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# A stability result for the diffusion coefficient of the heat operator defined on an unbounded guide 

Laure Cardoulis * Michel Cristofol* Morgan Morancey*


#### Abstract

In this article we consider the inverse problem of determining the diffusion coefficient of the heat operator in an unbounded guide using a finite number of localized observations. For this problem, we prove a stability estimate in any finite portion of the guide using an adapted Carleman inequality. The measurements are located on the boundary of a larger finite portion of the guide. A special care is required to avoid measurements on the cross-section boundaries which are inside the actual guide. This stability estimate uses a technical positivity assumption. Using arguments from control theory, we manage to remove this assumption for the inverse problem with a given non homogeneous boundary condition.


## 1 Introduction and main results

### 1.1 Introduction

Inverse problems associated with the heat operator have been frequently investigated both on theoretical aspects and for their applications. Such applications cover numerous domains such as medicine, ecology, biology, evolution of populations, physics ...
For these inverse problems there exist several methods involving different type of observations. In the case of a finite number of observations we can consider for example boundary data, observation on all the domain at one time, spectral data or pointwise observation. On the other hand, nice results have been obtained using an infinite number of measurements (e.g. methods using Dirichlet to Neumann map) but we don't consider such approaches here.
In this paper, we focus on the method based on Carleman estimates which allows to derive stability inequality of the following form:

$$
\|\left(\text { coefficient }_{1}\right)-\left(\text { coefficient }_{2}\right) \| \leq f\left(\| \text { observation }_{1}-\text { observation }_{2} \|\right)
$$

for a function $f$ satisfying $\lim _{s \rightarrow 0} f(s)=0$. This kind of inequality links the distance between two sets of coefficients to be reconstructed with the distance between two sets of observations in appropriate norms. Such stability inequalities lead to the uniqueness of the coefficient to be reconstructed. They are also useful to improve the numerical reconstruction of the coefficients using noise-free observations [12].

The strategy based on Carleman estimates has been initiated in [3]. Since then, there is a huge literature on the determination of nonlinear spatially homogeneous terms or source terms in reaction-difusion equations from boundary measurements, see for instance $[5,9,10,13,14,18$,

[^0]$20,22,28]$ or [27] for a survey on this topic. These works provide Lipschitz stability, in addition to the uniqueness of the coefficients. However, this method requires, among other measurements, the knowledge of the solution at some time $\theta>0$ and for all $x$ in the domain. Using a different approach, uniqueness results in the one dimensional case have been established by only one point measurement in space for different parabolic problems with strong nonlinearity [6, 24] or strongly coupled parabolic systems [25] or models with memory term [23] or simultaneous recovery of high order coefficients [11].

All the previous works concern bounded domain $\Omega$ and the case of unbounded domains for parabolic operators is less addressed (see for instance [4, 7]). Unbounded domain problems concentrate several difficulties: first, it is mandatory to rewrite or adapt carefully the existing Carleman estimates. In this paper we adapt ideas from [8] and [4] to write a new Carleman estimate for the parabolic operator in divergence form defined on an infinite guide. Second, this method involves some positivity assumption on the gradient of a solution of the problem. In this paper, we develop a control approach in Sobolev norms for parabolic operators in unbounded domains.

The model of an unbounded guide is relevant for numerous applications. We can consider for example the case of the detection of cracks or inhomogeneities for a long tube in underground or underwater or dangerous areas. On the other hand, as soon as the length of the studied domain is sufficiently large, the mathematical modelisation by an infinite domain is frequently considered to be relevant by engineers.

Our manuscript is organized as follows. In Section 2, we prove an adapted global Carleman estimate for our problem. In section 3 we prove our main result and in section 4, we carry out an adapted control in order to eliminate the positivity condition imposed by the inverse problem.

### 1.2 Settings and hypotheses

Let $\omega$ be a bounded domain in $\mathbb{R}^{n-1}, n \geq 2$ with $C^{2}$ boundary. Denote $\Omega=\mathbb{R} \times \omega$ and $Q=\Omega \times(0, T), \Sigma=\partial \Omega \times(0, T)$. We consider the following problem

$$
\left\{\begin{array}{l}
\partial_{t} u-\nabla \cdot(c \nabla u)=0 \text { in } Q  \tag{1.1}\\
u=0 \text { on } \Sigma \\
u(x, 0)=u_{0}(x) \text { in } \Omega
\end{array}\right.
$$

where $u_{0}$ is a sufficiently smooth function and $c$ is a bounded coefficient defined in $\Omega$ such that $c>0$ and $c \in C^{1}(\bar{\Omega})$. Our problem can be stated as follows.
Let $l>0$ and denote $\Omega_{*}=(-*, *) \times \omega$. We want to determine the coefficient $c$ in $\Omega_{l}$ from a finite number of measurements of the solution $u$ of the system (1.1) on a lateral subset of $\partial \Omega_{L}$ for some $L>l$ and from the knowledge of the solution at the time $\frac{T}{2}$. We stress out that the required measurements are not performed on all the boundary $\partial \Omega_{L}$ and that they avoid the cross-sections $\{ \pm L\} \times \omega$.

Our proof requires the strong positivity condition (1.16). This condition, involved in almost all inverse problems dedicated to coefficients of the principal part of operators, is removed by the construction of an adapted control. The main difficulty is to prove a controllability result in $H^{1}$ norm on the unbounded domain $\Omega$. The strategy is inspired by [ 1, Section 4$]$ where the authors deal with a system of coupled equations but on a bounded domain.

Before stating our results let us mention the work [16] where similar concerns and technics are at stake. In this paper, the authors prove null controllability of parabolic equations in unbounded domains by means of a Carleman estimate (with a distributed observability). The main advantage of their result is that they directly deal with unbounded domains. The price to pay is that the observation region (or control region in their setting) has to be unbounded which is not suitable for our purpose.

We denote by $Q_{*}$ the set $Q_{*}=\Omega_{*} \times(0, T)=(-*, *) \times \omega \times(0, T)$. For each $x=\left(x_{1}, \ldots, x_{n}\right) \in$ $\mathbb{R}^{n}$, let $x^{\prime}=\left(x_{2}, \ldots, x_{n}\right) \in \mathbb{R}^{n-1}$.

The cornerstone of Carleman estimates is to carry out special weight functions. In our inverse problem, the design of the weight function will allow us to eliminate observations on the cross-section of the wave guide. We follow some ideas from [8] mainly to eliminate the cross-section observations. Nevertheless, we are going to detail the different steps which allow to get the two main inequalities (1.4). For this, we choose $a \in \mathbb{R}^{n} \backslash \Omega_{L}$ such that, if

$$
d(x)=\left|x^{\prime}-a^{\prime}\right|^{2}-x_{1}^{2} \text { for } x \in \Omega_{L}
$$

then

$$
\begin{equation*}
d>0 \text { in } \Omega_{L}, \quad|\nabla d|>0 \text { in } \overline{\Omega_{L}} . \tag{1.2}
\end{equation*}
$$

From now on and for simplicity, we denote $\theta=\frac{T}{2}$. From [27, 28] we consider for $t \in(0, T)$ the following classical form of weight function

$$
\begin{equation*}
\psi(x, t)=d(x)-(t-\theta)^{2}+M_{1}, \quad \text { where } \quad M_{1}>\sup _{0<t<T}(t-\theta)^{2}=(T / 2)^{2} \tag{1.3}
\end{equation*}
$$

and

$$
\phi(x, t)=e^{\lambda \psi(x, t)}
$$

The constant $\lambda>0$ will be set in Proposition 2.1. In view of the Carleman inequalities for the parabolic operators with regular weights, we need to use cut-off functions in time. On the other hand, to manage our infinite wave guide we will need also to consider cut-off functions in space but only in the infinite direction $x_{1}$. These cut-off functions will induce additive terms coming from the commutator between the evolution operator and these cut-off functions. These residual terms will be estimated thanks to the following crucial properties of the weight function.

Proposition 1.1. There exist $T>0, L>l, a \in \mathbb{R}^{n} \backslash \Omega_{L}$ and $\epsilon>0$ such that (1.2) holds and, setting

$$
O_{L, \epsilon}=\left\{\Omega_{L} \times((0,2 \epsilon) \cup(T-2 \epsilon, T))\right\} \cup\{((-L,-L+2 \epsilon) \cup(L-2 \epsilon, L)) \times \omega \times(0, T)\},
$$

we have

$$
\begin{equation*}
d_{1}<d_{0}<d_{2} \tag{1.4}
\end{equation*}
$$

where

$$
d_{0}=\inf _{\Omega_{l}} \phi(\cdot, \theta), \quad d_{1}=\sup _{O_{L, \epsilon}} \phi, \quad d_{2}=\sup _{\overline{\Omega_{L}}} \phi(\cdot, \theta)
$$

These two estimates will be fruitful in Section 3 to solve our inverse problem.
Proof. First we define $\beta_{0}=\inf _{x \in \Omega_{l}} \psi(x, \theta)=\inf _{x \in \Omega_{l}}\left(\left|x^{\prime}-a^{\prime}\right|^{2}-x_{1}^{2}\right)+M_{1}$ and $\beta_{1}>0$ by

$$
\beta_{1}^{2}=\sup _{x \in \Omega_{L}}\left(\left|x^{\prime}-a^{\prime}\right|^{2}-x_{1}^{2}\right)-\inf _{x \in \Omega_{l}}\left(\left|x^{\prime}-a^{\prime}\right|^{2}-x_{1}^{2}\right)
$$

Notice that $\beta_{1}^{2}=\sup _{x^{\prime} \in \omega}\left|x^{\prime}-a^{\prime}\right|^{2}-\inf _{x^{\prime} \in \omega}\left|x^{\prime}-a^{\prime}\right|^{2}+l^{2}$.
Then, we consider $L$ and $T=2 L$ sufficiently large such that $L>l$ and $\beta_{2}=T-\theta-\beta_{1}>0$ i.e.

$$
\sup _{x^{\prime} \in \omega}\left|x^{\prime}-a^{\prime}\right|^{2}-\inf _{x^{\prime} \in \omega}\left|x^{\prime}-a^{\prime}\right|^{2}+l^{2}<L^{2}
$$

But simultaneously, we might have to push $a^{\prime}$ further away from $\omega$ to ensure the condition (1.2). Thus, we want

$$
\begin{equation*}
\sup _{x^{\prime} \in \omega}\left|x^{\prime}-a^{\prime}\right|^{2}-\inf _{x^{\prime} \in \omega}\left|x^{\prime}-a^{\prime}\right|^{2}+l^{2}<L^{2}<\inf _{x^{\prime} \in \omega}\left|x^{\prime}-a^{\prime}\right|^{2} . \tag{1.5}
\end{equation*}
$$

