# Spectral properties of the Neumann-Poincaré operator and uniformity of estimates for the conductivity equation with complex coefficients* 

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#### Abstract

We consider well-posedness of the boundary value problem in presence of an inclusion with complex conductivity $k$. We first consider the transmission problem in $\mathbb{R}^{d}$ and characterize solvability of the problem in terms of the spectrum of the Neumann-Poincare operator. We then deal with the boundary value problem and show that the solution is bounded in its $H^{1}$-norm uniformly in $k$ as long as $k$ is at some distance from a closed interval in the negative real axis. We then show with an estimate that the solution depends on $k$ in its $H^{1}$-norm Lipschitz continuously. We finally show that the boundary perturbation formula in presence of a diametrically small inclusion is valid uniformly in $k$ away from the closed interval mentioned before. The results for the single inclusion case are extended to the case when there are multiple inclusions with different complex conductivities: We first obtain a complete characterization of solvability when inclusions consist of two disjoint disks and then prove solvability and uniform estimates when imaginary parts of conductivities have the same signs. The results are obtained using the spectral property of the associated NeumannPoincaré operator and the spectral resolution.


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Key words. Neumann-Poincaré operator, Lipschitz domain, spectrum, spectral resolution, asymptotic formula, uniformity, complex conductivity, high contrast

## 1 Introduction

Let $\Omega$ be a bounded simply connected domain in $\mathbb{R}^{d}(d \geq 2)$ and let $D$ be a simply connected domain compactly contained in $\Omega$. We assume that boundaries of $D$ and $\Omega$ are Lipschitz continuous and the conductivity of $\Omega \backslash D$ is 1 while that of $D$ is $k$ so that the conductivity distribution is given by

$$
\begin{equation*}
\gamma_{k}=\chi(\Omega \backslash D)+k \chi(D) \tag{1.1}
\end{equation*}
$$

[^0]where $\chi(D)$ is the characteristic function of $D$. For a given Neumann data $g$ consider the following elliptic problem:
\[

\left\{$$
\begin{array}{l}
\nabla \cdot\left(\gamma_{k} \nabla u\right)=0 \quad \text { in } \Omega,  \tag{1.2}\\
\left.\partial_{\nu} u\right|_{\partial \Omega}=g, \\
\int_{\partial \Omega} u d \sigma=0
\end{array}
$$\right.
\]

Here and throughout this paper $\partial_{\nu} u$ denotes the outward normal derivative of $u$ and the conductivity $k=k^{\prime}+i k^{\prime \prime}$ is a complex number ( $k^{\prime}$ and $k^{\prime \prime}$ denote the real and imaginary parts of $k$, respectively). We emphasize that the problem (1.2) admits a unique solution if $k$ is not on the negative real axis.

The solution $u_{k}$ to (1.2) varies depending on the conductivity $k$ regarded as a parameter. For instance, we have a standard regularity estimate

$$
\begin{equation*}
\left\|u_{k}\right\|_{H^{1}(\Omega)} \leq C\|g\|_{H^{-1 / 2}(\partial \Omega)} \tag{1.3}
\end{equation*}
$$

for some constant $C$ which may depend on $k$. Here and throughout this paper $H^{s}(\Omega)$ (and $\left.H^{s}(\partial \Omega)\right)$ denotes the standard $L^{2}$-Sobolev space and $H^{-1 / 2}(\partial \Omega)$ is the dual space of $H^{1 / 2}(\partial \Omega)$. However, if $k$ is real $(0 \leq k \leq \infty)$, it is proved in [34] that the constant $C$ can be chosen independently of $k$. Moreover, if the inclusion $D$ is diametrically small, it is proved in the same paper that the asymptotic boundary perturbation formula (see section 5) is valid uniformly in $k$. We emphasize that these results were obtained using variational methods.

In this paper we develop a new method to investigate solvability of (1.2) and dependency of solution $u_{k}$ on the conductivity $k$ of the inclusion, and to extend above mentioned results to the case when $k$ is a complex number. We first prove existence and uniqueness of the solution to (1.2) when $k^{\prime \prime} \neq 0$. We then show that if $k=k^{\prime}+i k^{\prime \prime}$ satisfies

$$
\begin{equation*}
\left|k^{\prime \prime}\right| \geq-L k^{\prime} \tag{1.4}
\end{equation*}
$$

for any given constant $L>0$ (see the left figure in Figure 1.1), we prove that (1.3) holds for some $C$ independent of $k$. We then show (1.3) holds near $k=0$ and $k=\infty$ by a perturbation argument. The uniform estimates near 0 and $\infty$, and that in the sector imply that there is a closed interval such that uniform estimate holds away from the interval. We also show with a precise quantitative estimate that $u_{k}$ depends on $k$ in $H^{1}(\Omega)$-norm Lipschitz continuously. We finally prove uniform validity of the asymptotic boundary perturbation formula when $D$ is diametrically small, regardless of $k$ as long as $k$ satisfies (1.4).

Solvability and uniformity results for the single inclusion case are extended to the case when there are multiple inclusions. Namely, there are multiple inclusions $D_{j}, j=1, \ldots, M$, with complex conductivities $k_{j}$. If all the conductivities are the same, namely, $k_{1}=\ldots=k_{M}$, then results for the single inclusion case are valid for the multiple inclusion case without change. However, if we allow them to be different (and complex), finding meaningful conditions which guarantee solvability of the problem seems quite difficult (see subsection 6.2). We attempt to present one sufficient condition in this paper. The condition is basically that $k_{j}^{\prime \prime}$ have the same sign for all $j$.

There is growing interest in the complex conductivity, especially in relation to the electrical impedance tomography (EIT). For instance, the imaginary part of the complex conductivity changes depending on the frequency of the prescribed current, and by exploiting this property inclusion can be reconstructed with a high resolution (see, for example, [30]). Another example is the size estimation problem in the EIT, which is to derive bounds on the volume fraction of the inclusion via boundary measurements. This problem is recently considered when the conductivity of the inclusion is complex [13] [24] 37]. The bounds for the complex conductivity
case turned out to be quite tight as the numerical examples presented in the last two papers show.

The method of this paper is based on the spectral property of the Neumann-Poincaré (NP) operator related to the problem (1.2). The NP operator is a boundary integral operator which appears naturally when solving the Neumann (and Dirichlet) boundary value problems using single or double layer potentials. It is not self-adjoint with respect to the usual $L^{2}$-inner product. However, it can be symmetrized using a certain twisted inner product and Plemelj's symmetrization principle [29]. We show that if $\partial D$ is Lipschitz, then the spectrum of the NP operator defined on $\partial D$ on a Sobolev space lies in $\left[-b_{\partial D}, b_{\partial D}\right]$ for some $b_{\partial D}<1 / 2$. Here, the spectrum or resolvent $\lambda$ of the NP operator is related to the conductivity $k$ by the bilinear transformation

$$
\begin{equation*}
\lambda=\lambda(k)=\frac{k+1}{2(k-1)} . \tag{1.5}
\end{equation*}
$$

Since bilinear transformations map circles on the Riemann sphere onto circles, this transformation maps the region defined by (1.4) onto the region outside the oval shaped curve (see Figure 1.1). We emphasize that the boundary of the region in the $\lambda$-space intersects with the real axis at $1 / 2$ and $-1 / 2$. Since $b_{\partial D}<1 / 2$, there is some distance between the transformed region and the spectrum of the NP operator. Using this property and the spectral resolution of a self-adjoint operator we are able to obtain the results described above.

We also consider the following transmission problem in $\mathbb{R}^{d}$ :

$$
\left\{\begin{array}{l}
\nabla \cdot\left(\chi\left(\mathbb{R}^{d} \backslash D\right)+k \chi(D)\right) \nabla u=0 \quad \text { in } \mathbb{R}^{d},  \tag{1.6}\\
u(x)-h(x)=O\left(|x|^{1-d}\right) \quad \text { as }|x| \rightarrow \infty
\end{array}\right.
$$

where $h$ is a given harmonic function in $\mathbb{R}^{d}$. This problem is simpler than (1.2) because the function $h$ does not depend on $k$ (see the next section), and solvability of this problem is completely characterized by the spectrum of the NP operator. In fact, we prove that the problem is well-posed if $\lambda(k)$ does not belong to the spectrum of the NP operator (see section (3). We emphasize that the condition (1.4) is required for the boundary value problem to show that the Neumann-to-Dirichlet map for the problem (1.2) is bounded uniformly in $k$ (see Lemma 4.4).


Figure 1.1: Left: the region $\left|k^{\prime \prime}\right| \geq-L k^{\prime}$ (grey). Right: the region after the transformation $\lambda=\frac{k+1}{2(k-1)} \cdot b_{\partial D}$ is the spectral bound of the NP operator.

The program of estimating the spectrum of the NP operator goes back to Poincaré [36] as the name suggests, and was recently revisited in [29] with a new perspective. The spectrum of the

NP operator also has been studied using complex function theory such as the Beurling-Ahlfors transform and quasi-conformal mapping, for which we refer to [35] (and references therein) where a bound on the essential spectrum of the NP operator for two dimensional domains with corners is obtained. It is proved in [38] that the set $\mathbb{R} \backslash(-1 / 2,1 / 2)$ is contained in the resolvent set of the NP operator on $L_{0}^{2}(\partial \Omega)$ for a domain $\Omega$ with the Lipschtz boundary. Its spectrum on $H^{-1 / 2}$ lies in $(-1 / 2,1 / 2]$ (see [18]). Recently the spectral theory of the NP operator has been applied to analysis of cloaking by anomalous localized resonance on the plasmonic structure [2], analysis of high concentration of the gradient [3] [14] [15] [32], and a shape optimization problem [1].

This paper is organized as follows. In section 2 we review symmetrization and spectral properties of the Neumann-Poincaré operator and prove equivalence of norms. In section 3 we deal with the transmission problem in the whole space $\mathbb{R}^{d}$. In section 4 we deal with the boundary value problem and show that the solution exists and is unique, it is bounded in its $H^{1}$-norm uniformly in $k$ away from a closed interval in the negative real axis, and a quantitative estimate is derived which shows that the solution depends Lipschitz continuously on $k$. We then show in section 5 that the boundary perturbation formula when the inclusion is diametrically small is valid uniformly in $k$. The last section is to deal with the case when there are multiple inclusions.

## 2 Spectral properties of the Neumann-Poincaré operator

Let $\Gamma(x)$ be the fundamental solution to the Laplacian, i.e.,

$$
\Gamma(x)= \begin{cases}\frac{1}{2 \pi} \ln |x|, & d=2  \tag{2.1}\\ \frac{1}{(2-d) \omega_{d}}|x|^{2-d}, & d \geq 3\end{cases}
$$

where $\omega_{d}$ denotes the area of the unit sphere in $\mathbb{R}^{d}$. Suppose that the inclusion $D$ has a single simply connected component. The single layer potential $\mathcal{S}_{\partial D}[\varphi]$ of a density function $\varphi \in H^{-1 / 2}(\partial D)$ is defined by

$$
\begin{equation*}
\mathcal{S}_{\partial D}[\varphi](x):=\int_{\partial D} \Gamma(x-y) \varphi(y) d \sigma(y), \quad x \in \mathbb{R}^{d} . \tag{2.2}
\end{equation*}
$$

It is well known (see for example [7] [21]) that $\mathcal{S}_{\partial D}[\varphi]$ satisfies the jump relation

$$
\begin{equation*}
\left.\frac{\partial}{\partial \nu} \mathcal{S}_{\partial D}[\varphi]\right|_{ \pm}(x)=\left( \pm \frac{1}{2} I+\mathcal{K}_{\partial D}^{*}\right)[\varphi](x), \quad x \in \partial D \tag{2.3}
\end{equation*}
$$

where $\frac{\partial}{\partial \nu}$ denotes the outward normal derivative, the subscripts $\pm$ indicate the limit from outside and inside $D$, respectively, and the operator $\mathcal{K}_{\partial D}^{*}$ is defined by

$$
\begin{equation*}
\mathcal{K}_{\partial D}^{*}[\varphi](x)=\frac{1}{\omega_{d}} \int_{\partial D} \frac{(x-y) \cdot \nu_{x}}{|x-y|^{d}} \varphi(y) d \sigma(y), \quad x \in \partial D . \tag{2.4}
\end{equation*}
$$

Here $\nu_{x}$ denotes the outward unit normal vector to $\partial D$ at $x$.
The operator $\mathcal{K}_{\partial D}^{*}$ is called the Neumann-Poincaré (NP) operator associated with the domain $D$. If $\partial D$ is Lipschitz continuous, then $\mathcal{K}_{\partial D}^{*}$ is a singular integral operator and known to be bounded on $L^{2}(\partial D)$ (and on $H^{-1 / 2}(\partial D)$ ) [20]. Let $H_{0}^{-1 / 2}(\partial D)$ be the collection of $\varphi \in$ $H^{-1 / 2}(\partial D)$ with the mean zero, i.e.,

$$
\begin{equation*}
\langle\varphi, 1\rangle=0 \tag{2.5}
\end{equation*}
$$

where $\langle$,$\rangle is the H^{-1 / 2}-H^{1 / 2}$ pairing. For $\varphi, \psi \in H_{0}^{-1 / 2}(\partial D)$, define

$$
\begin{equation*}
\langle\varphi, \psi\rangle_{\mathcal{H}}:=-\left\langle\varphi, \mathcal{S}_{\partial D}[\psi]\right\rangle . \tag{2.6}
\end{equation*}
$$

Since $\mathcal{S}_{\partial D}$ maps $H^{-1 / 2}(\partial D)$ into $H^{1 / 2}(\partial D)$, the right hand side of (2.6) is well-defined. It is known (see, for example, [2] [15] [23] [29]) that $\langle,\rangle_{\mathcal{H}}$ is an inner product on $H_{0}^{-1 / 2}(\partial D)$ and $\mathcal{K}_{\partial D}^{*}$ is self-adjoint with respect to this inner product, which is a consequence of Plemelj's symmetrization principle (also known as Calderón's identity)

$$
\begin{equation*}
\mathcal{S}_{\partial D} \mathcal{K}_{\partial D}^{*}=\mathcal{K}_{\partial D} \mathcal{S}_{\partial D} . \tag{2.7}
\end{equation*}
$$

Here $\mathcal{K}_{\partial D}$ is the adjoint of $\mathcal{K}_{\partial D}^{*}$ with respect to the usual $L^{2}$-inner product. It is worth mentioning that $\mathcal{K}_{\partial D}^{*}$ is not self-adjoint in the usual inner product unless the domain $D$ is a disk or a ball [31]. Let $\mathcal{H}=\mathcal{H}(\partial D)$ be the space $H_{0}^{-1 / 2}(\partial D)$ equipped with the inner product $\langle,\rangle_{\mathcal{H}}$. We denote the norm associated with $\langle,\rangle_{\mathcal{H}}$ by $\|\cdot\|_{\mathcal{H}}$.

