# STABILITY RESULT FOR ELLIPTIC INVERSE PERIODIC COEFFICIENT PROBLEM BY PARTIAL DIRICHLET-TO-NEUMANN MAP 

MOURAD CHOULLI, YAVAR KIAN, ERIC SOCCORSI


#### Abstract

We study the inverse problem of identifying a periodic potential perturbation of the Dirichlet Laplacian acting in an infinite cylindrical domain, whose cross section is assumed to be bounded. We prove $\log -\log$ stable determination of the potential with respect to the partial Dirichlet-to-Neumann map, where the Neumann data is taken on slightly more than half of the boundary of the domain.


## 1. Introduction

Let $\Omega:=\mathbb{R} \times \omega$, where $\omega$ is a bounded domain of $\mathbb{R}^{2}$ which contains the origin, with $C^{2}$-boundary. Throughout the entire text we denote the generic point $x \in \Omega$ by $x=\left(x_{1}, x^{\prime}\right)$, where $x_{1} \in \mathbb{R}$ and $x^{\prime}:=$ $\left(x_{2}, x_{3}\right) \in \omega$. Given $V \in L^{\infty}(\Omega)$, real-valued and 1-periodic w.r.t. $x_{1}$, i.e.

$$
\begin{equation*}
V\left(x_{1}+1, x^{\prime}\right)=V\left(x_{1}, x^{\prime}\right), x^{\prime} \in \omega, x_{1} \in \mathbb{R} \tag{1.1}
\end{equation*}
$$

we consider the following boundary value problem (abbreviated as BVP):

$$
\left\{\begin{align*}
(-\Delta+V) v & =0, \quad \text { in } \Omega  \tag{1.2}\\
v & =f, \quad \text { on } \Gamma:=\partial \Omega
\end{align*}\right.
$$

Since $\Gamma=\mathbb{R} \times \partial \omega$, the outward unit vector $\nu$ normal to $\Gamma$ reads

$$
\nu\left(x_{1}, x^{\prime}\right)=\left(0, \nu^{\prime}\left(x^{\prime}\right)\right), x=\left(x_{1}, x^{\prime}\right) \in \Gamma
$$

where $\nu^{\prime}$ is the outer unit normal vector of $\partial \omega$. Therefore, for notational simplicity, we shall refer to $\nu$ for both exterior unit vectors normal to $\partial \omega$ and to $\Gamma$. Next for $\xi_{0} \in \mathbb{S}^{1}:=\left\{y \in \mathbb{R}^{2} ;|y|=1\right\}$ fixed, we introduce the $\xi_{0}$-illuminated (resp., $\xi_{0}$-shadowed) face of $\partial \omega$, as

$$
\begin{equation*}
\partial \omega_{\xi_{0}}^{-}:=\left\{x \in \partial \omega ; \xi_{0} \cdot \nu(x) \leqslant 0\right\}\left(\text { resp. }, \partial \omega_{\xi_{0}}^{+}=\left\{x \in \partial \omega ; \xi_{0} \cdot \nu(x) \geqslant 0\right\}\right) \tag{1.3}
\end{equation*}
$$

Here and in the remaining part of this text, we denote by $x \cdot y:=\sum_{j=1}^{k} x_{j} y_{j}$ the Euclidian scalar product of any two vectors $x:=\left(x_{1}, \ldots, x_{k}\right)$ and $y:=\left(y_{1}, \ldots, y_{k}\right)$ of $\mathbb{R}^{k}$, for $k \in \mathbb{N}^{*}$, and we put $|x|:=(x \cdot x)^{1 / 2}$.

Set $G:=\mathbb{R} \times G^{\prime}$, where $G^{\prime}$ is an arbitrary open neighbourhood containing the compact set $\partial \omega_{\xi_{0}}^{-}$in $\partial \omega$. In the present paper we seek stability in the determination of $V$ from the knowledge of the partial Dirichlet-to-Neumann (DN) map

$$
\begin{equation*}
\Lambda_{V}: f \mapsto \partial_{\nu} v_{\mid G} \tag{1.4}
\end{equation*}
$$

where $\partial_{\nu} v(x):=\nabla v(x) \cdot \nu(x)$ is the normal derivative of the solution $v$ to (1.2), computed at $x \in \Gamma$. Otherwise stated we aim for recovering the 1-periodic electric perturbation $V$ of the Dirichlet Laplacian in the waveguide $\Omega$, by probing the system with voltage $f$ at the boundary and measuring the current $\partial_{\nu} u$ on the sub-part $G$ of $\Gamma$. From a physics viewpoint, this amounts to estimating the impurity potential perturbing the guided propagation in periodic media such as crystals.
1.1. A short bibliography. Inverse coefficient problems in elliptic partial differential equations such as the celebrated Calderón problem have attracted many attention in recent years. In [SU], one of the first mathematical papers dealing with this problem, Sylvester and Uhlmann showed in dimension $n \geqslant 3$ that the full DN map (both the input and the output are taken on the whole boundary of the domain) uniquely determines a smooth conductivity coefficient. The case of $C^{1}$ conductivities or Lipschitz conductivities sufficiently close to the identity, is treated by [HT]. The identifiability of an unknown coefficient from partial knowledge of the DN map was first proved in [BU]. Assuming that the voltage (i.e. the Dirichlet data) is prescribed everywhere, Bukhgeim and Uhlmann claimed unique determination of the conductivity even when the current measurement (i.e. the Neumann data) is taken on slightly more than half of the boundary. Kenig, Sjöstrand and Uhlmann improved this result in [KSU] by taking both the Dirichlet and the Neumann observations on a neighbourhood of, respectively, the back and the front face illuminated by a point light-source lying outside of the convex hull of the domain. Identification for the corresponding two-dimensional inverse problem was treated by Bukhgeim in $[B]$ with the full data, and by Imanuvilov, Uhlmann and Yamamoto in [IUY1, IUY2] with partial data.

The stability issue for $n \geqslant 3$ was first addressed in [Al] by Alessandrini, who established a log-type stability estimate for the conductivity from the full DN map. Later on, in [HW], Heck and Wang proved a log-log-type stability estimate with respect to the partial DN map of [BU]. More recently, in [T], Tzou showed that both the magnetic field and the electric potential depend stably on the DN map, even when the Neumann boundary measurement is taken only on a subset that is slightly larger than half of the boundary. In [CDR1, CDR2], Caro, Dos Santos Ferreira and Ruiz derived a log-log stability estimate for the electric potential from the partial DN map of [KSU]. Notice that the derivation of the stability estimate of [CDR1] when the domain of observation is illuminated by a point at infinity, was revisited and simplified in [CKS1]. For the stability issue in the two-dimensional case we refer to [BIY, NSa, Sa].

All the above mentioned results were obtained in a bounded domain. It turns out that there is only a small number of mathematical papers dealing with inverse boundary measurements problems in an unbounded domain. Several of them are concerned with the slab geometry. This is precisely the case of [Ik, SW], where embedded objects are identified in an infinite slab. In [LU], Li and Uhlmann proved that the compactly supported electric potential of the stationary Schrödinger operator can be determined uniquely, when the Dirichlet and Neumann data are given either on the different boundary hyperplanes of the slab or on the same hyperplane. This result was generalized to the case of a magnetic Schrödinger operator in [KLU]. Let us mention that inverse boundary value problems in an infinite slab were addressed by Yang in [Y] for bi-harmonic operators. Recently, several stability results were derived in [CS, Ki, KPS1, KPS2, BKS] for non compactly supported coefficients inverse problems in an infinite waveguide with a bounded cross section. More specifically, we refer to [KKS, CKS] for the analysis of inverse problems in the framework of a periodic cylindrical domain examined in this paper.
1.2. Notations and admissible potentials. In this subsection we introduce some basic notations used throughout the section and define the set of admissible potentials under consideration in this paper.

Let $Y$ be either $\omega, \partial \omega$ or $G^{\prime}$. For $r$ and $s$ in $\mathbb{R}$, we denote by $\mathcal{H}^{r, s}(\mathbb{R} \times Y)$ the set $H^{r}\left(\mathbb{R} ; H^{s}(Y)\right)$. Evidently we write $\mathcal{H}^{r, s}(\Omega)$ (resp., $\mathcal{H}^{r, s}(\Gamma), \mathcal{H}^{r, s}(G)$ ) instead of $\mathcal{H}^{r, s}(\mathbb{R} \times \omega)$ (resp., $\mathcal{H}^{r, s}(\mathbb{R} \times \partial \omega), \mathcal{H}^{r, s}\left(\mathbb{R} \times G^{\prime}\right)$ ). Although this notation is reminiscent of the one used by Lions and Magenes in [LM1] for anisotropic Sobolev spaces $H^{r}\left(\mathbb{R} ; L^{2}(Y)\right) \cap L^{2}\left(\mathbb{R} ; H^{s}(Y)\right)$, it is worth noticing that they do not coincide with $\mathcal{H}^{r, s}(\mathbb{R} \times Y)$, unless we have $r=s=0$. Next, it is easy to see for each $r>0$ and $s>0$ that $\mathcal{H}^{-r,-s}(\mathbb{R} \times Y)$ is canonically identified with the space dual to $\mathcal{H}_{0}^{r, s}(\mathbb{R} \times Y)$, with respect to the pivot space $\mathcal{H}^{0,0}(\mathbb{R} \times Y)=L^{2}(\mathbb{R} \times Y)$. Here we have set $\mathcal{H}_{0}^{r, s}(\mathbb{R} \times Y):=H^{r}\left(\mathbb{R} ; H_{0}^{s}(Y)\right)$, where $H_{0}^{s}(Y)$ denotes the closure of $C_{0}^{\infty}(Y)$ in the topology of the Sobolev space $H^{s}(Y)$.

Further, $X_{1}$ and $X_{2}$ being two Hilbert spaces, we denote by $\mathcal{B}\left(X_{1}, X_{2}\right)$ the class of bounded operators $T: X_{1} \rightarrow X_{2}$.

Let us now introduce the set of admissible unknown potentials. To this end we denote by $C_{\omega}$ the Poincaré constant associated with $\omega$, i.e. the largest of those constants $c>0$ such that the Poincaré inequality

$$
\begin{equation*}
\left\|\nabla^{\prime} u\right\|_{L^{2}(\omega)} \geqslant c\|u\|_{L^{2}(\omega)}, u \in H_{0}^{1}(\omega) \tag{1.5}
\end{equation*}
$$

holds. Here $\nabla^{\prime}$ stands for the gradient with respect to $x^{\prime}=\left(x_{2}, x_{3}\right)$. Otherwise stated, we have

$$
\begin{equation*}
C_{\omega}:=\sup \{c>0 \text { satisfying (1.5) }\} . \tag{1.6}
\end{equation*}
$$

For $M_{-} \in\left(0, C_{\omega}\right)$ and $M_{+} \in\left[M_{-},+\infty\right)$, we define the set of admissible unknown potentials as

$$
\begin{equation*}
\mathscr{V}_{\omega}\left(M_{ \pm}\right):=\left\{V \in L^{\infty}(\Omega ; \mathbb{R}) \text { satisfying }(1.1),\|V\|_{L^{\infty}(\Omega)} \leqslant M_{+} \text {and }\|\max (0,-V)\|_{L^{\infty}(\Omega)} \leqslant M_{-}\right\} \tag{1.7}
\end{equation*}
$$

Notice that the constraint $\|\max (0,-V)\|_{L^{\infty}(\Omega)} \leqslant M_{-}$, imposed on admissible potentials $V$ in $\mathscr{V}_{\omega}\left(M_{ \pm}\right)$, guarantees that the perturbation by $V$ of the Dirichlet Laplacian in $\Omega$, is boundedly invertible in $L^{2}(\Omega)$, with norm not greater than $\left(C_{\omega}-M_{-}\right)^{-1}$. This condition could actually be weakened by only requiring that the distance of the spectrum of this operator to zero, be positive. Nevertheless, since the above mentioned condition on $V$ is more explicit than this latter, we stick with the definition (1.7) in the remaining part of this text.
1.3. Statement of the main result. Prior to stating the main result of this article we first examine in Proposition 1.1 below, the well-posedness of the BVP (1.2) in the space $H_{\Delta}(\Omega):=\left\{u \in L^{2}(\Omega) ; \Delta u \in L^{2}(\Omega)\right\}$ endowed with the norm

$$
\|u\|_{H_{\Delta}(\Omega)}^{2}:=\|u\|_{L^{2}(\Omega)}^{2}+\|\Delta u\|_{L^{2}(\Omega)}^{2}
$$

for suitable non-homogeneous Dirichlet boundary data $f$. Second, we rigorously define the DN map $\Lambda_{V}$ expressed in (1.4) and describe its main properties.

As a preamble, we introduce the two following trace maps by adapting the derivation of [LM1, Section 2, Theorem 6.5]. Namely, since $C_{0}^{\infty}(\bar{\Omega}):=\left\{u_{\mid \bar{\Omega}}, u \in C_{0}^{\infty}\left(\mathbb{R}^{3}\right)\right\}$ is dense in $H_{\Delta}(\Omega)$, by Lemma 2.1 below, we extend the mapping

$$
\mathcal{T}_{0} u:=u_{\mid \Gamma}\left(\text { resp., } \mathcal{T}_{1} u:=\partial_{\nu} u_{\mid \Gamma}\right), u \in C_{0}^{\infty}(\bar{\Omega})
$$

into a continuous function $\mathcal{T}_{0}: H_{\Delta}(\Omega) \rightarrow \mathcal{H}^{-2,-\frac{1}{2}}(\Gamma)$ (resp., $\mathcal{T}_{1}: H_{\Delta}(\Omega) \rightarrow \mathcal{H}^{-2,-\frac{3}{2}}(\Gamma)$ ). We refer to Lemma 2.2 and its proof, for more details.

Next we consider the space

$$
\mathscr{H}(\Gamma):=\mathcal{T}_{0} H_{\Delta}(\Omega)=\left\{\mathcal{T}_{0} u ; u \in H_{\Delta}(\Omega)\right\}
$$

and notice from Lemma 2.3 that $\mathcal{T}_{0}$ is bijective from $B:=\left\{u \in L^{2}(\Omega) ; \Delta u=0\right\}$ onto $\mathscr{H}(\Gamma)$. Therefore, with reference to $[\mathrm{BU}, \mathrm{NS}]$, we put

$$
\begin{equation*}
\|f\|_{\mathscr{H}(\Gamma)}:=\left\|\mathcal{T}_{0}^{-1} f\right\|_{H_{\Delta}(\Omega)}=\left\|\mathcal{T}_{0}^{-1} f\right\|_{L^{2}(\Omega)} \tag{1.8}
\end{equation*}
$$

where $\mathcal{T}_{0}^{-1}$ denotes the operator inverse to $\mathcal{T}_{0}: B \rightarrow \mathscr{H}(\Gamma)$.
We have the following existence and uniqueness result for the BVP (1.2).
Proposition 1.1. Pick $V \in \mathscr{V}_{\omega}\left(M_{ \pm}\right)$, where $M_{-} \in\left(0, C_{\omega}\right)$ and $M_{+} \in\left[M_{-},+\infty\right)$ are fixed.
(i) Then, for any $f \in \mathscr{H}(\Gamma)$, there exists a unique solution $v \in L^{2}(\Omega)$ to (1.2), such that the estimate

$$
\begin{equation*}
\|v\|_{L^{2}(\Omega)} \leqslant C\|f\|_{\mathscr{H}(\Gamma)} \tag{1.9}
\end{equation*}
$$

holds for some constant $C>0$ depending only on $\omega$ and $M_{ \pm}$.
(ii) The $D N \operatorname{map} \Lambda_{V}: f \mapsto \mathcal{T}_{1} v_{\mid G}$ is a bounded operator from $\mathscr{H}(\Gamma)$ into $\mathcal{H}^{-2,-\frac{3}{2}}(G)$.
(iii) Moreover, for each $W \in \mathscr{V}_{\omega}\left(M_{ \pm}\right)$, the operator $\Lambda_{V}-\Lambda_{W}$ is bounded from $\mathscr{H}(\Gamma)$ into $L^{2}(G)$.

Put $\Omega:=(0,1) \times \omega$. In view of Proposition 1.1, we now state the main result of this paper.