To validate this construction, we prove that for $\left|a^{\prime}\right|$ large enough

$$
\begin{equation*}
\sup _{x^{\prime} \in \omega}\left|x^{\prime}-a^{\prime}\right|^{2}+l^{2}<2 \inf _{x^{\prime} \in \omega}\left|x^{\prime}-a^{\prime}\right|^{2} \tag{1.6}
\end{equation*}
$$

Indeed, from the triangle inequality we have

$$
\sup _{x^{\prime} \in \omega}\left|x^{\prime}-a^{\prime}\right| \leq \inf _{x^{\prime} \in \omega}\left|x^{\prime}-a^{\prime}\right|+\operatorname{diam}(\omega)
$$

Thus,

$$
\begin{aligned}
\sup _{x^{\prime} \in \omega}\left|x^{\prime}-a^{\prime}\right|^{2}+l^{2} & \leq\left(\inf _{x^{\prime} \in \omega}\left|x^{\prime}-a^{\prime}\right|+\operatorname{diam}(\omega)\right)^{2}+l^{2} \\
& =\inf _{x^{\prime} \in \omega}\left|x^{\prime}-a^{\prime}\right|^{2}+2 \operatorname{diam}(\omega) \inf _{x^{\prime} \in \omega}\left|x^{\prime}-a^{\prime}\right|+\operatorname{diam}(\omega)^{2}+l^{2} .
\end{aligned}
$$

As

$$
\lim _{\left|a^{\prime}\right| \rightarrow+\infty} \inf _{x^{\prime} \in \omega}\left|x^{\prime}-a^{\prime}\right|=+\infty
$$

this proves (1.6). Thus, there exist $a \in \mathbb{R}^{n} \backslash \Omega$ and $L>l$ (and $T=2 L$ ) such that (1.5) holds.
Here we take advantage of the particular form of our weight function $d(x)$. With these definitions, we get

$$
\begin{aligned}
(T-\theta)^{2} \geq \beta_{1}^{2}+\beta_{2}^{2} & =\sup _{x \in \Omega_{L}}\left(\left|x^{\prime}-a^{\prime}\right|^{2}-x_{1}^{2}\right)-\inf _{x \in \Omega_{l}}\left(\left|x^{\prime}-a^{\prime}\right|^{2}-x_{1}^{2}\right)+\beta_{2}^{2} \\
& \geq \sup _{x \in \Omega_{L}}\left(\left|x^{\prime}-a^{\prime}\right|^{2}-x_{1}^{2}\right)+M_{1}-\beta_{0}+\beta_{2}^{2}
\end{aligned}
$$

Then, for all $x \in \Omega_{L}$,

$$
\psi(x, T) \leq\left|x^{\prime}-a^{\prime}\right|^{2}-x_{1}^{2}-\sup _{x \in \Omega_{L}}\left(\left|x^{\prime}-a^{\prime}\right|^{2}-x_{1}^{2}\right)+\beta_{0}-\beta_{2}^{2} \leq \beta_{0}-\beta_{2}^{2}
$$

As $\psi(x, 0)=\psi(x, T)$, we deduce that there exists $\epsilon>0$ such that $\epsilon<\frac{T}{4}$ and

$$
\text { for all } x \in \Omega_{L} \text { and } t \in((0,2 \epsilon) \cup(T-2 \epsilon, T)), \quad \psi(x, t)<\beta_{0}-\frac{\beta_{2}^{2}}{2} \text {. }
$$

Now, we choose $\epsilon$ small enough such that $l \leq L-2 \epsilon$. Due to the symmetric role played by $t-\theta$ and $x_{1}$ in the formulation of $\psi(x, t)$, by the same way, we have

$$
\forall x \in((-L,-L+2 \epsilon) \cup(L-2 \epsilon, L)) \times \omega \text { and } t \in(0, T), \quad \psi(x, t)<\beta_{0}-\frac{\beta_{2}^{2}}{2}
$$

These two estimates end the proof of Proposition 1.1.
An example of $\Omega_{l}$ and $\Omega_{L}$ is given in Figure 1.


Figure 1: Representation of $\Omega_{l}$ and $\Omega_{L}$ in a 3 D setting

Let us remark that the weight function is bounded by below (away from 0 ) et by above on $Q_{L}$. The fact that the weight function do not explode will be convenient in various estimates.

Finally, we define the lighted boundary from $a$ and the lighted lateral boundary from $a$ by

$$
\begin{equation*}
\Gamma_{L}=\left\{x \in \partial \Omega_{L},\langle x-a, \nu(x)\rangle \geq 0\right\} \quad \text { and } \quad \gamma_{L}=\Gamma_{L} \cap \partial \Omega \tag{1.7}
\end{equation*}
$$

Therefore $\gamma_{L}$ does not contain any cross-section (see Figure 2). Here $\langle\cdot, \cdot\rangle$ denotes the usual inner product in $\mathbb{R}^{n}$ and $\nu(x)$ is the outward unit normal vector to $\partial \Omega_{L}$ at $x$.

Notice that, due to the geometry, $\nu_{1}(x)$ the first component of $\nu(x)$ vanishes on $\gamma_{L}$. Thus, we obtain

$$
\begin{equation*}
\sum_{j=1}^{n} c(x) \partial_{j} d(x) \nu_{j}(x) \leq 0, \quad \forall x \in \partial \Omega_{L} \backslash \Gamma_{L} \tag{1.8}
\end{equation*}
$$

This is one of the conditions required for the weight function in the Carleman inequality proved in [28] that we will use in Section 2.

Let $V_{L}$ be a neighborhood of the lateral boundary of $\Omega_{L}$ i.e. an open set such that $V_{L} \subset \subset \Omega$ and

$$
\begin{equation*}
\partial V_{L} \supset\left(\partial \Omega_{L} \cap \partial \Omega\right) \tag{1.9}
\end{equation*}
$$

We also define

$$
\begin{equation*}
\widetilde{\Omega}_{L}=\Omega_{L} \backslash V_{L} \tag{1.10}
\end{equation*}
$$

Let $C_{\text {min }}>0$ and $\widetilde{M}>0$ be given constants. Let $c^{*} \in C^{1}(\bar{\Omega}) \cap L^{\infty}(\Omega), c^{*}>C_{\text {min }}>0$ on $\bar{\Omega}$. We will consider the following admissible set of diffusion coefficients

$$
\begin{equation*}
D=\left\{c \in C^{1}(\bar{\Omega}) \cap L^{\infty}(\Omega) \text { such that } \inf _{\bar{\Omega}} c>C_{\min }, c=c^{*} \text { on } V_{L} \text { and }\|c\|_{C^{1}(\bar{\Omega})}<\widetilde{M}\right\} \tag{1.11}
\end{equation*}
$$

Notice that this means that the diffusion coefficient is supposed to be known in a neighborhood of the lateral boundary of interest.

To ease the reading of Section 4, up to a restriction of $V_{L}$, we assume that there exists $r>0$ such that

$$
\begin{equation*}
V_{L} \cap \Omega_{L}=(-L, L) \times\left\{x^{\prime} \in \omega ; \operatorname{dist}\left(x^{\prime}, \mathbb{R}^{n-1} \backslash \omega\right)<r\right\} \tag{1.12}
\end{equation*}
$$

This is illustrated in Figure 3.


Figure 2: The lighted lateral boundary from $a$


Figure 3: A neighborhood of the lateral boundary

### 1.3 Regularity assumptions for the inverse problem

The method of Carleman estimate used in this paper requires that solutions of the problem (1.1) are sufficiently regular. Indeed the Buckgheim-Klibanov method [3] implies several time differentiations of system (1.1).

We will use the following notations. Let $\alpha=\left(\alpha_{1}, \cdots, \alpha_{n}\right)$ be a multi-index with $\alpha_{i} \in \mathbb{N}$. We set $\partial_{x}^{\alpha}=\partial_{1}^{\alpha_{1}} \cdots \partial_{n}^{\alpha_{n}},|\alpha|=\alpha_{1}+\cdots+\alpha_{n}$ and we define, for any $\ell \in \mathbb{N}^{*}$,

$$
H^{2 \ell, \ell}\left(Q_{*}\right)=\left\{u \in L^{2}\left(Q_{*}\right) ; \partial_{x}^{\alpha} \partial_{t}^{\alpha_{n+1}} u \in L^{2}\left(Q_{*}\right),|\alpha|+2 \alpha_{n+1} \leq 2 \ell\right\}
$$

endowed with its norm

$$
\|u\|_{H^{2 \ell, \ell}\left(Q_{*}\right)}^{2}=\sum_{|\alpha|+2 \alpha_{n+1} \leq 2 \ell}\left\|\partial_{x}^{\alpha} \partial_{t}^{\alpha_{n+1}} u\right\|_{L^{2}\left(Q_{*}\right)}^{2}
$$

In all what follows, we assume that $c \in D$ and that $u$ is an element of $\mathcal{H}=H^{6,3}\left(Q_{L}\right)$ satifying the a-priori bound

$$
\begin{equation*}
\|u\|_{\mathcal{H}}<M \quad \text { for a given } M>0 \tag{1.13}
\end{equation*}
$$

In order to prove such regularity, let us denote by $A_{0}$ the self-adjoint realization of $-\nabla(c \nabla \cdot)$ with domain

$$
\begin{equation*}
D\left(A_{0}\right)=H^{2}(\Omega) \cap H_{0}^{1}(\Omega) \tag{1.14}
\end{equation*}
$$

We derive (1.1) twice with respect to the time variable. Then, applying [17, Chapter IV, Theorem 9.1] (see also in the same reference Chapter IV, § 4 for the assumptions on the domain, Chapter IV, $\S 5$ for the assumptions on the operator and Chapter I, $\S 1$ for the definitions of functional spaces), we obtain that if

$$
\begin{equation*}
u_{0} \in \mathcal{H}_{0}=\left\{u_{0} \in D\left(A_{0}^{2}\right) ; A_{0}^{2} u_{0} \in H_{0}^{1}(\Omega)\right\} \tag{1.15}
\end{equation*}
$$

then there exists a unique solution $u$ to system (1.1) which satisfies $u \in \mathcal{H}$.
Remark 1.1. Let us mention that the space $\mathcal{H}$ is too regular for our purpose. For the proof of Theorem 1.1, the exact assumption we need in our proof is that $u$ solves (1.1) and enjoys the following regularity

$$
u \in C^{0}\left(0, T, H^{3}\left(\Omega_{L}\right)\right) \cap C^{1}\left(0, T, H^{1}\left(\Omega_{L}\right)\right) \cap H^{3}\left(0, T, L^{2}\left(\Omega_{L}\right)\right) \cap L^{2}\left(0, T, H^{3}\left(\Omega_{L}\right)\right)
$$