Let us write $\|\varphi\|_{H^{s}(\partial D)}$ as $\|\varphi\|_{s}$ from now on for ease of notation.
Lemma 2.1 Let $D$ be a bounded Lipschitz domain in $\mathbb{R}^{d}(d \geq 2)$. The operator $A: H^{-1 / 2}(\partial D) \times$ $\mathbb{C} \rightarrow H^{1 / 2}(\partial D) \times \mathbb{C}$ by

$$
\begin{equation*}
A(\psi, a):=\left(\mathcal{S}_{\partial D}[\psi]+a,\langle\psi, 1\rangle\right) \tag{2.8}
\end{equation*}
$$

is invertible.
Proof. In two dimensions, it is shown in [7, Theorem 2.26] that $A: L^{2}(\partial D) \times \mathbb{C} \rightarrow H^{1}(\partial D) \times \mathbb{C}$ is invertible. If $(\psi, a) \in L^{2}(\partial D) \times \mathbb{C}$ and $(\varphi, b) \in H^{-1}(\partial D) \times \mathbb{C}$, then

$$
\begin{aligned}
\langle(\varphi, b), A(\psi, a)\rangle & =\left\langle\varphi, \mathcal{S}_{\partial D}[\psi]+a\right\rangle+b \int_{\partial D} \bar{\psi} d \sigma \\
& =\left\langle\mathcal{S}_{\partial D}[\varphi], \psi\right\rangle+\langle\varphi, a\rangle+b \int_{\partial D} \bar{\psi} d \sigma \\
& =\left\langle\mathcal{S}_{\partial D}[\varphi]+b, \psi\right\rangle+\bar{a}\langle\varphi, 1\rangle=\langle A(\varphi, b),(\psi, a)\rangle .
\end{aligned}
$$

So, by duality, $A: H^{-1}(\partial D) \times \mathbb{C} \rightarrow L^{2}(\partial D) \times \mathbb{C}$ is invertible. By interpolation we infer that $A: H^{-1 / 2}(\partial D) \times \mathbb{C} \rightarrow H^{1 / 2}(\partial D) \times \mathbb{C}$ in invertible.

In three or higher dimensions, it is proved in [38] that $\mathcal{S}_{\partial D}: L^{2}(\partial D) \rightarrow H^{1}(\partial D)$ is invertible. So, $A$ is a Fredholm operator of index zero. Thus one can show by the exactly the same proof as in two dimensions that $A: L^{2}(\partial D) \times \mathbb{C} \rightarrow H^{1}(\partial D) \times \mathbb{C}$ is invertible in three dimensions. So we obtain the desired result by the same argument (duality and interpolation).

Lemma 2.1 says that for any $f \in H^{1 / 2}(\partial D)$ and $b \in \mathbb{C}$ there is a unique pair $(\psi, a) \in$ $H^{-1 / 2}(\partial D) \times \mathbb{C}$ such that $\langle\psi, 1\rangle=b$,

$$
\begin{equation*}
\mathcal{S}_{\partial D}[\psi]+a=f, \tag{2.9}
\end{equation*}
$$

and

$$
\begin{equation*}
\|\psi\|_{-1 / 2}+|a| \leq C\left(\|f\|_{1 / 2}+|b|\right) \tag{2.10}
\end{equation*}
$$

for some constant $C$. In particular, if we take $b=0$, then $\psi \in H_{0}^{-1 / 2}(\partial D)$ and (2.10) becomes

$$
\begin{equation*}
\|\psi\|_{-1 / 2}+|a| \leq C\|f\|_{1 / 2} . \tag{2.11}
\end{equation*}
$$

So, we obtain the following lemma.

Lemma 2.2 For any $f \in H^{1 / 2}(\partial D)$ there is a unique $\psi \in H_{0}^{-1 / 2}(\partial D)$ such that

$$
\begin{equation*}
\langle\varphi, f\rangle=\langle\varphi, \psi\rangle_{\mathcal{H}} \tag{2.12}
\end{equation*}
$$

for all $\varphi \in H_{0}^{-1 / 2}(\partial D)$. Moreover, there is a constant $C$ independent of $f$ such that

$$
\begin{equation*}
\|\psi\|_{-1 / 2} \leq C\|f\|_{1 / 2} \tag{2.13}
\end{equation*}
$$

Proof. For $f \in H^{1 / 2}(\partial D)$ let $\psi \in H_{0}^{-1 / 2}(\partial D)$ and $a$ be such that (2.9) and (2.11) hold. Then we have

$$
\langle\varphi, f\rangle=\left\langle\varphi, \mathcal{S}_{\partial D}[\psi]+a\right\rangle=-\langle\varphi, \psi\rangle_{\mathcal{H}} .
$$

Replace $\psi$ by $-\psi$ to have (2.12). Uniqueness of $\psi$ is obvious and the proof is complete.
As a consequence we obtain the following theorem.
Theorem 2.3 Let $D$ be a bounded Lipschitz domain. There are constants $C_{1}$ and $C_{2}$ (which may depend on $D$ ) such that

$$
\begin{equation*}
C_{1}\|\varphi\|_{H^{-1 / 2}(\partial D)} \leq\|\varphi\|_{\mathcal{H}} \leq C_{2}\|\varphi\|_{H^{-1 / 2}(\partial D)} \tag{2.14}
\end{equation*}
$$

for all $\varphi \in H_{0}^{-1 / 2}(\partial D)$.
Proof. Since $\mathcal{S}_{\partial D}$ is bounded from $H^{-1 / 2}(\partial D)$ into $H^{1 / 2}(\partial D)$, we have for $\varphi \in \mathcal{H}$

$$
\left|\langle\varphi, \varphi\rangle_{\mathcal{H}}\right| \leq\left\|\mathcal{S}_{\partial D}[\varphi]\right\|_{1 / 2}\|\varphi\|_{-1 / 2} \leq C\|\varphi\|_{-1 / 2}^{2}
$$

So, we have

$$
\begin{equation*}
\|\varphi\|_{\mathcal{H}} \leq C\|\varphi\|_{-1 / 2} \tag{2.15}
\end{equation*}
$$

for some constant $C$.
To prove the opposite inequality, for $f \in H^{1 / 2}(\partial D)$ choose $\psi \in H_{0}^{-1 / 2}(\partial D)$ satisfying (2.12) and (2.13). Then we have, by the Cauchy-Schwarz inequality and (2.13),

$$
\begin{equation*}
|\langle\varphi, f\rangle|=\left|\langle\varphi, \psi\rangle_{\mathcal{H}}\right| \leq\|\varphi\|_{\mathcal{H}}\|\psi\|_{\mathcal{H}} \leq C\|\varphi\|_{\mathcal{H}}\|\psi\|_{-1 / 2} \leq C\|\varphi\|_{\mathcal{H}}\|f\|_{1 / 2} . \tag{2.16}
\end{equation*}
$$

Since the above inequality holds for all $f \in H^{1 / 2}(\partial D)$, we obtain

$$
\begin{equation*}
\|\varphi\|_{-1 / 2} \leq C\|\varphi\|_{\mathcal{H}} . \tag{2.17}
\end{equation*}
$$

This completes the proof.
Since $\mathcal{K}_{\partial D}^{*}$ is self-adjoint on $\mathcal{H}$, its spectrum $\sigma\left(\mathcal{K}_{\partial D}^{*}\right)$ is real, consists of point and continuous spectra, and is a closed set. Moreover, by the spectral resolution theorem (see [39]) there is a family of projection operators $\mathcal{E}(t)$ on $\mathcal{H}$ (called a resolution of identity) such that

$$
\begin{equation*}
\mathcal{K}_{\partial D}^{*}=\int_{t \in \sigma\left(\mathcal{K}_{\partial D}^{*}\right)} t d \mathcal{E}(t) \tag{2.18}
\end{equation*}
$$

Let $b_{\partial D}$ be the spectral bound of $\mathcal{K}_{\partial D}^{*}$, namely,

$$
\begin{equation*}
b_{\partial D}:=\sup \left\{|\lambda|: \lambda \in \sigma\left(\mathcal{K}_{\partial D}^{*}\right)\right\} . \tag{2.19}
\end{equation*}
$$

It is well-known that for any Lipschitz domain $b_{\partial D} \leq \frac{1}{2}$. In fact, it is proved in [28] that

$$
\begin{equation*}
b_{\partial D}=\frac{1}{2} \sup _{\varphi \in \mathcal{H}} \frac{\left.\left|\int_{\mathbb{R}^{d} \backslash D}\right| \nabla \mathcal{S}_{\partial D}[\varphi]\right|^{2} d x-\int_{D}\left|\nabla \mathcal{S}_{\partial D}[\varphi]\right|^{2} d x \mid}{\int_{\mathbb{R}^{d}}\left|\nabla \mathcal{S}_{\partial D}[\varphi]\right|^{2} d x} . \tag{2.20}
\end{equation*}
$$

We obtain the following lemma.

Theorem 2.4 Let $D$ be a bounded Lipschitz domain. Then,

$$
\begin{equation*}
b_{\partial D}<\frac{1}{2} \tag{2.21}
\end{equation*}
$$

Proof. It is proved in [38] that $\lambda I-\mathcal{K}_{\partial D}^{*}$ is invertible on $L^{2}(\partial D)$ if $\lambda \in \mathbb{R} \backslash(-1 / 2,1 / 2]$, and invertible on $L_{0}^{2}(\partial D)$ if $\lambda \in \mathbb{R} \backslash(-1 / 2,1 / 2)$. Using this result, it is proved in [18] that $\lambda I-\mathcal{K}_{\partial D}^{*}$ is invertible on $H^{-1 / 2}(\partial D)$ if $\lambda \notin(-1 / 2,1 / 2]$. But, the proof there shows that $\lambda I-\mathcal{K}_{\partial D}^{*}$ is invertible on $H_{0}^{-1 / 2}(\partial D)$ if $\lambda \notin(-1 / 2,1 / 2)$. We then infer from (2.14) that the spectrum $\sigma\left(\mathcal{K}_{\partial D}^{*}\right)$ of $\mathcal{K}_{\partial D}^{*}$ on $\mathcal{H}$ lies in $(-1 / 2,1 / 2)$. Since $\sigma\left(\mathcal{K}_{\partial D}^{*}\right)$ is a closed set, we have (2.21).

Before completing this section let us make a few remarks on the spectrum of the NP operator. If $\partial D$ is $\mathcal{C}^{1, \alpha}$ for some $\alpha>0$, then $\mathcal{K}_{\partial D}^{*}$ is compact and has only a point spectrum accumulating to 0 . The point spectrum of $\mathcal{K}_{\partial D}^{*}$ is completely known when $D$ is a disk, a ball, or an ellipse (see [23]). For example, if $D$ is an ellipse of the long axis $a$ and short axis $b$, then eigenvalues of $\mathcal{K}_{\partial D}^{*}$ are

$$
\begin{equation*}
\pm \frac{1}{2}\left(\frac{a-b}{a+b}\right)^{n}, \quad n=1,2, \ldots \tag{2.22}
\end{equation*}
$$

See [11] [29]. So, if the eccentricity of the ellipse becomes large, the spectral bound approaches to $1 / 2$. Other than these examples, the complete spectrum of the NP operator on two discs is obtained in relation to the analysis of gradient concentration [14] 32]. On the other hand, a bound for the essential spectrum has been obtained [35].

## 3 Transmission problems in the free space

The solution $u_{k}$ to (1.6) can be represented as

$$
\begin{equation*}
u_{k}(x)=h(x)+\mathcal{S}_{\partial D}\left[\varphi_{k}\right](x), \quad x \in \mathbb{R}^{d} \tag{3.1}
\end{equation*}
$$

where $\varphi_{k} \in \mathcal{H}(\partial D)$ is the solution to

$$
\begin{equation*}
\left(\lambda I-\mathcal{K}_{\partial D}^{*}\right)\left[\varphi_{k}\right]=\left.\partial_{\nu} h\right|_{\partial D} \quad \text { on } \partial D \tag{3.2}
\end{equation*}
$$

where $\lambda=\lambda(k)$ is defined by (1.5). In fact, (3.2) is a consequence of the jump relation (2.3) and transmission conditions on $\partial D$ :

$$
\begin{equation*}
\left.u\right|_{+}=\left.u\right|_{-},\left.\quad \partial_{\nu} u\right|_{+}=\left.k \partial_{\nu} u\right|_{-} \quad \text { on } \partial D \tag{3.3}
\end{equation*}
$$

Note that the righthand side of (3.2) does not depend on $k$. So, in this case, if $\lambda(k) \notin \sigma\left(\mathcal{K}_{\partial D}^{*}\right)$, then (3.2) is solvable in $\mathcal{H}$, and hence (1.6) has a solution. This is already interesting since a negative $k$ can satisfy the condition $\lambda(k) \notin \sigma\left(\mathcal{K}_{\partial D}^{*}\right)$. For example, any $k$ with

$$
\begin{equation*}
k<-\frac{1+2 b_{\partial D}}{1-2 b_{\partial D}} \quad \text { or } \quad-\frac{1-2 b_{\partial D}}{1+2 b_{\partial D}}<k<0 \tag{3.4}
\end{equation*}
$$

satisfies the condition.
We obtain the following solvability result.
Theorem 3.1 If $\lambda(k) \notin \sigma\left(\mathcal{K}_{\partial D}^{*}\right)$, then for any harmonic function $h$ in $\mathbb{R}^{d}$ there is a unique solution $u_{k}$ to (1.6) satisfying

$$
\begin{equation*}
\left\|\nabla\left(u_{k}-h\right)\right\|_{L^{2}\left(\mathbb{R}^{d}\right)} \leq C_{k}\left\|\partial_{\nu} h\right\|_{-1 / 2(\partial D)} \tag{3.5}
\end{equation*}
$$

for some $C_{k}$ depending on $k$.

