_{loc}^{k,p}(\Omega)$ to be the local Sobolev space such that for any compact $K \subset \Omega$, $u \in W^{k,p}(K)$. For $0 < \alpha < 1$, we define the Sobolev space $W^{\alpha,p}(\Omega)$ as

$$W^{\alpha,p}(\Omega) = \{ u \in L^p(\Omega) : \|u\|_{W^{\alpha,p}(\Omega)}^p = \|u\|_{L^p(\Omega)}^p + \int_{\Omega} \int_{\Omega} \frac{|u(x) - u(y)|^p}{|x - y|^{n + p\alpha}} dy dx < \infty \}.$$

Let Ω be a bounded Lipschitz domain in \mathbb{R}^N , $N \ge 3$. We now introduce the following Sobolev spaces. For $p = 2, s \in \mathbb{R}$ and $0 \le \alpha \le 1$,

- 1. $H^{s}(\mathbb{R}^{N}) = \{u : \int_{\mathbb{R}^{N}} (1 + |\xi|^{2})^{s/2} \widehat{u}(\xi) e^{i2\pi\xi \cdot x} d\xi \in L^{2}(\mathbb{R}^{N})\}, \ \widehat{u} \text{ is the Fourier transform of } u.$ This space is a separable Hilbert space.
- 2. $H^{s}(\Omega) = \{u|_{\Omega} : u \in H^{s}(\mathbb{R}^{N})\}$

- 3. $\widetilde{H}^{s}(\Omega) =$ Closure of $C_{0}^{\infty}(\Omega)$ in $H^{s}(\mathbb{R}^{N})$. For further details on these Sobolev spaces one may refer to Hsiao and Wendland (2008) Chapter 4.
- 4. $H^{\alpha}(\Gamma) = \begin{cases} \{g|_{\Gamma} : g \in H^{\frac{1}{2} + \alpha}(\Omega)\}, & (0 < \alpha \le 1) \\ L^{2}(\Gamma) & (\alpha = 0), \end{cases}$
- 5. $H^{\alpha}(\Gamma_i) = \{g|_{\Gamma_i} : g \in H^{\alpha}(\Gamma)\},\$
- 6. $\widetilde{H}^{\alpha}(\Gamma_i) = \{g|_{\Gamma_i} : g \in H^{\alpha}(\Gamma), \text{ supp}(g) \subset \overline{\Gamma}_i\}, i = 1, 2.$

Let $H^{-\alpha}(\Gamma)$ is the dual space of $H^{\alpha}(\Gamma)$, i.e. $H^{-\alpha}(\Gamma) = (H^{\alpha}(\Gamma))^*$. Equivalently, for $i = 1, 2, H^{-\alpha}(\Gamma_i) = (\widetilde{H}^{\alpha}(\Gamma_i))^*$ and $\widetilde{H}^{-\alpha}(\Gamma_i) = (H^{\alpha}(\Gamma_i))^*$.

We denote $\langle \cdot, \cdot \rangle_{\Gamma}$ as the duality pairing between $H^{\alpha}(\Gamma)$ and $H^{-\alpha}(\Gamma)$ given by $\langle f, g \rangle_{\Gamma} = \int_{\Gamma} f(z)g(z)ds_z$ for any $f \in H^{\alpha}(\Gamma)$ and $g \in H^{-\alpha}(\Gamma)$. Similarly, $\langle \cdot, \cdot \rangle_{\Gamma_i}$ is the duality pairing between $H^{\alpha}(\Gamma_i)$ and $\widetilde{H}^{-\alpha}(\Gamma_i)$ (or $H^{-\alpha}(\Gamma_i)$ and $\widetilde{H}^{\alpha}(\Gamma_i)$), i = 1, 2.

Since $\widetilde{H}^{\alpha}(\Gamma_i)$, i = 1, 2, is a reflexive space, the operator

$$J: \widetilde{H}^{\alpha}(\Gamma_i) \to (\widetilde{H}^{\alpha}(\Gamma_i))^{**} = (H^{-\alpha}(\Gamma_i))^*$$

is a bijection. Hence, for any $f' \in (H^{-\alpha}(\Gamma_i))^*$ there exists a unique $f \in \widetilde{H}^{\alpha}(\Gamma_i)$ such that J(f) = f'. For $g \in H^{-\alpha}(\Gamma_i)$ we define $\langle \langle \cdot, \cdot \rangle \rangle$ by the duality pairing between $H^{-\alpha}(\Gamma_i)$ and $(H^{-\alpha}(\Gamma_i))^*$ such that

$$\langle\langle f',g\rangle\rangle = \langle g,f\rangle_{\Gamma_i}.$$

Definition 2.1 An open set $\Omega \subset \mathbb{R}^N$ is said to be a Lipschitz domain if for each $P \in \partial \Omega$ there exist a rectangular coordinate system, (x, z) such that $x \in \mathbb{R}^{N-1}, z \in \mathbb{R}$, a neighborhood $U(P) = U \subset \mathbb{R}^N$ and a function $\varphi : \mathbb{R}^{N-1} \to \mathbb{R}$ such that

1. $|\varphi(x) - \varphi(y)| \le C|x - y|, \forall x, y \in \mathbb{R}^{N-1},$ 2. $U \cap \Omega = \{(x, z) : z > \varphi(x)\} \cap U.$

Definition 2.2 The Marcinkiewicz space denoted as $M^r(\Omega)$ (or weak $L^r(\Omega)$ space), for every $0 < r < \infty$, consists of all measurable functions $g : \Omega \to \mathbb{R}$ such that

$$m(\{x \in \Omega : |g(x)| > b\}) \le \frac{C}{b^r}, \ b > 0, C < \infty,$$

where *m* is the Lebesgue measure. In fact in the case of bounded domain Ω , for any fixed $\bar{r} > 0$ we observe $M^r(\Omega) \subset M^{\bar{r}}(\Omega)$ for $r \geq \bar{r}$. Furthermore, the embeddings

$$L^{r}(\Omega) \hookrightarrow M^{r}(\Omega) \hookrightarrow L^{r-\epsilon}(\Omega),$$
 (2.1)

is continuous for every $1 < r < \infty$ and $0 < \epsilon < r - 1$.

Definition 2.3 (*Fredholm operator*) Let X and Y are two Banach spaces and A is a bounded linear operator from X to Y. Then A is said to be a Fredholm operator if its kernel (ker(A)) and cokernel (coker(A)) = Y/Range(A) are finite dimensional.

Remark 2.4 1. The "Fredholmness" of an operator *A* ensures that Range(*A*) is closed.

2. The index of a Fredholm operator A is given by ind(A)=dim(ker(A))-dim(coker(A)).

The following two theorems are borrowed from Driver (2003) which show the relationship between a Fredholm operator and a compact operator.

Theorem 2.5 For a bounded linear operator $A : X \to Y$, the following two statements are equivalent.

- 1. A is a Fredholm operator.
- 2. A is an invertible modulo compact operators, i.e. there exist compact operators C_1, C_2 and an operator B such that $AB = I + C_1$ and $BA = I + C_2$.

Theorem 2.6 If A is a Fredholm opertor then ind(A) = 0 iff $A = A_1 + A_2$, where A_1 is an invertible operator and A_2 is a compact operator.

Definition 2.7 The space of all finite Radon measures on $\Omega \subset \mathbb{R}^N$, is denoted as $\mathcal{M}(\Omega)$. For $\mu \in \mathcal{M}(\Omega)$ we define

$$\|\mu\|_{\mathcal{M}(\Omega)} = \int_{\Omega} d|\mu|,$$

which is called the 'Total variation' norm.

We now define the weak solution of the first problem (P1).

Definition 2.8 Let *X* and *Y* are two test function spaces defined as $X = \{\varphi \in C^1(\overline{\Omega}) : \varphi|_{\Gamma_1} = 0\}$ and $Y = \{\zeta \in C_c^1(\mathbb{R}^N \setminus (\Omega)) : \zeta|_{\Gamma_1} = 0 \text{ and satisfies (1.4) and (1.5)}\}$. A function $u \in W^{1,1}(\Omega)$ is a weak solution to the problem (1.2) if it satisfies

$$\int_{\Omega} \nabla u \cdot \nabla \varphi - \int_{\Omega} \lambda^2 u \varphi = \int_{\Omega} h \varphi + \int_{\Gamma_2} f_2 \varphi, \ \forall \varphi \in X.$$

Similarly a function $u \in W_{loc}^{1,1}(\mathbb{R}^N \setminus \overline{\Omega})$ is said to be a weak solution of (1.3) if

$$\int_{\mathbb{R}^N\setminus\bar{\Omega}}\nabla u\cdot\nabla\zeta \ -\int_{\mathbb{R}^N\setminus\bar{\Omega}}\lambda^2 u\zeta = -\int_{\Gamma_2}f_2\zeta, \ \forall\zeta\in Y.$$

Remark 2.9 Hereafter, a subsequence of a sequence will be denoted by the same notation as that of the sequence. Further a solution will always refer to a weak solution.

