# DETERMINING THE POTENTIAL AND THE GRADIENT COUPLING OF TWO-STATE QUANTUM SYSTEMS IN AN INFINITE WAVEGUIDE 

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#### Abstract

We consider the inverse coefficient problem of simultaneously determining the space dependent electric potential, the zero-th order coupling term and the first order coupling vector of a two-state Schrödinger equation in an infinite cylindrical domain of $\mathrm{R} \mathrm{n}, n[?] 2$, from finitely many partial boundary measurements of the solution. We prove that these $n+3$ unknown scalar coefficients can be Holder stably retrieved by ( $n+1$ )-times suitably changing the initial condition attached at the system.


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#### Abstract

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Keywords: Inverse problem, stability estimate, two-state Schrödinger equation.

Mathematics subject classification 2010: 35R30.

## 1. Introduction

This is the second of two papers dealing with the stability issue in the inverse problem of determining the electric potential and the coupling coefficients of a two-state quantum system, from local Neumann data. Two-state systems are the simplest of non-trivial quantum systems that occur in nature. Examples include the quantum bit, or qubit, which is the basic unit of quantum information in quantum computing, i.e. the quantum version of the classical binary bit. In a classical system a bit is one state or the other but in quantum mechanics the qubit can be in a superposition of both states simultaneously. This feature is fundamental to quantum computing. Another well known example is the spin of the electron, in which the two levels can be taken by spin-up and spin-down. The two state formalism is used to describe simple mixing of states, which leads to quantum phenomena with tremendous industrial applications such as maser or laser. It is made of Schrödinger equations bound together through a linear gradient coupling. When the two quantum states are constrained to a bounded spatial domain, it was proved in [19] that the electric potential and the coupling are stably determined by finitely many partial boundary observations of the system. In the present work we aim for the same type of identification result when the quantum motion is no longer bounded but may escape to infinity in one direction over the course of time.
1.1. Settings. Throughout this article, $\omega$ is a bounded domain of $\mathbb{R}^{n-1}, n \geq 2$, with smooth boundary $\gamma:=\partial \omega$, and $\Omega:=\omega \times \mathbb{R}$. For $T \in \mathbb{R}_{+}$, we consider the following initial-boundary value problem (IBVP) with initial states $u_{0}^{ \pm}$ and non-homogenous Dirichlet boundary conditions $g^{ \pm}$, for the coupled Schrödinger equations in the unknowns $u^{ \pm}$,

$$
\begin{cases}-i \partial_{t} u^{+}-\Delta u^{+}+q^{+} u^{+}+A \cdot \nabla u^{-}+p u^{-}=0 & \text { in } Q:=\Omega \times(0, T)  \tag{1.1}\\ -i \partial_{t} u^{-}-\Delta u^{-}+q^{-} u^{-}-A \cdot \nabla u^{+}+p u^{+}=0 & \text { in } Q \\ u^{+}(\cdot, 0)=u_{0}^{+}, u^{-}(\cdot, 0)=u_{0}^{-} & \text {in } \Omega \\ u^{+}=g^{+}, u^{-}=g^{-} & \text {on } \Sigma:=\Gamma \times(0, T),\end{cases}
$$

where $\Gamma:=\gamma \times \mathbb{R}$. Since $\Gamma$ is unbounded, let us make the above boundary condition more precise. For all $x \in \Omega$, we write $x=\left(x^{\prime}, x_{n}\right)$ where $x^{\prime}=\left(x_{1}, \ldots, x_{n-1}\right) \in \omega$ and $x_{n} \in \mathbb{R}$, and using a standard density argument we extend the mapping

$$
\begin{aligned}
C_{0}^{\infty}\left(\mathbb{R} \times(0, T), H^{2}(\omega)\right) & \rightarrow L^{2}\left(\mathbb{R} \times(0, T), H^{\frac{3}{2}}(\omega)\right) \\
w & \mapsto\left[\left(x_{n}, t\right) \in \mathbb{R} \times(0, T) \mapsto w\left(\cdot, x_{n}, t\right)_{\mid \gamma}\right]
\end{aligned}
$$

to a bounded operator $\gamma_{0}$ acting from $L^{2}\left(\mathbb{R} \times(0, T), H^{2}(\omega)\right)$ into $L^{2}\left(\mathbb{R} \times(0, T), H^{\frac{3}{2}}(\gamma)\right)$. Then, for all $u^{ \pm} \in$ $L^{2}\left(0, T, H^{2}(\Omega)\right)$, the boundary condition in (1.1) reads $\gamma_{0} u^{ \pm}=g^{ \pm}$.

In the present paper we aim to stably retrieve the electric potentials $q^{ \pm}: \Omega \rightarrow \mathbb{R}$, the zero-th order coupling term $p: \Omega \rightarrow \mathbb{R}$ and the first order coupling vector $A: \Omega \rightarrow \mathbb{R}^{n}$, by finitely many partial boundary measurements over the entire time-span $(0, T)$ of the solution $u^{ \pm}$to (1.1). In contrast with [19] where the spatial domain $\Omega$ is bounded, here we consider an infinitely extended cylindrical domain and we address the problem of simultaneous identification of non-compactly supported unknown coefficients $p, q^{ \pm}$and $A$. This requires a slightly different and technically more demanding approach than the one implemented in [19].
1.2. Motivations. The dynamics of the two states $u^{ \pm}$governed by (1.1) are bound together through linear gradient coupling $p u^{\mp} \pm A \cdot \nabla u^{\mp}$. We refer the reader to [24] and the references therein for the relevance of these processes in physics. Gradient coupling appears also naturally in quantum fields theory (see [2,25]) or quantum cosmology (see [9,14]), and it is sometimes a first-order approximation of nonlinear coupling (see [28]).

Quantum wires are extremely narrow cylindrical structures which have a length-to-diameter ratio up to $10^{8}$. Examples include carbon nanotubes, which are made of carbon a with diameter typically measured in nanometers. These structures are commonly modeled by infinite three-dimensional cylindrical domains such as $\Omega$, in which electrons are essentially free to move in one direction. Quantum wires exhibit valuable physical properties for electronics, optics and other fields of materials science and technology, see e.g., [1], and for this reason they have attracted a lot of attention from the scientific community.

The IBVP (1.1) can be interpreted as the time-evolution of the spin of a spin- $\frac{1}{2}$ particle such as an electron (whose spin can have values $\pm \frac{\hbar}{2}$, where $\hbar$ is the reduced Planck constant) confined in a carbon nanotube, see e.g., [13, 16]. Notice that for the sake of notational simplicity, the various physical constants such as $\hbar$, the charge and the mass of the electron, are all taken equal to one in (1.1).
1.3. Bibliography. The mathematical literature devoted to inverse coefficient problems for the dynamic Schrödinger equation is so extensive that this presentation is not intended to be exhaustive, but we can mention $[4,5,7,8,22]$ where zero-th or/and first order unknown coefficients of the Schrödinger equation are determined by the Dirichlet-to-Neumann map. These articles assume knowledge of infinitely many boundary data, but in [3, 27] the real-valued electric potential is stably retrieved by one partial lateral observation of the solution. This result was extended to complex-valued electric potentials in [17]. The boundary measurement in [3, 17, 27] is taken on a subpart of the boundary fulfilling a geometric condition related to geometric optics condition insuring observability. This condition was relaxed to arbitrarily small sub-boundaries in [5], provided the potential is known in the vicinity of the boundary. The inverse problem of determining the magnetic vector potential of the autonomous Schrödinger equation is addressed in [17]. The same problem for the space-varying part of the magnetic potential appearing in a nonautonomous Schrödinger equation is treated in [12]. In both cases, the $n$-th dimensional unknown magnetic vector potential, $n \geq 1$, is retrieved from $n$ partial Neumann data obtained by $n$-times suitably selecting the initial condition attached at the magnetic Schrödinger equation.

The strategy of [3, 12, 17, 27] relies on a Carleman inequality specifically designed for the Schrödinger equation, see [17, 26, 27] for actual examples of such weighted energy estimates. The idea of using a Carleman estimate for solving inverse problems goes back to 1981 and was introduced by A. L. Bukhgeim and M. V. Klibanov in their seminal article [11]. Since then, the Bukhgeim-Klibanov approach has been successfully applied to parabolic, hyperbolic and Schrödinger systems and even to coupled systems of partial differential equations. We refer the reader to [18] and references therein, for a complete survey of multidimensional inverse problems solved by the Bukhgeim-Klibanov method.