Proof. Note that existence of a solution is already proved. To show (3.5), we note that

$$
\begin{aligned}
\left\|\nabla\left(u_{k}-h\right)\right\|_{L^{2}\left(\mathbb{R}^{d}\right)}^{2} & =\int_{D}\left|\nabla \mathcal{S}_{\partial D}\left[\varphi_{k}\right](x)\right|^{2} d x+\int_{\mathbb{R}^{d} \backslash D}\left|\nabla \mathcal{S}_{\partial D}\left[\varphi_{k}\right](x)\right|^{2} d x \\
& =\left.\int_{\partial D} \partial_{\nu} \mathcal{S}_{\partial D}\left[\varphi_{k}\right]\right|_{-} \mathcal{S}_{\partial D}\left[\varphi_{k}\right] d \sigma-\left.\int_{\partial D} \partial_{\nu} \mathcal{S}_{\partial D}\left[\varphi_{k}\right]\right|_{+} \mathcal{S}_{\partial D}\left[\varphi_{k}\right] d \sigma \\
& =-\int_{\partial D} \varphi_{k} \mathcal{S}_{\partial D}\left[\varphi_{k}\right] d \sigma=\left\|\varphi_{k}\right\|_{\mathcal{H}}^{2} .
\end{aligned}
$$

So, (3.5) follows from (3.2).
To show uniqueness of the solution, assume that $u_{k}^{1}$ and $u_{k}^{2}$ are solutions to (1.6). Let $v=u_{k}^{1}-u_{k}^{2}$. Then $v$ is the solution to (1.6) with $h=0$. So we have

$$
\begin{aligned}
0 & =\int_{\mathbb{R}^{d}}\left(\chi\left(\mathbb{R}^{d} \backslash D\right)+k \chi(D)\right)|\nabla v(x)|^{2} d x \\
& =\int_{\mathbb{R}^{d} \backslash D}|\nabla v(x)|^{2} d x+k^{\prime} \int_{D}|\nabla v(x)|^{2} d x+i k^{\prime \prime} \int_{D}|\nabla v(x)|^{2} d x,
\end{aligned}
$$

where $k=k^{\prime}+i k^{\prime \prime}$. So, if $k^{\prime}>0$, or if $k^{\prime} \leq 0$ and $k^{\prime \prime} \neq 0$, then

$$
\int_{D}|\nabla v(x)|^{2} d x=\int_{\mathbb{R}^{d} \backslash D}|\nabla v(x)|^{2} d x=0 .
$$

So, $v$ is constant. Since $v(x) \rightarrow 0$ as $|x| \rightarrow \infty$, we conclude that $v=0$.
Uniqueness for the case $k \leq 0$ (and $\lambda(k) \notin \sigma\left(\mathcal{K}_{\partial D}^{*}\right)$ ) can be proved as a limiting case of $k+i \delta$ as $\delta \rightarrow 0$. For that we need estimate (3.5) with the constant $C$ independent of $\delta$. So the proof will be presented after establishing uniformity results in the following.

For $\epsilon>0$ let

$$
\begin{equation*}
R_{\epsilon}:=\left\{k \in \mathbb{C} \mid \operatorname{dist}\left(\lambda(k), \sigma\left(\mathcal{K}_{\partial D}^{*}\right)\right) \geq \epsilon\right\} . \tag{3.6}
\end{equation*}
$$

It follows from the spectral resolution (2.18) that

$$
\begin{equation*}
\varphi_{k}=\left(\lambda I-\mathcal{K}_{\partial D}^{*}\right)^{-1}\left[\left.\partial_{\nu} h\right|_{\partial D}\right]=\int_{t \in \sigma\left(\mathcal{K}_{\partial D}^{*}\right)} \frac{1}{\lambda-t} d \mathcal{E}(t)\left[\left.\partial_{\nu} h\right|_{\partial D}\right], \tag{3.7}
\end{equation*}
$$

where $\lambda=\lambda(k)$. If $k \in R_{\epsilon}$, then $|\lambda-t| \geq \epsilon$ for all $t \in \sigma\left(\mathcal{K}_{\partial D}^{*}\right)$, and hence we obtain from (3.7) that

$$
\begin{align*}
\left\|\varphi_{k}\right\|_{\mathcal{H}}^{2} & =\int_{t \in \sigma\left(\mathcal{K}_{\partial D}^{*}\right)} \frac{1}{|\lambda-t|^{2}} d\left\langle\mathcal{E}(t)\left[\partial_{\nu} h\right], \partial_{\nu} h\right\rangle_{\mathcal{H}} \\
& \leq \frac{1}{\epsilon^{2}} \int_{t \in \sigma\left(\mathcal{K}_{\partial D}^{*}\right)} d\left\langle\mathcal{E}(t)\left[\partial_{\nu} h\right], \partial_{\nu} h\right\rangle_{\mathcal{H}}=\frac{1}{\epsilon^{2}}\left\|\partial_{\nu} h\right\|_{\mathcal{H}}^{2} . \tag{3.8}
\end{align*}
$$

Thus we obtain the following theorem from (3.1).
Theorem 3.2 For each $\epsilon>0$ there is a constant $C=C_{\epsilon}>0$ such that

$$
\begin{equation*}
\left\|\nabla\left(u_{k}-h\right)\right\|_{L^{2}\left(\mathbb{R}^{d}\right)} \leq C\left\|\partial_{\nu} h\right\|_{-1 / 2(\partial D)} \tag{3.9}
\end{equation*}
$$

for all $k \in R_{\epsilon}$ and harmonic functions $h$ in $\mathbb{R}^{d}$.
We also have the following theorem on Lipschitz dependency of the solution on $k$.

Theorem 3.3 For any $\epsilon>0$ there is a constant $C=C_{\epsilon}>0$ such that

$$
\begin{equation*}
\left\|\nabla\left(u_{k}-u_{s}\right)\right\|_{L^{2}\left(\mathbb{R}^{d}\right)} \leq \frac{C|k-s|}{(1+|k|)(1+|s|)}\left\|\partial_{\nu} h\right\|_{-1 / 2(\partial D)} \tag{3.10}
\end{equation*}
$$

for all $k, s \in R_{\epsilon}$ and harmonic functions $h$ in $\mathbb{R}^{d}$.
Proof. We use the spectral resolution (3.7) to have

$$
\begin{aligned}
\varphi_{k}-\varphi_{s} & =\int_{\sigma\left(\mathcal{K}_{\partial D}^{*}\right)} \frac{1}{\lambda(k)-t} d \mathcal{E}(t)\left[\left.\partial_{\nu} h\right|_{\partial D}\right]-\int_{\sigma\left(\mathcal{K}_{\partial D}^{*}\right)} \frac{1}{\lambda(s)-t} d \mathcal{E}(t)\left[\left.\partial_{\nu} h\right|_{\partial D}\right] \\
& =\int_{\sigma\left(\mathcal{K}_{\partial D}^{*}\right)} \frac{\lambda(s)-\lambda(k)}{(\lambda(k)-t)(\lambda(s)-t)} d \mathcal{E}(t)\left[\left.\partial_{\nu} h\right|_{\partial D}\right]
\end{aligned}
$$

Note that

$$
\frac{\lambda(s)-\lambda(k)}{(\lambda(k)-t)(\lambda(s)-t)}=\frac{4(k-s)}{[(1-2 t) s+1+2 t][(1-2 t) k+1+2 t]}
$$

So, if $k, s \in R_{\epsilon}$, we have

$$
\begin{equation*}
\left|\frac{\lambda(s)-\lambda(k)}{(\lambda(k)-t)(\lambda(s)-t)}\right| \leq C \frac{|k-s|}{(1+|k|)(1+|s|)} \tag{3.11}
\end{equation*}
$$

for all $t \in \sigma\left(\mathcal{K}_{\partial D}^{*}\right)$. So, we have

$$
\left\|\varphi_{k}-\varphi_{s}\right\|_{-1 / 2(\partial D)} \leq \frac{C|k-s|}{(1+|k|)(1+|s|)}\left\|\partial_{\nu} h\right\|_{-1 / 2(\partial D)}
$$

Since

$$
u_{k}(x)-u_{s}(x)=\mathcal{S}_{\partial D}\left[\varphi_{k}-\varphi_{s}\right](x), \quad x \in \Omega
$$

we obtain (3.10).
Continuation of proof of Theorem 3.1. Suppose $k \leq 0$ and $\lambda(k) \notin \sigma\left(\mathcal{K}_{\partial D}^{*}\right)$, and let $v=u_{k}^{1}-u_{k}^{2}$ as before. Choose $\epsilon>0$ so that $k \in R_{\epsilon}$. For $\delta>0$ let $\psi_{\delta}$ be the solution to

$$
\left(\frac{k+1+i \delta}{2(k-1+i \delta)} I-\mathcal{K}_{\partial D}^{*}\right)\left[\psi_{\delta}\right]=\left.\frac{i \delta}{k-1+i \delta} \partial_{\nu} v\right|_{-} \quad \text { on } \partial D
$$

Observe that

$$
\frac{k+1+i \delta}{2(k-1+i \delta)}=\lambda(k+i \delta)
$$

and

$$
\left|\frac{k+1+i \delta}{2(k-1+i \delta)}-t\right| \geq \epsilon \quad \text { for all } t \in \sigma\left(\mathcal{K}_{\partial D}^{*}\right)
$$

In other words, $k+i \delta \in R_{\epsilon}$. Then we see from (3.8) that

$$
\begin{equation*}
\left\|\psi_{\delta}\right\|_{\mathcal{H}} \leq C \delta\left\|\left.\partial_{\nu} v\right|_{-}\right\|_{\mathcal{H}} \tag{3.12}
\end{equation*}
$$

Let

$$
\begin{equation*}
w(x)=v(x)+\mathcal{S}_{\partial D}\left[\psi_{\delta}\right](x), \quad x \in \mathbb{R}^{d} \tag{3.13}
\end{equation*}
$$

Then we can see that $w$ satisfies

$$
\left.w\right|_{+}=\left.w\right|_{-},\left.\quad \partial_{\nu} w\right|_{+}=\left.(k+i \delta) \partial_{\nu} w\right|_{-} \quad \text { on } \partial D
$$

So $w$ is the solution to

$$
\left\{\begin{array}{l}
\nabla \cdot\left(\chi\left(\mathbb{R}^{d} \backslash D\right)+(k+i \delta) \chi(D)\right) \nabla w=0 \quad \text { in } \mathbb{R}^{d} \\
w(x)=O\left(|x|^{1-d}\right) \quad \text { as }|x| \rightarrow \infty
\end{array}\right.
$$

Since the solution for the complex coefficient $k+i \delta$ is unique, we obtain from (3.9) that

$$
\|\nabla w\|_{L^{2}\left(\mathbb{R}^{d}\right)}=0
$$

Thus we obtain from (3.12) that

$$
\|\nabla v\|_{L^{2}\left(\mathbb{R}^{d}\right)}=\left\|\nabla \mathcal{S}_{\partial D}\left[\psi_{\delta}\right]\right\|_{L^{2}\left(\mathbb{R}^{d}\right)} \leq C\left\|\psi_{\delta}\right\|_{\mathcal{H}} \leq C \delta\left\|\left.\partial_{\nu} v\right|_{-}\right\|_{\mathcal{H}}
$$

Letting $\delta \rightarrow 0$, we infer that $v$ is constant in $\mathbb{R}^{d}$. Since $v(x) \rightarrow 0$ as $|x| \rightarrow \infty$, we conclude that $v=0$. This completes the proof.

## 4 Boundary value problem

In this section we consider the boundary value problem (1.2). We first recall a representation formula for a solution $u$ to (1.2). Suppose that (1.2) admits a solution $u \in H^{1}(\Omega)$. Define the harmonic function $h$ in $\Omega$ by

$$
\begin{equation*}
h(x)=-\mathcal{S}_{\partial \Omega}[g](x)+\mathcal{D}_{\partial \Omega}\left[\left.u\right|_{\partial \Omega}\right](x), \quad x \in \Omega \tag{4.1}
\end{equation*}
$$

where the double layer potential $\mathcal{D}_{\partial \Omega}[\varphi]$ is defined by

$$
\begin{equation*}
\mathcal{D}_{\partial \Omega}[\varphi](x):=\frac{1}{\omega_{d}} \int_{\partial \Omega} \frac{(y-x) \cdot \nu_{y}}{|x-y|^{d}} \varphi(y) d \sigma(y), \quad x \in \mathbb{R}^{d} \backslash \partial \Omega \tag{4.2}
\end{equation*}
$$

It is proved in [25] [27] that $u$ is represented as

$$
\begin{equation*}
u(x)=h(x)+\mathcal{S}_{\partial D}[\varphi](x), \quad x \in \Omega \tag{4.3}
\end{equation*}
$$

for some $\varphi \in \mathcal{H}$. In fact, $\varphi$ is given by

$$
\begin{equation*}
\varphi=\left.\partial_{\nu} u\right|_{+}-\left.\partial_{\nu} u\right|_{-} \tag{4.4}
\end{equation*}
$$

and satisfies

$$
\begin{equation*}
\left(\lambda I-\mathcal{K}_{\partial D}^{*}\right)[\varphi]=\left.\partial_{\nu} h\right|_{\partial D} \quad \text { on } \partial D \tag{4.5}
\end{equation*}
$$

where $\lambda=\lambda(k)$ is defined by (1.5).
There is yet another representation for a solution $u$. For $y \in \Omega$ let $N_{y}(x)=N(x, y)$ be the Neumann function on $\Omega$, which is the solution to

$$
\begin{cases}\Delta N_{y}=-\delta_{y} & \text { in } \Omega  \tag{4.6}\\ \partial_{\nu} N_{y}=-|\partial \Omega|^{-1} & \text { on } \partial \Omega \\ \int_{\partial \Omega} N_{y}(x) d \sigma=0\end{cases}
$$

where $\delta_{y}$ is the Dirac mass at $y$ and $|\partial \Omega|$ denotes the area (or the length) of $\partial \Omega$. Using $N(x, y)$ we define

$$
\begin{equation*}
\mathcal{N}_{\partial D}[\varphi](x)=\int_{\partial D} N(x, y) \varphi(y) d \sigma(y), \quad x \in \Omega \tag{4.7}
\end{equation*}
$$

Let $U$ be the solution in absence of an inclusion, that is, the solution to

$$
\left\{\begin{array}{l}
\Delta U=0 \quad \text { in } \Omega,  \tag{4.8}\\
\left.\partial_{\nu} U\right|_{\partial \Omega}=g, \\
\int_{\partial \Omega} U d \sigma=0 .
\end{array}\right.
$$

Then it is proved in [5] (see also [7) that with the same $\varphi$ in (4.4) the following holds:

$$
\begin{equation*}
u(x)=U(x)-\mathcal{N}_{\partial D}[\varphi](x), \quad x \in \partial \Omega . \tag{4.9}
\end{equation*}
$$

We prove the following representation theorem.
Theorem 4.1 Let $u \in H^{1}(\Omega)$ be a solution to (1.2). Then with $\varphi$ in (4.4) it holds that

$$
\begin{equation*}
u(x)=U(x)-\mathcal{N}_{\partial D}[\varphi](x), \quad x \in \Omega . \tag{4.10}
\end{equation*}
$$

Proof. Let

$$
v(x):=U(x)-\mathcal{N}_{\partial D}[\varphi](x), \quad x \in \Omega .
$$

Note that

$$
\partial_{\nu} \mathcal{N}_{\partial D}[\varphi]=-|\partial \Omega|^{-1} \int_{\partial D} \varphi=0 \quad \text { on } \partial \Omega .
$$

So, we have $\partial_{\nu} u=\partial_{\nu} v$ on $\partial \Omega$. So by (4.9) and unique continuation of harmonic functions, we have

$$
u(x)=v(x), \quad x \in \Omega \backslash \bar{D} .
$$

Define $R(x, y)$ by

$$
R(x, y)=N(x, y)+\Gamma(x-y),
$$

and let

$$
\begin{equation*}
\mathcal{R}_{\partial D}[\varphi](x)=\int_{\partial D} R(x, y) \varphi(y) d \sigma(y), \quad x \in \Omega, \tag{4.11}
\end{equation*}
$$

so that

$$
\begin{equation*}
\mathcal{N}_{\partial D}=-\mathcal{S}_{\partial D}+\mathcal{R}_{\partial D} . \tag{4.12}
\end{equation*}
$$

Since $R(x, y)$ is smooth for $x, y \in \Omega$, we have $\left.\partial_{\nu} \mathcal{R}_{\partial D}[\varphi]\right|_{+}=\left.\partial_{\nu} \mathcal{R}_{\partial D}[\varphi]\right|_{-}$on $\partial D$. Therefore we have on $\partial D$

$$
\begin{aligned}
\left.\partial_{\nu} v\right|_{+}-\left.\partial_{\nu} v\right|_{-} & =-\left.\partial_{\nu} \mathcal{N}_{\partial D}[\varphi]\right|_{+}+\left.\partial_{\nu} \mathcal{N}_{\partial D}[\varphi]\right|_{-} \\
& =\left.\partial_{\nu} \mathcal{S}_{\partial D}[\varphi]\right|_{+}-\left.\partial_{\nu} \mathcal{S}_{\partial D}[\varphi]\right|_{-}=\left.\partial_{\nu} u\right|_{+}-\left.\partial_{\nu} u\right|_{-} .
\end{aligned}
$$

So, we have $\left.\partial_{\nu} v\right|_{-}=\left.\partial_{\nu} u\right|_{-}$on $\partial D$. We then infer from unique continuation that $u=v$ in $D$. This completes the proof.