Theorem 1.2. Given $M_{-} \in\left(0, C_{\omega}\right)$ and $M_{+} \in\left[M_{-},+\infty\right)$, let $V_{j} \in \mathscr{V}_{\omega}\left(M_{ \pm}\right)$for $j=1,2$. Then, there exist two constants $C>0$ and $\gamma_{*} \in(0,1)$, both of them depending only on $\omega, M_{ \pm}$and $G^{\prime}$, such that the estimate

$$
\begin{equation*}
\left\|V_{1}-V_{2}\right\|_{H^{-1}(\check{\Omega})} \leqslant C \Phi\left(\left\|\Lambda_{V_{1}}-\Lambda_{V_{2}}\right\|\right) \tag{1.10}
\end{equation*}
$$

holds for

$$
\Phi(\gamma):=\left\{\begin{array}{cl}
\gamma & \text { if } \gamma \geqslant \gamma^{*}  \tag{1.11}\\
(\ln |\ln \gamma|)^{-1} & \text { if } \gamma \in\left(0, \gamma^{*}\right) \\
0 & \text { if } \gamma=0
\end{array}\right.
$$

Here $\left\|\Lambda_{V_{1}}-\Lambda_{V_{2}}\right\|$ denotes the norm of $\Lambda_{V_{1}}-\Lambda_{V_{2}}$ in $\mathcal{B}\left(\mathscr{H}(\Gamma), L^{2}(G)\right)$.
The statement of Theorem 1.2 remains valid for any periodic potential $V \in L^{\infty}(\Omega)$, provided 0 is in the resolvent set of $A_{V}$, the self-adjoint realization in $L^{2}(\Omega)$ of the Dirichlet Laplacian $-\Delta+V$. In this case, the multiplicative constants $C$ and $\gamma_{*}$, appearing in (1.10)-(1.11), depend on (the inverse of) the distance $d>0$, between 0 and the spectrum of $A_{V}$. In the particular case where $V \in \mathscr{V}_{\omega}\left(M_{ \pm}\right)$, with $M_{-} \in\left(0, C_{\omega}\right)$, we have $d \geqslant C_{\omega}-M_{-}$, and the implicit condition $d>0$ imposed on $V$, can be replaced by the explicit one on the negative part of the potential, i.e. $\|\max (0,-V)\|_{L^{\infty}(\Omega)} \leqslant M_{-}$.
1.4. Outline. The remaining part of the paper is organized as follows. Section 2 contains the proof of Proposition 1.1. In Section 3, we decompose (1.2) with the aid of the Floquet-Bloch-Gel'fand (FBG) transform, into a family of BVP with quasiperiodic boundary conditions,

$$
\left\{\begin{align*}
(-\Delta+V) v & =0, & & \text { in } \check{\Omega},  \tag{1.12}\\
v & =g_{\theta}, & & \text { on } \check{\Gamma}:=(0,1) \times \partial \omega, \\
v(1, \cdot)-e^{i \theta} v(0, \cdot) & =0, & & \text { in } \omega, \\
\partial_{x_{1}} v(1, \cdot)-e^{i \theta} \partial_{x_{1}} v(0, \cdot) & =0, & & \text { in } \omega,
\end{align*}\right.
$$

indexed by the real parameter $\theta \in[0,2 \pi)$. Here $g_{\theta}$ stands for the FBG transform of $f$, computed at $\theta$. We study the direct problem associated with (1.12) and reformulate the inverse problem under consideration as to whether the unknown function $V$ may be stably determined from the partial DN map associated with (1.12), for any arbitrary $\theta \in[0,2 \pi)$. We state in Theorem 3.3 that the answer is positive and establish that this claim entails Theorem 1.2. There are two key ingredients in the proof of Theorem 3.3. The first one is a sufficiently rich set of suitable complex geometric optics (CGO) solutions to (1.12), built in Section 4. The second one is a specifically designed Carleman estimate for quasiperiodic Laplace operators, derived in Section 5. Finally, the proof of Theorem 3.3 is displayed in Section 6.

Let us now briefly comment on the strategy of the proof of Theorems 1.2 and 3.3. Our approach is similar to the one of $[\mathrm{HW}]$ as it combines CGO solutions to the quasiperiodic Laplace equation in $\Omega$ with a suitable Carleman estimate. Nevertheless, in contrast to [CDR1, CDR2], the Carleman estimate of [KSU, Proposition 3.2 ] is not adapted to the framework of this paper. This is due to the quasiperiodic boundary conditions imposed on the CGO solutions employed in the context of inverse periodic coefficients problems. Therefore, in view of taking the Neumann measurements on $G$ only, we shall rather use the Carleman estimate with linear weights introduced in [BU].

## 2. Proof of Proposition 1.1

In this section we prove the claim of Proposition 1.1. As a preliminary we introduce in Subsection 2.1 the trace operators $\mathcal{T}_{j}, j=0,1$, and establish some useful properties that are needed for the proof of Proposition 1.1, which can be found in Subsection 2.2.
2.1. The trace operators. The rigorous definition of the trace operators $\mathcal{T}_{j}, j=0,1$, boils down to the coming lemma. Such a density result is rather classical for bounded domains (see e.g. [LM1, Section 2, Theorem 6.4]), but it has to be justified here since $\Omega$ is infinitely extended in the $x_{1}$ direction.
Lemma 2.1. The space $C_{0}^{\infty}(\bar{\Omega})$ is dense in $H_{\Delta}(\Omega)$.

Proof. Let $f \in H_{\Delta}(\Omega)^{\prime}$, the space of linear continuous forms on $H_{\Delta}(\Omega)$, satisfy

$$
\begin{equation*}
\langle f, w\rangle_{H_{\Delta}(\Omega)^{\prime}, H_{\Delta}(\Omega)}=0, w \in C_{0}^{\infty}(\bar{\Omega}) \tag{2.13}
\end{equation*}
$$

In order to establish the claim of Lemma 2.1, it is enough to show that $f$ is identically zero.
To do that we put $\jmath: u \in H_{\Delta}(\Omega) \mapsto(u, \Delta u) \in L^{2}(\Omega)^{2}$, notice that $H_{\Delta}(\Omega)$ is isometrically isomorphic to the closed subspace $Y:=\jmath\left(H_{\Delta}(\Omega)\right)$ of $L^{2}(\Omega)^{2}$, and introduce the following linear continuous form $g$ on $Y$ :

$$
\langle g, v\rangle_{Y^{\prime}, Y}:=\left\langle f, \jmath^{-1} v\right\rangle_{H_{\Delta}(\Omega)^{\prime}, H_{\Delta}(\Omega)}, v \in Y
$$

Here and henceforth, $Y^{\prime}$ denotes the space dual to $Y$ and $\langle\cdot, \cdot\rangle_{Y^{\prime}, Y}$ stands for the duality pairing between $Y^{\prime}$ and $Y$. Since $g$ can be extended by Hahn Banach theorem to a linear continuous form $\tilde{g}$ on $L^{2}(\Omega)^{2}$, we may find $\left(g_{1}, g_{2}\right) \in L^{2}(\Omega)^{2}$ such that

$$
\begin{equation*}
\langle f, u\rangle_{H_{\Delta}(\Omega)^{\prime}, H_{\Delta}(\Omega)}=\langle\tilde{g}, j u\rangle_{L^{2}(\Omega)^{2}}=\left\langle g_{1}, u\right\rangle_{L^{2}(\Omega)}+\left\langle g_{2}, \Delta u\right\rangle_{L^{2}(\Omega)}, u \in H_{\Delta}(\Omega), \tag{2.14}
\end{equation*}
$$

according to Riesz representation theorem. Upon extending $g_{1}$ and $g_{2}$ by zero outside $\Omega$, we deduce from (2.14) that

$$
\int_{\mathbb{R}^{3}}\left(g_{1} \bar{w}+g_{2} \overline{\Delta w}\right) d x=0, w \in C_{0}^{\infty}\left(\mathbb{R}^{3}\right)
$$

whence

$$
\begin{equation*}
-\Delta g_{2}=g_{1} \text { in } \mathbb{R}^{3} \tag{2.15}
\end{equation*}
$$

Since $g_{1} \in L^{2}(\Omega)$, then (2.15) yields that $g_{2} \in H^{2}(\Omega)$ by the classical elliptic regularity property. Further, as $g_{2}$ vanishes in $\mathbb{R}^{2} \backslash \bar{\Omega}$, we get that $g_{2} \in H_{0}^{2}(\Omega)$, the closure of $C_{0}^{\infty}(\Omega)$ in the topology of the second-order Sobolev space $H^{2}(\Omega)$. As a consequence we have

$$
\left\langle g_{2}, \Delta u\right\rangle_{L^{2}(\Omega)}=\left\langle\Delta g_{2}, u\right\rangle_{L^{2}(\Omega)}=-\left\langle g_{1}, u\right\rangle_{L^{2}(\Omega)}, u \in H_{\Delta}(\Omega)
$$

by (2.15). In view of (2.14) this entails that $f=0$, hence the result.
Armed with Lemma 2.1, we now define the trace maps $\mathcal{T}_{j}, j=0,1$, on the space $H_{\Delta}(\Omega)$, in the following manner.

Lemma 2.2. The mapping $w \in C_{0}^{\infty}(\bar{\Omega}) \mapsto w_{\mid \partial \Omega}$ (resp., $w \in C_{0}^{\infty}(\bar{\Omega}) \mapsto \partial_{\nu} w_{\mid \partial \Omega}$ ) can be extended over $H_{\Delta}(\Omega)$ to a bounded operator $\mathcal{T}_{0}: H_{\Delta}(\Omega) \rightarrow \mathcal{H}^{-2,-\frac{1}{2}}(\Gamma)$ (resp., $\mathcal{T}_{1}: H_{\Delta}(\Omega) \rightarrow \mathcal{H}^{-2,-\frac{3}{2}}(\Gamma)$ ).
Proof. It is well known that $u \in C^{\infty}(\bar{\omega}) \mapsto\left(u_{\mid \partial \omega}, \partial_{\nu} u_{\mid \partial \omega}\right)$ extends continuously to a bounded operator from $H^{2}(\omega)$ onto $H^{\frac{3}{2}}(\partial \omega) \times H^{\frac{1}{2}}(\partial \omega)$, so there exists $L: H^{\frac{3}{2}}(\partial \omega) \times H^{\frac{1}{2}}(\partial \omega) \rightarrow H^{2}(\omega)$, linear and bounded, such that

$$
L\left(h_{1}, h_{2}\right)_{\mid \partial \omega}=h_{1}, \partial_{\nu} L\left(h_{1}, h_{2}\right)_{\mid \partial \omega}=h_{2},\left(h_{1}, h_{2}\right) \in H^{\frac{3}{2}}(\partial \omega) \times H^{\frac{1}{2}}(\partial \omega) .
$$

Let us define the operator $\mathcal{L}: \mathcal{H}^{2, \frac{3}{2}}(\Gamma) \times \mathcal{H}^{2, \frac{1}{2}}(\Gamma) \rightarrow H^{2}(\Omega)$, by setting for a.e. $\left(x_{1}, x^{\prime}\right) \in \Omega$,

$$
\mathcal{L}\left(h_{1}, h_{2}\right)\left(x_{1}, x^{\prime}\right):=L\left(h_{1}\left(x_{1}, \cdot\right), h_{2}\left(x_{1}, \cdot\right)\right)\left(x^{\prime}\right),\left(h_{1}, h_{2}\right) \in \mathcal{H}^{2, \frac{3}{2}}(\Gamma) \times \mathcal{H}^{2, \frac{1}{2}}(\Gamma) .
$$

Using that $\|f\|_{H^{2}(\Omega)}^{2}=\sum_{k=0}^{2}\left\|\partial_{x_{1}}^{k} f\right\|_{L^{2}\left(\mathbb{R}, H^{2-k}(\omega)\right)}^{2}$, it is easy to see that $\mathcal{L}$ is bounded, and we check for every $\left(h_{1}, h_{2}\right) \in \mathcal{H}^{2, \frac{3}{2}}(\Gamma) \times H^{2, \frac{1}{2}}(\Gamma)$ that

$$
\mathcal{L}\left(h_{1}, h_{2}\right)_{\mid \partial \Omega}=h_{1}, \partial_{\nu} L\left(h_{1}, h_{2}\right)_{\mid \partial \Omega}=h_{2} .
$$

For $h \in \mathcal{H}^{2, \frac{1}{2}}(\Gamma)$, we put $w:=\mathcal{L}(0, h)$ in such a way that $w \in H^{2}(\Omega)$ satisfies

$$
\begin{equation*}
w_{\mid \partial \Omega}=0, \partial_{\nu} w_{\mid \partial \Omega}=h \text { and }\|w\|_{H^{2}(\Omega)} \leqslant C\|h\|_{\mathcal{H}^{2, \frac{1}{2}}(\Gamma)} \tag{2.16}
\end{equation*}
$$

the constant $C>0$ being independent of $h$. Next, for notational simplicity, we denote by $\mathcal{T}_{0} v$ the trace $v_{\mid \partial \Omega}$ of any function $v \in C_{0}^{\infty}(\bar{\Omega})$. Then, applying Green formula twice with respect to $x^{\prime} \in \omega$, and integrating by parts with respect to $x_{1}$ over $\mathbb{R}$, we find that

$$
\int_{\partial \Omega}\left(\mathcal{T}_{0} v\right) \bar{h} d \sigma(x)=\langle v, \Delta w\rangle_{L^{2}(\Omega)}-\langle\Delta v, w\rangle_{L^{2}(\Omega)}, v \in C_{0}^{\infty}(\bar{\Omega})
$$

Therefore, we get $\left|\int_{\partial \Omega}\left(\mathcal{T}_{0} v\right) \bar{h} d \sigma(x)\right| \leqslant 2\|v\|_{H_{\Delta}(\Omega)}\|w\|_{H^{2}(\Omega)}$ by Cauchy-Schwarz inequality, and hence

$$
\left|\left\langle\mathcal{T}_{0} v, h\right\rangle_{\mathcal{H}^{-2,-\frac{1}{2}}(\Gamma), \mathcal{H}^{2, \frac{1}{2}}(\Gamma)}\right| \leqslant 2 C\|v\|_{H_{\Delta}(\Omega)}\|h\|_{\mathcal{H}^{2, \frac{1}{2}}(\Gamma)},
$$

with the help of (2.16). Since $h$ is arbitrary in $\mathcal{H}^{2, \frac{1}{2}}(\Gamma)$, this entails that

$$
\left\|\mathcal{T}_{0} v\right\|_{\mathcal{H}^{-2,-\frac{1}{2}}(\Gamma)} \leqslant 2 C\|v\|_{H_{\Delta}(\Omega)}, v \in C_{0}^{\infty}(\bar{\Omega})
$$

which together with Lemma 2.1, proves that $\mathcal{T}_{0}$ can be extended over $H_{\Delta}(\Omega)$ to a bounded operator into $\mathcal{H}^{-2,-\frac{1}{2}}(\Gamma)$.

To prove the second part of the claim, we pick $g \in \mathcal{H}^{2, \frac{3}{2}}(\Gamma)$ and set $w:=\mathcal{L}(g, 0)$, so we have $w \in H^{2}(\Omega)$ and

$$
\begin{equation*}
w_{\mid \partial \Omega}=g, \partial_{\nu} w_{\mid \partial \Omega}=0 \text { and }\|w\|_{H^{2}(\Omega)} \leqslant C\|g\|_{\mathcal{H}^{2, \frac{3}{2}}(\Gamma)} \tag{2.17}
\end{equation*}
$$

for some constant $C>0$ that is independent of $w$. Next we denote by $\mathcal{T}_{1} v$ the trace $\partial_{\nu} v_{\mid \partial \Omega}$ of any $v \in C_{0}^{\infty}(\bar{\Omega})$ on $\partial \Omega$. Then, arguing as before, we find that the following estimate

$$
\left|\left\langle\mathcal{T}_{1} v, g\right\rangle_{\mathcal{H}^{-2,-\frac{3}{2}}(\Gamma), \mathcal{H}^{2, \frac{3}{2}}(\Gamma)}\right| \leqslant 2 C\|v\|_{H_{\Delta}(\Omega)}\|g\|_{\mathcal{H}^{2, \frac{3}{2}}(\Gamma)},
$$

holds uniformly in $g \in \mathcal{H}^{2, \frac{3}{2}}(\Gamma)$ and $v \in C_{0}^{\infty}(\bar{\Omega})$. This and Lemma 2.1 yield that $\mathcal{T}_{1}$ is extendable to a linear bounded operator from $H_{\Delta}(\Omega)$ into $\mathcal{H}^{-2,-\frac{3}{2}}(\Gamma)$.

We end this subsection by establishing the following result which was required by (1.8) to define the topology of the space $\mathscr{H}(\Gamma)$.
Lemma 2.3. $\mathcal{T}_{0}$ is bijective from $B=\left\{u \in L^{2}(\Omega) ; \Delta u=0\right\}$ onto $\mathcal{H}^{-2,-\frac{1}{2}}(\Gamma)$.
Proof. Given $u$ and $v$ in $B$ obeying $\mathcal{T}_{0} u=\mathcal{T}_{0} v$, we see that $w:=u-v$ satisfies the system

$$
\left\{\begin{align*}
-\Delta w=0, & \text { in } \Omega  \tag{2.18}\\
w=0, & \text { on } \Gamma
\end{align*}\right.
$$

Therefore we have $w=0$, by uniqueness (see e.g. [KPS2, Lemma 2.4]) of the solution to (2.18), proving that $\mathcal{T}_{0}$ is injective on $B$.

