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CONTRIBUTION TO INVERSE PROBLEMS AND CONTROLLABILITY ISSUES OF HYPERBOLIC AND PARABOLIC PARTIAL DIFFERENTIAL EQUATIONS

TESIS PARA OPTAR AL GRADO DE DOCTOR EN CIENCIAS DE LA INGENIERÍA, MENCIÓN MODELACIÓN MATEMÁTICA

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## CONTRIBUTION TO INVERSE PROBLEMS AND CONTROLLABILITY ISSUES OF HYPERBOLIC AND PARABOLIC PARTIAL DIFFERENTIAL EQUATIONS

El objetivo de esta tesis consiste principalmente en el estudio teórico de algunos resultados de problemas inversos y de controlabilidad en ecuaciones hiperbólicas y parabólicas.

En el Capítulo 1 presentamos una breve introducción de los tópicos tratados en este trabajo. Principalmente, centramos nuestra atención en las definiciones clásicas de controlabilidad y problemas inversos. Posteriormente, indicamos cuáles son los resultados generales obtenidos en esta tesis.

En el Capítulo 2, describimos los resultados de estabilidad obtenida para la reconstrucción de potenciales en un sistema de ecuaciones hiperbólicas acopladas en cascada. Para probar este resultado, nos inspiramos en el método de Bukhgeim-Klibanov combinado con un tipo especial de desigualdades conocidas como estimaciones de Carleman. Estas dos herramientas, junto con el hecho que las ecuaciones del sistema están acopladas en cascada, nos permiten obtener un resultado de estabilidad Lipschitz para la recuperación de todos los potenciales del sistema utilizando mediciones de algunas componentes accesibles de él.

En el Capítulo 3, nos centramos en el estudio de la controlabilidad a cero de una ecuación del calor con condiciones de borde dinámicas. Este problema se puede ver como una ecuación del calor acoplada con una ecuación diferencial ordinaria actuando en un extremo del borde. Nuestros resultados apuntan en dos direcciones. En primer lugar, probamos que este tipo de problemas se puede controlar a cero en una región que está lejos de la interacción entre las dos dinámicas. Usando la dualidad entre observabilidad y controlabilidad, la prueba de este resultado está basado en la construcción de una estimación de Carleman adecuada. En segundo lugar, probamos que una modificación de este tipo de problemas puede ser visto como el problema límite de una familia de problemas parabólicos con coeficientes de difusión discontinuos en donde la difusión es muy alta en una parte del dominio. Adicionalmente, estudiamos el efecto que tiene el control del problema límite en la sucesión de problemas aproximados.

Finalmente, en el Capítulo 4 desarrollamos una manera de obtener una estimación de tipo Carleman para una ecuación del calor con coeficientes de difusión discontinuos. La novedad en esta estrategia están basadas en las ideas del análisis microlocal desarrollado por L. Robbiano y J. Le Rousseau et al. para ecuaciones parabólicas, con la ventaja de que podemos obtener información de la constante de observabilidad.

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## CONTRIBUTION TO INVERSE PROBLEMS AND CONTROLLABILITY ISSUES OF HYPERBOLIC AND PARABOLIC PARTIAL DIFFERENTIAL EQUATIONS

The goal of this thesis is the theoretical study of some controllability and inverse problems for hyperbolic and parabolic equations.

In Chapter 1, we present a brief introduction of the general topics of this thesis. We focus on the classical definitions of controllability and inverse problems in partial differential equations. Then, we present the main results of this work.

In Chapter 2, we deal with the potential reconstruction for hyperbolic systems in cascade where measurements of the last component are not available. Roughly speaking, the novelty of this work consists in the Lipschitz stability of this inverse problem from partial measurements of the components of the system. More precisely, we measure all components except the last one. The main tool to achieve this result is a global Carleman estimate for a system of wave equations in cascade where the last component is not accessible.

In Chapter 3, the null controllability of a parabolic equation with dynamic boundary conditions is studied. This problem can be seen as a heat equation with an ordinary differential equation coupling through the boundary. We present our results in two directions. Firstly, we prove that these kind of problems are null-controllable at any time when control acts on a subset which is far from the coupling region. Following the well-known duality between controllability and observability, we prove the associated observability inequality for the adjoint system. Secondly, we prove that a slight modification of this problem can be seen as a limit of a family of parabolic equations with discontinuous diffusion coefficients where the diffusivity is very high in a part of the domain. Additionally, we study the effect of controls for the limit problem in the approximate system.

Finally, in Chapter 4 we develop a suitable Carleman estimate for the heat equation in the presence of an interface. The novelty in this strategy is based on the ideas of microlocal analysis by L. Robbiano and J. Le Rousseau in the context of parabolic equations, with the advantage that we can track the observability constant.

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## Chapter 1

## General Introduction

In this chapter, we state the elementary notions concerning inverse problems and controllability issues in partial differential equations (PDE's for short). In order to get a self-contained exposition, we divide this chapter into three sections.

Firstly, in Section 1.1 we restrict our attention to inverse problems for hyperbolic equations and systems. More precisely, we focus on stability results of potential reconstruction for this kind of equations where the observation is on a part of the boundary satisfying suitable geometric and time conditions. Moreover, some other inverse problems concerning hyperbolic equations are considered. In particular, we state results on the stability for some coefficients of an acoustic equation studied by M. Yamamoto and M. Bellassoued.

On the other hand, in Section 1.2 controllability results on parabolic equations are studied. Since the literature is very rich concerning this topic, we reduce our scope to the basic results for the heat equation and their variants. In addition, recent results of controllability for parabolic equations with dynamic boundary conditions are considered.

Finally, in Section 1.3 we state the main theoretical contributions of each topic in this thesis.

### 1.1 Inverse problems for hyperbolic equations

Intuitively, the observation of an effect in a physic phenomenon may not be sufficient to determine its cause. In fact, if we go inside a room and notice that the temperature is (approximately) uniform, it is difficult for us to know what the distribution of the temperature was four hours ago. Even more, one can think about if there is two different distributions of the temperature which provides the same observation.

The so-called direct problems in PDE's try to describe various physical phenomena such as the propagation of sound, heat, seismic waves, electromagnetic waves, etc. Here, the media properties, the initial state or its conditions on the boundary are assumed to be known.

For example, we can formulate the following direct problem for the acoustic equation:

In the domain $\Omega \subset \mathbb{R}^{d}$, with $d \geq 1$ with smooth boundary $\partial \Omega$, let $y=y(x, t)$ be a solution of

$$
\begin{cases}\rho(x) \partial_{t}^{2} y-c^{2}(x) \Delta y=f(x, t), & \text { in } \Omega \times(0, T)  \tag{1.1.1}\\ y(x, 0)=y_{0}, \partial_{t} y(x, 0)=y_{1}, & \text { in } \Omega \\ y=g, & \text { on } \partial \Omega \times(0, T)\end{cases}
$$

Here, $y=y(x, t)$ is the acoustic pressure, $\rho=\rho(x)$ and $c=c(x)$ stand for the density and sound speed of the medium and $f=f(x, t)$ is the source. Under suitable assumptions, system (1.1.1) is well posed, i.e., it has a unique solution and is stable with respect to small perturbations in the data.

Generally speaking, the unknowns in inverse problems include some functions given in the formulation of a direct problem, which are the solution of our inverse problem. In order to compute these unknowns, the direct problem is supplied with some additional information about the solution to the direct problem. This one represents the data of our inverse problem. For example, in (1.1.1) one can consider the partial data on the flux $\frac{\partial u}{\partial \nu}=h$ on $\partial \Omega \times(0, T)$.

There is not a universal definition for inverse problems. Indeed, given a direct problem (sometimes called forward problem), one can define several inverse problems. For example we say the inverse problem for (1.1.1) is a source term inverse problem if it is required to determine the function $f=f(x, t)$. In a similar way, we say that an inverse problem is a coefficient inverse problem if it is required to reconstruct the coefficients $c=c(x)$ and/or $\rho=\rho(x)$ in 1.1.1. There exist other classifications based on the additional information, on equations on the structure of the operator, etc. For more details about this topic, we refer to [68].

In contrast to direct problems, inverse problems are ill posed. Mathematically speaking, this means that this kind of problems has no solution in the desired class, or has several solutions, or the solution procedure is unstable, i.e., arbitrarily small errors in the data of the inverse problems may lead to indefinitely large errors in the solutions. For this reason, three questions arise naturally: uniqueness, stability and reconstruction of the coefficients studied. In this thesis, we focus only on uniqueness and stability issues.

Concerning uniqueness, we discuss whether the adopted extra data on the solution can uniquely determine an unknown coefficient or source term. On the other hand, in the stability issue, we are interested on getting the so-called stability estimates. Roughly speaking, these ones determine if it is possible to obtain the norm of the unknown coefficients by partial measurements. Of course, it follows from this that a stability result implies uniqueness. We refer to [16] and [65] for a complete description of these problems.

In general, concerning theoretical methods for coefficient inverse problems, we can consider two types of formulations:

- Infinitely many measurements by Dirichlet-to-Neumann map: In this case, the data are all the pairs of Dirichlet boundary inputs and the corresponding Neumann boundary values. For example, given $g$, we solve (1.1.1 with $f=0$ in
$\Omega \times(0, T), y(\cdot, 0)=\partial_{t} y(\cdot, 0)=0$ in $\Omega$ and $y=g$ on $\partial \Omega \times(0, T)$, we define the map

$$
\left.g \mapsto c \frac{\partial y}{\partial \nu}\right|_{\partial \Omega \times(0, T)},
$$

which is called the Dirichlet-to-Neumann map. Then, in this case, the problem is to determine $c$ from the Dirichlet-to-Neumann map, which means that we have to repeat measurements of $c \frac{\partial y}{\partial \nu}$ on $\partial \Omega \times(0, T)$ after choosing all possible $g$. For this reason, we say that this is an inverse problem with infinitely many measurements.

- Finitely many measurements by Carleman estimates: In contrast to the above formulation, in this one it is sufficient to observe boundary or distributed data of the solution after suitably choosing initial values at finitely many times or a single time. Concerning uniqueness and stability, in 1981 Bukhgeim and Klibanov [26] proposed a fundamental method to obtain uniqueness and stability of the inverse problem based on Global Carleman estimates.

There exists a huge literature on these topics. For more details in other contexts and equations, we refer to [65], [93], [16], and the references therein.

### 1.1.1 On the well-posedness of the wave equation with potential

In this section, we present the classical inverse problem related to the wave equation. In order to get an idea, let $\Omega \subset \mathbb{R}^{N}$ with $N \geq 1$ be a domain with smooth boundary $\partial \Omega$ and $T>0$. Then, let $u=u(x, t)$ be the solution of the following problem:

$$
\begin{cases}\partial_{t}^{2} u-\Delta u+p u=f, & \text { in } \Omega \times(0, T)  \tag{1.1.2}\\ u(x, 0)=u_{0}(x), \partial_{t} u(x, 0)=u_{1}(x), & \text { in } \Omega, \\ u=0, & \text { on } \partial \Omega \times(0, T)\end{cases}
$$

Here $u=u(x, t)$ denotes the evolution of the amplitude of the waves, $p=p(x, t)$ is a bounded potential and $f=f(x, t)$ denotes a source term acting according to the equation 1.1.2 $)_{1}$. Moreover, $\left(u_{0}, u_{1}\right)$ denotes the initial state of the waves. In addition, equation 1.1.2 $3_{3}$ is a Dirichlet boundary condition which states that the amplitude of the waves vanishes at the boundary.

The following result asserts that (1.1.2) is well posed in the sense of Hadamard.
Theorem 1.1 Suppose that $p \in L^{\infty}(\Omega \times(0, T))$, $f \in L^{1}\left(0, T ; L^{2}(\Omega)\right)$, $u_{0} \in H_{0}^{1}(\Omega)$ and $u_{1} \in L^{2}(\Omega)$. Then, the problem (1.1.2) admits a unique (weak) solution satisfying

$$
u \in C^{0}\left([0, T] ; H_{0}^{1}(\Omega)\right), \quad \partial_{t} u \in C^{0}\left([0, T] ; L^{2}(\Omega)\right)
$$

Moreover, there exists a constant $C=C(\Omega, T)>0$ such that the solution $u$ of (1.1.2) satisfies the following estimate:

$$
\begin{equation*}
\|u\|_{C^{0}\left([0, T] ; H_{0}^{1}(\Omega)\right)}+\left\|\partial_{t} u\right\|_{C^{0}\left([0, T] ; L^{2}(\Omega)\right)} \leq C\left(\left\|u_{0}\right\|_{H_{0}^{1}(\Omega)}+\left\|u_{1}\right\|_{L^{2}(\Omega)}+\|f\|_{L^{1}\left(0, T ; L^{2}(\Omega)\right)}\right) \tag{1.1.3}
\end{equation*}
$$

Notice that Theorem 1.1 gives us the regularity of the solution in the presence of a source term $f \in L^{1}\left(0, T ; L^{2}(\Omega)\right)$. Alternatively, one can establish the regularity of the solution $u$ when $f$ belongs in a different functional space:

Theorem 1.2 Let us assume that $p \in L^{\infty}(\Omega \times(0, T)), f \in W^{1,1}\left(0, T ; H^{-1}(\Omega)\right), u_{0} \in$ $H_{0}^{1}(\Omega)$ and $u_{1} \in L^{2}(\Omega)$. Then, there exists a unique solution $u$ of (1.1.2) with the following properties

$$
u \in C^{0}\left([0, T] ; H_{0}^{1}(\Omega)\right), \quad \partial_{t} u \in C^{0}\left([0, T] ; L^{2}(\Omega)\right)
$$

Moreover, there exists a positive constant $C$ which depends at least of $\Omega$ and $T>0$ such that the unique solution of (1.1.2) fulfills

$$
\begin{equation*}
\|u\|_{C^{0}\left([0, T] ; H_{0}^{1}(\Omega)\right)}+\left\|\partial_{t} u\right\|_{C^{0}\left([0, T] ; L^{2}(\Omega)\right)} \leq C\left(\left\|u_{0}\right\|_{H_{0}^{1}(\Omega)}+\left\|u_{1}\right\|_{L^{2}(\Omega)}+\|f\|_{W^{1,1}\left(0, T ; L^{2}(\Omega)\right)}\right) . \tag{1.1.4}
\end{equation*}
$$

Let us mention that the inequalities $(\sqrt{1.1 .3})$ and $(1.1 .4)$ assert the continuous dependence of the solution $u$ with respect to the initial data and source terms, see [80].

We remark that Theorems 1.1 or 1.2 do not provide information about the normal derivative $\partial_{\nu} u$ of the solution of 1.1 .2 . To be more precise, if a function belongs to $C^{0}\left([0, T] ; H_{0}^{1}(\Omega)\right)$ the normal derivative could not be well defined. However, under the assumptions of Theorem 1.1, we will see that the solution $u$ has an extra regularity. This is given in the following:

Theorem 1.3 (Hidden regularity of wave equation [79]) Assume the same hypotheses of Theorem 1.1. Then, the solution $u$ of (1.1.2) fulfills

$$
\partial_{\nu} u \in L^{2}(\partial \Omega \times(0, T)) .
$$

Furthermore, the application

$$
\Lambda: L^{1}\left(0, T ; L^{2}(\Omega)\right) \times H_{0}^{1}(\Omega) \times L^{2}(\Omega) \mapsto L^{2}(\partial \Omega \times(0, T))
$$

defined by $\Lambda\left(f, u_{0}, u_{1}\right)=\partial_{\nu} u$ is well defined and is a linear continuous map, that is to say, there exists a constant $C>0$ such that

$$
\left\|\partial_{\nu} u\right\|_{L^{2}(\partial \Omega \times(0, T))} \leq C\left(\|f\|_{L^{1}\left(0, T ; L^{2}(\Omega)\right)}+\left\|u_{0}\right\|_{H_{0}^{1}(\Omega)}+\left\|u_{1}\right\|_{L^{2}(\Omega)}\right) .
$$

Remark 1.4 Under the assumptions of Theorem 1.2 we cannot get any regularity result on the normal derivative $\partial_{\nu} u$.

### 1.1.2 Potential reconstruction for the wave equation

Now we will introduce an inverse problem for (1.1.2). Suppose that the potential $p$ in (1.1.2) is a time-independent function. Then, one can consider the following

Inverse problem: Can we retrieve the potential $p=p(x)$ of (1.1.2) from the knowledge of the flux $\partial_{\nu} u$ on $\Gamma_{0} \times(0, T)$ with $\Gamma_{0} \subset \partial \Omega$ or partial measurements of $u$ in $\omega \times(0, T)$ with $\omega \subset \Omega$ ?

Notice that the above problem is interesting since the partial data is defined only on a part of the boundary. Of course, we point out that this makes sense thanks to the hidden regularity result of the wave equation.

In the following, we are interested in the dependence of the solution $u$ of (1.1.2) with respect to the potential $p$. For this reason, here and subsequently, $u(p)$ and $u(q)$ stand for the corresponding solution of $u$ associated to the potentials $p$ and $q$ respectively, for fixed initial data $\left(u_{0}, u_{1}\right)$ and source term $f$.

Then, thanks to this notation, we can formulate questions related to the above inverse problem in three directions:

- Uniqueness: Does the equality

$$
\partial_{\nu} u(p)=\partial_{\nu} u(q) \text { on } \Gamma_{0} \times(0, T)
$$

imply $p=q$ in $\Omega$ ?

- Stability: Is it possible to estimate $\|q-p\|_{L^{2}(\Omega)}$ or better yet, a stronger norm of $(q-p)$, by a suitable norm of $\partial_{\nu} u(q)-\partial_{\nu} u(p)$ in $\Gamma_{0} \times(0, T)$ ?
- Reconstruction: Can we find a formula or an algorithm to retrieve the potential $p$ from the knowledge of $\partial_{\nu} u(p)$ on $\Gamma_{0} \times(0, T)$ ?

Of course, the same three questions can be formulated in the case of partial data of interior observations, i.e., $u(p)$ in $\omega \times(0, T)$, with $\omega \subset \Omega$.

Now we focus on the stability problem. To this end, we shall introduce the so-called geometric and time conditions:

- Geometric condition: there exists $x_{0} \notin \bar{\Omega}$ such that $\Gamma_{0} \subset \partial \Omega$ fulfills

$$
\begin{equation*}
\left\{x \in \partial \Omega ; \nu(x) \cdot\left(x-x_{0}\right) \geq 0\right\} \subset \overline{\Gamma_{0}} . \tag{1.1.5}
\end{equation*}
$$

- Time condition: $T$ is chosen such that $x_{0} \notin \bar{\Omega}$ given in the geometric condition satisfies:

$$
\begin{equation*}
\sup _{x \in \Omega}\left|x-x_{0}\right| \leq T . \tag{1.1.6}
\end{equation*}
$$

In the case of interior observations, the geometric condition reads as follows: there exists $x_{0} \notin \bar{\Omega}$ such that $\omega \subset \Omega$ satisfies

$$
\begin{equation*}
\left\{x \in \partial \Omega ; \nu(x) \cdot\left(x-x_{0}\right)>0\right\} \subset \partial \omega \cap \partial \Omega . \tag{1.1.7}
\end{equation*}
$$

Let us emphasize that $\Gamma_{0} \subset \partial \Omega$ and $\omega \subset \Omega$ satisfying the geometric condition are not arbitrary subsets. Indeed, in the particular case of $\Omega$ being a ball in $\mathbb{R}^{2}$, the length of $\Gamma_{0}$ is larger than half of the ball. In the same way, $\omega$ is a boundary neighborhood of $\Gamma_{0}$, see figure 1.1:


Figure 1.1: $\Gamma_{0}$ and $\omega$ satisfying the geometric condition
Roughly speaking, the Geometric and Time conditions assert that, thanks to the Snell law, all rays of geometric optics in $\Omega$, which are simply straight lines reflected on the boundary, should meet the observation region $\Gamma_{0}$ (or $\omega$ ) at a non-diffractive point in a time less than $T$.

Let us also introduce the admissible sets of potentials for the above inverse problem. For $m>0$, we define the set

$$
L_{\leq m}^{\infty}(\Omega)=\left\{p \in L^{\infty}(\Omega) ;\|p\|_{L^{\infty}(\Omega)} \leq m\right\}
$$

Theorem 1.5 (see [11]) Let $m>0, K>0$ and $r>0$. Let $p \in L_{\leq m}^{\infty}(\Omega)$. Assume that the solution $u(p)$ of (1.1.2) is such that

$$
\|u(p)\|_{H^{1}\left(0, T ; L^{\infty}(\Omega)\right)} \leq K
$$

and assume also that the initial datum $u_{0}$ satisfies the following positivity condition:

$$
\inf _{x \in \Omega}\left|u_{0}(x)\right| \geq r>0
$$

Additionally, suppose that $\Gamma_{0} \subset \partial \Omega$ and $T>0$ satisfy the geometric and time condition. Then, for all $q \in L_{\leq m}^{\infty}(\Omega), \partial_{\nu} u(q)-\partial_{\nu} u(p)$ belongs to $H^{1}\left(0, T ; L^{2}\left(\Gamma_{0}\right)\right)$ and there exists a constant $C=C(m, T, K, r)>0$ such that for any $q \in L_{\leq m}^{\infty}(\Omega)$, the following inequalities hold:

$$
\begin{equation*}
\left\|\partial_{\nu} u(q)-\partial_{\nu} u(p)\right\|_{H^{1}\left(0, T ; L^{2}\left(\Gamma_{0}\right)\right)} \leq C\|q-p\|_{L^{2}(\Omega)} \tag{1.1.8}
\end{equation*}
$$

and

$$
\begin{equation*}
\|q-p\|_{L^{2}(\Omega)} \leq C\|u(q)-u(p)\|_{H^{1}\left(0, T ; L^{2}\left(\Gamma_{0}\right)\right)} . \tag{1.1.9}
\end{equation*}
$$

We emphasize that the estimate (1.1.9) asserts the Lipschitz stability of the inverse problem while (1.1.8) gives the continuous dependence of the normal derivative in the norm $H^{1}\left(0, T ; L^{2}\left(\Gamma_{0}\right)\right)$ of the solution with respect to the potentials in the $L^{2}(\Omega)$-norm.

We remark that the geometric condition is a restrictive assumption on the observations. Without this condition, M. Bellassoued proved in [15] that the same inverse problem has logarithmic stability. More precisely, the author achieved the following result:

Theorem 1.6 (see [15], [16]) Let $\tilde{\Gamma}$ be an arbitrary subset of $\partial \Omega$. Assume that $u_{0} \in$ $H^{3}(\Omega) \cap H_{0}^{1}(\Omega)$ and $u_{1} \in H^{2}(\Omega)$, and there exists a constant $m_{0}>0$ such that

$$
\left|u_{0}(x)\right| \geq m_{0}>0, \quad \forall x \in \overline{\Omega \backslash \omega} .
$$

Then, there exist $T>0$ sufficiently large and a constant $C>0$ such that

$$
\begin{equation*}
\|p-q\|_{L^{2}(\Omega)} \leq C\left[\log \left(1+\frac{C}{\left\|\partial_{\nu}\left(u_{p}-u_{q}\right)\right\|_{H^{1}\left(0, T ; L^{2}(\tilde{\Gamma})\right)}}\right)\right]^{-1 / 2}, \tag{1.1.10}
\end{equation*}
$$

for all $p, q \in \Lambda(\Omega)$, where $\Lambda(\Omega)$ is the set of admissible potentials given by

$$
\Lambda(\Omega)=\left\{p \in W^{1, \infty}(\Omega) ;\|p\|_{W^{1, \infty}(\Omega)} \leq M, p=p_{0} \text { in } \Omega \backslash \omega\right\}
$$

with $M>0$ and $p_{0} \in C^{\infty}\left(\mathbb{R}^{n}\right)$ are given arbitrarily. Moreover, the constant $C$ in 1.1.10 is dependent of $\Omega, \omega, T, M, u_{0}, u_{1}$, but independent of $p, q$.

According to the definition of $\Lambda(\Omega)$, we remark that this stability result comes from the fact that $p$ is known in a part of $\Omega$. Let us also mention that the main ingredient of the proof of the above theorem is the Fourier-Bros-Iagolnitzer (FBI) transform. This one is crucially used in order to prove a sharp unique continuation property for hyperbolic equations (see [87], [86]).