There is a significant difference between representations (4.3) and (4.9): $U$ in (4.9) is independent of $k$, but $h$ in (4.3) depends on $k$ through $\left.u_{k}\right|_{\partial \Omega}$, the Dirichlet data of the solution.

### 4.1 Existence and uniqueness of the solution

The purpose of this subsection is to prove the following theorem.
Theorem 4.2 There is a compact interval on the negative real axis, say $[a, b]$ with $-\infty<a<$ $b<0$, such that for any $k \in \mathbb{C} \backslash[a, b]$ there is a unique solution $u_{k} \in H^{1}(\Omega)$ to (1.2) such that

$$
\begin{equation*}
\left\|u_{k}\right\|_{H^{1}(\Omega)} \leq C_{k}\|g\|_{-1 / 2(\partial \Omega)} \tag{4.13}
\end{equation*}
$$

for some $C_{k}$ independent of $g$ (which may depend on $k$ ).
Proof. Let $U$ be the solution to (4.8). According to Theorem 4.1, unique existence of the solution $u$ to (1.2) amounts to that of $\varphi$ in (4.10), which in turn amounts to unique solvability of the integral equation

$$
\left.k \partial_{\nu} \mathcal{N}_{\partial D}[\varphi]\right|_{-}-\left.\partial_{\nu} \mathcal{N}_{\partial D}[\varphi]\right|_{+}=(k-1) \partial_{\nu} U \quad \text { on } \partial D .
$$

Using (2.3) and (4.12) this equation can be written as

$$
k\left(1 / 2 I-\mathcal{K}_{\partial D}^{*}\right)[\varphi]+k \partial_{\nu} \mathcal{R}_{\partial D}[\varphi]+\left(1 / 2 I+\mathcal{K}_{\partial D}^{*}\right)[\varphi]-\partial_{\nu} \mathcal{R}_{\partial D}[\varphi]=(k-1) \partial_{\nu} U,
$$

or equivalently

$$
\begin{equation*}
\left(\lambda(k) I-\mathcal{K}_{\partial D}^{*}\right)[\varphi]+\partial_{\nu} \mathcal{R}_{\partial D}[\varphi]=\partial_{\nu} U . \tag{4.14}
\end{equation*}
$$

It is convenient to write the above equation as

$$
\begin{equation*}
\left(\lambda(k) I-\widetilde{\mathcal{K}}_{\partial D}^{*}\right)[\varphi]=\partial_{\nu} U \tag{4.15}
\end{equation*}
$$

by putting

$$
\widetilde{\mathcal{K}}_{\partial D}^{*}[\varphi]=\mathcal{K}_{\partial D}^{*}[\varphi]-\partial_{\nu} \mathcal{R}_{\partial D}[\varphi] .
$$

Suppose that $\lambda(k) \in \mathbb{C} \backslash\left[-b_{\partial D}, b_{\partial D}\right]$ where $b_{\partial D}$ be the spectral bound of $\mathcal{K}_{\partial D}^{*}$. Since $\lambda(k) I-$ $\mathcal{K}_{\partial D}^{*}$ is invertible on $\mathcal{H}(\partial D)$ and $\varphi \mapsto \partial_{\nu} \mathcal{R}_{\partial D}[\varphi]$ is a compact operator on $\mathcal{H}(\partial D)$, unique solvability of (4.14) follows from injectivity. To prove injectivity, suppose that

$$
\left(\lambda(k) I-\mathcal{K}_{\partial D}^{*}\right)[\varphi]+\partial_{\nu} \mathcal{R}_{\partial D}[\varphi]=0 \quad \text { on } \partial D .
$$

Let $u(x):=-\mathcal{N}_{\partial D}[\varphi](x), x \in \Omega$. Then, $u$ is a solution to (1.2) with $g=0$. If $k=0$, then $\left.\partial_{\nu} u\right|_{+}=0$ on $\partial D$, and hence $u=0$ in $\Omega \backslash D$. It then follows that $\mathcal{N}_{\partial D}[\varphi]=0$ in $D$. So we have

$$
\varphi=\left.\partial_{\nu} \mathcal{N}_{\partial D}[\varphi]\right|_{-}-\left.\partial_{\nu} \mathcal{N}_{\partial D}[\varphi]\right|_{+}=0
$$

If $k=\infty$, then $u=$ const on $\partial D$. So one can see similarly that $\varphi=0$. In general, it can be seen that

$$
k^{\prime} \int_{D}|\nabla u(x)|^{2} d x+\int_{\Omega \backslash D}|\nabla u(x)|^{2} d x+i k^{\prime \prime} \int_{D}|\nabla u(x)|^{2} d x=0 .
$$

So, if $k^{\prime} \geq 0(k \neq 0)$ or $k^{\prime}<0$ and $k^{\prime \prime} \neq 0$, then $u=0$, and hence $\varphi=0$.
So far, we have shown that (4.1) has a unique solution if $k \in \mathbb{C} \backslash(-\infty, 0)$. Suppose that $|k|>R$ for some $R$ to be determined. Note that $\lambda(\infty)=1 / 2$. We then have

$$
|\lambda(k)-1 / 2|=\left|\frac{k+1}{2(k-1)}-\frac{1}{2}\right|=\frac{1}{|k-1|} \leq \frac{1}{R-1} .
$$

If $R$ is sufficiently large, it follows that

$$
\begin{aligned}
\lambda(k) I-\widetilde{\mathcal{K}}_{\partial D}^{*} & =(\lambda(k)-1 / 2) I+\left(1 / 2 I-\widetilde{\mathcal{K}}_{\partial D}^{*}\right) \\
& =\left(1 / 2 I-\widetilde{\mathcal{K}}_{\partial D}^{*}\right)\left[I+(\lambda(k)-1 / 2)\left(1 / 2 I-\widetilde{\mathcal{K}}_{\partial D}^{*}\right)^{-1}\right],
\end{aligned}
$$

and hence there is $C$ such that

$$
\begin{equation*}
\left\|\left(\lambda(k) I-\widetilde{\mathcal{K}}_{\partial D}^{*}\right)^{-1}[\varphi]\right\|_{1 / 2(\partial D)} \leq C\|\varphi\|_{1 / 2(\partial D)} \tag{4.16}
\end{equation*}
$$

for all $|k|>R$. Similarly, one can show that there is $r<1$ such that (4.16) holds for all $|k|<r$.
Let $a=-R$ and $b=-r$. We have shown that (4.1) has a unique solution if $k \in \mathbb{C} \backslash[a, b]$ for some $-\infty<a<b<0$. Let $\varphi_{k}$ be the solution, and define

$$
u_{k}(x)=U(x)-\mathcal{N}_{\partial D}\left[\varphi_{k}\right](x), \quad x \in \Omega .
$$

Then, $u_{k}$ is the solution to (1.2) satisfying (4.13).
During the course of proof above we proved the following lemma.
Lemma 4.3 There is $C$ (independent of $k$ and $g$ ) such that

$$
\begin{equation*}
\left\|u_{k}\right\|_{H^{1}(\Omega)} \leq C\|g\|_{-1 / 2(\partial \Omega)} \tag{4.17}
\end{equation*}
$$

for all $k$ satisfying $|k|>-a$ or $|k|<-b$.
The following example shows that for the boundary value problem uniqueness may fail, unlike the free space problem in the previous section, even if $\lambda(k) \notin \sigma\left(\mathcal{K}_{\partial D}^{*}\right)$.

Example 1 Let $D$ be the disk centered at 0 of radius $r_{i}$. Then $\mathcal{K}_{\partial D}^{*}=0$. So, the free space problem (1.6) admits a unique solution for all $k \neq-1(k \in \mathbb{C})$. But it is not the case for the boundary value problem. Let $\Omega$ be the disk centered at 0 of radius $r_{e}\left(r_{i}<r_{e}\right)$. Let $\rho:=r_{i} / r_{e}$. For a positive integer n let

$$
k_{n}:=\frac{\rho^{2 n}-1}{\rho^{2 n}+1} .
$$

Then the function $v$, defined by

$$
v(x):= \begin{cases}\left(r^{n}+r_{e}^{2 n} r^{-n}\right) e^{i n \theta}, & \text { if } r_{i}<r=|x| \leq r_{e}, \\ \left(1+\rho^{-2 n}\right) r^{n} e^{i n \theta}, & \text { if } r \leq r_{i},\end{cases}
$$

is a solution to (1.2) with $k=k_{n}$. Note that $\partial_{\nu} v=0$ on $\partial \Omega$. It is quite interesting to observe that $\frac{1}{2} \rho^{2 n}$ is the eigenvalues of the NP operator associated with two interfaces $\partial D$ and $\partial \Omega$ (see [2]).

### 4.2 Uniformity of regularity estimates

In this section and sections to follow we use the condition (1.4). This condition is also required for estimation of the Neumann-to-Dirichlet map (see Lemma4.4). For a positive constant $L>0$ let $S_{L}$ be the infinite sector of $k$ satisfying (1.4) and let $O_{L}$ be the image of $S_{L}$ under the transformation (1.5) (the grey region in the right figure in Figure 1.1). Note that for a given $L>0$ there is $\epsilon>0$ such that $S_{L} \subset R_{\epsilon}$. In fact, we have

$$
\begin{equation*}
\operatorname{dist}\left(O_{L},\left[-b_{\partial D}, b_{\partial D}\right]\right)=2\left(\frac{1}{4}-b_{\partial D}^{2}\right)\left(\sqrt{1+L^{-2}}+\sqrt{\left(2 b_{\partial D}\right)^{2}+L^{-2}}\right)^{-1} \tag{4.18}
\end{equation*}
$$

It is worth mentioning that $\operatorname{dist}\left(O_{L},\left[-b_{\partial D}, b_{\partial D}\right]\right) \rightarrow 0$ if either the spectral bound $b_{\partial D}$ tends to $1 / 2$ or $L \rightarrow 0$.

Lemma 4.4 For each $L>0$ there is a constant $C=C_{L}$ such that

$$
\begin{equation*}
\left\|u_{k}\right\|_{1 / 2(\partial \Omega)} \leq C\|g\|_{-1 / 2(\partial \Omega)} \tag{4.19}
\end{equation*}
$$

for all $k \in S_{L}$ and $g \in H_{0}^{-1 / 2}(\partial \Omega)$.
Proof. Let $\langle,\rangle_{\partial \Omega}$ be the $H^{-1 / 2}-H^{1 / 2}$ pairing on $\partial \Omega$. Note that

$$
\left\langle g, u_{k}\right\rangle_{\partial \Omega}=\int_{\Omega} \gamma_{k}\left|\nabla u_{k}\right|^{2} d x=\int_{\Omega \backslash D}\left|\nabla u_{k}\right|^{2} d x+k^{\prime} \int_{D}\left|\nabla u_{k}\right|^{2} d x+i k^{\prime \prime} \int_{D}\left|\nabla u_{k}\right|^{2} d x
$$

So we have

$$
\begin{equation*}
\left|\left\langle g, u_{k}\right\rangle_{\partial \Omega}\right|^{2}=\left(\left\|\nabla u_{k}\right\|_{L^{2}(\Omega \backslash D)}^{2}+k^{\prime}\left\|\nabla u_{k}\right\|_{L^{2}(D)}^{2}\right)^{2}+k^{\prime \prime 2}\left\|\nabla u_{k}\right\|_{L^{2}(D)}^{4} \tag{4.20}
\end{equation*}
$$

If $k^{\prime} \geq 0$, then

$$
\begin{equation*}
\left\|\nabla u_{k}\right\|_{L^{2}(\Omega \backslash D)}^{2} \leq\left|\left\langle g, u_{k}\right\rangle_{\partial \Omega}\right| \leq\|g\|_{-1 / 2(\partial \Omega)}\left\|u_{k}\right\|_{1 / 2(\partial \Omega)} \tag{4.21}
\end{equation*}
$$

If $k^{\prime}<0$, we rewrite (4.20) as

$$
\begin{equation*}
\left|\left\langle g, u_{k}\right\rangle_{\partial \Omega}\right|^{2}=|k|^{2}\left(\left\|\nabla u_{k}\right\|_{L^{2}(D)}^{2}+\frac{k^{\prime}}{|k|^{2}}\left\|\nabla u_{k}\right\|_{L^{2}(\Omega \backslash D)}^{2}\right)^{2}+\frac{k^{\prime \prime 2}}{|k|^{2}}\left\|\nabla u_{k}\right\|_{L^{2}(\Omega \backslash D)}^{4} \tag{4.22}
\end{equation*}
$$

Since $k^{\prime}<0$, the first term on the righthand side can be 0 . On the other hand, $k^{\prime \prime 2} /|k|^{2}>C$ for some constant $C>0$ if and only if $k^{\prime \prime 2} \geq L^{2} k^{\prime 2}$ for some $L$. So, we obtain that if $k \in S_{L}$, then

$$
\begin{equation*}
\left\|\nabla u_{k}\right\|_{L^{2}(\Omega \backslash D)}^{2} \leq C\|g\|_{-1 / 2(\partial \Omega)}\left\|u_{k}\right\|_{1 / 2(\partial \Omega)} \tag{4.23}
\end{equation*}
$$

for some constant $C$ independent of $k$.
Choose a smooth subdomain $\Omega_{0}$ of $\Omega$ containing $\bar{D}$. For $\eta \in H_{0}^{-1 / 2}(\partial \Omega)$ let $w$ be the solution to

$$
\left\{\begin{array}{l}
\Delta w=0 \quad \text { in } \Omega \backslash \overline{\Omega_{0}} \\
\left.\partial_{\nu} w\right|_{\partial \Omega}=\eta,\left.\quad \partial_{\nu} w\right|_{\partial \Omega_{0}}=0 \\
\int_{\Omega \backslash \Omega_{0}} w(x) d x=0
\end{array}\right.
$$

Then there is a constant $C$ such that

$$
\|\nabla w\|_{L^{2}\left(\Omega \backslash \Omega_{0}\right)} \leq C\|\eta\|_{-1 / 2(\partial \Omega)}
$$

Since

$$
\left\langle\eta, u_{k}\right\rangle_{\partial \Omega}=\int_{\Omega \backslash \Omega_{0}} \nabla w \cdot \nabla \overline{u_{k}} d x
$$

we have

$$
\left|\left\langle\eta, u_{k}\right\rangle_{\partial \Omega}\right| \leq C\left\|\nabla u_{k}\right\|_{L^{2}\left(\Omega \backslash \Omega_{0}\right)}\|\eta\|_{-1 / 2(\partial \Omega)}
$$

and hence

$$
\begin{equation*}
\left\|u_{k}\right\|_{1 / 2(\partial \Omega)} \leq C\left\|\nabla u_{k}\right\|_{L^{2}\left(\Omega \backslash \Omega_{0}\right)} \tag{4.24}
\end{equation*}
$$

Now, (4.19) follows from (4.23) and (4.24).
As a consequence we obtain the following lemma.