Further, we know from the definition of $\mathscr{H}(\Gamma)$ that for any $f \in \mathscr{H}(\Gamma)$, there exists $u \in H_{\Delta}(\Omega)$ such that $\mathcal{T}_{0} u=f$. Next, since $\Delta u \in L^{2}(\Omega)$, then the BVP

$$
\left\{\begin{array}{rlrl}
-\Delta v & =\Delta u, & & \text { in } \Omega \\
v & =0, & \text { on } \Gamma
\end{array}\right.
$$

admits a unique solution $v \in H^{2}(\Omega) \cap H_{0}^{1}(\Omega)$. Now, it is apparent that $w=u+v \in L^{2}(\Omega)$ satisfies $\Delta w=0$ and $\mathcal{T}_{0} w=\mathcal{T}_{0} u+\mathcal{T}_{0} v=f$.
2.2. Completion of the proof. Firstly, we introduce $A_{V}$, the self-adjoint operator in $L^{2}(\Omega)$, generated by the closed quadratic form

$$
u \mapsto a_{V}[u]:=\int_{\Omega}\left(|\nabla u(x)|^{2}+V(x)|u(x)|^{2}\right) d x, u \in D\left(a_{V}\right):=H_{0}^{1}(\Omega)
$$

We recall from [CKS, Lemma 2.2] that $A_{V}$ acts as $-\Delta+V$ on its domain $D\left(A_{V}\right)=H_{0}^{1}(\Omega) \cap H^{2}(\Omega)$. Moreover, since $V \in \mathscr{V}_{\omega}\left(M_{ \pm}\right)$we have

$$
a_{V}[u] \geqslant\left(C_{\omega}-M_{-}\right)\|u\|_{L^{2}(\Omega)}^{2}, u \in H_{0}^{1}(\Omega)
$$

by (1.6)-(1.5), hence $A_{V}$ is boundedly invertible in $L^{2}(\Omega)$ and

$$
\begin{equation*}
\left\|A_{V}^{-1}\right\|_{\mathcal{B}\left(L^{2}(\Omega)\right)} \leqslant \frac{1}{C_{\omega}-M_{-}} \tag{2.19}
\end{equation*}
$$

In order to establish the first statement of Proposition 1.1, i.e. (i), we notice that $v$ is solution to (1.2) if and only if $u:=v-\mathcal{T}_{0}^{-1} f$ solves the system

$$
\left\{\begin{align*}
(-\Delta+V) u & =-V \mathcal{T}_{0}^{-1} f & & \text { in } \Omega  \tag{2.20}\\
u & =0 & & \text { on } \Gamma .
\end{align*}\right.
$$

Since $V \mathcal{T}_{0}^{-1} f \in L^{2}(\Omega)$ and $A_{V}$ is boundedly invertible in $L^{2}(\Omega)$, then $u:=-A_{V}^{-1} V \mathcal{T}_{0}^{-1} f$ is the unique solution to (2.20). As a consequence we have

$$
\begin{equation*}
\|u\|_{L^{2}(\Omega)} \leqslant M_{+}\left\|A_{V}^{-1}\right\|_{\mathcal{B}\left(L^{2}(\Omega)\right)}\left\|\mathcal{T}_{0}^{-1} f\right\|_{L^{2}(\Omega)} \tag{2.21}
\end{equation*}
$$

Evidently $v:=u+\mathcal{T}_{0}^{-1} f$ is the unique solution to (1.2), and (1.9) follows readily from this, (1.8), (2.19) and (2.21).

We turn now to proving (ii)-(iii). For $f \in \mathscr{H}(\Gamma)$ fixed, we still denote by $v$ the solution to (1.2) associated with $f$. Since $v \in L^{2}(\Omega)$ and $\Delta v=V v \in L^{2}(\Omega)$ it holds true that $v \in H_{\Delta}(\Omega)$ and that

$$
\|v\|_{H_{\Delta}(\Omega)}^{2}=\|v\|_{L^{2}(\Omega)}^{2}+\|V v\|_{L^{2}(\Omega)}^{2} \leqslant\left(1+M_{+}^{2}\right)\|v\|_{L^{2}(\Omega)}^{2}
$$

From the continuity of $\mathcal{T}_{1}: H_{\Delta}(\Omega) \rightarrow \mathcal{H}^{-2,-\frac{3}{2}}(\Gamma)$, it then follows that

$$
\left\|\mathcal{T}_{1} v\right\|_{\mathcal{H}^{-2,-\frac{3}{2}}(\Gamma)}^{2} \leqslant\left(1+M_{+}^{2}\right)\left\|\mathcal{T}_{1}\right\|_{\mathcal{B}\left(H_{\Delta}(\Omega), \mathcal{H}^{-2,-\frac{3}{2}}(\Gamma)\right)}^{2}\|v\|_{L^{2}(\Omega)}^{2} .
$$

This and (1.9) yield that the DN map $\Lambda_{V}$ is bounded from $\mathscr{H}(\Gamma)$ into $\mathcal{H}^{-2,-\frac{3}{2}}(G)$.
Let us now establish (iii). We denote by $w$ the solution to (1.2) where $W$ is substituted for $V$, and notice that $u:=v-w$ is solution to the BVP

$$
\left\{\begin{aligned}
(-\Delta+V) u & =(W-V) w & & \text { in } \Omega \\
u & =0 & & \text { on } \Gamma .
\end{aligned}\right.
$$

Since $(W-V) w \in L^{2}(\Omega)$ and 0 is in the resolvent set of $A_{V}$, we have

$$
\begin{equation*}
u=A_{V}^{-1}(W-V) w, \tag{2.22}
\end{equation*}
$$

whence $u \in D\left(A_{V}\right)=H^{2}(\Omega) \cap H_{0}^{1}(\Omega)$. As the usual norm in $H^{2}(\Omega)$ is equivalent to the one associated with the domain of $A_{V}$, by [CKS, Lemma 2.2], we have

$$
\|u\|_{H^{2}(\Omega)} \leqslant C\left(\left\|A_{V} u\right\|_{L^{2}(\Omega)}+\|u\|_{L^{2}(\Omega)}\right)
$$

for some constant $C>0$, depending only on $\omega$ and $M_{+}$. This, (2.19) and (2.22) yield that

$$
\begin{equation*}
\|u\|_{H^{2}(\Omega)} \leqslant C\left(1+\left\|A_{V}^{-1}\right\|_{\mathcal{B}\left(L^{2}(\Omega)\right)}\right)\|(W-V) w\|_{L^{2}(\Omega)} \leqslant 2 M_{+} C\left(1+\frac{1}{C_{\omega}-M_{-}}\right)\|w\|_{L^{2}(\Omega)} \tag{2.23}
\end{equation*}
$$

Bearing in mind that $\|w\|_{L^{2}(\Omega)} \leqslant c\|f\|_{\mathscr{H}(\Gamma)}$, by (1.9), and using the continuity of the trace operator $u \mapsto$ $\partial_{\nu} u_{\mid \Gamma}$ from $H^{2}(\Omega)$ into $L^{2}(\Gamma)$, we deduce from (2.23) that

$$
\left\|\partial_{\nu} u\right\|_{L^{2}(\Gamma)} \leqslant C\|u\|_{H^{2}(\Omega)} \leqslant C\|f\|_{\mathscr{H}(\Gamma)}
$$

where $C$ is another positive constant depending only on $\omega$ and $M_{ \pm}$. Now, the desired result follows from this and the identity $\partial_{\nu} u_{\mid G}=\mathcal{T}_{1} u_{\mid G}=\mathcal{T}_{1} v_{\mid G}-\mathcal{T}_{1} w_{\mid G}=\left(\Lambda_{V}-\Lambda_{W}\right) f$.

## 3. Fiber Decomposition

In this section we decompose (1.2) into the family of $\operatorname{BVP}(1.12)$ indexed by $\theta \in[0,2 \pi)$. This is by means of the FBG transform, introduced in Subsection 3.1, which decomposes the operator $A_{V}$. In Subsection 3.2 we examine the direct problem associated with (1.12) for each $\theta \in[0,2 \pi)$ and study the corresponding DN $\operatorname{map} \Lambda_{V, \theta}$. Finally, in Subsection 3.3 we reformulate the inverse problem under consideration as to whether $V$ can be stably retrieved from partial knowledge of $\Lambda_{V, \theta}$. We state in Theorem 3.3 that this is actually the case provided the Neumann data are measured on $\check{G}:=(0,1) \times G^{\prime}$.
3.1. FBG transform and fiber decomposition of $A_{V}$. Let $Y$ be either $\omega$ or $\partial \omega$. The main tool for the analysis of 1-periodic waveguides $\mathbb{R} \times Y$ is the partial FBG transform defined for every $f \in C_{0}^{\infty}(\mathbb{R} \times Y)$ as

$$
\begin{equation*}
\left(\mathcal{U}_{Y} f\right)_{\theta}\left(x_{1}, y\right):=\sum_{k \in \mathbb{Z}} e^{-i k \theta} f\left(x_{1}+k, y\right),\left(x_{1}, y\right) \in \mathbb{R} \times Y, \theta \in[0,2 \pi) \tag{3.24}
\end{equation*}
$$

For notational simplicity, we systematically drop the $Y$ in $\mathcal{U}_{Y}$ and write $\mathcal{U}$ instead of $\mathcal{U}_{Y}$. In view of [RS2, Section XIII.16], the above operator can be extended to a unitary operator, still denoted by $\mathcal{U}$, from $L^{2}(\mathbb{R} \times Y)$ onto the Hilbert space $\int_{(0,2 \pi)}^{\oplus} L^{2}((0,1) \times Y) \frac{d \theta}{2 \pi}$. Otherwise stated, $\mathcal{U}$ decomposes any $f \in L^{2}(\mathbb{R} \times Y)$ into a family of $L^{2}((0,1) \times Y)$-functions $(\mathcal{U} f)_{\theta}$, for $\theta \in[0,2 \pi)$, such that $\int_{0}^{2 \pi}\left\|(\mathcal{U} f)_{\theta}\right\|_{L^{2}((0,1) \times Y)}^{2} \frac{d \theta}{2 \pi}=\|f\|_{L^{2}(\mathbb{R} \times Y)}^{2}$. Although the direct integral sum $\int_{(0,2 \pi)}^{\oplus} L^{2}((0,1) \times Y) \frac{d \theta}{2 \pi}$ has constant (i.e., $\theta$-independent) fibers $L^{2}((0,1) \times$ $Y$ ), each function $(\mathcal{U} f)_{\theta}$ satisfies $\theta$-quasiperiodic boundary conditions on $\{0,1\} \times \partial Y$. For the sake of clarity, we point out that since $\check{\Omega}=(0,1) \times \omega$ (resp., $\check{\Gamma}=(0,1) \times \partial \omega$ ), then $\mathcal{U}$ maps $L^{2}(\Omega)$ onto $\int_{(0,2 \pi)}^{\oplus} L^{2}(\check{\Omega}) \frac{d \theta}{2 \pi}$ if $Y=\omega$, (resp., $L^{2}(\partial \Omega)$ onto $\int_{(0,2 \pi)}^{\oplus} L^{2}(\check{\Gamma}) \frac{d \theta}{2 \pi}$ if $Y=\partial \omega$ ).

Let $V \in L^{\infty}(\Omega ; \mathbb{R})$ satisfy (1.1). Since $V$ is 1 -periodic with respect to $x_{1}$, the operator $A_{V}$ defined in Section 1, is decomposed by $\mathcal{U}$. To make this claim more precise, we first introduce the following functional spaces. For $\theta \in[0,2 \pi)$ fixed, we put with reference to [CKS, Section 6.1],

$$
\mathcal{H}_{\theta}^{s}((0,1) \times Y):=\left\{u \in H^{s}((0,1) \times Y) ; \partial_{x_{1}}^{j} u(1, \cdot)-e^{i \theta} \partial_{x_{1}}^{j} u(0, \cdot)=0, j<s-\frac{1}{2}\right\} \text { if } s>\frac{1}{2}
$$

and

$$
\mathcal{H}_{\theta}^{s}((0,1) \times Y):=H^{s}((0,1) \times Y) \text { if } s \in\left[0, \frac{1}{2}\right]
$$

Evidently, we shall write $\mathcal{H}_{\theta}^{s}(\check{\Omega})\left(\right.$ resp., $\left.\mathcal{H}_{\theta}^{s}(\check{\Gamma})\right)$ instead of $\mathcal{H}_{\theta}^{s}((0,1) \times Y)$ for $Y=\omega($ resp., $Y=\partial \omega)$.
Next we introduce the operator

$$
\mathcal{A}_{V, \theta}:=-\Delta+V, D\left(\mathcal{A}_{V, \theta}\right):=\left\{(\mathcal{U} \psi)_{\theta}, \psi \in D\left(A_{V}\right)\right\}=\mathcal{H}_{\theta}^{2}(\check{\Omega}) \cap L^{2}\left(0,1 ; H_{0}^{1}(\omega)\right)
$$

self-adjoint in $L^{2}(\check{\Omega})$. Then, taking into account that $(\mathcal{U} \psi)_{\theta}$ is quasiperiodic in $x_{1}$,

$$
(\mathcal{U} \psi)_{\theta}\left(x_{1}+1, x^{\prime}\right)=e^{i \theta}(\mathcal{U} \psi)_{\theta}\left(x_{1}, x^{\prime}\right), x_{1} \in \mathbb{R}, x^{\prime} \in \omega, \theta \in[0,2 \pi)
$$

that $\mathcal{U}$ commutes with the partial derivatives with respect to the variables $x_{j}, j=1,2,3$,

$$
\begin{equation*}
\left(\mathcal{U} \frac{\partial^{m} \psi}{\partial x_{j}^{m}}\right)_{\theta}=\frac{\partial^{m}(\mathcal{U} \psi)_{\theta}}{\partial x_{j}^{m}}, j=1,2,3, m \in \mathbb{N}^{*}, \theta \in[0,2 \pi), \tag{3.25}
\end{equation*}
$$

and that the potential $V$ obeys (1.1), we find for any $\psi \in D\left(A_{V}\right)$ that $\left(\mathcal{U} A_{V} \psi\right)_{\theta}=\mathcal{A}_{V, \theta}(\mathcal{U} \psi)_{\theta}$ for a.e. $\theta \in(0,2 \pi)$, and that $\int_{0}^{2 \pi}\left\|\mathcal{A}_{V, \theta}(\mathcal{U} \psi)_{\theta}\right\|_{L^{2}(\check{\Omega})}^{2} \frac{d \theta}{2 \pi}=\left\|A_{V} \psi\right\|_{L^{2}(\Omega)}^{2}$, i.e.,

$$
\begin{equation*}
\mathcal{U} A_{V} \mathcal{U}^{-1}=\int_{(0,2 \pi)}^{\oplus} \mathcal{A}_{V, \theta} \frac{d \theta}{2 \pi} \tag{3.26}
\end{equation*}
$$

Having seen this we turn now to analysing the direct problem associated with (1.12).
3.2. Analysis of the fibered problems. For $\theta \in[0, \pi)$ and $n \in \mathbb{N} \cup\{\infty\}$, we set
$\mathcal{C}_{\theta}^{n}([0,1] \times \bar{\omega}):=\left\{u \in \mathcal{C}^{n}([0,1] \times \bar{\omega}) ; \partial_{x_{1}}^{j} u(1, \cdot)-e^{i \theta} \partial_{x_{1}}^{j} u(0, \cdot)=0\right.$ for all $j \in \mathbb{N}$ such that $\left.j \leqslant n\right\}$.
A slight modification of the proof of [LM1, Section 2, Theorem 6.4] shows that $C^{\infty}([0,1] \times \bar{\omega})$ is dense in $H_{\Delta}(\check{\Omega}):=\left\{u \in L^{2}(\check{\Omega}) ; \Delta u \in L^{2}(\check{\Omega})\right\}$ endowed with the norm $\|u\|_{H_{\Delta}(\check{\Omega})}:=\left(\|u\|_{L^{2}(\check{\Omega})}^{2}+\|\Delta u\|_{L^{2}(\check{\Omega})}^{2}\right)^{\frac{1}{2}}$.

Lemma 3.1. For $j=0$ and $j=1$, each of the two following mappings

$$
u \mapsto \partial_{x_{1}}^{j} u(0, \cdot) \text { and } u \mapsto \partial_{x_{1}}^{j} u(1, \cdot), u \in C^{\infty}([0,1] \times \bar{\omega})
$$

can be extended to a bounded operator from $H_{\Delta}(\check{\Omega})$ into $H^{-2}(\omega)$.
Proof. Let us prove the claim for $u \mapsto \partial_{x_{1}} u(0, \cdot)$, the three other cases being treated in the same way. To this purpose we consider a function $\chi \in C_{0}^{\infty}(\mathbb{R})$ satisfying $\chi\left(x_{1}\right)=1$ for $x_{1} \in[-1 / 3,1 / 3], \chi\left(x_{1}\right) \in[0,1]$ for $x_{1} \in(1 / 3,2 / 3)$ and $\chi\left(x_{1}\right)=0$ for $x_{1} \in[2 / 3,4 / 3]$. We pick $g \in H_{0}^{2}(\omega)$, the closure of $C_{0}^{\infty}(\omega)$ in $H^{2}(\omega)$, put $v\left(x_{1}, x^{\prime}\right):=\chi\left(x_{1}\right) g\left(x^{\prime}\right)$, and then notice that

$$
\begin{equation*}
v=\partial_{\nu} v=0 \text { on } \check{\Gamma}, \partial_{x_{1}} v(0, \cdot)=v(1, \cdot)=\partial_{x_{1}} v(1, \cdot)=0 \text { and } v(0, \cdot)=g \text { in } \omega . \tag{3.27}
\end{equation*}
$$

Moreover, there exists a positive constant $C=C(\omega, \chi)$ such that we have

$$
\begin{equation*}
\|v\|_{H^{2}(\check{\Omega})} \leqslant C\|g\|_{H^{2}(\omega)} . \tag{3.28}
\end{equation*}
$$

For all $u \in C^{\infty}([0,1] \times \bar{\omega})$, we deduce from (3.27) and the Green formula that

$$
\left\langle\partial_{x_{1}} u(0, \cdot), g\right\rangle_{L^{2}(\omega)}=\langle u, \Delta v\rangle_{\check{\Omega}}-\langle\Delta u, v\rangle_{\check{\Omega}}
$$

This and (3.28) yield

$$
\left|\left\langle\partial_{x_{1}} u(0, \cdot), g\right\rangle_{H^{-2}(\omega), H_{0}^{2}(\omega)}\right| \leqslant 2\|u\|_{H_{\Delta}(\check{\Omega})}\|v\|_{H^{2}(\check{\Omega})} \leqslant C\|u\|_{H_{\Delta}(\check{\Omega})}\|g\|_{H^{2}(\omega)} .
$$

Therefore, since $C^{\infty}([0,1] \times \bar{\omega})$ is dense in $H_{\Delta}(\check{\Omega})$, we may extend the mapping $u \mapsto \partial_{x_{1}} u(0, \cdot)$ to a continuous map from $H_{\Delta}(\Omega)$ into $H^{-2}(\omega)$.