### 1.4 Main results

The first main result of this article is the following global stability estimate. For any $c_{1}, c_{2} \in D$, and any $u_{0,1}, u_{0,2} \in L^{2}(\Omega)$ we denote, for any $j \in\{1,2\}$, by $u_{j}$ the solution of (1.1) where $c_{j}$ and $u_{0, j}$ are substituted respectively to $c$ and $u_{0}$.
Theorem 1.1. Let $l>0$. Let $T>0, L>l$ and $a \in \mathbb{R}^{n} \backslash \Omega_{L}$ satisfying the conditions of Proposition 1.1. Assume that $c_{1}, c_{2} \in D$ and that $u_{0,1}, u_{0,2}$ are sufficiently regular (for instance satisfying (1.15)) to ensure that the associated solutions $u_{1}, u_{2}$ of (1.1) satisfy (1.13). Suppose that there exists $\delta>0$ such that

$$
\begin{equation*}
\left|\nabla d \cdot \nabla u_{2}\left(x, \frac{T}{2}\right)\right| \geq \delta, \quad \text { for a.e. } x \in \widetilde{\Omega}_{L} \tag{1.16}
\end{equation*}
$$

where $\widetilde{\Omega}_{L}$ is defined by (1.10). Then there exist constants $K$ and $\kappa$ such that

$$
\left\|c_{1}-c_{2}\right\|_{H^{1}\left(\Omega_{l}\right)}^{2} \leq K\left(\left\|\left(u_{1}-u_{2}\right)(\cdot, T / 2)\right\|_{H^{3}\left(\Omega_{L}\right)}^{2}+\int_{\gamma_{L} \times(0, T)} \sum_{k=1}^{2}\left|\partial_{\nu}\left(\partial_{t}^{k}\left(u_{1}-u_{2}\right)\right)\right|^{2}\right)^{\kappa}
$$

where $\partial_{\nu} u=\nu \cdot \nabla u$.
Here, $K>0$ and $\kappa \in(0,1)$ are two constants depending only on $\omega, l, L, T, a, \epsilon, M, \widetilde{M}$, $C_{\text {min }}$ and $\delta .{ }^{1}$
Remark 1.2. A careful inspection of the proof of Theorem 1.1 shows that the same result holds if we only assume that the diffusion coefficients $c_{1}$ and $c_{2}$ as well as their gradients are known on the boundary of $\Omega_{L}$ (instead of in a neighborhood of this boundary) i.e. if $c_{1}, c_{2} \in C^{1}(\bar{\Omega}) \cap L^{\infty}(\Omega)$ are such that

$$
\inf _{\bar{\Omega}} c_{j}>C_{\text {min }}, \quad c_{j}=c^{*} \text { and } \nabla c_{j}=\nabla c^{*} \text { on } \partial \Omega \cap \partial \Omega_{L}, \quad \text { and } \quad\left\|c_{j}\right\|_{C^{1}(\bar{\Omega})}<\widetilde{M}
$$

and one replaces (1.16) by

$$
\left|\nabla d \cdot \nabla u_{2}\left(x, \frac{T}{2}\right)\right| \geq \delta, \quad \text { for a.e. } x \in \Omega_{L} .
$$

Remark 1.3. Notice that, as already highlighted, the boundary observation is on $\gamma_{L}$ and thus does not contain any cross-section in $\Omega_{L}$. Though the equation is set on the unbounded guide $\Omega$, the stability estimate on $\Omega_{l}$ is obtained with measurements on the finite portion $\Omega_{L}$.
Remark 1.4. We look at the homogeneous Dirichlet case in order to simplify the regularity required for the solution $u$ of (1.1) but we could obtain the same inverse result in the case of the inhomogeneous Dirichlet case

$$
\left\{\begin{array}{l}
\partial_{t} u-\nabla \cdot(c \nabla u)=0 \text { in } Q,  \tag{1.17}\\
u=h \text { on } \Sigma, \\
u(x, 0)=u_{0}(x) \text { in } \Omega .
\end{array}\right.
$$

The proof of Theorem 1.1 follows the ideas of [3] and is proved in Section 3. It strongly relies on the Carleman estimate proved in Section 2. This strategy is quite classical.

The drawback is that, to develop the Carleman machinery, one needs the technical (not so easy to verify) assumption (1.16). In the following we propose a strategy to remove this assumption:

- we prove the existence of a boundary control $h$ such that the associated solution of (1.17) satisfies (1.16);
- we fix this boundary condition for all the required measurements that is we deal only with system (1.17) for a fixed $h$.
More precisely, the second main result of this article reads as follows. For $j=1,2$ and $h \in$ $L^{2}(0, T ; \partial \Omega)$ we denote by $u_{j}$ the solutions of (1.17) where $c_{j}$ and $u_{0, j}$ are substituted respectively to $c$ and $u_{0}$.
Theorem 1.2. Let $l>0$. Let $T>0, L>l$ and $a \in \mathbb{R}^{n} \backslash \Omega_{L}$ satisfying the conditions of Proposition 1.1.

Let $u_{0,1}, u_{0,2} \in \mathcal{H}_{0}$. Let $c_{2} \in D$ be such that $c_{2}$ and for every $j \in\{1, \ldots, n\}$, $\partial_{j} c_{2}$, are uniformly continuous and bounded functions in $\Omega$. Then there exists a control $h \in L^{2}\left(\gamma_{L} \times(0, T)\right)$ depending on $c_{2}$ and $u_{0,2}$ such that for any $c_{1} \in D$,

$$
\left\|c_{1}-c_{2}\right\|_{H^{1}\left(\Omega_{l}\right)}^{2} \leq K\left(\left\|\left(u_{1}-u_{2}\right)(., T / 2)\right\|_{H^{3}\left(\Omega_{L}\right)}^{2}+\int_{\gamma_{L} \times(0, T)} \sum_{k=1}^{2}\left|\partial_{\nu}\left(\partial_{t}^{k}\left(u_{1}-u_{2}\right)\right)\right|^{2}\right)^{\kappa}
$$

[^1]where $\partial_{\nu} \bullet=\nu \cdot \nabla \bullet$.
Here, $K>0$ and $\kappa \in(0,1)$ are two constants depending only on $\omega, l, L, T, a, \epsilon, \widetilde{M}, C_{\min }{ }^{2}$, $u_{0,2}$ and $M$ where
$$
M>\left\|u_{1}-u_{2}\right\|_{\mathcal{H}}
$$

As noticed the boundary condition $h$ exerted for the measurements associated to $c_{1}$ only depends on $c_{2}$. Morally, this means that the coefficient $c_{2}$ is known and that Theorem 1.2 is a local stability estimate around $c_{2}$.

Remark 1.5. With the same proof, for any $\widetilde{\gamma}_{L} \subset \gamma_{L}$, one can obtain the exact same statement but with a control $h$ supported in $\widetilde{\gamma}_{L}$.

Remark 1.6. The existence of the control $h$ stated in Theorem 1.2 is obtained in Section 4 using a non-constructive proof. For the reader interested in the actual computation of (an approximation of) such a control we mention the book [15, Chapters 1 and 2] where this question is studied.

This result is proved in Section 4. The main difficulty is to prove the existence of a control $h$ such that the associated solution satisfies (1.16). This part of the proof is inspired by the strategy developed in [1] for systems on bounded domains.

## 2 Global Carleman estimate for a parabolic equation in a cylindrical domain

We start with a global Carleman-type estimate proved by Yuan and Yamamoto [28] in a bounded domain. Its validity is ensured by the estimates (1.2) and (1.8) satisfied by our weight function. For more informations about Carleman estimates in a parabolic setting we refer to Yamamoto [27].

Let $s>0$ and denote

$$
L H S(u)=\int_{Q_{L}}\left(\frac{1}{s \phi}\left(\left|\partial_{t} u\right|^{2}+|\Delta u|^{2}\right)+s \phi|\nabla u|^{2}+s^{3} \phi^{3}|u|^{2}\right) e^{2 s \phi}
$$

In the following parts, $C$ will be a generic positive constant. When needed, we will specify its dependency with respect to the different parameters.

Proposition 2.1 ([28, Theorem 2.1, (2)]). There exist a value of $\lambda>0$ and positive constants $s_{0}$ and $C$ such that

$$
\begin{equation*}
L H S(u) \leq C\left\|e^{s \phi} f\right\|_{L^{2}\left(Q_{L}\right)}^{2}+C s \int_{\Gamma_{L} \times(0, T)}\left|\partial_{\nu} u\right|^{2} e^{2 s \phi} \tag{2.1}
\end{equation*}
$$

for all $s>s_{0}$, and all $u \in H^{2,1}\left(Q_{L}\right)$ satisfying

$$
\left\{\begin{array}{l}
\partial_{t} u-\nabla \cdot(c \nabla u)=f, \quad \text { in } \Omega_{L} \\
u(\cdot, 0)=u(\cdot, T)=0 \quad \text { in } \Omega_{L} \\
u=0 \quad \text { on } \partial \Omega_{L} \times(0, T)
\end{array}\right.
$$

Here, the constants $\lambda$, $s_{0}$ and $C$ depend on $\omega, L, T, a, \widetilde{M}$ and $C_{m i n} .{ }^{3}$

[^2]Let us remark that this Carleman inequality uses also $\lambda$ as a second large parameter. As we will not use it, we now consider $\lambda$ fixed in all the rest of the article such that Proposition 2.1 holds.