Lemma 4.5 For any $L>0$ there is a constant $C=C_{L}$ such that

$$
\begin{equation*}
\left\|u_{k}\right\|_{H^{1}(\Omega)} \leq C\|g\|_{-1 / 2(\partial \Omega)} \tag{4.25}
\end{equation*}
$$

for all $k \in S_{L}$ and $g \in H_{0}^{-1 / 2}(\partial \Omega)$.
Proof. By the same way as in (3.8) we obtain

$$
\begin{equation*}
\left\|\varphi_{k}\right\|_{\mathcal{H}} \leq C\left\|\partial_{\nu} h_{k}\right\|_{\mathcal{H}} \tag{4.26}
\end{equation*}
$$

We then obtain from (2.14) that

$$
\begin{equation*}
\left\|\varphi_{k}\right\|_{-1 / 2(\partial D)} \leq C\left\|\partial_{\nu} h_{k}\right\|_{-1 / 2(\partial D)} \tag{4.27}
\end{equation*}
$$

for some constant $C$ independent of $k \in S_{L}$. Here (and throughout this paper) the constant $C$ may differ at each occurrence.

It follows from (4.1) and (4.19) that

$$
\begin{align*}
\left\|h_{k}\right\|_{H^{1}(\Omega)} & \leq\left\|\mathcal{S}_{\partial \Omega}[g]\right\|_{H^{1}(\Omega)}+\left\|\mathcal{D}_{\partial \Omega}\left[\left.u_{k}\right|_{\partial \Omega}\right]\right\|_{H^{1}(\Omega)} \\
& \leq C\left(\|g\|_{-1 / 2(\partial \Omega)}+\left\|u_{k}\right\|_{1 / 2(\partial \Omega)}\right) \leq C\|g\|_{-1 / 2(\partial \Omega)} \tag{4.28}
\end{align*}
$$

for some $C$ independent of $k$. As a consequence, we obtain

$$
\begin{equation*}
\left\|\partial_{\nu} h_{k}\right\|_{-1 / 2(\partial D)} \leq C\left\|h_{k}\right\|_{H^{1}(\Omega)} \leq C\|g\|_{-1 / 2(\partial \Omega)} \tag{4.29}
\end{equation*}
$$

So, we have from (4.27)

$$
\begin{equation*}
\left\|\mathcal{S}_{\partial D}\left[\varphi_{k}\right]\right\|_{H^{1}(\Omega)} \leq C\left\|\varphi_{k}\right\|_{-1 / 2(\partial D)} \leq C\|g\|_{-1 / 2(\partial \Omega)} \tag{4.30}
\end{equation*}
$$

Then, (4.25) follows from (4.3), (4.28) and (4.30).
Let $a$ and $b$ be numbers appearing in Theorem 4.2 and Lemma 4.3. For $\epsilon>0$, define

$$
\begin{equation*}
Q_{\epsilon}:=\{k \in \mathbb{C}: \operatorname{dist}(k,[a, b])>\epsilon\} \tag{4.31}
\end{equation*}
$$

For a given $\epsilon$ there is $L$ such that

$$
Q_{\epsilon} \subset\{|k|>-a\} \cup\{|k|<-b\} \cup S_{L}
$$

So we obtain the following theorem from Lemma 4.3 and Lemma 4.5 .
Theorem 4.6 For any $\epsilon>0$ there is a constant $C=C_{\epsilon}$ such that

$$
\begin{equation*}
\left\|u_{k}\right\|_{H^{1}(\Omega)} \leq C\|g\|_{-1 / 2(\partial \Omega)} \tag{4.32}
\end{equation*}
$$

for all $k \in Q_{\epsilon}$ and $g \in H_{0}^{-1 / 2}(\partial \Omega)$.

### 4.3 Lipschitz continuity estimate with respect to the conductivity

In this section we investigate the continuous dependency of $u_{k}$ on the conductivity parameter $k$. There has been some work on this problem. If $k$ is real and approaches to $\infty$ (or 0 ), then it was proved that $u_{k}$ converges to $u_{\infty}\left(\right.$ or $\left.u_{0}\right)$ in $H^{1}$-norm $([22])$ and in $L^{\infty}$-norm $([26])$.

The solution $u_{k}$ to (1.2) depends on $k$ analytically. In fact, one can see easily that $\bar{\partial}_{k} u_{k}=\frac{\partial u_{k}}{\partial k}$ is the solution to

$$
\left\{\begin{array}{l}
\nabla \cdot\left(\gamma_{k} \nabla u\right)=0 \quad \text { in } \Omega  \tag{4.33}\\
\left.\partial_{\nu} u\right|_{\partial \Omega}=0 \\
\int_{\partial \Omega} u d \sigma=0
\end{array}\right.
$$

and hence $\bar{\partial}_{k} u_{k}=0$ for $k \in \mathbb{C} \backslash[a, b]$. The following theorem quantifies Lipschitz dependency of the solution on $k$.

Theorem 4.7 For each $\epsilon>0$ there is a constant $C=C_{\epsilon}$ such that

$$
\begin{equation*}
\left\|u_{k}-u_{s}\right\|_{H^{1}(\Omega)} \leq C \frac{|k-s|}{(1+|k|)(1+|s|)}\|g\|_{-1 / 2} \tag{4.34}
\end{equation*}
$$

for all $k, s \in Q_{\epsilon}$ and $g \in H_{0}^{-1 / 2}(\partial \Omega)$.
Before proving Theorem 4.7, we make a brief remark on the expression of (4.34). If $k$ and $s$ are bounded, then it is a Lipschitz continuity estimate. But if $k$ and $s$ are large, then it means more than Lipschitz continuity. It shows, for example, that $u_{k}-u_{s}$ can be arbitrarily small in $H^{1}$-norm if $k$ and $s$ are large.

Proof of Theorem 4.7. Suppose $k, s \in Q_{\epsilon}$. Let us define $\tilde{u}_{s}$ by

$$
\begin{equation*}
\tilde{u}_{s}(x):=h_{k}(x)+\mathcal{S}_{\partial D}\left[\tilde{\varphi}_{s}\right](x), \quad x \in \Omega, \tag{4.35}
\end{equation*}
$$

where $\tilde{\varphi}_{s} \in H_{0}^{-1 / 2}(\partial D)$ is the solution to

$$
\begin{equation*}
\left(\lambda(s) I-\mathcal{K}_{\partial D}^{*}\right)\left[\tilde{\varphi}_{s}\right]=\left.\partial_{\nu} h_{k}\right|_{\partial D} \quad \text { on } \partial D . \tag{4.36}
\end{equation*}
$$

It is worth mentioning that $\tilde{u}_{s}$ is a solution to the equation $\nabla \cdot\left(\gamma_{s} \nabla u\right)=0$ in $\Omega$.
We now compare $u_{k}$ with $\tilde{u}_{s}$, and then $\tilde{u}_{s}$ with $u_{s}$. We use the spectral resolution (3.7) to have

$$
\varphi_{k}-\tilde{\varphi}_{s}=\int_{\sigma\left(\mathcal{K}_{\partial D}^{*}\right)} \frac{1}{\lambda(k)-t} d \mathcal{E}(t)\left[\left.\partial_{\nu} h_{k}\right|_{\partial D}\right]-\int_{\sigma\left(\mathcal{K}_{\partial D}^{*}\right)} \frac{1}{\lambda(s)-t} d \mathcal{E}(t)\left[\left.\partial_{\nu} h_{k}\right|_{\partial D}\right] .
$$

So using the same argument as in the proof of Theorem 3.3 and (4.29), we see that

$$
\begin{equation*}
\left\|\varphi_{k}-\tilde{\varphi}_{s}\right\|_{-1 / 2(\partial D)} \leq \frac{C|k-s|}{(1+|k|)(1+|s|)}\left\|\partial_{\nu} h_{k}\right\|_{-1 / 2(\partial D)} \leq \frac{C|k-s|}{(1+|k|)(1+|s|)}\|g\|_{-1 / 2(\partial \Omega)} \tag{4.37}
\end{equation*}
$$

Since

$$
u_{k}(x)-\tilde{u}_{s}(x)=\mathcal{S}_{\partial D}\left[\varphi_{k}-\tilde{\varphi}_{s}\right](x), \quad x \in \Omega,
$$

we have

$$
\begin{equation*}
\left\|u_{k}-\tilde{u}_{s}\right\|_{H^{1}(\Omega)} \leq \frac{C|k-s|}{(1+|k|)(1+|s|)}\|g\|_{-1 / 2(\partial \Omega)} \tag{4.38}
\end{equation*}
$$

On the other hand, we obtain from Theorem 4.5

$$
\begin{equation*}
\left\|\tilde{u}_{s}-u_{s}\right\|_{H^{1}(\Omega)} \leq C\left\|\partial_{\nu} \tilde{u}_{s}-\partial_{\nu} u_{k}\right\|_{-1 / 2(\partial \Omega)} \tag{4.39}
\end{equation*}
$$

and

$$
\left\|\partial_{\nu} \tilde{u}_{s}-\partial_{\nu} u_{k}\right\|_{-1 / 2(\partial \Omega)}=\left\|\partial_{\nu} \mathcal{S}_{\partial D}\left[\tilde{\varphi}_{s}-\varphi_{k}\right]\right\|_{-1 / 2(\partial \Omega)} \leq C\left\|\tilde{\varphi}_{s}-\varphi_{k}\right\|_{-1 / 2(\partial D)}
$$

where the last inequality holds since there is a distance between $\partial D$ and $\partial \Omega$. We then infer from (4.37) that

$$
\begin{equation*}
\left\|\tilde{u}_{s}-u_{s}\right\|_{H^{1}(\Omega)} \leq \frac{C|k-s|}{(1+|k|)(1+|s|)}\|g\|_{-1 / 2(\partial \Omega)} . \tag{4.40}
\end{equation*}
$$

We obtain (4.34) from (4.38) and (4.40).

### 4.4 Dirichlet problem

Let us briefly mention on the Dirichlet boundary value problem:

$$
\left\{\begin{array}{l}
\nabla \cdot\left(\gamma_{k} \nabla v\right)=0 \quad \text { in } \Omega  \tag{4.41}\\
\left.v\right|_{\partial \Omega}=f \\
\int_{\partial \Omega} v d \sigma=0
\end{array}\right.
$$

The representation formula (4.3) for the solution $v_{k}$ is still valid if we use $h_{k}$ with $g$ replaced with $\left.\partial_{\nu} v_{k}\right|_{\partial \Omega}$ and $\left.u_{k}\right|_{\partial \Omega}$ replaced with $f$ in (4.1). Similarly to Lemma 4.4 one can show that

$$
\left\|\partial_{\nu} v_{k}\right\|_{-1 / 2(\partial \Omega)} \leq C\|f\|_{1 / 2(\partial \Omega)}
$$

for all $k \in Q_{\epsilon}$. So following the same lines of proofs we can obtain results for the Dirichlet problem similar to Theorem 4.5 and Theorem 4.7.

## 5 Uniform validity of the boundary perturbation formula

Let $U$ be the solution to (4.8). Then $u_{k}-U$ can be regarded as a perturbation due to presence of the inclusion with the conductivity $k$. If $D$ is diametrically small, then the perturbation is small and the asymptotic formula as the diameter tends to zero is known. We assume that the inclusion $D$ is represented as

$$
\begin{equation*}
D=\delta B+z \tag{5.1}
\end{equation*}
$$

where $\delta$ represents the small diameter of $D, B$ is a reference domain containing 0 , and $z$ is the location of $D$. Here we also assume that $D$ is away from $\partial \Omega$, namely, there is $c_{0}>0$ such that

$$
\begin{equation*}
\operatorname{dist}(D, \partial \Omega) \geq c_{0} \tag{5.2}
\end{equation*}
$$

This condition is required so that the interaction between $D$ and $\partial \Omega$ does not appear in the boundary perturbation formula in the following. Then, the following asymptotic expansion of the boundary perturbation holds:

$$
\begin{equation*}
u_{k}(x)-U(x)=-\delta^{d} \nabla U(z) \cdot M \nabla_{z} N(x, z)+E_{k}^{1}(x), \quad x \in \partial \Omega, \tag{5.3}
\end{equation*}
$$

where $M=M(k, B)$ is the polarization tensor associated with the conductivity $k$ and the domain $B$ (see (5.8) below for the definition of the polarization tensor), $N(x, z)$ is the Neumann function (4.6), and $E_{k}^{1}(x)$ indicates the error of the approximation and satisfies

$$
\begin{equation*}
\left|E_{k}^{1}(x)\right| \leq C_{k} \delta^{d+1} \tag{5.4}
\end{equation*}
$$

for some $C_{k}$ which may depend on $k$. The formula (5.3) was first discovered in [22] and used effectively for the inverse problem to find the location and/or some geometric properties (especially the equivalent ellipse of the inclusion) of the inclusion using the boundary measurements [10] (16] (19].