We further we denote Φ as the fundamental solution of Helmholtz equation for $N \ge 3$ which satisfies $-\Delta \Phi - \lambda^2 \Phi = \delta$, where δ is the Dirac distribution and Φ is

$$\Phi(x, y) = \begin{cases} \frac{1}{(N-2)w_N} \frac{1}{|x-y|^{N-2}}, & \text{for } \lambda = 0\\ \frac{e^{i\lambda|x-y|}}{4\pi|x-y|}, & \text{for } \lambda \neq 0, \ N = 3\\ \frac{i}{4} \left(\frac{\lambda}{2\pi(|x-y|)}\right)^{\frac{N-2}{2}} H^{(1)}_{\frac{N-2}{2}}(\lambda|x-y|), & \text{for } \lambda \neq 0, \ N \ge 3 \end{cases}$$

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for every $x, y \in \mathbb{R}^N, x \neq y$. Here w_N is the measure of the unit sphere in \mathbb{R}^N and $H_m^{(1)}$ denotes the Hankel function of the first kind of order m. We next define, boundary layer potentials (single layer and double layer) to solve the homogeneous Helmholtz equation in \mathbb{R}^N . Let $g_1 \in H^{\alpha}(\Gamma), g_2 \in H^{-\alpha}(\Gamma)$ for some $0 \leq \alpha \leq 1$, then the single layer potential is given by,

$$v_1(x) = \mathbb{S}_{\lambda} g_2(x)$$

= $\int_{\Gamma} g_2(y) \Phi(x - y) dy, \ \forall x \in \mathbb{R}^N \setminus \Gamma$ (2.2)

and the double layer potential is by

$$v_2(x) = \mathbb{K}_{\lambda} g_1(x)$$

= $\int_{\Gamma} g_1(y) \frac{\partial}{\partial \hat{n}_y} \Phi(x - y) dy, \ \forall x \in \mathbb{R}^N \setminus \Gamma$ (2.3)

where \hat{n}_y denotes the unit outward normal to the boundary Γ . We can see that for $x \in \mathbb{R}^N \setminus \Gamma$ the above two kernels are C^{∞} functions on Γ .

If $P \in \Gamma$, then X(P) denotes a cone with vertex at P such that one component is in Ω which is denoted by $X_i(P)$ and the other is in $\mathbb{R}^N \setminus \overline{\Omega}$ denoted by $X_e(P)$.

Definition 2.10 Let $P \in \Gamma$, then we define

$$S_{\lambda}g_2(P) = \int_{\Gamma} g_2(y)\Phi(P-y)dy$$

and

$$K_{\lambda}g_1(P) = \int_{\Gamma} g_1(y) \frac{\partial}{\partial \hat{n}_y} \Phi(P-y) dy.$$

According to the Lemma 3.8 of Costabel and Stephan (1985) the boundary values of the two potentials in (2.2) and (2.3) are given by

$$v_1^i(P) = \lim_{X_i(P), x \to P} \mathbb{S}_{\lambda} g_2(x)$$

= $S_{\lambda} g_2(P),$ (2.4)

$$v_1^e(P) = \lim_{X_e(P), x \to P} \mathbb{S}_{\lambda} g_2(x)$$

= $S_{\lambda} g_2(P)$ (2.5)

and

$$v_2^i(P) = \lim_{X_i(P), x \to P} \mathbb{K}_{\lambda} g_1(x)$$
$$= \left(-\frac{1}{2}I + K_{\lambda}\right) g_1(P), \qquad (2.6)$$

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$$w_2^e(P) = \lim_{X_e(P), x \to P} \mathbb{K}_{\lambda} g_1(x)$$
$$= \left(\frac{1}{2}I + K_{\lambda}\right) g_1(P).$$
(2.7)

In case of inhomogeneous Helmholtz equation $-\Delta u - \lambda^2 u = h$ in Ω , where $h \in \widetilde{H}^{-1}(\Omega)$, the Newton potential (or Volume potential) appears in the form,

$$\mathbb{N}_{\lambda}h(x) = \int_{\Omega} \Phi(x-y)h(y)dy, \ x \in \mathbb{R}^{N}.$$

It is well known that the Newton potential $\mathbb{N}_{\lambda} : \widetilde{H}^{-1}(\Omega) \to H^{1}(\Omega)$ is a continuous map by McLean and McLean (2000); Steinbach (2007). From Nečas (2011) we know the Dirichlet trace operator, $\gamma_{D} : H^{1}(\Omega) \to H^{\frac{1}{2}}(\Gamma)$ and the Neumann trace operator, $\gamma_{N} : H^{1}(\Omega) \to H^{-\frac{1}{2}}(\Gamma)$, are continuous operators. The Dirichlet trace operator of \mathbb{N}_{λ} denoted as $\gamma_{D}\mathbb{N}_{\lambda}$ is given by

$$\psi_D \mathbb{N}_{\lambda}(h(P)) = \lim_{x \to P} \mathbb{N}_{\lambda}h(x), \ \forall P \in \Gamma.$$

Thus

$$\begin{aligned} \|\gamma_D \mathbb{N}_{\lambda}(h)\|_{H^{\frac{1}{2}}(\Gamma)} &\leq C \|\mathbb{N}_{\lambda}(h)\|_{H^{1}(\Omega)} \\ &\leq C \|h\|_{\widetilde{H}^{-1}(\Omega)} \,. \end{aligned}$$

The Neumann trace of \mathbb{N}_{λ} is denoted as $\gamma_N \mathbb{N}_{\lambda}$ and hence it satisfies

$$\|\gamma_N \mathbb{N}_{\lambda}(h)\|_{H^{-\frac{1}{2}}(\Gamma)} \leq C \|h\|_{\widetilde{H}^{-1}(\Omega)}$$

Let us fix $\alpha = \frac{1}{2}$. Consider the single layer potential $v_1(x) = \mathbb{S}_{\lambda}g_2(x)$, for $g_2 \in H^{-\frac{1}{2}}(\Gamma)$. Then v_1 solves the Helmholtz equation in $\mathbb{R}^N \setminus \Gamma$. Thus $v_1 \in H^1(\Omega)$ for (IP1), $v_1 \in H^1_{loc}(\mathbb{R}^N \setminus \overline{\Omega})$ for (EP1) and satisfies (1.4)–(1.5) at infinity. We now define the ouward normal derivative of v_1 , i.e. $\frac{\partial v_1}{\partial \hat{n}}$ that belongs to $H^{-\frac{1}{2}}(\Gamma)$. Let us choose $h_1, h_2 \in H^{\frac{1}{2}}(\Gamma)$. We will denote h_1^* , h_2^* to be the extensions of h_1 , h_2 respectively such that

$$\|h_1^*\|_{H^1(\Omega)} \le C \|h_1\|_{H^{\frac{1}{2}}(\Gamma)}, \ \|h_2^*\|_{H^1(\mathbb{R}^N \setminus \bar{\Omega})} \le C \|h_2\|_{H^{\frac{1}{2}}(\Gamma)}$$
(2.8)

for some constant C > 0 which does not depend on h_1 and h_2 by Jonsson and Wallin (1978). Define

$$\left\langle \frac{\partial v_1}{\partial \hat{n}}, h_1 \right\rangle_{\Gamma} = \int_{\Omega} \nabla v_1 \cdot \nabla h_1^* - \int_{\Omega} \lambda^2 v_1 h_1^*, \left\langle \frac{\partial v_1}{\partial \hat{n}}, h_2 \right\rangle_{\Gamma} = -\int_{\mathbb{R}^N \setminus \bar{\Omega}} \nabla v_1 \cdot \nabla h_2^* + \int_{\mathbb{R}^N \setminus \bar{\Omega}} \lambda^2 v_1 h_2^*.$$

$$(2.9)$$

We have from Costabel and Stephan (1985) that for every $P \in \Gamma$,

$$\frac{\partial v_1^l(P)}{\partial \hat{n}} = \lim_{X_i(P), x \to P} \frac{\partial v_1(x)}{\partial \hat{n}}$$

$$= \left(\frac{1}{2}I + K_{\lambda}^*\right)g_2(P)$$
(2.10)

and

$$\frac{\partial v_1^e(P)}{\partial \hat{n}} = \lim_{X_e(P), x \to P} \frac{\partial v_1(x)}{\partial \hat{n}}$$
$$= \left(-\frac{1}{2}I + K_{\lambda}^* \right) g_2(P)$$
(2.11)

where K_{λ}^* is the adjoint operator of K_{λ} defined as

$$K_{\lambda}^* g_1(P) = \int_{\Gamma} \frac{\partial}{\partial \hat{n}_P} \Phi(P - y) g_1(y) dy.$$

Similarly, in case of double layer potential $v_2(x) = \mathbb{K}_{\lambda}g_1(x)$ for $g_1 \in H^{\frac{1}{2}}(\Gamma)$, $\frac{\partial v_2}{\partial \hat{n}} \in H^{-\frac{1}{2}}(\Gamma)$ which satisfies (2.8) and (2.9). Let us define an operator $D_{\lambda} : H^{\frac{1}{2}}(\Gamma) \to H^{-\frac{1}{2}}(\Gamma)$ as in Costabel and Stephan (1985) such that for every $P \in \Gamma$,

$$D_{\lambda}g_{1}(P) = \frac{\partial}{\partial \hat{n}_{P}} \mathbb{K}_{\lambda}g_{1}(P)$$
(2.12)

and

$$\lim_{X_i(P), x \to P} \frac{\partial}{\partial \hat{n}_x} \mathbb{K}_{\lambda} g_1(x) = \lim_{X_e(P), x \to P} \frac{\partial}{\partial \hat{n}_x} \mathbb{K}_{\lambda} g_1(x)$$
$$= D_{\lambda} g_1(P).$$
(2.13)

Lemma 2.11 The operators
1.
$$S_{\lambda} : H^{-\frac{1}{2}}(\Gamma) \to H^{\frac{1}{2}}(\Gamma),$$

2. $\left(\pm \frac{1}{2}I + K_{\lambda}\right) : H^{\frac{1}{2}}(\Gamma) \to H^{\frac{1}{2}}(\Gamma),$
3. $\left(\pm \frac{1}{2}I + K_{\lambda}^{*}\right) : H^{-\frac{1}{2}}(\Gamma) \to H^{-\frac{1}{2}}(\Gamma),$
4. $D_{\lambda} : H^{\frac{1}{2}}(\Gamma) \to H^{-\frac{1}{2}}(\Gamma)$
are continuous by Costabel and Stephan (1985).