In all the aforementioned papers, the Schrödinger equation under study is posed in a bounded spatial domain. The inverse problem of determining the electric potential of the Schrödinger equation stated in an infinite waveguide is examined in [6,21]. This is achieved by mean of a specifically designed Carleman estimate for the Schrödinger equation in an unbounded cylindrical domain, which is established in [20]. All the articles listed above are concerned with the "one state" Schrödinger equation. In [23], assuming that the gradient coupling vector is known, the authors show that the zero-th order coupling term of a two state magnetic Schrödinger equation is uniquely determined by one partial Neumann data. Recently in [29], the electric potential of a strongly coupled Schrödinger equations in a bounded spatial domain was Lipschitz stably retrieved by one partial (internal or boundary) measurement of the solution to the
system. In [19], the zero-th and first order coefficients of the coupling are Lispchitz stably recovered by finitely many partial boundary observations of the solution. The coupled Schrödinger equations under study in [19, 23, 29] are posed on a bounded spatial domain. In the present paper, we aim to extend the result of [19] to the case of an unbounded waveguide.
1.4. Notations. Throughout this text $x=\left(x_{1}, \ldots, x_{n}\right)$ is a generic point of $\bar{\Omega}$ that is sometimes written $x=\left(x^{\prime}, x_{n}\right)$ where $x^{\prime}=\left(x_{1}, \cdots, x_{n-1}\right) \in \bar{\omega}$ is the variable of the transverse section of $\Omega$ and $x_{n} \in \mathbb{R}$ is the longitudinal variable. For all $x=\left(x^{\prime}, x_{n}\right) \in \Gamma$, the outward unit normal $\nu$ to $\Gamma$ reads $\nu(x)=\nu\left(x^{\prime}\right)=\left(\nu^{\prime}\left(x^{\prime}\right), 0\right)^{T}$, where $\nu^{\prime}\left(x^{\prime}\right) \in \mathbb{R}^{n-1}$ is the outgoing normal vector to $\gamma$ at $x^{\prime}$ and $a^{T}$ denotes the transpose of the row vector $a$.

For all $i=1, \ldots, n$ we set $\partial_{i}:=\frac{\partial}{\partial x_{i}}$ in such a way that $\nabla:=\left(\partial_{1}, \ldots, \partial_{n}\right)^{T}\left(\right.$ resp., $\left.\nabla^{\prime}:=\left(\partial_{1}, \ldots, \partial_{n-1}\right)^{T}\right)$ is the gradient operator with respect to $x=\left(x_{1}, \ldots, x_{n}\right)$ (resp., $x^{\prime}=\left(x_{1}, \ldots, x_{n-1}\right)$ ). Similarly, we write $\partial_{t}=\frac{\partial}{\partial t}$. For the sake of shortness we write $\partial_{i j}^{2}, i, j=1, \ldots, n$, instead of $\partial_{i} \partial_{j}$ and as usual we denote by $\Delta$ the Laplace operator $\partial_{1}^{2}+\ldots+\partial_{n}^{2}$. Next, for any multi-index $k=\left(k_{1}, \ldots, k_{n}\right) \in \mathbb{N}_{0}^{n}$, where $\mathbb{N}_{0}:=\{0\} \cup \mathbb{N}$, we put $|k|:=k_{1}+\ldots+k_{n}$ and $\partial_{x}^{k}=\partial_{1}^{k_{1}} \ldots \partial_{n}^{k_{n}}$.

Further, the symbol $\cdot$ denotes the scalar product in $\mathbb{C}^{m}, m \in \mathbb{N}$, and we set $|\zeta|:=\sqrt{\zeta \cdot \zeta}$ for all $\zeta \in \mathbb{C}^{m}$. We simply write $\nabla \cdot$ for the divergence operator in $\mathbb{R}^{n}$ and we set $\partial_{\nu} u:=\nabla u \cdot \nu=\nabla^{\prime} \cdot \nu^{\prime}$.

Finally, for all $r>0$ and $s>0$, we introduce $H^{r, s}(\Sigma):=L^{2}\left(0, T ; H^{r}(\Gamma)\right) \cap H^{s}\left(0, T ; L^{2}(\Gamma)\right)$ where $H^{s}(\Gamma)$ denotes the usual Sobolev space on $\Gamma$ of order $s$.
1.5. Main results. Prior to investigating the inverse problem under study in this article, we examine the wellposedness issue for the forward problem associated with (1.1). For this purpose we introduce the Hamiltonian operator acting on $\left(C_{0}^{\infty}(Q)^{\prime}\right)^{2}$,

$$
\mathcal{H}\left(A, p, q^{ \pm}\right):=\left(\begin{array}{cc}
-\Delta+q^{+} & A \cdot \nabla+p \\
-A \cdot \nabla+p & -\Delta+q^{-}
\end{array}\right)
$$

and state the following existence, uniqueness and regularity result for the solution to the IBVP (1.1).
Proposition 1.1. Let $m \in \mathbb{N}$ and assume that $\gamma$ is $\mathcal{C}^{2(m+1)}$. Let $A \in W^{2 m+1, \infty}\left(\Omega, \mathbb{R}^{n}\right) \cap C^{2(m-1)}\left(\bar{\Omega}, \mathbb{R}^{n}\right)$ be such that $\nabla \cdot A=0$ a.e. in $\Omega$, let $p \in W^{2 m+1, \infty}(\Omega, \mathbb{R}) \cap C^{2(m-1)}(\bar{\Omega}, \mathbb{R})$ and let $q^{ \pm} \in W^{2 m+1, \infty}(\Omega, \mathbb{R}) \cap C^{2(m-1)}(\bar{\Omega}, \mathbb{R})$ satisfy

$$
\|A\|_{W^{2 m+1, \infty}(\Omega)}+\|p\|_{W^{2 m+1, \infty}(\Omega)}+\left\|q^{+}\right\|_{W^{2 m+1, \infty}(\Omega)}+\left\|q^{-}\right\|_{W^{2 m+1, \infty}(\Omega)} \leq M,
$$

for some a priori fixed positive constant $M$. Then, for all $g=\left(g^{+}, g^{-}\right)^{T} \in H^{2(m+7 / 4), m+7 / 4}(\Sigma)^{2}$ and all $u_{0}=$ $\left(u_{0}^{+}, u_{0}^{-}\right)^{T} \in H^{2 m+3}(\Omega)^{2}$ fulfilling the following compatibility conditions

$$
\begin{equation*}
\partial_{t}^{\ell} g(\cdot, 0)=(-i)^{\ell} \mathcal{H}\left(A, p, q^{ \pm}\right)^{\ell} u_{0} \text { on } \Gamma, \ell=0, \cdots, m \tag{1.2}
\end{equation*}
$$

the IBVP (1.1) admits a unique solution $u=\left(u^{+}, u^{-}\right)^{T} \in \cap_{\ell=0}^{m+1} H^{m+1-\ell}\left(0, T ; H^{2 \ell}(\Omega)^{2}\right)$. Moreover, there exists a positive constant $C$, depending only on $\omega, T$ and $M$ such that

$$
\begin{equation*}
\sum_{\ell=0}^{m+1}\|u\|_{H^{m+1-\ell}\left(0, T ; H^{2 \ell}(\Omega)^{2}\right)} \leq C\left(\left\|u_{0}\right\|_{H^{2 m+3}(\Omega)^{2}}+\|g\|_{H^{2(m+7 / 4), m+7 / 4}(\Sigma)^{2}}\right) \tag{1.3}
\end{equation*}
$$

Notice that the divergence-free condition on $A$ requested by Proposition 1.1 is to guarantee that $\mathcal{H}\left(A, p, q^{ \pm}\right)$endowed with homogeneous Dirichlet boundary condition on $\Gamma$, has a self-adjoint realization $H\left(A, p, q^{ \pm}\right)$in $L^{2}(\Omega)^{2}$, see [19, Lemma 2.1]. As a consequence the operator $-i H\left(A, p, q^{ \pm}\right)$is $m$-dissipative in $L^{2}(\Omega)^{2}$, and since the IBVP (1.1) is equivalently rewritten as

$$
\begin{cases}-i \partial_{t} u+\mathcal{H}\left(A, p, q^{ \pm}\right) u=0 & \text { in } Q \\ u(\cdot, 0)=u_{0} & \text { in } \Omega \\ u=g & \text { on } \Sigma\end{cases}
$$

the statement of Proposition 1.1 follows by arguing in the same way as in the proof of [19, Lemma 2.3].
We point out that the regularity assumptions on the coefficients $A, p$ and $q^{ \pm}$, the initial states $u_{0}^{ \pm}$and the boundary conditions $g^{ \pm}$, in Proposition 1.1, are only sufficient conditions ensuring a higher order of regularity of the solution
$u^{ \pm}$to (1.1), as requested by the analysis of the inverse problem under study in this article. As a matter of fact the Bukhgeim-Klibanov method requires $\partial_{t} u^{ \pm}$and $\partial_{t} \nabla u^{ \pm}$to be bounded in $Q$, which can be achieved upon taking $m$ in Proposition 1.1, sufficiently large relative to $n$. Namely, we choose

$$
\begin{equation*}
N \in \mathbb{N} \cap\left(\frac{n+2}{4}+1, \frac{n+2}{4}+2\right], \tag{1.4}
\end{equation*}
$$

pick $M, \kappa, \varrho, \mathfrak{a}, \mathfrak{p}$ and $\mathfrak{q}$ in $\mathbb{R}_{+}$, and for $A_{0} \in W^{2 N+1, \infty}\left(\Omega, \mathbb{R}^{n}\right) \cap C^{2(N-1)}\left(\bar{\Omega}, \mathbb{R}^{n}\right), p_{0} \in W^{2 N+1, \infty}(\Omega, \mathbb{R}) \cap$ $C^{2(N-1)}(\bar{\Omega}, \mathbb{R})$ and $q_{0}^{ \pm} \in W^{2 N+1, \infty}(\Omega, \mathbb{R}) \cap C^{2(N-1)}(\bar{\Omega}, \mathbb{R})$, we introduce the set of unknown electric potentials as