The proof of the Theorem 1.5 is based on the Bukhgeim-Klibanov's method and the so-called global Carleman inequalities, which is an interesting result by itself. In order to state this result, we shall introduce some weight functions. For $\beta \in(0,1)$, we define

$$
\begin{equation*}
\psi(x, t)=\left|x-x_{0}\right|^{2}-\beta t^{2}+C_{0}, \quad \varphi(x, t)=e^{\lambda \psi(x, t)}, \quad \forall(x, t) \in \Omega \times(-T, T) \tag{1.1.11}
\end{equation*}
$$

Then, the Carleman estimate for the wave operator reads as follows:
Theorem 1.7 (see [11], [96]) Let us assume the geometric and time conditions. Let $\psi$ and $\varphi$ weight functions defined by (1.1.11. Then, there exist three positive constants $C$, $\lambda_{0}$ and $s_{0}$ such that for all $\lambda \geq \lambda_{0}$ and $s \geq s_{0}$, the following inequality holds

$$
\begin{aligned}
& \quad s \lambda \int_{-T}^{T} \int_{\Omega} e^{2 s \varphi} \varphi\left(\left|\partial_{t} v\right|^{2}+|\nabla v|^{2}\right) d x d t+s^{3} \lambda^{3} \int_{-T}^{T} \int_{\Omega} e^{2 s \varphi} \varphi^{3}|v|^{2} d x d t \\
& \quad+\int_{-T}^{T} \int_{\Omega} e^{2 s \varphi}\left|P_{1}\left(e^{s \varphi} v\right)\right|^{2} d x d t \\
& \leq \\
& \leq \int_{-T}^{T} \int_{\Omega} e^{2 s \varphi}|\square v|^{2} d x d t+C s \lambda \int_{-T}^{T} \int_{\Gamma_{0}} e^{2 s \varphi} \varphi\left|\partial_{\nu} v\right|^{2} d \sigma d t,
\end{aligned}
$$

for all $v \in L^{2}\left(-T, T ; H_{0}^{1}(\Omega)\right)$ satisfying $\square v:=\partial_{t}^{2} v-\Delta v \in L^{2}(\Omega \times(-T, T))$ with normal derivative $\partial_{\nu} v \in L^{2}\left(\Gamma_{0} \times(-T, T)\right), v( \pm T)=\partial_{t} v( \pm T)=0$ in $\Omega$ and $P_{1}$ defined by

$$
P_{1} w=\partial_{t}^{2} w-\Delta w+s^{2} \lambda^{2} \varphi^{2} w\left(\left|\partial_{t} \psi\right|^{2}-|\nabla \psi|^{2}\right) .
$$

Theorem 1.7 ensures the existence of a 2-parameter Carleman estimate (i.e. $\lambda$ and $s$ ) for the wave operator $\square=\partial_{t}^{2}-\Delta$. As we shall see in Section 2 and 3, often we just need a one-parameter Carleman estimate for our purposes. On the other hand, Carleman estimates provide another interesting result called Unique Continuation Property (UCP for short). Then, thanks to Theorem 1.7, we obtain the following:

Corollary 1.8 (UCP for the wave operator) Suppose that $v \in L^{2}\left(-T, T ; H_{0}^{1}(\Omega)\right)$ is a function which verifies $\square v=0$ in $\Omega \times(-T, T), v( \pm T)=\partial_{t} v( \pm T)=0$ in $\Omega$, and additionally $\partial_{\nu} u=0$ on $\Gamma_{0} \times(-T, T)$, where $\Gamma_{0}$ and $T>0$ satisfy the geometric and time conditions 1.1.5 and 1.1.6, respectively. Then, v vanishes in $\Omega \times(-T, T)$.

### 1.1.3 Inverse problem for the wave speed of the wave equation

Let $\Omega$ be a bounded domain with smooth boundary ( $C^{2}$ at least) and $T>0$. Let $u=u(x, t)$ be the solution of

$$
\begin{cases}\partial_{t}^{2} u-\operatorname{div}(p(x) \nabla u)=0, & \text { in } \Omega \times(0, T)  \tag{1.1.12}\\ u(\cdot, 0)=a, \partial_{t} u(\cdot, 0)=0, & \text { in } \Omega, \\ u=b, & \text { on } \partial \Omega \times(0, T)\end{cases}
$$

Here, $p$ denotes the bulk-modulus of the acoustic equation considered in a non homogeneous medium. Under smooth assumptions on $a, b$ and $p$ the problem has a unique (weak) solution. We consider the following:

Inverse problem: Determine the coefficient $p$ from the knowledge of partial measurements of $u$ in $\omega \times(0, T)$, with $\omega \subset \Omega$.

In order to formulate the results, we consider $\omega \subset \Omega$ satisfying the geometric condition (1.1.5). In addition, set

$$
\begin{equation*}
\mathcal{D}:=\sqrt{\sup _{x \in \Omega}\left|x-x_{0}\right|^{2}-\inf _{x \in \Omega}\left|x-x_{0}\right|^{2}} \tag{1.1.13}
\end{equation*}
$$

Given $\eta \in C^{1}(\partial \Omega)$ and constants $M_{0} \geq 0, M_{1}>0,0<\theta_{0} \leq 1$ and $\theta_{1}>0$, we define the following admissible sets:

$$
\begin{align*}
\mathcal{U}_{1}=\{ & p \in C^{2}(\bar{\Omega}) ;\|p\|_{C^{2}(\bar{\Omega})}<M_{1},\|\nabla p\|_{C(\bar{\Omega})}<M_{0}, p(x)>\theta_{1}, \forall x \in \bar{\Omega} \\
& \left.\left|\frac{\nabla p(x) \cdot\left(x-x_{0}\right)}{2 p(x)}\right|<1-\theta_{0}, \forall x \in \overline{\Omega \backslash \omega},\|u\|_{W^{4, \infty}(\Omega \times(0, T))}<M_{1}\right\}, \tag{1.1.14}
\end{align*}
$$

and

$$
\begin{align*}
\mathcal{U}_{2}=\{ & p \in C^{2}(\bar{\Omega}) ;\|p\|_{C^{2}(\bar{\Omega})}<M_{1},\|\nabla p\|_{C(\bar{\Omega})}<M_{0}, p(x)>\theta_{1}, \forall x \in \bar{\Omega} \\
& \left.\frac{\nabla p(x) \cdot\left(x-x_{0}\right)}{2 p(x)}<1-\theta_{0}, \forall x \in \overline{\Omega \backslash \omega},\|u\|_{W^{4, \infty}(\Omega \times(0, T))}<M_{1}\right\} . \tag{1.1.15}
\end{align*}
$$

It is possible to replace $<$ and $>$ by $\leq$ and $\geq$ respectively. We choose $\beta_{1}>0$ and $\beta_{2}>0$ such that

$$
\begin{equation*}
\beta_{1}+\frac{M_{0} \mathcal{D}}{\sqrt{\theta_{1}}} \sqrt{\beta_{1}}<\theta_{0} \theta_{1} \tag{1.1.16}
\end{equation*}
$$

and

$$
\begin{equation*}
\beta_{2}+\frac{M_{0} \mathcal{D}}{\sqrt{\theta_{1}}} \sqrt{\beta_{2}}<\theta_{0} \theta_{1}, \quad \theta_{1} \inf _{x \in \Omega}\left|x-x_{0}\right|^{2}-\beta_{2} \mathcal{D}^{2}>0 \tag{1.1.17}
\end{equation*}
$$

Now we have all the ingredients to state a stability result for the inverse problem with interior observations:

Theorem 1.9 Let us consider $\omega$ and $x_{0}$ satisfying the geometric condition. Let $a \in$ $W^{1, \infty}(\Omega)$ such that

$$
\nabla a(x) \cdot\left(x-x_{0}\right) \neq 0, \quad \forall x \in \bar{\Omega}
$$

Let $k=1$ or $k=2$. We choose the observation time $T>0$ such that

$$
\begin{equation*}
T>\frac{1}{\sqrt{\beta_{k}}} \mathcal{D} . \tag{1.1.18}
\end{equation*}
$$

Then, there exist constants $\kappa \in(0,1)$ and $C>0$ such that for all $(p, q) \in \mathcal{U}_{k}$, the associated solutions $u(p)$ and $u(q)$ fulfill the following inequality:

$$
\begin{equation*}
\|p-q\|_{L^{2}(\Omega)} \leq C\left(\sum_{j=2}^{3}\left\|\partial_{t}^{j}(u(p)-u(q))\right\|_{L^{2}(\omega \times(0, T)}\right) . \tag{1.1.19}
\end{equation*}
$$

The main tool to prove Theorem 1.9 is a Carleman estimate in $H^{-1}(\Omega \times(0, T))$. In order to state this result, let us recall that $x_{0}$ is defined by the Geometric condition 1.1.7) and $\beta_{1}$ and $\beta_{2}$ given by (1.1.16) and 1.1.17, respectively. We define the functions $\psi_{k}=\psi_{k}(x, t)$ and $\varphi=\varphi_{k}(x, t)$ by

$$
\begin{equation*}
\psi_{k}(x, t)=\left|x-x_{0}\right|^{2}-\beta_{k} t^{2}, \tag{1.1.20}
\end{equation*}
$$

and

$$
\begin{equation*}
\varphi_{k}(x, t)=e^{\lambda \psi_{k}(x, t)}, \quad \lambda>0 . \tag{1.1.21}
\end{equation*}
$$

Now we are ready to state the Carleman estimate with source lying in $H^{-1}(\Omega \times$ $(-T, T))$ :

Theorem 1.10 Let $k \in\{1,2\}$. We assume that $p \in \mathcal{U}_{k}, x_{0}$ and $\omega$ satisfying 1.1.7). Let $y \in H^{1}(\Omega \times(-T, T))$ satisfy

$$
\begin{cases}\partial_{t}^{2} y-\operatorname{div}(p(x) \nabla y)=\tilde{f}+\partial_{t} f_{0}+\sum_{j=1}^{N} \partial_{j} f_{j}, & \text { in } \Omega \times\left(-T_{k}, T_{k}\right),  \tag{1.1.22}\\ y\left(\cdot, \pm T_{k}\right)=0, & \text { in } \Omega, \\ y=0, & \text { on } \partial \Omega \times\left(-T_{k}, T_{k}\right),\end{cases}
$$

where $\tilde{f} \in H^{-1}\left(\Omega \times\left(-T_{k}, T_{k}\right)\right), f_{j} \in L^{2}\left(\Omega \times\left(-T_{k}, T_{k}\right)\right)$, with $0 \leq j \leq N$. Then, there exists a constant $\mu>0$ such that for each $T_{k}$ satisfying

$$
T_{k} \in\left(\frac{\mathcal{D}}{\sqrt{\beta_{k}}}, \frac{\mathcal{D}}{\sqrt{\beta_{k}}}+\mu\right)
$$

there exists $\lambda_{0}>0$ such that there exist constants $s_{0}=s_{0}(\lambda)>0$ and $C_{1}>0$ such that

$$
\begin{align*}
s \int_{-T_{k}}^{T_{k}} \int_{\Omega} e^{2 s \varphi_{k}}|y|^{2} d x d t \leq & C_{1} s \int_{-T_{k}}^{T_{k}} \int_{\Omega} e^{2 s \varphi_{k}}\left|f_{j}\right|^{2} d x d t+C_{1}\left\|e^{s \varphi_{k}} \tilde{f}\right\|_{H^{-1}\left(\Omega \times\left(-T_{k}, T_{k}\right)\right)}^{2} \\
& +C_{1} s \int_{-T_{k}}^{T_{k}} \int_{\omega} e^{2 s \varphi_{k}}|y|^{2} d x d t \tag{1.1.23}
\end{align*}
$$

for all $s \geq s_{0}$.

### 1.2 Controllability issues in PDE's

Control theory is the area of mathematics concerning dynamical systems whose behavior can be changed by means of controls applied through actuators. This is also a rich interdisciplinary branch of mathematics, with applications in areas such as biology, chemistry, engineering, economics and seismic prospection. For more details about this theory, we refer to the books [79], [91] and the reviews [48], [22] and the references therein.

Roughly speaking, a control system can be written in the following abstract form

$$
\left\{\begin{array}{l}
\frac{d y}{d t}=L(y, u), \quad 0<t<T  \tag{1.2.1}\\
y(0)=y_{0}
\end{array}\right.
$$

where $y \in Y$ and $u \in \mathcal{U}_{a d}$. Here $y$ is the state, the unknown of the problem that we want to control, $y_{0}$ is the initial state, $u$ is the control, the variable that can be chosen appropriately to act on the system and $\mathcal{U}_{a d}$, and $Y$ stand for the set of admissible controls and the state space, respectively.

Given a control system like 1.2.1), we can formulate the so-called controllabiliy problem, which can be stated as follows:

Controllability problem: find a control $u \in \mathcal{U}_{a d}$ such that the associated state behaves in a appropriate manner in a given final time $T>0$.

We distinguish four different notions of controllability. We say that the system 1.2.1) is approximately controllable if, for any initial state $y_{0}$, it is possible to steer the solution to a state arbitrarily close to any given target (in an appropriate topology). The exact controlability of (1.2.1) asserts that the system can be driven from any initial data to a prescribed target. On the other hand, we say that the system (1.2.1) has the nullcontrollability property if, for any initial data, the solution can be driven to zero. Finally, another interesting concept of controllability is the exact controllability to trajectories, which means that it is possible to steer the state of the system to join a control-free prescribed trajectory, i.e., a given solution of the system without control.

From a mathematical viewpoint, the literature is very rich on controllability problems, see for instance [92], [32, [41, [5] and the references therein. The control theory started to be developed in the beginning of the 1960's for finite dimensional systems. The linear case of this problem is by now completely understood thanks to the Kalman rank condition, and moreover the four notions of controllability introduced before are equivalent. Furthermore, the case of nonlinear finite dimensional systems has been intensely studied in the last two decades and there are many powerful sufficient conditions for local and global controllability, see [32].

Nevertheless, in the context of PDE's the situation is more delicate, even in the linear case. The main reason is that a linear PDE governing the evolution of a process may be of hyperbolic type (wave equation, Maxwell equations), of dispersive type (plate equation, Schrödinger equations, Korteweg-de Vries equation), or of parabolic type (heat equation, Stokes equation). Each equation induces specific properties on the trajectories: propagation of singularities with finite velocity for hyperbolic equations, infinite speed propagation property together with a weak (resp. strong) smoothing effect for dispersive (resp. parabolic) equations, and time irreversibility for parabolic equations.

Accordingly to the above description of the evolution of a linear PDE, we cannot expect equivalence between the different notions of controllability in general. For instance, the regularizing effect of the heat equation asserts that the associated solution of a $L^{2}$ initial state is a smooth function. Thus, it is difficult to ensure exact controllability for the heat equation when the control acts in a small part of the domain. On the other hand, the location and the duration of the control play an important role in the controllability of the wave equation. This role may be completely hidden in the finite dimension setting.

Generally, the study of controllability of linear PDE's is equivalent to a suitable observability inequality for the adjoint problem. This means that we need full knowledge of the solution of the adjoint problem at a given time using only local measurements of it. Nevertheless, we emphasize that the proof of such inequalities are a challenging issue and requires tools such as Ingham inequalities [64] [72], multiplier methods [71] [58], [79], [84], microlocal analysis [10], [77] [27], or Carleman estimates [61, [50], [48, [45].

### 1.2.1 Classical results on controllability of parabolic equations

In this subsection, we follow the presentation given in [83]. Let $\Omega$ be a bounded open set with boundary of class $C^{2}$ and $\omega \subset \Omega$ a non-empty open subset of $\Omega$. Given $T>0$ we consider the following non-homogeneous heat equation:

$$
\begin{cases}\partial_{t} u-\Delta u=\chi_{\omega} f, & \text { in } \Omega \times(0, T),  \tag{1.2.2}\\ u(\cdot, 0)=u_{0}, & \text { in } \Omega, \\ u=0, & \text { on } \partial \Omega \times(0, T)\end{cases}
$$

In (1.2.2), $u=u(x, t)$ is the state and $f=f(x, t)$ is the control function with a support localized in $\omega$.

Theorem 1.11 For any $f \in L^{2}((0, T) \times \omega)$ and $u_{0} \in L^{2}(\Omega)$ problem 1.2.2 has a unique
weak solution $u \in C\left([0, T] ; L^{2}(\Omega)\right)$ given by the variation of constants formula

$$
u(t)=S(t) u_{0}+\int_{0}^{t} S(t-s) \chi_{\omega} f(s) d s
$$

where $(S(t))_{t \geq 0}$ is the semigroup of constractions generated by $-\Delta$ in $L^{2}(\Omega)$.
Moreover, if $f \in W^{1,1}\left(0, T ; L^{2}(\omega)\right)$ and $u_{0} \in H^{2}(\Omega) \cap H_{0}^{1}(\Omega)$, problem 1.2.2 has a classical solution

$$
u \in C^{1}\left([0, T] ; L^{2}(\Omega)\right) \cap C\left([0, T] ; H^{2}(\Omega) \cap H_{0}^{1}(\Omega)\right)
$$

and (1.2.2) is verified in $L^{2}(\Omega)$ for all $t>0$.
From the fact that $u \in L^{\infty}\left(0, T ; L^{2}(\Omega)\right)$, it follows that $u \in L^{2}\left(0, T ; H_{0}^{1}(\Omega)\right)$. Consequently, whenever $u_{0} \in L^{2}(\Omega)$ and $f \in L^{2}\left(0, T ; L^{2}(\omega)\right)$ the solution verifies

$$
u \in L^{\infty}\left(0, T ; L^{2}(\Omega)\right) \cap L^{2}\left(0, T ; H_{0}^{1}(\Omega)\right)
$$

and we have the following energy estimate:

$$
\|u(\cdot, t)\|_{L^{2}(\Omega)}^{2}+\int_{0}^{t}\|\nabla u\|_{L^{2}(\Omega)} d t \leq C \int_{0}^{t}\|f\|_{L^{2}(\Omega)}^{2} d t+C\left\|u_{0}\right\|_{L^{2}(\Omega)}^{2}
$$

Now, we focus on the null-controllability results on the heat equation with Dirichlet boundary conditions where $\omega \subset \mathbb{R}$ is arbitrary. More precisely, we wonder if for each $T>0$ and $u_{0} \in L^{2}(\Omega)$ in 1.2 .2 , there exists a control $f \in L^{2}(\omega \times(0, T))$ such that

$$
u(T)=0, \quad \text { in } \Omega
$$

First, we notice that one of the most important properties of the heat equation is its regularizing effect. When $\Omega \backslash \omega \neq \emptyset$, the solutions of 1.2 .2 belong to $C^{\infty}(\Omega \backslash \omega)$ at time $t=T$. Hence, the restriction of the elements of $R\left(T ; u_{0}\right)$ to $\Omega \backslash \omega$ are smooth functions. Then, unless the trivial case $\omega=\Omega$, that is to say, when the control function acts on the entire domain $\Omega$, exact controllability may not hold. In this sense, the notion of exact controllability is not very relevant for the heat equation. This is due to its strong time irreversibility of the system under consideration.

Moreover, it is not difficult to see that if null controllability holds, then any initial data may be let to any final state of the form $S(T) v_{0}$ with $v_{0} \in L^{2}(\Omega)$, i.e., to the range of the semigroup in time $t=T$. Null controllability implies approximate controllability. Indeed, this is a consequence that the eigenfunctions of the laplacian operator belong to $S(T)\left[L^{2}(\Omega)\right]$. Then we deduce that $R\left(T ; u_{0}\right)$ is dense in $L^{2}(\Omega)$, which is the definition of approximate controllability.

On the other hand, approximate controllability together with uniform estimates on the approximate controls as $\varepsilon$ goes to zero may lead to null controllability properties. More precisely, given $u^{1}$, we have that $u^{1} \in R\left(T ; u_{0}\right)$ if and only if there exists a sequence of controls $\left(f_{\varepsilon}\right)_{\varepsilon>0}$ such that

$$
\left\|u(T, \cdot)-u_{1}\right\|_{L^{2}(\Omega)} \leq \varepsilon
$$

and $\left(f_{\varepsilon}\right)_{\varepsilon>0}$ is bounded in $L^{2}((0, T) \times \omega)$. Indeed, in this case any weak limit in $L^{2}(\omega \times$ $(0, T))$ of the sequence of controls $\left(f_{\varepsilon}\right)_{\varepsilon>0}$ gives an exact control which makes that $u(\cdot, T)=$ $u_{1}$ in $\Omega$.

### 1.2.2 Null controllability of the heat equation for parabolic equations with dynamic boundary conditions

In this part, we follow the presentation of [81]. Let $\Omega \subset \mathbb{R}^{N}$ be a bounded domain with smooth boundary $\Gamma=\partial \Omega, N \geq 2$ and $T>0$. Let $y=y(x, t)$ be the solution of the following problem

$$
\begin{cases}\partial_{t} y-d \Delta y+a(x, t) y=\chi_{\omega} v(x, t), & \text { in } \Omega \times(0, T),  \tag{1.2.3}\\ \partial_{t} y_{\Gamma}-\Delta_{\Gamma} y_{\Gamma}+d \partial_{\nu} y+b(x, t) y_{\Gamma}=0, & \text { on } \Gamma \times(0, T), \\ y_{\Gamma}(x, t)=\left.y\right|_{\Gamma}(x, t), & \text { on } \Gamma \times(0, T), \\ \left.\left(y, y_{\Gamma}\right)\right|_{t=0}=\left(y_{0}, y_{0, \Gamma}\right), & \text { in } \Omega \times \Gamma .\end{cases}
$$

Here, $\omega \subset \subset \Omega$ is an arbitrary nonempty open subset of $\Omega, y_{0} \in L^{2}(\Omega)$ and $y_{0, \Gamma} \in L^{2}(\Gamma)$ are the initial data, the constants $\delta, d$ are positive, $a \in L^{\infty}(\Omega \times(0, T))$ and $b \in L^{\infty}(\Gamma \times$ $(0, T))$. In addition, $\left.y\right|_{\Gamma}$ denotes the trace of a function $y: \Omega \rightarrow \mathbb{R}, \nu$ is the outer unit normal field, $\partial_{\nu} y=\left.(\nu \cdot \nabla y)\right|_{\Gamma}$ stands for the normal derivative at $\Gamma$, and $\Delta_{\Gamma}$ denotes the Laplace-Beltrami operator on $\Gamma$.