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2.1 Derivation of Representation Formulae

Let Ω be a bounded Lipschitz domain in \mathbb{R}^N and $\overline{\Gamma}_1 \cup \Gamma_2 = \Gamma_1 \cup \overline{\Gamma}_2 = \Gamma, \Gamma_1 \cap \Gamma_2 = \emptyset$. For $g_1 \in \widetilde{H}^{\frac{1}{2}}(\Gamma_i)$, i = 1, 2, we denote the zero extension function \widetilde{g}_1 of g_1 by

$$\widetilde{g}_1 = \begin{cases} g_1 & in \ \Gamma_i \\ 0 & in \ \Gamma \setminus \Gamma_i, \ i = 1, 2 \end{cases}$$

Clearly, $\tilde{g}_1 \in H^{\frac{1}{2}}(\Gamma)$. Similarly, for $g_2 \in \tilde{H}^{-\frac{1}{2}}(\Gamma_i)$, i = 1, 2, we extend g_2 to a function $\tilde{g}_2 \in H^{-\frac{1}{2}}(\Gamma)$. We now introduce the following operators.

$$\begin{split} S_{ij} &: \widetilde{H}^{-\frac{1}{2}}(\Gamma_i) \to H^{\frac{1}{2}}(\Gamma_j), \ S_{ij}g_2 = S_{\lambda}\widetilde{g}_2|_{\Gamma_j} \text{ for } g_2 \in \widetilde{H}^{-\frac{1}{2}}(\Gamma_i), \\ K_{ij} &: \widetilde{H}^{\frac{1}{2}}(\Gamma_i) \to H^{\frac{1}{2}}(\Gamma_j), \ K_{ij}g_1 = K_{\lambda}\widetilde{g}_1|_{\Gamma_j} \text{ for } g_1 \in \widetilde{H}^{\frac{1}{2}}(\Gamma_i), \\ K_{ij}^* &: \widetilde{H}^{-\frac{1}{2}}(\Gamma_i) \to H^{-\frac{1}{2}}(\Gamma_j), \ K_{ij}^*g_2 = K_{\lambda}^*\widetilde{g}_2|_{\Gamma_j} \text{ for } g_2 \in \widetilde{H}^{-\frac{1}{2}}(\Gamma_i), \\ D_{ij} &: \widetilde{H}^{\frac{1}{2}}(\Gamma_i) \to H^{-\frac{1}{2}}(\Gamma_j), \ D_{ij}g_1 = D_{\lambda}\widetilde{g}_1|_{\Gamma_j} \text{ for } g_1 \in \widetilde{H}^{\frac{1}{2}}(\Gamma_i). \end{split}$$

Let us consider $u \in H^1(\Omega)$ be a solution to the Helmholtz equation $-\Delta u - \lambda^2 u = h$ in Ω and $u \in H^1_{loc}(\mathbb{R}^N \setminus \overline{\Omega})$ satisfies $-\Delta u - \lambda^2 u = 0$ in $\mathbb{R}^N \setminus \overline{\Omega}$ along with (1.4)-(1.5). On using Lemma 3.1 of Costabel and Stephan (1985), we obtain the Green's second identity

$$\int_{\Omega} u \Delta v - v \Delta u = \int_{\Gamma} u \frac{\partial v}{\partial \hat{n}} - v \frac{\partial u}{\partial \hat{n}}$$

When we replace v with Φ , the fundamental solution of Helmholtz equation, we obtain the following.

$$\int_{\Omega} u(y) \Delta \Phi(x, y) - \Delta u(y) \Phi(x, y) = \int_{\Gamma} u(y) \frac{\partial \Phi(x, y)}{\partial \hat{n}} - \Phi(x, y) \frac{\partial u(y)}{\partial \hat{n}}$$
$$u(x) = \int_{\Omega} \Phi(x, y) h(y) - \int_{\Gamma} u(y) \frac{\partial \Phi(x, y)}{\partial \hat{n}}$$
$$- \Phi(x, y) \frac{\partial u(y)}{\partial \hat{n}} . \tag{2.14}$$

Let $B_r = \{z \in \mathbb{R}^N : |z| = r\}$ and $D_r = \{x \in \mathbb{R}^N \setminus \overline{\Omega} : |x| < r\}$. On applying the Green's second identity in the domain D_r we get

$$u(x) = -\int_{D_r} u(y)\Delta\Phi(x, y) - \Delta u(y)\Phi(x, y)$$

= $-\int_{B_r} u(y)\frac{\partial\Phi(x, y)}{\partial\hat{n}} - \Phi(x, y)\frac{\partial u(y)}{\partial\hat{n}} + \int_{\Gamma} u(y)\frac{\partial\Phi(x, y)}{\partial\hat{n}} - \Phi(x, y)\frac{\partial u(y)}{\partial\hat{n}}$
(2.15)

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On passing the limit $r \to \infty$ and by using (1.4)–(1.5) we see that

$$\int_{B_r} u(y) \frac{\partial \Phi(x, y)}{\partial \hat{n}} - \Phi(x, y) \frac{\partial u(y)}{\partial \hat{n}} \to 0.$$

Let us denote the Cauchy data as $(\phi, \psi) \in H^{\frac{1}{2}}(\Gamma) \times H^{-\frac{1}{2}}(\Gamma)$, where $u|_{\Gamma} = \phi$ and $\frac{\partial u}{\partial \hat{n}}|_{\Gamma} = \psi$. On combining (2.14) and (2.15), we can express u as

$$u(x) = \begin{cases} \mathbb{N}_{\lambda}h(x) - \mathbb{K}_{\lambda}\phi(x) + \mathbb{S}_{\lambda}\psi(x), & \text{if } x \in \Omega\\ \mathbb{K}_{\lambda}\phi(x) - \mathbb{S}_{\lambda}\psi(x), & \text{if } x \in \mathbb{R}^{N} \setminus \bar{\Omega}. \end{cases}$$
(2.16)

Consider (P1), with the boundary data $u|_{\Gamma_1} = f_1 \in H^{\frac{1}{2}}(\Gamma_1)$ and $\frac{\partial u}{\partial \hat{n}}|_{\Gamma_2} = f_2 \in H^{-\frac{1}{2}}(\Gamma_2)$. For simplicity, we restrict ourselves to the interior mixed boundary value problem (1.2). Obviously the corresponding results for the exterior problem (1.3) are obtained by only slight modifications. Furthermore, we say \mathring{f}_1 , \mathring{f}_2 are the extensions of f_1 and f_2 respectively, which satisfy

$$\left\| \mathring{f}_{1} \right\|_{H^{\frac{1}{2}}(\Gamma)} \le C \left\| f_{1} \right\|_{H^{\frac{1}{2}}(\Gamma_{1})}$$
(2.17)

and

$$\left\| f_{2}^{*} \right\|_{H^{-\frac{1}{2}}(\Gamma)} \leq C \left\| f_{2} \right\|_{H^{-\frac{1}{2}}(\Gamma_{2})}.$$
(2.18)

The above extension is possible since we know $\partial \Gamma_1 = \partial \Gamma_2$ and Γ is Lipschitz [3]. Let us define $\phi = \mathring{f}_1 + \widetilde{g}_1$ and $\psi = \mathring{f}_2 + \widetilde{g}_2$, where \widetilde{g}_1 and \widetilde{g}_2 are arbitrary functions in $H^{\frac{1}{2}}(\Gamma)$ and $H^{-\frac{1}{2}}(\Gamma)$, respectively. Here \widetilde{g}_1 is the zero extension of $g_1 \in \widetilde{H}^{\frac{1}{2}}(\Gamma_2)$ and \widetilde{g}_2 is the zero extension of $g_2 \in \widetilde{H}^{-\frac{1}{2}}(\Gamma_1)$. The representation (2.16) is used to express the solutions of problem (1.2) as

$$u(x) = \mathbb{N}_{\lambda}h(x) - \mathbb{K}_{\lambda}(\mathring{f}_1 + \widetilde{g}_1)(x) + \mathbb{S}_{\lambda}(\mathring{f}_2 + \widetilde{g}_2)(x).$$
(2.19)

On restricting the Eq. (2.19) to Γ we get,

$$\mathring{f}_1 + \widetilde{g}_1 = \gamma_D \mathbb{N}_{\lambda} h - \left(-\frac{1}{2}I + K_{\lambda}\right)(\mathring{f}_1 + \widetilde{g}_1) - S_{\lambda}(\mathring{f}_2 + \widetilde{g}_2).$$