$$
\begin{align*}
\mathcal{P}_{\mathfrak{p}}\left(p_{0}\right):= & \left\{p \in W^{2 N+1, \infty}(\Omega, \mathbb{R}) \cap C^{2(N-1)}(\bar{\Omega}, \mathbb{R}) \text { s.t. }\|p\|_{W^{2 N+1, \infty}(\Omega)} \leq M,\right.  \tag{1.5}\\
& \left.\partial_{x}^{k} p=\partial_{x}^{k} p_{0} \text { on } \Gamma, k=0, \ldots, 2(N-1) \text { and }\left|\left(p-p_{0}\right)\left(\cdot, x_{n}\right)\right| \leq \mathfrak{p} e^{-\kappa\left\langle x_{n}\right\rangle^{e}}, x_{n} \in \mathbb{R}\right\},
\end{align*}
$$

the set of unknown zero-th order coupling coefficients as $\mathcal{P}_{\mathfrak{q}}\left(q_{0}^{ \pm}\right)$, and the set of unknown first order coupling vectors as

$$
\begin{aligned}
(1.6) \mathcal{A}_{\mathfrak{a}}\left(A_{0}\right):= & \left\{A \in W^{2 N+1, \infty}\left(\Omega, \mathbb{R}^{n}\right) \cap C^{2(N-1)}(\bar{\Omega}, \mathbb{R}) \text { s.t. }\|A\|_{W^{2 N+1, \infty}(\Omega)^{n}} \leq M, \nabla \cdot A=0 \text { in } \Omega\right. \\
& \left.\partial_{x}^{k} A=\partial_{x}^{k} A_{0} \text { on } \Gamma,|k|=0, \ldots, 2(N-1) \text { and }\left|\left(A-A_{0}\right)\left(\cdot, x_{n}\right)\right| \leq \mathfrak{a} e^{-\kappa\left\langle x_{n}\right\rangle^{e}}, x_{n} \in \mathbb{R}\right\} .
\end{aligned}
$$

Here, the notation $\partial_{x}^{k}$ for $|k|=m \in \mathbb{N}_{0}$ is a shorthand for $\partial_{1}^{k_{1}} \ldots \partial_{n}^{k_{n}}$ where $k=\left(k_{1}, \ldots, k_{n}\right) \in \mathbb{N}_{0}^{n}$ satisfies $|k|=k_{1}+\ldots+k_{n}=m$.

Then, the main result of this article can be stated as follows.
Theorem 1.2. Assume that $\gamma$ is $\mathcal{C}^{2(N+1)}$. For $j=1,2$, let $A_{j} \in \mathcal{A}_{\mathfrak{a}}\left(A_{0}\right)$ satisfy

$$
\begin{equation*}
\exists y_{*} \in \mathbb{R}_{+}, a_{1, n}\left(x^{\prime}, x_{n}\right)=a_{2, n}\left(x^{\prime}, x_{n}\right), x^{\prime} \in \omega, x_{n} \in\left(-y_{*}, y_{*}\right) \tag{1.7}
\end{equation*}
$$

let $p_{j} \in \mathcal{P}_{\mathfrak{p}}\left(p_{0}\right)$ and let $q_{j}^{ \pm} \in \mathcal{P}_{\mathfrak{q}}\left(q_{0}^{ \pm}\right)$.
Then, there exist a sub-boundary $\gamma_{*} \subset \partial \omega$ and a set of $n+1$ initial states $u_{0}^{k}=\left(u_{0}^{+, k}, u_{0}^{-, k}\right)^{T} \in H^{2 N+3}(\Omega)^{2}$ and boundary conditions $g^{k}=\left(g^{+, k}, g^{-, k}\right)^{T} \in H^{2(N+7 / 4), N+7 / 4}(\Sigma)^{2}, k=1, \ldots, n+1$, fulfilling the compatibility conditions

$$
\begin{equation*}
\partial_{t}^{\ell} g^{k}(\cdot, 0)=(-i)^{\ell} \mathcal{H}\left(A_{0}, p_{0}, q_{0}^{ \pm}\right)^{\ell} u_{0}^{k} \text { on } \Gamma, \ell=0, \cdots, N, \tag{1.8}
\end{equation*}
$$

such that for all $\theta \in\left(0, \frac{1}{2}\right)$, the following estimate

$$
\begin{align*}
& \left\|A_{1}-A_{2}\right\|_{L^{2}(\Omega)}^{2}+\left\|p_{1}-p_{2}\right\|_{L^{2}(\Omega)}^{2}+\left\|q_{1}^{+}-q_{2}^{+}\right\|_{L^{2}(\Omega)}^{2}+\left\|q_{1}^{-}-q_{2}^{-}\right\|_{L^{2}(\Omega)}^{2}  \tag{1.9}\\
\leq & C \sum_{k=1}^{n+1}\left(\left\|\partial_{\nu} \partial_{t} u_{1}^{-, k}-\partial_{\nu} \partial_{t} u_{2}^{-, k}\right\|_{L^{2}\left(\Sigma_{*}\right)}^{\theta}+\left\|\partial_{\nu} \partial_{t} u_{1}^{+, k}-\partial_{\nu} \partial_{t} u_{2}^{+, k}\right\|_{L^{2}\left(\Sigma_{*}\right)}^{\theta}\right),
\end{align*}
$$

holds for some positive constant $C$ depending only on $\omega, T, \gamma^{*}, M, y_{*}, \theta, \kappa, \varrho, \mathfrak{a}, \mathfrak{p}, \mathfrak{q}$ and $\left(u_{0}^{ \pm, k}, g^{ \pm, k}\right), k=$ $1, \ldots, n+1$. Here, $\Sigma_{*}:=\gamma_{*} \times \mathbb{R} \times(0, T)$ and $u_{j}^{k}=\left(u_{j}^{+, k}, u_{j}^{-, k}\right)^{T}$, for $j=1,2$, is the solution to (1.1) given by Proposition 1.1, where $\left(A_{j}, p_{j}, q_{j}^{ \pm}, u_{0}^{ \pm, k}, g^{ \pm, k}\right)$ is substituted for $\left(A, p, q^{ \pm}, u_{0}^{ \pm}, g^{ \pm}\right)$.
1.6. Brief comments. Theorem 1.2 claims that $n+1$ Neumann data stably determine $n+2$ unknown scalar coefficients (strictly speaking there are $n+3$ unknown scalar coefficients in the inverse problem that Theorem 1.2 is dealing with, but since the $n$ components of the gradient coupling vector are bound together through the divergence free condition, they only amount for $n-1$ free unknown scalar coefficients). This may seem surprising from the viewpoint of the analysis of inverse problems, but it should be noticed that Assumption (1.7) implies full knowledge of the $n$-th component of $A$ on a bounded subpart of $\Omega$.

The statement and the strategy of the proof of Theorem 1.2 are very similar to the ones of [19, Theorem 1.2], which holds for a bounded spatial domain $\Omega$. Nevertheless, there are two major differences in the derivation of Theorem 1.2 as compared to the one of [19, Theorem 1.2]. Firstly, the Carleman estimate that is used in Section 3 below is designed for a Schrödinger equation in an unbounded cylindrical domain, and it is slightly different from the one used in [19], which is specific to the Schrödinger equation in a bounded domain. Secondly, the construction of the initial states $u_{0}$
used for probing the system in the analysis of the inverse problem under examination in this article, is more delicate than in [19]. This is due to the fact that it is technically more challenging to design a suitable set of $L^{2}(\Omega)$-initial states $u_{0}$ when the domain $\Omega$ is infinitely extended, than in the case where $\Omega$ is bounded. As can be seen from Section 3 , this $L^{2}$-condition will be satisfied by introducing an additional decay with respect to the infinite direction of the waveguide.
1.7. Outline. The paper is designed as follows: In the following section we collect several technical results needed for the proof of Theorem 1.2, which is given in Section 3.

## 2. Preliminaries

We first establish that the solution to (1.1) is bounded in $Q$.
2.1. Boundedness of the solution. The result we have in mind is as follows.

Lemma 2.1. Assume that conditions of Proposition 1.1 are satisfied with $m=N$, where $N$ is the same as in (1.4). Then, the solution $u$ to (1.1) lies in $W^{1, \infty}\left(0, T ; W^{1, \infty}(\Omega)^{2}\right)$ and satisfies

$$
\|u\|_{W^{1, \infty}\left(0, T ; W^{1, \infty}(\Omega)^{2}\right)} \leq C
$$

for some positive constant $C$ depending only on $\omega, T, M, u_{0}$ and $g$.
Proof. We have $u \in H^{2}\left(0, T, H^{2(N-1)}(\Omega)^{2}\right)$ by Proposition 1.1, with $2(N-1)>\frac{n}{2}+1$ from (1.4). Since $H^{k}(\Omega)$ is continuously embedded in $L^{\infty}(\Omega)$ for all $k>\frac{n}{2}$, according to [21, Lemma 2.7] (which extends the corresponding well-known Sobolev embedding theorem in $\mathbb{R}^{n}$, see e.g. [10, Corollary IX.13] or [15, Section 5.10, Problem 18], to the case of the unbounded cylindrical domain $\Omega$ ), the result follows from this and (1.3).
2.2. Global Carleman estimate for the Schrödinger equation in $\omega \times \mathbb{R}$. For further use we introduce a global Carleman estimate specifically designed for the Schrödinger equation in the unbounded cylindrical domain $\Omega$, which is borrowed from [20, Proposition 3.3 and Lemma 4.2].