With reference to Lemma 3.1, we introduce the following closed subset of $H_{\Delta}(\bar{\Omega})$ :

$$
H_{\Delta, \theta}(\check{\Omega}):=\left\{u \in H_{\Delta}(\check{\Omega}) ; \partial_{x_{1}}^{j} u(1, \cdot)-e^{i \theta} \partial_{x_{1}}^{j} u(0, \cdot)=0 \text { on } \omega \text { for } j=0,1\right\}
$$

Notice for further use that

$$
\begin{equation*}
\mathcal{U} H_{\Delta}(\Omega)=\int_{(0,2 \pi)}^{\oplus} H_{\Delta, \theta}(\check{\Omega}) \frac{d \theta}{2 \pi} \tag{3.29}
\end{equation*}
$$

and that $H_{\Delta, \theta}(\check{\Omega}) \cap C^{\infty}([0,1] \times \bar{\omega}) \subset \mathcal{H}_{\theta}^{2}(\check{\Omega})$. Moreover, since $\mathcal{H}_{\theta}^{2}(\check{\Omega})$ is continuously embedded in $H_{\Delta, \theta}(\check{\Omega})$ and that $\mathcal{C}_{\theta}^{\infty}([0,1] \times \bar{\omega})$ is dense in $\mathcal{H}_{\theta}^{2}(\Omega)$, then the space $\mathcal{C}_{\theta}^{\infty}([0,1] \times \bar{\omega})$ is dense in $H_{\Delta, \theta}(\Omega)$ as well. Therefore, by reasoning in the same way as in the proof of Lemma 2.2, we extend the mapping $w \in$ $\mathcal{C}_{\theta}^{\infty}([0,1] \times \bar{\omega}) \mapsto w_{\mid \check{\Gamma}}\left(\right.$ resp., $\left.w \in \mathcal{C}_{\theta}^{\infty}([0,1] \times \bar{\omega}) \mapsto \partial_{\nu} w_{\mid \check{\Gamma}}\right)$ to a bounded operator

$$
\mathcal{T}_{0, \theta}: H_{\Delta, \theta}(\check{\Omega}) \rightarrow \mathcal{H}_{\theta}^{-2}\left(0,1 ; H^{-\frac{1}{2}}(\partial \omega)\right)\left(\text { resp., } \mathcal{T}_{1, \theta}: H_{\Delta, \theta}(\check{\Omega}) \rightarrow \mathcal{H}_{\theta}^{-2}\left(0,1 ; H^{-\frac{3}{2}}(\partial \omega)\right)\right)
$$

Here we denote by $\mathcal{H}_{\theta}^{s}(0,1 ; X)$, where $s \in \mathbb{R}$ and $X$ is a Banach space for the norm $\|\cdot\|_{X}$, the set of functions

$$
t \in(0,1) \mapsto \varphi(t):=\sum_{k \in \mathbb{Z}} \varphi_{k} e^{i(\theta+2 \pi k) t} \text { with }\left(\varphi_{k}\right)_{k} \in X^{\mathbb{Z}} \text { satisfying } \sum_{k \in \mathbb{Z}}\left(1+k^{2}\right)^{s}\left\|\varphi_{k}\right\|_{X}^{2}<\infty
$$

Endowed with the norm $\|\varphi\|_{\mathcal{H}_{\theta}^{s}(0,1 ; X)}:=\left(\sum_{k \in \mathbb{Z}}\left(1+k^{2}\right)^{s}\left\|\varphi_{k}\right\|_{X}^{2}\right)^{\frac{1}{2}}, \mathcal{H}_{\theta}^{s}(0,1 ; X)$ is a Banach space. Notice that if $X^{\prime}$ is the dual space of $X$, then $\mathcal{H}_{\theta}^{-s}\left(0,1 ; X^{\prime}\right)$ is the space dual to $\mathcal{H}_{\theta}^{s}(0,1 ; X)$.

Let us next introduce the set

$$
\mathscr{H}_{\theta}(\check{\Gamma}):=\left\{\mathcal{T}_{0, \theta} u ; u \in H_{\Delta, \theta}(\check{\Omega})\right\}
$$

Arguing as in the derivation of Lemma 2.3, we check that $\mathcal{T}_{0, \theta}$ is a bijection from $B_{\theta}:=\left\{u \in H_{\Delta, \theta}(\check{\Omega}) ; \Delta u=\right.$ $0\}$ onto $\mathscr{H}_{\theta}(\check{\Gamma})$. Put

$$
\begin{equation*}
\|f\|_{\mathscr{H}_{\theta}(\check{\Gamma})}:=\left\|\mathcal{T}_{0, \theta}^{-1} f\right\|_{H_{\Delta}(\check{\Omega})}=\left\|\mathcal{T}_{0, \theta}^{-1} f\right\|_{L^{2}(\check{\Omega})} \tag{3.30}
\end{equation*}
$$

where $\mathcal{T}_{0, \theta}^{-1}$ denotes the operator inverse to $\mathcal{T}_{0, \theta}: B_{\theta} \rightarrow \mathscr{H}_{\theta}(\check{\Gamma})$.
We may now establish the following technical result, which is inspired by Proposition 1.1.

Proposition 3.2. Pick $\theta \in[0,2 \pi)$ and assume that the conditions of Proposition 1.1 are satisfied. Then the three following statements are true.
(i) For any $f \in \mathscr{H}_{\theta}(\check{\Gamma})$, there exists a unique solution $v \in H_{\Delta, \theta}(\check{\Omega})$ to (1.12), satisfying

$$
\begin{equation*}
\|v\|_{L^{2}(\check{\Omega})} \leqslant C\|f\|_{\mathscr{H}_{\theta}(\check{\Gamma})} \tag{3.31}
\end{equation*}
$$

for some positive constant $C=C\left(\omega, M_{ \pm}\right)$, independent of $\theta$.
(ii) The $D N \operatorname{map} \Lambda_{V, \theta}: f \mapsto \mathcal{T}_{1, \theta} u_{\mid \check{G}}$, where we recall that $\check{G}=(0,1) \times G^{\prime}$, is a bounded operator from $\mathscr{H}_{\theta}(\check{\Gamma})$ into $\mathcal{H}_{\theta}^{-2}\left(0,1 ; H^{-\frac{3}{2}}\left(G^{\prime}\right)\right)$.
(iii) For every $W \in \mathscr{V}_{\omega}\left(M_{ \pm}\right)$, the operator $\Lambda_{V, \theta}-\Lambda_{W, \theta}$ is bounded from $\mathscr{H}_{\theta}(\check{\Gamma})$ into $L^{2}(\check{G})$.

Proof. To prove (i) we use the fact that $v$ is solution to (1.12) if and only if the function $u:=v-\mathcal{T}_{0, \theta}^{-1} f$ satisfies the BVP

$$
\left\{\begin{align*}
(-\Delta+V) u & =-V \mathcal{T}_{0, \theta}^{-1} f, & & \text { in } \check{\Omega},  \tag{3.32}\\
u & =0, & & \text { on } \check{\Gamma} \\
u(1, \cdot)-e^{i \theta} u(0, \cdot) & =0, & & \text { in } \omega . \\
\partial_{x_{1}} u(1, \cdot)-e^{i \theta} \partial_{x_{1}} u(0, \cdot) & =0, & & \text { in } \omega .
\end{align*}\right.
$$

Due to (3.26) we know from [RS2, Theorem XIII.98] that $\sigma\left(\mathcal{A}_{V, \theta}\right) \subset \sigma\left(A_{V}\right)$. Hence $0 \notin \sigma\left(\mathcal{A}_{V, \theta}\right)$ and (3.32) admits a unique solution $u=-\mathcal{A}_{V, \theta}^{-1} V \mathcal{T}_{0, \theta}^{-1} f$ as $V \mathcal{T}_{0, \theta}^{-1} f \in L^{2}(\Omega \check{\Omega})$. As a consequence $v=\left(I-\mathcal{A}_{V, \theta}^{-1} V\right) \mathcal{T}_{0, \theta}^{-1} f$ is the unique solution to the BVP (1.12) and (3.31) follows readily from this, (3.30) and the estimate

$$
\left\|\mathcal{A}_{V, \theta}^{-1}\right\|_{\mathcal{B}\left(L^{2}(\check{\Omega})\right)} \leqslant \frac{1}{C_{\omega}-M_{-}}
$$

arising from (1.5)-(1.6) and the fact that the operator $\mathcal{A}_{V, \theta}$ is generated by the following quadratic form:

$$
\begin{equation*}
u \mapsto a_{V, \theta}[u]:=\int_{\check{\Omega}}\left(|\nabla u(x)|^{2}+V(x)|u(x)|^{2}\right) d x, u \in D\left(a_{V, \theta}\right):=\mathcal{H}_{\theta}^{1}(\check{\Omega}) \cap L^{2}\left(0,1 ; H_{0}^{1}(\omega)\right) . \tag{3.33}
\end{equation*}
$$

The rest of the proof follows the same lines as the derivation of (ii)-(iii) in Proposition 1.1.
For further reference, we now establish for every $\theta \in[0,2 \pi)$ that the estimate

$$
\begin{equation*}
\left\|\mathcal{T}_{0, \theta} u\right\|_{\mathscr{H}_{\theta}(\check{\Gamma})} \leqslant C\|u\|_{H_{\Delta}(\check{\Omega})}, u \in H_{\Delta, \theta}(\check{\Omega}), \tag{3.34}
\end{equation*}
$$

holds for some constant $C>0$ depending only on $\omega$. Indeed, for each $u \in H_{\Delta, \theta}(\Omega)$, where $\theta \in[0,2 \pi)$ is fixed, put $f:=\mathcal{T}_{0, \theta} u \in \mathscr{H}_{\theta}(\check{\Gamma})$ and let $v:=\mathcal{T}_{0, \theta}^{-1} f \in B_{\theta}$. Evidently $w:=v-u$ satisfies

$$
\left\{\begin{aligned}
-\Delta w & =\Delta u, & & \text { in } \check{\Omega}, \\
w & =0, & & \text { on } \check{\Gamma} \\
w(1, \cdot)-e^{i \theta} w(0, \cdot) & =0, & & \text { in } \omega \\
\partial_{x_{1}} w(1, \cdot)-e^{i \theta} \partial_{x_{1}} w(0, \cdot) & =0, & & \text { in } \omega
\end{aligned}\right.
$$

hence $w=\mathcal{A}_{0, \theta}^{-1} \Delta u$. Since $\mathcal{A}_{0, \theta}$ is bounded from below by the Poincaré constant $C_{\omega}>0$ defined in (1.6), in virtue of (3.33), it holds true that $\|w\|_{L^{2}(\check{\Omega})} \leqslant C_{\omega}^{-1}\|\Delta u\|_{L^{2}(\check{\Omega})}$. Therefore we have

$$
\|f\|_{\mathscr{H}_{\theta}(\check{\Gamma})}=\|v\|_{L^{2}(\check{\Omega})} \leqslant\|u\|_{L^{2}(\check{\Omega})}+\|w\|_{L^{2}(\check{\Omega})} \leqslant\left(1+C_{\omega}^{-1}\right)\left(\|u\|_{L^{2}(\check{\Omega})}+\|\Delta u\|_{L^{2}(\check{\Omega})}\right)
$$

which yields (3.34).
3.3. Linking up (1.2) with (1.12). Let $Y$ be either $\omega$ or $\partial \omega$, and let $s \in \mathbb{R}_{+}$and $k \in \mathbb{N}$. In light of (3.25) we extend $\mathcal{U}$ to a unitary operator from $\mathcal{H}^{k, s}(\mathbb{R} \times Y)$ onto the Hilbert space $\int_{(0,2 \pi)}^{\oplus} \mathcal{H}_{\theta}^{k}\left(0,1 ; H^{s}(Y)\right) \frac{d \theta}{2 \pi}$. Next, for $f \in \mathcal{H}^{-k,-s}(\mathbb{R} \times Y)$, we define $\mathcal{U} f$ by setting

$$
\langle\mathcal{U} f, g\rangle:=\left\langle f, \mathcal{U}^{-1} g\right\rangle_{\mathcal{H}^{-k,-s}(\mathbb{R} \times Y), \mathcal{H}_{0}^{k, s}(\mathbb{R} \times Y)} \text { for all } g \in \int_{(0,2 \pi)}^{\oplus} \mathcal{H}_{\theta}^{k}\left(0,1 ; H_{0}^{s}(Y)\right) \frac{d \theta}{2 \pi}
$$

Here $\langle.,$.$\rangle denotes the duality pairing between \int_{(0,2 \pi)}^{\oplus} \mathcal{H}_{\theta}^{-k}\left(0,1 ; H^{-s}(Y)\right) \frac{d \theta}{2 \pi}$ and $\int_{(0,2 \pi)}^{\oplus} \mathcal{H}_{\theta}^{k}\left(0,1 ; H_{0}^{s}(Y)\right) \frac{d \theta}{2 \pi}$, and we recall that $\mathcal{H}_{0}^{k, s}(\mathbb{R} \times Y)=H^{k}\left(\mathbb{R} ; H_{0}^{s}(Y)\right)$, where $H_{0}^{s}(Y)$ is the closure of $C_{0}^{\infty}(Y)$ in the topology of $H^{s}(Y)$. It is clear that the above defined operator $\mathcal{U}$ is unitary from $\mathcal{H}^{-k,-s}(\mathbb{R} \times Y)$ onto the Hilbert space $\int_{(0,2 \pi)}^{\oplus} \mathcal{H}_{\theta}^{-k}\left(0,1 ; H^{-s}(Y)\right) \frac{d \theta}{2 \pi}$. Moreover, for every $f \in C_{0}^{\infty}\left(\mathbb{R} ; C^{\infty}(\bar{Y})\right)$ and $g \in \int_{(0,2 \pi)}^{\oplus} \mathcal{H}_{\theta}^{k}\left(0,1 ; H_{0}^{s}(Y)\right) \frac{d \theta}{2 \pi}$, we have

$$
\langle\mathcal{U} f, g\rangle=\langle\mathcal{U} f, g\rangle_{L^{2}\left((0,2 \pi) \frac{d \theta}{2 \pi} ; \mathcal{H}_{\theta}^{k}\left(0,1 ; H^{s}(Y)\right)\right.}=\left\langle f, \mathcal{U}^{-1} g\right\rangle_{L^{2}(\mathbb{R} \times Y)},
$$

and hence $\langle\mathcal{U} f, g\rangle=\left\langle f, \mathcal{U}^{-1} g\right\rangle_{\mathcal{H}^{-k,-s}(\mathbb{R} \times Y), \mathcal{H}_{0}^{k, s}(\mathbb{R} \times Y)}$, which shows that the above definition of $\mathcal{U}$ is more general than the one of Subsection 3.1.