On Γ_1 we have the following,

$$f_{1} = \gamma_{D} \mathbb{N}_{\lambda} h|_{\Gamma_{1}} - \left(-\frac{1}{2}I + K_{\lambda}\right) (\mathring{f}_{1} + \widetilde{g}_{1})\Big|_{\Gamma_{1}} + S_{\lambda}(\mathring{f}_{2} + \widetilde{g}_{2})\Big|_{\Gamma_{1}}$$

$$= \gamma_{D} \mathbb{N}_{\lambda} h|_{\Gamma_{1}} - K_{21}g_{1} - \left(-\frac{1}{2}I + K_{\lambda}\right) \mathring{f}_{1}\Big|_{\Gamma_{1}} + S_{11}g_{2} + S_{\lambda}\mathring{f}_{2}\Big|_{\Gamma_{1}}$$

$$K_{21}g_{1} - S_{11}g_{2} = -f_{1} + \gamma_{D} \mathbb{N}_{\lambda} h|_{\Gamma_{1}} - \left(-\frac{1}{2}I + K_{\lambda}\right) \mathring{f}_{1}\Big|_{\Gamma_{1}} + S_{\lambda}\mathring{f}_{2}\Big|_{\Gamma_{1}}$$

$$= F^{*}(f_{1}, f_{2}, h) \text{ (say).}$$
(2.20)

Taking the Neumann trace of (2.19) we have

$$\mathring{f}_2 + \widetilde{g}_2 = \gamma_N \mathbb{N}_\lambda h - D_\lambda (\mathring{f}_1 + \widetilde{g}_1) + \left(\frac{1}{2}I + K_\lambda^*\right) (\mathring{f}_2 + \widetilde{g}_2).$$

Similarly on Γ_2 ,

$$f_{2} = \gamma_{N} \mathbb{N}_{\lambda} h|_{\Gamma_{2}} - D_{\lambda} (\mathring{f}_{1} + \widetilde{g}_{1})|_{\Gamma_{2}} + \left(\frac{1}{2}I + K_{\lambda}^{*}\right) (\mathring{f}_{2} + \widetilde{g}_{2})|_{\Gamma_{2}}$$

$$= \gamma_{N} \mathbb{N}_{\lambda} h|_{\Gamma_{2}} - D_{22}g_{1} - D\mathring{f}_{1}|_{\Gamma_{2}} + K_{12}^{*}g_{2} + \left(\frac{1}{2}I + K_{\lambda}^{*}\right) \mathring{f}_{2}|_{\Gamma_{2}}$$

$$D_{22}g_{1} - K_{12}^{*}g_{2} = -f_{2} + \gamma_{N} \mathbb{N}_{\lambda} h|_{\Gamma_{2}} - D_{\lambda} \mathring{f}_{1}|_{\Gamma_{2}} + \left(\frac{1}{2}I + K_{\lambda}^{*}\right) \mathring{f}_{2}|_{\Gamma_{2}}$$

$$= G^{*}(f_{1}, f_{2}, h) \text{ (say).}$$
(2.21)

Clearly $F^* \in H^{\frac{1}{2}}(\Gamma_1)$ and $G^* \in H^{-\frac{1}{2}}(\Gamma_2)$. Combining Eqs. (2.20) and (2.21) we get

$$\begin{pmatrix} K_{21} & -S_{11} \\ D_{22} & -K_{12}^* \end{pmatrix} \begin{pmatrix} g_1 \\ g_2 \end{pmatrix} = \begin{pmatrix} F^* \\ G^* \end{pmatrix}$$

We now define a matrix operator A as

$$A = \begin{pmatrix} K_{21} & -S_{11} \\ D_{22} & -K_{12}^* \end{pmatrix}$$

where, $A: \widetilde{H}^{\frac{1}{2}}(\Gamma_2) \times \widetilde{H}^{-\frac{1}{2}}(\Gamma_1) \to H^{\frac{1}{2}}(\Gamma_1) \times H^{-\frac{1}{2}}(\Gamma_2).$

2.2 Invertibility of Layer Potentials

For the homogeneous Helmholtz equation with $\lambda = 0$, the two boundary layer operators $S_0 : H^{-\frac{1}{2}}(\Gamma) \to H^{\frac{1}{2}}(\Gamma)$ and $\left(-\frac{1}{2}I + K_0\right) : H^{\frac{1}{2}}(\Gamma) \to H^{\frac{1}{2}}(\Gamma)$ are bijective by Chang and Choe (2008).

Proposition 2.12 *This Proposition is from* Torres and Welland (1993) *which concludes that for* $Im(\lambda) > 0$

1. $S_{\lambda} : L^{2}(\Gamma) \to H^{1}(\Gamma)$ is invertible. 2. $\left(\pm \frac{1}{2}I + K_{\lambda}\right) : L^{2}(\Gamma) \to L^{2}(\Gamma)$ is invertible. 3. $\left(\pm \frac{1}{2}I + K_{\lambda}^{*}\right) : L^{2}(\Gamma) \to L^{2}(\Gamma)$ is invertible.

Theorem 2.13 Let $Im(\lambda) > 0$. Then $D_{\lambda} : H^1(\Gamma) \to L^2(\Gamma)$ is an invertible operator.

Proof Let us consider a $g \in L^2(\Gamma)$. From the above Proposition 2.12, $(\frac{1}{2}I + K_{\lambda}^*)$ is bijective from $L^2(\Gamma)$ to itself. Hence, there exists a $g' \in L^2(\Gamma)$ such that $(\frac{1}{2}I + K_{\lambda}^*)g' = g$.

Let $v(x) = S_{\lambda} \left(-\frac{1}{2}I + K_{\lambda}^*\right)^{-1} g'(x)$. Then v satisfies the homogenous Helmholtz equation in $\mathbb{R}^N \setminus \overline{\Omega}$. Using the properties (2.5), (2.11) and the decay conditions at infinity (1.4)–(1.5) in the exterior domain we have the following representation for v.

$$v(x) = \mathbb{K}_{\lambda} f(x) - \mathbb{S}_{\lambda} g'(x)$$

where, $f = S_{\lambda}(-\frac{1}{2}I + K_{\lambda}^*)^{-1}g' \in H^1(\Gamma)$. Taking the Neumann trace of v we get

$$g' = D_{\lambda}f - \left(-\frac{1}{2}I + K_{\lambda}^*\right)g'$$

which implies

$$D_{\lambda}f = \left(\frac{1}{2}I + K_{\lambda}^{*}\right)g'.$$
(2.22)

Hence, for any $g \in L^2(\Gamma)$, there exists $f \in H^1(\Gamma)$ such that $D_{\lambda}f = (\frac{1}{2}I + K_{\lambda}^*)g' = g$.

Claim: D_{λ} is injective.

Suppose there exists $f \in H^1(\Gamma)$ such that $D_{\lambda} f = 0$ on Γ . We write $v(x) = \mathbb{K}_{\lambda} f(x)$, for all x in $\mathbb{R}^N \setminus \Gamma$. Hence, $v \in H^1(\Omega)$ is a solution of $-\Delta v - \lambda^2 v = 0$ in Ω and $v \in H^1_{loc}(\mathbb{R}^N \setminus \overline{\Omega})$ satisfies $-\Delta v - \lambda^2 v = 0$ in $\mathbb{R}^N \setminus \overline{\Omega}$ along with (1.4)–(1.5). From the Eqs. (2.6)–(2.7) we get

$$v^{i} - v^{e} = \left(-\frac{1}{2}I + K_{\lambda}\right)f - \left(\frac{1}{2}I + K_{\lambda}\right)f$$
$$= -f$$

and from Torres and Welland (1993) we have

$$D_{\lambda}f(P) = \lim_{X_{i}(P), x \to P} \frac{\partial}{\partial \hat{n}}v(x)$$

=
$$\lim_{X_{e}(P), x \to P} \frac{\partial}{\partial \hat{n}}v(x).$$
 (2.23)

Thus,

$$0 = \langle D_{\lambda}f, -f \rangle_{\Gamma}$$

= $\langle D_{\lambda}f, v^{i} \rangle_{\Gamma} - \langle D_{\lambda}f, v^{e} \rangle_{\Gamma}$
= $\int_{\mathbb{R}^{N}} |\nabla v|^{2} - \int_{\mathbb{R}^{N}} \lambda^{2} |v|^{2}$ (2.24)

where the last term in (2.24) is due to the fact that v is a solution to the homogeneous Helmholtz equation. As per our assumption λ^2 is not an eigenvalue of $(-\Delta)$. Hence, using the conditions (1.4)–(1.5) we have v = 0 a.e. in \mathbb{R}^N . Since v is continuous across the boundary, we have $-f = v^i - v^e = 0$. This implies f = 0 on Γ . So, D_{λ} is injective.

Remark 2.14 The operators S_{λ} and D_{λ} are self-adjoint operators, i.e. $S_{\lambda} = S_{\lambda}^*$, $D_{\lambda} = D_{\lambda}^*$ [refer Lemma 3.9(a) of Costabel and Stephan (1985)], where S_{λ}^* , D_{λ}^* are the adjoint operators of S_{λ} , D_{λ} respectively. Hence, using Proposition 2.12, Theorem 2.13 we obtain $S_{\lambda}^* : H^{-1}(\Gamma) \to L^2(\Gamma)$ and $D_{\lambda}^* : L^2(\Gamma) \to H^{-1}(\Gamma)$ are invertible operators. Using the properties of real interpolation from Appendix B (Theorem B.2) of McLean and McLean (2000) on S_{λ} , D_{λ} we have

1. $S_{\lambda}: H^{-\frac{1}{2}}(\Gamma) \to H^{\frac{1}{2}}(\Gamma),$ 2. $D_{\lambda}: H^{\frac{1}{2}}(\Gamma) \to H^{-\frac{1}{2}}(\Gamma)$

are invertible operators.