Then, the main question is: find a control $v \in L^{2}(\omega \times(0, T))$ such that the solution $y$ of (1.2.3) satisfies

$$
y(\cdot, T)=0, \text { in } \bar{\Omega} .
$$

Following the classical equivalence between controllability and observability, we introduce the following adjoint system:

$$
\begin{cases}-\partial_{t} \varphi-d \Delta \varphi+a(x, t) \varphi=0, & \text { in } \Omega \times(0, T),  \tag{1.2.4}\\ -\partial_{t} \varphi_{\Gamma}-\delta \Delta_{\Gamma} \varphi+d \partial_{\nu} \varphi+b(x, t) \varphi_{\Gamma}=0, & \text { on } \Gamma \times(0, T), \\ \varphi_{\Gamma}(x, t)=\gamma_{\Gamma}(x, t), & \text { on } \Gamma \times(0, T), \\ (\varphi(\cdot, T), \varphi(\cdot, T))=\left(\varphi_{T}, \varphi_{T, \Gamma}\right), & \text { in } \Omega \times \Gamma .\end{cases}
$$

Proposition 1.12 There is a constant $C>0$ such that for all $\left(\varphi_{T}, \varphi_{T, \Gamma}\right) \in L^{2}(\Omega) \times L^{2}(\Gamma)$ the mild solution $\left(\varphi, \varphi_{\Gamma}\right)$ of the backward problem (1.2.4) satisfies

$$
\begin{equation*}
\|\varphi(\cdot, 0)\|^{2}+\left\|\varphi_{\Gamma}(\cdot, 0)\right\|^{2} \leq C \int_{0}^{T} \int_{\omega}|\varphi|^{2} d x d t \tag{1.2.5}
\end{equation*}
$$

Given $R>0$, the constant $C=C(R)$ can be chosen independently of $a$, $b$ with

$$
\|a\|_{\infty},\|b\|_{\infty} \leq R
$$

To prove Proposition 1.12, the authors prove a Carleman estimate for the problem 1.2.4). Let us emphasize that weights appearing in such estimate are the same in [48] for the case Dirichlet boundary conditions and in the classical text of A. V. Fursikov and O. Yu. Imanuvilov [50] for mixed boundary conditions. Of course, such functions are based on an auxiliary function $\eta^{0}$ whose existence is guaranteed in the following result:

Lemma 1.13 Given a nonempty open set $\omega^{\prime} \subset \subset \Omega$, there is a function $\eta^{0} \in C^{2}(\bar{\Omega})$ such that

$$
\eta^{0}>0 \text { in } \Omega, \quad \eta^{0}=0, \text { on } \Gamma, \quad\left|\nabla \eta^{0}\right|>0, \text { in } \overline{\Omega \backslash \omega^{\prime}} .
$$

We emphasize that, since $\left|\nabla \eta^{0}\right|^{2}=\left|\nabla_{\Gamma} \eta^{0}\right|^{2}+\left|\partial_{\nu} \eta^{0}\right|^{2}$ on $\Gamma$, the function $\eta^{0}$ in the above lemma satisfies

$$
\nabla_{\Gamma} \eta^{0}=0, \quad\left|\nabla \eta^{0}\right|=\left|\partial_{\nu} \eta^{0}\right|, \quad \partial_{\nu} \eta^{0} \leq-c<0, \text { on } \Gamma,
$$

for some constant $c>0$. Now let us define the Carleman weight functions. For $\lambda, m>1$, we set

$$
\begin{aligned}
& \alpha(x, t)=(t(T-t))^{-1}\left(e^{2 \lambda m\left\|\eta^{0}\right\|_{\infty}}-e^{\lambda\left(m\left\|\eta^{0}\right\|_{\infty}+\eta^{0}(x)\right)}\right), \\
& \xi(x, t)=(t(T-t))^{-1} e^{\lambda\left(m\left\|\eta^{0}\right\|_{\infty}+\eta^{0}(x)\right)},
\end{aligned}
$$

for $x \in \bar{\Omega}$ and $t \in(0, T)$. Notice that $\alpha$ and $\xi$ are smooth and strictly positive on $\bar{\Omega} \times(0, T)$ and blow up as $t \rightarrow 0$ and as $t \rightarrow T$. Moreover, such functions are constant on the boundary $\Gamma$ so that

$$
\nabla_{\Gamma} \alpha=0, \quad \text { and } \nabla_{\Gamma} \xi=0, \quad \text { on } \Gamma .
$$

Lemma 1.14 Let $T>0, \omega \subset \subset \Omega$ be a nonempty and open subset of $\Omega, d, \delta>0$, $a \in L^{\infty}(\Omega \times(0, T))$ and $b \in L^{\infty}\left(\Gamma_{T}\right)$. Let $\omega^{\prime} \subset \subset \omega$. Define $\eta^{0}$, $\alpha$ and $\xi$ as above with respect to $\omega^{\prime}$. Then, there exist constants $C>0, \lambda_{1} \geq 1$ and $s_{1} \geq 1$ such that

$$
\begin{align*}
& s^{-1} \int_{0}^{T} \int_{\Omega} e^{-2 s \alpha} \xi^{-1}\left(\left|\partial_{t} \varphi\right|^{2}+|\Delta \varphi|^{2}\right) d x d t+s^{-1} \int_{0}^{T} \int_{\Gamma} e^{-2 s \alpha} \xi^{-1}\left(\left|\partial_{t} \varphi_{\Gamma}\right|^{2}+\left|\Delta_{\Gamma} \varphi\right|^{2}\right) d S d t \\
& +s \lambda^{2} \int_{0}^{T} \int_{\Omega} e^{-2 s \alpha} \xi|\nabla \varphi|^{2} d x d t+s \lambda \int_{0}^{T} \int_{\Gamma} e^{-2 s \alpha} \xi\left|\nabla_{\Gamma} \varphi_{\Gamma}\right|^{2} d S d t \\
& +s^{3} \lambda^{4} \int_{0}^{T} \int_{\Omega} e^{-2 s \alpha} \xi^{3}|\varphi|^{2} d x d t+s^{3} \lambda^{3} \int_{0}^{T} \int_{\Gamma} e^{-2 s \alpha} \xi^{3}\left|\varphi_{\Gamma}\right|^{2} d S d t \\
& +s \lambda \int_{0}^{T} \int_{\Gamma} e^{-2 s \alpha} \xi\left|\partial_{\nu} \varphi\right|^{2} d S d t \\
& \leq C s^{3} \lambda^{4} \int_{0}^{T} \int_{\omega} e^{-2 s \alpha} \xi^{3}|\varphi|^{2} d x d t+C \int_{0}^{T} \int_{\Omega} e^{-2 s \alpha}\left|\partial_{t} \varphi+d \Delta \varphi-a \varphi\right|^{2} d x d t \\
& +C \int_{0}^{T} \int_{\Gamma} e^{-2 s \alpha}\left|\partial_{t} \varphi_{\Gamma}+\delta \Delta_{\Gamma} \varphi_{\Gamma}-d \partial_{\nu} \varphi-b \varphi_{\Gamma}\right|^{2} d S d t \tag{1.2.6}
\end{align*}
$$

for all $\lambda \geq \lambda_{1}$, and for all

$$
\left(\varphi, \varphi_{T}\right) \in H^{1}\left(0, T ; L^{2}(\Omega) \times L^{2}(\Gamma)\right) \cap L^{2}\left(0, T ; \mathbb{H}^{2}\right),
$$

where $\mathbb{H}^{2}=\left\{\left(y, y_{\Gamma}\right) \in H^{2}(\Omega) \times H^{2}(\Gamma)\right.$ with $\left.\left.y\right|_{\Gamma}=y_{\Gamma}\right\}$. Furthermore, given $R>0$, the constant $C=C(R)$ can be chosen independently of all $a, b$ with

$$
\|a\|_{\infty},\|b\|_{\infty} \leq R
$$

### 1.3 Main results of the thesis

In this section we briefly introduce the problems and results obtained in this thesis. The main topics covered here are:

- Potential reconstruction for a class of hyperbolic systems in cascade: In Chapter 2, we analyze the simultaneous reconstruction of each potentials $q_{1}, \ldots, q_{n}$ defined in $\Omega \subset \mathbb{R}^{N}, N \geq 1$ in a linear hyperbolic system of the form

$$
\begin{cases}\square u_{1}+q_{1} u_{1}=a_{1} u_{2}+g_{1}, & \text { in } \Omega \times(0, T),  \tag{1.3.1}\\ \square u_{2}+q_{2} u_{2}=a_{2} u_{3}+g_{2}, & \text { in } \Omega \times(0, T), \\ \quad \vdots & \vdots \\ \square u_{n-1}+q_{n-1} u_{n-1}=a_{n-1} u_{n}+g_{n-1}, & \text { in } \Omega \times(0, T), \\ \square u_{n}+q_{n} u_{n}=g_{n}, & \text { in } \Omega \times(0, T), \\ \partial_{t}^{k} u_{j}(0)=u_{j}^{k}, k=0,1, j=1, \ldots, n, & \text { in } \Omega, \\ u_{j}=0, j=1, \ldots, n, & \text { on } \partial \Omega \times(0, T) .\end{cases}
$$

from a reduced number of controls of $\left(u_{j}\right)_{j \in \mathcal{I}}$ with $j \in \mathcal{I} \subset\{1, \ldots, n\}$. Inspired in the Bukhgeim-Klibanov method, we prove a Lipschitz stability result for these coefficients of the form:

$$
\begin{equation*}
\sum_{j=1}^{n}\left\|q_{j}-\tilde{q}_{j}\right\|_{L^{2}(\Omega)}^{2} \leq C \sum_{j=1}^{n-2}\left\|u_{j}-\tilde{u}_{j}\right\|_{H^{3}\left(0, T ; L^{2}(\omega)\right)}^{2}+C\left\|u_{n-1}-\tilde{u}_{n-1}\right\|_{H^{4}\left(0, T ; L^{2}(\omega)\right)}^{2} \tag{1.3.2}
\end{equation*}
$$

where $\tilde{u}_{j}$ with $j=1, \ldots, n$ is the solution of (1.3.1) associated to $\left(\tilde{q}_{1}, \ldots, \tilde{q}_{n}\right)$ (source and initial conditions are fixed), $C$ is a positive constant independent of these potentials and $\omega \subset \Omega$. We point out that in 1.3.2, measurements of $u_{n}-\tilde{u}_{n}$ do not appear in our results. This means that we can reconstruct the potentials $q_{1}, \ldots, q_{n}$ of (1.3.1) without any knowledge of the last component of the system.

The main ingredient to prove this result is a Carleman estimate for problems having the form:

$$
\begin{cases}\square v_{1}+r_{1} v_{1}=v_{2}+h_{1}, & \text { in } \Omega \times(-T, T),  \tag{1.3.3}\\ \square v_{2}+r_{2} v_{2}=v_{3}+h_{2}, & \text { in } \Omega \times(-T, T), \\ \multicolumn{1}{c}{\quad} & \vdots \\ \square v_{n-1}+r_{n-1} v_{n-1}=v_{n}+h_{n-1}, & \text { in } \Omega \times(-T, T), \\ \square v_{n}+r_{n} v_{n}=h_{n}, & \text { in } \Omega \times(-T, T), \\ \partial_{t}^{k} v_{j}( \pm T)=0, k=0,1, j=1, \ldots, n, & \text { in } \Omega, \\ v_{j}=0, j=1, \ldots, n, & \text { on } \partial \Omega \times(-T, T) .\end{cases}
$$

More precisely, using the abbreviation

$$
\begin{aligned}
I(\alpha, v, \Omega)= & s^{\alpha} \int_{\Omega} e^{2 s \varphi(0)}\left(s^{2}|v(0)|^{2}+\left|\partial_{t} v(0)\right|^{2}+|\nabla v(0)|^{2}\right) d x \\
& +s^{\alpha+1} \int_{-T}^{T} \int_{\Omega} e^{2 s \varphi}\left(s^{2}|v|^{2}+\left|\partial_{t} v\right|^{2}+|\nabla v|^{2}\right) d x d t
\end{aligned}
$$

with

$$
\psi(x, t)=\left|x-x_{0}\right|^{2}-\beta t^{2}+C_{0}, \quad \varphi(x, t)=e^{\lambda \psi(x, t)}, \quad \forall(x, t) \in \Omega \times(-T, T)
$$

and suitable assumptions on $x_{0} \in \mathbb{R}^{N}, \alpha, \beta, C_{0}, \lambda>0$ and $T>0$ we get

$$
\begin{align*}
& \sum_{j=1}^{n-1} I\left(\alpha, v_{j}, \Omega\right)+I\left(0, v_{n}, \Omega\right) \\
\leq & C_{2} s^{\alpha+1} \sum_{j=1}^{n-2} \int_{-T}^{T} \int_{\omega} e^{2 s \varphi}\left(s^{2}\left|v_{j}\right|^{2}+\left|\partial_{t} v_{j}\right|^{2}\right) d x d t+C_{2} s^{\alpha} \sum_{j=1}^{n-2} \int_{-T}^{T} \int_{\Omega} e^{2 s \varphi}\left|h_{j}\right|^{2} d x d t  \tag{1.3.4}\\
& +C_{2} s^{3} \int_{-T}^{T} \int_{\omega} e^{2 s \varphi}\left(s^{5}\left|v_{n-1}\right|^{2}+s^{3}\left|\partial_{t} v_{n-1}\right|^{2}+\left|\partial_{t}^{2} v_{n-1}\right|^{2}\right) d x d t \\
& +C_{2} \int_{-T}^{T} \int_{\Omega} e^{2 s \varphi}\left(s^{3}\left|h_{n-1}\right|^{2}+\left|h_{n}\right|^{2}+s\left|\partial_{t} h_{n-1}\right|^{2}+\left|\partial_{t} h_{n}\right|^{2}\right) d x d t .
\end{align*}
$$

for all $s \geq s_{0}>0$, where $v_{1}, \ldots, v_{n}$ is a solution of (1.3.3) and $C$ being a positive constant. The results of Chapter 2 are based on the article 29 in collaboration with Nicolás Carreño y Axel Osses.

- Controllability properties of a class of heat equation with dynamic boundary conditions: In Chapter 3, the null controllability of a suitable class of 1-D parabolic equations with dynamic boundary conditions is studied. The prototype of such problems is

$$
\begin{cases}\partial_{t} u(x, t)-\partial_{x}^{2} u(x, t)=\chi_{\omega}(x) v(x, t), & \forall(x, t) \in \Omega_{L} \times(0, T),  \tag{1.3.5}\\ \left(u(x, 0), u_{\Gamma}(0)\right)=\left(u_{0}(x), u_{\Gamma, 0}\right), & \forall x \in \Omega_{L}, \\ u_{\Gamma}(t)=u(0, t), & \forall t \in(0, T), \\ u\left(-L_{1}, t\right)=0, & \forall t \in(0, T), \\ u_{\Gamma}^{\prime}(t)+\partial_{x} u(0, t)=0, & \forall t \in(0, T) .\end{cases}
$$

Here $\Omega_{L}=\left(-L_{1}, 0\right) \subset \mathbb{R}$ with $L_{1}>0$ and $\omega \subset \Omega_{L}$. In other words, the goal is to steer the state $u$ of (1.3.5) to a null final target by a suitable choice of the control function, i.e., given $\left(u_{0}, u_{0, \Gamma}\right) \in L^{2}\left(\Omega_{L}\right) \times \mathbb{R}$ and $T>0$ we want to find a control $v \in L^{2}(\omega \times(0, T))$ such that the associated solution of (1.3.5) satisfies

$$
u(x, T)=0, \quad \forall x \in \overline{\Omega_{L}} .
$$

This means that the first equation is controlled directly by the action of $v$, while the ODE at $x=0$ is being controlled indirectly through the coupling.

Concerning this question, our results provide that (1.3.5) is null-controllable at any time $T>0$ with $\omega=\left(-L_{1},-a\right)$, with $a>0$. However, some discussions are presented in the case $\omega \subset \subset \Omega_{L}$.

Following the duality between controllability and observability, the proof of this result consists in obtaining an observability estimate of the form

$$
\begin{equation*}
\|z(\cdot, 0)\|_{L^{2}\left(\Omega_{L}\right)}^{2}+\left|z_{\Gamma}(0)\right|^{2} \leq C \int_{0}^{T} \int_{\omega}|z|^{2} d x d t \tag{1.3.6}
\end{equation*}
$$

for each $\left(z_{T}, z_{T, \Gamma}\right) \in L^{2}\left(\Omega_{L}\right) \times \mathbb{R}$, where $\left(z, z_{\Gamma}\right)$ is the solution of the adjoint system

$$
\begin{cases}\partial_{t} z(x, t)+\partial_{x}^{2} z(x, t)=0, & \forall(x, t) \in \Omega_{L} \times(0, T),  \tag{1.3.7}\\ \left.z(x, T), z_{\Gamma}(T)\right)=\left(z_{T}(x), z_{\Gamma, T}\right), & \forall x \in \Omega_{L}, \\ z\left(-L_{1}, t\right)=0, & \forall t \in(0, T), \\ z_{\Gamma}^{\prime}(t)-\partial_{x} z(0, t)=0, & \forall t \in(0, T) .\end{cases}
$$

In order to prove 1.3.6 we use a Carleman estimate of the form

$$
\begin{align*}
& s^{3} \int_{0}^{T} \int_{\Omega_{L}} e^{-2 s \varphi}(t(T-t))^{-3 \alpha}|y|^{2} d x d t+s \int_{0}^{T} \int_{\Omega_{L}} e^{-2 s \varphi}(t(T-t))^{-\alpha}\left|\partial_{x} y\right|^{2} d x d t \\
& +s^{-1} \int_{0}^{T} \int_{\Omega_{L}} e^{-2 s \varphi}(t(T-t))^{\alpha}\left(\left|\partial_{x}^{2} y\right|^{2}+\left|\partial_{t} y\right|^{2}\right) d x d t \\
& +s^{3} \int_{0}^{T} e^{-2 s \varphi(0, t)}(t(T-t))^{-3 \alpha}\left|y_{\Gamma}(t)\right|^{2} d t \\
& +s \int_{0}^{T} e^{-2 s \varphi(0, t)}(t(T-t))^{-\alpha}\left|\partial_{x} y(0, t)\right|^{2} d t+\int_{0}^{T} e^{-2 s \varphi(0, t)}\left|y_{\Gamma}^{\prime}(t)\right|^{2} d t  \tag{1.3.8}\\
& \leq C \int_{0}^{T} \int_{\Omega_{L}} e^{-2 s \varphi}\left|\partial_{t} y+\partial_{x}^{2} y\right|^{2} d x d t+C \int_{0}^{T} e^{-2 s \varphi(0, t)}\left|y_{\Gamma}^{\prime}(t)-\partial_{x} y(0, t)\right|^{2} d t \\
& +C s \int_{0}^{T} e^{-2 s \varphi\left(-L_{1}, t\right)}(t(T-t))^{-\alpha}\left|\partial_{x} y\left(-L_{1}, t\right)\right|^{2} d t
\end{align*}
$$

for all $s \geq s_{1} \geq 1$ and for all $\left(y, y_{\Gamma}\right)$ smooth enough, where $\varphi=\theta(t) \psi(x)$ and

$$
\begin{aligned}
\theta(t) & =(t(T-t))^{-\alpha}, \quad \forall t \in(0, T), \\
\psi(x) & =-\frac{1}{4 L_{1}} x^{2}+x+2 L_{1}, \quad \forall x \in \Omega_{L}
\end{aligned}
$$

On the other hand, we prove that a similar problem to 1.3 .5 appears as limit of a sequence of solutions for parabolic problems with discontinuous diffusion coefficients. In this context, functional setting of both problems play an important role.

In addition, according to the above result of convergence, questions arise naturally. One of them is: can we employ the limit control of the problem to drive the solutions of the approximate system too? Under suitable assumptions of the initial conditions, we prove that the last system is approximately controllable at any time $T>0$.

These results are based on a joint work with Jéremi Dardé and Sylvain Ervedoza.

- Controllability of a 1-D heat equation with discontinuous diffusion coefficients: In Chapter 4, we study controllability properties of the following class of problems:

$$
\begin{cases}\partial_{t} u(x, t)-\partial_{x}\left(\sigma(x) \partial_{x} u(x, t)\right)=0, & \forall(x, t) \in \Omega \times(0, T),  \tag{1.3.9}\\ u(x, 0)=u^{0}(x), & \forall x \in \Omega, \\ u\left(-L_{1}, t\right)=v(t), & \forall t \in(0, T), \\ \partial_{x} u\left(L_{2}, t\right)=0, & \forall t \in(0, T) .\end{cases}
$$

Here, $\Omega=\left(-L_{1}, L_{2}\right) \subset \mathbb{R}$ with $L_{1}, L_{2}>0, T>0$ and $\sigma$ is given by

$$
\sigma(x)= \begin{cases}\sigma_{1}^{2}, & \forall x \in\left(-L_{1}, 0\right) \\ \sigma_{2}^{2}, & \forall x \in\left(0, L_{2}\right)\end{cases}
$$

Moreover, the control $v=v(t)$ acts only in the left-hand side of the domain. Then, once again, in order to obtain the null controllability of such systems, we look for a Carleman estimate of the form

$$
\begin{aligned}
& s^{3} \int_{0}^{T} \int_{\Omega} \rho|z|^{2} d x d t+s \int_{0}^{T} \int_{\Omega} \mu\left|\partial_{x} z\right|^{2} d x d t \\
\leq & C \int_{0}^{T} \int_{\Omega} \nu\left|\partial_{t} \pm \partial_{x}\left(\sigma \partial_{x} z\right)\right|^{2} d x d t+C s \int_{0}^{T} \mu\left(t,-L_{1}\right)\left|\partial_{x} z\left(t,-L_{1}\right)\right|^{2} d t
\end{aligned}
$$

for some positive functions $\rho, \mu$ and $\nu$. In order to prove it, we use similar arguments based on [86] and [75] and suitable localization in time functions. In spite of microlocal techniques, this choice allows us to keep tracking of the observability constant.

## Chapter 2

# Potential reconstruction for a class of hyperbolic systems in cascade 

In this chapter we present new results concerning the potential reconstruction of wave systems in cascade when some components of this ones (that is to say, some variables of the system) are not available to get partial measurements. We adapt the BukhgeimKlibanov method to the case of hyperbolic systems and we use Carleman estimates for the scalar wave equation to achieve a new Carleman estimates for a hyperbolic system in cascade with missing components. Let us mention that the main results of this chapter were published in [29] in collaboration with Nicolás Carreño and Axel Osses.

The outline of this chapter is as follows. In Section 2.1, we introduce the basic notation, the inverse problem that we will consider along this chapter. Additionally, we give a literature discussion about this subject and we state the main result obtained, i.e. the Theorem 2.1. In Section 2.2, we adapt the Carleman estimate for the scalar wave equation to deduce a new Carleman inequality to our problem (see Theorem 2.8). In Section 2.3 , we modify the Bukhgeim-Klibanov's method to proof of the Theorem 2.1.

### 2.1 General Setting

In this section, we devote to introduce the main results about an inverse problem for a hyperbolic system in cascade. Let us start giving basic notations. Let $\Omega$ be a smooth open set in $\mathbb{R}^{d}$ with boundary $\partial \Omega, d \geq 1$ and $T>0$.

Before going further, let us mention that the results available in this section could be formulated under weak smoothness assumptions of the boundary of $\Omega$. Indeed, for instance we can take $\Omega$ be a bounded, connected and open subset of $\mathbb{R}^{d}$ with boundary of class $C^{4}$, but the goal of these hypothesis is to simplify the presentation.

Then, according to the notation previously introduced, let us consider the following
coupled hyperbolic system in cascade:

$$
\begin{cases}\square u_{1}+q_{1} u_{1}=a_{1} u_{2}+g_{1}, & \text { in } \Omega \times(0, T),  \tag{2.1.1}\\ \square u_{2}+q_{2} u_{2}=a_{2} u_{3}+g_{2}, & \text { in } \Omega \times(0, T), \\ \quad \vdots & \vdots \\ \square u_{n-1}+q_{n-1} u_{n-1}=a_{n-1} u_{n}+g_{n-1}, & \text { in } \Omega \times(0, T), \\ \square u_{n}+q_{n} u_{n}=g_{n}, & \text { in } \Omega \times(0, T), \\ \partial_{t}^{k} u_{j}(0)=u_{j}^{k}, k=0,1, j=1, \ldots, n, & \text { in } \Omega, \\ u_{j}=0, j=1, \ldots, n, & \text { on } \partial \Omega \times(0, T) .\end{cases}
$$

Here, $\square:=\partial_{t}^{2}-\Delta$ is the D'Alembertian operator, $a_{j}$ are non-zero constants, $u_{j}^{k} \in L^{2}(\Omega)$ are the initial conditions and $q_{j} \in L^{2}(\Omega)$ are the potentials and $g_{j} \in L^{2}(\Omega \times(0, T))$ are the source terms, for every $k=0,1$ and $j=1, \ldots, n$.

Suppose, for instance, that $g_{j} \in L^{1}\left(0, T ; L^{2}(\Omega)\right), u_{j}^{0} \in H_{0}^{1}(\Omega)$ and $u_{j}^{1} \in L^{2}(\Omega), j=$ $1, \ldots, n$. Then, according to the results presented in Section 1.1.1 on the well-posedness of the scalar wave equation, it is not difficult to deduce that the system $(2.1 .1)$ is well posed in the sense of Hadamard and moreover the normal derivative of each component $u_{j}$ with $1 \leq j \leq n$ belongs in $L^{2}(\partial \Omega \times(0, T))$. We recall that in the case of source terms $g_{j}$ lying in $W^{1,1}\left(0, T ; H^{-1}(\Omega)\right)$ with $1 \leq j \leq n$, we still have the same well-posedness result but we do not have any regularity result on the normal derivative on each component.