As in [4, Proposition 4.2], we deduce the following Carleman inequality and we detail the proof for better understanding. The key difference with the Carleman inequality of Proposition 2.1 is to remove, on the cross-sections of $\Omega_{L}$, the boundary condition and the observation.
Proposition 2.2. Let $s_{0}$ be the constant given by Proposition 2.1. There exists a positive constant $C$ such that

$$
\begin{equation*}
L H S(u) \leq C\left\|e^{s \phi} f\right\|_{L^{2}\left(Q_{L}\right)}^{2}+C s^{3} e^{2 s d_{1}}\|u\|_{H^{2,1}\left(Q_{L}\right)}^{2}+C s \int_{\gamma_{L} \times(0, T)}\left|\partial_{\nu} u\right|^{2} e^{2 s \phi} \tag{2.2}
\end{equation*}
$$

for all $s>s_{0}$ and all $u \in H^{2,1}\left(Q_{L}\right)$ satisfying

$$
\left\{\begin{array}{l}
\partial_{t} u-\nabla \cdot(c \nabla u)=f, \quad \text { in } \Omega_{L}, \\
u(\cdot, 0)=u(\cdot, T)=0 \quad \text { in } \Omega_{L} \\
u=0 \quad \text { on }\left(\partial \Omega \cap \partial \Omega_{L}\right) \times(0, T) .
\end{array}\right.
$$

Here, the constant $C$ depends on $\omega, L, T, a, \epsilon, \widetilde{M}$ and $C_{m i n} .^{4}$
Proof. Let $\chi, \eta$ be $C^{\infty}$ cut-off functions satisfying $0 \leq \chi \leq 1,0 \leq \eta \leq 1$ and

$$
\eta(t)=\left\{\begin{array}{l}
0 \text { if } t \in[0, \epsilon] \cup[T-\epsilon, T],  \tag{2.3}\\
1 \text { if } t \in[2 \epsilon, T-2 \epsilon],
\end{array} \quad \chi(x)=\left\{\begin{array}{l}
0 \text { if } x_{1} \in(-\infty,-L+\epsilon] \cup[L-\epsilon,+\infty), \\
1 \text { if } x_{1} \in[-L+2 \epsilon, L-2 \epsilon],
\end{array}\right.\right.
$$

with $\epsilon$ defined in Proposition 1.1.
Recall that $\partial_{t} u-\nabla \cdot(c \nabla u)=f$. We consider $y=\eta \chi u$ and we get

$$
\partial_{t} y-\nabla \cdot(c \nabla y)=h \quad \text { with } \quad h=\eta \chi f+\eta R(u)+\left(\partial_{t} \eta\right) \chi u
$$

where $R$ is the first order differential operator defined by $R(u)=-\nabla \cdot(c u \nabla \chi)-c \nabla \chi \cdot \nabla u$.
Then, applying the previous Carleman estimate (2.1) we deduce that there exists a positive constant $C$ such that

$$
L H S(y) \leq C\left\|e^{s \phi} h\right\|_{L^{2}\left(Q_{L}\right)}^{2}+C s \int_{\Gamma_{L} \times(0, T)}\left|\partial_{\nu} y\right|^{2} e^{2 s \phi}
$$

Recall that $\gamma_{L}$ is defined in (1.7). Then,

$$
\int_{\Gamma_{L} \times(0, T)}\left|\partial_{\nu} y\right|^{2} e^{2 s \phi} \leq \int_{\gamma_{L} \times(0, T)}\left|\partial_{\nu} u\right|^{2} e^{2 s \phi}
$$

Moreover

$$
\begin{aligned}
& \quad\left\|e^{s \phi} \eta R(u)\right\|_{L^{2}\left(Q_{L}\right)}^{2} \leq C e^{2 s d_{1}}\|u\|_{L^{2}\left(0, T, H^{1}\left(\Omega_{L}\right)\right)}^{2} \\
& \text { and }\left\|e^{s \phi}\left(\partial_{t} \eta\right) \chi u\right\|_{L^{2}\left(Q_{L}\right)}^{2} \leq C e^{2 s d_{1}}\|u\|_{L^{2}\left(0, T, L^{2}\left(\Omega_{L}\right)\right)}^{2}
\end{aligned}
$$

This implies

$$
\begin{equation*}
L H S(y) \leq C\left\|e^{s \phi} f\right\|_{L^{2}\left(Q_{L}\right)}^{2}+C e^{2 s d_{1}}\|u\|_{L^{2}\left(0, T, H^{1}\left(\Omega_{L}\right)\right)}^{2}+C s \int_{\gamma_{L} \times(0, T)}\left|\partial_{\nu} u\right|^{2} e^{2 s \phi} \tag{2.4}
\end{equation*}
$$

[^3]Now we deal with $L H S(y)$. For $j=0,1,2$, (with $\left.\nabla^{0} u=u, \nabla^{1} u=\nabla u, \nabla^{2} u=\Delta u\right)$ since $\chi u=(1-\eta) \chi u+y$, it comes that
$\left\|(s \phi)^{3 / 2-j} e^{s \phi} \nabla^{j}(\chi u)\right\|_{L^{2}\left(Q_{L}\right)} \leq\left\|(s \phi)^{3 / 2-j} e^{s \phi}(1-\eta) \nabla^{j}(\chi u)\right\|_{L^{2}\left(Q_{L}\right)}+\left\|(s \phi)^{3 / 2-j} e^{s \phi} \nabla^{j} y\right\|_{L^{2}\left(Q_{L}\right)}$.
Thus,

$$
\left\|(s \phi)^{3 / 2-j} e^{s \phi} \nabla^{j}(\chi u)\right\|_{L^{2}\left(Q_{L}\right)} \leq C\left(s^{3 / 2} e^{s d_{1}}\|u\|_{H^{2,1}\left(Q_{L}\right)}+\left\|(s \phi)^{3 / 2-j} e^{s \phi} \nabla^{j} y\right\|_{L^{2}\left(Q_{L}\right)}\right)
$$

Doing the same for the term

$$
\partial_{t}(\chi u)=\left(-\partial_{t} \eta\right) \chi u+(1-\eta) \chi \partial_{t} u+\partial_{t} y
$$

we deduce that

$$
\left\|(s \phi)^{-1 / 2} e^{s \phi} \partial_{t}(\chi u)\right\|_{L^{2}\left(Q_{L}\right)} \leq e^{s d_{1}}\|u\|_{H^{2,1}\left(Q_{L}\right)}+\left\|(s \phi)^{-1 / 2} e^{s \phi} \partial_{t} y\right\|_{L^{2}\left(Q_{L}\right)}
$$

Thus since

$$
L H S(\chi u)=\sum_{j=0}^{2}\left\|(s \phi)^{3 / 2-j} e^{s \phi} \nabla^{j} u\right\|_{L^{2}\left(Q_{L}\right)}^{2}+\left\|(s \phi)^{-1 / 2} e^{s \phi} \partial_{t}(\chi u)\right\|_{L^{2}\left(Q_{L}\right)}^{2}
$$

we have

$$
L H S(\chi u) \leq C\left(s^{3} e^{2 s d_{1}}\|u\|_{H^{2,1}\left(Q_{L}\right)}^{2}+L H S(y)\right)
$$

Then, by the identities

$$
\left\{\begin{array}{l}
\partial_{t} u=\partial_{t}(\chi u)+(1-\chi) \partial_{t} u \\
\nabla u=\nabla(\chi u)+(1-\chi) \nabla u-u \nabla \chi \\
\Delta u=\Delta(\chi u)+(1-\chi) \Delta u-2 \nabla \chi \cdot \nabla u-u \Delta \chi
\end{array}\right.
$$

we get

$$
L H S(u) \leq C\left(L H S(\chi u)+s^{3} e^{2 s d_{1}}\|u\|_{H^{2,1}\left(Q_{L}\right)}\right) \leq C\left(s^{3} e^{2 s d_{1}}\|u\|_{H^{2,1}\left(Q_{L}\right)}^{2}+L H S(y)\right)
$$

Then, from (2.4), this ends the proof.

## 3 Inverse Problem

### 3.1 Preliminary lemmas

We denote in the following $b(\theta)=b(\cdot, \theta)$ for any function $b$. As for the reconstruction of zeroorder coefficients (see [4, p.6]), we need to assume some hypothesis on the solution $u$ at time $\theta$ on $\Omega_{L}$. For the diffusion coefficient, the assumption is more involved and is exactly (1.16).

Now following an idea developed for example in [2, Lemma 2.4], we obtain the following result. Note that this key lemma is the one that requires the assumption (1.16). Thus, for the sake of completeness, we will give its proof.

Lemma 3.1. Assume that (1.16) is satisfied and consider the first order partial differential operator $P f=\nabla \cdot\left(f \nabla u_{2}(\theta)\right)$. Then there exist positive constants $s_{1}>0$ and $C>0$ such that for all $s \geq s_{1}$

$$
s^{2} \int_{\Omega_{L}} e^{2 s \phi(\theta)}|f|^{2} \leq C \int_{\Omega_{L}} e^{2 s \phi(\theta)}|P f|^{2}
$$

for all $f \in H_{0}^{1}\left(\Omega_{L}\right)$ such that $f=0$ in $V_{L}$ (see (1.9)).
Here, the constants $s_{1}$ and $C$ depend on $\omega, L, T, a, M$ and $\delta .{ }^{5}$

[^4]Proof. Let $f \in H_{0}^{1}\left(\Omega_{L}\right)$ be such that $f=0$ in $V_{L}$. We denote $w=e^{s \phi(\theta)} f$ and $Q w=$ $e^{s \phi(\theta)} P\left(e^{-s \phi(\theta)} w\right)$. So we get $Q w=P w-s w \nabla \phi(\theta) \cdot \nabla u_{2}(\theta)$. Therefore, we have

$$
\begin{aligned}
\int_{\Omega_{L}}|Q w|^{2} & \geq s^{2} \int_{\Omega_{L}} w^{2}\left|\nabla \phi(\theta) \cdot \nabla u_{2}(\theta)\right|^{2}-2 s \int_{\Omega_{L}} P w\left(w \nabla \phi(\theta) \cdot \nabla u_{2}(\theta)\right) \\
& =s^{2} \lambda^{2} \int_{\Omega_{L}} w^{2} \phi(\theta)^{2}\left|\nabla d \cdot \nabla u_{2}(\theta)\right|^{2}-2 s \lambda \int_{\Omega_{L}}(P w) w \phi(\theta)\left(\nabla d \cdot \nabla u_{2}(\theta)\right) .
\end{aligned}
$$

As, $P w=\nabla \cdot\left(w \nabla u_{2}(\theta)\right)=\nabla w \cdot \nabla u_{2}(\theta)+w \Delta u_{2}(\theta)$, we obtain,

$$
\begin{aligned}
\int_{\Omega_{L}}|Q w|^{2} & \geq s^{2} \lambda^{2} \int_{\Omega_{L}} w^{2} \phi(\theta)^{2}\left|\nabla d \cdot \nabla u_{2}(\theta)\right|^{2}-s \lambda \int_{\Omega_{L}} \phi(\theta)\left(\nabla\left(w^{2}\right) \cdot \nabla u_{2}(\theta)\right)\left(\nabla d \cdot \nabla u_{2}(\theta)\right) \\
& -2 s \lambda \int_{\Omega_{L}} w^{2} \phi(\theta) \Delta u_{2}(\theta)\left(\nabla d \cdot \nabla u_{2}(\theta)\right)
\end{aligned}
$$