3 Existence and Uniqueness Results of (P1)

Theorem 3.1 Let $\Gamma_1 \subset \Gamma$, then $S_{11} : \widetilde{H}^{-\frac{1}{2}}(\Gamma_1) \to H^{\frac{1}{2}}(\Gamma_1)$ is a bijective operator.

Proof We break the proof into three steps.

Step 1. *The operator* S_{11} *is injective.*

Assume that there exists $g_2 \in \widetilde{H}^{-\frac{1}{2}}(\Gamma_1)$ such that $S_{11}g_2 = 0$ on Γ_1 . We write $v_1(x) = S_{\lambda}\widetilde{g}_2(x)$, for all $x \in \mathbb{R}^N \setminus \Gamma$, where $\widetilde{g}_2 \in H^{-\frac{1}{2}}(\Gamma)$ is the zero extension of g_2 . Hence, from the Eqs. (2.4)–(2.5) we have $v_1^i = v_1^e \in H^{\frac{1}{2}}(\Gamma)$ and from (2.10), $\frac{\partial v_1^i}{\partial \hat{n}} - \frac{\partial v_1^e}{\partial \hat{n}} = \widetilde{g}_2$. On replacing h_1, h_2 with v_1^i, v_1^e , respectively in the Eq. (2.9) we have

$$0 = \langle g_2, S_{11}g_2 \rangle_{\Gamma_1}$$

$$= \langle \tilde{g}_2, S_\lambda \tilde{g}_2 \rangle_{\Gamma}$$

$$= \left\langle \frac{\partial v_1^i}{\partial \hat{n}}, v_1^i \right\rangle_{\Gamma} - \left\langle \frac{\partial v_1^e}{\partial \hat{n}}, v_1^e \right\rangle_{\Gamma}$$

$$= \int_{\mathbb{R}^N} |\nabla v_1|^2 - \int_{\mathbb{R}^N} \lambda^2 |v_1|^2.$$
(3.1)

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Thus, on using the conditions (1.4)–(1.5) we conclude that $v_1 = 0$ a.e. in \mathbb{R}^N . By the continuity of v_1 on Γ , we have $\tilde{g}_2 = \frac{\partial v_1^i}{\partial \hat{n}} - \frac{\partial v_1^e}{\partial \hat{n}} = 0$. This implies $g_2 = 0$ in Γ_1 and hence S_{11} is injective.

Step 2. *S*₁₁ *is bounded below.*

Suppose there exists a sequence $(g_2^n) \in \widetilde{H}^{-\frac{1}{2}}(\Gamma_1)$ such that $S_{11}g_2^n \to f$, for some $f \in H^{\frac{1}{2}}(\Gamma_1)$.

Case 1. Assume that (g_2^n) is a bounded sequence in $\widetilde{H}^{-\frac{1}{2}}(\Gamma_1)$. Therefore, there exists a subsequence (g_2^n) and $g_2 \in \widetilde{H}^{-\frac{1}{2}}(\Gamma_1)$ such that (g_2^n) converges weakly to g_2 , i.e. $g_2^n \xrightarrow{w} g_2$ in $\widetilde{H}^{-\frac{1}{2}}(\Gamma_1)$. Let $l \in \widetilde{H}^{-\frac{1}{2}}(\Gamma_1)$. Then we have

$$\langle l, f \rangle_{\Gamma_1} = \langle l, \lim_{n \to \infty} S_{11} g_2^n \rangle_{\Gamma_1}$$

$$= \lim_{n \to \infty} \langle l, S_{11} g_2^n \rangle_{\Gamma_1}$$

$$= \lim_{n \to \infty} \langle \langle S_{11}^* l, g_2^n \rangle \rangle$$

$$= \langle \langle S_{11}^* l, g_2 \rangle \rangle$$

$$= \langle l, S_{11} g_2 \rangle_{\Gamma_1}.$$

Since every reflexive space has a unique predual, hence $S_{11}g_2 = f$. Therefore, S_{11} has a closed range.

Case 2. Assume that (g_2^n) is an unbounded sequence in $\widetilde{H}^{-\frac{1}{2}}(\Gamma_1)$. Let us denote

$$G^{n} = \frac{g_{2}^{n}}{\|g_{2}^{n}\|_{\widetilde{H}^{-\frac{1}{2}}(\Gamma_{1})}}$$

Hence, $\|G^n\|_{\widetilde{H}^{-\frac{1}{2}}(\Gamma_1)} = 1$. Therefore, there exists a subsequence (G^n) and $G \in \widetilde{H}^{-\frac{1}{2}}(\Gamma_1)$ such that $G^n \xrightarrow{w} G$ in $\widetilde{H}^{-\frac{1}{2}}(\Gamma_1)$. Since $S_{11}g_2^n \to f$ and $\|g_2^n\|_{\widetilde{H}^{-\frac{1}{2}}(\Gamma_1)} \to \infty$, we have $S_{11}G^n \to 0$ in $H^{\frac{1}{2}}(\Gamma_1)$. From Case 1 it easily follows that $S_{11}G = 0$, which further implies G = 0 by the injectivity of S_{11} . Using the invertibility of S_{λ} (refer Remark 2.14) we obtain

$$1 = \left\| G^{n} \right\|_{\widetilde{H}^{-\frac{1}{2}}(\Gamma_{1})}$$

$$\leq \left\| \widetilde{G}^{n} \right\|_{H^{-\frac{1}{2}}(\Gamma)}$$

$$\leq C \left\| S_{\lambda}(\widetilde{G}^{n}) \right\|_{H^{\frac{1}{2}}(\Gamma)} \text{ (for } C > 0\text{).}$$

$$(3.2)$$

We know that $S_{11}G^n = S_{\lambda}(\widetilde{G}^n)|_{\Gamma_1}$ and $S_{12}G^n = S_{\lambda}(\widetilde{G}^n)|_{\Gamma_2}$. For $x \neq y$, $\Phi(x-y)$ is a C^{∞} function. This implies $S_{12}G^n \to 0$ in $H^{\frac{1}{2}}(\Gamma_2)$, since $G^n \stackrel{w}{\to} 0$ in $\widetilde{H}^{-\frac{1}{2}}(\Gamma_1)$. Hence, $S_{\lambda}(\widetilde{G}^n) \to 0$ in $H^{\frac{1}{2}}(\Gamma)$, which is a contradiction to (3.2). Therefore, we conclude that S_{11} has closed range. Thus, S_{11} is bounded below since S_{11} is injective and its range is closed.

Step 3. *S*₁₁ *has dense range.*

Assume that $S_{11}^*g_2 = 0$ for some $g_2 \in \widetilde{H}^{-\frac{1}{2}}(\Gamma_1)$. Hence, for any $l \in \widetilde{H}^{-\frac{1}{2}}(\Gamma_1)$,

$$0 = \langle \langle S_{11}^* g_2, l \rangle \rangle$$
$$= \langle g_2, S_{11} l \rangle_{\Gamma_1}.$$

Choose $l = g_2$. Then by proceeding on similar lines as in step 1 we get $g_2 = 0$. Since Kernel $(S_{11}^*) = \text{Range}(S_{11})^{\perp} = \overline{\text{Range}(S_{11})}^{\perp}$, the injectivity of S_{11}^* implies S_{11} has dense range.

Combining the results from the above three steps we conclude that the operator $S_{11}: \widetilde{H}^{-\frac{1}{2}}(\Gamma_1) \to H^{\frac{1}{2}}(\Gamma_1)$ is bijective.

Theorem 3.2 Let $\Gamma_2 \subset \Gamma$, then the operator $D_{22} : \widetilde{H}^{\frac{1}{2}}(\Gamma_2) \to H^{-\frac{1}{2}}(\Gamma_2)$ is invertible. **Proof** Similar to the steps in Theorem 3.1, we will show that D_{22} is injective and bounded below with a dense range. Assume that there exists $g_1 \in \widetilde{H}^{\frac{1}{2}}(\Gamma_2)$ such that $D_{22}g_1 = 0$ on Γ_2 . We now express $v_2(x) = \mathbb{K}_{\lambda}\widetilde{g}_1(x)$, for all $x \in \mathbb{R}^N \setminus \Gamma$. From the Eqs. (2.6) and (2.7) we get

$$v_2^i - v_2^e = \left(-\frac{1}{2}I + K_\lambda\right)\widetilde{g}_1 - \left(\frac{1}{2}I + K_\lambda\right)\widetilde{g}_1$$
$$= -\widetilde{g}_1.$$

Thus,

$$D = \langle D_{22}g_1, -g_1 \rangle_{\Gamma_2}$$

= $\langle D_{\lambda}\widetilde{g}_1, -\widetilde{g}_1 \rangle_{\Gamma}$
= $\left\langle D_{\lambda}\widetilde{g}_1, v_2^i \right\rangle_{\Gamma} - \left\langle D_{\lambda}\widetilde{g}_1, v_2^e \right\rangle_{\Gamma}$ (from the Eq. (2.13))
= $\int_{\mathbb{R}^N} |\nabla v_2|^2 - \int_{\mathbb{R}^N} \lambda^2 |v_2|^2.$ (3.3)

Hence, using the conditions (1.4)–(1.5) we have $v_2 = 0$ a.e. in \mathbb{R}^N , since λ^2 is not an eigenvalue of $(-\Delta)$. By the continuity of v_2 in $x \in \mathbb{R}^N \setminus \Gamma_2$ we have $v_2^i - v_2^e = -\widetilde{g}_1 = 0$. This implies $g_1 = 0$ in Γ_2 . So, D_{22} is injective.