The formula has been extended in [5] to include the higher order terms as

$$
\begin{equation*}
u_{k}(x)-U(x)=-\sum_{n=2}^{d+1} \sum_{|\alpha|+|\beta|=n} \frac{\delta^{n+d-2}}{\alpha!\beta!} \partial^{\alpha} U(z) m_{\alpha \beta} \partial_{z}^{\beta} N(x, z)+E_{k}(x), \quad x \in \partial \Omega \tag{5.5}
\end{equation*}
$$

where $m_{\alpha \beta}=m_{\alpha \beta}(k, B)$ is a series of tensors, called generalized polarization tensors (GPTs), associated with $(k, B)$ (see (5.8) below). Here $\alpha$ and $\beta$ are multi-indices. The error of this approximation satisfies

$$
\begin{equation*}
\left|E_{k}(x)\right| \leq C_{k} \delta^{2 d} \tag{5.6}
\end{equation*}
$$

The formula (5.5) has been also used to solve inverse problems for which we refer readers to [6] [7] and references therein. It is worth mentioning that the asymptotic expansion method has been applied in various contexts such as multi-static imaging and bio-medical imaging for which we refer to [4] [8] and references in. The first order formula (5.3) was also generalized to the case when the inclusion is an arbitrary subset of $\Omega$ of low volume fraction [17]. As pointed out in the paper, there may not be a higher order formula in such a case.

Since the conductivity of the inclusion can be extreme (close or equal to 0 or $\infty$ ) or complex, it is important to clarify dependence of the constant $C_{k}$ appearing in (5.4) and (5.6). Recently it is proved in [34] that the approximation formula (5.3) is valid uniformly in $k$ for $k$ real and $0 \leq k \leq \infty$, namely, $C_{k}$ in (5.4) can be chosen independently of $k$. The purpose of this section is to show that the approximation formula (5.5) is valid uniformly in complex $k \in Q_{\epsilon}$. It is worth mentioning that an asymptotic formula for elasticity similar to (5.5) was obtained in [12] and its uniform validity for the real Lamé parameters was proved in [9] (using a variational method).

Let us recall the definition of GPTs associated to the inclusion $B$ with the conductivity $k$. For a given multi-index $\alpha \in \mathbb{N}^{d}$, let $\varphi_{k, \alpha}$ be the solution to

$$
\left(\lambda(k) I-\mathcal{K}_{\partial D}^{*}\right)\left[\varphi_{k, \alpha}\right]=\partial_{\nu} x^{\alpha} \quad \text { on } \partial D .
$$

Here $x^{\alpha}=x_{1}^{\alpha_{1}} \cdots x_{d}^{\alpha_{d}}$. We see through the same estimates as in (3.8) that there is a constant $C$ such that

$$
\begin{equation*}
\left\|\varphi_{k, \alpha}\right\|_{H^{-1 / 2}(\partial B)} \leq C \tag{5.7}
\end{equation*}
$$

for all $k \in Q_{\epsilon}$. The GPT associated with $(k, B)$ is defined by

$$
\begin{equation*}
m_{\alpha \beta}(k, B)=\int_{\partial B} x^{\beta} \varphi_{k, \alpha}(x) d \sigma \tag{5.8}
\end{equation*}
$$

for $\alpha, \beta \in \mathbb{N}^{d}$. Here the integral is understood as the $H^{1 / 2}-H^{-1 / 2}$ pairing. Thus GPTs for the Lipschitz domain $B$ are bounded independently of $k \in Q_{\epsilon}$.

We obtain the following theorem. The proof is nothing but repetition of the arguments to derive (5.5) in [5] except that we need to keep track of the dependency on $k$ of the constants appearing during the derivation using results of this paper. So we omit the proof.

Theorem 5.1 Let $E_{k}$ be the error defined in (5.5). For each $\epsilon$ there is a constant $C=C_{\epsilon}$ such that

$$
\begin{equation*}
\left\|E_{k}\right\|_{L^{\infty}(\partial \Omega)} \leq C \delta^{2 d}\|g\|_{-1 / 2(\partial \Omega)} . \tag{5.9}
\end{equation*}
$$

for all $k \in Q_{\epsilon}$ and $g \in H_{0}^{-1 / 2}(\partial \Omega)$.

## 6 Multiple inclusions

In this section, we consider the case when there are multiple inclusions $D_{1}, \ldots, D_{M}$ whose closures are mutually disjoint. As before, $D_{j}$ 's are simply connected and have Lipschitz boundaries. We assume that the inclusions are at some distance from $\partial \Omega$, that is, there is a constant $c_{0}$ such that

$$
\begin{equation*}
\operatorname{dist}\left(D_{j}, \partial \Omega\right) \geq c_{0} \tag{6.1}
\end{equation*}
$$

Suppose that $D_{j}$ has the complex conductivity $k_{j}(j=1, \ldots, M)$ so that the conductivity distribution is given by

$$
\begin{equation*}
\gamma_{\mathbf{k}}=\chi\left(\Omega \backslash \cup_{j=1}^{M} D_{j}\right)+\sum_{j=1}^{M} k_{j} \chi\left(D_{j}\right) \tag{6.2}
\end{equation*}
$$

Here $\mathbf{k}=\left(k_{1}, \ldots, k_{M}\right)$. We seek conditions on $\mathbf{k}$ which is sufficient for solvability and uniformity of the estimates for the problem

$$
\left\{\begin{array}{l}
\nabla \cdot\left(\gamma_{\mathbf{k}} \nabla u\right)=0 \quad \text { in } \Omega  \tag{6.3}\\
\left.\partial_{\nu} u\right|_{\partial \Omega}=g \\
\int_{\partial \Omega} u d \sigma=0
\end{array}\right.
$$

### 6.1 The NP operator and symmetrization

Let us use notation $\mathcal{S}_{j}, \mathcal{K}_{j}^{*}$ and $\mathcal{N}_{j}$ for $\mathcal{S}_{\partial D_{j}}, \mathcal{K}_{\partial D_{j}}^{*}$ and $\mathcal{N}_{\partial D_{j}}$ to make expressions simpler. Similarly to (4.3), the solution $u_{\mathbf{k}}$ to (6.3) can be uniquely represented as

$$
\begin{equation*}
u_{\mathbf{k}}(x)=h_{\mathbf{k}}(x)+\sum_{j=1}^{M} \mathcal{S}_{j}\left[\varphi_{j}^{\mathbf{k}}\right](x), \quad x \in \Omega \tag{6.4}
\end{equation*}
$$

where the harmonic function $h_{\mathbf{k}}$ is given by (4.1) and $\varphi_{j}^{\mathbf{k}} \in H_{0}^{-1 / 2}\left(\partial D_{j}\right), j=1, \cdots, M$, is the solution to the system of the integral equations

$$
\left(\lambda_{j} I-\mathcal{K}_{j}^{*}\right)\left[\varphi_{j}^{\mathbf{k}}\right]-\sum_{l \neq j} \partial_{\nu_{j}} \mathcal{S}_{l}\left[\varphi_{l}^{\mathbf{k}}\right]=\partial_{\nu_{j}} h_{\mathbf{k}} \quad \text { on } \partial D_{j}, \quad \lambda_{j}=\lambda\left(k_{j}\right)
$$

where $\partial_{\nu_{j}}$ denotes the outward normal derivative on $\partial D_{j}$. Set

$$
\begin{equation*}
\mathbb{D}(\Lambda)=\mathbb{D}\left(\lambda_{1}, \ldots, \lambda_{M}\right):=\operatorname{diag}\left[\lambda_{1}, \ldots, \lambda_{M}\right] \tag{6.5}
\end{equation*}
$$

and

$$
\Phi_{\mathbf{k}}:=\left[\varphi_{1}^{\mathbf{k}}, \ldots, \varphi_{M}^{\mathbf{k}}\right]^{T}, \quad \partial h_{\mathbf{k}}=\left[\left.\partial_{\nu_{1}} h_{\mathbf{k}}\right|_{\partial D_{1}}, \ldots,\left.\partial_{\nu_{M}} h_{\mathbf{k}}\right|_{\partial D_{M}}\right]^{T}(T \text { for transpose })
$$

Then the above system of integral equations can be rewritten in a matrix form

$$
\begin{equation*}
\left(\mathbb{D}(\Lambda)-\mathbb{K}^{*}\right)\left[\Phi_{\mathbf{k}}\right]=\partial h_{\mathbf{k}} \tag{6.6}
\end{equation*}
$$

where

$$
\mathbb{K}^{*}:=\left[\begin{array}{cccc}
\mathcal{K}_{1}^{*} & \partial_{\nu_{1}} \mathcal{S}_{2} & \cdots & \partial_{\nu_{1}} \mathcal{S}_{M}  \tag{6.7}\\
\partial_{\nu_{2}} \mathcal{S}_{1} & \mathcal{K}_{2}^{*} & \cdots & \partial_{\nu_{2}} \mathcal{S}_{M} \\
\vdots & \vdots & \ddots & \vdots \\
\partial_{\nu_{M}} \mathcal{S}_{1} & \partial_{\nu_{M}} \mathcal{S}_{2} & \cdots & \mathcal{K}_{M}^{*}
\end{array}\right]
$$

The operator $\mathbb{K}^{*}$ is the NP operator corresponding to multiple inclusions $\left(D_{1}, \ldots, D_{M}\right)$.
The equation (6.6) holds in the space $\mathbf{H}:=\prod_{j=1}^{M} H_{0}^{-1 / 2}\left(\partial D_{j}\right)$. As in the single interface case, $\mathbb{K}^{*}$ is not self-adjoint on $\mathbf{H}$, but it can be symmetrized by introducing a new inner product on $\mathbf{H}$. Define, for $i, j=1, \ldots, M, \mathcal{S}_{i j}: H_{0}^{-1 / 2}\left(\partial D_{j}\right) \rightarrow H_{0}^{1 / 2}\left(\partial D_{i}\right)$ by

$$
\begin{equation*}
\mathcal{S}_{i j}[\varphi](x):=\mathcal{S}_{j}[\varphi](x), \quad x \in \partial D_{i} \tag{6.8}
\end{equation*}
$$

and $\mathbb{S}$ on $\mathbf{H}$ by

$$
\mathbb{S}:=\left[\begin{array}{cccc}
\mathcal{S}_{11} & \mathcal{S}_{12} & \cdots & \mathcal{S}_{1 M} \\
\mathcal{S}_{21} & \mathcal{S}_{22} & \cdots & \mathcal{S}_{2 M} \\
\vdots & \vdots & \ddots & \vdots \\
\mathcal{S}_{M 1} & \mathcal{S}_{M 2} & \cdots & \mathcal{S}_{M M}
\end{array}\right] .
$$

The twisted inner product (2.6) is now extended to

$$
\begin{equation*}
\langle\Phi, \Psi\rangle_{\mathcal{H}}:=-\langle\Phi, \mathbb{S}[\Psi]\rangle, \quad \Phi, \Psi \in \mathbf{H} . \tag{6.9}
\end{equation*}
$$

Then one can show that
(i) $\langle\Phi, \Psi\rangle_{\mathcal{H}}$ is an inner product on $\mathbf{H}$.
(ii) The symmetrization principle (2.7) can be extended as

$$
\begin{equation*}
\mathbb{S K}^{*}=\mathbb{K} \mathbb{S}, \tag{6.10}
\end{equation*}
$$

where $\mathbb{K}$ is the adjoint of $\mathbb{K}^{*}$ with respect to $L^{2}$ inner product.
(iii) The spectrum of $\mathbb{K}^{*}$ on $\mathbf{H}$ lies in $(-1 / 2,1 / 2)$.

In fact, these facts were proved in [2] where there are two interfaces. Following the same lines of proofs there one can extend them to the case of multiple inclusions. Let $\|\cdot\|_{\mathcal{H}}$ be the norm induced by the inner product $\langle\cdot, \cdot\rangle_{\mathcal{H}}$. Then we have the following lemma.

Lemma 6.1 There are constants $C_{1}$ and $C_{2}$ such that

$$
\begin{equation*}
C_{1}\|\Phi\|_{\mathbf{H}} \leq\|\Phi\|_{\mathcal{H}} \leq C_{2}\|\Phi\|_{\mathbf{H}} \quad \text { for all } \Phi \in \mathbf{H} . \tag{6.11}
\end{equation*}
$$

Proof. Similarly to the operator $A$ in Lemma 2.1 we define the operator

$$
A: \prod_{j=1}^{M} H^{-1 / 2}\left(\partial D_{j}\right) \times \mathbb{C}^{M} \rightarrow \prod_{j=1}^{M} H^{1 / 2}\left(\partial D_{j}\right) \times \mathbb{C}^{M}
$$

by

$$
\begin{equation*}
A\left(\Phi,\left(a_{1}, \ldots, a_{M}\right)\right):=\left(\mathbb{S}[\Phi]+\left(a_{1}, \ldots, a_{M}\right),\left(\left\langle\varphi_{1}, 1\right\rangle_{1}, \ldots,\left\langle\varphi_{M}, 1\right\rangle_{M}\right)\right) \tag{6.12}
\end{equation*}
$$

where $\Phi=\left[\varphi_{1}, \ldots, \varphi_{M}\right]^{T}$ and $\left\langle\varphi_{j}, 1\right\rangle_{j}$ denotes $H^{-1 / 2}-H^{1 / 2}$ duality pairing on $\partial D_{j}$. One can show following the same lines of the proof of Lemma 2.1 that $A$ is invertible. Then (6.11) follows by the same argument as in section 2,

### 6.2 A complete characterization of solvability for two disks

Suppose that multiple inclusions consist of two disks $D_{1}$ and $D_{2}$ of the same radius. Suppose that the conductivity of $D_{j}$ is $k_{j}$ for $j=1,2$ so that the conductivity distribution is given by

$$
\begin{equation*}
\gamma=k_{1} \chi\left(D_{1}\right)+k_{2} \chi\left(D_{2}\right)+\chi\left(\mathbb{R}^{2} \backslash\left(D_{1} \cup D_{2}\right)\right) . \tag{6.13}
\end{equation*}
$$

We consider solvability in terms of $k_{1}$ and $k_{2}$ of the problem

$$
\begin{cases}\nabla \cdot \gamma \nabla u=0 & \text { in } \mathbb{R}^{2}  \tag{6.14}\\ u(\mathbf{x})-h(\mathbf{x})=O\left(|\mathbf{x}|^{-1}\right) & \text { as }|\mathbf{x}| \rightarrow \infty\end{cases}
$$

where $h$ is a given harmonic function in $\mathbb{R}^{2}$.
For that purpose we use the bipolar coordinates. The bipolar coordinate system $(\xi, \theta) \in$ $\mathbb{R} \times(-\pi, \pi]$ is defined as follows (see [33]): for a fixed $\alpha>0$

$$
\begin{equation*}
x=\alpha \frac{\sinh \xi}{\cosh \xi-\cos \theta} \quad \text { and } \quad y=\alpha \frac{\sin \theta}{\cosh \xi-\cos \theta} . \tag{6.15}
\end{equation*}
$$