For $j=0,1$, it is apparent that $\left(\mathcal{U} \mathcal{T}_{j} u\right)_{\theta}=\mathcal{T}_{j, \theta}(\mathcal{U} u)_{\theta}$ for all $u \in C_{0}^{\infty}(\bar{\Omega})$ and $\theta \in[0,2 \pi)$. Since the operators $\mathcal{T}_{j}: H_{\Delta}(\Omega) \rightarrow \mathcal{H}^{-2,-\left(j+\frac{1}{2}\right)}(\Gamma)$ and $\mathcal{T}_{j, \theta}: H_{\Delta, \theta}(\check{\Omega}) \rightarrow \mathcal{H}_{\theta}^{-2}\left(0,1 ; H^{-\left(j+\frac{1}{2}\right)}(\partial \omega)\right)$ are bounded, we deduce from (3.29) and the density of $C_{0}^{\infty}(\bar{\Omega})$ in $H_{\Delta}(\Omega)$ that $\mathcal{U} \mathcal{T}_{j} \mathcal{U}^{-1}=\int_{(0,2 \pi)}^{\oplus} \mathcal{T}_{j, \theta} \frac{d \theta}{2 \pi}$. This entails that $\mathcal{U} \mathscr{H}(\Gamma)=\int_{(0,2 \pi)}^{\oplus} \mathscr{H}_{\theta}(\check{\Gamma}) \frac{d \theta}{2 \pi}$. Moreover, arguing as in the proof of [CKS, Proposition 3.1], we find for any $f \in \mathscr{H}(\Gamma)$ that $u$ is the $H_{\Delta}(\Omega)$-solution to (1.2) if and only if each function $(\mathcal{U} u)_{\theta} \in H_{\Delta, \theta}(\Omega)$, for a.e. $\theta \in[0,2 \pi)$, satisfies (1.12) where $(\mathcal{U} f)_{\theta}$ is substituted for $g$. This is provided $V \in L^{\infty}(\Omega ; \mathbb{R})$ verifies (1.1) in such a way that (3.26) holds. Consequently we have:

$$
\begin{equation*}
\mathcal{U} \Lambda_{V} \mathcal{U}^{-1}=\int_{(0,2 \pi)}^{\oplus} \Lambda_{V, \theta} \frac{d \theta}{2 \pi} \tag{3.35}
\end{equation*}
$$

Let $V_{j}, j=1,2$, be the two potentials introduced in Theorem 1.2. Then, in light of (3.35) and (iii) in Propositions 1.1 and 3.2, we have

$$
\begin{equation*}
\left\|\Lambda_{V_{1}}-\Lambda_{V_{2}}\right\|_{\mathcal{B}\left(\mathscr{H}(\Gamma), L^{2}(G)\right)}=\sup _{\theta \in[0,2 \pi)}\left\|\Lambda_{V_{1}, \theta}-\Lambda_{V_{2}, \theta}\right\|_{\mathcal{B}\left(\mathscr{H}_{\theta}(\check{\Gamma}), L^{2}(\check{G})\right)} \tag{3.36}
\end{equation*}
$$

from [Di, Section II.2, Proposition 2]. As the function $\Phi$, given by (1.11), is non decreasing, then Theorem 1.2 follows readily from (3.36) and the following statement whose proof is given in Section 6.

Theorem 3.3. Under the conditions of Theorem 1.2, the following estimate

$$
\begin{equation*}
\left\|V_{1}-V_{2}\right\|_{H^{-1}(\check{\Omega})} \leqslant C_{\theta} \Phi\left(\left\|\Lambda_{V_{1}, \theta}-\Lambda_{V_{2}, \theta}\right\|\right) \tag{3.37}
\end{equation*}
$$

holds for every $\theta \in[0,2 \pi)$. Here $C_{\theta}$ is a positive constant depending only on $T, \omega, G^{\prime}$ and possibly on $\theta, \Phi$ is defined by (1.11), and $\left\|\Lambda_{V_{1}, \theta}-\Lambda_{V_{2}, \theta}\right\|$ denotes the usual norm of $\Lambda_{V_{1}, \theta}-\Lambda_{V_{2}, \theta}$ in $\mathcal{B}\left(\mathscr{H}_{\theta}(\check{\Gamma}), L^{2}(\check{G})\right)$.

Remark 3.4. Notice from (3.36)-(3.37) and the non decreasing behavior of $\Phi$ on $[0,+\infty)$, that the multiplicative constant $C$ appearing in the right hand side of the stability estimate (1.10), reads

$$
C:=\inf _{\theta \in[0,2 \pi)} C_{\theta}
$$

and hence is independent of $\theta$.

## 4. Complex geometric optics solutions

In this section we aim for building CGO solutions to the system

$$
\left\{\begin{align*}
&(-\Delta+V) u=0,  \tag{4.38}\\
& u(1, \cdot)-e^{i \theta} u(0, \cdot)=0, \\
& \text { in } \omega \\
& \partial_{x_{1}} u(1, \cdot)-e^{i \theta} \partial_{x_{1}} u(0, \cdot)=0, \\
& \text { in } \omega
\end{align*}\right.
$$

associated with $V \in L^{\infty}(\check{\Omega} ; \mathbb{R})$ and $\theta \in[0,2 \pi)$. Namely, given a sufficiently large $\tau>0$, we seek solutions of the form

$$
\begin{equation*}
u(x)=(1+w(x)) e^{\zeta \cdot x}, x \in \check{\Omega} \tag{4.39}
\end{equation*}
$$

to (4.38), where $\zeta \in i(\theta+2 \pi \mathbb{Z}) \times \mathbb{C}^{2}$ is chosen in such a way that $\Delta e^{\zeta \cdot x}=0$ for every $x \in \check{\Omega}$, and $w \in H^{2}(\check{\Omega})$ satisfies the periodic boundary conditions

$$
\frac{\partial^{k} w}{\partial_{x_{1}}^{k}}\left(1, x^{\prime}\right)=\frac{\partial^{k} w}{\partial_{x_{1}}^{k}}\left(0, x^{\prime}\right), x^{\prime} \in \omega, k=0,1
$$

together with the estimate

$$
\begin{equation*}
\|w\|_{H^{s}(\check{\Omega})} \leqslant C \tau^{s-1}, s \in[0,2] \tag{4.40}
\end{equation*}
$$

for some positive constant $C$, independent of $\tau$.
To do that we proceed as follows. We pick $k \in \mathbb{Z}$ and for $\xi \in \mathbb{S}^{1}$ we choose $\eta \in \mathbb{R}^{2} \backslash\{0\}$ such that $\eta \cdot \xi=0$. For $r>0$, we set

$$
\ell=\ell(k, \eta, r, \theta):=\left\{\begin{array}{cl}
(\theta+2 \pi([r]+1))\left(1,-2 \pi k \frac{\eta}{|\eta|^{2}}\right), & \text { if } k \text { is even }  \tag{4.41}\\
\left(\theta+2 \pi\left([r]+\frac{3}{2}\right)\right)\left(1,-2 \pi k \frac{\eta}{|\eta|^{2}}\right), & \text { if } k \text { is odd }
\end{array}\right.
$$

in such a way that $\ell \cdot(2 \pi k, \eta)=\ell^{\prime} \cdot \xi=0$, where we used the notation $\ell=\left(\ell_{1}, \ell^{\prime}\right) \in \mathbb{R} \times \mathbb{R}^{2}$. Here [r] stands for the integer part of $r$, that is the unique integer fulfilling $[r] \leqslant r<[r]+1$. Next, we introduce

$$
\begin{equation*}
\tau=\tau(k, \eta, r, \theta):=\sqrt{\frac{|\eta|^{2}}{4}+\pi^{2} k^{2}+|\ell|^{2}} \tag{4.42}
\end{equation*}
$$

and notice that

$$
\begin{equation*}
2 \pi r<\tau \leqslant \frac{|(2 \pi k, \eta)|}{2}+4 \pi(r+1)\left(1+\frac{|2 \pi k|}{|\eta|}\right) . \tag{4.43}
\end{equation*}
$$

Then, putting

$$
\begin{equation*}
\zeta_{1}:=\left(i \pi k,-\tau \xi+i \frac{\eta}{2}\right)+i \ell \text { and } \zeta_{2}:=\left(-i \pi k, \tau \xi-i \frac{\eta}{2}\right)+i \ell \tag{4.44}
\end{equation*}
$$

it is easy to check for $j=1,2$, that we have

$$
\begin{equation*}
\zeta_{1}+\overline{\zeta_{2}}=i(2 \pi k, \eta) \text { and } \zeta_{j} \cdot \zeta_{j}=\Im \zeta_{j} \cdot \mathfrak{R} \zeta_{j}=0, \text { with } \zeta_{j} \in i(\theta+2 \pi \mathbb{Z}) \times \mathbb{C}^{2} \text { for } j=1,2 \tag{4.45}
\end{equation*}
$$

Further, we argue as in the derivation of [CKS, Proposition 4.1 and Lemma 4.2], and obtain the following:
Lemma 4.1. Fix $V \in L^{\infty}(\check{\Omega})$. Let $\theta \in[0,2 \pi)$, $k \in \mathbb{Z}$, and let $\xi \in \mathbb{S}^{1}$. Pick $\eta \in \mathbb{R}^{2} \backslash\{0\}$ such that $\xi \cdot \eta=0$. Then we may find $\tau_{0}>1$ and $w \in \mathcal{H}_{0}^{2}(\Omega)$ fulfilling (4.40), such that for every $\tau \geqslant \tau_{0}$, the function $u$ given by (4.39) with $\zeta=\zeta_{1}$, where $\zeta_{1}$ is defined in (4.44), is solution to the system (4.38).

We recall from Subsection 3.3 that the space $\mathcal{H}_{0}^{2}(\Omega)$ appearing in Lemma 4.1 denotes the set $\{u \in$ $H^{2}(\check{\Omega}), u(1, \cdot)-u(0, \cdot)=\partial_{x_{1}} u(1, \cdot)-\partial_{x_{1}} u(0, \cdot)=0$ in $\left.\omega\right\}$.

## 5. A CARLEMAN EStimate for the quasiperiodic Laplace operator in $\check{\Omega}$

In this section we derive a Carleman estimate for the Laplace operator in $\check{\Omega}$ with quasiperiodic boundary conditions. We proceed by adapting the Carleman inequality of [BU, Lemma 2.1] to quasiperiodic functions (with respect to $x_{1}$ ) on $\check{\Omega}$.

Proposition 5.1. Let $\xi \in \mathbb{S}^{1}$ and pick $a, b$ in $\mathbb{R}$, with $a<b$, in such $a$ way that we have

$$
\omega \subset\left\{x^{\prime} \in \mathbb{R}^{2} ; \xi \cdot x^{\prime} \in(a, b)\right\}
$$

Put $d:=b-a$. Then for all $\theta \in[0,2 \pi)$ and all $\tau>0$, the estimate

$$
\begin{align*}
& \frac{8 \tau^{2}}{d}\left\|e^{-\tau \xi \cdot x^{\prime}} u\right\|_{L^{2}(\check{\Omega})}^{2}+2 \tau\left\|e^{-\tau \xi \cdot x^{\prime}}(\xi \cdot \nu)^{1 / 2} \partial_{\nu} u\right\|_{L^{2}\left(\check{\Gamma}_{\xi}^{+}\right)}^{2} \\
\leqslant & \left\|e^{-\tau \xi \cdot x^{\prime}} \Delta u\right\|_{L^{2}(\check{\Omega})}^{2}+2 \tau\left\|e^{-\tau \xi \cdot x^{\prime}}|\xi \cdot \nu|^{1 / 2} \partial_{\nu} u\right\|_{L^{2}\left(\check{\Gamma}_{\xi}^{-}\right)}^{2}, \tag{5.46}
\end{align*}
$$

holds for every $u \in \mathcal{C}_{\theta}^{2}([0,1] \times \bar{\omega})$ satisfying $u_{\mid \check{\Gamma}}=0$. Here we used the notations $\check{\Gamma}_{\xi}^{ \pm}:=(0,1) \times \partial \omega_{\xi}^{ \pm}$.
Proof. The operator $e^{-\tau \xi \cdot x^{\prime}} \Delta e^{\tau \xi \cdot x^{\prime}}$ decomposes into the sum $P_{+}^{1}+P_{+}^{\prime}+P_{-}$, with

$$
P_{+}^{1}:=\partial_{x_{1}}^{2}, P_{+}^{\prime}:=\Delta^{\prime}+\tau^{2} \text { and } P_{-}:=2 \tau \xi \cdot \nabla^{\prime}
$$

where the symbol $\Delta^{\prime}$ (resp., $\nabla^{\prime}$ ) stands for the Laplace (resp., gradient) operator with respect to $x^{\prime} \in \omega$. Thus we get upon setting $v(x):=e^{-\tau \xi \cdot x^{\prime}} u(x)$ that

$$
\begin{aligned}
\left\|e^{-\tau \xi \cdot x^{\prime}} \Delta u\right\|_{L^{2}(\check{\Omega})}^{2} & =\left\|e^{-\tau \xi \cdot x^{\prime}} \Delta e^{\tau \xi \cdot x^{\prime}} v\right\|_{L^{2}(\check{\Omega})}^{2} \\
& =\left\|\left(P_{+}^{1}+P_{+}^{\prime}+P_{-}\right) v\right\|_{L^{2}(\check{\Omega})}^{2} \\
& =\left\|\left(P_{+}^{1}+P_{+}^{\prime}\right) v\right\|_{L^{2}(\check{\Omega})}^{2}+\left\|P_{-} v\right\|_{L^{2}(\check{\Omega})}^{2}+2 \Re\left\langle P_{+}^{1} v, P_{-} v\right\rangle_{L^{2}(\check{\Omega})}+2 \Re\left\langle P_{+}^{\prime} v, P_{-} v\right\rangle_{L^{2}(\check{\Omega})}
\end{aligned}
$$

and hence

$$
\begin{equation*}
\left\|P_{-} v\right\|_{L^{2}(\check{\Omega})}^{2}+2 \mathfrak{R}\left\langle P_{+}^{\prime} v, P_{-} v\right\rangle_{L^{2}(\check{\Omega})} \leqslant\left\|e^{-\tau \xi \cdot x^{\prime}} \Delta u\right\|_{L^{2}(\check{\Omega})}^{2}-2 \Re\left\langle P_{+}^{1} v, P_{-} v\right\rangle_{L^{2}(\check{\Omega})} . \tag{5.47}
\end{equation*}
$$

Moreover, we find upon integrating by parts that

$$
\begin{align*}
\mathfrak{\Re}\left\langle P_{+}^{1} v, P_{-} v\right\rangle_{L^{2}(\check{\Omega})}= & -\tau \int_{0}^{1} \int_{\omega} \nabla^{\prime} \cdot\left(\left|\partial_{x_{1}} v(x)\right|^{2} \xi\right) d x^{\prime} d x_{1} \\
& +2 \tau \int_{\omega}\left(\partial_{x_{1}} v\left(1, x^{\prime}\right) \xi \cdot \overline{\nabla^{\prime} v\left(1, x^{\prime}\right)}-\partial_{x_{1}} v\left(0, x^{\prime}\right) \xi \cdot \overline{\nabla^{\prime} v\left(0, x^{\prime}\right)}\right) d x^{\prime} \tag{5.48}
\end{align*}
$$

Since $v$ is quasiperiodic in $x_{1}$, then it holds true that

$$
\partial_{x_{1}} v\left(1, x^{\prime}\right) \xi \cdot \overline{\nabla^{\prime} v\left(1, x^{\prime}\right)}=e^{i \theta} \partial_{x_{1}} v\left(0, x^{\prime}\right) \xi \cdot \overline{e^{i \theta} \nabla^{\prime} v\left(0, x^{\prime}\right)}=\partial_{x_{1}} v\left(0, x^{\prime}\right) \xi \cdot \overline{\nabla^{\prime} v\left(0, x^{\prime}\right)}, x^{\prime} \in \omega,
$$

and the last term in the right-hand-side of (5.48) cancels, so we have

$$
\begin{equation*}
\mathfrak{R}\left\langle P_{+}^{1} v, P_{-} v\right\rangle_{L^{2}(\check{\Omega})}=-\tau \int_{0}^{1} \int_{\omega} \nabla^{\prime} \cdot\left(\left|\partial_{x_{1}} v(x)\right|^{2} \xi\right) d x^{\prime} d x_{1}=-\tau \int_{\check{\Gamma}}\left|\partial_{x_{1}} v(x)\right|^{2} \xi \cdot \nu(x) d \sigma(x)=0 \tag{5.49}
\end{equation*}
$$

Here we used the fact, arising from the homogeneous boundary data $v_{\mid \check{\Gamma}}=0$, that $\partial_{x_{1}} v$ vanishes on $\check{\Gamma}$.
Next, as the function $w:=v\left(x_{1}, \cdot\right) \in C^{2}(\bar{\omega})$ satisfies $w_{\mid \partial \omega}=0$ for a.e. $x_{1} \in(0,1)$, we deduce from the following Carleman estimate

$$
\frac{8 \tau^{2}}{d^{2}}\|w\|_{L^{2}(\omega)}^{2}+2 \tau \int_{\partial \omega} e^{-2 \tau \xi \cdot x^{\prime}} \xi \cdot \nu\left(x^{\prime}\right)\left|\partial_{\nu} e^{\tau \xi \cdot x^{\prime}} w\left(x^{\prime}\right)\right|^{2} d \sigma\left(x^{\prime}\right) \leqslant\left\|P_{-} w\right\|_{L^{2}(\omega)}^{2}+2 \Re\left\langle P_{+}^{\prime} w, P_{-} w\right\rangle_{L^{2}(\omega)}
$$

which is borrowed from [BU, Lemma 2.1], that

$$
\begin{aligned}
& \frac{8 \tau^{2}}{d^{2}}\left\|e^{-\tau \xi \cdot x^{\prime}} u\left(x_{1}, \cdot\right)\right\|_{L^{2}(\omega)}^{2}+2 \tau \int_{\partial \omega} e^{-2 \tau \xi \cdot x^{\prime}} \xi \cdot \nu(x)\left|\partial_{\nu} u\left(x_{1}, x^{\prime}\right)\right|^{2} d \sigma\left(x^{\prime}\right) \\
\leqslant & \left\|P_{-} v\left(x_{1}, \cdot\right)\right\|_{L^{2}(\omega)}^{2}+2 \mathfrak{R}\left\langle P_{+}^{\prime} v\left(x_{1}, \cdot\right), P_{-} v\left(x_{1}, \cdot\right)\right\rangle_{L^{2}(\omega)}
\end{aligned}
$$