On using arguments from Theorem 3.1, we can show that D_{22} has a closed range and hence it is bounded below. We suppose that $D_{22}^*g'_1 = 0$ for some $g_1 \in \widetilde{H}^{\frac{1}{2}}(\Gamma_2)$. Then for $f \in \widetilde{H}^{\frac{1}{2}}(\Gamma_2)$,

$$0 = \langle -D_{22}^*g'_1, f \rangle_{\Gamma_2}$$

= $\langle \langle -g'_1, D_{22}f \rangle \rangle$
= $\langle -g_1, D_{22}f \rangle_{\Gamma_2}.$

Taking $f = g_1$, then from (3.3) we obtain $g_1 = 0$ in Γ_2 . Hence, D_{22}^* is injective which implies D_{22} has dense range. Therefore, D_{22} is an invertible operator.

Theorem 3.3 The matrix operator $A : \widetilde{H}^{\frac{1}{2}}(\Gamma_2) \times \widetilde{H}^{-\frac{1}{2}}(\Gamma_1) \to H^{\frac{1}{2}}(\Gamma_1) \times H^{-\frac{1}{2}}(\Gamma_2)$ is invertible.

Proof For any $g_1 \in \widetilde{H}^{\frac{1}{2}}(\Gamma_2)$ and $P \in \Gamma_1$, the operator $K_{21} : \widetilde{H}^{\frac{1}{2}}(\Gamma_2) \to H^{\frac{1}{2}}(\Gamma_1)$ is defined as

$$K_{21}g_1(P) = K_{\lambda}\widetilde{g}_1(P)$$

= $\int_{\Gamma} \frac{\partial}{\partial \widehat{n}_y} \Phi(P, y)\widetilde{g}_1(y)dy$
= $\int_{\Gamma_2} \frac{\partial}{\partial \widehat{n}_y} \Phi(P, y)g_1(y)dy.$ (3.4)

We can see that the kernel $\frac{\partial \Phi(P,y)}{\partial \hat{n}_y}$ in (3.4) is a C^{∞} function. Let (g_1^n) be a bounded sequence in $\widetilde{H}^{\frac{1}{2}}(\Gamma_2)$, then there exists a subsequence (g_1^n) and g_1 in $\widetilde{H}^{\frac{1}{2}}(\Gamma_2)$ such that (g_1^n) converges weakly to g_1 . Hence,

$$\lim_{n \to \infty} K_{21} g_1^n(P) = \lim_{n \to \infty} \int_{\Gamma_2} \frac{\partial}{\partial \hat{n}_y} \Phi(P, y) g_1^n(y) dy$$
$$= \int_{\Gamma_2} \frac{\partial}{\partial \hat{n}_y} \Phi(P, y) g_1(y) dy$$
$$= K_{21} g_1(P).$$

Thus, K_{21} is a compact operator. Similarly we can show that the operator K_{12}^* is also compact.

We have

$$A = \begin{pmatrix} K_{21} & -S_{11} \\ D_{22} & -K_{12}^* \end{pmatrix}$$
$$= \begin{pmatrix} K_{21} & 0 \\ 0 & -K_{12}^* \end{pmatrix} + \begin{pmatrix} 0 & -S_{11} \\ D_{22} & 0 \end{pmatrix}$$
$$= A_1 + A_2$$

where $A_1 = \begin{pmatrix} K_{21} & 0 \\ 0 & -K_{12}^* \end{pmatrix}$ and $A_2 = \begin{pmatrix} 0 & -S_{11} \\ D_{22} & 0 \end{pmatrix}$. The matrix A_2 is invertible, since S_{11} and D_{22} are invertible operators by Theorems 3.1 and Theorem 3.2 respectively. As the operators D_{λ} and S_{λ} are also continuous by Costabel and Stephan (1985), the inverse of A_2 , i.e. A_2^{-1} is also bounded. We know the operators K_{21} and K_{12}^* are compact operators and hence A_1 is also a compact operator.

compact operators and hence A_1 is also a compact operator. Thus, we can write $A_2^{-1}A = A_2^{-1}A_1 + I = C_1 + I$ and $AA_2^{-1} = A_1A_2^{-1} + I = C_2 + I$ where C_1 , C_2 are compact operators. Using Theorem 2.5, it is equivalent to say that A is a Fredholm operator. This implies ind(A) = 0 (by Theorem 2.6). Now to show A is bijective it is sufficient to show A is injective, i.e. dim ker(A) = 0. *Claim:* A is injective. Let us assume that there exist some $g_1 \in \widetilde{H}^{\frac{1}{2}}(\Gamma_2)$ and $g_2 \in \widetilde{H}^{-\frac{1}{2}}(\Gamma_1)$ such that $A(g_1, g_2) = 0$. Now for $x \in \mathbb{R}^N \setminus \Gamma$, we write

$$v(x) = \begin{cases} \mathbb{S}_{\lambda} \widetilde{g}_{2}(x) - \mathbb{K}_{\lambda} \widetilde{g}_{1}(x), & \text{if } x \in \Omega \\ -\mathbb{S}_{\lambda} \widetilde{g}_{2}(x) + \mathbb{K}_{\lambda} \widetilde{g}_{1}(x), & \text{if } x \in \mathbb{R}^{N} \setminus \bar{\Omega} \end{cases}$$

Then v satisfies the following problems

$$-\Delta v - \lambda^2 v = 0 \text{ in } \Omega,$$

$$v = 0 \text{ on } \Gamma_1,$$

$$\frac{\partial v}{\partial \hat{n}} = 0 \text{ on } \Gamma_2,$$

$$v \in H^1(\Omega)$$
(3.5)

and

$$-\Delta v - \lambda^2 v = 0 \text{ in } \mathbb{R}^N \setminus \bar{\Omega},$$

$$v = 0 \text{ on } \Gamma_1,$$

$$\frac{\partial v}{\partial \hat{n}} = 0 \text{ on } \Gamma_2,$$

$$v \in H^1_{loc}(\mathbb{R}^N \setminus \bar{\Omega}).$$
(3.6)

This implies $v|_{\Gamma} \in H^{\frac{1}{2}}(\Gamma)$, $\frac{\partial v}{\partial \hat{n}} \in H^{-\frac{1}{2}}(\Gamma)$ and

$$0 = \left\langle \frac{\partial v}{\partial \hat{n}}, v \right\rangle$$
$$= \int_{\Omega} |\nabla v|^2 - \int_{\Omega} \lambda^2 |v|^2$$

and

$$0 = \left\langle \frac{\partial v}{\partial \hat{n}}, v \right\rangle$$

= $-\int_{\mathbb{R}^N \setminus \bar{\Omega}} |\nabla v|^2 + \int_{\mathbb{R}^N \setminus \bar{\Omega}} \lambda^2 |v|^2.$

Thus, v = 0 a.e. in \mathbb{R}^N , since λ^2 is not an eigenvalue of $(-\Delta)$ and v satisfies the radiation conditions at infinity. On Γ_1 , v = 0 and hence $v^i - v^e = \pm \tilde{g}_1 = 0$ and $\frac{\partial v^i}{\partial \hat{n}} - \frac{\partial v^e}{\partial \hat{n}} = \pm \tilde{g}_2 = 0$. Thus, $g_1 = 0$ and $g_2 = 0$.

Theorem 3.4 *The mixed boundary value problem* (1.2) *with* $Im(\lambda) > 0$ *possesses a unique solution u which is represented as*

$$u(x) = \mathbb{N}_{\lambda}h(x) - \mathbb{K}_{\lambda}(\mathring{f}_1 + \widetilde{g}_1)(x) + \mathbb{S}_{\lambda}(\mathring{f}_2 + \widetilde{g}_2)(x)$$

for unique $g_1 \in \widetilde{H}^{\frac{1}{2}}(\Gamma_2)$ and $g_2 \in \widetilde{H}^{-\frac{1}{2}}(\Gamma_1)$. This solution u also satisfies (1.4)–(1.5) and

$$\|u\|_{H^{\frac{1}{2}}(\Gamma)} + \left\|\frac{\partial u}{\partial \hat{n}}\right\|_{H^{-\frac{1}{2}}(\Gamma)} \le C\{\|f_1\|_{H^{\frac{1}{2}}(\Gamma_1)} + \|f_2\|_{H^{-\frac{1}{2}}(\Gamma_2)} + \|h\|_{\widetilde{H}^{-1}(\Omega)}\}. (3.7)$$

Proof The solvability and uniqueness of problem (1.2) depend on the invertibility of the operator A. Due to Theorem 3.3 we know that A is invertible. Hence, there exists a unique pair $(g_1, g_2) \in \tilde{H}^{\frac{1}{2}}(\Gamma_2) \times \tilde{H}^{-\frac{1}{2}}(\Gamma_1)$ such that

$$A(g_1, g_2) = (F^*, G^*),$$

where, $F^* \in H^{\frac{1}{2}}(\Gamma_1)$ and $G^* \in H^{-\frac{1}{2}}(\Gamma_2)$ as defined in Eqs. (2.20) and (2.21). Then we can represent

$$u(x) = \mathbb{N}_{\lambda}h(x) - \mathbb{K}_{\lambda}\phi(x) + \mathbb{S}_{\lambda}\psi(x),$$

where, $\phi = \mathring{f}_1 + \widetilde{g}_1$ and $\psi = \mathring{f}_2 + \widetilde{g}_2$ are the Cauchy data for the problem (1.2). Since

$$\begin{pmatrix} K_{21} & -S_{11} \\ D_{22} & -K_{12}^* \end{pmatrix} \begin{pmatrix} g_1 \\ g_2 \end{pmatrix} = \begin{pmatrix} F^* \\ G^* \end{pmatrix},$$
(3.8)

so we can write

$$g_1 = D_{22}^{-1}(K_{12}^*g_2 + G^*).$$
(3.9)

Substituting the value of g_1 in the Eq. (3.8) we get

$$(S_{11} - K_{21}D_{22}^{-1}K_{12}^*)g_2 = K_{21}D_{22}^{-1}G^* - F^*.$$
(3.10)

We will now show that the operator $H := S_{11} - K_{21}D_{22}^{-1}K_{12}^* : \widetilde{H}^{-\frac{1}{2}}(\Gamma_1) \to H^{\frac{1}{2}}(\Gamma_1)$ is invertible. We can then represent g_1, g_2 as follows.

$$g_1 = D_{22}^{-1}G^* + D_{22}^{-1}K_{12}^*H^{-1}(K_{21}D_{22}^{-1}G^* - F^*)$$

and

$$g_2 = H^{-1}(K_{21}D_{22}^{-1}G^* - F^*).$$

Claim: $H = S_{11} - K_{21}D_{22}^{-1}K_{12}^*$ is invertible.