With an appropriately chosen $\alpha$ two disks $D_{1}$ and $D_{2}$ are given in terms of bipolar coordinates by

$$
\begin{equation*}
D_{1}=\left\{\xi<-\xi_{0}\right\}, \quad D_{2}=\left\{\xi>\xi_{0}\right\} \tag{6.16}
\end{equation*}
$$

for some $\xi_{0}>0$. Recall that the normal derivative of a function $u$ on $\partial D_{j}$ are given by

$$
\begin{equation*}
\left.\frac{\partial u}{\partial \nu}\right|_{\xi=c}=-\left.\operatorname{sgn}(c) J(c, \theta) \frac{\partial u}{\partial \xi}\right|_{\xi=c}, \tag{6.17}
\end{equation*}
$$

where

$$
J(\xi, \theta):=\frac{\cosh \xi-\cos \theta}{\alpha} .
$$

The bipolar coordinate system admits a general separation of variables solution to the harmonic function $u$ as follows:

$$
\begin{equation*}
u(\xi, \theta)=a_{0}+b_{0} \xi+c_{0} \theta+\sum_{n=1}^{\infty}\left[\left(a_{n} e^{n \xi}+b_{n} e^{-n \xi}\right) \cos n \theta+\left(c_{n} e^{n \xi}+d_{n} e^{-n \xi}\right) \sin n \theta\right] \tag{6.18}
\end{equation*}
$$

where $a_{n}, b_{n}, c_{n}$ and $d_{n}$ are constants. Suppose that the expansion of the harmonic function $h$ is given by

$$
h(\xi, \theta)=f_{0}+ \begin{cases}\sum_{n \neq 0} f_{n} e^{-|n| \xi} e^{i n \theta}, & \xi>0,  \tag{6.19}\\ \sum_{n \neq 0} g_{n} e^{|n| \xi} e^{i n \theta}, & \xi<0\end{cases}
$$

for some coefficients $f_{n}$ and $g_{n}$. Then one can see using the transmission conditions $\left(\left.u\right|_{+}=\left.u\right|_{-}\right.$ and $\left.\partial_{\nu} u\right|_{+}=\left.k_{j} \partial_{\nu} u\right|_{-}$) on the interfaces that the solution to (6.14) is given by

$$
(u-h)(\xi, \theta)=C+ \begin{cases}\sum_{n \neq 0}\left(a_{n} e^{|n| \xi}+b_{n} e^{-|n| \xi}\right) e^{i n \theta}, & x \in \mathbb{R}^{2} \backslash\left(D_{1} \cup D_{2}\right),  \tag{6.20}\\ \sum_{n \neq 0}\left(a_{n} e^{|n| \xi}+b_{n} e^{|n|\left(2 \xi_{0}+\xi\right)}\right) e^{i n \theta}, & x \in D_{1}, \\ \sum_{n \neq 0}\left(a_{n} e^{|n|\left(2 \xi_{0}-\xi\right)}+b_{n} e^{-|n| \xi}\right) e^{i n \theta}, & x \in D_{2},\end{cases}
$$

where

$$
\begin{equation*}
a_{n}=\frac{g_{n}-\tau_{1}^{-1} e^{2|n| \xi_{0}} f_{n}}{\tau^{-1} e^{4|n| \xi_{0}}-1}, \quad b_{n}=\frac{f_{n}-\tau_{2}^{-1} e^{2|n| \xi_{0}} g_{n}}{\tau^{-1} e^{4|n| \xi_{0}}-1} \tag{6.21}
\end{equation*}
$$

and

$$
C=-\sum_{n \neq 0}\left(a_{n}+b_{n}\right)
$$

Here, we set

$$
\begin{equation*}
\tau=\tau_{1} \tau_{2}, \quad \tau_{j}=\frac{k_{j}-1}{k_{j}+1}, \quad j=1,2 \tag{6.22}
\end{equation*}
$$

One can see that the equation (6.14) is solvable if and only if the denominator in (6.21) is non-zero, namely,

$$
\tau^{-1} e^{4 n \xi_{0}}-1 \neq 0, \quad n=1,2, \ldots
$$

This condition can be rewritten as

$$
\begin{equation*}
\lambda\left(k_{1}\right) \lambda\left(k_{2}\right) \neq\left( \pm \frac{1}{2} e^{-2 n \xi_{0}}\right)^{2}, \quad n=1,2, \ldots \tag{6.23}
\end{equation*}
$$

We emphasize that $\pm \frac{1}{2} e^{-2 n \xi_{0}}$ are eigenvalues of the NP operator for two disks (see [32]).
Note that

$$
\begin{equation*}
V_{n}:=\left\{\left(k_{1}, k_{2}\right) \in \mathbb{C}^{2} \left\lvert\, \lambda\left(k_{1}\right) \lambda\left(k_{2}\right)=\left(\frac{1}{2} e^{-2 n \xi_{0}}\right)^{2}\right.\right\}, \quad n=1,2, \ldots, \tag{6.24}
\end{equation*}
$$

is a two dimensional algebraic variety in $\mathbb{C}^{2}$. So, the insolvability region in presence of two disks is a countable family of two dimensional varieties in $\mathbb{C}^{2}$. It suggests that it is difficult to find a solvability condition for the free space problem or the boundary value problem when complex conductivities of inclusions are different. However, since $\lambda$ maps the right half plane onto the outside of the disk of radius $1 / 2$ (centered at 0 ) and $\left(\frac{1}{2} e^{-2 n \xi_{0}}\right)^{2}<1 / 4$, we see that if $k_{1}^{\prime} \geq 0$ and $k_{2}^{\prime} \geq 0$ then $\left(k_{1}, k_{2}\right) \notin V_{n}$. This is the standard case of elliptic equations. On the other hand, $\lambda$ maps the real line onto the real line and $\left(\frac{1}{2} e^{-2 n \xi_{0}}\right)^{2}>0$, and hence if $k_{1}^{\prime \prime}$ and $k_{2}^{\prime \prime}$ have the same signs, then $\left(k_{1}, k_{2}\right) \notin V_{n}$. In the next section we show the problem (6.3) in presence of multiple inclusions is well-posed if the imaginary parts of $k_{j}$ have the same signs.

### 6.3 Uniqueness and existence

The purpose of this section is to present a condition on $k_{j}$ under which the problem (6.3) is well-posed. The condition is that the imaginary parts of $k_{j}$ 's are either all positive or all negative (they can be zero if the real part is non-negative). It is worth mentioning that if the real parts of $k_{j}$ are all non-negative, (6.3) is elliptic and well-posed.

Let

$$
S_{L}^{+}:=S_{L} \cap\left\{k^{\prime}+i k^{\prime \prime}: k^{\prime} \in \mathbb{R}, k^{\prime \prime} \geq 0\right\}
$$

where $S_{L}$ is the sector defined in subsection 4.2, Let

$$
\begin{equation*}
\Theta^{+}:=\bigcup_{L>0}\left(S_{L}^{+}\right)^{M} . \tag{6.25}
\end{equation*}
$$

It is worth mentioning that $\Theta^{+}$is the collection of all $\mathbf{k}=\left(k_{1}, \ldots, k_{M}\right) \in \mathbb{C}^{M}$ with $k_{j} \in$ $\mathbb{C} \backslash(-\infty, 0)$ and $k_{j}^{\prime \prime} \geq 0$ (for all $j$ ). We define $\Theta^{-}$likewise, and let

$$
\begin{equation*}
\Theta:=\Theta^{+} \cup \Theta^{-} . \tag{6.26}
\end{equation*}
$$

We have the following well-posedness result for $\mathbf{k} \in \Theta$.
Theorem 6.2 If $\mathbf{k}=\left(k_{1}, \ldots, k_{M}\right) \in \Theta$, then $\mathbb{D}(\Lambda)-\mathbb{K}^{*}$ is invertible on $\mathbf{H}$, where $\Lambda=$ $\left(\lambda_{1}, \ldots, \lambda_{M}\right)=\left(\lambda\left(k_{1}\right), \ldots, \lambda\left(k_{M}\right)\right)$.

Proof. We realize $\mathbb{D}(\Lambda)-\mathbb{K}^{*}$ as a compact perturbation of invertible operator. In fact, we write

$$
\mathbb{D}(\Lambda)-\mathbb{K}^{*}=\operatorname{diag}\left[\lambda_{1} I-\mathcal{K}_{1}^{*}, \ldots, \lambda_{M} I-\mathcal{K}_{M}^{*}\right]-\left[\begin{array}{cccc}
0 & \partial_{\nu_{1}} \mathcal{S}_{2} & \cdots & \partial_{\nu_{1}} \mathcal{S}_{M} \\
\partial_{\nu_{2}} \mathcal{S}_{1} & 0 & \cdots & \partial_{\nu_{2}} \mathcal{S}_{M} \\
\vdots & \vdots & \ddots & \vdots \\
\partial_{\nu_{M}} \mathcal{S}_{1} & \partial_{\nu_{M}} \mathcal{S}_{2} & \cdots & 0
\end{array}\right] .
$$

Since there is a positive distance between inclusions, the second operator on the righthand side above is a compact operator on $\mathbf{H}$, while the first operator is invertible. So to show invertibility of $\mathbb{D}(\Lambda)-\mathbb{K}^{*}$ it suffices to prove its injectivity by the Fredholm alternative. Suppose that $\Phi=\left[\varphi_{1}, \ldots, \varphi_{M}\right]^{T}$ satisfies $\left(\mathbb{D}(\Lambda)-\mathbb{K}^{*}\right)[\Phi]=0$. Then we can see that the function $u$, defined by

$$
u(x)=\sum_{j=1}^{M} \mathcal{S}_{j}\left[\varphi_{j}\right](x), \quad x \in \mathbb{R}^{d}
$$

is a solution to

$$
\left\{\begin{array}{l}
\nabla \cdot\left(\sum_{j=1}^{M} k_{j} \chi\left(D_{j}\right)+\chi\left(\mathbb{R}^{d} \backslash D_{j}\right)\right) \nabla u=0 \quad \text { in } \mathbb{R}^{d} \\
u(x)=O\left(|x|^{-d+1}\right) \quad \text { as }|x| \rightarrow \infty
\end{array}\right.
$$

So, we have

$$
\begin{equation*}
\int_{\mathbb{R}^{d} \backslash \cup_{j=1}^{M} D_{j}}|\nabla u(x)|^{2} d x+\sum_{k_{j}^{\prime} \neq 0} k_{j}^{\prime} \int_{D_{j}}|\nabla u(x)|^{2} d x+i \sum_{j=1}^{M} k_{j}^{\prime \prime} \int_{D_{j}}|\nabla u(x)|^{2} d x=0 . \tag{6.27}
\end{equation*}
$$

In particular, we have

$$
\begin{equation*}
\sum_{j=1}^{M} k_{j}^{\prime \prime} \int_{D_{j}}|\nabla u(x)|^{2} d x=0 \tag{6.28}
\end{equation*}
$$

Suppose that $\mathbf{k} \in \Theta^{+}$. Since $k_{j}^{\prime \prime} \geq 0$ for all $j$ and $k_{j}^{\prime \prime}>0$ if $k_{j}^{\prime}<0$, we infer from (6.28) that $u$ is constant in $D_{j}$ if $k_{j}^{\prime}<0$. Then (6.27) becomes

$$
\int_{\mathbb{R}^{d} \backslash \cup_{j=1}^{M} D_{j}}|\nabla u(x)|^{2} d x+\sum_{k_{j}^{\prime}>0} k_{j}^{\prime} \int_{D_{j}}|\nabla u(x)|^{2} d x=0,
$$

from which we infer that $u$ is constant in $\mathbb{R}^{d} \backslash \cup_{j=1}^{M} D_{j}$. Since $u(x) \rightarrow 0$ as $|x| \rightarrow \infty$, we have $u(x)=0$ in $\mathbb{R}^{d} \backslash \cup_{j=1}^{M} D_{j}$. By the transmission condition (continuity of the potential) on $\partial D_{j}$, we have $u=0$ in $D_{j}, j=1, \ldots, M$. Thus we have $\varphi_{j}=\left.\partial_{\nu} u\right|_{+}-\left.\partial_{\nu} u\right|_{-}=0$ on $\partial D_{j}$. This completes the proof.

The following is the main theorem of this subsection.
Theorem 6.3 For all $\mathbf{k}=\left(k_{1}, \ldots, k_{M}\right) \in \Theta$ there is a unique solution $u_{\mathbf{k}} \in H^{1}(\Omega)$ to (6.3), and there is a unique $\psi_{j}^{\mathbf{k}} \in H_{0}^{-1 / 2}\left(\partial D_{j}\right), j=1, \ldots, M$, such that $u_{\mathbf{k}}$ is represented as

$$
\begin{equation*}
u_{\mathbf{k}}(x)=U(x)-\sum_{j=1}^{M} \mathcal{N}_{j}\left[\psi_{j}^{\mathbf{k}}\right](x), \quad x \in \Omega, \tag{6.29}
\end{equation*}
$$

where $U$ is the solution to (4.8).
Proof. Uniqueness of the solution $u_{\mathbf{k}}$ can be proved similarly to Theorem 6.2 Existence and uniqueness of $\psi_{j}^{\mathbf{k}}$ can be proved in a similar way to the proof of Theorem 4.2 using Theorem 6.2

### 6.4 Uniformity of estimates and Lipschitz dependency

We want to extend uniformity results in previous sections to the case of multiple inclusions. If $k_{1}, \ldots, k_{M}$ are identical, so are $\lambda_{1}, \ldots, \lambda_{M}$, say $\lambda$. Then, the equation (6.6) takes the form

$$
\left(\lambda \mathbf{I}-\mathbb{K}^{*}\right)\left[\Phi_{\mathbf{k}}\right]=\partial h_{\mathbf{k}},
$$

where $\mathbf{I}$ is the identity operator, and the spectral resolution can be applied as in section 4.2, However, if $\lambda_{i} \neq \lambda_{j}$ for some $i, j$, we do not have the spectral resolution and arguments of the previous sections cannot be applied. We overcome this difficulty by induction.

For $\mathbf{k}=\left(k_{1}, \ldots, k_{M}\right)$, let $\lambda(\mathbf{k}):=\left(\lambda\left(k_{1}\right), \ldots, \lambda\left(k_{M}\right)\right)$ where $\lambda$ is the bilinear transform defined by (1.5). For $L>0$ set $\Theta_{L}:=\left(S_{L}^{+}\right)^{M}$ and

$$
\begin{equation*}
\Sigma_{L}:=\left\{\Lambda=\lambda(\mathbf{k}) \mid \mathbf{k} \in \Theta_{L}\right\} . \tag{6.30}
\end{equation*}
$$

We first prove the following lemma.
Lemma 6.4 For any $R>0$ and $L>0$ there is a constant $C=C(L, R)$ such that

$$
\begin{equation*}
\left\|\left(\mathbf{I}-\mathbb{D}(\Lambda)^{-1} \mathbb{K}^{*}\right)^{-1}\right\| \leq C, \quad\left\|\left(\mathbf{I}-\mathbb{K}^{*} \mathbb{D}(\Lambda)^{-1}\right)^{-1}\right\| \leq C \tag{6.31}
\end{equation*}
$$

for any $\Lambda=\left(\lambda_{1}, \ldots, \lambda_{M}\right) \in \Sigma_{L}$ with $|\Lambda| \leq R$. Here $\|\cdot\|$ is the operator norm on $\mathbf{H}$.
Proof. For $\Lambda \in \Sigma_{L}, \mathbb{D}(\Lambda)-\mathbb{K}^{*}$ is invertible on $\mathbf{H}$. So,

$$
\begin{equation*}
f(\Lambda):=\left\|\left(\mathbf{I}-\mathbb{D}(\Lambda)^{-1} \mathbb{K}^{*}\right)^{-1}\right\| \tag{6.32}
\end{equation*}
$$

is well-defined. Since $\Sigma_{L}$ is a closed set, it is sufficient to show that $f$ is continuous to obtain the first estimate in (6.31).