Let us point out that the hyperbolic and parabolic systems play an important role in mathematical models which come from biological, chemical, engineering, mechanical and medical applications. Nevertheless, some components of such models are not accessible in practice. Motivated for this kind of limitations, some natural questions arise: Can we observe such systems from incomplete measurements? Can we retrieve information of the inaccessible components of such systems from information of the accesible ones? These questions has been studied recently by several authors for different kind of PDE models, see for instance [1], [18] and [6] and the bibliographic discusion below.

In this chapter, we are interested in the following inverse problem associated to the system (2.1.1):

Inverse problem: Is it possible to retrieve the potentials $q_{1}, \ldots, q_{n}$ in system (2.1.1) from incomplete data, that is to say, from a reduced number of measurements of the solution, saying $\left(u_{j}\right)$, with $j \in \mathcal{I} \subset\{1, \ldots, n\}$ ?

We point out that our goal is the study of dependence of the solutions $u_{1}, \ldots, u_{n}$ with respect to the potentials $q_{1}, \ldots, q_{n}$. Then, in order to understand this we shall write $u_{j}[Q]$ where $Q=\left(q_{1}, \ldots, q_{n}\right)$ and $1 \leq j \leq n$. For simplicity, we ignore for instance the dependence of the solutions with respect to the initial conditions and source terms.

Concerning to the inverse problem previously stated, we can formulate questions in three directions: Let $Q=\left(q_{1}, \ldots, q_{n}\right)$ and $\tilde{Q}=\left(\tilde{q}_{1}, \ldots, \tilde{q}_{n}\right)$ be two sets of potentials for the system (2.1.1),

- Uniqueness: Suppose that the available measurements of the system coincide in a
part of the domain $\omega \subset \Omega$ for two sets of potentials, i.e.,

$$
u_{j}[Q]=u_{j}[\tilde{Q}], \quad \text { in } \omega \times(0, T),
$$

for each $j \in \mathcal{I}$. Then, can we conclude that $Q=\tilde{Q}$ in $\Omega$ ?

- Stability: is it possible to estimate

$$
\|Q-\tilde{Q}\|_{\left(L^{2}(\Omega)\right)^{n}}:=\sum_{j=1}^{n}\left\|q_{j}-\tilde{q}_{j}\right\|_{L^{2}(\Omega)}
$$

or better, a strong norm of $Q-\tilde{Q}$ by a suitable norm of $u_{j}[Q]-u_{j}[\tilde{Q}], j \in \mathcal{I}$ in $\omega \times(0, T)$ ?

- Reconstruction: Can we find a formula or an algorithm to rebuild the potentials $Q=\left(q_{1}, \ldots, q_{n}\right)$ from the knowledge of $u_{j}[Q]$ with $j \in \mathcal{I}$ ?

In this chapter, we are interested in the stability of the potentials in terms of the observations of the solution of system (2.1.1). In particular, we restrict our attention in the case when the last component of the system is missing.

### 2.1.1 Literature review

Before to state the main results of this topic, let us briefly discuss the available literature about inverse problems of wave equation and systems and their relation with exact controllability.

In 1992, Bukhgeim and Klibanov dealt for the first time with uniqueness issues in inverse problems for the wave equation in [69] using local Carleman estimates. Then, the first results about the stability of inverse problems for hyperbolic equations were obtained using local Carleman estimates (see e.g. [97], 61] and [65]). Concerning other inverse problems for the wave equation with a single observation, we refer to [63], [62], [70], and [11] and the references therein. In these articles, the authors consider the case of interior or Dirichlet boundary data observation satisfying stronger geometric conditions and they use global Carleman estimates.

For an arbitrary set of observation, we refer to [14] and [15] for logarithmic stability results. Roughly speaking, these results are connected with stability results of elliptic thanks to the Fourier-Bros-Iagolnitzer (FBI) transform. Let us also mention the work 66] where the authors proved the uniqueness of the inverse problem of recovering a spatial component of the source term of the wave equation from the final observation data.

However, to the best of the author's knowledge, there exist few works concerning inverse problems for coupled parabolic or hyperbolic systems with incomplete measurements of their components. In the recent work [6], the authors study the reconstruction of the spatial distribution of external forces only from data of one component of a 2 coupled hyperbolic system in cascade. The proof is based on an observability property of such system, following the approach of 96 .

Similar inverse problems for linear and semilinear parabolic systems like reactiondiffusion systems has been studied in [34], [18], [17], [35] and [33]. In these articles,
the authors deal with identification and stability of the inverse problem of recovering parameters and initial conditions of such systems from a finite number of measurements of one component using appropriate Carleman estimates for parabolic equations.

Furthermore, hyperbolic-parabolic systems are considered in [51] with different kinds of observations. Another relevant work is [56] for the Stokes system, where the authors give a reconstruction algorithm for a source of the form $F(x, t)=f(x) \sigma(t)$ from incomplete velocity measurements.

Exact controllability properties of hyperbolic systems with a reduced number of controls has been extensively studied and there exist many works published on this topic. In [1], a strategy called Two-Level Energy method is developed to prove positive results in the case of wave-type systems (see also [2, ,3, [7] and the references therein). Moreover these results allow to deduce null-controllability results for the heat or the Schrödinger equations satisfying the geometric control condition using the transmutation method.

Furthermore, the literature is also very rich concerning controllability results for coupled parabolic systems with a reduced numbers of controls in the one or multidimensional setting. We refer to the survey article [8], [45], and the references therein.

Coupled systems are also connected with insensitizing control problems, notion introduced by Lions in [79]. Indeed, these problems are equivalent to the null-controllability of a cascade system. We reference to [36], [90], and [4] for some results about this subject in the case of wave-type equations, [23], [28] and [37] in the case of parabolic equations and systems.

### 2.1.2 Main result

Now, we will state a Lipschitz stability result for system (2.1.1), from observations in all the components of the system except the last one. In order to state this result, we shall introduce some geometrical and time conditions which are classical in the context of control and inverse problems for hyperbolic equations. Specifically, let $x_{0} \notin \bar{\Omega}, \Gamma_{0} \subset \partial \Omega$ and $T>0$.

- Geometric condition: $x_{0}$ and $\Gamma_{0}$ satisfy the following inclusion:

$$
\begin{equation*}
\left\{x \in \partial \Omega ;\left(x-x_{0}\right) \cdot \nu(x) \geq 0\right\} \subset \Gamma_{0} \subset \partial \Omega . \tag{2.1.2}
\end{equation*}
$$

- Time condition: There exists $\beta \in(0,1)$ such that

$$
\begin{equation*}
\sup _{x \in \Omega}\left|x-x_{0}\right|<\sqrt{\beta} . \tag{2.1.3}
\end{equation*}
$$

We emphasize that the Geometric condition given above is the same in 1.1.1. However, the Time condition is slightly different from the mentioned in the above chapter. This change is just for technical reasons to state our results in a simple way.

Now, let us introduce the admissible set of the unknown potentials. For a positive number $m$, we define the set

$$
L_{\leq m}^{\infty}(\Omega)=\left\{p \in L^{\infty}(\Omega) ;\|p\|_{L^{\infty}(\Omega)} \leq m\right\}
$$

Now we have all the ingredients to state the main result of this chapter:
Theorem 2.1 Suppose that $\Gamma_{0} \subset \partial \Omega, T>0$ and $x_{0} \notin \bar{\Omega}$ satisfy the geometric and time conditions 2.1.2 and 2.1.3. Let $\omega \subset \Omega$ such that $\overline{\Gamma_{0}} \subset \partial \omega \cap \partial \Omega$. Let $\left(u_{1}, \ldots, u_{n}\right)$ and $\left(\tilde{u}_{1}, \ldots, \tilde{u}_{n}\right)$ be the solutions of the system (2.1.1) associated to the potentials $q_{1}, \ldots, q_{n} \in$ $L_{\leq m}^{\infty}(\Omega)$ and $\tilde{q}_{1}, \ldots, \tilde{q}_{n} \in L_{\leq m}^{\infty}(\Omega)$, respectively, with $m>0$. Assume that there exists a constant $c>0$ such that

$$
\begin{equation*}
\left\|\tilde{u}_{j}(0)\right\|_{L^{2}(\Omega)}^{2} \geq c, \quad \forall j=1,2, \ldots, n \tag{2.1.4}
\end{equation*}
$$

Furthermore, suppose that

$$
\left\{\begin{array}{l}
u_{j}, \tilde{u}_{j} \in H^{3}\left(0, T ; H^{2}(\Omega) \cap H_{0}^{1}(\Omega) \cap L^{\infty}(\Omega)\right), \quad j=1, \ldots, n, \\
u_{n-1}, \tilde{u}_{n-1} \in H^{4}\left(0, T ; H^{2}(\Omega) \cap H_{0}^{1}(\Omega) \cap L^{\infty}(\Omega)\right) .
\end{array}\right.
$$

Then, there exists a constant $C=C(\beta, c, T, \Omega, \omega)$ such that

$$
\begin{equation*}
\sum_{j=1}^{n}\left\|q_{j}-\tilde{q}_{j}\right\|_{L^{2}(\Omega)}^{2} \leq C \sum_{j=1}^{n-2}\left\|u_{j}-\tilde{u}_{j}\right\|_{H^{3}\left(0, T ; L^{2}(\omega)\right)}^{2}+C\left\|u_{n-1}-\tilde{u}_{n-1}\right\|_{H^{4}\left(0, T ; L^{2}(\omega)\right)}^{2} \tag{2.1.5}
\end{equation*}
$$

Remark 2.2 Let us emphasize that inequality (2.1.5 establishes the Lipschitz stability of the hyperbolic system (2.1.1) with incomplete measurements in the sense that $u_{n}$ is missing. Moreover, notice that the estimate (2.1.5) does not depend on the observations of the gradients.

Remark 2.3 Theorem 2.1 is also valid if we suppose that the coupling coefficients $a_{j}$ are not constants satisfying

$$
a_{j}(x) \geq c>0, \quad \text { in } \omega^{\prime},
$$

where $\omega^{\prime} \subset \Omega$ such that $\overline{\Gamma_{0}} \subset \partial \omega^{\prime}$ and $\omega^{\prime} \cap \omega \neq \emptyset$. In other words, the inequality (2.1.5) holds if the coupled and the observations regions of each components meet.

As we said before, the main tool of the proof of Theorem 2.1 is a Carleman estimate for a hyperbolic system in cascade where we do not have access to the observations associated to the last component. This inequality depends on a suitable Carleman estimate for the scalar wave equation in the spirit of the work of Imanuvilov and Yamamoto [62] (see also [12]).

### 2.2 Carleman estimates

The goal of this section is to prove a Carleman estimate for a system of wave equations in cascade. In order to do that, our starting point is a suitable Carleman estimate for the scalar wave equation. Nevertheless, before doing that, we will give some technical results.

### 2.2.1 Technical results

We start with the following

Lemma 2.4 Let $z \in L^{2}\left(-T, T ; H_{0}^{1}(\Omega)\right)$ be a function such that $\square z+p z \in L^{2}(\Omega \times(-T, T))$, $\partial_{\nu} z \in L^{2}(\partial \Omega \times(-T, T))$ and $z( \pm T)=0$ in $\Omega$, with $p \in L^{\infty}(\Omega)$. Let $\gamma \in \mathbb{R}$. Let $\omega_{1}, \omega_{2} \subset \Omega$ be two open sets such that $\overline{\omega_{1}} \subset \omega_{2}$.
a) If $\tilde{\varphi} \in C^{1}([-T, T] \times \bar{\Omega})$, then there exists a constant $C>0$ such that

$$
\begin{gather*}
\int_{-T}^{T} \int_{\omega_{1}} e^{2 s \tilde{\varphi}}|\nabla z|^{2} d x d t \leq C s^{\max \{2, \gamma\}} \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \tilde{\varphi}}|z|^{2} d x d t+C \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \tilde{\varphi}}\left|\partial_{t} z\right|^{2} d x d t \\
+C s^{-\gamma} \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \tilde{\varphi}}|\square z+p z|^{2} d x d t \tag{2.2.1}
\end{gather*}
$$

for all $s \geq 1$.
b) If the function $\tilde{\varphi} \in C^{1}\left([-T, T] ; C^{2}(\bar{\Omega})\right)$ satisfies

$$
\inf _{x \in \bar{\Omega}}|\nabla \tilde{\varphi}(t)| \geq c_{0}>0, \quad \forall t \in[-T, T]
$$

then, there exist two positive constants $C$ and $s_{0}$ independent of s such that

$$
\begin{align*}
& s^{2} \int_{-T}^{T} \int_{\omega_{1}} e^{2 s \tilde{\varphi}}|z|^{2} d x d t+\int_{-T}^{T} \int_{\omega_{1}} e^{2 s \tilde{\varphi}}\left|\partial_{t} z\right|^{2} d x d t \\
\leq & C s^{\max \{0, \gamma-2\}} \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \tilde{\varphi}}|\nabla z|^{2} d x d t+C s^{-\gamma} \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \tilde{\varphi}}|\square z+p z|^{2} d x d t \tag{2.2.2}
\end{align*}
$$

for all $s \geq s_{0}$.
Remark 2.5 The principal significance of part a) on Lemma 2.4 is that it allows to drop the local term of the gradient. This fact plays an important role in some steps of the proof of the Carleman estimate for the wave system in cascade in Subsection 2.2.

Proof of Lemma 2.4. Let us consider a function $\xi \in C^{\infty}(\Omega, \mathbb{R})$ such that

$$
\begin{cases}0 \leq \xi \leq 1, & \text { in } \Omega \\ \xi \equiv 1, & \text { in } \omega_{1}, \\ \xi \equiv 0, & \text { in } \Omega \backslash \overline{\omega_{2}}\end{cases}
$$

Additionally, we suppose that $\xi$ has the form $\xi=e^{\phi}$ in $\omega_{2} \backslash \overline{\omega_{1}}$, for some smooth function $\phi$. We have the following identity:

$$
\begin{align*}
& \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \tilde{\varphi}} \xi z \partial_{t}^{2} z d x d t-\int_{-T}^{T} \int_{\omega_{2}} e^{2 s \tilde{\varphi}} \xi z \Delta z d x d t+\int_{-T}^{T} \int_{\omega_{2}} e^{2 s \tilde{\varphi}} \xi p|z|^{2} d x d t \\
= & \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \tilde{\varphi}} \xi z(\square z+p z) d x d t . \tag{2.2.3}
\end{align*}
$$

Integration by parts yields

$$
\begin{equation*}
\int_{-T}^{T} \int_{\omega_{2}} e^{2 s \tilde{\varphi}} \xi z \partial_{t}^{2} z d x d t=-2 s \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \tilde{\varphi}} \partial_{t} \tilde{\varphi} \xi z \partial_{t} z d x d t-\int_{-T}^{T} \int_{\omega_{2}} e^{2 s \varphi} \xi\left|\partial_{t} z\right|^{2} d x d t \tag{2.2.4}
\end{equation*}
$$

and

$$
\begin{align*}
-\int_{-T}^{T} \int_{\omega_{2}} e^{2 s \tilde{\varphi}} \xi z \Delta z d x d t= & 2 s \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \tilde{\varphi}} \xi z \nabla \tilde{\varphi} \cdot \nabla z d x d t+\int_{-T}^{T} \int_{\omega_{2}} e^{2 s \varphi} z \nabla \xi \cdot \nabla z d x d t \\
& +\int_{-T}^{T} \int_{\omega_{2}} e^{2 s \varphi} \xi|\nabla z|^{2} d x d t . \tag{2.2.5}
\end{align*}
$$

Substituting (2.2.4) and (2.2.5) into (2.2.3), we have

$$
\begin{align*}
& \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \varphi} \xi|\nabla z|^{2} d x d t \\
= & 2 s \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \varphi} \xi \partial_{t} \tilde{\varphi} z \partial_{t} z d x d t+\int_{-T}^{T} \int_{\omega_{2}} e^{2 s \tilde{\varphi}} \xi\left|\partial_{t} z\right|^{2} d x d t+\int_{-T}^{T} \int_{\omega_{2}} e^{2 s \tilde{\varphi}} \xi p|z|^{2} d x d t \\
& +\int_{-T}^{T} \int_{\omega_{2}} e^{2 s \tilde{\varphi}} \xi z(\square z+p z) d x d t-2 s \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \tilde{\varphi}} \xi z \nabla \tilde{\varphi} \cdot \nabla z d x d t  \tag{2.2.6}\\
& -\int_{-T}^{T} \int_{\omega_{2}} e^{2 s \tilde{\varphi}} z \nabla \xi \cdot \nabla z d x d t=J_{1}+J_{2} .
\end{align*}
$$

Here, $J_{1}$ is the sum of the first four terms of $(2.2 .6)$ and $J_{2}$ is the sum of the fifth and sixth terms of the same equation. Straightforward computations show that

$$
\begin{align*}
\left|J_{1}\right| \leq & 2\|\xi\|_{C^{0}\left(\overline{\omega_{2}}\right)} \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \tilde{\varphi}}\left|\partial_{t} z\right|^{2} d x d t+\frac{1}{3} s^{-\gamma} \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \tilde{\varphi}}|\square z+p z|^{2} d x d t \\
& +\Lambda \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \tilde{\varphi}}|z|^{2} d x d t \tag{2.2.7}
\end{align*}
$$

where $\Lambda$ is defined by

$$
\Lambda=\left(\left\|\partial_{t} \tilde{\varphi}\right\|_{C^{0}\left(\overline{\left.\omega_{2} \times(-T, T)\right)}\right.}^{2} s^{2}+\frac{3}{4}\|\xi\|_{C^{0}\left(\overline{\omega_{2}}\right)} s^{\gamma}+\|p\|_{L^{\infty}\left(\omega_{2}\right)}\right)\|\xi\|_{C^{0}\left(\overline{\omega_{2}}\right)}
$$

and

$$
\begin{align*}
\left|J_{2}\right| \leq & \left(3\|\nabla \tilde{\varphi}\|_{L^{\infty}\left(\omega_{2} \times(-T, T)\right)}^{2} s^{2}+\frac{3}{4}\|\nabla \phi\|_{L^{\infty}\left(\omega_{2} \backslash \overline{\omega_{1}}\right)}\right)\|\xi\|_{L^{\infty}\left(\omega_{2}\right)} \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \tilde{\varphi}}|z|^{2} d x d t \\
& +\frac{2}{3} \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \tilde{\varphi}} \xi|\nabla z|^{2} d x d t \tag{2.2.8}
\end{align*}
$$

Combining (2.2.7), (2.2.8) with (2.2.6) we obtain (2.2.1), which completes the part a) of Lemma 2.4. The rest of the proof runs as before but additionally we have to estimate the local term $|z|^{2}$ by using the weighted Poincaré inequality (see [11], Lemma 2.4).

Now, we introduce the classical Carleman weights for the scalar wave equation. Suppose that $\Gamma_{0}, x_{0}$ and $T>0$ satisfy the Geometric and Time condition (2.1.2) and 2.1.3). Let $\beta \in(0, T)$. For $(x, t) \in \Omega \times(-T, T)$, we define the following functions:

$$
\begin{equation*}
\psi(x, t)=\left|x-x_{0}\right|^{2}-\beta t^{2}+C_{0}, \quad \varphi(x, t)=e^{\lambda \psi(x, t)} \tag{2.2.9}
\end{equation*}
$$

where $\lambda>0$ and $C_{0}>0$ is chosen such that $\psi \geq 0$ (and therefore $\varphi \geq 1$ ) in $\Omega \times(-T, T)$.
For brevity, we shall use the following notation

$$
\begin{aligned}
I(\alpha, v, \Omega)= & s^{\alpha} \int_{\Omega} e^{2 s \varphi(0)}\left(s^{2}|v(0)|^{2}+\left|\partial_{t} v(0)\right|^{2}+|\nabla v(0)|^{2}\right) d x \\
& +s^{\alpha+1} \int_{-T}^{T} \int_{\Omega} e^{2 s \varphi}\left(s^{2}|v|^{2}+\left|\partial_{t} v\right|^{2}+|\nabla v|^{2}\right) d x d t
\end{aligned}
$$

In the remainder of this section, $C$ denotes a generic positive constant which depends at least on $\Gamma_{0}, T$ and $x_{0}$ and may change from line to line.

Proposition 2.6 Assume that $\Gamma_{0}, T$ and $x_{0}$ satisfy the Geometric condition and Time condition (2.1.2 and 2.1.3) and let $p \in L_{\leq m}^{\infty}(\Omega)$ with $m>0$. Let us consider the Carleman weight functions defined in (2.2.9). Let $\omega_{0} \subset \Omega$ be an open subset such that $\overline{\Gamma_{0}} \subset \partial \omega_{0} \cap \partial \Omega$. Then, there exist two positive constants $C_{1}=C_{1}\left(\Gamma_{0}, T, x_{0}, \omega_{2}\right)$ and $s_{0}=s_{0}\left(\Gamma_{0}, T, x_{0}, \omega_{2}\right)$ independent of $s$ such that for all $s \geq s_{0}$, we have

$$
\begin{equation*}
I(0, v, \Omega) \leq C_{1} \int_{-T}^{T} \int_{\Omega} e^{2 s \varphi}|\square v+p v|^{2} d x d t+C_{1} s \int_{-T}^{T} \int_{\omega_{0}} e^{2 s \varphi}\left(s^{2}|v|^{2}+\left|\partial_{t} v\right|^{2}\right) d x d t \tag{2.2.10}
\end{equation*}
$$

for all $v \in L^{2}\left(-T, T ; H_{0}^{1}(\Omega)\right)$ such that $\square v+p v \in L^{2}(\Omega \times(-T, T)), \partial_{\nu} v \in L^{2}(\partial \Omega \times(-T, T))$ and $v( \pm T)=\partial_{t} v( \pm T)=0$ in $\Omega$.

Remark 2.7 In contrast to the Theorem 2.5 in [11] in the case of the Carleman estimate of the scalar wave equation with a single boundary observation, we emphasize that Proposition 2.6 requires the assumptions $z( \pm T)=\partial_{t} z( \pm T)=0$ in $\Omega$. This point becomes important if we want to eliminate more components in the inequality (2.1.5) of Theorem 2.1 .

Let us point out that the proof of the Proposition 2.6 is straightforward and many of the ingredients of the proof are already available in the literature (see for instance [62] and [12]). Nevertheless, for our purposes, it is convenient to write the Carleman estimate for wave equation under the form of Proposition 2.6. For the sake of completeness, we will give the proof of this result.