Thus, by integrating both sides of the above inequality with respect to $x_{1} \in(0,1)$, we get that

$$
\begin{equation*}
\frac{8 \tau^{2}}{d^{2}}\left\|e^{-\tau \xi \cdot x^{\prime}} u\right\|_{L^{2}(\check{\Omega})}^{2}+2 \tau \int_{\check{\Gamma}} e^{-2 \tau \xi \cdot x^{\prime}} \xi \cdot \nu(x)\left|\partial_{\nu} u(x)\right|^{2} d \sigma(x) \leqslant\left\|P_{-} v\right\|_{L^{2}(\check{\Omega})}^{2}+2 \Re\left\langle P_{+}^{\prime} v, P_{-} v\right\rangle_{L^{2}(\check{\Omega})} \tag{5.50}
\end{equation*}
$$

Finally, putting (5.47)-(5.50) together and recalling (1.3), we end up getting (5.46).
Let us now perturb the Laplacian in (5.46) by the multiplier by $V \in L^{\infty}(\check{\Omega})$. Since

$$
|\Delta u|^{2} \leqslant 2\left(|(-\Delta+V) u|^{2}+\|V\|_{L^{\infty}(\Omega)}^{2}|u|^{2}\right)
$$

we find through elementary computations that

$$
\begin{aligned}
& \left(\frac{4 \tau^{2}}{d}-\|V\|_{L^{\infty}(\check{\Omega})}^{2}\right)\left\|e^{-\tau \xi \cdot x^{\prime}} u\right\|_{L^{2}(\check{\Omega})}^{2}+\tau\left\|e^{-\tau \xi \cdot x^{\prime}}(\xi \cdot \nu)^{\frac{1}{2}} \partial_{\nu} u\right\|_{L^{2}\left(\check{\Gamma}_{\xi}^{+}\right)}^{2} \\
\leqslant & \left\|e^{-\tau \xi \cdot x^{\prime}}(-\Delta+V) u\right\|_{L^{2}(\check{\Omega})}^{2}+\tau\left\|e^{-\tau \xi \cdot x^{\prime}}|\xi \cdot \nu|^{\frac{1}{2}} \partial_{\nu} u\right\|_{L^{2}\left(\check{\Gamma}_{\xi}^{-}\right)}^{2}
\end{aligned}
$$

As a consequence we have obtained the:
Corollary 5.2. For $M>0$, let $V \in L^{\infty}(\Omega)$ satisfy (1.1) and $\|V\|_{L^{\infty}(\Omega)} \leqslant M$. Then, under the conditions of Proposition 5.1, we have

$$
\begin{aligned}
& \frac{2 \tau^{2}}{d}\left\|e^{-\tau \xi \cdot x^{\prime}} u\right\|_{L^{2}(\check{\Omega})}^{2}+\tau\left\|e^{-\tau \xi \cdot x^{\prime}}(\xi \cdot \nu)^{\frac{1}{2}} \partial_{\nu} u\right\|_{L^{2}\left(\check{\Gamma}_{\xi}^{+}\right)}^{2} \\
\leqslant & \left\|e^{-\tau \xi \cdot x^{\prime}}(-\Delta+V) u\right\|_{L^{2}(\check{\Omega})}^{2}+\tau\left\|e^{-\tau \xi \cdot x^{\prime}}|\xi \cdot \nu|^{\frac{1}{2}} \partial_{\nu} u\right\|_{L^{2}\left(\check{\Gamma}_{\xi}^{-}\right)}^{2},
\end{aligned}
$$

provided $\tau \geqslant \tau_{1}:=M(d / 2)^{\frac{1}{2}}$.
Notice from the density of $\left\{u \in \mathcal{C}_{\theta}^{2}([0,1] \times \bar{\omega}), u_{\mid \check{\Gamma}}=0\right\}$ in $\left\{u \in \mathcal{H}_{\theta}^{2}(\check{\Omega}), u_{\mid \check{\Gamma}}=0\right\}$ that the Carleman estimate of Corollary 5.2 remains valid for all $u \in \mathcal{H}_{\theta}^{2}(\check{\Omega})$ such that $u_{\mid \check{\Gamma}}=0$.

## 6. Proof of Theorem 3.3

In this section we prove the stability estimate (3.37). To this purpose we set for a.e. $x \in(0,1) \times \mathbb{R}^{2}$,

$$
V(x):=\left\{\begin{array}{cl}
\left(V_{2}-V_{1}\right)(x) & \text { if } x \in \check{\Omega} \\
0 & \text { if } x \in(0,1) \times\left(\mathbb{R}^{2} \backslash \omega\right),
\end{array}\right.
$$

and start by establishing several technical results that are useful for the proof of (3.37).
Since we aim for proving the stability estimate (1.10) with the aid of (3.37), we recall from Remark 3.4 that we may completely leave aside the question of how the constant $C_{\theta}$ involved in (3.37) depends on $\theta$. Therefore we shall not specify the possible dependence with respect to $\theta$ of the various constants appearing in this section.
6.1. Preliminary estimate. Bearing in mind that $G^{\prime}$ is an open neighbourhood containing $\partial \omega_{\xi_{0}}^{-}=\{x \in$ $\left.\partial \omega ; \xi_{0} \cdot \nu(x) \leqslant 0\right\}$, we pick $\varepsilon>0$ so small that

$$
\begin{equation*}
\forall \xi \in \mathbb{S}^{1},\left(\left|\xi-\xi_{0}\right| \leqslant \varepsilon\right) \Rightarrow\left(\{x \in \partial \omega ; \xi \cdot \nu(x) \leqslant \varepsilon\} \subset G^{\prime}\right) \tag{6.51}
\end{equation*}
$$

and we establish the following technical result with the help of the CGO solutions introduced in Section 4 and the Carleman estimate derived in Section 5.

Lemma 6.1. Let $\varepsilon$ be as in (6.51). Then we may find $r_{1}>0$ such that the estimate

$$
\begin{equation*}
\left|\int_{\mathbb{R}^{2}} \int_{0}^{1} V\left(x_{1}, x^{\prime}\right) e^{-i\left(2 \pi k x_{1}+\eta \cdot x^{\prime}\right)} d x_{1} d x^{\prime}\right|^{2} \leqslant C\left(\frac{1}{\tau}+e^{C^{\prime} \tau}\left\|\Lambda_{V_{2}, \theta}-\Lambda_{V_{1}, \theta}\right\|^{2}\right) \tag{6.52}
\end{equation*}
$$

holds for all $r \geqslant r_{1}, \xi \in\left\{z \in \mathbb{S}^{1} ;\left|z-\xi_{0}\right| \leqslant \varepsilon\right\}, \eta \in \mathbb{R}^{2} \backslash\{0\}$ satisfying $\xi \cdot \eta=0$, and $k \in \mathbb{Z}$. Here $\tau \geqslant 1$ is defined by (4.41)-(4.42), and the positive constants $C$ and $C^{\prime}$ depend only on $\omega, M_{+}$, $\varepsilon$ and $\xi_{0}$.

Proof. We first introduce the sets

$$
\begin{equation*}
\partial \omega_{\xi, \varepsilon}^{+}:=\{x \in \partial \omega ; \xi \cdot \nu(x)>\varepsilon\} \text { and } \partial \omega_{\xi, \varepsilon}^{-}:=\{x \in \partial \omega ; \xi \cdot \nu(x) \leqslant \varepsilon\} \tag{6.53}
\end{equation*}
$$

and we establish the orthogonality identity (6.59) below, with the aid of the CGO solutions of Lemma 4.1. To this end we choose $r$ sufficiently large, namely $r \geqslant r_{1}:=\max \left(1+\tau_{0}, \tau_{1}\right)$, where $\tau_{0}$ (resp., $\tau_{1}$ ) is the constant introduced in Lemma 4.1 (resp., Corollary 5.2), in such a way that we have

$$
\begin{equation*}
\tau=\sqrt{\frac{|\eta|^{2}}{4}+\pi^{2} k^{2}+|\ell(r, k, \eta, \theta)|^{2}} \geqslant \max \left(\tau_{0}, \tau_{1}\right) \tag{6.54}
\end{equation*}
$$

by (4.41)-(4.42). Next, for $j=1,2$, we define $\zeta_{j}$ as in (4.44) and we denote by

$$
\begin{equation*}
u_{j}(x)=\left(1+w_{j}(x)\right) e^{\zeta_{j} \cdot x}, x \in \check{\Omega} \tag{6.55}
\end{equation*}
$$

the $\mathcal{H}_{\theta}^{2}(\Omega)$-solution to (4.38) associated with $V=V_{j}$, which is given by Lemma 4.1. For further reference we recall that $w_{j}$ satisfies the condition (4.40) with $s=1$, entailing that

$$
\begin{equation*}
\left\|w_{j}\right\|_{H^{1}(\check{\Omega})} \leqslant C \tag{6.56}
\end{equation*}
$$

for some constant $C=C\left(\omega, M_{+}\right)>0$. Next, if $v_{1}$ satisfies the system

$$
\left\{\begin{align*}
\left(-\Delta+V_{1}\right) v_{1} & =0, & & \text { in } \check{\Omega},  \tag{6.57}\\
v_{1} & =\mathcal{T}_{0, \theta} u_{2}, & & \text { on } \check{\Gamma}, \\
v_{1}(1, \cdot)-e^{i \theta} v_{1}(0, \cdot) & =0, & & \text { on } \omega, \\
\partial_{x_{1}} v_{1}(1, \cdot)-e^{i \theta} \partial_{x_{1}} v_{1}(0, \cdot) & =0, & & \text { on } \omega,
\end{align*}\right.
$$

then it is easy to check that the function $u:=v_{1}-u_{2}$ is solution to the following BVP:

$$
\left\{\begin{align*}
\left(-\Delta+V_{1}\right) u & =V u_{2}, & & \text { in } \check{\Omega},  \tag{6.58}\\
u & =0, & & \text { on } \check{\Gamma} \\
u(1, \cdot)-e^{i \theta} u(0, \cdot) & =0, & & \text { on } \omega \\
\partial_{x_{1}} u(1, \cdot)-e^{i \theta} \partial_{x_{1}} u(0, \cdot) & =0, & & \text { on } \omega .
\end{align*}\right.
$$

Moreover, as $V u_{2} \in L^{2}(\check{\Omega})$ and 0 belongs to the resolvent set of $\mathcal{A}_{V_{1}, \theta}$, then $u=\mathcal{A}_{V_{1}, \theta}^{-1} V u_{2} \in \mathcal{H}_{\theta}^{2}(\check{\Omega})$. Further, bearing in mind that $\left(-\Delta+V_{1}\right) u_{1}=0$ in $\check{\Omega}$, from the first line of (4.38) with $V=V_{1}$, we deduce from (6.58) and the Green formula that

$$
\int_{\check{\Omega}} V u_{2} \overline{u_{1}} d x=\int_{\check{\Omega}}\left(-\Delta+V_{1}\right) u \overline{u_{1}} d x=\int_{\check{\Gamma}}\left(\partial_{\nu} u\right) \overline{u_{1}} d \sigma(x) .
$$

In view of (6.53), this can be equivalently rewritten as

$$
\begin{equation*}
\int_{\check{\Omega}} V u_{2} \overline{u_{1}} d x=\int_{\check{\Gamma}_{\xi, \varepsilon}^{+}}\left(\partial_{\nu} u\right) \overline{u_{1}} d \sigma(x)+\int_{\check{\Gamma}_{\xi, \varepsilon}^{-}}\left(\partial_{\nu} u\right) \overline{u_{1}} d \sigma(x), \tag{6.59}
\end{equation*}
$$

where $\check{\Gamma}_{\xi, \varepsilon}^{ \pm}:=(0,1) \times \partial \omega_{\xi, \varepsilon}^{ \pm}$. With reference to (6.55), we deduce from (6.56) and the continuity of the trace from $H^{1}(\check{\Omega})$ into $L^{2}(\check{\Gamma})$, that

$$
\begin{align*}
\left|\int_{\check{\Gamma}_{\xi, \varepsilon}^{ \pm}}\left(\partial_{\nu} u\right) \overline{u_{1}} d \sigma(x)\right| & \leqslant \int_{\check{\Gamma}_{\xi, \varepsilon}^{ \pm}}\left|\left(\partial_{\nu} u\right) e^{-\tau \xi \cdot x^{\prime}}\left(1+w_{1}(x)\right)\right| d \sigma\left(x^{\prime}\right) d x_{1} \\
& \leqslant C\left\|e^{-\tau \xi \cdot x^{\prime}} \partial_{\nu} u\right\|_{L^{2}\left(\check{\Gamma}_{\xi, \varepsilon}^{ \pm}\right)} \tag{6.60}
\end{align*}
$$

where $C$ is another positive constant depending only on $\omega$ and $M_{+}$. Moreover, we have

$$
\varepsilon\left\|e^{-\tau \xi \cdot x^{\prime}} \partial_{\nu} u\right\|_{L^{2}\left(\check{\Gamma}_{\xi, \varepsilon}^{+}\right)}^{2} \leqslant\left\|e^{-\tau \xi \cdot x^{\prime}}(\xi \cdot \nu)^{\frac{1}{2}} \partial_{\nu} u\right\|_{L^{2}\left(\check{\Gamma}_{\xi, \varepsilon}^{+}\right)}^{2} \leqslant\left\|e^{-\tau \xi \cdot x^{\prime}}(\xi \cdot \nu)^{\frac{1}{2}} \partial_{\nu} u\right\|_{L^{2}\left(\check{\Gamma}_{\xi, \varepsilon}^{+}\right)}^{2}
$$

from the very definition of $\partial \omega_{\xi, \varepsilon}^{+}$and the imbedding $\partial \omega_{\xi, \varepsilon}^{+} \subset \partial \omega_{\xi}^{+}$. Therefore, applying the Carleman estimate of Corollary 5.2 to the $\mathcal{H}_{\theta}^{2}(\check{\Omega})$-solution $u$ of (6.58), which is permitted since $\tau \geqslant \tau_{1}$, we get that

$$
\begin{align*}
\tau \varepsilon\left\|e^{-\tau \xi \cdot x^{\prime}} \partial_{\nu} u\right\|_{L^{2}\left(\check{\Gamma}_{\xi, \varepsilon}^{+}\right)}^{2} & \leqslant C\left(\left\|e^{-\tau \xi \cdot x^{\prime}}\left(-\Delta+V_{1}\right) u\right\|_{L^{2}(\check{\Omega})}^{2}+\tau\left\|e^{-\tau \xi \cdot x^{\prime}}|\xi \cdot \nu|^{\frac{1}{2}} \partial_{\nu} u\right\|_{L^{2}\left(\check{\Gamma}_{\xi}^{-}\right)}^{2}\right) \\
& \leqslant C\left(\left\|e^{-\tau \xi \cdot x^{\prime}} V u_{2}\right\|_{L^{2}(\check{\Omega})}^{2}+\tau\left\|e^{-\tau \xi \cdot x^{\prime}}|\xi \cdot \nu|^{\frac{1}{2}} \partial_{\nu} u\right\|_{L^{2}\left(\check{\Gamma}_{\xi}^{-}\right)}^{2}\right) \tag{6.61}
\end{align*}
$$

Next, as $e^{-\tau \xi \cdot x^{\prime}} V u_{2}(x)=e^{-\tau \xi \cdot x^{\prime}} V e^{\zeta_{2} \cdot x}\left(1+w_{2}(x)\right)=e^{-i\left(\pi k x_{1}+\frac{\eta \cdot x^{\prime}}{2}-\ell \cdot x\right)} V\left(1+w_{2}(x)\right)$ for a.e. $x \in \check{\Omega}$, from (4.44) and (6.55), we have $\left|e^{-\tau \xi \cdot x^{\prime}} V u_{2}(x)\right|=|V(x)|\left|1+w_{2}(x)\right|$ by (4.41), so it follows from (6.56) that

$$
\begin{equation*}
\left\|e^{-\tau \xi \cdot x^{\prime}} V u_{2}\right\|_{L^{2}(\check{\Omega})} \leqslant M_{+}(|\omega|+C) \tag{6.62}
\end{equation*}
$$