For this purpose we pick a function $\alpha \in C^{4}\left(\bar{\omega}, \mathbb{R}_{+}\right)$and an open subset $\gamma_{*} \subset \partial \omega$ satisfying the following conditions:

## Assumption 2.2.

(i) $\exists c \in \mathbb{R}_{+}$s.t. $\left|\nabla^{\prime} \alpha\left(x^{\prime}\right)\right| \geq$ cfor all $x^{\prime} \in \omega$.
(ii) $\forall x^{\prime} \in \gamma \backslash \gamma_{*}, \partial_{\nu} \alpha\left(x^{\prime}\right)=\nabla^{\prime} \alpha\left(x^{\prime}\right) \cdot \nu^{\prime}\left(x^{\prime}\right)<0$.
(iii) $\exists \lambda_{0} \in \mathbb{R}_{+}, \exists c \in \mathbb{R}_{+}$s.t.

$$
\lambda\left|\nabla^{\prime} \alpha\left(x^{\prime}\right) \cdot \zeta\right|^{2}+D^{2} \alpha\left(x^{\prime}, \zeta\right) \geq c|\zeta|^{2}, \zeta \in \mathbb{R}^{n-1}, x^{\prime} \in \omega, \lambda \geq \lambda_{0}
$$

where $D^{2} \alpha\left(x^{\prime}\right):=\left(\partial_{i, j}^{2} \alpha\left(x^{\prime}\right)\right)_{1 \leq i, j \leq n-1}$ and $D^{2} \alpha\left(x^{\prime}, \zeta\right)$ denotes the $\mathbb{R}^{n-1}$-scalar product of $D^{2} \alpha\left(x^{\prime}\right) \zeta$ with $\zeta$.
We point out that there exist $\alpha$ and $\gamma_{*}$ fulfilling the above conditions (i), (ii) and (iii). As a matter of fact, for all $x_{0}^{\prime} \in \mathbb{R}^{n-1} \backslash \bar{\omega}$ fixed, this is the case of the function $\alpha\left(x^{\prime}\right)=\left|x^{\prime}-x_{0}^{\prime}\right|^{2}$ and any open subset $\gamma_{*} \subset \gamma$ such that $\left\{x^{\prime} \in \gamma ;\left(x^{\prime}-x_{0}^{\prime}\right) \cdot \nu\left(x^{\prime}\right) \geq 0\right\} \subset \gamma_{*}$.

Next, putting $K:=r\|\alpha\|_{L^{\infty}(\omega)}$ for some $r \in(1,+\infty)$, we set

$$
\begin{equation*}
\beta(x):=\alpha\left(x^{\prime}\right)+K, x=\left(x^{\prime}, x_{n}\right) \in \Omega \tag{2.10}
\end{equation*}
$$

and we introduce the following weight functions on $\tilde{Q}:=\Omega \times(-T, T)$ :

$$
\begin{equation*}
\varphi(x, t):=\frac{e^{2 \beta(x)}}{(T+t)(T-t)} \text { and } \eta(x, t):=\frac{e^{2 K}-e^{\beta(x)}}{(T+t)(T-t)},(x, t) \in \tilde{Q} \tag{2.11}
\end{equation*}
$$

Let us notice for further use that

$$
\begin{equation*}
\eta(x, t) \geq \eta_{0}(x)>0,(x, t) \in \tilde{Q} \tag{2.12}
\end{equation*}
$$

where $\eta_{0}(x):=\eta(0, x)$ for all $x \in \Omega$. This being said, we may now state the global Carleman estimate established in [20, Proposition 3.3 and Lemma 4.2].

Proposition 2.3. Suppose that $\alpha$ and $\gamma_{*}$ fulfill Assumption 2.2. Let $\beta$ be as in (2.10) and let $\varphi$ and $\eta$ be defined by (2.11). Then, there exist two constants $s_{0}>0$ and $C>0$, depending only on $T, \omega$ and $\gamma_{*}$, such that the estimate

$$
\begin{aligned}
& s^{-1 / 2}\left\|e^{-s \eta} \nabla^{\prime} w\right\|_{L^{2}(\tilde{Q})}^{2}+s^{-1 / 2}\left\|e^{-s \eta} w\right\|_{L^{2}(\tilde{Q})}^{2}+\left\|e^{-s \eta_{0}} w(\cdot, 0)\right\|_{L^{2}(\Omega)}^{2} \\
\leq & C s^{-3 / 2}\left(s\left\|e^{-s \eta} \varphi^{1 / 2}\left|\partial_{\nu} \beta\right|^{\frac{1}{2}} \partial_{\nu} w\right\|_{L^{2}\left(\tilde{\Sigma}_{*}\right)}^{2}+\left\|e^{-s \eta} L w\right\|_{L^{2}(\tilde{Q})}^{2}\right)
\end{aligned}
$$

holds whenever $s \geq s_{0}$ and $w \in L^{2}\left(-T, T ; H_{0}^{1}(\Omega)\right)$ satisfies $L w \in L^{2}(\tilde{Q})$ and $\partial_{\nu} w \in L^{2}\left(\tilde{\Sigma}_{*}\right)$. Here, $\tilde{\Sigma}_{*}:=$ $(-T, T) \times \Gamma_{*}$ and $\Gamma_{*}:=\gamma_{*} \times \mathbb{R}$.

Armed with Proposition 2.3, we turn now to proving the main result of this article.

## 3. Proof of Theorem 1.2

3.1. Linearization, time-differentiation and all that. We start by linearizing the system (1.1). For this purpose we consider the two solutions $u_{j}=\left(u_{j}^{+}, u_{j}^{-}\right)^{T}, j=1,2$, to the $\operatorname{IBVP}(1.1)$ where $\left(A_{j}, p_{j}, q_{j}^{ \pm}\right)$is substituted for $\left(A, p, q^{ \pm}\right)$. Then, $u^{ \pm}:=u_{1}^{ \pm}-u_{2}^{ \pm}$solves

$$
\begin{cases}-i \partial_{t} u^{+}-\Delta u^{+}+q_{1}^{+} u^{+}=-A_{1} \cdot \nabla u^{-}-A \cdot \nabla u_{2}^{-}-q^{+} u_{2}^{+}-p_{1} u^{-}-p u_{2}^{-} & \text {in } Q  \tag{3.13}\\ -i \partial_{t} u^{-}-\Delta u^{-}+q_{1}^{-} u^{-}=A_{1} \cdot \nabla u^{+}+A \cdot \nabla u_{2}^{+}-q^{-} u_{2}^{-}-p_{1} u^{+}-p u_{2}^{+} & \text {in } Q \\ u^{+}(\cdot, 0)=0, u^{-}(\cdot, 0)=0 & \text { in } \Omega \\ u^{+}=0, u^{-}=0 & \text { on } \Sigma,\end{cases}
$$

where $A:=A_{1}-A_{2}, p:=p_{1}-p_{2}$ and $q^{ \pm}:=q_{1}^{ \pm}-q_{2}^{ \pm}$. Further, $u^{ \pm}$lies in $H^{2}\left(0, T ; L^{2}(\Omega)\right) \cap H^{1}\left(0, T ; H^{2}(\Omega) \cap\right.$ $H_{0}^{1}(\Omega)$ ), we differentiate (3.13) with respect to the time-variable and find that

$$
\begin{cases}-i \partial_{t} v^{+}-\Delta v^{+}+q_{1}^{+} v^{+}=-A_{1} \cdot \nabla v^{-}-A \cdot \nabla \partial_{t} u_{2}^{-}-q^{+} \partial_{t} u_{2}^{+}-p_{1} v^{-}-p \partial_{t} u_{2}^{-} & \text {in } Q \\ -i \partial_{t} v^{-}-\Delta v^{-}+q_{1}^{-} v^{-}=A_{1} \cdot \nabla v^{+}+A \cdot \nabla \partial_{t} u_{2}^{+}-q^{-} \partial_{t} u_{2}^{-}-p_{1} v^{+}-p \partial_{t} u_{2}^{+} & \text {in } Q \\ v^{+}(\cdot, 0)=-i\left(A \cdot \nabla u_{0}^{-}+q^{+} u_{0}^{+}+p u_{0}^{-}\right) & \text {in } \Omega \\ v^{-}(\cdot, 0)=-i\left(-A \cdot \nabla u_{0}^{+}+q^{-} u_{0}^{-}+p u_{0}^{+}\right) & \text {in } \Omega \\ v^{+}=0, v^{-}=0 & \text { on } \Sigma,\end{cases}
$$

where $v^{ \pm}:=\partial_{t} u^{ \pm}$. The next step is to extend $u_{2}^{ \pm}$to $\tilde{Q}=\Omega \times(-T, T)$ by setting $u_{2}^{ \pm}(x, t):=\overline{u_{2}^{ \pm}(x,-t)}$ for a.e. $(x, t) \in \Omega \times(-T, 0)$. Since $u_{0}^{ \pm}, A, p$ and $q^{ \pm}$are-real valued, it is not hard to see that the function $v^{ \pm}$, extended to $\Omega \times(-T, 0)$ as $v^{ \pm}(x, t):=-\overline{v^{ \pm}(x,-t)}$, satisfies
(3.14)