Proof of Proposition 2.4. For $s \geq 1$, let us define

$$
E_{s}(t)=\frac{1}{2} \int_{\Omega} e^{2 s \varphi(t)}\left(\left|\partial_{t} v(t)\right|^{2}+|\nabla v(t)|^{2}\right) d x, \quad \forall t \in(-T, T) .
$$

Differentiation with respect to $t$ and integration by parts in space yields

$$
\begin{aligned}
\frac{d E_{s}}{d t}(t)= & s \int_{\Omega} e^{2 s \varphi(t)} \partial_{t} \varphi(t)\left(\left|\partial_{t} v(t)\right|^{2}+|\nabla v(t)|^{2}\right) d x+\int_{\Omega} e^{2 s \varphi(t)} \partial_{t} v(t) \square v(t) d x \\
& -2 s \int_{\Omega} e^{2 s \varphi(t)} \partial_{t} v(t) \nabla v(t) \cdot \nabla \varphi(t) d x, \quad \forall t \in(-T, T) .
\end{aligned}
$$

After integration on $(-T, 0)$ in time we obtain

$$
\begin{aligned}
E_{s}(0)= & s \int_{-T}^{0} \int_{\Omega} e^{2 s \varphi} \partial_{t} \varphi\left(\left|\partial_{t} v\right|^{2}+|\nabla v|^{2}\right) d x d t+\int_{-T}^{0} \int_{\Omega} e^{2 s \varphi} \partial_{t} v \square v d x d t \\
& -2 s \int_{-T}^{0} \int_{\Omega} e^{2 s \varphi} \partial_{t} v \nabla v \cdot \nabla \varphi d x d t,
\end{aligned}
$$

where we have used $v(-T)=\partial_{t} v(-T)=0$ in $\Omega$. Applying Young's inequality and the weighted Poincaré inequality to $v$ (see [11], Lemma 2.4) we obtain

$$
\begin{align*}
& \int_{\Omega} e^{2 s \varphi(0)}\left(s^{2}|v(0)|^{2}+\left|\partial_{t} v(0)\right|^{2}+|\nabla v(0)|^{2}\right) d x \\
\leq & C s \int_{-T}^{T} \int_{\Omega} e^{2 s \varphi}\left(s^{2}|v|^{2}+\left|\partial_{t} v\right|^{2}+|\nabla v|^{2}\right) d x d t+C \int_{-T}^{T} \int_{\Omega} e^{2 s \varphi}|\square v+p v|^{2} d x d t, \quad \forall s \geq s_{0} . \tag{2.2.11}
\end{align*}
$$

On the other hand, let us recall the classical Carleman estimate for the wave equation with $\lambda=\lambda_{0}$ fixed applied to $v$ :

$$
\begin{align*}
& s \int_{-T}^{T} \int_{\Omega} e^{2 s \varphi}\left(s^{2}|v|^{2}+\left|\partial_{t} v\right|^{2}+|\nabla v|^{2}\right) d x d t  \tag{2.2.12}\\
\leq & C \int_{-T}^{T} \int_{\Omega} e^{2 s \varphi}|\square v+p v|^{2} d x d t+C s \int_{-T}^{T} \int_{\Gamma_{0}} e^{s \varphi}\left|\partial_{\nu} v\right|^{2} d \sigma d t
\end{align*}
$$

Let us consider an open subset $\omega_{0}^{\prime} \subset \omega_{0}$ such that $\overline{\omega_{0}^{\prime}} \subset \omega$ and $\partial \omega_{0}^{\prime} \cap \partial \Omega \subset \partial \omega_{0} \cap \partial \Omega$. Consider the function $\eta \in C^{\infty}(\bar{\Omega}, \mathbb{R})$ satisfying

$$
\begin{cases}0 \leq \eta \leq 1, & \text { in } \Omega, \\ \eta \equiv 1, & \text { in } \Omega \backslash \overline{\omega_{0}^{\prime}}, \\ \eta=\partial_{\nu} \eta \equiv 0, & \text { on } \Gamma_{0}\end{cases}
$$

Replacing $v$ by $\eta v$ in 2.2.12, we have

$$
\begin{align*}
& s \int_{-T}^{T} \int_{\Omega} e^{2 s \varphi}\left(s^{2}|v|^{2}+\left|\partial_{t} v\right|^{2}+|\nabla v|^{2}\right) d x d t \\
\leq & C \int_{-T}^{T} \int_{\Omega} e^{2 s \varphi}|\square v+p v|^{2} d x d t+C \int_{-T}^{T} \int_{\Omega} e^{2 s \varphi}\left(|v|^{2}+|\nabla v|^{2}\right) d x d t  \tag{2.2.13}\\
& +s \int_{-T}^{T} \int_{\omega_{0}^{\prime}} e^{2 s \varphi}\left(s^{2}|v|^{2}+\left|\partial_{t} v\right|^{2}+|\nabla v|^{2}\right) d x d t
\end{align*}
$$

where we have used that $\square(\eta v)=\eta \square v-\Delta \eta v-2 \nabla \eta \cdot \nabla v$ in $\Omega \times(-T, T)$ and $\nabla \eta \equiv 0$ in $\Omega \backslash \overline{\omega_{1}}$. Notice that the second term of the right-hand side of 2.2 .13 ) can be absorbed taking $s$ large enough. Finally, combining the previous estimate obtained with (2.2.11) and applying the estimate (2.2.1) with $\tilde{\varphi}=\varphi, \omega_{1}=\omega_{0}^{\prime}, \omega_{2}=\omega_{0}$ and $\gamma=1$, the proof of (2.2.10) is complete.

### 2.2.2 A new Carleman estimate for a hyperbolic system

The aim of this section is to prove a Carleman estimate for a wave-type system with potentials. In order to formulate our result, let us consider the following system:

$$
\begin{cases}\square v_{1}+r_{1} v_{1}=v_{2}+h_{1}, & \text { in } \Omega \times(-T, T),  \tag{2.2.14}\\ \square v_{2}+r_{2} v_{2}=v_{3}+h_{2}, & \text { in } \Omega \times(-T, T), \\ \quad \vdots & \vdots \\ \square v_{n-1}+r_{n-1} v_{n-1}=v_{n}+h_{n-1}, & \text { in } \Omega \times(-T, T), \\ \square v_{n}+r_{n} v_{n}=h_{n}, & \text { in } \Omega \times(-T, T), \\ \partial_{t}^{k} v_{j}( \pm T)=0, k=0,1, j=1, \ldots, n, & \text { in } \Omega, \\ v_{j}=0, j=1, \ldots, n, & \text { on } \partial \Omega \times(-T, T)\end{cases}
$$

Here, $r_{j} \in L^{\infty}(\Omega)$ are the potentials and $h_{j} \in L^{2}(\Omega \times(-T, T))$ are the source terms, for each $j=1, \ldots, n$.

Now, we are in position to state the Carleman estimate for system (2.2.14), which is one of the main results of this article:

Theorem 2.8 Let us consider the Carleman weights defined in (2.2.9), where $\Gamma_{0} \subset \partial \Omega$, $T>0$ and $x_{0} \notin \bar{\Omega}$ satisfy the geometric and time conditions (2.1.2) and (2.1.3). For $m>0$, suppose that $r_{j} \in L_{\leq m}^{\infty}(\Omega), j=1, \ldots, n$, and let $\omega \subset \Omega$ be an open set such that $\overline{\Gamma_{0}} \subset \partial \omega \cap \partial \Omega$. In addition, consider $h_{j} \in L^{2}(\Omega \times(-T, T))$ for each $j=1, \ldots, n-2$ and $h_{n-1}, h_{n} \in H^{1}\left(-T, T ; L^{2}(\Omega)\right)$ such that

$$
\left\{\begin{array}{l}
v_{j} \in H^{1}\left(-T, T ; H^{2}(\Omega) \cap H_{0}^{1}(\Omega)\right), j=1, \ldots, n, \\
v_{n-1} \in H^{2}\left(-T, T ; H^{2}(\Omega) \cap H_{0}^{1}(\Omega)\right) .
\end{array}\right.
$$

Furthermore, we choose $1<\alpha<2$. Then, there exist two positive constants $C_{2}$ and $s_{0}$ which depends at least on $\left.\Gamma_{0}, \Omega, \omega, T, x_{0}\right)>0$ such that for all $s \geq s_{0}$, the solution $\left(v_{1}, \ldots, v_{n}\right)$ of system (2.2.14) satisfies

$$
\begin{align*}
& \sum_{j=1}^{n-1} I\left(\alpha, v_{j}, \Omega\right)+I\left(0, v_{n}, \Omega\right) \\
\leq & C_{2} s^{\alpha+1} \sum_{j=1}^{n-2} \int_{-T}^{T} \int_{\omega} e^{2 s \varphi}\left(s^{2}\left|v_{j}\right|^{2}+\left|\partial_{t} v_{j}\right|^{2}\right) d x d t+C_{2} s^{\alpha} \sum_{j=1}^{n-2} \int_{-T}^{T} \int_{\Omega} e^{2 s \varphi}\left|h_{j}\right|^{2} d x d t  \tag{2.2.15}\\
& +C_{2} s^{3} \int_{-T}^{T} \int_{\omega} e^{2 s \varphi}\left(s^{5}\left|v_{n-1}\right|^{2}+s^{3}\left|\partial_{t} v_{n-1}\right|^{2}+\left|\partial_{t}^{2} v_{n-1}\right|^{2}\right) d x d t \\
& +C_{2} \int_{-T}^{T} \int_{\Omega} e^{2 s \varphi}\left(s^{3}\left|h_{n-1}\right|^{2}+\left|h_{n}\right|^{2}+s\left|\partial_{t} h_{n-1}\right|^{2}+\left|\partial_{t} h_{n}\right|^{2}\right) d x d t
\end{align*}
$$

Remark 2.9 We emphasize that the Carleman estimate (2.2.15) depends only on $h_{n}$ and $\partial_{t} h_{n}$ in the last component.

Proof of Theorem 2.8. Let $\omega_{1}$ and $\omega_{2}$ be two subsets of $\omega$ be two open sets such that $\overline{\Gamma_{0}} \subset \partial \omega_{j} \cap \partial \Omega$ for each $j=1,2$ and $\overline{\omega_{1}} \subset \omega_{2}$ and $\overline{\omega_{2}} \subset \omega$.

We start applying the Carleman inequality of Proposition 2.6 to $v_{1}, \ldots, v_{n}$ in system (2.2.14) with $\omega_{0}=\omega_{1}$. We have:

$$
\begin{aligned}
& \sum_{j=1}^{n-1} I\left(\alpha, v_{j}, \Omega\right)+I\left(0, v_{n}, \Omega\right) \\
\leq & C s^{\alpha} \sum_{j=2}^{n} \int_{-T}^{T} \int_{\Omega} e^{2 s \varphi}\left|v_{j}\right|^{2} d x d t+C s^{\alpha} \sum_{j=1}^{n-1} \int_{-T}^{T} \int_{\Omega} e^{2 s \varphi}\left|h_{j}\right|^{2} d x d t+C \int_{-T}^{T} \int_{\Omega} e^{2 s \varphi}\left|h_{n}\right|^{2} d x d t \\
& +C s^{\alpha+1} \sum_{j=1}^{n-1} \int_{-T}^{T} \int_{\omega_{1}} e^{2 s \varphi}\left(s^{2}\left|v_{j}\right|^{2}+\left|\partial_{t} v_{j}\right|^{2}\right) d x d t+C s \int_{-T}^{T} \int_{\omega_{1}} e^{2 s \varphi}\left(s^{2}\left|v_{n}\right|^{2}+\left|\partial_{t} v_{n}\right|^{2}\right) d x d t .
\end{aligned}
$$

Note that the first term of the right-hand side of the inequality above can be absorbed by taking $s$ large enough since $1<\alpha<2$. Therefore, we can rewrite this inequality as follows:

$$
\begin{align*}
& \sum_{j=1}^{n-1} I\left(\alpha, v_{j}, \Omega\right)+I\left(0, v_{n}, \Omega\right) \\
\leq & C s^{\alpha} \sum_{j=1}^{n-1} \int_{-T}^{T} \int_{\Omega} e^{2 s \varphi}\left|h_{j}\right|^{2} d x d t+C \int_{-T}^{T} \int_{\Omega} e^{2 s \varphi}\left|h_{n}\right|^{2} d x d t  \tag{2.2.16}\\
& +C s^{\alpha+1} \sum_{j=1}^{n-1} \int_{-T}^{T} \int_{\omega_{1}} e^{2 s \varphi}\left(s^{2}\left|v_{j}\right|^{2}+\left|\partial_{t} v_{j}\right|^{2}\right) d x d t \\
& +C s \int_{-T}^{T} \int_{\omega_{1}} e^{2 s \varphi}\left(s^{2}\left|v_{n}\right|^{2}+\left|\partial_{t} v_{n}\right|^{2}\right) d x d t .
\end{align*}
$$

Now we are going to estimate the local term of $v_{n}$ and $\partial_{t} v_{n}$ in 2.2.16). To do this, we consider a cut-off function $\xi \in C^{\infty}(\Omega, \mathbb{R})$ such that

$$
\begin{cases}0 \leq \xi \leq 1 & \text { in } \Omega \\ \xi \equiv 1 & \text { in } \omega_{1} \\ \xi \equiv 0 & \text { in } \Omega \backslash \overline{\omega_{2}} .\end{cases}
$$

Using the equation of $v_{n-1}$ in 2.2.14, we see that:

$$
\begin{align*}
s^{3} \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \varphi} \xi\left|v_{n}\right|^{2} d x d t= & s^{3} \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \varphi} \xi v_{n}\left(\square v_{n-1}+r_{n-1} v_{n-1}\right) d x d t  \tag{2.2.17}\\
& -s^{3} \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \varphi} \xi v_{n} h_{n-1} d x d t
\end{align*}
$$

Let us estimate each term of the equation above. First, by Young's inequality for every $\delta>0$, there exists a constant $C=C(\delta)$ such that

$$
\begin{equation*}
s^{3} \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \varphi} \xi h_{n-1} v_{n} d x d t \leq \delta s^{3} \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \varphi}\left|v_{n}\right|^{2} d x d t+C s^{3} \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \varphi}\left|h_{n-1}\right|^{2} d x d t, \tag{2.2.18}
\end{equation*}
$$

On the other hand, integration by parts yields

$$
\begin{align*}
& \quad s^{2} \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \varphi} \xi \partial_{t}^{2} v_{n-1} v_{n} d x d t \\
& =s^{3} \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \varphi} \xi\left(4 s^{2}\left|\partial_{t} \varphi\right|^{2}+2 s \partial_{t}^{2} \varphi\right) v_{n-1} v_{n} d x d t+4 s^{4} \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \varphi} \xi \partial_{t} \varphi v_{n-1} \partial_{t} v_{n} d x d t \\
& \quad+s^{3} \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \varphi} \xi v_{n-1} \partial_{t}^{2} v_{n} d x d t, \tag{2.2.19}
\end{align*}
$$

and

$$
\begin{align*}
& -s^{3} \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \varphi} \xi \Delta v_{n-1} v_{n} d x d t \\
= & 2 s^{4} \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \varphi} \xi \nabla \varphi\left(v_{n} \nabla v_{n-1}-v_{n-1} \nabla v_{n}\right) d x d t-s^{3} \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \varphi} \xi v_{n-1} \Delta v_{n} d x d t \\
& +s^{3} \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \varphi} \nabla \xi\left(v_{n} \nabla v_{n-1}-v_{n-1} \nabla v_{n}\right) d x d t . \tag{2.2.20}
\end{align*}
$$

By (2.2.19, 2.2.20) and $\square v_{n}+r_{n} v_{n}=h_{n}$, we have

$$
\begin{align*}
& s^{3} \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \varphi} \xi v_{n}\left(\square v_{n-1}+r_{n-1} v_{n-1}\right) d x d t \\
\leq & C \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \varphi}\left(s^{3}\left|h_{n-1}\right|^{2}+\left|h_{n}\right|^{2}\right) d x d t  \tag{2.2.21}\\
& +C s^{5} \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \varphi} \xi\left(s^{2}\left|v_{n-1}\right|^{2}+\left|\partial_{t} v_{n-1}\right|^{2}+\left|\nabla v_{n-1}\right|^{2}\right) d x d t+\delta I\left(0, v_{n}, \omega_{2}\right),
\end{align*}
$$

for every $\delta>0$, where we have used the Young inequality. Moreover, by part a) of Lemma 2.4 applied to $v_{n-1}, \omega_{2}$ and $\omega$ with $\gamma=3$ one has

$$
\begin{align*}
s^{5} \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \varphi}\left|\nabla v_{n-1}\right|^{2} d x d t \leq & C s^{2} \int_{-T}^{T} \int_{\omega} e^{2 s \varphi}\left(\left|v_{n}\right|^{2}+\left|h_{n-1}\right|^{2}\right) d x d t \\
& +C s^{5} \int_{-T}^{T} \int_{\omega} e^{2 s \varphi}\left(s^{3}\left|v_{n-1}\right|^{2}+\left|\partial_{t} v_{n-1}\right|^{2}\right) d x d t \tag{2.2.22}
\end{align*}
$$

Substituting (2.2.22) into (2.2.21) and substituting the obtained estimate into (2.2.18), we conclude that

$$
\begin{align*}
& s^{3} \int_{-T}^{T} \int_{\omega_{1}} e^{2 s \varphi}\left|v_{n}\right|^{2} d x d t \\
& \leq  \tag{2.2.23}\\
& \leq \int_{-T}^{T} \int_{\omega} e^{2 s \varphi}\left(s^{3}\left|h_{n-1}\right|^{2}+\left|h_{n}\right|^{2}\right) d x d t \\
& \quad+C s^{8} \int_{-T}^{T} \int_{\omega} e^{2 s \varphi}\left|v_{n-1}\right|^{2} d x d t+C s^{5} \int_{-T}^{T} \int_{\omega} e^{2 s \varphi}\left|\partial_{t} v_{n-1}\right|^{2} d x d t+\delta I\left(0, v_{n}, \omega\right),
\end{align*}
$$

for every $\delta>0$, where we have included the integral term of $\left|v_{n}\right|^{2}$ which has in front $s^{2}$ in $\delta I\left(0, v_{n}, \omega\right)$ by taking $s$ large enough. In the same manner, we can estimate the local term of $\partial_{t} v_{n}$. In fact, let us consider the function $\xi$ defined above. Then,

$$
\begin{align*}
& s \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \varphi} \xi\left|\partial_{t} v_{n}\right|^{2} d x d t \\
= & -s \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \varphi} \xi \partial_{t} v_{n} \partial_{t} h_{n-1} d x d t+s \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \varphi} \xi \partial_{t} v_{n}\left(\square \partial_{t} v_{n-1}+r_{n-1} \partial_{t} v_{n-1}\right) d x d t . \tag{2.2.24}
\end{align*}
$$

First, notice that for all $\delta>0$, there is a positive constant $C=C(\delta)$ such that

$$
\begin{align*}
-s \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \varphi} \xi \partial_{t} v_{n} \partial_{t} h_{n-1} d x d t \leq & \delta s \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \varphi} \xi\left|\partial_{t} v_{n}\right|^{2} d x d t  \tag{2.2.25}\\
& +C(\delta) s \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \varphi} \xi\left|\partial_{t} h_{n-1}\right|^{2} d x d t
\end{align*}
$$

On the other hand,

$$
\begin{align*}
& s \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \varphi} \xi \partial_{t} v_{n}\left(\square \partial_{t} v_{n-1}+r_{n-1} \partial_{t} v_{n-1}\right) d x d t \\
= & s \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \varphi} \xi \partial_{t} v_{n}\left(\partial_{t}^{3} v_{n-1}-\Delta \partial_{t} v_{n-1}+r_{n-1} \partial_{t} v_{n-1}\right) d x d t \\
= & J_{1}+J_{2}+J_{3} . \tag{2.2.26}
\end{align*}
$$

Then, we compute the terms $J_{k}$, with $k=1,2,3$. We start with $J_{1}$. Integration by parts yields

$$
\begin{equation*}
J_{1}=-2 s^{2} \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \varphi} \xi \partial_{t} \varphi \partial_{t} v_{n} \partial_{t}^{2} v_{n-1} d x d t-s \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \varphi} \xi \partial_{t}^{2} v_{n} \partial_{t}^{2} v_{n-1} d x d t \tag{2.2.27}
\end{equation*}
$$

The last term can be estimated by using integration by parts again. After straightforward computations, we get

$$
\begin{align*}
& -s \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \varphi} \xi \partial_{t}^{2} v_{n} \partial_{t}^{2} v_{n-1} d x d t \\
= & -2 s^{2} \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \varphi} \xi\left(2 s\left|\partial_{t} \varphi\right|^{2}+\partial_{t}^{2} \varphi\right) \partial_{t} v_{n-1} \partial_{t} v_{n} d x d t+s \int_{-T}^{T} \int_{\omega_{1}} e^{2 s \varphi} \xi \partial_{t}^{3} v_{n} \partial_{t} v_{n-1} d x d t . \tag{2.2.28}
\end{align*}
$$

Substituting (2.2.28) into 2.2.27 we deduce that

$$
\begin{align*}
J_{1}= & -4 s^{2} \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \varphi} \xi \partial_{t} \varphi \partial_{t} v_{n} \partial_{t}^{2} v_{n-1} d x d t \\
& -2 s^{2} \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \varphi} \xi\left(2 s\left|\partial_{t} \varphi\right|^{2}+\partial_{t}^{2} \varphi\right) \partial_{t} v_{n-1} \partial_{t} v_{n} d x d t+s \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \varphi} \xi \partial_{t}^{3} v_{n} \partial_{t} v_{n-1} d x d t \tag{2.2.29}
\end{align*}
$$

Now, we estimate $J_{2}$. To do this, we proceed integrating by parts in space:

$$
\begin{align*}
J_{2}= & s \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \varphi}\left(2 s \xi \partial_{t} v_{n} \nabla \varphi+\partial_{t} v_{n} \nabla \xi\right) \nabla \partial_{t} v_{n-1} d x d t  \tag{2.2.30}\\
& +s \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \varphi} \xi \nabla \partial_{t} v_{n-1} \cdot \nabla \partial_{t} v_{n} d x d t .
\end{align*}
$$

Moreover, the second term of the right-hand side of 2.2 .30 can be estimated as follows: first we integrate by parts in space to get

$$
\begin{aligned}
& s \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \varphi} \xi \nabla \partial_{t} v_{n-1} \cdot \nabla \partial_{t} v_{n} d x d t \\
= & -s \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \varphi}(2 s \xi \nabla \varphi+\nabla \xi) \nabla \partial_{t} v_{n} \partial_{t} v_{n-1} d x d t \\
& -s \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \varphi} \xi \Delta \partial_{t} v_{n} \partial_{t} v_{n-1} d x d t .
\end{aligned}
$$

However, notice that the term $\nabla \partial_{t} v_{n}$ cannot be absorbed by using the classical Carleman approach. To solve this, we integrate by parts in time the first term of the right-hand side of the above equation and therefore we obtain

$$
\begin{align*}
& s \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \varphi} \xi \nabla \partial_{t} v_{n-1} \cdot \nabla \partial_{t} v_{n} d x d t \\
= & s \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \varphi} \partial_{t}\left(e^{2 s \varphi}(2 s \xi \nabla \varphi+\nabla \xi)\right) \nabla v_{n} \partial_{t} v_{n-1} d x d t \\
& +s \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \varphi}\left[(2 s \xi \nabla \varphi+\nabla \xi) \cdot \nabla v_{n-1}\right] \partial_{t}^{2} v_{n-1} d x d t  \tag{2.2.31}\\
& -s \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \varphi} \xi \Delta \partial_{t} v_{n} \partial_{t} v_{n-1} d x d t .
\end{align*}
$$

Substituting (2.2.31) into 2.2 .30 we get the following estimate for $J_{2}$ :

$$
\begin{align*}
J_{2}= & s \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \varphi}\left(2 s \xi \partial_{t} v_{n} \nabla \varphi+\partial_{t} v_{n} \nabla \xi\right) \cdot \nabla \partial_{t} v_{n-1} d x d t \\
& +s \int_{-T}^{T} \int_{\omega_{2}} \partial_{t}\left[e^{2 s \varphi}(2 s \xi \nabla \varphi+\nabla \xi)\right] \cdot \nabla v_{n} \partial_{t} v_{n-1} d x d t \\
& +s \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \varphi}(2 s \xi \nabla \varphi+\nabla \xi) \nabla v_{n-1} \partial_{t}^{2} v_{n-1} d x d t  \tag{2.2.32}\\
& -s \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \varphi} \xi \Delta \partial_{t} v_{n} \partial_{t} v_{n-1} d x d t .
\end{align*}
$$

On the other hand, $J_{3}$ can be written as follows:

$$
\begin{align*}
J_{3} & =s \int_{-T}^{T} \int_{\omega_{1}} e^{2 s \varphi} \xi r_{n-1} \partial_{t} v_{n-1} \partial_{t} v_{n} d x d t \\
& =s \int_{-T}^{T} \int_{\omega_{1}} e^{2 s \varphi} \xi r_{n} \partial_{t} v_{n-1} \partial_{t} v_{n} d x d t+s \int_{-T}^{T} \int_{\omega_{1}} e^{2 s \varphi} \xi\left(r_{n-1}-r_{n}\right) \partial_{t} v_{n-1} \partial_{t} v_{n} d x d t . \tag{2.2.33}
\end{align*}
$$