Fix $\Lambda \in \Sigma_{L}$ and let $\left\{\Lambda_{n}\right\}$ be a sequence in $\Sigma_{L}$ converging to $\Lambda$. Note that

$$
\begin{align*}
\left|f\left(\Lambda_{n}\right)-f(\Lambda)\right| & \leq\left\|\left(\mathbf{I}-\mathbb{D}\left(\Lambda_{n}\right)^{-1} \mathbb{K}^{*}\right)^{-1}-\left(\mathbf{I}-\mathbb{D}(\Lambda)^{-1} \mathbb{K}^{*}\right)^{-1}\right\| \\
& \leq\left\|\left(\mathbf{I}-\mathbb{D}\left(\Lambda_{n}\right)^{-1} \mathbb{K}^{*}\right)^{-1}\right\|\left\|\left(\mathbb{D}(\Lambda)^{-1}-\mathbb{D}\left(\Lambda_{n}\right)^{-1}\right) \mathbb{K}^{*}\right\|\left\|\left(\mathbf{I}-\mathbb{D}(\Lambda)^{-1} \mathbb{K}^{*}\right)^{-1}\right\| \\
& =f\left(\Lambda_{n}\right)\left\|\left(\mathbb{D}(\Lambda)^{-1}-\mathbb{D}\left(\Lambda_{n}\right)^{-1}\right) \mathbb{K}^{*}\right\| f(\Lambda) . \tag{6.33}
\end{align*}
$$

Since $\left\|\left(\mathbb{D}(\Lambda)^{-1}-\mathbb{D}\left(\Lambda_{n}\right)^{-1}\right) \mathbb{K}^{*}\right\| \rightarrow 0$ as $n \rightarrow \infty$, we have

$$
\left\|\left(\mathbb{D}(\Lambda)^{-1}-\mathbb{D}\left(\Lambda_{n}\right)^{-1}\right) \mathbb{K}^{*}\right\| f(\Lambda)<\frac{1}{2}
$$

for all sufficiently large $n$, and hence we obtain from (6.33) and the triangular inequality that

$$
f\left(\Lambda_{n}\right)<2 f(\Lambda)
$$

Inserting it into (6.33), it now follows

$$
\left|f\left(\Lambda_{n}\right)-f(\Lambda)\right| \leq 2\left\|\left(\mathbb{D}(\Lambda)^{-1}-\mathbb{D}\left(\Lambda_{n}\right)^{-1}\right) \mathbb{K}^{*}\right\| f(\Lambda)^{2}
$$

Therefore $f$ is continuous on $\Sigma_{L}$. Similarly, we obtain the second estimate in (6.31).
Theorem 6.5 For any $L>0$ there is a constant $C=C(L)$ such that

$$
\begin{equation*}
\left\|\left(\mathbf{I}-\mathbb{D}(\Lambda)^{-1} \mathbb{K}^{*}\right)^{-1}\right\| \leq C, \quad\left\|\left(\mathbf{I}-\mathbb{K}^{*} \mathbb{D}(\Lambda)^{-1}\right)^{-1}\right\| \leq C \tag{6.34}
\end{equation*}
$$

for all $\Lambda \in \Sigma_{L}$.

Proof. We only prove the first estimate in (6.34) because the second one can be proved similarly. We use induction on $M$, the number of inclusions.

If $M=1$, (2.18) shows

$$
\begin{equation*}
\left(I-\lambda_{1}^{-1} \mathcal{K}_{1}^{*}\right)^{-1}=\int_{\sigma\left(\mathcal{K}_{1}^{*}\right)} \frac{1}{1-\lambda_{1}^{-1} t} d \mathcal{E}_{1}(t) \tag{6.35}
\end{equation*}
$$

where $\mathcal{E}_{1}(t)$ is the resolution of identity for $\mathcal{K}_{1}^{*}$ on $\partial D_{1}$. Since $\left|1-\lambda_{1}^{-1} t\right| \geq C$, we see that $\left\|\left(I-\lambda_{1}^{-1} \mathcal{K}_{\partial D_{1}}^{*}\right)^{-1}\right\| \leq C$ uniformly in $\lambda_{1} \in O_{L}$. Assume $M>1$. Without loss of generality, we may assume that $\left|\lambda_{1}\right| \leq \cdots \leq\left|\lambda_{M}\right|$ by reordering $\lambda_{j}$ 's. Let

$$
\Lambda_{M}:=\left(\lambda_{1}, \ldots, \lambda_{M}\right) \quad \text { and } \quad \Lambda_{M-1}:=\left(\lambda_{1}, \ldots, \lambda_{M-1}\right) .
$$

Denote by $\mathbb{K}_{M}^{*}$ the NP operator given in (6.7) when there are $M$ inclusions. Then it takes the following form:

$$
\mathbb{K}_{M}^{*}=\left[\begin{array}{cccc} 
& & & \partial_{\nu_{1}} \mathcal{S}_{M} \\
& \mathbb{K}_{M-1}^{*} & & \partial_{\nu_{2}} \mathcal{S}_{M} \\
& & & \vdots \\
\partial_{\nu_{M}} \mathcal{S}_{1} & \partial_{\nu_{M}} \mathcal{S}_{2} & \cdots & \mathcal{K}_{M}^{*}
\end{array}\right]=:\left[\begin{array}{cc}
\mathbb{K}_{M-1}^{*} & \mathbb{E} \\
\mathbb{F} & \mathcal{K}_{M}^{*}
\end{array}\right]
$$

We emphasize that $\mathbb{E}$ is a bounded operator from $H^{-1 / 2}\left(\partial D_{M}\right)$ into $H^{-1 / 2}\left(\partial D_{1}\right) \times \ldots \times$ $H^{-1 / 2}\left(\partial D_{M-1}\right)$ and $\mathbb{F}$ is from $H^{-1 / 2}\left(\partial D_{1}\right) \times \ldots \times H^{-1 / 2}\left(\partial D_{M-1}\right)$ into $H^{-1 / 2}\left(\partial D_{M}\right)$. Let

$$
\Phi_{M}=\left[\begin{array}{c}
\varphi_{1} \\
\vdots \\
\varphi_{M}
\end{array}\right], \quad \Phi_{M-1}=\left[\begin{array}{c}
\varphi_{1} \\
\vdots \\
\varphi_{M-1}
\end{array}\right], \quad \Psi_{M}=\left[\begin{array}{c}
\psi_{1} \\
\vdots \\
\psi_{M}
\end{array}\right], \quad \Psi_{M-1}=\left[\begin{array}{c}
\psi_{1} \\
\vdots \\
\psi_{M-1}
\end{array}\right]
$$

and suppose that the following integral equation holds:

$$
\begin{equation*}
\left(\mathbf{I}-\mathbb{D}\left(\Lambda_{M}\right)^{-1} \mathbb{K}_{M}^{*}\right)\left[\Phi_{M}\right]=\Psi_{M} . \tag{6.36}
\end{equation*}
$$

Then we may rewrite this equation as

$$
\begin{align*}
\left(\mathbf{I}-\mathbb{D}\left(\Lambda_{M-1}\right)^{-1} \mathbb{K}_{M-1}^{*}\right)\left[\Phi_{M-1}\right]-\mathbb{D}\left(\Lambda_{M-1}\right)^{-1} \mathbb{E}\left[\varphi_{M}\right] & =\Psi_{M-1},  \tag{6.37}\\
-\lambda_{M}^{-1} \mathbb{F}\left[\Phi_{M-1}\right]+\left(I-\lambda_{M}^{-1} \mathcal{K}_{M}^{*}\right)\left[\varphi_{M}\right] & =\psi_{M} .
\end{align*}
$$

Solving the first equation for $\Phi_{M-1}$ and substituting it into the second equation yield

$$
\begin{aligned}
& \left(I-\lambda_{M}^{-1} \mathcal{K}_{M}^{*}\right)\left[\varphi_{M}\right]-\lambda_{M}^{-1} \mathbb{F}\left(\mathbf{I}-\mathbb{D}\left(\Lambda_{M-1}\right)^{-1} \mathbb{K}_{M-1}^{*}\right)^{-1} \mathbb{D}\left(\Lambda_{M-1}\right)^{-1} \mathbb{E}\left[\varphi_{M}\right] \\
& =\psi_{M}+\lambda_{M}^{-1} \mathbb{F}\left(\mathbf{I}-\mathbb{D}\left(\Lambda_{M-1}\right)^{-1} \mathbb{K}_{M-1}^{*}\right)^{-1}\left[\Psi_{M-1}\right] .
\end{aligned}
$$

The induction hypothesis is

$$
\begin{equation*}
\left\|\left(I-\lambda_{M}^{-1} \mathcal{K}_{M}^{*}\right)^{-1}\right\| \leq C, \quad\left\|\left(\mathbf{I}-\mathbb{D}\left(\Lambda_{M-1}\right)^{-1} \mathbb{K}_{M-1}^{*}\right)^{-1}\right\| \leq C \tag{6.38}
\end{equation*}
$$

So, we infer that

$$
\left\|\varphi_{M}\right\| \leq C_{1} \lambda_{M}^{-1}\left\|\varphi_{M}\right\|+C_{2}\left\|\psi_{M}\right\|+C_{3} \lambda_{M}^{-1}\left\|\Psi_{M-1}\right\|
$$

for some constants $C_{1}, C_{2}$, and $C_{3}$ independent of $\Lambda$. Here the norms $\|\cdot\|$ are $H^{-1 / 2}$-norms on appropriate boundaries $\partial D_{j}$. For example, $\left\|\varphi_{M}\right\|$ is on $\partial D_{M}$ and $\left\|\Psi_{M-1}\right\|$ is on $\partial D_{1} \times \cdots \times$ $\partial D_{M-1}$. So, if $\left|\lambda_{M}\right| \geq R$ for some large $R$, then we have

$$
\left\|\varphi_{M}\right\| \leq C_{4}\left\|\Psi_{M}\right\|
$$

We then see from (6.37) and (6.38) that

$$
\left\|\Phi_{M}\right\| \leq C\left\|\Psi_{M}\right\|
$$

for some constant $C$. So we obtain

$$
\left\|\left(\mathbf{I}-\mathbb{D}(\lambda)^{-1} \mathbb{K}^{*}\right)^{-1}\right\| \leq C
$$

when $\left|\lambda_{M}\right| \geq R$. Combined with Lemma 6.4, the proof is completed.
As a consequence we have the following corollary.
Corollary 6.6 For each $L>0$ there is a constant $C=C(L)$ such that

$$
\begin{equation*}
\left\|\left(\mathbb{D}(\lambda(\mathbf{k}))-\mathbb{K}^{*}\right)^{-1}-\left(\mathbb{D}(\lambda(\mathbf{s}))-\mathbb{K}^{*}\right)^{-1}\right\| \leq C \sum_{j=1}^{M} \frac{\left|k_{j}-s_{j}\right|}{\left(1+\left|k_{j}\right|\right)\left(1+\left|s_{j}\right|\right)} \tag{6.39}
\end{equation*}
$$

for all $\mathbf{k}, \mathbf{s} \in \Theta_{L}$.
Proof. Using (6.34) and the identity

$$
\begin{aligned}
& \left(\mathbb{D}(\lambda(\mathbf{k}))-\mathbb{K}^{*}\right)^{-1}-\left(\mathbb{D}(\lambda(\mathbf{s}))-\mathbb{K}^{*}\right)^{-1} \\
& =\left(\mathbf{I}-\mathbb{D}(\lambda(\mathbf{k}))^{-1} \mathbb{K}^{*}\right)^{-1}\left(\mathbb{D}(\lambda(\mathbf{k}))^{-1}-\mathbb{D}(\lambda(\mathbf{s}))^{-1}\right)\left(\mathbf{I}-\mathbb{K}^{*} \mathbb{D}(\lambda(\mathbf{s}))^{-1}\right)^{-1}
\end{aligned}
$$

we have

$$
\begin{aligned}
\left\|\left(\mathbb{D}(\lambda(\mathbf{k}))-\mathbb{K}^{*}\right)^{-1}-\left(\mathbb{D}(\lambda(\mathbf{s}))-\mathbb{K}^{*}\right)^{-1}\right\| & \leq C\left\|\mathbb{D}(\lambda(\mathbf{k}))^{-1}-\mathbb{D}(\lambda(\mathbf{s}))^{-1}\right\| \\
& \leq C \sum_{j=1}^{M} \frac{\left|k_{j}-s_{j}\right|}{\left(1+\left|k_{j}\right|\right)\left(1+\left|s_{j}\right|\right)}
\end{aligned}
$$

for all $\mathbf{k}, \mathbf{s} \in \Theta_{L}$.
Now we obtain the uniformity of the regularity estimate and the Lipschitz continuity estimate for the solution $u_{\mathrm{k}}$ to (1.2) when there are multiple inclusions. The following theorem can be proved in a similar way to the proof of Theorem 4.5) and Theorem 4.7 using (6.34) and (6.39).

Theorem 6.7 For any $L>0$ there is a constant $C=C(L)$ such that

$$
\begin{equation*}
\left\|u_{\mathbf{k}}\right\|_{H^{1}(\Omega)} \leq C\|g\|_{-1 / 2(\partial \Omega)} \tag{6.40}
\end{equation*}
$$

and

$$
\begin{equation*}
\left\|u_{\mathbf{k}}-u_{\mathbf{s}}\right\|_{H^{1}(\Omega)} \leq C \sum_{j=1}^{M} \frac{\left|k_{j}-s_{j}\right|}{\left(1+\left|k_{j}\right|\right)\left(1+\left|s_{j}\right|\right)}\|g\|_{-1 / 2(\partial \Omega)} \tag{6.41}
\end{equation*}
$$

for all $\mathbf{k}, \mathbf{s} \in \Theta_{L}$ and $g \in H_{0}^{-1 / 2}(\partial \Omega)$.
The results of this section can also be applied to extend Theorem 3.2 and the uniformity of the boundary perturbation formula (Theorem 5.1) to the multiple inclusion case. But we do not state the extended results. We only mention that there are two different boundary perturbation formulas: one for the case of well-separated diametrically small inclusions and another for the case of a cluster of closely located small inclusions (see [10]). The uniform validity of the boundary perturbation formula for the first case can be proved using arguments in section 5, and that for the second case can be proved using results of this section.

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