We know K_{21} and K_{12}^* are compact operators. So $K_{21}D_{22}^{-1}K_{12}^*$ is also compact. Using Theorem 2.6 we get ind(H) = 0, since S_{11} is bijective. Thus we only need to show that H is injective.

Suppose $H(g_2) = 0$, for some $g_2 \in \widetilde{H}^{-\frac{1}{2}}(\Gamma_1)$. We express $w(x) = \mathbb{S}_{\lambda} \widetilde{g_2}(x) - \mathbb{K}_{\lambda} D_{22}^{-1} K_{12}^* g_2(x)$. We observe that on Γ_1 ,

$$w = (S_{11} - K_{21}D_{22}^{-1}K_{12}^*)g_2$$

= H(g_2)
= 0

and on Γ_2 ,

$$\frac{\partial w}{\partial \hat{n}} = K_{12}^* g_2 - D_{22} \{ D_{22}^{-1} K_{12}^* \} g_2$$
$$= 0.$$

Therefore, w satisfies (3.5) and (3.6). Following the proof of Theorem 3.3 we conclude that $g_2 = 0$. So H is injective hence invertible. Furthermore,

$$\begin{split} \|u\|_{H^{\frac{1}{2}}(\Gamma)} &= \left\| \mathring{f}_{1}^{1} + \widetilde{g}_{1} \right\|_{H^{\frac{1}{2}}(\Gamma)} \\ &\leq \left\| \mathring{f}_{1} \right\|_{H^{\frac{1}{2}}(\Gamma)} + \left\| \widetilde{g}_{1} \right\|_{H^{\frac{1}{2}}(\Gamma)} \\ &\leq C \left\{ \|f_{1}\|_{H^{\frac{1}{2}}(\Gamma_{1})} + \left\| K_{12}^{*}g_{2} + G^{*} \right\|_{H^{-\frac{1}{2}}(\Gamma_{2})} \right\} \text{ (by (2.17) and (3.9))} \\ &\leq C \left\{ \|f_{1}\|_{H^{\frac{1}{2}}(\Gamma_{1})} + \left\| G^{*} \right\|_{H^{-\frac{1}{2}}(\Gamma_{2})} + \left\| g_{2} \right\|_{H^{-\frac{1}{2}}(\Gamma_{1})} \right\} \\ &\leq C \left\{ \|f_{1}\|_{H^{\frac{1}{2}}(\Gamma_{1})} + \left\| G^{*} \right\|_{H^{-\frac{1}{2}}(\Gamma_{2})} + \left\| K_{21}D_{22}^{-1}G^{*} - F^{*} \right\|_{H^{\frac{1}{2}}(\Gamma_{1})} \right\} \text{ (by (3.10))} \\ &\leq C \left\{ \|f_{1}\|_{H^{\frac{1}{2}}(\Gamma_{1})} + \|f_{2}\|_{H^{-\frac{1}{2}}(\Gamma_{2})} + \|h\|_{\widetilde{H}^{-1}(\Omega)} 1 \right\} \text{ (by (2.20) and (2.21))} \end{aligned}$$

$$(3.11)$$

Similarly

$$\left\|\frac{\partial u}{\partial \hat{n}}\right\|_{H^{-\frac{1}{2}}(\Gamma)} \le C\left\{\|f_1\|_{H^{\frac{1}{2}}(\Gamma_1)} + \|f_2\|_{H^{-\frac{1}{2}}(\Gamma_2)} + \|h\|_{\widetilde{H}^{-1}(\Omega)}\right\}.$$
 (3.12)

On combining inequalities (3.11) and (3.12) we get

$$\|u\|_{H^{\frac{1}{2}}(\Gamma)} + \left\|\frac{\partial u}{\partial \hat{n}}\right\|_{H^{-\frac{1}{2}}(\Gamma)} \le C\{\|f_1\|_{H^{\frac{1}{2}}(\Gamma_1)} + \|f_2\|_{H^{-\frac{1}{2}}(\Gamma_2)} + \|h\|_{\widetilde{H}^{-1}(\Omega)}\}.$$

Proceeding similarly for the exterior problem (1.3) of (P1), the following Theorem can be established.

Theorem 3.5 For $Im(\lambda) > 0$, problem (P1) with given boundary data $f_1 \in H^{\frac{1}{2}}(\Gamma_1)$ and $f_2 \in H^{-\frac{1}{2}}(\Gamma_2)$ has a unique solution u which is represented as

$$u(x) = \begin{cases} \mathbb{N}_{\lambda}h(x) - \mathbb{K}_{\lambda}(\mathring{f}_{1} + \widetilde{g}_{1})(x) + \mathbb{S}_{\lambda}(\mathring{f}_{2} + \widetilde{g}_{2})(x), & \text{if } x \in \Omega\\ \mathbb{K}_{\lambda}(\mathring{f}_{1} + \widetilde{g}_{1})(x) - \mathbb{S}_{\lambda}(\mathring{f}_{2} + \widetilde{g}_{2})(x), & \text{if } x \in \mathbb{R}^{N} \setminus \bar{\Omega} \end{cases}$$

for unique $g_1 \in \widetilde{H}^{\frac{1}{2}}(\Gamma_2)$ and $g_2 \in \widetilde{H}^{-\frac{1}{2}}(\Gamma_1)$. The solution u belongs to $H^1(\Omega)$ for (IP1) and belongs to $H^1_{loc}(\mathbb{R}^N \setminus \overline{\Omega})$ for (EP1) satisfying the conditions (1.4)–(1.5) at infinity. Furthermore, u satisfies (3.7).

4 Existence Results of (P2)

Problem (P2) is a mixed boundary value problem of Poisson equation where $\mu \in \mathcal{M}(\overline{\Omega})$ is a bounded Radon measure supported on Ω and the boundary data are $f_1 \in H^{\frac{1}{2}}(\Gamma_1)$ and $f_2 \in H^{-\frac{1}{2}}(\Gamma_2)$.

Definition 4.1 We say a function $u \in W^{1,1}(\Omega)$ is a weak solution to the problem (1.6) if

$$\int_{\Omega} \nabla u \cdot \nabla \varphi = \int_{\Omega} \varphi d\mu + \int_{\Gamma_2} f_2 \varphi, \ \forall \varphi \in X$$

where, $X = \{\varphi \in C^1(\overline{\Omega}) : \varphi|_{\Gamma_1} = 0\}$ is the test function space. Similarly a function $u \in W_{loc}^{1,1}(\mathbb{R}^N \setminus \overline{\Omega})$ is said to be a weak solution of (1.7) if

$$\int_{\mathbb{R}^N \setminus \bar{\Omega}} \nabla u \cdot \nabla \varphi = - \int_{\Gamma_2} f_2 \varphi, \; \forall \varphi \in Z$$

where, $Z = \{\zeta \in C_c^1(\mathbb{R}^N \setminus \Omega) : \zeta|_{\Gamma_1} = 0 \text{ and satisfies (1.8)}\}.$

We will now approximate $\mu \in \mathcal{M}(\overline{\Omega})$ by a smooth sequence $(\mu_n) \subset L^{\infty}(\Omega)$, in the weak* topology, i.e.

$$\int_{\Omega} g \, d\mu_n \to \int_{\Omega} g \, d\mu, \; \forall g \in C(\bar{\Omega}). \tag{4.1}$$

In order to show the existence of solutions to (P2), we consider the '*approximating*' problems to (1.6)–(1.7) which are as follows.

$$-\Delta u_n = \mu_n \text{ in } \Omega,$$

$$u_n = f_1 \text{ on } \Gamma_1,$$

$$\frac{\partial u_n}{\partial \hat{n}} = f_2 \text{ on } \Gamma_2$$
(4.2)

and

$$-\Delta u_n = 0 \text{ in } \mathbb{R}^N \setminus \overline{\Omega},$$

$$u_n = f_1 \text{ on } \Gamma_1,$$

$$\frac{\partial u_n}{\partial \hat{n}} = f_2 \text{ on } \Gamma_2.$$
(4.3)