Thus, substituting (2.2.29) 2.2.32 and 2.2.33 into (2.2.26 and using Young's inequality and the fact that $\partial_{t}^{3} v_{n}-\Delta \partial_{t} v_{n}+r_{n} \partial_{t} v_{n}=\partial_{t} h_{n}$ we deduce that for all $\delta>0$, there exists a positive constant $C=C(\delta)$ such that

$$
\begin{align*}
& s \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \varphi} \xi \partial_{t} v_{n}\left(\square \partial_{t} v_{n-1}+r_{n-1} \partial_{t} v_{n-1}\right) d x d t \\
= & C(\delta) \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \varphi}\left(s^{2+\gamma^{*}}\left|\partial_{t} v_{n-1}\right|^{2}+\left|\nabla v_{n-1}\right|^{2}+\left|\partial_{t}^{2} v_{n-1}\right|^{2}\right) d x d t  \tag{2.2.34}\\
& +C(\delta) s^{-\gamma^{*}} \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \varphi}\left|\partial_{t} h_{n}\right|^{2} d x d t+\delta s \int_{-T}^{T} \int_{\omega_{1}} e^{2 s \varphi}\left|\partial_{t} v_{n}\right|^{2}+\left|\nabla v_{n}\right|^{2} d x d t .
\end{align*}
$$

for every $\gamma^{*}>0$. Furthermore, the local term $\nabla v_{n}$ in $\omega_{1} \times(-T, T)$ can be estimated by using the technical lemmas introduced in the above section and the weighted Poincaré inequality as follows:

$$
\begin{align*}
& \quad s^{3} \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \varphi}\left|\nabla v_{n-1}\right|^{2} d x d t \\
& \leq C s^{3} \int_{-T}^{T} \int_{\omega} e^{2 s \varphi}\left(s^{2}\left|v_{n-1}\right|^{2}+\left|\partial_{t} v_{n-1}\right|^{2}\right) d x d t  \tag{2.2.35}\\
& \quad+C s^{2} \int_{-T}^{T} \int_{\omega} e^{2 s \varphi}\left(\left|v_{n}\right|^{2}+\left|h_{n-1}\right|^{2}\right) d x d t
\end{align*}
$$

Substituting (2.2.35) into 2.2.34 we have

$$
\begin{align*}
& s \int_{-T}^{T} \int_{\omega_{1}} e^{2 s \varphi} \xi \partial_{t} v_{n}\left(\square \partial_{t} v_{n-1}+r_{n-1} \partial_{t} v_{n-1}\right) d x d t \\
\leq & C(\delta) s^{3} \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \varphi}\left(s^{2+\gamma^{*}}\left|\partial_{t} v_{n-1}\right|^{2}+\left|\partial_{t}^{2} v_{n-1}\right|^{2}\right) d x d t \\
& +C(\delta) s^{-\gamma^{*}} \int_{-T}^{T} \int_{\omega} e^{2 s \varphi}\left|\partial_{t} h_{n-1}\right|^{2} d x d t+\delta \int_{-T}^{T} \int_{\omega} e^{2 s \varphi}\left(\left|\partial_{t} v_{n}\right|^{2}+\left|\nabla v_{n}\right|^{2}\right) d x d t  \tag{2.2.36}\\
& +C s^{2} \int_{-T}^{T} \int_{\omega} e^{2 s \varphi}\left|h_{n-1}\right|^{2} d x d t+C s^{2} \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \varphi}\left|v_{n}\right|^{2} d x d t .
\end{align*}
$$

for each $\delta>0$ and $\gamma^{*}>0$. Finally substituting 2.2.25 (2.2.36) into 2.2.24) and using $\xi=1$ in $\omega_{0}$ we have

$$
\begin{align*}
& s \int_{-T}^{T} \int_{\omega_{1}} e^{2 s \varphi}\left|\partial_{t} v_{n}\right|^{2} d x d t \\
\leq & C(\delta) s^{3} \int_{-T}^{T} \int_{\omega} e^{2 s \varphi}\left(s^{2}\left|v_{n-1}\right|^{2}+s^{2+\gamma^{*}}\left|\partial_{t} v_{n-1}\right|^{2}\left|\partial_{t}^{2} v_{n-1}\right|^{2}\right) d x d t  \tag{2.2.37}\\
& +C(\delta) \int_{-T}^{T} \int_{\omega} e^{2 s \varphi}\left(s^{2}\left|h_{n-1}\right|^{2}+s\left|\partial_{t} h_{n-1}\right|^{2}+s^{-\gamma^{*}}\left|\partial_{t} h_{n}\right|^{2}\right) d x d t \\
& +\delta s \int_{-T}^{T} \int_{\omega} e^{2 s \varphi}\left(d x d t\left|\partial_{t} v_{n}\right|^{2}+\left|\nabla v_{n}\right|^{2}\right) .
\end{align*}
$$

Finally, by substituting (2.2.23) and (2.2.37) into (2.2.16), by taking the Carleman parameter $s \geq 1$ large enough and by choosing $\delta>0$ sufficiently small, we obtain

$$
\begin{aligned}
& \sum_{j=1}^{n-1} I\left(\alpha, v_{j}, \Omega\right)+I\left(0, v_{n}, \Omega\right) \\
\leq & C s^{\alpha} \sum_{j=1}^{n-2} \int_{-T}^{T} \int_{\Omega} e^{2 s \varphi}\left|h_{j}\right|^{2} d x d t+C \int_{-T}^{T} \int_{\Omega} e^{2 s \varphi}\left(s^{3}\left|h_{n-1}\right|^{2}+\left|h_{n}\right|^{2}\right) d x d t \\
& +C \int_{-T}^{T} \int_{\Omega} e^{2 s \varphi}\left(s\left|\partial_{t} h_{n-1}\right|^{2}+\left|\partial_{t} h_{n}\right|^{2}\right) d x d t \\
& +C s^{\alpha+1} \sum_{j=1}^{n-2} \int_{-T}^{T} \int_{\omega} e^{2 s \varphi}\left(s^{2}\left|v_{j}\right|^{2}+\left|\partial_{t} v_{j}\right|^{2}\right) d x d t \\
& +C s^{3} \int_{-T}^{T} \int e^{2 s \varphi}\left(s^{5}\left|v_{n-1}\right|^{2}+s^{3}\left|\partial_{t} v_{n-1}\right|^{2}+\left|\partial_{t}^{2} v_{n-1}\right|^{2}\right) d x d t
\end{aligned}
$$

which completes the proof of Theorem 2.8 .

### 2.3 Proof of Theorem 2.1

The plan of the proof of Theorem 2.1 contains three parts:
Step 1 In the same spirit of the Bukhgeim-Klibanov method, we rewrite appropriately system 2.1.1 to apply the estimate 2.2.15 in Theorem 2.8.
Step 2 After applying the Carleman estimate of Theorem 2.8 to the new system, we estimate the residual and source terms.

Step 3 We conclude the proof gathering the estimates of the previous steps and eliminating the small order terms.

## - Step 1: Setting

For each $j=1, \ldots, n$, let us denote by $y_{j}=u_{j}-\tilde{u}_{j}, p_{j}=q_{j}, f_{j}=q_{j}-\tilde{q}_{j}$ and $R_{j}=\tilde{u}_{j}$.

Then, following this notation, $y_{1}, \ldots, y_{n}$ solves:

$$
\begin{cases}\square y_{1}+p_{1} y_{1}=y_{2}+f_{1} R_{1}, & \text { in } \Omega \times(0, T),  \tag{2.3.1}\\ \square y_{2}+p_{2} y_{2}=y_{2}+f_{2} R_{2}, & \text { in } \Omega \times(0, T), \\ \quad \vdots & \vdots \\ \square y_{n-1}+p_{n-1} y_{n-1}=y_{n}+f_{n-1} R_{n-1}, & \text { in } \Omega \times(0, T), \\ \square y_{n}+p_{n} y_{n}=f_{n} R_{n}, & \text { in } \Omega \times(0, T), \\ \partial_{t}^{k} y(0)=0, k=0,1, j=1, \ldots, n, & \text { in } \Omega, \\ y_{j}=0, j=1, \ldots, n, & \text { on } \partial \Omega \times(0, T)\end{cases}
$$

For each $j=1, \ldots, n$, we set $w_{j}=\partial_{t}^{2} y_{j}$. Then, the new variables solve the following system:

$$
\begin{cases}\square w_{1}+p_{1} w_{1}=w_{2}+f_{1} \partial_{t}^{2} R_{1}, & \text { in } \Omega \times(0, T),  \tag{2.3.2}\\ \square w_{2}+p_{2} w_{2}=w_{3}+f_{2} \partial_{t}^{2} R_{2}, & \text { in } \Omega \times(0, T), \\ \quad \vdots & \vdots \\ \square w_{n-1}+p_{n-1} w_{n-1}=w_{n}+f_{n-1} \partial_{t}^{2} R_{n-1}, & \text { in } \Omega \times(0, T), \\ \square w_{n}+p_{n} w_{n}=f_{n} \partial_{t}^{2} R_{n}, & \text { in } \Omega \times(0, T), \\ \partial_{t}^{k} w_{j}(0)=f_{j} \partial_{t}^{k} R_{j}(0), k=0,1, j=1, \ldots, n, & \text { in } \Omega, \\ w_{j}=0, j=1, \ldots, n, & \text { on } \partial \Omega \times(0, T) .\end{cases}
$$

Now, we want to apply Theorem 2.8 to a suitable system. In order to do that, we extend system (2.3.2) in an even way, setting $w_{j}(x, t)=w_{j}(x,-t)$ for all $(x, t) \in \Omega \times(-T, 0)$. We also extend the functions $R_{j}, \partial_{t} R_{j}$ and $\partial_{t}^{2} R_{j}$ in an even way and keep the same notations for the new system.

To be able to apply the Carleman estimate (2.2.15), the functions $w_{j}$ must satisfy $\partial_{t}^{k} w( \pm T)=0$ in $\Omega$, for $k=0,1$. However, this condition does not hold. To avoid this difficulty, we consider a cut-off function $\theta \in C_{c}^{\infty}((-T, T), \mathbb{R})$ defined as follows:

$$
\begin{cases}0 \leq \theta \leq 1, & \text { in }(-T, T) \\ \theta \equiv 1, & \text { in }(-T+\tau, T-\tau)\end{cases}
$$

According to the definition of $\theta$, it is clear that $z_{j}=\theta w_{j}$, for $j=1, \ldots, n$, solves

$$
\begin{cases}\square z_{1}+p_{1} z_{1}=z_{2}+F_{1}, & \text { in } \Omega \times(-T, T),  \tag{2.3.3}\\ \square z_{2}+p_{2} z_{2}=z_{3}+F_{2}, & \text { in } \Omega \times(-T, T), \\ \quad \vdots & \vdots \\ \square z_{n-1}+p_{n-1} z_{n-1}=z_{n}+F_{n-1}, & \text { in } \Omega \times(-T, T), \\ \square z_{n}+p_{n} z_{n}=F_{n}, & \text { in } \Omega \times(-T, T), \\ \partial_{t}^{k} z(0)=f_{j} \partial_{t}^{k} R(0), k=0,1, j=1, \ldots, n, & \text { in } \Omega, \\ \partial_{t}^{k} z( \pm T)=0, k=0,1,2, j=1, \ldots, n, & \text { in } \Omega, \\ z_{j}=0, j=1, \ldots, n, & \text { on } \partial \Omega \times(-T, T)\end{cases}
$$

Here, the functions $F_{j}$ are defined by

$$
F_{j}=\theta f_{j} \partial_{t}^{2} R_{j}+\partial_{t}^{2} \theta w_{j}+2 \partial_{t} \theta \partial_{t} w_{j}, \quad \text { in } \Omega \times(-T, T),
$$

for each $j=1, \ldots, n$.

- Step 2: Applying Carleman estimate for hyperbolic systems In this step, we denote by $C$ a generic positive constant which depends at least of $\Gamma_{0}, m, T, \omega$ and $x_{0}$ and may change from line to line.

Applying the Carleman estimate of Theorem 2.8 to the system 2.3.3 with $v_{j}=z_{j}$, $r_{j}=p_{j}$ and $h_{j}=F_{j}$, we see that

$$
\begin{align*}
& \sum_{j=1}^{n-1} I\left(\alpha, z_{j}, \Omega\right)+I\left(0, z_{n}, \Omega\right) \\
\leq & C s \sum_{j=1}^{n-2} \int_{-T}^{T} \int_{\Omega} e^{2 s \varphi}\left|F_{j}\right|^{2} d x d t+C \int_{-T}^{T} \int_{\Omega} e^{2 s \varphi}\left(s^{3}\left|F_{n-1}\right|^{2}+\left|F_{n}\right|^{2}\right) d x d t \\
& +C \int_{-T}^{T} \int_{\Omega} e^{2 s \varphi}\left(s\left|\partial_{t} F_{n-1}\right|^{2}+\left|\partial_{t} F_{n}\right|^{2}\right) d x d t  \tag{2.3.4}\\
& +C s^{\alpha+1} \sum_{j=1}^{n-2} \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \varphi}\left(s^{2}\left|z_{j}\right|^{2}+\left|\partial_{t} z_{j}\right|^{2}\right) d x d t \\
& +C s^{3} \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \varphi}\left(s^{5}\left|z_{n-1}\right|^{2}+s^{3}\left|\partial_{t} z_{n-1}\right|^{2}+\left|\partial_{t}^{2} z_{n-1}\right|^{2}\right) d x d t
\end{align*}
$$

Note that the assumption (2.1.4) implies

$$
c<\left|R_{j}(0)\right|^{2}, \quad \forall j=1,2, \ldots, n
$$

Then, the following estimate holds:

$$
\begin{equation*}
c \int_{\Omega} e^{2 s \varphi(0)}\left|f_{j}\right|^{2} d x \leq \int_{\Omega} e^{2 s \varphi(0)}\left|f_{j} R_{j}(0)\right|^{2} d x=\int_{\Omega} e^{2 s \varphi(0)}\left|z_{j}(0)\right|^{2} d x \tag{2.3.5}
\end{equation*}
$$

Hence, summing 2.3.5 over $j$, we deduce that

$$
\frac{1}{c} \sum_{j=1}^{n-1} I\left(\alpha, z_{j}\right)+\frac{1}{c} I\left(0, z_{n}\right) \geq s^{\alpha+2} \sum_{j=1}^{n-1} \int_{\Omega} e^{2 s \varphi(0)}\left|f_{j}\right|^{2} d x+s^{2} \int_{\Omega} e^{2 s \varphi(0)}\left|f_{n}\right|^{2} d x
$$

Now we estimate the global terms of $F_{j}$ and its derivatives, for each $j=1, \ldots, n$. By
definition,

$$
\begin{align*}
& \int_{-T}^{T} \int_{\Omega} e^{2 s \varphi}\left|F_{j}\right|^{2} d x d t \\
\leq & 2 \int_{-T}^{T} \int_{\Omega} e^{2 s \varphi}\left|\theta f_{j} \partial_{t}^{2} R_{j}\right|^{2} d x d t+2 \int_{-T}^{T} \int_{\Omega} e^{2 s \varphi}\left|2 \partial_{t} \theta \partial_{t} w_{j}+\partial_{t}^{2} \theta w_{j}\right|^{2} d x d t  \tag{2.3.6}\\
\leq & C \int_{\Omega} e^{2 s \varphi(0)}\left|f_{j}\right|^{2} d x+2 \int_{-T}^{T} \int_{\Omega} e^{2 s \varphi}\left|2 \partial_{t} \theta \partial_{t} w_{j}+\partial_{t}^{2} \theta w_{j}\right|^{2} d x d t .
\end{align*}
$$

Now, we focus our attention on estimating the global term of $2 \partial_{t} \theta \partial_{t} w_{j}+\partial_{t}^{2} \theta w_{j}$, for $j=1, \ldots, n$. Notice that if the Time condition (2.1.3) holds, the Carleman weight $\psi$ defined in (2.2.9) satisfies

$$
\psi(x, \pm T)=\left|x-x_{0}\right|^{2}-\beta T^{2}+C_{0}<C_{0}, \quad \text { in } \Omega .
$$

Then, we choose $\tau>0$ such that

$$
\psi(x, t) \leq C_{0}, \quad \text { in } \Omega \times([-T,-T+\tau] \cup[T-\tau, T]),
$$

and therefore,

$$
\varphi(x, t)=e^{\lambda C_{0}}<e^{\lambda \psi(x, 0)}=\varphi(x, 0), \quad \text { in } \Omega \times([-T,-T+\tau] \cup[T-\tau, T])
$$

Since the derivatives of $\theta$ vanish in $[-T+\tau, T-\tau]$ we see that

$$
\begin{aligned}
& \int_{-T}^{T} \int_{\Omega} e^{2 s \varphi}\left|2 \partial_{t} \theta \partial_{t} w_{j}+\partial_{t}^{2} \theta w_{j}\right|^{2} d x d t \\
\leq & C\left(\int_{-T}^{-T+\tau}+\int_{T-\tau}^{T}\right) e^{2 s e^{\lambda C_{0}}} \int_{\Omega}\left(\left|\partial_{t} w_{j}\right|^{2}+\left|w_{j}\right|^{2}\right) d x d t .
\end{aligned}
$$

Now, we will estimate the last term of the above inequality. To do this, we will use the following energy estimates of 2.3 .2 :

$$
\int_{\Omega}\left|\partial_{t} w_{j}(t)\right|^{2} d x+\int_{\Omega}\left|\nabla w_{j}(t)\right|^{2} d x \leq C \int_{\Omega}\left|f_{j}(t)\right|^{2} d x+C \int_{\Omega}\left|w_{j+1}(t)\right|^{2} d x, \quad \forall t \in(-T, T),
$$

for each $j=1, \ldots, n$ and

$$
\int_{\Omega}\left|\partial_{t} w_{n}(t)\right|^{2} d x+\int_{\Omega}\left|\nabla w_{n}(t)\right|^{2} d x \leq C \int_{\Omega}\left|f_{n}\right|^{2} d x, \forall t \in(-T, T)
$$

where have used that $R_{j} \in H^{2}\left(-T, T ; L^{\infty}(\Omega)\right)$ for each $j=1, \ldots, n$. Integrating on $(-T,-T+\tau) \cup(T-\tau, T)$ the estimate above and using the Poincaré inequality to $w_{j}$, we see that

$$
\begin{aligned}
& \left(\int_{-T}^{-T+\tau}+\int_{T-\tau}^{T}\right) e^{2 s e^{\lambda C_{0}}} \int_{\Omega}\left(\left|\partial_{t} w_{j}\right|^{2}+\left|w_{j}\right|^{2}\right) d x d t \\
\leq & \int_{\Omega} e^{2 s \varphi(0)}\left|f_{j}\right|^{2} d x+C e^{s e^{\lambda C_{0}}} \int_{-T}^{T} \int_{\Omega}\left|w_{j+1}\right|^{2} d x d t
\end{aligned}
$$

for each $j=1, \ldots, n-1$. Furthermore, due to the structure in cascade of system (2.3.3), we obtain

$$
e^{s e^{\lambda C_{0}}} \int_{-T}^{T} \int_{\Omega} e^{2 s \varphi}\left|w_{j+1}\right|^{2} d x d t \leq C e^{s e^{\lambda C_{0}}} \sum_{i=j+1}^{n} \int_{\Omega}\left|f_{i}\right|^{2} d x \leq C \sum_{i=j+1}^{n} \int_{\Omega} e^{2 s \varphi(0)}\left|f_{i}\right|^{2} d x
$$

for each $j=2, \ldots, n-1$. Therefore, for every $j=1, \ldots, n$, we deduce that

$$
\begin{equation*}
\int_{-T}^{T} \int_{\Omega} e^{2 s \varphi}\left|2 \partial_{t} \theta \partial_{t} w_{j}+\partial_{t}^{2} \theta w_{j}\right|^{2} d x d t \leq C \sum_{i=j}^{n} \int_{\Omega} e^{2 s \varphi(0)}\left|f_{i}\right|^{2} d x \tag{2.3.7}
\end{equation*}
$$

Substituting (2.3.7) in (2.3.6), we see that

$$
\begin{align*}
& s^{\alpha} \sum_{j=1}^{n-2} \int_{-T}^{T} \int_{\Omega} e^{2 s \varphi}\left|F_{j}\right|^{2} d x d t+\int_{-T}^{T} \int_{\Omega} e^{2 s \varphi}\left(s^{3}\left|F_{n-1}\right|^{2}+\left|F_{n}\right|^{2}\right) d x d t \\
& \quad \leq C s^{\alpha} \sum_{j=1}^{n-2} \int_{\Omega} e^{2 s \varphi(0)}\left|f_{j}\right|^{2} d x+C \int_{\Omega} e^{2 s \varphi(0)}\left(s^{3}\left|f_{n-1}\right|^{2}+\left|f_{n}\right|^{2}\right) d x \tag{2.3.8}
\end{align*}
$$

In the same manner we can see that

$$
\begin{equation*}
\int_{-T}^{T} \int_{\Omega} e^{2 s \varphi}\left(s\left|\partial_{t} F_{n-1}\right|^{2}+\left|\partial_{t} F_{n}\right|^{2}\right) d x d t \leq C \int_{-T}^{T} \int_{\Omega} e^{2 s \varphi}\left(s\left|f_{n-1}\right|^{2}+\left|f_{n}\right|^{2}\right) d x d t \tag{2.3.9}
\end{equation*}
$$

Thus, substituting (2.3.8) and (2.3.9) into (2.3.4), and taking $s$ large enough, we have

$$
\begin{align*}
& \quad s^{\alpha+2} \sum_{j=1}^{n-1} \int_{\Omega} e^{2 s \varphi(0)}\left|f_{j}\right|^{2} d x+s^{2} \int_{\Omega}\left|f_{n}\right|^{2} d x \\
& \leq  \tag{2.3.10}\\
& \leq s^{\alpha+1} \sum_{j=1}^{n-2} \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \varphi}\left(s^{2}\left|z_{j}\right|^{2}+\left|\partial_{t} z_{j}\right|^{2}\right) d x d t \\
& \quad+C s^{3} \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \varphi}\left(s^{5}\left|z_{n-1}\right|^{2}+s^{3}\left|\partial_{t} z_{n-1}\right|^{2}+\left|\partial_{t}^{2} z_{n-1}\right|^{2}\right) d x d t .
\end{align*}
$$

## - Step 3: Last arrangements and conclusion

From (2.3.10), we fix the parameter $s$ and put it into the constant $C$ :

$$
\begin{aligned}
\sum_{j=1}^{n} \int_{\Omega}\left|f_{j}\right|^{2} d x \leq & C \sum_{j=1}^{n} \int_{-T}^{T} \int_{\omega_{\omega}} e^{2 s \varphi}\left(\left|z_{j}\right|^{2}+\left|\partial_{t} z_{j}\right|^{2}\right) d x d t \\
& +C \int_{-T}^{T} \int_{\omega}\left(\left|z_{n-1}\right|^{2}+\left|\partial_{t} z_{n-1}\right|^{2}+\left|\partial_{t}^{2} z_{n-1}\right|^{2}\right) d x d t
\end{aligned}
$$

where we have used that the Carleman weights defined in (2.2.9) are bounded. Moreover, by definition of each $z_{j}$ we see that

$$
\begin{equation*}
\sum_{j=1}^{n}\left\|f_{j}\right\|_{L^{2}(\Omega)}^{2} \leq C \sum_{j=1}^{n-1}\left\|y_{j}\right\|_{H^{3}\left(-T, T ; L^{2}(\omega)\right)}^{2}+\left\|y_{n}\right\|_{H^{4}\left(-T, T ; L^{2}(\omega)\right)}^{2} \tag{2.3.11}
\end{equation*}
$$

Finally, replacing $f_{j}=q_{j}-\tilde{q}_{j}$ and $y_{j}=u_{j}-\tilde{u}_{j}$ by (2.3.11) for each $j=1, \ldots, n$, we conclude the proof of Theorem 2.1.