$$
\begin{cases}-i \partial_{t} v^{+}-\Delta v^{+}+q_{1}^{+} v^{+}=-A_{1} \cdot \nabla v^{-}-A \cdot \nabla \partial_{t} u_{2}^{-}-q^{+} \partial_{t} u_{2}^{+}-p_{1} v^{-}-p \partial_{t} u_{2}^{-} & \text {in } \tilde{Q} \\ -i \partial_{t} v^{-}-\Delta v^{-}+q_{1}^{-} v^{-}=A_{1} \cdot \nabla v^{+}+A \cdot \nabla \partial_{t} u_{2}^{+}-q^{-} \partial_{t} u_{2}^{-}-p_{1} v^{+}-p \partial_{t} u_{2}^{+} & \text {in } \tilde{Q} \\ v^{+}(\cdot, 0)=-i\left(A \cdot \nabla u_{0}^{-}+q^{+} u_{0}^{+}+p u_{0}^{-}\right) & \text {in } \Omega \\ v^{-}(\cdot, 0)=-i\left(-A \cdot \nabla u_{0}^{+}+q^{-} u_{0}^{-}+p u_{0}^{+}\right) & \text {in } \Omega \\ v^{+}=0, v^{-}=0 & \text { on } \tilde{\Sigma}:=\Gamma \times(-T, T)\end{cases}
$$

Put $\mu^{ \pm}:=\left\|e^{-s \eta_{0}} \varphi^{1 / 2}\left|\partial_{\nu} \beta\right|^{1 / 2} \partial_{\nu} v^{ \pm}\right\|_{L^{2}\left(\tilde{\Sigma}_{*}\right)}^{2}$. Then, applying Proposition 2.3 to (3.14), we get for all $s \geq s_{0}$ that

$$
\begin{align*}
& s^{-1 / 2}\left\|e^{-s \eta} \nabla^{\prime} v^{ \pm}\right\|_{L^{2}(\tilde{Q})}^{2}+s^{-1 / 2}\left\|e^{-s \eta} v^{ \pm}\right\|_{L^{2}(\tilde{Q})}^{2}+\left\|e^{-s \eta_{0}} v^{ \pm}(\cdot, 0)\right\|_{L^{2}(\Omega)}^{2}  \tag{3.15}\\
\leq & C s^{-3 / 2}\left(s \mu^{ \pm}+\left\|e^{-s \eta}\left( \pm A_{1} \cdot \nabla v^{\mp} \pm A \cdot \nabla \partial_{t} u_{2}^{\mp}+q^{ \pm} \partial_{t} u_{2}^{ \pm}+p_{1} v^{\mp}+p \partial_{t} u_{2}^{\mp}\right)\right\|_{L^{2}(\tilde{Q})}^{2}\right),
\end{align*}
$$

for some positive constant $C$ depending only on $\omega, T$ and $\gamma_{*}$. Taking into account that $\left\|A_{1}\right\|_{L^{\infty}(\Omega)} \leq M,\left\|p_{1}\right\|_{L^{\infty}(\Omega)} \leq$ $M$, and that the two functions $\partial_{t} u_{2}^{ \pm}$and $\nabla \partial_{t} u_{2}^{ \pm}$are bounded on $\tilde{Q}$ by some positive constant depending only on $\omega$, $T, M, u_{0}$ and $g$ according to Lemma 2.1, (2.12) and (3.15) then yield that

$$
\begin{aligned}
& s^{-1 / 2}\left\|e^{-s \eta} \nabla^{\prime} v^{ \pm}\right\|_{L^{2}(\tilde{Q})}^{2}+s^{-1 / 2}\left\|e^{-s \eta} v^{ \pm}\right\|_{L^{2}(\tilde{Q})}^{2}+\left\|e^{-s \eta_{0}} v^{ \pm}(\cdot, 0)\right\|_{L^{2}(\Omega)}^{2} \\
\leq & C s^{-3 / 2}\left(s \mu^{ \pm}+\left\|e^{-s \eta} \nabla_{x^{\prime}} v^{\mp}\right\|_{L^{2}(\tilde{Q})}^{2}+\left\|e^{-s \eta} v^{\mp}\right\|_{L^{2}(\tilde{Q})}^{2}+\left\|e^{-s \eta_{0}} A\right\|_{L^{2}(\Omega)^{n}}^{2}+\left\|e^{-s \eta_{0}} q^{ \pm}\right\|_{L^{2}(\Omega)}^{2}+\left\|e^{-s \eta_{0}} p\right\|_{L^{2}(\Omega)}^{2}\right)
\end{aligned}
$$

provided $s \geq s_{0}$. Here and in the remaining part of this proof, $C$ denotes a generic positive constant which may change from line to line. Although the constant $C$ depends only on $\omega, T, \gamma_{*}, M, u_{0}$ and $g$ in the above estimate, in the sequel it might also depend on one or several of the parameters $n, y_{*}, \kappa, \varrho, \mathfrak{a}, \mathfrak{p}, \mathfrak{q}$ and $\theta$ of the problem, as well. Nevertheless, we shall not systematically specify the dependence of $C$ with respect to the above mentioned parameters.

As a consequence we have

$$
\begin{aligned}
& s^{-\frac{1}{2}}\left(1-C s^{-1}\right) \sum_{\ell= \pm}\left(\left\|e^{-s \eta} \nabla^{\prime} v^{\ell}\right\|_{L^{2}(\tilde{Q})}^{2}+\left\|e^{-s \eta} v^{\ell}\right\|_{L^{2}(\tilde{Q})}^{2}\right)+\sum_{\ell= \pm}\left\|e^{-s \eta_{0}} v^{\ell}(\cdot, 0)\right\|_{L^{2}(\Omega)}^{2} \\
\leq & C s^{-\frac{3}{2}}\left(\left\|e^{-s \eta_{0}} A\right\|_{L^{2}(\Omega)^{n}}^{2}+\left\|e^{-s \eta_{0}} p\right\|_{L^{2}(\Omega)}^{2}+\left\|e^{-s \eta_{0}} q^{+}\right\|_{L^{2}(\Omega)}^{2}+\left\|e^{-s \eta_{0}} q^{-}\right\|_{L^{2}(\Omega)}^{2}+s\left(\mu^{+}+\mu^{-}\right)\right),
\end{aligned}
$$

provided $s \geq s_{0}$. Thus, taking $s_{1}:=\max \left(s_{0}, 2 C\right)$ in the above estimate, we infer from (3.14) that

$$
\begin{align*}
& \left\|e^{-s \eta_{0}}\left(q^{+} u_{0}^{+}+A \cdot \nabla u_{0}^{-}+p u_{0}^{-}\right)\right\|_{L^{2}(\Omega)}^{2}+\left\|e^{-s \eta_{0}}\left(q^{-} u_{0}^{-}-A \cdot \nabla u_{0}^{+}+p u_{0}^{+}\right)\right\|_{L^{2}(\Omega)}^{2}  \tag{3.16}\\
\leq & C s^{-\frac{3}{2}}\left(\left\|e^{-s \eta_{0}} A\right\|_{L^{2}(\Omega)^{n}}^{2}+\left\|e^{-s \eta_{0}} p\right\|_{L^{2}(\Omega)}^{2}+\left\|e^{-s \eta_{0}} q^{+}\right\|_{L^{2}(\Omega)}^{2}+\left\|e^{-s \eta_{0}} q^{-}\right\|_{L^{2}(\Omega)}^{2}+s\left(\mu^{+}+\mu^{-}\right)\right),
\end{align*}
$$

whenever $s \geq s_{1}$.
The rest of the proof is to adequately choose $n+1$ initial states $u_{0}^{k}:=\left(u_{0}^{+, k}, u_{0}^{-, k}\right)^{T}, k=1, \ldots, n+1$, in order to estimate each of the four unknown functions $A, p$ and $q^{ \pm}$separately, in terms of the corresponding boundary data $\mu_{k}^{ \pm}:=\left\|e^{-s \eta_{0}} \varphi^{1 / 2}\left|\partial_{\nu} \beta\right|^{1 / 2} \partial_{\nu} v^{ \pm, k}\right\|_{L^{2}\left(\tilde{\Sigma}_{*}\right)}^{2}$, where $v^{ \pm, k}$ is the solution to (3.14) with $u_{0}^{ \pm}=u_{0}^{ \pm, k}$.
3.2. Building $n+1$ suitable initial data. We proceed in two steps.