Thus, integrating by parts the second term of the right-hand side, we obtain

$$
\begin{aligned}
\int_{\Omega_{L}}|Q w|^{2} & \geq s^{2} \lambda^{2} \int_{\Omega_{L}} w^{2} \phi(\theta)^{2}\left|\nabla d \cdot \nabla u_{2}(\theta)\right|^{2}+s \lambda \int_{\Omega_{L}} w^{2} \nabla \cdot\left(\phi(\theta)\left(\nabla d \cdot \nabla u_{2}(\theta)\right) \nabla u_{2}(\theta)\right) \\
& -2 s \lambda \int_{\Omega_{L}} w^{2} \phi(\theta) \Delta u_{2}(\theta)\left(\nabla d \cdot \nabla u_{2}(\theta)\right) .
\end{aligned}
$$

Getting back to the original variables we obtain

$$
\begin{aligned}
\int_{\Omega_{L}} e^{2 s \phi(\theta)}|P f|^{2}=\int_{\Omega_{L}}|Q w|^{2} & \geq s^{2} \lambda^{2} \int_{\Omega_{L}} e^{2 s \phi(\theta)} f^{2} \phi(\theta)^{2}\left|\nabla d \cdot \nabla u_{2}(\theta)\right|^{2} \\
& +s \lambda \int_{\Omega_{L}} e^{2 s \phi(\theta)} f^{2} \nabla \cdot\left(\phi(\theta)\left(\nabla d \cdot \nabla u_{2}(\theta)\right) \nabla u_{2}(\theta)\right) \\
& -2 s \lambda \int_{\Omega_{L}} e^{2 s \phi(\theta)} f^{2} \phi(\theta) \Delta u_{2}(\theta)\left(\nabla d \cdot \nabla u_{2}(\theta)\right)
\end{aligned}
$$

Thus, using assumption (1.16), there exists a constant $C>0$ such that

$$
\int_{\Omega_{L}} e^{2 s \phi(\theta)}|P f|^{2} \geq C s^{2} \int_{\tilde{\Omega}_{L}} e^{2 s \phi(\theta)} f^{2}-C s \int_{\Omega_{L}} e^{2 s \phi(\theta)} f^{2}
$$

where $\widetilde{\Omega}_{L}$ is defined by (1.10). As $f=0$ in $V_{L}$,

$$
\int_{\Omega_{L}} e^{2 s \phi(\theta)}|P f|^{2} \geq C s^{2} \int_{\Omega_{L}} e^{2 s \phi(\theta)} f^{2}-C s \int_{\Omega_{L}} e^{2 s \phi(\theta)} f^{2}
$$

We conclude taking $s$ sufficiently large.
Moreover we recall the following classical result (see [8]).
Lemma 3.2. There exist positive constants $s_{2}$ and $C$ such that

$$
\int_{\Omega_{L}} e^{2 s \phi(\theta)}|z(\theta)|^{2} \leq C s \int_{Q_{L}} e^{2 s \phi}|z|^{2}+\frac{C}{s} \int_{Q_{L}} e^{2 s \phi}\left|\partial_{t} z\right|^{2}
$$

for all $s \geq s_{2}$ and $z \in H^{1}\left(0, T ; L^{2}\left(\Omega_{L}\right)\right)$.
Here the constants $s_{2}$ and $C$ depend only on $T$ and $\epsilon .{ }^{6}$

[^5]Proof. Recall that $\eta$ is defined by (2.3). Consider any $w \in H^{1}\left(0, T ; L^{2}\left(\Omega_{L}\right)\right)$. We have

$$
\begin{aligned}
\int_{\Omega_{L}}|w(x, \theta)|^{2} & =\int_{\Omega_{L}}|\eta(\theta) w(x, \theta)|^{2} d x \\
& =\int_{\Omega_{L}} \int_{0}^{\theta} \partial_{t}\left(\eta^{2}(t)|w(x, t)|^{2}\right) d t d x \\
& =2 \int_{0}^{\theta} \int_{\Omega_{L}} \eta^{2}(t) w(x, t) \partial_{t} w(x, t) d x d t+2 \int_{0}^{\theta} \int_{\Omega_{L}} \eta(t) \eta^{\prime}(t)|w(x, t)|^{2} d x d t
\end{aligned}
$$

As $0 \leq \eta \leq 1$, using Young's inequality, it comes that for any $s>0$,

$$
\int_{\Omega_{L}}|w(x, \theta)|^{2} \leq C(s+1) \int_{Q_{L}}|w|^{2}+\frac{C}{s} \int_{Q_{L}}\left|\partial_{t} w\right|^{2}
$$

Then we can conclude replacing $w$ by $e^{s \phi} z$.

### 3.2 Proof of Theorem 1.1

Now we study the linearized inverse problem associated with (1.1) and we evaluate some Sobolev norm of the conductivity in terms of suitable observations of the solution of (1.1) on a part of the lateral boundary $\gamma_{L}$. For this consider the following systems with $c_{1}, c_{2} \in D$

$$
\left\{\begin{array}{l}
\partial_{t} u_{1}-\nabla \cdot\left(c_{1} \nabla u_{1}\right)=0 \text { in } Q, \\
u_{1}=0 \text { on } \Sigma \\
u_{1}(x, 0)=u_{0,1}(x) \text { in } \Omega,
\end{array}\right.
$$

and

$$
\left\{\begin{array}{l}
\partial_{t} u_{2}-\nabla \cdot\left(c_{2} \nabla u_{2}\right)=0 \text { in } Q \\
u_{2}=0 \text { on } \Sigma \\
u_{2}(x, 0)=u_{0,2}(x) \text { in } \Omega .
\end{array}\right.
$$

Let

$$
\begin{equation*}
y=u_{1}-u_{2}, \quad c=c_{1}-c_{2} \tag{3.1}
\end{equation*}
$$

We obtain

$$
\partial_{t} y-\nabla \cdot\left(c_{1} \nabla y\right)=\nabla \cdot\left(c \nabla u_{2}\right) \text { in } Q
$$

Remark 3.1. In the rest of the proof we will deal with $y$ only. Thus, even if we had started with (1.17), y would still enjoy homogeneous Dirichlet boundary condition. This remark will be crucial in Section 4 to prove Theorem 1.2.

We decompose the proof in three steps.

- First step: as in Proposition 2.2, we derive the equation satisfied by $z=\chi \eta y$ where the cut-off functions are defined by (2.3). This allows to work on the bounded domain $\Omega_{L} \times(0, T)$.
As $\eta$ only depends on $t$, we have

$$
\partial_{t}(\eta y)-\nabla \cdot\left(c_{1} \nabla(\eta y)\right)=\nabla \cdot\left(c \nabla\left(\eta u_{2}\right)\right)+y \partial_{t} \eta .
$$

Thus,

$$
\chi\left(\partial_{t}(\eta y)-\nabla \cdot\left(c_{1} \nabla(\eta y)\right)\right)=\chi\left(\nabla \cdot\left(c \nabla\left(\eta u_{2}\right)\right)+y \partial_{t} \eta\right) .
$$

Moreover

$$
\begin{aligned}
\nabla \cdot\left(c_{1} \nabla z\right) & =\nabla \cdot\left(c_{1} \nabla(\chi \eta y)\right) \\
& =\nabla \cdot\left(c_{1} \chi \nabla(\eta y)\right)+\nabla \cdot\left(c_{1} \eta y \nabla \chi\right) \\
& =\chi \nabla \cdot\left(c_{1} \nabla(\eta y)\right)+2 c_{1} \nabla(\eta y) \cdot \nabla \chi+\eta y \nabla \cdot\left(c_{1} \nabla \chi\right)
\end{aligned}
$$

and

$$
\nabla \cdot\left(c \chi \nabla\left(\eta u_{2}\right)\right)=\chi \nabla \cdot\left(c \nabla\left(\eta u_{2}\right)\right)+c \nabla\left(\eta u_{2}\right) \cdot \nabla \chi
$$

So we get
$\partial_{t} z-\nabla \cdot\left(c_{1} \nabla z\right)=\nabla \cdot\left(c \chi \nabla\left(\eta u_{2}\right)\right)+\chi y \partial_{t} \eta-2 c_{1} \nabla(\eta y) \cdot \nabla \chi-\eta y \nabla \cdot\left(c_{1} \nabla \chi\right)-c \nabla\left(\eta u_{2}\right) \cdot \nabla \chi$.

- Second step: to estimate $\partial_{t} z$ (appearing in Lemma 3.2) we differentiate (3.2) with respect to $t$.
Let

$$
z_{1}=\partial_{t} z, \quad z_{2}=\partial_{t}^{2} z
$$

Then,

$$
\begin{equation*}
\partial_{t} z_{1}-\nabla \cdot\left(c_{1} \nabla z_{1}\right)=f_{1} \tag{3.3}
\end{equation*}
$$

where

$$
\begin{aligned}
f_{1}: & \left.=\eta \nabla \cdot\left(c \chi \nabla\left(\partial_{t} u_{2}\right)\right)+\partial_{t} \eta \nabla \cdot\left(c \chi \nabla u_{2}\right)\right)-c \nabla\left(\partial_{t}\left(\eta u_{2}\right)\right) \cdot \nabla \chi+\chi \partial_{t}\left(y \partial_{t} \eta\right) \\
& -2 c_{1} \nabla\left(\partial_{t}(\eta y)\right) \cdot \nabla \chi-\partial_{t}(\eta y) \nabla \cdot\left(c_{1} \nabla \chi\right) .
\end{aligned}
$$

And in the same way

$$
\partial_{t} z_{2}-\nabla \cdot\left(c_{1} \nabla z_{2}\right)=f_{2}
$$

where

$$
\begin{aligned}
f_{2}: & =\eta \nabla \cdot\left(c \chi \nabla\left(\partial_{t}^{2}\left(u_{2}\right)\right)\right)+2 \partial_{t} \eta \nabla \cdot\left(c \chi \nabla\left(\partial_{t} u_{2}\right)\right)+\partial_{t}^{2} \eta \nabla \cdot\left(c \chi \nabla u_{2}\right) \\
& +\chi \partial_{t}^{2}\left(y \partial_{t} \eta\right)-2 c_{1} \nabla\left(\partial_{t}^{2}(\eta y)\right) \cdot \nabla \chi-\partial_{t}^{2}(\eta y) \nabla \cdot\left(c_{1} \nabla \chi\right)-c \nabla\left(\partial_{t}^{2}\left(\eta u_{2}\right)\right) \cdot \nabla \chi
\end{aligned}
$$