Further, bearing in mind that $\partial \omega_{\xi}^{-} \subset \partial \omega_{\xi, \varepsilon}^{-}$and $|\xi \cdot \nu| \leqslant 1$ on $\partial \omega_{\xi, \varepsilon}^{-}$, we get that

$$
\left\|e^{-\tau \xi \cdot x^{\prime}}|\xi \cdot \nu|^{\frac{1}{2}} \partial_{\nu} u\right\|_{L^{2}\left(\check{\Gamma}_{\xi}^{-}\right)} \leqslant\left\|e^{-\tau \xi \cdot x^{\prime}} \partial_{\nu} u\right\|_{L^{2}\left(\check{\Gamma}_{\xi, \varepsilon}^{-}\right)}
$$

which, together with (6.61)-(6.62), yield the estimate

$$
\left\|e^{-\tau \xi \cdot x^{\prime}} \partial_{\nu} u\right\|_{L^{2}\left(\check{\Gamma}_{\xi, \varepsilon}^{+}\right)}^{2} \leqslant \frac{C}{\varepsilon}\left(\frac{1}{\tau}+\left\|e^{-\tau \xi \cdot x^{\prime}} \partial_{\nu} u\right\|_{L^{2}\left(\check{\Gamma}_{\xi, \varepsilon}^{-}\right)}^{2}\right)
$$

where $C=C\left(\omega, M_{ \pm}\right)>0$. From this and (6.59)-(6.60), it follows that

$$
\begin{equation*}
\left|\int_{\check{\Omega}} V u_{2} \overline{u_{1}} d x\right| \leqslant C\left(\frac{1}{\tau}+\left\|e^{-\tau \xi \cdot x^{\prime}} \partial_{\nu} u\right\|_{L^{2}\left(\check{\Gamma}_{\xi, \varepsilon}^{-}\right)}^{2}\right)^{\frac{1}{2}} \tag{6.63}
\end{equation*}
$$

the positive constant $C$ depending this time on $M_{ \pm}$and $G^{\prime}$.
On the other hand, with reference to (4.44)-(4.45) and (6.55), we find through direct calculation that

$$
\begin{equation*}
\int_{\check{\Omega}} V u_{2} \overline{u_{1}} d x=\int_{(0,1) \times \mathbb{R}^{2}} e^{-i\left(2 \pi k x_{1}+\eta \cdot x^{\prime}\right)} V\left(x_{1}, x^{\prime}\right) d x_{1} d x^{\prime}+\int_{\check{\Omega}} R(x) d x \tag{6.64}
\end{equation*}
$$

where

$$
R(x):=e^{-i\left(2 \pi k x_{1}+\eta \cdot x^{\prime}\right)} V(x)\left(w_{2}(x)+\overline{w_{1}(x)}+w_{2}(x) \overline{w_{1}(x)}\right), x=\left(x_{1}, x^{\prime}\right) \in \check{\Omega}
$$

Since $\left\|w_{j}\right\|_{L^{2}(\Omega)}$, for $j=1,2$, is bounded (up to some multiplicative constant) from above by $\tau^{-1}$, according to (4.40), we obtain that

$$
\left|\int_{\check{\Omega}} R(x) d x\right| \leqslant M_{+}\left(|\omega|^{\frac{1}{2}}\left(\left\|w_{1}\right\|_{L^{2}(\check{\Omega})}+\left\|w_{2}\right\|_{L^{2}(\check{\Omega})}\right)+\left\|w_{1}\right\|_{L^{2}(\check{\Omega})}\left\|w_{2}\right\|_{L^{2}(\check{\Omega})}\right) \leqslant C \tau^{-1}
$$

where $C$ is independent of $\tau$. It follows from this and (6.63)-(6.64) that

$$
\begin{align*}
\left|\int_{(0,1) \times \mathbb{R}^{2}} e^{-i\left(2 \pi k x_{1}+\eta \cdot x^{\prime}\right)} V(x) d x\right|^{2} & \leqslant C\left(\frac{1}{\tau}+\left\|e^{-\tau \xi \cdot x^{\prime}} \partial_{\nu} u\right\|_{L^{2}\left(\check{\Gamma}_{\xi, \varepsilon}^{-}\right)}^{2}\right) \\
& \leqslant C\left(\frac{1}{\tau}+e^{c_{\omega} \tau}\left\|\partial_{\nu} u\right\|_{L^{2}\left(\check{\Gamma}_{\xi, \varepsilon}^{-}\right)}^{2}\right) \tag{6.65}
\end{align*}
$$

where $c_{\omega}:=\max \left\{\left|x^{\prime}\right| ; x^{\prime} \in \bar{\omega}\right\}$ and $C=C\left(\omega, M_{ \pm}, G^{\prime}\right)>0$. Finally, upon recalling that $u=v_{1}-u_{2}$, where $v_{1}$ is solution to (6.57) and $u_{2}$ satisfies (4.38) with $V=V_{2}$, we see that

$$
\partial_{\nu} u=\left(\Lambda_{V_{2}, \theta}-\Lambda_{V_{1}, \theta}\right) f, f=\mathcal{T}_{0, \theta} u_{2}
$$

Since $\partial \omega_{\xi, \varepsilon}^{-} \subset G^{\prime}$, by (6.51), we have $\left\|\partial_{\nu} u\right\|_{L^{2}\left(\check{\Gamma}_{\xi, \varepsilon}^{-}\right)} \leqslant\left\|\Lambda_{V_{2}, \theta}-\Lambda_{V_{1}, \theta}\right\|\left\|\mathcal{T}_{0, \theta} u_{2}\right\|_{\mathscr{H}_{\theta}(\check{\Gamma})}$, and hence

$$
\left\|\partial_{\nu} u\right\|_{L^{2}\left(\check{\Gamma}_{\xi, \varepsilon}^{-}\right)} \leqslant C\left\|\Lambda_{V_{2}, \theta}-\Lambda_{V_{1}, \theta}\right\|\left\|u_{2}\right\|_{H_{\Delta}(\check{\Omega})} \leqslant C \tau e^{c_{\omega} \tau}\left\|\Lambda_{V_{2}, \theta}-\Lambda_{V_{1}, \theta}\right\|
$$

by (3.34) and (4.40), the constant $C>0$ depending only on $\omega, M_{ \pm}$and $G^{\prime}$. This and (6.65) entail (6.52).
6.2. Two technical results. Prior to completing the derivation of Theorem 1.2 we collect two technical results that are needed in the remaining part of proof.

The first statement, which makes use of the following notation

$$
D_{s}:=\left\{x^{\prime} \in \mathbb{R}^{2} ;\left|x^{\prime}\right|<s\right\}, s \in(0,+\infty)
$$

is borrowed from [V, Theorem 1] and [AE, Theorem 3].
Lemma 6.2. For $R>0$, let $f: D_{2 R} \subset \mathbb{R}^{2} \rightarrow \mathbb{C}$ be real analytic and satisfy the condition

$$
\exists c>0, \exists \varrho \in(0,1], \forall \beta \in \mathbb{N}^{2},\left\|\partial^{\beta} f\right\|_{L^{\infty}\left(D_{2 R}\right)} \leqslant \frac{c|\beta|!}{(\varrho R)^{|\beta|}}
$$

Then for any $E \subset D_{\frac{R}{2}}$ with positive Lebesgue measure, we may find two constants $N=N(\varrho,|E|, R)>0$ and $v=v(\varrho,|E|, R) \in(0,1)$, such that we have:

$$
\|f\|_{L^{\infty}\left(D_{R}\right)} \leqslant N c^{1-v}\|f\|_{L^{1}(E)}^{v}
$$

Let us denote by $\hat{u}$ the Fourier transform with respect to $x^{\prime} \in \mathbb{R}^{2}$ of $u$, i.e.

$$
\hat{u}(\eta):=\frac{1}{2 \pi} \int_{\mathbb{R}^{2}} u\left(x^{\prime}\right) e^{-i \eta \cdot x^{\prime}} d x^{\prime}, \eta \in \mathbb{R}^{2}
$$

Then the second result is as follows.
Lemma 6.3. There exists $C>0$ such that the estimate

$$
\begin{equation*}
\|G\|_{H^{-1}\left((0,1) \times \mathbb{R}^{2}\right)} \leqslant C\left\|\sum_{k \in \mathbb{Z}}\left(1+|(k, \cdot)|^{2}\right)^{-\frac{1}{2}} \widehat{G_{k}}\right\|_{L^{2}\left(\mathbb{R}^{2}\right)} \tag{6.66}
\end{equation*}
$$

holds for every $G \in L^{2}\left((0,1) \times \mathbb{R}^{2}\right)$, with $G_{k}\left(x^{\prime}\right):=\int_{0}^{1} G\left(x_{1}, x^{\prime}\right) e^{-i 2 \pi k x_{1}} d x_{1}, x^{\prime} \in \mathbb{R}^{2}, k \in \mathbb{Z}$.
Proof. Let $C$ be any positive constant, satisfying

$$
\begin{equation*}
\left\|\sum_{k \in \mathbb{Z}}\left(1+|(k, \cdot)|^{2}\right)^{\frac{1}{2}} \widehat{w_{k}}\right\|_{L^{2}\left(\mathbb{R}^{2}\right)} \leqslant C\|w\|_{H^{1}\left((0,1) \times \mathbb{R}^{2}\right)} \tag{6.67}
\end{equation*}
$$

for all $w \in H_{0}^{1}\left((0,1) \times \mathbb{R}^{2}\right)$. Since $G \in L^{2}\left((0,1) \times \mathbb{R}^{2}\right)$, we have

$$
\langle G, w\rangle_{H^{-1}\left((0,1) \times \mathbb{R}^{2}\right), H_{0}^{1}\left((0,1) \times \mathbb{R}^{2}\right)}=\langle G, w\rangle_{L^{2}\left((0,1) \times \mathbb{R}^{2}\right)}=\int_{\mathbb{R}^{2}}\left(\sum_{k \in \mathbb{Z}} \widehat{G_{k}}(\eta) \widehat{\widehat{w_{k}}(\eta)}\right) d \eta
$$

by Plancherel formula. It follows from this and (6.67) that

$$
\left|\langle G, w\rangle_{H^{-1}\left((0,1) \times \mathbb{R}^{2}\right), H_{0}^{1}\left((0,1) \times \mathbb{R}^{2}\right)}\right| \leqslant C\left\|\sum_{k \in \mathbb{Z}}\left(1+|(k, \cdot)|^{2}\right)^{-\frac{1}{2}} \widehat{G_{k}}\right\|_{L^{2}\left(\mathbb{R}^{2}\right)}\|w\|_{H^{1}\left((0,1) \times \mathbb{R}^{2}\right)}
$$

entailing (6.66).
6.3. Completion of the proof. Let us express $\xi_{0} \in \mathbb{S}_{1}$ as $\xi_{0}=\left(\cos \alpha_{0}, \sin \alpha_{0}\right)$, where $\alpha_{0}$ is uniquely defined in $[0,2 \pi)$, and for $\varepsilon>0$ satisfying (6.51), pick $\alpha_{1} \in(0, \pi)$ such that

$$
\begin{equation*}
\left\{(\cos \alpha, \sin \alpha) ; \alpha \in\left(\alpha_{0}-\alpha_{1}, \alpha_{0}+\alpha_{1}\right)\right\} \subset\left\{z \in \mathbb{S}^{1} ;\left|z-\xi_{0}\right| \leqslant \varepsilon\right\} \tag{6.68}
\end{equation*}
$$

Next, for $\rho>0$ fixed, and for all $\eta \in \mathbb{R}^{2}$ and $k \in \mathbb{Z}$, we introduce

$$
\begin{equation*}
H_{k}(\eta):=\widehat{v_{k}}(\rho \eta)=\frac{1}{2 \pi} \int_{\mathbb{R}^{2}} v_{k}\left(x^{\prime}\right) e^{-i \rho \eta \cdot x^{\prime}} d x^{\prime} \text { with } v_{k}\left(x^{\prime}\right):=\int_{0}^{1} V\left(x_{1}, x^{\prime}\right) e^{-i 2 \pi k x_{1}} d x_{1} \tag{6.69}
\end{equation*}
$$

Since $v_{k}$ is supported in $\bar{\omega}$ and $0 \in \omega$ by assumption, the function $H_{k}$ is analytic in $\mathbb{R}^{2}$, and we get through elementary computations that $2 \pi\left|\partial^{\beta} H_{k}(\eta)\right| \leqslant\left\|v_{k}\right\|_{L^{1}(\omega)} c_{\omega}^{|\beta|} \rho^{|\beta|} \leqslant|\omega|^{\frac{1}{2}}\left\|v_{k}\right\|_{L^{2}(\omega)} c_{\omega}^{|\beta|} \rho^{|\beta|}$ for each $\beta \in \mathbb{N}^{2}$, where we recall from Subsection 6.1 that $c_{\omega}=\max \left\{\left|x^{\prime}\right|, x^{\prime} \in \omega\right\}$. As $\left\|v_{k}\right\|_{L^{2}(\omega)} \leqslant\|V\|_{L^{2}(\check{\Omega})} \leqslant M_{+}|\omega|^{\frac{1}{2}}$ from the Plancherel theorem, and $\rho^{|\beta|} \leqslant|\beta|!e^{\rho}$, this entails that

$$
\left|\partial^{\beta} H_{k}(\eta)\right| \leqslant \frac{M_{+}|\omega|}{2 \pi} e^{\rho} c_{\omega}^{|\beta|}|\beta|!, \beta \in \mathbb{N}^{2}
$$

Thus, applying Lemma 6.2 with $f=H_{k}, R=c_{\omega}^{-1}+1, c=M_{+}|\omega| e^{\rho} /(2 \pi), \varrho=\left(R c_{\omega}\right)^{-1}=\left(1+c_{\omega}\right)^{-1} \in(0,1]$, and

$$
\begin{equation*}
E:=\left\{t(-\sin \alpha, \cos \alpha) ; \alpha \in\left(\alpha_{0}-\alpha_{1}, \alpha_{0}+\alpha_{1}\right), t \in[0, \min (1, R / 2))\right\} \tag{6.70}
\end{equation*}
$$

we find that

$$
\begin{equation*}
\left\|H_{k}\right\|_{L^{\infty}\left(D_{R}\right)} \leqslant C e^{\rho(1-v)}\left\|H_{k}\right\|_{L^{\infty}(E)}^{v}, \tag{6.71}
\end{equation*}
$$

where $C=C\left(\omega, M_{+}, G^{\prime}\right)>0$ and $v=v\left(\omega, M_{+}, G^{\prime}\right) \in(0,1)$. Moreover, in view of (6.68)-(6.70), Lemma 6.1 tells us that the estimate (6.52) holds uniformly in $\eta \in E$, i.e. $\left\|H_{k}\right\|_{L^{\infty}(E)}^{2} \leqslant C\left(\tau^{-1}+e^{C^{\prime} \tau}\left\|\Lambda_{V_{2}, \theta}-\Lambda_{V_{1}, \theta}\right\|^{2}\right)$, provided $r \geqslant r_{1}$. It follows from this and (6.71) that

$$
\begin{equation*}
\left|\widehat{v_{k}}(\eta)\right|^{2} \leqslant C e^{2 \rho(1-v)}\left(\frac{1}{\tau}+e^{C^{\prime} \tau}\left\|\Lambda_{V_{2}, \theta}-\Lambda_{V_{1}, \theta}\right\|^{2}\right)^{v},|(k, \eta)|<\rho, r \geqslant r_{1} \tag{6.72}
\end{equation*}
$$

The next step is to apply Lemma 6.3 with $G=V$ : With reference to (6.69), we obtain that

$$
\begin{align*}
\|V\|_{H^{-1}\left((0,1) \times \mathbb{R}^{2}\right)}^{2} & \leqslant C \int_{\mathbb{R}^{2}} \sum_{k=-\infty}^{+\infty}\left(1+|(k, \eta)|^{2}\right)^{-1}\left|\widehat{v_{k}}(\eta)\right|^{2} d \eta \\
& \leqslant C \int_{\mathbb{R}^{3}}\left(1+|(k, \eta)|^{2}\right)^{-1}\left|\widehat{v_{k}}(\eta)\right|^{2} d \mu(k) d \eta \tag{6.73}
\end{align*}
$$

where $\mu:=\sum_{n \in \mathbb{Z}} \delta_{n}$. Putting $B_{\rho}:=\left\{(k, \eta) \in \mathbb{R}^{3},|(k, \eta)| \leqslant \rho\right\}$ we now examine the two integrals $\int_{B_{\rho}}(1+$ $\left.|(k, \eta)|^{2}\right)^{-1}\left|\widehat{v_{k}}(\eta)\right|^{2} d \mu(k) d \eta$ and $\int_{\mathbb{R}^{3} \backslash B_{\rho}}\left(1+|(k, \eta)|^{2}\right)^{-1}\left|\widehat{v_{k}}(\eta)\right|^{2} d \mu(k) d \eta$ separately. The last one, is easily treated, as we have

$$
\int_{\mathbb{R}^{3} \backslash B_{\rho}}\left(1+|(k, \eta)|^{2}\right)^{-1}\left|\widehat{v_{k}}(\eta)\right|^{2} d \mu(k) d \eta \leqslant \frac{1}{\rho^{2}} \int_{\mathbb{R}^{3} \backslash B_{\rho}}\left|\widehat{v_{k}}(\eta)\right|^{2} d \mu(k) d \eta \leqslant \frac{1}{\rho^{2}} \int_{\mathbb{R}^{3}}\left|\widehat{v_{k}}(\eta)\right|^{2} d \mu(k) d \eta
$$

and hence

$$
\begin{equation*}
\int_{\mathbb{R}^{3} \backslash B_{\rho}}\left(1+|(k, \eta)|^{2}\right)^{-1}\left|\widehat{v_{k}}(\eta)\right|^{2} d \mu(k) d \eta \leqslant \frac{1}{\rho^{2}} \int_{(0,1) \times \mathbb{R}^{2}}\left|V\left(x_{1}, x^{\prime}\right)\right|^{2} d x_{1} d x^{\prime} \leqslant \frac{|\omega| M_{+}^{2}}{\rho^{2}} \tag{6.74}
\end{equation*}
$$

by Parseval-Plancherel theorem. We turn now to studying the first integral. To this end we notice upon setting $\mathcal{C}_{\rho}:=\mathbb{R} \times D_{1 / \rho}$, where we recall that $D_{1 / \rho}=\left\{\eta \in \mathbb{R}^{2} ;|\eta|<1 / \rho\right\}$, that

$$
\begin{align*}
& \int_{B_{\rho} \cap \mathcal{C}_{\rho}}\left(1+|(k, \eta)|^{2}\right)^{-1}\left|\widehat{v_{k}}(\eta)\right|^{2} d \eta d \mu(k) \\
\leqslant & \int_{\mathcal{C}_{\rho}}\left|\widehat{v_{k}}(\eta)\right|^{2} d \mu(k) d \eta \leqslant \int_{D_{1 / \rho}}\left(\sum_{k \in \mathbb{Z}}\left|\widehat{v_{k}}(\eta)\right|^{2}\right) d \eta \leqslant \frac{\pi}{\rho^{2}}\left\|\sum_{k \in \mathbb{Z}}\left|\widehat{v_{k}}\right|^{2}\right\|_{L^{\infty}\left(D_{1 / \rho}\right)} \tag{6.75}
\end{align*}
$$