## Chapter 3

## Controllability properties of a class of heat equations with dynamic boundary conditions

### 3.1 Introduction and main results

In this chapter, the null controllability property for a suitable class of parabolic equations with dynamic boundary conditions is studied. In order to state the main results of this article, we shall introduce some notation. Let $\Omega_{L}=\left(-L_{1}, 0\right)$ be a bounded interval with $L_{1}>0, \omega \subset \Omega_{L}$, and $T>0$. Then, let us consider $\left(u, u_{\Gamma}\right) \in L^{2}\left(\Omega_{L} \times(0, T)\right) \times L^{2}(0, T)$ be a solution of

$$
\begin{cases}\partial_{t} u(x, t)-\partial_{x}^{2} u(x, t)=\chi_{\omega}(x) v(x, t), & \forall(x, t) \in \Omega_{L} \times(0, T),  \tag{3.1.1}\\ \left(u(x, 0), u_{\Gamma}(0)\right)=\left(u_{0}(x), u_{\Gamma, 0}\right), & \forall x \in \Omega_{L}, \\ u_{\Gamma}(t)=u(0, t), & \forall t \in(0, T), \\ u\left(-L_{1}, t\right)=0, & \forall t \in(0, T), \\ u_{\Gamma}^{\prime}(t)+\partial_{x} u(0, t)=0, & \forall t \in(0, T) .\end{cases}
$$

Here, the pair $\left(u, u_{\Gamma}\right)$ stands for the state of (3.1.1), $\left(u_{0}, u_{0, \Gamma}\right) \in L^{2}\left(\Omega_{L}\right) \times \mathbb{R}$ is the initial state and $v \in L^{2}(\omega \times(0, T))$ denotes the control acting on $\omega \subset \Omega_{L}$. Notice that (3.1.1) can be treated as a coupled system of dynamic equations for $u$ and $u_{\Gamma}$ with side condition $\left.u\right|_{\Gamma}=u_{\Gamma}$ at the boundary $x=0$. Moreover, if $u$ is smooth enough, then $u_{\Gamma}^{\prime}(t)=\partial_{t} u(0, t)$.

In this chapter, we analyze the null controllability of (3.1.1), i.e., given any data $\left(u_{0}, u_{0, \Gamma}\right) \in L^{2}\left(\Omega_{L}\right) \times \mathbb{R}$ and $T>0$, we want to find a control $v \in L^{2}(\omega \times(0, T))$ such that the associated solution satisfies

$$
u(x, T)=0, \quad \forall x \in \overline{\Omega_{L}} .
$$

In other words, the goal is to steer the state of (3.1.1) to a null final target by a suitable choice of the control function. In addition, we point out that in our model the control is applied in a (small) subset of $\Omega_{L}$. This means that the first equation is controlled directly
by the action of $v$, while the ODE at $x=0$ is being controlled indirectly, through the coupling.

Sometimes, boundary conditions like $(3.1 .1)_{5}$ are called of Wentzell type and in the unidimensional case has the general form $\partial_{t} u-a \Delta u+b \partial_{\nu} u=0$ with $a \geq 0$ and $b \in \mathbb{R}$. In our case, dynamic surface diffusion contributions are neglected (i.e. $a=0$ ).

Physically, equation (3.1.1 $5_{5}$ can be viewed as a transport equation acting on a neighborhood of the boundary at $x=0$. Then, the unidirectionally heat wave travels into the region $\Omega_{L}$ and this wave lives only for an infinitesimally short time. Of course, once the heat wave is inside the region, diffusion is the primary process. For a complete description of physical interpretation for Wentzell boundary conditions in linear and nonlinear models we refer to [31, [57, [53].

Parabolic models with general Wentzell boundary conditions were introduced in the context of the heat equation by A. Favini, G. Goldstein, J. Goldstein and S. Romanelli 44 and subsequently have been intensely studied in the last two decades by many authors, see for instance [54] [94, [95], [55], [52] and the references therein.

The study of controllability properties of parabolic equations are well-known in the case of Dirichlet and Neumann boundary conditions (see for instance [50, [78], [48]), as well as for Robin or Fourier boundary conditions [38], [46], 47]. Moreover, controllability properties of parabolic equations with discontinuous diffusion coefficients are recently studied in [39], [19], [20], [22], [21] and [82].

However, to the best of our knowledge, there are a few work concerning controllability for parabolic equations with dynamic boundary conditions. In particular, optimal control problem and approximate controllability have been considered in [59] and [13] in the case of global controls, i.e, $\omega=\Omega$. Moreover, in 73 the authors studied approximate controllability of a one-dimensional heat equation with dynamical boundary conditions by using the theory of one-sided coupled operator matrices developed by K.-J. Engel in [40].

Recently, in [81] null controllability for parabolic equations with dynamic boundary conditions with surface diffusion (i.e. $a>0$ ) is studied. In particular, the authors consider Generalized Wentzell conditions with surface diffusion acting on the boundary. In this case, the result was proved by applying Carleman estimates for the homogeneous dual problem. The authors used the weight functions defined in the work of O. Yu Imanuvilov et al [62] (see also [61] and [48]). In this context, the novelty relies in the fact that several boundary terms appears in the deduction of the Carleman estimate. Some of these enter in the final estimate, a few cancel, and others can be controlled using the smoothing effect of the surface diffusion. Let us also mention the papers [60] and [24] where the authors studied the local null controllability of some two-dimensional fluid-structure interaction problems where parabolic equations are coupled with some dynamics (typically an ODE) in a part of the boundary. Again, the main ingredient in the proof of this result is a suitable Carleman estimate, which will be applied to the corresponding adjoint system.

In our case, we obtain a null controllability result for 3.1.1 when the control region is far from the right-hand side of the interval $\Omega_{L}$ even if initial data $u_{0}$ and $u_{\Gamma, 0}$ are not
related. More precisely, the main result of this paper is the following
Theorem 3.1 Suppose that $\omega=\left(-L_{1},-a\right)$ with $a>0$. Then, the system defined by (3.1.1) is null-controllable at time $T>0$.

Using the well-known equivalence between null controllability and observability (see e.g. [48, [79]), the proof of Theorem 3.1 consists in obtaining a suitable observability inequality for the corresponding adjoint system. This will be done by obtaining an auxiliary Carleman inequality.

In order to state the second main result of this paper, let $\Omega_{R}=(0,1)$ and set $\Omega=$ $\left(-L_{1}, 1\right)$. From now on, we shall use the following notation: for a function $h: \Omega \rightarrow \mathbb{R}, h_{L}$ and $h_{R}$ stand for the restriction of $\Omega_{L}$ and $\Omega_{R}$, respectively.

Now, let $\left(u, u_{\Gamma}\right) \in L^{2}\left(\Omega_{L} \times(0, T)\right) \times L^{2}(0, T)$ be a solution of

$$
\begin{cases}\partial_{t} y_{L}(x, t)-\partial_{x}^{2} y_{L}(x, t)=f(x, t), & \forall(x, t) \in \Omega_{L} \times(0, T),  \tag{3.1.2}\\ y_{R}(x, t)=y_{L}(0, t)=y_{\Gamma}(t), & \forall(x, t) \in \Omega_{R} \times(0, T), \\ \left(y(x, 0), y_{\Gamma}(0)\right)=\left(y_{0}(x), y_{\Gamma, 0}\right), & \forall x \in \Omega_{L}, \\ y_{L}\left(-L_{1}, t\right)=0, & \forall t \in(0, T), \\ y_{\Gamma}^{\prime}(t)+\partial_{x} y_{L}(0, t)=0, & \forall t \in(0, T) .\end{cases}
$$

We emphasize that (3.1.2) is an extension of (3.1.1) where $u_{R}$ is only a time-dependent function in $\Omega_{R}$. This problem is well-posed in the sense of Hadamard. In addition, thanks to Theorem 3.1 problem (3.1.2) is null-controllable at time $T>0$.

As we mentioned above, the problem (3.1.2) appears as the limit case of a heat equation with discontinuous diffusion coefficient of the form

$$
\sigma^{K}(x)= \begin{cases}1, & \text { if } x \in \Omega_{L}  \tag{3.1.3}\\ K^{2}, & \text { if } x \in \Omega_{R}\end{cases}
$$

with $K \geq 1$. More precisely, the second result of this paper is the following:
Theorem 3.2 Let $u_{0} \in L^{2}(\Omega), \omega \subset \Omega, K>0$ and $T>0$. For $v^{k} \in L^{2}(\omega \times(0, T))$, let $u^{K}$ be the solution of

$$
\begin{cases}\partial_{t} u^{K}-\partial_{x}\left(\sigma^{K} \partial_{x} u^{K}\right)=\chi_{\omega} v^{K}, & \forall(x, t) \in \Omega \times(0, T),  \tag{3.1.4}\\ u^{K}(x, 0)=u_{0}(x), & \forall x \in \Omega, \\ u^{K}\left(-L_{1}, t\right)=\partial_{x} u^{K}\left(L_{2}, t\right)=0, & \forall t \in(0, T)\end{cases}
$$

Suppose that

$$
\begin{equation*}
v^{K} \rightharpoonup v \text { weakly in } L^{2}(\omega \times(0, T)) . \tag{3.1.5}
\end{equation*}
$$

Then, there exists a subsequence of the associated family of solutions $\left(u^{K}\right)_{K>0}$ defined
by (3.1.4) with initial data $u_{0} \in L^{2}(\Omega)$ which converges as $K \rightarrow+\infty$ in the sense that

$$
\begin{align*}
u^{K} \rightharpoonup u & \text { weakly in } L^{2}\left(0, T ; H^{1}(\Omega)\right)  \tag{3.1.6}\\
u^{k} \rightharpoonup u & \text { weakly-star in } L^{\infty}\left(0, T ; L^{2}(\Omega)\right),  \tag{3.1.7}\\
u^{K} \rightarrow u & \text { strongly in } L^{2}(\Omega \times(0, T)) \tag{3.1.8}
\end{align*}
$$

Moreover, $u$ is a weak solution of (3.1.2) with source term $f=\chi_{\omega} v$ and initial condition

$$
y_{L}(x, 0)=u_{0, L}(x), \quad \forall x \in \Omega_{L}, \quad y_{\Gamma}(0)=y_{R}(x, 0)=\int_{\Omega_{R}} u_{0, R} d x, \quad \forall x \in \Omega_{R}
$$

On the other hand, another interesting question concerning models (3.1.2) and (3.1.4) can be considered. In fact, for each $K \geq 1$, let $u^{K}$ be the solution of

$$
\begin{cases}\partial_{t} u_{L}^{K}(x, t)-\partial_{x}^{2} u_{L}^{K}(x, t)=\chi_{\omega}(x) v(x, t), & \forall(x, t) \in \Omega_{L} \times(0, T),  \tag{3.1.9}\\ \partial_{t} u_{R}^{K}(x, t)-K^{2} \partial_{x}^{2} u_{R}^{K}(x, t)=0, & \forall(x, t) \in \Omega_{R} \times(0, T), \\ u_{L}^{K}(x, 0)=u_{0, L}(x), & \forall x \in \Omega_{L}, \\ u_{R}^{K}(x, 0)=u_{0, R}(x), & \forall x \in \Omega_{R}, \\ u_{R}^{K}\left(0^{+}, t\right)=u_{L}^{K}\left(0^{-}, t\right), & \forall t \in(0, T), \\ K^{2} \partial_{x} u_{R}\left(0^{+}, t\right)=\partial_{x} u_{L}\left(0^{-}, t\right), & \forall t \in(0, T), \\ u_{L}^{K}\left(-L_{1}, t\right)=\partial_{x} u_{R}^{K}(1, t)=0, & \forall t \in(0, T) .\end{cases}
$$

In addition, let us consider $\left(u, u_{\Gamma}\right)$ be the solution of

$$
\begin{cases}\partial_{t} u_{L}(x, t)-\partial_{x}^{2} u_{L}(x, t)=\chi_{\omega} v(x, t), & \forall(x, t) \in \Omega_{L} \times(0, T),  \tag{3.1.10}\\ u_{R}(x, t)=u_{L}(0, t)=u_{\Gamma}(t), & \forall(x, t) \in \Omega_{R} \times(0, T), \\ \left(u(x, 0), u_{\Gamma}(0)\right)=\left(\tilde{u}_{0}(x), \tilde{u}_{0, \Gamma}\right), & \forall x \in \Omega, \\ u_{L}\left(-L_{1}, t\right)=0, & \forall t \in(0, T), \\ u_{\Gamma}^{\prime}(t)+\partial_{x} u_{L}(0, t)=0, & \forall t \in(0, T) .\end{cases}
$$

Suppose that $v \in L^{2}(\omega \times(0, T))$ is a control which drives the initial state $u_{0} \in L^{2}(\Omega)$ in (3.1.9) to zero at time $T>0$. Then, we can formulate the following question: can we employ the limit control $v$ of the problem (3.1.10) to control (in a suitable sense) system (3.1.9) too?

Concerning the above question, we have the following result
Theorem 3.3 Let $\varepsilon>0$ and choose $u_{0} \in L^{2}(\Omega)$ such that

$$
\begin{equation*}
\left\|u_{0, R}-\overline{u_{0, R}}\right\|_{L^{2}(\Omega)}^{2} \leq \frac{\varepsilon^{2}}{3}, \quad \text { where } \quad \overline{u_{0, R}}=\int_{\Omega_{R}} u_{0, R} d x . \tag{3.1.11}
\end{equation*}
$$

Let $\tilde{u}$ be the solution of (3.1.10) with control $v \in L^{2}(\omega \times(0, T))$ and initial condition

$$
\tilde{u}_{0, L}(x)=u_{0, L}(x), \quad \forall x \in \Omega_{L}, \quad \tilde{u}_{\Gamma, 0}=\tilde{u}_{0, R}(x)=\overline{u_{0, R}}, \quad \forall x \in \Omega_{R} .
$$

Then, there exists $K_{0}>0$ such that for all $K \geq K_{0}$ the associated solution $u^{K}$ of (3.1.9) fulfills

$$
\left\|u^{K}(\cdot, T)\right\|_{L^{2}(\Omega)} \leq \varepsilon .
$$

Roughly speaking, Theorem 3.3 asserts that for suitable initial conditions under consideration (in particular where $u_{0, R}$ is constant), problem (3.1.9) is approximately controllable at time $T>0$.

To end this section, we give the outline of this paper. First, in Section 3.2 we fix the functional setting to study the problem (3.1.1). In particular, well-posedness in the sense of Hadamard is given. In Section 3.3 we prove the Theorem 3.1. This will be done by using a suitable Carleman estimate to the adjoint system. In Section 3.4 we prove Theorem 3.2 by using well-known arguments combined the definition of weak solution of such problems. In Section 3.5 we prove Theorem 3.3 by using the ideas developed in the proof of Theorem 3.2.

### 3.2 Well-posedness of the heat equation with dynamic boundary conditions

In this section, we study well-posedness results for parabolic equations with dynamic boundary conditions like (3.1.1). In particular, we restrict our attention to definitions and properties for this kind of problem.

Before going further, we shall point out that the used in this section are well-known in the literature, see for instance [25], [41] and [80] and for parabolic problems with dynamic boundary conditions we refer to [44], [31], [55] and [81]. However, for completeness we give these notions and results (some of them without proof).

### 3.2.1 Variational approach

In this section, well-posedness result in the appropriate functional spaces of problems in the form (3.1.1) is considered. First, for $T>0$, let us consider ( $u, u_{\Gamma}$ ) being a solution of the following problem:

$$
\begin{cases}\partial_{t} u(x, t)-\partial_{x}^{2} u(x, t)=f(x, t), & \forall(x, t) \in \Omega_{L} \times(0, T),  \tag{3.2.1}\\ \left(u(x, 0), u_{\Gamma}(0)\right)=\left(u_{0}(x), u_{\Gamma, 0}\right), & \forall x \in \Omega_{L}, \\ u_{\Gamma}(t)=u(0, t), & \forall t \in(0, T), \\ u\left(-L_{1}, t\right)=0, & \forall t \in(0, T), \\ \partial_{t} u_{\Gamma}(t)+\partial_{x} u(0, t)=g(t), & \forall t \in(0, T),\end{cases}
$$

with $f \in L^{2}\left(\Omega_{L} \times(0, T)\right),\left(u_{0}, u_{\Gamma, 0}\right) \in L^{2}\left(\Omega_{L}\right) \times \mathbb{R}$ and $g \in L^{2}(0, T)$. To state the result, we introduce the following space:

$$
H_{L}^{1}\left(\Omega_{L}\right)=\left\{v \in H^{1}\left(\Omega_{L}\right) ; v\left(-L_{1}\right)=0\right\},
$$

endowed by the usual norm of $H^{1}\left(\Omega_{L}\right)$. It is clear that $H_{L}^{1}\left(\Omega_{L}\right)$ is a Hilbert space. Now, we introduce the bilinear form $a: H_{L}^{1}\left(\Omega_{L}\right) \times H_{L}^{1}\left(\Omega_{L}\right) \rightarrow \mathbb{R}$ given by

$$
a(u, v)=\int_{\Omega_{L}} \partial_{x} u \partial_{x} v d x, \quad \forall u, v \in V .
$$

Additionally, let $H$ be the completion of $H_{L}^{1}\left(\Omega_{L}\right)$ with respect to the norm induced by the inner product

$$
\begin{equation*}
(u, v)_{H}=\int_{\Omega_{L}} u v d x+u(0) v(0) . \tag{3.2.2}
\end{equation*}
$$

In this sense, it is clear that $H$ is isomorphic to $L^{2}\left(\Omega_{L}, d x\right) \times \mathbb{R}$, where $d x$ denotes the Lebesgue measure in $\Omega_{L}$. In the same manner, $H_{L}^{1}\left(\Omega_{L}\right)$ is isomorphic to the space

$$
V=\left\{\left(v, v_{\Gamma}\right) \in H^{1}(\Omega) \times \mathbb{R},\left.v\right|_{\Gamma}=v_{\Gamma} \text { and } v\left(-L_{1}\right)=0\right\},
$$

endowed by the norm $\left\|\left(v, v_{\Gamma}\right)\right\|_{V}=\left\|\partial_{x} v\right\|_{L^{2}\left(\Omega_{L}\right)}$. For this reason, from now one we shall write

$$
\left(\left(u, u_{\Gamma}\right),\left(v, v_{\Gamma}\right)\right)_{H}=\int_{\Omega_{L}} u v d x+u_{\Gamma}(0) v_{\Gamma}(0)
$$

for each $\left(u, u_{\Gamma}\right),\left(v, v_{\Gamma}\right) \in H$ to represent the inner product (3.2.2). Similarly, the bilinear form $a: V \times V \rightarrow \mathbb{R}$ can be defined by

$$
a\left(\left(u, u_{\Gamma}\right),\left(v, v_{\Gamma}\right)\right)=\int_{\Omega_{L}} \partial_{x} u \partial_{x} v d x, \quad \forall\left(u, u_{\Gamma}\right),\left(v, v_{\Gamma}\right) \in V .
$$

Now, we are interested in the following:
Problem: find $\left(u, u_{\Gamma}\right) \in C^{0}([0, T] ; H) \cap L^{2}(0, T ; V)$ such that for all $\left(v, v_{\Gamma}\right) \in V$ the following identity holds:

$$
\begin{equation*}
\left(\left(\partial_{t} u(t), u_{\Gamma}^{\prime}(t)\right),\left(v, v_{\Gamma}\right)\right)_{H}+a\left(\left(u, u_{\Gamma}\right),\left(v, v_{\Gamma}\right)\right)=\left((f, g),\left(v, v_{\Gamma}\right)\right)_{H}, \tag{3.2.3}
\end{equation*}
$$

in the sense of distributions on $(0, T)$ with $\left(u(x, 0), u_{\Gamma}(0)\right)=\left(u_{0}(x), u_{\Gamma, 0}\right)$ for each $x \in \Omega_{L}$.
It is not difficult to see that after integration by parts and well-known arguments (see [80], [85]), the above problem is equivalent to find a solution to (3.2.1). Then, concerning the above variational problem, we have the following result:

Proposition 3.4 For each $\left(u_{0}, u_{\Gamma, 0}\right) \in H$ and $(f, g) \in L^{2}(0, T ; H)$, the problem (3.2.3) admits a unique solution

$$
\left(u, u_{\Gamma}\right) \in C^{0}([0, T] ; H) \cap L^{2}(0, T ; V) .
$$

Moreover, the following energy estimate holds true

$$
\begin{equation*}
\left\|\left(u, u_{\Gamma}\right)\right\|_{C^{0}([0, T] ; H)}^{2}+\left\|\left(u, u_{\Gamma}\right)\right\|_{L^{2}(0, T ; V)}^{2} \leq C\left(\left\|\left(u_{0}, u_{\Gamma, 0}\right)\right\|_{H}^{2}+\|(f, g)\|_{L^{2}(0, T ; H)}^{2}\right) \tag{3.2.4}
\end{equation*}
$$

for some positive constant $C=C\left(\Omega_{L}, T\right)$.
It is clear that $a: V \times V \rightarrow \mathbb{R}$ is coercive and continuous on $V \times V$. Then, by standard arguments concerning parabolic problems (see e.g. 80) the existence of ( $u, u_{\Gamma}$ ) is guaranteed.

### 3.2.2 Semigroup approach

Using the notation previously introduced in the above section, let $A: D(A) \subset H \rightarrow H$ be the linear operator defined by

$$
\begin{equation*}
A\left(v, v_{\Gamma}\right)=\left(\partial_{x}^{2} v,-\partial_{x} v(0)\right), \tag{3.2.5}
\end{equation*}
$$

with domain

$$
D(A)=\left\{\left(v, v_{\Gamma}\right) \in V ; \partial_{x}^{2} v \in L^{2}\left(\Omega_{L}\right)\right\} .
$$

We have the following result
Proposition 3.5 The operator A given by (3.2.5) is densely defined, self-adjoint and generates a contraction semigroup $\left(e^{t A}\right)_{t \geq 0}$ on $H$.

Proof. It is easy to check that $\left\{\left(y, y_{\Gamma}\right) \in V ; y \in C^{\infty}\left(\overline{\Omega_{L}}\right)\right\} \subset D(A)$ is dense in $L^{2}\left(\Omega_{L}\right) \times \mathbb{R}$. Hence, $A$ is densely defined. In addition, for each $\left(v, v_{\Gamma}\right) \in D(A)$, we have

$$
\left(A\left(v, v_{\Gamma}\right),\left(v, v_{\Gamma}\right)\right)_{H}=\int_{\Omega_{L}} v \partial_{x}^{2} v d x-v_{\Gamma}(0) \partial_{x} v(0),
$$

and integration by parts shows that

$$
\begin{equation*}
\left(A\left(v, v_{\Gamma}\right),\left(v, v_{\Gamma}\right)\right)_{H}=-\int_{\Omega_{L}}\left|\partial_{x} v\right|^{2} d x \leq 0 . \tag{3.2.6}
\end{equation*}
$$

In the same manner, after integration by parts twice, we can assert that for each $\left(v, v_{\Gamma}\right),\left(w, w_{\Gamma}\right) \in D(A)$

$$
\left(A\left(v, v_{\Gamma}\right),\left(w, w_{\Gamma}\right)\right)_{H}=\left(\left(v, v_{\Gamma}\right), A\left(w, w_{\Gamma}\right)\right)_{H} .
$$

Thus, by Hille-Yosida's Theorem (see, for example, [30]), we conclude that $A$ is the generator of a contraction semigroup on $H$.