We evaluate (3.2) at $t=\theta$

$$
\begin{equation*}
\partial_{t} z(\theta)-\nabla \cdot\left(c_{1} \nabla z(\theta)\right)=P(c \chi)-c \nabla u_{2}(\theta) \cdot \nabla \chi-2 c_{1} \nabla y(\theta) \cdot \nabla \chi-y(\theta) \nabla \cdot\left(c_{1} \nabla \chi\right) \tag{3.4}
\end{equation*}
$$

with $P$ the operator defined in Lemma 3.1.
Using Lemma 3.1 with $f=c \chi$, we obtain an upper bound for a weighted $L^{2}$-norm of $c$. Notice that indeed $c \chi \in H_{0}^{1}\left(\Omega_{L}\right)$ : the homogeneous Dirichlet boundary condition on the cross-section is ensured by the presence of the cut-off function $\chi$. As $c_{1}, c_{2} \in D$, we also get that $c \chi=0$ in $V_{L}$. Thus,

$$
s^{2} \int_{\Omega_{L}} e^{2 s \phi(\theta)} c^{2} \chi^{2} \leq C \int_{\Omega_{L}} e^{2 s \phi(\theta)}(P(c \chi))^{2}
$$

Then, using also (3.4), there exists a positive constant $C$ such that for $s$ sufficiently large

$$
\begin{aligned}
s^{2} \int_{\Omega_{L}} e^{2 s \phi(\theta)} c^{2} \chi^{2} & \leq C \int_{\Omega_{L}} e^{2 s \phi(\theta)} c^{2}|\nabla \chi|^{2}+C \int_{\Omega_{L}} e^{2 s \phi(\theta)}\left|\partial_{t} z(\theta)\right|^{2} \\
& +C \int_{\Omega_{L}} e^{2 s \phi(\theta)}\left(|\Delta z(\theta)|^{2}+|\nabla z(\theta)|^{2}+|\nabla y(\theta)|^{2}+|y(\theta)|^{2}\right)
\end{aligned}
$$

Recall that $\int_{\Omega_{L}} e^{2 s \phi(\theta)}\left|\partial_{t} z(\theta)\right|^{2}=\int_{\Omega_{L}} e^{2 s \phi(\theta)}\left|z_{1}(\theta)\right|^{2}$. Thus, applying Lemma 3.2 to $z_{1}$ yields

$$
\begin{equation*}
s^{2} \int_{\Omega_{L}} e^{2 s \phi(\theta)} c^{2} \chi^{2} \leq C e^{2 s d_{1}}+C e^{2 s d_{2}} B_{1}(\theta)+C s \int_{Q_{L}} e^{2 s \phi}\left|z_{1}\right|^{2}+\frac{C}{s} \int_{Q_{L}} e^{2 s \phi}\left|z_{2}\right|^{2} \tag{3.5}
\end{equation*}
$$

with $B_{1}(\theta)=\|z(\theta)\|_{H^{2}\left(\Omega_{L}\right)}^{2}+\|y(\theta)\|_{H^{1}\left(\Omega_{L}\right)}^{2}$.
Moreover by the Carleman inequality (2.2), for $s$ sufficiently large, we have for $i=1,2$,

$$
\begin{align*}
\int_{Q_{L}} e^{2 s \phi}\left|z_{i}\right|^{2} & \leq \int_{Q_{L}} \phi^{3} e^{2 s \phi}\left|z_{i}\right|^{2} \\
& \leq \frac{C}{s^{3}} \int_{Q_{L}} e^{2 s \phi}\left|f_{i}\right|^{2}+C e^{2 s d_{1}}\left\|z_{i}\right\|_{H^{2,1}\left(Q_{L}\right)}^{2}+\frac{C}{s^{2}} \int_{\gamma_{L} \times(0, T)}\left|\partial_{\nu} z_{i}\right|^{2} e^{2 s \phi} \tag{3.6}
\end{align*}
$$

Thus, from (3.5) and (3.6), we obtain

$$
\begin{aligned}
& s^{2} \int_{\Omega_{L}} e^{2 s \phi(\theta)} c^{2} \chi^{2} \leq C e^{2 s d_{1}}+C e^{2 s d_{2}} B_{1}(\theta)+\frac{C}{s^{2}} \int_{Q_{L}} e^{2 s \phi}\left(\left|f_{1}\right|^{2}+\left|f_{2}\right|^{2}\right) \\
& +C s e^{2 s d_{1}}\left(\left\|z_{1}\right\|_{H^{2,1}\left(Q_{L}\right)}^{2}+\left\|z_{2}\right\|_{H^{2,1}\left(Q_{L}\right)}^{2}\right)+\frac{C}{s} \int_{\gamma_{L} \times(0, T)}\left(\left|\partial_{\nu} z_{1}\right|^{2}+\left|\partial_{\nu} z_{2}\right|^{2}\right) e^{2 s \phi}
\end{aligned}
$$

Notice that, as $\left\|u_{i}\right\|_{\mathcal{H}}<M$, the term $\left\|z_{1}\right\|_{H^{2,1}\left(Q_{L}\right)}^{2}+\left\|z_{2}\right\|_{H^{2,1}\left(Q_{L}\right)}^{2}$ is bounded.
Now let us deal with the terms $\int_{Q_{L}} e^{2 s \phi}\left|f_{i}\right|^{2}$. The first term in $f_{1}$ (see (3.3)) can be controlled by $(c \chi)^{2}+|\nabla(c \chi)|^{2}$. The other terms involve derivatives of the cut-off functions (and bounded quantities depending on $M$ and $\widetilde{M})$. Thus, since $e^{2 s \phi} \leq e^{2 s \phi(\theta)}$, we get

$$
\int_{Q_{L}}\left|f_{1}\right|^{2} e^{2 s \phi} \leq C e^{2 s d_{1}}+C \int_{\Omega_{L}}\left((c \chi)^{2}+|\nabla(c \chi)|^{2}\right) e^{2 s \phi(\theta)}
$$

The same computations on $f_{2}$ finally imply

$$
\begin{align*}
s^{2} \int_{\Omega_{L}} e^{2 s \phi(\theta)} c^{2} \chi^{2} & \leq C e^{2 s d_{2}} B_{1}(\theta)+\frac{C}{s^{2}} \int_{\Omega_{L}}\left((c \chi)^{2}+|\nabla(c \chi)|^{2}\right) e^{2 s \phi(\theta)}+C s e^{2 s d_{1}} \\
& +\frac{C}{s} e^{2 s d_{2}} \int_{\gamma_{L} \times(0, T)}\left(\left|\partial_{\nu} z_{1}\right|^{2}+\left|\partial_{\nu} z_{2}\right|^{2}\right) \tag{3.7}
\end{align*}
$$

- Third step: we apply the same strategy to derive an estimate of a weighted $H^{1}$-norm of $c$. For any integer $1 \leq i \leq n$, taking the space derivative with respect to $x_{i}$ in (3.4), we obtain

$$
\begin{gather*}
\partial_{t}\left(\partial_{x_{i}} z(\theta)\right)-\nabla \cdot\left(c_{1} \nabla \partial_{x_{i}} z(\theta)\right)=\nabla \cdot\left(\left(\partial_{x_{i}} c_{1}\right) \nabla z(\theta)\right)+\nabla \cdot\left(\partial_{x_{i}}(c \chi) \nabla u_{2}(\theta)\right) \\
+\nabla \cdot\left(c \chi \nabla\left(\partial_{x_{i}} u_{2}(\theta)\right)\right)-\partial_{x_{i}}\left(c \nabla u_{2}(\theta) \cdot \nabla \chi\right)-\partial_{x_{i}}\left(2 c_{1} \nabla y(\theta) \cdot \nabla \chi\right)-\partial_{x_{i}}\left(y(\theta) \nabla \cdot\left(c_{1} \nabla \chi\right)\right) . \tag{3.8}
\end{gather*}
$$

Notice that the second term of the right-hand side can be expressed as

$$
\nabla \cdot\left(\partial_{x_{i}}(c \chi) \nabla u_{2}(\theta)\right)=P\left(\partial_{x_{i}}(c \chi)\right)
$$

As previously, the fact that $c_{1}, c_{2} \in D$ and the definition of $\chi$ imply that we can apply again Lemma 3.1: there exists a positive constant $C$ such that for $s$ sufficiently large,

$$
s^{2} \int_{\Omega_{L}} e^{2 s \phi(\theta)}\left(\partial_{x_{i}}(c \chi)\right)^{2} \leq C \int_{\Omega_{L}} e^{2 s \phi(\theta)}\left(P\left(\partial_{x_{i}}(c \chi)\right)\right)^{2} .
$$

Thus, expressing $P\left(\partial_{x_{i}}(c \chi)\right)$ from (3.8) we obtain

$$
\begin{aligned}
s^{2} \int_{\Omega_{L}} e^{2 s \phi(\theta)}\left(\partial_{x_{i}}(c \chi)\right)^{2} & \leq C e^{2 s d_{2}} B_{2}(\theta)+C \int_{\Omega_{L}} e^{2 s \phi(\theta)}\left|\partial_{t}\left(\partial_{x_{i}} z(\theta)\right)\right|^{2}+C \int_{O_{L, \epsilon}} e^{2 s \phi(\theta)} \\
& +C \int_{\Omega_{L}} e^{2 s \phi(\theta)}\left(|c \chi|^{2}+|\nabla(c \chi)|^{2}\right)
\end{aligned}
$$

with $B_{2}(\theta)=\|z(\theta)\|_{H^{3}\left(\Omega_{L}\right)}^{2}+\|y(\theta)\|_{H^{2}\left(\Omega_{L}\right)}^{2}$.
Then,

$$
\begin{aligned}
s^{2} \int_{\Omega_{L}} e^{2 s \phi(\theta)}\left(\partial_{x_{i}}(c \chi)\right)^{2} & \left.\leq C e^{2 s d_{2}} B_{2}(\theta)+C \int_{\Omega_{L}} e^{2 s \phi(\theta)} \mid \partial_{x_{i}} z_{1}(\theta)\right)\left.\right|^{2} \\
& +C \int_{\Omega_{L}} e^{2 s \phi(\theta)}\left(|c \chi|^{2}+|\nabla(c \chi)|^{2}\right)+C e^{2 s d_{1}} .
\end{aligned}
$$