Step 1: Estimation of $p, q^{ \pm}$and $a_{n}$. We pick $\epsilon \in(0,1)$, put $u_{0}^{+, 1}\left(x^{\prime}, x_{n}\right):=0, u_{0}^{-, 1}\left(x^{\prime}, x_{n}\right):=\left\langle x_{n}\right\rangle^{-\frac{1+\epsilon}{2}}$ for all $\left(x^{\prime}, x_{n}\right) \in \Omega$ and take $u_{0}^{ \pm}=u_{0}^{ \pm, 1}$ in (3.16). For all $s \geq s_{1}$, we get that

$$
\begin{aligned}
& \left\|e^{-s \eta_{0}}\left(2\left\langle x_{n}\right\rangle^{-\frac{1+\epsilon}{2}} p-(1+\epsilon)\left\langle x_{n}\right\rangle^{-\frac{5+\epsilon}{2}} x_{n} a_{n}\right)\right\|_{L^{2}(\Omega)}^{2}+4\left\|e^{-s \eta_{0}}\left\langle x_{n}\right\rangle^{-\frac{1+\epsilon}{2}} q^{-}\right\|_{L^{2}(\Omega)}^{2} \\
\leq & C s^{-\frac{3}{2}}\left(\left\|e^{-s \eta_{0}} A\right\|_{L^{2}(\Omega)^{n}}^{2}+\left\|e^{-s \eta_{0}} p\right\|_{L^{2}(\Omega)}^{2}+\left\|e^{-s \eta_{0}} q^{+}\right\|_{L^{2}(\Omega)}^{2}+\left\|e^{-s \eta_{0}} q^{-}\right\|_{L^{2}(\Omega)}^{2}+s\left(\mu^{+, 1}+\mu^{-, 1}\right)\right)
\end{aligned}
$$

which entails that

$$
\begin{align*}
& \left\|e^{-s \eta_{0}}\left\langle x_{n}\right\rangle^{-\frac{1+\epsilon}{2}} q^{-}\right\|_{L^{2}(\Omega)}^{2}  \tag{3.17}\\
\leq & C s^{-\frac{3}{2}}\left(\left\|e^{-s \eta_{0}} A\right\|_{L^{2}(\Omega)^{n}}^{2}+\left\|e^{-s \eta_{0}} p\right\|_{L^{2}(\Omega)}^{2}+\left\|e^{-s \eta_{0}} q^{+}\right\|_{L^{2}(\Omega)}^{2}+\left\|e^{-s \eta_{0}} q^{-}\right\|_{L^{2}(\Omega)}^{2}+s\left(\mu^{+, 1}+\mu^{-, 1}\right)\right)
\end{align*}
$$

and

$$
\begin{align*}
\text { 8) } & \left\|e^{-s \eta_{0}}\left(2\left\langle x_{n}\right\rangle^{-\frac{1+\epsilon}{2}} p-(1+\epsilon)\left\langle x_{n}\right\rangle^{-\frac{5+\epsilon}{2}} x_{n} a_{n}\right)\right\|_{L^{2}(\Omega)}^{2}  \tag{3.18}\\
\leq & C s^{-\frac{3}{2}}\left(\left\|e^{-s \eta_{0}} A\right\|_{L^{2}(\Omega)^{n}}^{2}+\left\|e^{-s \eta_{0}} p\right\|_{L^{2}(\Omega)}^{2}+\left\|e^{-s \eta_{0}} q^{+}\right\|_{L^{2}(\Omega)}^{2}+\left\|e^{-s \eta_{0}} q^{-}\right\|_{L^{2}(\Omega)}^{2}+s\left(\mu^{+, 1}+\mu^{-, 1}\right)\right) .
\end{align*}
$$

Doing the same with $u_{0}^{ \pm}=u_{0}^{ \pm, 2}:=u_{0}^{\mp, 1}$, we obtain for all $s \geq s_{1}$ that

$$
\begin{align*}
& \left\|e^{-s \eta_{0}}\left\langle x_{n}\right\rangle^{-\frac{1+\epsilon}{2}} q^{+}\right\|_{L^{2}(\Omega)}^{2}  \tag{3.19}\\
\leq & C s^{-\frac{3}{2}}\left(\left\|e^{-s \eta_{0}} A\right\|_{L^{2}(\Omega)^{n}}^{2}+\left\|e^{-s \eta_{0}} p\right\|_{L^{2}(\Omega)}^{2}+\left\|e^{-s \eta_{0}} q^{+}\right\|_{L^{2}(\Omega)}^{2}+\left\|e^{-s \eta_{0}} q^{-}\right\|_{L^{2}(\Omega)}^{2}+s\left(\mu^{+, 2}+\mu^{-, 2}\right)\right),
\end{align*}
$$

and

$$
\begin{align*}
& \left\|e^{-s \eta_{0}}\left(2\left\langle x_{n}\right\rangle^{-\frac{1+\epsilon}{2}} p+(1+\epsilon)\left\langle x_{n}\right\rangle^{-\frac{5+\epsilon}{2}} x_{n} a_{n}\right)\right\|_{L^{2}(\Omega)}^{2}  \tag{3.20}\\
\leq & C s^{-\frac{3}{2}}\left(\left\|e^{-s \eta_{0}} A\right\|_{L^{2}(\Omega)^{n}}^{2}+\left\|e^{-s \eta_{0}} p\right\|_{L^{2}(\Omega)}^{2}+\left\|e^{-s \eta_{0}} q^{+}\right\|_{L^{2}(\Omega)}^{2}+\left\|e^{-s \eta_{0}} q^{-}\right\|_{L^{2}(\Omega)}^{2}+s\left(\mu^{+, 2}+\mu^{-, 2}\right)\right)
\end{align*}
$$

Since $8\left\|e^{-s \eta_{0}}\left\langle x_{n}\right\rangle^{-\frac{1+\epsilon}{2}} p\right\|_{L^{2}(\Omega)}^{2}$ is upper-bounded by the sum of $\left\|e^{-s \eta_{0}}\left(2\left\langle x_{n}\right\rangle^{-\frac{1+\epsilon}{2}} p+(1+\epsilon)\left\langle x_{n}\right\rangle^{-\frac{5+\epsilon}{2}} x_{n} a_{n}\right)\right\|_{L^{2}(\Omega)}^{2}$ and $\left\|e^{-s \eta_{0}}\left(2\left\langle x_{n}\right\rangle^{-\frac{1+\epsilon}{2}} p-(1+\epsilon)\left\langle x_{n}\right\rangle^{-\frac{5+\epsilon}{2}} x_{n} a_{n}\right)\right\|_{L^{2}(\Omega)}^{2}$, it follows from (3.18) and (3.20) that

$$
\begin{equation*}
\left\|e^{-s \eta_{0}}\left\langle x_{n}\right\rangle^{-\frac{1+\epsilon}{2}} p\right\|_{L^{2}(\Omega)}^{2} \tag{3.21}
\end{equation*}
$$

$$
\leq C s^{-\frac{3}{2}}\left(\left\|e^{-s \eta_{0}} A\right\|_{L^{2}(\Omega)^{n}}^{2}+\left\|e^{-s \eta_{0}} p\right\|_{L^{2}(\Omega)}^{2}+\left\|e^{-s \eta_{0}} q^{+}\right\|_{L^{2}(\Omega)}^{2}+\left\|e^{-s \eta_{0}} q^{-}\right\|_{L^{2}(\Omega)}^{2}+s \sum_{i=1}^{2}\left(\mu^{+, i}+\mu^{-, i}\right)\right)
$$

whenever $s \geq s_{1}$. Similarly, upon estimating $\left\|e^{-s \eta_{0}}\left(2\left\langle x_{n}\right\rangle^{-\frac{1+\epsilon}{2}} p+(1+\epsilon)\left\langle x_{n}\right\rangle^{-\frac{5+\epsilon}{2}} x_{n} a_{n}\right)\right\|_{L^{2}(\Omega)}^{2}$ from below by the difference $\frac{(1+\epsilon)^{2}}{2}\left\|e^{-s \eta_{0}}\left\langle x_{n}\right\rangle^{-\frac{5+\epsilon}{2}} x_{n} a_{n}\right\|_{L^{2}(\Omega)}^{2}-4\left\|e^{-s \eta_{0}}\left\langle x_{n}\right\rangle^{-\frac{1+\epsilon}{2}} p\right\|_{L^{2}(\Omega)}^{2}$, we get from (3.20)-(3.21) that
(3.22) $\left\|e^{-s \eta_{0}}\left\langle x_{n}\right\rangle^{-\frac{5+\epsilon}{2}} x_{n} a_{n}\right\|_{L^{2}(\Omega)}^{2}$

$$
\leq C s^{-\frac{3}{2}}\left(\left\|e^{-s \eta_{0}} A\right\|_{L^{2}(\Omega)^{n}}^{2}+\left\|e^{-s \eta_{0}} p\right\|_{L^{2}(\Omega)}^{2}+\left\|e^{-s \eta_{0}} q^{+}\right\|_{L^{2}(\Omega)}^{2}+\left\|e^{-s \eta_{0}} q^{-}\right\|_{L^{2}(\Omega)}^{2}+s \sum_{i=1}^{2}\left(\mu^{+, i}+\mu^{-, i}\right)\right)
$$

for all $s \geq s_{1}$. Bearing in mind that $\left|x_{n} a_{n}\right| \geq y_{*}\left|a_{n}\right|$ in $\Omega$, by virtue of the assumption (1.7), it follows from (3.22) that