Moreover, since each $v_{k}$, for $k \in \mathbb{Z}$, is supported in $\omega$ by (6.69), we have

$$
\left|\widehat{v_{k}}(\eta)\right|^{2} \leqslant \frac{1}{4 \pi^{2}}\left\|v_{k}\right\|_{L^{1}(\omega)}^{2} \leqslant \frac{|\omega|}{4 \pi^{2}}\left\|v_{k}\right\|_{L^{2}(\omega)}^{2}, \quad \eta \in \mathbb{R}^{2}
$$

which entails

$$
\sum_{k \in \mathbb{Z}}\left|\widehat{v_{k}}(\eta)\right|^{2} \leqslant \frac{|\omega|}{4 \pi^{2}}\left(\sum_{k \in \mathbb{Z}}\left\|v_{k}\right\|_{L^{2}(\omega)}^{2}\right) \leqslant \frac{|\omega|}{4 \pi^{2}}\|V\|_{L^{2}(\check{\Omega})}^{2} \leqslant \frac{M_{+}^{2}|\omega|^{2}}{4 \pi^{2}}, \eta \in \mathbb{R}^{2}
$$

upon applying the Parseval formula and the dominated convergence theorem. Putting this together with (6.75), we obtain that

$$
\begin{equation*}
\int_{B_{\rho} \cap \mathcal{C}_{\rho}}\left(1+|(k, \eta)|^{2}\right)^{-1}\left|\widehat{v_{k}}(\eta)\right|^{2} d \eta d \mu(k) \leqslant \frac{M_{+}^{2}|\omega|^{2}}{4 \pi \rho^{2}} \tag{6.76}
\end{equation*}
$$

Further, if $(k, \eta) \in B_{\rho} \cap\left(\mathbb{R}^{3} \backslash \mathcal{C}_{\rho}\right)$, then there exists a positive constant $C$ such that $\tau(k, \eta, \theta, r) \in\left(r, C \rho^{2} r\right]$ for all $\rho \geqslant 1$ and $r \geqslant r_{1}$, according to (4.43). As a consequence we have

$$
\left|\widehat{v_{k}}(\eta)\right|^{2} \leqslant C e^{2 \rho(1-v)}\left(\frac{1}{r}+e^{C \rho^{2} r}\left\|\Lambda_{V_{2}, \theta}-\Lambda_{V_{1}, \theta}\right\|^{2}\right)^{v}, r \geqslant r_{1},(k, \eta) \in B_{\rho} \cap\left(\mathbb{R}^{3} \backslash \mathcal{C}_{\rho}\right)
$$

by (6.72), and hence

$$
\begin{equation*}
\int_{B_{\rho} \cap\left(\mathbb{R}^{3} \backslash \mathcal{C}_{\rho}\right)}\left(1+|(k, \eta)|^{2}\right)^{-1}\left|\widehat{v_{k}}(\eta)\right|^{2} d \mu(k) d \eta \leqslant C \rho^{3} e^{2 \rho(1-v)}\left(\frac{1}{r}+e^{C \rho^{2} r}\left\|\Lambda_{V_{2}, \theta}-\Lambda_{V_{1}, \theta}\right\|^{2}\right)^{v}, r \geqslant r_{1} \tag{6.77}
\end{equation*}
$$

upon eventually substituting $C$ for some suitable algebraic expression of $C$.
Therefore, putting (6.73)-(6.74) and (6.76)-(6.77) together, we find for all $\rho \geqslant 1$ and $r \geqslant r_{1}$ that

$$
\begin{align*}
\|V\|_{H^{-1}\left((0,1) \times \mathbb{R}^{2}\right)}^{\frac{2}{v}} & \leqslant C\left(\frac{1}{\rho^{2}}+\rho^{3} e^{2 \rho(1-v)}\left(\frac{1}{r}+e^{C \rho^{2} r}\left\|\Lambda_{V_{2}, \theta}-\Lambda_{V_{1}, \theta}\right\|^{2}\right)^{v}\right)^{\frac{1}{v}} \\
& \leqslant C\left(\rho^{-\frac{2}{v}}+\rho^{\frac{3}{v}} e^{2 \rho\left(\frac{1}{v}-1\right)} r^{-1}+\rho^{\frac{3}{v}} e^{2 \rho\left(\frac{1}{v}-1\right)} e^{C \rho^{2} r}\left\|\Lambda_{V_{2}, \theta}-\Lambda_{V_{1}, \theta}\right\|^{2}\right) \tag{6.78}
\end{align*}
$$

Here we substituted $C=C\left(\omega, M_{+}, F^{\prime}, v\right)>0$ for $2^{\frac{1}{v}} C$ in the last line. Thus, taking $r:=\rho^{\frac{5}{v}} e^{2 \rho\left(\frac{1}{v}-1\right)}$ in (6.78), with $\rho \geqslant \rho_{1}$, in such a way that $r \geqslant r_{1}$ and $\rho^{\frac{3}{v}} e^{2 \rho\left(\frac{1}{v}-1\right)} r^{-1}=\rho^{-\frac{2}{v}}$, we obtain that

$$
\begin{equation*}
\|V\|_{H^{-1}\left((0,1) \times \mathbb{R}^{2}\right)} \leqslant C\left(\rho^{-\frac{2}{v}}+\rho^{\frac{3}{v}} e^{2 \rho\left(\frac{1}{v}-1\right)} \exp \left(C \rho^{\left(2+\frac{5}{v}\right)} e^{2 \rho\left(\frac{1}{v}-1\right)}\right)\left\|\Lambda_{V_{2}, \theta}-\Lambda_{V_{1}, \theta}\right\|^{2}\right)^{\frac{v}{2}} \tag{6.79}
\end{equation*}
$$

Let $C^{\prime}>0$ be so large that $\rho^{\frac{3}{v}} e^{2 \rho\left(\frac{1}{v}-1\right)} \exp \left(C \rho^{\left(2+\frac{5}{v}\right)} e^{2 \rho\left(\frac{1}{v}-1\right)}\right) \leqslant \exp \left(e^{C^{\prime} \rho}\right)$ for every $\rho \geqslant \rho_{1}$. Notice that $C^{\prime}$ depends only on $C$ and $v$, hence on $\omega, M_{+}$and $G^{\prime}$. Then (6.79) entails that

$$
\begin{equation*}
\|V\|_{H^{-1}\left((0,1) \times \mathbb{R}^{2}\right)} \leqslant C\left(\rho^{-\frac{2}{v}}+\exp \left(e^{C^{\prime} \rho}\right)\left\|\Lambda_{V_{2}, \theta}-\Lambda_{V_{1}, \theta}\right\|^{2}\right)^{\frac{v}{2}}, \rho \geqslant \rho_{1} \tag{6.80}
\end{equation*}
$$

Set $\gamma^{*}:=\exp \left(-e^{C^{\prime} \rho_{1}}\right) \in(0,1)$. If $\gamma:=\left\|\Lambda_{V_{2}, \theta}-\Lambda_{V_{1}, \theta}\right\| \in\left(0, \gamma^{*}\right]$, we find upon taking $\rho=\frac{1}{C^{\prime}} \ln (|\ln \gamma|) \geqslant \rho_{1}$ in (6.80), that

$$
\begin{equation*}
\|V\|_{H^{-1}(\check{\Omega})} \leqslant\|V\|_{H^{-1}\left((0,1) \times \mathbb{R}^{2}\right)} \leqslant C\left(C^{\frac{2}{v}}+\gamma(\ln |\ln \gamma|)^{\frac{2}{v}}\right)^{\frac{v}{2}}(\ln |\ln \gamma|)^{-1} \tag{6.81}
\end{equation*}
$$

Further, since $\sup _{\gamma \in\left(0, \gamma_{*}\right]}\left(C^{\prime \frac{2}{v}}+\gamma(\ln |\ln \gamma|)^{\frac{2}{v}}\right.$ is just another positive constant depending only on $\omega, M$ and $G^{\prime}$, we end up getting from (6.81) that

$$
\begin{equation*}
\|V\|_{H^{-1}(\check{\Omega})} \leqslant C(\ln |\ln \gamma|)^{-1}, \gamma \in\left(0, \gamma^{*}\right] . \tag{6.82}
\end{equation*}
$$

Finally, as $\|V\|_{H^{-1}(\check{\Omega})} \leqslant C\|V\|_{L^{\infty}(\check{\Omega})} \leqslant\left(C M_{+} / \gamma^{*}\right) \gamma$ for $\gamma>\gamma^{*}$, where $C>0$ is a constant depending only on $\omega$, we deduce (3.37) from this and (6.82).

Acknowledgments. The authors are much indebted to G. Uhlmann for helpful comments and suggestions.

## References

[Al] G. Alessandrini, Stable determination of conductivity by boundary measurements, Appl. Anal., 27 (1988), no. 1-3, 153-172.
[AE] J. Apraiz, L. Escauriaza, Null-control and measurable sets, ESAIM: Control, Optimisation and Calculus of Variations, 19 (2013), 239-254.
[BKS] M. Bellassoued, Y. Kian, E. Soccorsi, An inverse stability result for non compactly supported potentials by one arbitrary lateral Neumann observation, preprint, arXiv:1507.06864.
[BIY] E. Blasten, O. Imanuvilov, M. Yamamoto, Stability and uniqueness for a two-dimensional inverse boundary value problem for less regular potentials, Inverse Prob. Imaging, 9 (2015), no 3, 709-723.
[B] A. Bukhgeim, Recovering the potential from Cauchy data in two dimensions, J. Inverse Ill-Posed Probl., 16 (2008), 19-34.
[BU] A. L. Bukhgeim, G. Uhlmann, Recovering a potential from partial Cauchy data, Comm. PDE, 27 (2002), no 3-4, 653-668.
[CDR1] P. Caro, D. Dos Santos Ferreira, A. Ruiz, Stability estimates for the Radon transform with restricted data and applications, Advances in Math., 267 (2014), 523-564.
[CDR2] P. Caro, D. Dos Santos Ferreira, A. Ruiz, Stability estimates for the Calderón problem with partial data, preprint, arXiv:1405.1217.
[Ch] M. Choulli, Une introduction aux problèmes inverses elliptiques et paraboliques, Mathématiques et Applications, Vol. 65, Springer-Verlag, Berlin, 2009.
[CKS] M. Choulli, Y. Kian, E. Soccorsi, Stable determination of time-dependent scalar potential from boundary measurements in a periodic quantum waveguide, SIAM Jour. Math. Analysis, 47 (2015), no 6, 4536-4558.
[CKS1] M. Choulli, Y. Kian, E. Soccorsi, Double logarithmic stability estimate in the identification of a scalar potential by a partial elliptic Dirichlet-to-Neumann map, Bulletin of the South Ural State University, Ser. Mathematical Modelling, Programming and Computer Software (SUSU MMCS) 8 (2015), no 3, 78-95.
[CS] M. Choulli, E. Soccorsi, An inverse anisotropic conductivity problem induced by twisting a homogeneous cylindrical domain, J. Spec. Theory, 5 (2015), 295-329.
[Di] J. Dixmier, Von Neumann Algebras, North Holland, Amsterdam, 1981.
[HT] B. Haberman, D. Tataru, Uniqueness in Calderón's problem with Lipschitz conductivities, Duke Math. Journal, 162 (2013), no. 3, 497-516.
[Ha] P. HÄhner, A periodic Faddeev-type operator, J. Diff. Equat., 128 (1996), 300-308.
[HW] H. Heck, J. N. WAng, Stability estimate for the inverse boundary value problem by partial Cauchy data, Inverse Problems, 22 (2006), 1787-1797.
[IUY1] O. Imanuvilov, G. Uhlmann, M. Yamamoto, The Calderón problem with partial data in two dimensions, Journal American Math. Society, 23 (2010), 655-691.
[IUY2] O. Imanuvilov, G. Uhlmann, M. Yamamoto, Partial Cauchy data for general second order elliptic operators in two dimensions, Publ. Research Institute Math. Sci., 48 (2012), 971-1055.
[Ik] M. Ikehata, Inverse conductivity problem in the infinite slab, Inverse Problems 17 (2001), 437-454.
[KKS] O. Kavian, Y. Kian, E. Soccorsi, Uniqueness and stability results for an inverse spectral problem in a periodic waveguide, Jour. Math. Pures Appl., 104 (2015), no. 6, 1160-1189.
[KSU] C.E. Kenig, J. Sjöstrand, G. Uhlmann, The Calderon problem with partial data, Ann. of Math., 165 (2007), $567-591$.
[Ki] Y. KIAN, Stability of the determination of a coefficient for wave equations in an infinite waveguide, Inverse Probl. Imaging, 8 (3) (2014), 713-732.
[KPS1] Y. Kian, Q. S. Phan, E. Soccorsi, Carleman estimate for infinite cylindrical quantum domains and application to inverse problems, Inverse Problems 30, 5 (2014), 055016.
[KPS2] Y. Kian, Q. S. Phan, E. Soccorsi, Hölder stable determination of a quantum scalar potential in unbounded cylindrical domains, Jour. Math. Anal. Appl. 426, 1 (2015), 194-210.
[KV] R. Kohn, M. Vogelius, Determining conductivity by boundary measurements, Comm. Pure Appl. Math., 37 (1984), 289298.
[KLU] K. Krupchyk, M. Lassas, G. Uhlmann, Inverse Problems with Partial Data for a Magnetic Schrödinger Operator in an Infinite Slab or Bounded Domain, Comm. Math. Phys., 312 (2012), 87-126.
[LU] X. Li, G. Uhlmann, Inverse Problems on a Slab, Inverse Problems and Imaging, 4(2010), 449-462.
[LM1] J.-L. Lions, E. Magenes, Problèmes aux limites non homogènes et applications, vol. 1, Dunod (1968).
[NS] A. Nachman, B. Street, Reconstruction in the Calderón problem with partial data, Comm. PDE, 35 (2010), 375-390.
[NSa] R.G. Novikov, M. Santacesaria, A global stability estimate for the Gel'fand-Calderón inverse problem in two dimensions, J. Inverse Ill Posed Probl., 18 (2010), 765-785.
[RS1] M. Reed, B. Simon, Methods of Modern Mathematical Physics II: Fourier Analysis, Self-adjointness, Academic Press, 1978.
[RS2] M. Reed, B. Simon, Methods of Modern Mathematical Physics IV: Analysis of Operators, Academic Press, 1978.
[SW] M. Salo, J. N. Wang, Complex spherical waves and inverse problems in unbounded domains, Inverse Problems, 22 (2006), 2299-2309.
[Sa] M. Santacesaria, New global stability estimates for the Calderón problem in two dimensions, J. Inst. Math. Jussieu, 12 (2013), 553-569.
[Sc] L. Schwartz, Topologie générale et analyse fonctionnelle, Hermann, Paris, 1970.
[SU] J. Sylvester, G. Uhlmann, A global uniqueness theorem for an inverse boundary value problem, Ann. of Math., $\mathbf{1 2 5}$ (1987), 153-169.
[T] L. Tzou, Stability Estimate for the coefficients of magnetic Schrödinger equation from full and partial boundary measurements, Commun. PDE, 11 (2008), 1911-1952.
[V] S. Vessella, A continuous dependence result in the analytic continuation problem, Forum math., 11 (1999), no. 6, 695-703.
[Y] Y. Yang, Determining the first order perturbation of a bi-harmonic operator on bounded and unbounded domains from partial data, J. Differential Equations, 257 (2014), 3607-3639.

Mourad Choulli, Université de Lorraine, Institut Elie Cartan de Lorraine, CNRS, UMR 7502, Boulevard des Aiguillettes, BP 70239, 54506 Vandoeuvre les Nancy cedex - Ile du Saulcy, 57045 Metz cedex 01, France. E-mail: mourad.choulli@univ-lorraine.fr.

Yavar Kian, Aix-Marseille Université, CNRS, CPT UMR 7332, 13288 Marseille, and Université de Toulon, CNRS, CPT UMR 7332, 83957 La Garde, France.
E-mail: yavar.kian@univ-amu.fr.

Eric Soccorsi, Aix-Marseille Université, CNRS, CPT UMR 7332, 13288 Marseille, and Université de Toulon, CNRS, CPT UMR 7332, 83957 La Garde, France.
E-mail: eric.soccorsi@univ-amu.fr.