So, using Lemma 3.2,

$$
\begin{align*}
s^{2} \int_{\Omega_{L}} e^{2 s \phi(\theta)}\left(\partial_{x_{i}}(c \chi)\right)^{2} & \leq C e^{2 s d_{2}} B_{2}(\theta)+C \int_{\Omega_{L}} e^{2 s \phi(\theta)}\left(|c \chi|^{2}+|\nabla(c \chi)|^{2}\right)+C e^{2 s d_{1}} \\
& +C s \int_{Q_{L}} e^{2 s \phi}\left(\partial_{x_{i}} z_{1}\right)^{2}+\frac{C}{s} \int_{Q_{L}} e^{2 s \phi}\left(\partial_{x_{i}} z_{2}\right)^{2} \tag{3.9}
\end{align*}
$$

Moreover by the Carleman inequality (2.2), we have for $j=1,2$,

$$
\begin{aligned}
s \int_{Q_{L}} e^{2 s \phi}\left|\partial_{x_{i}} z_{j}\right|^{2} & \leq C \int_{Q_{L}} s \phi e^{2 s \phi}\left|\partial_{x_{i}} z_{j}\right|^{2} \\
& \leq C \int_{Q_{L}} e^{2 s \phi}\left|f_{j}\right|^{2}+C s^{3} e^{2 s d_{1}}\left\|z_{j}\right\|_{H^{2,1}\left(Q_{L}\right)}^{2}+C s \int_{\gamma_{L} \times(0, T)}\left|\partial_{\nu} z_{j}\right|^{2} e^{2 s \phi}
\end{aligned}
$$

The term in $f_{j}$ in the right-hand side is estimated as in the previous step. This leads to

$$
\begin{equation*}
s \int_{Q_{L}} e^{2 s \phi}\left|\partial_{x_{i}} z_{j}\right|^{2} \leq C \int_{Q_{L}} e^{2 s \phi}\left((c \chi)^{2}+|\nabla(c \chi)|^{2}\right)+C s^{3} e^{2 s d_{1}}+C s \int_{\gamma_{L} \times(0, T)}\left|\partial_{\nu} z_{j}\right|^{2} e^{2 s \phi} \tag{3.10}
\end{equation*}
$$

From (3.9) and (3.10) we deduce

$$
\begin{align*}
s^{2} \int_{\Omega_{L}} e^{2 s \phi(\theta)}\left(\partial_{x_{i}}(c \chi)\right)^{2} & \leq C e^{2 s d_{2}} B_{2}(\theta)+C \int_{\Omega_{L}} e^{2 s \phi(\theta)}\left(|c \chi|^{2}+|\nabla(c \chi)|^{2}\right)+C s^{3} e^{2 s d_{1}} \\
& +C s \int_{\gamma_{L} \times(0, T)}\left(\left|\partial_{\nu} z_{1}\right|^{2}+\left|\partial_{\nu} z_{2}\right|^{2}\right) e^{2 s \phi} \tag{3.11}
\end{align*}
$$

which ends this step.
Using inequalities (3.11) for $1 \leq i \leq n$ and gathering with (3.7), we get for $s$ sufficiently large

$$
s^{2} \int_{\Omega_{L}} e^{2 s \phi(\theta)}\left((c \chi)^{2}+|\nabla(c \chi)|^{2}\right) \leq C s B_{3}(\theta)+C s^{3} e^{2 s d_{1}}
$$

with $B_{3}(\theta)=B_{1}(\theta)+B_{2}(\theta)+\int_{\gamma_{L} \times(0, T)}\left(\left|\partial_{\nu} z_{1}\right|^{2}+\left|\partial_{\nu} z_{2}\right|^{2}\right)$.
Therefore

$$
\|c\|_{H^{1}\left(\Omega_{l}\right)}^{2} \leq C\left(e^{2 s\left(d_{2}-d_{0}\right)} B_{3}(\theta)+s e^{2 s\left(d_{1}-d_{0}\right)}\right)
$$

As $d_{1}-d_{0}<0$ and $d_{2}-d_{0}>0$, optimizing this last inequality with respect to $s$, we complete the proof.

## 4 Removing the technical assumption by an adapted control

In Theorem 1.1, we assume that (1.16) satisfied. However since this hypothesis concerns the gradient of one solution of (1.1) in $\widetilde{\Omega}_{L}$, it is difficult to verify it in the case of real applications. Thus, we develop ideas based on control theory to ensure such strong hypothesis.

The goal of this section is to prove that for any $c \in D$ sufficiently regular, there exists a control $h$ such that the associated solution of (1.17) satisfies (1.16). Then, the proof of Theorem 1.2 will follow directly (see Section 4.4). As the condition (1.16) involves a property for the gradient, we will seek for controllability results in stronger norms, namely $H^{1}\left(\Omega_{L}\right)$. The strategy is inspired by [1] and is as follows:

- Assume that there is a function $u_{b}$ sufficiently smooth and satisfying (1.16) i.e. there exists $\delta>0$ such that

$$
\left|\nabla u_{b} \cdot \nabla d\right| \geq \delta>0 \quad \text { a.e. in } \widetilde{\Omega}_{L}
$$

The existence and regularity of such function is detailed in Section 4.3.

- Let $\beta$ be a positive constant such that $|\nabla d| \leq \beta$ in $\Omega_{L}$. Then, if there exists $h$ such that the associated solution of (1.17), denoted by $u(x, t ; h)$, satisfies

$$
\left\|\nabla u(\cdot, \theta, h)-\nabla u_{b}\right\|_{L^{\infty}\left(\Omega_{L}\right)} \leq \frac{\delta}{2 \beta}
$$

it comes that $u(\cdot, \cdot ; h)$ satisfies (1.16).

- To do so, we will prove the following result: for any $u_{b}$ sufficiently regular, any $\varepsilon>0$ and any $\tau>0$, there exists $h \in L^{2}\left(\gamma_{L} \times(0, T)\right)$ such that

$$
\left\|u(\cdot, \tau ; h)-u_{b}\right\|_{H^{2}(\Omega)} \leq \varepsilon
$$

A more precise statement is given in Proposition 4.5. This allows to construct a sequence of controls $h_{n}$ such that

$$
\nabla u\left(x, \theta, h_{n}\right) \underset{n \rightarrow+\infty}{\longrightarrow} \nabla u_{b}(x), \quad \text { for a.e. } x \in \Omega
$$

and thus proves the previous item. To prove this result, first, we use classical results from control theory to establish this approximate controllability in $L^{2}(\Omega)$-norm. Then, using the regularity properties of system (1.1), we extend this result to the $H^{2}(\Omega)$-norm. To do so, we extend to unbounded domains the strategy given in [1]: after reaching approximately (in the weak norm) the target, let the system evolves freely to benefit from the regularizing properties but not too long to stay close to the target (in strong norms). We start this section recalling those regularization properties we will need.
Assume in the rest of this section that $c$ satisfies the assumptions given for $c_{2}$ in Theorem 1.2.

### 4.1 Analytic properties of the elliptic operator

First, notice that the $H^{2}(\Omega)$-norm we are interested in is related to the operator $A_{0}$ (see (1.14)). Indeed from classical elliptic results (see for instance [19, Theorem 3.1.1]), one has the following result.

Proposition 4.1. There exists a positive constant $C$ such that for any $u \in D\left(A_{0}\right)$,

$$
\begin{equation*}
\|u\|_{H^{2}(\Omega)} \leq C\left\|A_{0} u\right\|_{L^{2}(\Omega)} \tag{4.1}
\end{equation*}
$$

Then, we notice that the following propositions (see [19, Theorems 3.1.2 ii) and 3.1.3 ii)]), imply that $-A_{0}$ generates an analytic semigroup on $L^{2}(\Omega)$.
Proposition 4.2. There exists $\omega_{0} \in \mathbb{R}$ such that for any $\Re(\lambda) \geq \omega_{0}$, for any $f \in L^{2}(\Omega)$ the problem

$$
\left\{\begin{array}{l}
\lambda u-\nabla(c \nabla u)=f \text { in } \Omega \\
u=0 \text { on } \partial \Omega
\end{array}\right.
$$

has a unique solution $u \in H^{2}(\Omega)$ which continuously depends on $f$. Moreover, the resolvent set $\rho\left(A_{0}\right)$ satisfies $\left\{\lambda \in \mathbb{C}, \Re(\lambda) \geq \omega_{0}\right\} \subset \rho\left(A_{0}\right)$.

Proposition 4.3. There exist $\omega \geq \omega_{0}$ and $C>0$ such that if $\Re(\lambda) \geq \omega$, for any $u \in D\left(A_{0}\right)$, then

$$
|\lambda|\|u\|_{L^{2}(\Omega)} \leq C\left\|\lambda u-A_{0} u\right\|_{L^{2}(\Omega)} .
$$

From Proposition 4.2 and Proposition 4.3 we deduce that the operator $-A_{0}$ is sectorial and thus generates an analytic semigroup. Therefore (see for instance $[21,(6.7) \mathrm{p} .70]$ ) there exists a positive constant $C$ such that for all $t>0$ and for all $u \in D\left(A_{0}\right)$

$$
\begin{equation*}
\left\|A_{0} e^{-t A_{0}} u\right\|_{L^{2}(\Omega)} \leq \frac{C}{t}\|u\|_{L^{2}(\Omega)} . \tag{4.2}
\end{equation*}
$$

Finally, from (4.1) and (4.2) we get that there exists a positive constant $C$ such that for all $t>0$ and for all $u \in D\left(A_{0}\right)$

$$
\begin{equation*}
\left\|e^{-t A_{0}} u\right\|_{H^{2}(\Omega)} \leq \frac{C}{t}\|u\|_{L^{2}(\Omega)} \tag{4.3}
\end{equation*}
$$

This is the key inequality we will use.