$$
\begin{align*}
& \text { 23) }\left\|e^{-s \eta_{0}}\left\langle x_{n}\right\rangle^{-\frac{5+\epsilon}{2}} a_{n}\right\|_{L^{2}(\Omega)}^{2}  \tag{3.23}\\
& \leq C s^{-\frac{3}{2}}\left(\left\|e^{-s \eta_{0}} A\right\|_{L^{2}(\Omega)^{n}}^{2}+\left\|e^{-s \eta_{0}} p\right\|_{L^{2}(\Omega)}^{2}+\left\|e^{-s \eta_{0}} q^{+}\right\|_{L^{2}(\Omega)}^{2}+\left\|e^{-s \eta_{0}} q^{-}\right\|_{L^{2}(\Omega)}^{2}+s \sum_{i=1}^{2}\left(\mu^{+, i}+\mu^{-, i}\right)\right)
\end{align*}
$$

provided we have $s \geq s_{1}$.
Step 2: Estimation of the $n-1$ first components $a_{j}, j=1, \ldots, n-1$, of $A$. For all $k=1, \cdots, n-1$ and all $x=$ $\left(x_{1}, \ldots, x_{n}\right) \in \Omega$, we put $u_{0}^{ \pm, k+2}(x):=x_{k}\left\langle x_{n}\right\rangle^{-\frac{1+\epsilon}{2}}$, substitute $u_{0}^{ \pm, k+2}$ for $u_{0}^{ \pm}$in (1.1) and then apply Proposition 2.3 to (3.14). We get for all $s \geq s_{1}$ that

$$
\begin{aligned}
& \left\|e^{-s \eta_{0}}\left(p u_{0}^{-, k+2}+A \cdot \nabla u_{0}^{-, k+2}+q^{+} u_{0}^{+, k+2}\right)\right\|_{L^{2}(\Omega)}^{2}+\left\|e^{-s \eta_{0}}\left(p u_{0}^{+, k+2}-A \cdot \nabla u_{0}^{+, k+2}+q^{-} u_{0}^{-, k+2}\right)\right\|_{L^{2}(\Omega)}^{2} \\
\leq & C s^{-3 / 2}\left(\left\|e^{-s \eta_{0}} A\right\|_{L^{2}(\Omega)^{n}}^{2}+\left\|e^{-s \eta_{0}} p\right\|_{L^{2}(\Omega)}^{2}+\left\|e^{-s \eta_{0}} q^{+}\right\|_{L^{2}(\Omega)}^{2}+\left\|e^{-s \eta_{0}} q^{-}\right\|_{L^{2}(\Omega)}^{2}+s\left(\mu^{+, k+2}+\mu^{-, k+2}\right)\right) .
\end{aligned}
$$

Since $\left|p u_{0}^{\mp, k+2} \pm A \cdot \nabla u_{0}^{\mp, k+2}+q^{ \pm} u_{0}^{ \pm, k+2}\right|^{2} \geq \frac{\left|A \cdot \nabla u_{0}^{\mp, k+2}\right|^{2}}{2}-\left|p u_{0}^{\mp, k+2}+q^{ \pm} u_{0}^{ \pm, k+2}\right|^{2}$, this entails that (3.24) $\left\|e^{-s \eta_{0}} A \cdot \nabla u_{0}^{+, k+2}\right\|_{L^{2}(\Omega)}^{2}+\left\|e^{-s \eta_{0}} A \cdot \nabla u_{0}^{-, k+2}\right\|_{L^{2}(\Omega)}^{2}$

$$
\begin{aligned}
\leq & C s^{-3 / 2}\left(\left\|e^{-s \eta_{0}} A\right\|_{L^{2}(\Omega)^{n}}^{2}+\left\|e^{-s \eta_{0}} p\right\|_{L^{2}(\Omega)}^{2}+\left\|e^{-s \eta_{0}} q^{+}\right\|_{L^{2}(\Omega)}^{2}+\left\|e^{-s \eta_{0}} q^{-}\right\|_{L^{2}(\Omega)}^{2}+s\left(\mu^{+, k+2}+\mu^{-, k+2}\right)\right) \\
& +\left\|e^{-s \eta_{0}}\left(p u_{0}^{+, k+2}+q^{-} u_{0}^{-, k+2}\right)\right\|_{L^{2}(\Omega)}^{2}+\left\|e^{-s \eta_{0}}\left(p u_{0}^{-, k+2}+q^{+} u_{0}^{+, k+2}\right)\right\|_{L^{2}(\Omega)}^{2}
\end{aligned}
$$

Moreover, $\left\|e^{-s \eta_{0}}\left(p u_{0}^{ \pm, k+2}+q^{\mp} u_{0}^{\mp, k+2}\right)\right\|_{L^{2}(\Omega)}^{2}=\left\|e^{-s \eta_{0}} x_{k}\left\langle x_{n}\right\rangle^{-\frac{1+\epsilon}{2}}\left(p+q^{\mp}\right)\right\|_{L^{2}(\Omega)}^{2}$ being upper-bounded by $2|\omega|^{2}\left(\left\|e^{-s \eta_{0}}\left\langle x_{n}\right\rangle^{-\frac{1+\epsilon}{2}} p\right\|_{L^{2}(\Omega)}^{2}+\left\|e^{-s \eta_{0}}\left\langle x_{n}\right\rangle^{-\frac{1+\epsilon}{2}} q^{\mp}\right\|_{L^{2}(\Omega)}^{2}\right)$, (3.17), (3.19), (3.21) and (3.24) then yield

$$
\begin{aligned}
& \left\|e^{-s \eta_{0}} A \cdot \nabla u_{0}^{+, k+2}\right\|_{L^{2}(\Omega)}^{2}+\left\|e^{-s \eta_{0}} A \cdot \nabla u_{0}^{-, k+2}\right\|_{L^{2}(\Omega)}^{2} \\
\leq & C s^{-3 / 2}\left(\left\|e^{-s \eta_{0}} A\right\|_{L^{2}(\Omega)^{n}}^{2}+\left\|e^{-s \eta_{0}} p\right\|_{L^{2}(\Omega)}^{2}+\left\|e^{-s \eta_{0}} q^{+}\right\|_{L^{2}(\Omega)}^{2}+\left\|e^{-s \eta_{0}} q^{-}\right\|_{L^{2}(\Omega)}^{2}\right. \\
& \left.+s\left(\sum_{i=1}^{2}\left(\mu^{+, i}+\mu^{-, i}\right)+\mu^{+, k+2}+\mu^{-, k+2}\right)\right), s \geq s_{1},
\end{aligned}
$$

From this, (3.22) and the estimates $\left|A . \nabla u_{0}^{ \pm, k+2}\right|^{2} \geq \frac{1}{2}\left|\left\langle x_{n}\right\rangle^{-\frac{1+\epsilon}{2}} a_{k}\right|^{2}-\frac{(1+\epsilon)^{2}}{4}\left|\left\langle x_{n}\right\rangle^{-\frac{5+\epsilon}{2}} x_{k} x_{n} a_{n}\right|^{2}$ and $\left\|e^{-s \eta_{0}}\left\langle x_{n}\right\rangle^{-\frac{5+\epsilon}{2}} x_{k} x_{n} a_{n}\right\|_{L^{2}(\Omega)} \leq|\omega|\left\|e^{-s \eta_{0}}\left\langle x_{n}\right\rangle^{-\frac{5+\epsilon}{2}} x_{n} a_{n}\right\|_{L^{2}(\Omega)}$, it then follows that

$$
\begin{aligned}
& \left\|e^{-s \eta_{0}}\left\langle x_{n}\right\rangle^{-\frac{1+\epsilon}{2}} a_{k}\right\|_{L^{2}(\Omega)}^{2} \\
\leq & C s^{-3 / 2}\left(\left\|e^{-s \eta_{0}} A\right\|_{L^{2}(\Omega)^{n}}^{2}+\left\|e^{-s \eta_{0}} p\right\|_{L^{2}(\Omega)}^{2}+\left\|e^{-s \eta_{0}} q^{+}\right\|_{L^{2}(\Omega)}^{2}+\left\|e^{-s \eta_{0}} q^{-}\right\|_{L^{2}(\Omega)}^{2}\right. \\
& \left.+s\left(\sum_{i=1}^{2}\left(\mu^{+, i}+\mu^{-, i}\right)+\mu^{+, k+2}+\mu^{-, k+2}\right)\right), s \geq s_{1}
\end{aligned}
$$