These '*approximating*' problems are special cases of (P1) with $\lambda = 0$. The weak formulation to (4.2) is

$$\int_{\Omega} \nabla u_n \cdot \nabla \varphi = \int_{\Omega} \varphi \mu_n + \int_{\Gamma_2} f_2 \varphi, \ \forall \varphi \in X.$$
(4.4)

Theorem 4.2 The problems (4.2) and (4.3) admit a unique solution u_n which is represented as

$$u_{n}(x) = \begin{cases} \mathbb{N}_{0}\mu_{n}(x) - \mathbb{K}_{0}(\mathring{f}_{1} + \widetilde{g}_{1_{n}})(x) + \mathbb{S}_{0}(\mathring{f}_{2} + \widetilde{g}_{2_{n}})(x), & \text{if } x \in \Omega\\ \mathbb{K}_{0}(\mathring{f}_{1} + \widetilde{g}_{1_{n}})(x) - \mathbb{S}_{0}(\mathring{f}_{2} + \widetilde{g}_{2_{n}})(x), & \text{if } x \in \mathbb{R}^{N} \setminus \bar{\Omega} \end{cases}$$
(4.5)

for a unique pair $(g_{1_n}, g_{2_n}) \in \widetilde{H}^{\frac{1}{2}}(\Gamma_2) \times \widetilde{H}^{-\frac{1}{2}}(\Gamma_1)$. The solution u_n belongs to $H^1(\Omega)$ for the problem (4.2) and belongs to $H^1_{loc}(\mathbb{R}^N \setminus \overline{\Omega})$ for the problem (4.3) satisfying the radiation condition (1.8) at infinity.

Proof The invertibility of the operators S_{11} and D_{22} follows from Theorem 3.1 and Theorem 3.2 of Chang and Choe (2008), respectively. Further, following the proof of Theorem 3.3, Theorem 3.4 one can see that the matrix operator A is invertible and the problems (4.2)–(4.3) have a unique solution denoted as u_n . The solution u_n can be represented as in (4.5) by Theorem 3.5 and satisfies the condition (1.8) at infinity. \Box

Now to show that problem in (1.6), involving measure, possesses a solution u we need to pass the limit $n \to \infty$ in the weak formulation (4.4).

Lemma 4.3 Let us suppose that u_n is a solution of problem (4.2) with $f_1 \in H^{\frac{1}{2}}(\Gamma_1)$ and $f_2 \in H^{-\frac{1}{2}}(\Gamma_2)$. Then the sequence (u_n) is bounded in $W^{1,q}(\Omega)$ for all $q < \frac{N}{N-1}$.

Proof From the continuous embedding (2.1) we have

$$L^{\frac{N}{N-1}}(\Omega) \hookrightarrow M^{\frac{N}{N-1}}(\Omega) \hookrightarrow L^{\frac{N}{N-1}-\epsilon}(\Omega).$$
 (4.6)

If we can show that (u_n) is bounded in $M^{\frac{N}{N-1}}(\Omega)$, then this will also imply (u_n) to be bounded in $L^{\frac{N}{N-1}-\epsilon}(\Omega)$. More precisely, we can say (u_n) to be bounded in $L^q(\Omega)$, for every $q < \frac{N}{N-1}$.

Claim. The sequences (u_n) and (∇u_n) are bounded in $M^{\frac{N}{N-1}}(\Omega)$.

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Let us fix a constant a > 0. The truncation function of u_n is defined as

$$T_a(u_n) = \max\left\{-a, \min\{a, u_n\}\right\}.$$

From Theorem 4.2, we have $\frac{\partial u_n}{\partial \hat{n}}\Big|_{\Gamma} = \mathring{f}_2 + \widetilde{g}_{2_n}$, for a unique $g_{2_n} \in \widetilde{H}^{-\frac{1}{2}}(\Gamma_1)$. For a change, let us test (4.2) with functions from $H^1(\Omega)$. Thus for $v \in H^1(\Omega)$ we have

$$\int_{\Omega} \nabla u_n \cdot \nabla \upsilon = \int_{\Omega} \mu_n \upsilon + \int_{\Gamma} \left(\mathring{f}_2 + \widetilde{g}_{2_n} \right) \upsilon.$$

Note that the integrals are all finite. Thus in particular, choose $v = T_a(u_n)$, to get

$$\begin{split} \int_{\Omega} |\nabla T_a(u_n)|^2 &\leq \int_{\Omega} \nabla u_n \cdot \nabla (T_a(u_n)) \\ &= \int_{\Omega} T_a(u_n) \mu_n + \int_{\Gamma} (\mathring{f}_2 + \widetilde{g}_{2_n}) T_a(u_n) \\ &\leq a \int_{\Omega} \mu_n + a \int_{\Gamma} (\mathring{f}_2 + \widetilde{g}_{2_n}) \\ &\leq Ca. \end{split}$$
(4.7)

Consider

$$\{|\nabla u_n| \ge k\} = \{|\nabla u_n| \ge k, u_n < a\} \cup \{|\nabla u_n| \ge k, u_n \ge a\}$$
$$\subset \{|\nabla u_n| \ge k, u_n < a\} \cup \{u_n \ge a\}$$
$$= B_1 \cup B_2 \subset \Omega$$

where $B_1 = \{|\nabla u_n| \ge k, u_n < a\}$ and $B_2 = \{u_n \ge a\}$. Hence, due to the subadditivity property of the Lebesgue measure '*m*' we have

$$m\left(\{|\nabla u_n| \ge k\}\right) \le m(B_1) + m(B_2). \tag{4.8}$$

Using the Sobolev inequality, we have

$$\left(\int_{\Omega} |T_a(u_n)|^{2^*}\right)^{\frac{2^*}{2^*}} \le \frac{1}{\lambda_1} \int_{\Omega} |\nabla T_a(u_n)|^2 \le Ca$$
(4.9)

where λ_1 is the first eigenvalue of $(-\Delta)$ with Dirichlet boundary condition. Now we restrict the above inequality (4.9) on B_2 to get

$$a^{2}m (\{u_{n} \ge a\})^{\frac{2}{2^{*}}} \le Ca, \text{ (Since } T_{a}(u_{n}) = a \text{ in } B_{2}).$$

Thus,

$$m\left(\{u_n \ge a\}\right) \le \frac{C}{a^{\frac{N}{N-2}}}, \ \forall a \ge 1.$$

Hence, (u_n) is bounded in $M^{\frac{N}{N-2}}(\Omega)$ and also bounded in $M^{\frac{N}{N-1}}(\Omega)$. Similarly on restricting (4.7) on B_1 , we have

$$m\left(\{|\nabla u_n| \ge k, u_n < a\}\right) \le \frac{1}{k^2} \int_{\Omega} |\nabla T_a(u_n)|^2$$
$$\le \frac{Ca}{k^2}, \ \forall a > 1.$$

Now the inequality (4.8) becomes

$$m\left(\{|\nabla u_n| \ge k\}\right) \le m(\{B_1\}) + m(B_2) \\ \le \frac{Ca}{k^2} + \frac{C}{a^{\frac{N}{N-2}}}, \ \forall a > 1.$$

On choosing $a = k^{\frac{N-2}{N-1}}$ we get

$$m\left(\{|\nabla u_n| \ge k\}\right) \le \frac{C}{k^{\frac{N}{N-1}}}, \ \forall k \ge 1.$$

So (∇u_n) is bounded in $M^{\frac{N}{N-1}}(\Omega)$. Hence, we conclude that (u_n) is bounded in $W^{1,q}(\Omega)$ for every $q < \frac{N}{N-1}$.

Theorem 4.4 There exists a weak solution u of (1.6) in $W^{1,q}(\Omega)$, for all $q < \frac{N}{N-1}$.

Proof According to Lemma 4.3, (u_n) is bounded in $W^{1,q}(\Omega)$ which is a reflexive space. This implies that there exists a function $u \in W^{1,q}(\Omega)$ such that u_n converges weakly to u, i.e. $u_n \stackrel{w}{\rightharpoonup} u$ in $W^{1,q}(\Omega)$, for all $q < \frac{N}{N-1}$. Thus for $\varphi \in X$,

$$\lim_{n \to +\infty} \int_{\Omega} \nabla u_n \cdot \nabla \varphi = \int_{\Omega} \nabla u \cdot \nabla \varphi.$$

The sequence (μ_n) converges to μ in the weak* topology in the sense given in (4.1). On passing the limit $n \to \infty$ in the weak formulation (4.4) involving μ_n we obtain

$$\int_{\Omega} \nabla u \cdot \nabla \varphi = \int_{\Omega} \varphi d\mu + \int_{\Gamma_2} f_2 \varphi, \ \forall \varphi \in X.$$

Hence, a weak solution of (1.6) in $W^{1,q}(\Omega)$ for every $q < \frac{N}{N-1}$ is guaranteed. \Box

Now with the consideration of Theorems 4.2 and 4.4 we state our main result which is as follows.

Theorem 4.5 There exists a weak solution u of (P2) with $\mu \in \mathcal{M}(\bar{\Omega})$, with support is in Ω , as a nonhomogeneous term, $f_1 \in H^{\frac{1}{2}}(\Gamma_1)$ and $f_2 \in H^{-\frac{1}{2}}(\Gamma_2)$. The solution ubelongs to $W^{1,q}(\Omega)$, for every $q < \frac{N}{N-1}$, for (IP2) and belongs to $H^1_{loc}(\mathbb{R}^N \setminus \bar{\Omega})$ for (EP2) satisfying Eq. (1.8) at infinity.

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