### 4.2 Approximate controllability

First we give a classical result of approximate controllability in $L^{2}(\Omega)$.
Proposition 4.4. For any $u_{b} \in L^{2}(\Omega)$, any $u_{0} \in L^{2}(\Omega)$, any $\varepsilon>0$ and any $\tau>0$, there exists $h \in L^{2}\left(\gamma_{L} \times(0, \tau)\right)$ such that

$$
\left\|u(\cdot, \tau ; h)-u_{b}\right\|_{L^{2}(\Omega)} \leq \varepsilon
$$

Proof. This result is quite classical in control theory. Let us sketch its proof for the sake of completeness. Let $\widehat{\Omega}$ be an open set in $\mathbb{R}^{n}$ satisfying

$$
\Omega \subset \widehat{\Omega} \quad \text { and } \quad\left(\partial \Omega \cap\left(\mathbb{R}^{n} \backslash \partial \widehat{\Omega}\right)\right) \subset \gamma_{L}
$$

and let $\hat{\omega}$ be a domain such that $\hat{\omega} \subset \subset\left(\widehat{\Omega} \cap\left(\mathbb{R}^{n} \backslash \Omega\right)\right)$. From [26] and a classical duality argument, we have the approximate controllability in $L^{2}(\widehat{\Omega} \times(0, \tau))$ with a localized control $g \in L^{2}(\hat{\omega} \times(0, \tau))$ for the auxiliary problem

$$
\left\{\begin{array}{l}
\partial_{t} u-\nabla \cdot(c \nabla u)=1_{\hat{\omega}} g \text { in } \widehat{\Omega} \times(0, \tau) \\
u=0 \text { on } \partial \widehat{\Omega} \times(0, \tau) \\
u(., 0)=u_{0} \text { in } \widehat{\Omega}
\end{array}\right.
$$

We conclude the proof of this proposition by taking $h$ as the trace of $u$ on $\gamma_{L}$.
We now turn to approximate controllability in more regular norms.

Proposition 4.5. For any $u_{b} \in D\left(A_{0}^{2}\right)$, any $u_{0} \in L^{2}(\Omega)$, any $\varepsilon>0$ and any $\tau>0$, there exists $h \in L^{2}\left(\gamma_{L} \times(0, \tau)\right)$ such that

$$
\left\|u(\cdot, \tau ; h)-u_{b}\right\|_{H^{2}(\Omega)} \leq \varepsilon
$$

Proof. Let $u_{b} \in D\left(A_{0}^{2}\right)$ and $\varepsilon>0$. From Proposition 4.4, for any $\tau_{1} \in(0, \tau)$ and any $\delta>0$, there exists $h \in L^{2}\left(\gamma_{L} \times\left(0, \tau_{1}\right)\right)$ such that

$$
\left\|u\left(\cdot, \tau_{1} ; h\right)-u_{b}\right\|_{L^{2}(\Omega)} \leq \delta
$$

In what follows, we extend $h$ by 0 on $\left(\tau_{1}, \tau\right)$. Thus,

$$
u(\cdot, \tau ; h)=e^{-\left(\tau-\tau_{1}\right) A_{0}} u\left(\cdot, \tau_{1} ; h\right)
$$

First, the regularizing properties allows to obtain estimates in stronger norms. Indeed, from (4.3) we get

$$
\begin{align*}
\left\|u(\cdot, \tau ; h)-e^{-\left(\tau-\tau_{1}\right) A_{0}} u_{b}\right\|_{H^{2}(\Omega)} & =\left\|e^{-\left(\tau-\tau_{1}\right) A_{0}}\left(u\left(\cdot, \tau_{1} ; h\right)-u_{b}\right)\right\|_{H^{2}(\Omega)} \\
& \leq \frac{C}{\left(\tau-\tau_{1}\right)}\left\|u\left(\cdot, \tau_{1} ; h\right)-u_{b}\right\|_{L^{2}(\Omega)} \\
& \leq \frac{C}{\left(\tau-\tau_{1}\right)} \delta . \tag{4.4}
\end{align*}
$$

Second, we prove that for $\tau_{1}$ small enough, $e^{-\left(\tau-\tau_{1}\right) A_{0}} u_{b}$ is not far from $u_{b}$. Indeed, from [21, Theorem 6.13], we have

$$
\begin{align*}
\left\|A_{0}\left(e^{-\left(\tau-\tau_{1}\right) A_{0}} u_{b}-u_{b}\right)\right\|_{L^{2}(\Omega)} & =\left\|\left(e^{-\left(\tau-\tau_{1}\right) A_{0}}-I\right) A_{0} u_{b}\right\|_{L^{2}(\Omega)} \\
& =\left\|\int_{0}^{\tau-\tau_{1}} \frac{d}{d \eta}\left(e^{-\eta A_{0}}\right) A_{0} u_{b} \mathrm{~d} \eta\right\|_{L^{2}(\Omega)} \\
& =\left\|\int_{0}^{\tau-\tau_{1}} e^{-\eta A_{0}} A_{0}^{2} u_{b} \mathrm{~d} \eta\right\|_{L^{2}(\Omega)} \\
& \leq\left(\tau-\tau_{1}\right)\left\|A_{0}^{2} u_{b}\right\|_{L^{2}(\Omega)} \tag{4.5}
\end{align*}
$$

From (4.1) and (4.5), we deduce that

$$
\begin{equation*}
\left\|e^{-\left(\tau-\tau_{1}\right) A_{0}} u_{b}-u_{b}\right\|_{H^{2}(\Omega)} \leq C\left(\tau-\tau_{1}\right)\left\|A_{0}^{2} u_{b}\right\|_{L^{2}(\Omega)} \tag{4.6}
\end{equation*}
$$

So using the triangle inequality, estimates (4.4) and (4.6) we obtain

$$
\left\|u(\cdot, \tau ; h)-u_{b}\right\|_{H^{2}(\Omega)} \leq \frac{\mathrm{C}}{\left(\tau-\tau_{1}\right)} \delta+\mathrm{C}\left(\tau-\tau_{1}\right)\left\|A_{0}^{2} u_{b}\right\|_{L^{2}(\Omega)}
$$

for a specific constant C. Finally, choosing $\tau_{1}$ sufficiently close to $\tau$ such that

$$
\mathrm{C}\left(\tau-\tau_{1}\right)\left\|A_{0}^{2} u_{b}\right\|_{L^{2}(\Omega)} \leq \frac{\varepsilon}{2}
$$

then, choosing $\delta$ sufficiently small such that

$$
\frac{\mathrm{C}}{\left(\tau-\tau_{1}\right)} \delta \leq \frac{\varepsilon}{2}
$$

we obtain the desired inequality.

### 4.3 Construction of an appropriate target

In the above proof, to obtain approximate controllability in $H^{2}(\Omega)$-norms we need targets that are sufficiently regular, namely in $D\left(A_{0}^{2}\right)$. To apply our strategy, we now prove that there exists some function $u_{b} \in D\left(A_{0}^{2}\right)$ and a constant $\delta>0$ such that

$$
\left|\nabla d(x) \cdot \nabla u_{b}(x)\right| \geq \delta, \quad \text { for a.e. } x \in \widetilde{\Omega}_{L} .
$$

Let $\xi_{1} \in C^{\infty}(\mathbb{R} ; \mathbb{R})$ be such that $0 \leq \xi_{1} \leq 1$ and

$$
\begin{cases}\xi_{1}\left(x_{1}\right)=0, & \text { if }\left|x_{1}\right| \geq 2 L \\ \xi_{1}\left(x_{1}\right)=1, & \text { if }\left|x_{1}\right| \leq L\end{cases}
$$

Let $\xi^{\prime} \in C^{\infty}\left(\mathbb{R}^{n-1} ; \mathbb{R}\right)$ be such that $0 \leq \xi^{\prime} \leq 1$ and

$$
\begin{cases}\xi^{\prime}\left(x^{\prime}\right)=0, & \text { if } \operatorname{dist}\left(x^{\prime}, \mathbb{R}^{n-1} \backslash \omega\right) \leq \frac{r}{2} \\ \xi^{\prime}\left(x^{\prime}\right)=1, & \text { if } \operatorname{dist}\left(x^{\prime}, \mathbb{R}^{n-1} \backslash \omega\right) \geq r\end{cases}
$$

where $r$ is defined in (1.12). Let

$$
u_{b}(x)=d(x) \xi_{1}\left(x_{1}\right) \xi^{\prime}\left(x^{\prime}\right), \quad \forall x=\left(x_{1}, x^{\prime}\right) \in \Omega
$$

Thus, $u_{b} \in D\left(A_{0}\right)$ and $\nabla u_{b}(x)=\nabla d(x)$, for any $x \in \widetilde{\Omega}_{L}$. Due to (1.2), this implies that (1.16) is satisfied. As $u_{b} \in C^{\infty}(\Omega)$ and $u_{b}$ identically vanishes near the boundary we also obtain that $u_{b} \in D\left(A_{0}^{2}\right)$.
Remark 4.1. Finding $u_{b}$ in $C^{\infty}(\Omega)$ such that

$$
\inf _{\Omega_{L}}\left|\nabla u_{b} \cdot \nabla d\right|>0
$$

can easily be done taking for instance $u_{b}=d$. The main difficulty is to ensure all the boundary conditions so that $u_{b} \in D\left(A_{0}^{2}\right)$. This is precisely why we assumed that c is known in $V_{L}$. With this assumption, there is no requirement on $\nabla u_{b}$ near the boundary which allows to use cut-off functions to design $u_{b}$.

### 4.4 Conclusion

We now have all the ingredients to prove Theorem 1.2. Let $u_{b}$ be the target constructed in Section 4.3. Then, from Proposition 4.5, there exists $h$ depending on $c_{2}$ such that the solution $u_{2}$ satisfies (1.16).

As proved in Section 3, the stability estimate comes from the Carleman inequality applied to $y$ defined by (3.1). As, $u_{1}$ and $u_{2}$ are both solutions of (1.17) with the same boundary condition $h$ it comes that $y$ solves (1.1) that is with homogeneous Dirichlet boundary conditions. The rest of the proof remains unchanged.

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[^0]:    *Aix Marseille Univ, CNRS, Centrale Marseille, I2M, Marseille, France. email: laure.cardoulis@univ-amu.fr, michel.cristofol@univ-amu.fr, morgan.morancey@univ-amu.fr

[^1]:    ${ }^{1}$ Recall that $\epsilon$ is defined in Proposition 1.1, $M$ is defined in (1.13) and $\widetilde{M}, C_{m i n}$ are defined in (1.11).

[^2]:    ${ }^{2}$ Recall that $\epsilon$ is defined in Proposition 1.1 and $\widetilde{M}, C_{\text {min }}$ are defined in (1.11).
    ${ }^{3}$ Recall that $L, T, a$ are defined in Proposition 1.1 and $\widetilde{M}, C_{m i n}$ are defined in (1.11).

[^3]:    ${ }^{4}$ Recall that $L, T, a, \epsilon$ are defined in Proposition 1.1 and $\widetilde{M}, C_{m i n}$ are defined in (1.11).

[^4]:    ${ }^{5}$ Recall that $L, T, a$ are defined in Proposition 1.1, $M$ is defined in (1.13) and $\delta$ is defined in (1.16).

[^5]:    ${ }^{6}$ Recall that $T, \epsilon$ are defined in Proposition 1.1.