Summing up the above inequality over $k=1, \ldots, n-1$ and remembering (3.23), we obtain

$$
\begin{align*}
& \left\|e^{-s \eta_{0}}\left\langle x_{n}\right\rangle^{-\frac{5+\epsilon}{2}} A\right\|_{L^{2}(\Omega)^{n}}^{2}  \tag{3.25}\\
\leq & C s^{-3 / 2}\left(\left\|e^{-s \eta_{0}} A\right\|_{L^{2}(\Omega)^{n}}^{2}+\left\|e^{-s \eta_{0}} p\right\|_{L^{2}(\Omega)}^{2}+\left\|e^{-s \eta_{0}} q^{+}\right\|_{L^{2}(\Omega)}^{2}+\left\|e^{-s \eta_{0}} q^{-}\right\|_{L^{2}(\Omega)}^{2}+s \xi\right),
\end{align*}
$$

for $s \geq s_{1}$, where $\xi:=\sum_{i=1}^{n+1}\left(\mu^{+, i}+\mu^{-, i}\right)$.
3.3. End of the proof. For all $y>0$ we have

$$
\text { 6) } \begin{align*}
& \left(\langle y\rangle^{-(5+\epsilon)}-C s^{-\frac{3}{2}}\right)\left(\left\|e^{-s \eta_{0}} A\right\|_{L^{2}(\Omega)^{n}}^{2}+\left\|e^{-s \eta_{0}} p\right\|_{L^{2}\left(\Omega_{y}\right)}^{2}+\left\|e^{-s \eta_{0}} q^{+}\right\|_{L^{2}\left(\Omega_{y}\right)}^{2}+\left\|e^{-s \eta_{0}} q^{-}\right\|_{L^{2}\left(\Omega_{y}\right)}^{2}\right)  \tag{3.26}\\
\leq & C s^{-\frac{3}{2}}\left(\left\|e^{-s \eta_{0}} A\right\|_{L^{2}(\Omega \backslash \Omega)^{n}}^{2}+\left\|e^{-s \eta_{0}} p\right\|_{L^{2}\left(\Omega \backslash \Omega_{y}\right)}^{2}+\left\|e^{-s \eta_{0}} q^{+}\right\|_{L^{2}\left(\Omega \backslash \Omega_{y}\right)}^{2}+\left\|e^{-s \eta_{0}} q^{-}\right\|_{L^{2}\left(\Omega \backslash \Omega_{y}\right)}^{2}+s \xi\right), \\
\leq & C s^{-\frac{3}{2}}\left(\|A\|_{L^{2}(\Omega \backslash \Omega)^{n}}^{2}+\|p\|_{L^{2}\left(\Omega \backslash \Omega_{y}\right)}^{2}+\left\|q^{+}\right\|_{L^{2}\left(\Omega \backslash \Omega_{y}\right)}^{2}+\left\|q^{-}\right\|_{L^{2}\left(\Omega \backslash \Omega_{y}\right)}^{2}+s \xi\right), s \geq s_{1},
\end{align*}
$$

by (3.17), (3.19), (3.21) and (3.25), where $\Omega_{y}:=\omega \times(-y, y)$. Notice that in the last line of (3.26), we used that $\eta_{0}$ is non-negative in $\Omega$. Moreover, for all $y \geq y_{1}:=\left((2 C)^{-\frac{2}{3}} s_{1}\right)^{\frac{3}{2(5+\epsilon)}}$ we have $s_{y}:=(2 C)^{\frac{2}{3}}\langle y\rangle^{\frac{2(5+\epsilon)}{3}} \geq s_{1}$ and $2 C s_{y}^{-\frac{3}{2}} \leq\langle y\rangle^{-(5+\epsilon)}$. Therefore, applying (3.26) with $s=s_{y}$ and using that $\eta_{0}(x) \leq \frac{e^{2 K}}{T^{2}}$ for all $x \in \Omega$, we obtain that

$$
\begin{equation*}
\Theta_{\Omega_{y}} \leq C\left(\Theta_{\Omega \backslash \Omega_{y}}+\langle y\rangle^{\frac{2(5+\epsilon)}{3}} \xi\right), y \geq y_{1} \tag{3.27}
\end{equation*}
$$

where we set $\Theta_{X}:=\|A\|_{0, X}^{2}+\|p\|_{0, X}^{2}+\left\|q^{+}\right\|_{0, X}^{2}+\left\|q^{-}\right\|_{0, X}^{2}$ for any subset $X \subset \Omega$. Next, using that $p_{j} \in \mathcal{P}_{\mathfrak{p}}\left(p_{0}\right)$ for $j=1,2$, we infer from (1.5) upon writing $\|p\|_{L^{2}\left(\Omega \backslash \Omega_{y}\right)} \leq \sum_{j=1,2}\left\|p_{j}-p_{0}\right\|_{L^{2}\left(\Omega \backslash \Omega_{y}\right)}$, that

$$
\begin{align*}
\|p\|_{L^{2}\left(\Omega \backslash \Omega_{y}\right)}^{2} & \leq 4 \mathfrak{p}^{2} \int_{\Omega \backslash \Omega_{y}} e^{-2 \kappa\left\langle x_{n}\right\rangle^{\varrho}} d x^{\prime} d x_{n}  \tag{3.28}\\
& \leq 4 \mathfrak{p}^{2}|\omega| \int_{\left|x_{n}\right|>y} e^{-2 \kappa\left\langle x_{n}\right\rangle^{\varrho}} d x_{n} \\
& \leq 4 \mathfrak{p}^{2}|\omega|\left(\int_{\mathbb{R}} e^{-\delta\left\langle x_{n}\right\rangle^{\varrho}} d x_{n}\right) e^{-(2 \kappa-\delta)\langle y\rangle^{\varrho}}, \delta \in(0,2 \varrho) .
\end{align*}
$$

Similarly, since $q_{j}^{ \pm} \in \mathcal{P}_{\mathfrak{q}}\left(q_{0}^{ \pm}\right)$and $A_{j} \in \mathcal{A}_{\mathfrak{a}}\left(A_{0}\right)$ for $j=1$, 2, we obtain

$$
\begin{equation*}
\Theta_{\Omega \backslash \Omega_{y}} \leq C e^{-(2 \kappa-\delta)\langle y\rangle^{\varrho}}, \delta \in(0,2 \varrho) \tag{3.29}
\end{equation*}
$$

from (1.6) and (3.28), where $C=4|\omega|\left(\mathfrak{a}^{2}+\mathfrak{p}^{2}+2 \mathfrak{q}^{2}\right) \int_{\mathbb{R}} e^{-\delta\left\langle x_{n}\right\rangle^{\varrho}} d x_{n}$. It follows from this and (3.27) that

$$
\begin{equation*}
\Theta_{\Omega_{y}} \leq C\left(e^{-(2 \kappa-\delta)\langle y\rangle^{\varrho}}+\langle y\rangle^{\frac{2(5+\epsilon)}{3}} \xi\right), y \geq y_{1}, \delta \in(0,2 \varrho) \tag{3.30}
\end{equation*}
$$

Put $\xi_{1}:=e^{-(2 \kappa-\delta)\left\langle y_{1}\right\rangle^{\varrho}}$. We shall examine the two cases $\xi \in\left(0, \xi_{1}\right]$ and $\xi \in\left(\xi_{1},+\infty\right)$ separately. Let us start with $\xi \in\left(0, \xi_{1}\right]$. In this case, we pick $y \in\left[y_{1},+\infty\right)$ so large that $e^{-(2 \kappa-\delta)\langle y\rangle^{e}}=\xi$, i.e., $y=\left(\left(-\frac{\ln \xi}{2 \kappa-\delta}\right)^{\frac{2}{e}}-1\right)^{\frac{1}{2}}$. Thus, with reference to (3.29)-(3.30) we get for all $\xi \in\left(0, \xi_{1}\right]$ that $\Theta_{\Omega \backslash \Omega_{y}} \leq C \xi_{1}^{1-2 \theta} \xi^{2 \theta}$ and that $\Theta_{\Omega_{y}} \leq C\left(\xi_{1}^{1-2 \theta}+C_{1}(\theta)\right) \xi^{2 \theta}$, where $C_{1}(\theta):=\sup _{\xi \in\left(0, \xi_{1}\right]}\left(\xi^{1-2 \theta}\left(\frac{-\ln \xi}{2 \kappa-\delta}\right)^{\frac{2(5+\epsilon)}{3 \varrho}}\right)<\infty$ from the assumption $\varrho>0$. As a consequence we have

$$
\begin{equation*}
\Theta_{\Omega} \leq C\left(2 \xi_{1}^{1-2 \theta}+C_{1}(\theta)\right) \xi^{2 \theta}, \xi \in\left(0, \xi_{1}\right] \tag{3.31}
\end{equation*}
$$

and the desired result follows. Now, when $\xi \in\left(\xi_{1},+\infty\right)$, we infer from (1.5) upon majorizing $\|p\|_{L^{2}(\Omega)}^{2}$ by $2 \sum_{j=1,2}\left\|p_{j}-p_{0}\right\|_{L^{2}(\Omega)}^{2}$, that $\|p\|_{L^{2}(\Omega)}^{2} \leq 4 \mathfrak{p}^{2}|\omega|\left(\int_{\mathbb{R}} e^{-2 \kappa\left\langle x_{n}\right\rangle^{\varrho}} d x_{n}\right) \xi_{1}^{-2 \theta} \xi^{2 \theta}$. Doing the same with $q^{ \pm}$and $A$, with the aid of, respectively, (1.5) and (1.6), we find that $\Theta_{\Omega} \leq \tilde{C}_{1}(\theta) \xi^{2 \theta}$, where the notation $\tilde{C}_{1}(\theta)$ stands for the constant $4\left(\mathfrak{a}^{2}+\mathfrak{p}^{2}+2 \mathfrak{q}^{2}\right)|\omega|\left(\int_{\mathbb{R}} e^{-2 \kappa\left\langle x_{n}\right\rangle^{e}} d x_{n}\right) \xi_{1}^{-2 \theta}$. This, (3.31) and the estimates $\mu_{k}^{ \pm} \leq C\left\|\partial_{\nu} v^{ \pm, k}\right\|_{L^{2}\left(\tilde{\Sigma}_{*}\right)}^{2}$ for all $k=1, \ldots, n+1$, yield (1.9), which completes the proof of Theorem 1.2.

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