Next we introduce different classes of solutions of (3.2.1).
Definition 3.6 Let $f \in L^{2}\left(\Omega_{L} \times(0, T)\right), g \in L^{2}(0, T)$ and $Y_{0}=\left(y_{0}, y_{0, \Gamma}\right) \in H$.
(a) A strong solution of (3.2.1) is a function $U=\left(u, u_{\Gamma}\right) \in H^{1}(0, T ; H) \cap L^{2}(0, T ; D(A))$ fulfilling (3.2.1) in $L^{2}(0, T ; H)$.
(b) A mild solution of (3.2.1) is a function $U=\left(u, u_{\Gamma}\right) \in C^{0}([0, T] ; H)$ satisfying

$$
\begin{equation*}
U(t)=e^{t A} U_{0}+\int_{0}^{t} e^{(t-\tau) A}(f(\tau), g(\tau)) d \tau \tag{3.2.7}
\end{equation*}
$$

(c) A distributional solution of (3.2.1) is a function $U=\left(u, u_{\Gamma}\right) \in C^{0}([0, T] ; H)$ such that for all $\tau \in[0, T]$ we have

$$
\begin{aligned}
& \int_{0}^{\tau} \int_{\Omega_{L}} u\left(\partial_{t} \phi+\partial_{x}^{2} \phi\right) d x d t+\int_{\Omega_{L}}\left(u(x, \tau) \phi(x, \tau)-u_{0}(x) \phi(x, 0)\right) d x \\
& +u_{\Gamma}(\tau) \phi(0, \tau)-u_{\Gamma, 0} \phi(0,0)+\int_{0}^{\tau} u_{\Gamma}(t)\left(\partial_{x} \phi(0, t)-\partial_{t} \phi(0, t)\right) d t \\
= & \int_{0}^{\tau} \int_{\Omega_{L}} f \phi d x d t+\int_{0}^{\tau} g(t) \phi(0, t) d t,
\end{aligned}
$$

for all $\phi \in C^{\infty}\left(\overline{\Omega_{L}} \times[0, T]\right)$ such that $\phi\left(-L_{1}, t\right)=0$ for all $t \in(0, T)$.
The next result asserts the existence of strong solutions for regular initial data:
Proposition 3.7 Let $f \in L^{2}\left(\Omega_{L} \times(0, T)\right), g \in L^{2}(0, T)$ and $\left(u_{0}, u_{0, \Gamma}\right) \in V$. Then, there exists a unique strong solution

$$
U=\left(u, u_{\Gamma}\right) \in E_{1}(0, t):=H^{1}(0, T ; H) \cap L^{2}(0, T ; D(A))
$$

of (3.2.1), which is also a mild solution. In addition, there exists a constant $C>0$ such that

$$
\left\|\left(u, u_{\Gamma}\right)\right\|_{E_{1}(0, T)} \leq C\left(\left\|\left(u, u_{0, \Gamma}\right)\right\|_{V}+\|f\|_{L^{2}\left(\Omega_{L}\right) \times(0, T)}+\|g\|_{L^{2}(0, T)}\right) .
$$

However, for our purposes (specially for controllability results) we shall consider initial data $\left(u_{0}, u_{0, \Gamma}\right) \in H$. Then, the next result gives necessary conditions to get the uniqueness of a mild solution and describes the regularity of such solutions.

Proposition 3.8 Let $f \in L^{2}\left(\Omega_{L} \times(0, T)\right), g \in L^{2}(0, T)$ and $\left(u, u_{0, \Gamma}\right) \in H$. Then,

1. there exists a unique mild solution $U \in C([0, T] ; H)$ of (3.2.1) and the following energy estimate holds:

$$
\left\|\left(u, u_{\Gamma}\right)\right\|_{C^{0}([0, T] ; H)} \leq C\left(\left\|\left(u_{0}, u_{0, \Gamma}\right)\right\|_{H}+\|f\|_{L^{2}\left(\Omega_{L} \times(0, T)\right)}+\|g\|_{L^{2}(0, T)}\right),
$$

for some positive constant $C=C\left(\Omega_{L}, T\right)$. Moreover, for each $\tau \in(0, T)$ we get

$$
\left(u, u_{\Gamma}\right) \in H^{1}(\tau, T ; H) \cap L^{2}(\tau, T ; D(A)) .
$$

2. If $\left(u, u_{0, \Gamma}\right) \in V$, then there mild solution of (3.2.1) given by the first item is a strong one.
3. A function $\left(u, u_{\Gamma}\right)$ is a distributional solution of (3.2.1) if and only if it is a mild solution.

### 3.3 Controllability properties of the original problem

In this section, we devote to prove Theorem 3.1. First, using the well-known relation between null controllability and observability, we introduce the adjoint system

$$
\begin{cases}\partial_{t} z(x, t)+\partial_{x}^{2} z(x, t)=0, & \forall(x, t) \in \Omega_{L} \times(0, T),  \tag{3.3.1}\\ \left(z(x, T), z_{\Gamma}(T)\right)=\left(z_{T}(x), z_{\Gamma, T}\right), & \forall x \in \Omega_{L}, \\ z\left(-L_{1}, t\right)=0, & \forall t \in(0, T), \\ z_{\Gamma}^{\prime}(t)-\partial_{x} z(0, t)=0, & \forall t \in(0, T)\end{cases}
$$

Using the change of variables $t^{\prime}=T-t$ and applying Proposition 3.4 we deduce that system (3.3.1) has a unique solution $z \in C^{0}([0, T] ; H) \cap L^{2}(0, T ; V)$, where $H$ and $V$ are the spaces defined in the previous section. Moreover, we have the following energy estimate:

$$
\begin{equation*}
\left\|\left(z, z_{\Gamma}\right)\right\|_{C^{0}([0, T] ; H)}+\left\|\left(z, z_{\Gamma}\right)\right\|_{L^{2}(0, T ; V)} \leq C\left\|\left(z_{T}, z_{\Gamma, T}\right)\right\|_{H}, \tag{3.3.2}
\end{equation*}
$$

for some constant $C=C\left(\Omega_{L}, T\right)$ Moreover, by Proposition 3.8 for all $\tau \in(0, T)$, we have

$$
\left(z, z_{\Gamma}\right) \in H^{1}(0, \tau ; H) \cap L^{2}(0, \tau ; D(A))
$$

As we said before, the proof of (3.3.2) (or equivalently Theorem 3.1) is based on the observability inequality for the adjoint system (3.3.1):

$$
\begin{equation*}
\|z(\cdot, 0)\|_{L^{2}\left(\Omega_{L}\right)}^{2}+\left|z_{\Gamma}(0)\right|^{2} \leq C \int_{0}^{T} \int_{\omega}|z|^{2} d x d t \tag{3.3.3}
\end{equation*}
$$

for all $\left(z_{T}, z_{\Gamma, T}\right) \in L^{2}\left(\Omega_{L}\right) \times \mathbb{R}$, where $z$ is the associated solution to (3.3.1), and for some positive constant $C=C\left(\Omega_{L}, T\right)$. This will be done by using a suitable Carleman estimate for (3.3.1).

In order to formulate next result, we shall introduce weight functions. For $\alpha \geq 1$, we define

$$
\begin{aligned}
\theta(t) & =(t(T-t))^{-\alpha}, \quad \forall t \in(0, T), \\
\psi(x) & =-\frac{1}{4 L_{1}} x^{2}+x+2 L_{1}, \quad \forall x \in \overline{\Omega_{L}}
\end{aligned}
$$

with $\varphi(x, t)=\theta(t) \psi(x)$, for each $(x, t) \in \Omega_{L} \times(0, T)$. Notice that $\varphi$ is a smooth positive function which blows up as $t \rightarrow 0^{+}$and as $t \rightarrow T^{-}$.

Now we state the one-parameter Carleman estimate:
Lemma 3.9 Let $T>0$ and for $\alpha \geq 1$ define the function $\varphi$ as above. Then, there exist positive constants $C=C\left(\alpha, \Omega_{L}, T\right)$ and $s_{0}=s_{0}\left(\alpha, \Omega_{L}, T\right)$ such that for all $s \geq s_{0}$ the

$$
\begin{align*}
& s^{3} \int_{0}^{T} \int_{\Omega_{L}} e^{-2 s \varphi}(t(T-t))^{-3 \alpha}|y|^{2} d x d t+s \int_{0}^{T} \int_{\Omega_{L}} e^{-2 s \varphi}(t(T-t))^{-\alpha}\left|\partial_{x} y\right|^{2} d x d t \\
& +s^{-1} \int_{0}^{T} \int_{\Omega_{L}} e^{-2 s \varphi}(t(T-t))^{\alpha}\left(\left|\partial_{x}^{2} y\right|^{2}+\left|\partial_{t} y\right|^{2}\right) d x d t \\
& +s^{3} \int_{0}^{T} e^{-2 s \varphi(0, t)}(t(T-t))^{-3 \alpha}\left|y_{\Gamma}(t)\right|^{2} d t  \tag{3.3.4}\\
& \left.+s \int_{0}^{T} e^{-2 s \varphi(0, t)} t(T-t)\right)^{-\alpha}\left|\partial_{x} y(0, t)\right|^{2} d t+\int_{0}^{T} e^{-2 s \varphi(0, t)}\left|y_{\Gamma}^{\prime}(t)\right|^{2} d t \\
& \leq C \int_{0}^{T} \int_{\Omega_{L}} e^{-2 s \varphi}\left|\partial_{t} y+\partial_{x}^{2} y\right|^{2} d x d t+C \int_{0}^{T} e^{-2 s \varphi(0, t)}\left|y_{\Gamma}^{\prime}(t)-\partial_{x} y(0, t)\right|^{2} d t \\
& +C s \int_{0}^{T} e^{-2 s \varphi\left(-L_{1}, t\right)}(t(T-t))^{-\alpha}\left|\partial_{x} y\left(-L_{1}, t\right)\right|^{2} d t
\end{align*}
$$

for all functions $\left(y, y_{\Gamma}\right) \in H^{1}(0, T ; H) \cap L^{2}(0, T ; D(A))$.
The proof is based on the works of A. Fursikov and O. Imanuvilov [50, 62], 61] in the case of Dirichlet and mixed boundary conditions. In our setting, we will see that some new boundary terms arise from the Wentzell dynamic condition. Then, the main difficulty is to prove that these new terms can be absorbed by choosing the parameter $s$ large enough in the spirit of Carleman estimates.

Remark 3.10 By using a cut-off function $\eta$ localized close to $x=-L_{1}$, and standard observability properties for parabolic equations (see for example [48]), we can prove easily (3.3.3) when $\omega=\left(-L_{1},-a\right)$ with $a>0$.

Remark 3.11 One can consider Classical weight functions introduced by A. Fursikov and O. Imanuvilov [50], [62] in the unidimensional case. Indeed, these details are given in the appendix. However, in the context of multidimensional setting this strategy fails, see Remark 3.3 of [81]. In fact, there are some boundary terms depending on $\nabla_{T} y$ which we cannot absorb it.

Proof of Lemma 3.9. Since all the terms in (3.3.4) are continuous with respect to the norm of $E_{1}(0, T)$, it suffices to consider smooth functions $y \in C^{\infty}\left(\overline{\Omega_{L}} \times[0, T]\right)$. In fact, the general case follows by the classical approximation by convolution with mollifiers in space and time by the density of $C^{\infty}\left([0, T] \times \overline{\Omega_{L}}\right)$ in $E_{1}(0, T)$. This allow us to deduce that $y_{\Gamma}^{\prime}(t)=y(0, t)$ for all $t \in(0, T)$.

In what follows, $C$ denotes a generic constant depending on $\alpha, \Omega_{L}, \omega$ and $T>0$ that may change from line to line. For an easier comprehension, we divide the proof into four steps:

- Step 1: Setting. Let us introduce the conjugate variable

$$
z(x, t)=e^{-s \varphi(x, t)} y(x, t), \quad \forall(x, t) \in \Omega_{L} \times(0, T)
$$

Then, direct computations show that the space and time derivatives of $z$ are given by

$$
\begin{array}{r}
\partial_{t} z=-s \partial_{t} \varphi z+e^{-s \varphi} \partial_{t} y, \quad \partial_{x} z=-s \partial_{x} \varphi z+e^{-s \varphi} \partial_{x} y, \\
\partial_{x}^{2} z=-s \partial_{x}^{2} \varphi z-s^{2}\left|\partial_{x} \varphi\right|^{2} z-2 s \partial_{x} \varphi \partial_{x} \varphi z+e^{-s \varphi} \partial_{x}^{2} y,
\end{array}
$$

for all $(x, t) \in \Omega_{L} \times(0, T)$. In order to simplify the computations, let us define the operators

$$
\begin{array}{cl}
M_{1}=s \partial_{t} \varphi+\partial_{x}^{2}+s^{2}\left|\partial_{x} \varphi\right|^{2}, & M_{2}=\partial_{t}+2 s \partial_{x} \varphi+s \partial_{x}^{2} \varphi \\
N_{1}=s \partial_{t} \varphi-\partial_{x}, & N_{2}=\partial_{t}-s \partial_{x} \varphi .
\end{array}
$$

In addition, due to the regularity of $\left(y, y_{\Gamma}\right)$ we set

$$
f(x, t)=\partial_{t} y(x, t)+\partial_{x}^{2} y(x, t) \quad \forall(x, t) \in \Omega_{L} \times(0, T), \quad g(t)=\partial_{t} y(0, t)-\partial_{x} y(0, t) \quad \forall t \in(0, T)
$$

Then, according to the above computations, it is clear that

$$
\begin{align*}
e^{-s \varphi(x, t)} f(x, t) & =M_{1} z(x, t)+M_{2} z(x, t), \quad \forall(x, t) \in \Omega_{L} \times(0, T) \quad \text { and }  \tag{3.3.5}\\
e^{-s \varphi(0, t)} g(t) & =N_{1}(z)(0, t)+N_{2}(z)(0, t), \quad \forall t \in(0, T) .
\end{align*}
$$

Taking $\|\cdot\|_{L^{2}\left(\Omega_{L} \times(0, T)\right)}$ and $\|\cdot\|_{L^{2}(0, T)}$ to the equations in 3.3.5) we have

$$
\begin{align*}
& \quad\left\|e^{-s \varphi} f\right\|_{L^{2}\left(\Omega_{L} \times(0, T)\right)}^{2}+\left\|e^{-s \varphi(0, t)} g\right\|_{L^{2}(0, T)}^{2} \\
& =\left\|M_{1} z\right\|_{L^{2}\left(\Omega_{L} \times(0, T)\right)}^{2}+\left\|M_{2} z\right\|_{L^{2}\left(\Omega_{L} \times(0, T)\right)}^{2}+\left\|N_{1} z(0, \cdot)\right\|_{L^{2}(0, T)}^{2}+\left\|N_{2} z(0, \cdot)\right\|_{L^{2}(0, T)}^{2}  \tag{3.3.6}\\
& \quad+2\left\langle M_{1} z, M_{2} z\right\rangle_{L^{2}\left(\Omega_{L} \times(0, T)\right)}+2\left\langle N_{1} z(0, \cdot), N_{2} z(0, \cdot)\right\rangle_{L^{2}(0, T)} .
\end{align*}
$$

Our next task is to compute the inner products in $L^{2}\left(\Omega_{L} \times(0, T)\right)$ and $L^{2}(0, T)$. This will be done using integration by parts and applying boundary conditions at $x=-L_{1}$ and $x=0$.

- Step 2. Now let us compute the terms of the first inner product. In order to do that, let us use the following notation:

$$
\left\langle M_{1} z, M_{2} z\right\rangle_{L^{2}\left(\Omega_{L} \times(0, T)\right)}=\sum_{i, j=1}^{3} I_{j, k}
$$

where $J_{j k}$ stands for the scalar product in $L^{2}\left(\Omega_{L} \times(0, T)\right)$ between the $j^{\text {th }}$-term of $M_{1} z$ and the $k^{\text {th }}$-term of $M_{2} z$. Then, we start with $J_{11}$. Using the fact that

$$
z \partial_{t} z=\frac{1}{2} \partial_{t}|z|^{2}, \text { in } \Omega_{L} \times(0, T)
$$

and integrating by parts in time leads

$$
\begin{equation*}
I_{11}=s \int_{0}^{T} \int_{\Omega_{L}} \partial_{t} \varphi z \partial_{t} z d x d t=-\frac{1}{2} s \int_{0}^{T} \int_{\Omega_{L}} \partial_{t}^{2} \varphi|z|^{2} d x d t \tag{3.3.7}
\end{equation*}
$$

where we have used the fact that $z(\cdot, t) \rightarrow 0$ as $t \rightarrow 0^{+}$and $t \rightarrow T^{-}$. In the same manner, $J_{12}$ can be estimated in the following way:

$$
\begin{align*}
I_{12} & =2 s^{2} \int_{0}^{T} \int_{\Omega_{L}} \partial_{t} \varphi \partial_{x} \varphi z \partial_{x} z d x d t \\
& =-s^{2} \int_{0}^{T} \int_{\Omega_{L}}\left(\partial_{t} \partial_{x} \varphi \partial_{x} \varphi+\partial_{t} \varphi \partial_{x}^{2} \varphi\right)|z|^{2} d x d t+s^{2} \int_{0}^{T} \partial_{t} \varphi(0, t) \partial_{x} \varphi(0, t)|z(0, t)|^{2} d t, \tag{3.3.8}
\end{align*}
$$

where we have used the homogeneous Dirichlet boundary condition on $x=-L_{1}$. Furthermore, by definition $J_{13}$ reads as follows

$$
\begin{equation*}
I_{13}=s^{2} \int_{0}^{T} \int_{\Omega_{L}} \partial_{t} \varphi \partial_{x}^{2} \varphi|z|^{2} d x d t \tag{3.3.9}
\end{equation*}
$$

Let us compute $J_{21}$. Integration by parts in space yields

$$
\begin{aligned}
I_{21} & =\int_{0}^{T} \int_{\Omega_{L}} \partial_{t} z \partial_{x}^{2} z d x d t \\
& =-\int_{0}^{T} \int_{\Omega_{L}} \partial_{t} \partial_{x} z \partial_{x} z d x d t+\int_{0}^{T} \partial_{t} z(0, t) \partial_{x} z(0, t) d t
\end{aligned}
$$

Let us compute the above terms. First, notice that

$$
\partial_{x} z(\cdot, t) \rightarrow 0 \text { as } t \rightarrow 0^{+} \text {and } t \rightarrow T^{-} .
$$

Then, using the fact that $\partial_{t} \partial_{x} \partial_{x} z=\frac{1}{2} \partial_{t}\left|\partial_{x} z\right|^{2}$ in $\Omega_{L} \times(0, T)$ it follows that

$$
\int_{0}^{T} \int_{\Omega_{L}} \partial_{t}\left(\left|\partial_{x} z\right|^{2}\right) d x d t=0
$$

On the other hand, according to the definitions of $N_{1}$ and $N_{2}$, we have

$$
\begin{aligned}
& \int_{0}^{T} \partial_{t} z(0, t) \partial_{x} z(0, t) d t \\
= & \int_{0}^{T}\left|\partial_{t} z(0, t)\right|^{2} d t-s \int_{0}^{T} \partial_{x} \varphi(0, t) z(0, t) \partial_{t} z(0, t) d t+s \int_{0}^{T} \partial_{t} \varphi(0, t) z(0, t) \partial_{t} z(0, t) d t \\
& -\int_{0}^{T} e^{-s \varphi(0, t)} g(t) \partial_{t} z(0, t) d t .
\end{aligned}
$$

Integration by parts yields

$$
-s \int_{0}^{T} \partial_{x} \varphi(0, t) z(0, t) \partial_{t} z(0, t) d t=\frac{1}{2} s \int_{0}^{T} \partial_{t} \partial_{x} \varphi(0, t)|z(0, t)|^{2} d t
$$

and

$$
s \int_{0}^{T} \partial_{t} \varphi(0, t) z(0, t) \partial_{t} z(0, t) d t=-\frac{1}{2} \int_{0}^{T} \partial_{t}^{2} \varphi(0, t)|z(0, t)|^{2} d t .
$$

Therefore, $I_{21}$ is given by

$$
\begin{align*}
I_{21}= & \int_{0}^{T}\left|\partial_{t} z(0, t)\right|^{2} d t+\frac{1}{2} s \int_{0}^{T} \partial_{t} \partial_{x} \varphi(0, t)|z(0, t)|^{2} d t  \tag{3.3.10}\\
& -\frac{1}{2} s \int_{0}^{T} \partial_{t}^{2} \varphi(0, t)|z(0, t)|^{2} d t
\end{align*}
$$

Let us emphasize that $I_{21}$ contains the boundary term $\left|\partial_{t} z(0, t)\right|^{2}$, which plays an important role to eliminate the boundary terms from the next step. Once again, integration by parts in space we have

$$
\begin{align*}
I_{22}= & -s \int_{0}^{T} \int_{\Omega_{L}} \partial_{x}^{2} \varphi\left|\partial_{x} z\right|^{2} d x d t+s \int_{0}^{T} \partial_{x} \varphi(0, t)\left|\partial_{x} z(0, t)\right|^{2} d t  \tag{3.3.11}\\
& -s \int_{0}^{T} \partial_{x} \varphi\left(-L_{1}, t\right)\left|\partial_{x} z\left(-L_{1}, t\right)\right|^{2} d t .
\end{align*}
$$

We point out that the third term of the right-hand side of $I_{22}$ will be considered as an observation. In the same manner, $I_{23}$ is given by

$$
\begin{align*}
I_{23} & =s \int_{0}^{T} \int_{\Omega_{L}} \partial_{x}^{2} \varphi \partial_{x}^{2} z d x d t \\
& =-s \int_{0}^{T} \int_{\Omega_{L}} \partial_{x}^{2} \varphi\left|\partial_{x} z\right|^{2} d x d t+s \int_{0}^{T} \partial_{x}^{2} \varphi(0, t) z(0, t) \partial_{x} z(0, t) d t, \tag{3.3.12}
\end{align*}
$$

where we have used $z\left(-L_{1}, t\right)=0$ for all $t \in(0, T)$ and the fact that $\partial_{x}^{3} \varphi=0$ in $\Omega_{L} \times(0, T)$. In addition, the term $I_{31}$ reads as follows

$$
\begin{equation*}
I_{31}=-\frac{1}{2} s^{2} \int_{0}^{T} \int_{\Omega_{L}} \partial_{t}\left(\left|\partial_{x} \varphi\right|^{2}|z|^{2}\right) d x d t . \tag{3.3.13}
\end{equation*}
$$

Moreover, integration by parts yields

$$
\begin{align*}
I_{32} & =2 s^{3} \int_{0}^{T} \int_{\Omega_{L}}\left|\partial_{x} \varphi\right|^{3} z \partial_{x} z d x d t  \tag{3.3.14}\\
& =-3 s^{2} \int_{0}^{T} \int_{\Omega_{L}}\left|\partial_{x} \varphi\right|^{2} \partial_{x}^{2} \varphi|z|^{2} d x d t+s^{3} \int_{0}^{T}\left|\partial_{x} \varphi(0, t)\right|^{3}|z(0, t)|^{2} d t . \tag{3.3.15}
\end{align*}
$$

Finally, by definition $I_{33}$ is given by

$$
\begin{equation*}
I_{33}=s^{3} \int_{0}^{T} \int_{\Omega_{L}}\left|\partial_{x} \varphi\right|^{2} \partial_{x}^{2} \varphi|z|^{2} d x d t \tag{3.3.16}
\end{equation*}
$$

In the next step, we compute the the second inner product in the equation (3.3.6). We emphasize that this step plays an important role on the proof of Lemma 3.9. In fact, we must ensure that all these new terms can be controlled or absorbed by taking the parameter $s$ large enough.

- Step 3: Boundary terms. We introduce the notation

$$
\left\langle N_{1}(z)(0, t), N_{2}(z)(0, t)\right\rangle_{L^{2}(0, T)}=\sum_{j, k=1}^{2} J_{j, k},
$$

where $J_{j k}$ stands for the scalar product in $L^{2}(0, T)$ between the $j^{\text {th }}$-term of $N_{1}(z)(0, \cdot)$ and the $k^{\text {th }}$-term of $N_{2}(z)(0, \cdot)$. Then, $J_{1} 1$ can be estimated in the following way:

$$
\begin{equation*}
J_{11}=s \int_{0}^{T} \partial_{t} \varphi(0, t) z(0, t) \partial_{t} z(0, t) d t=-\frac{1}{2} s \int_{0}^{T} \partial_{t}^{2} \varphi(0, t)|z(0, t)|^{2} d t \tag{3.3.17}
\end{equation*}
$$

Moreover, the other terms are given by

$$
\begin{align*}
& J_{12}=-s^{2} \int_{0}^{T} \partial_{t} \varphi(0, t) \partial_{x} \varphi(0, t)|z(0, t)|^{2} d t  \tag{3.3.18}\\
& J_{21}=-\int_{0}^{T} \partial_{x} z(0, t) \partial_{t} z(0, t) d t  \tag{3.3.19}\\
& J_{22}=-s \int_{0}^{T} \partial_{x} \varphi(0, t) z(0, t) \partial_{x} z(0, t) d t \tag{3.3.20}
\end{align*}
$$

As we shall see in the next step, due to the definitions of Carleman weights $\psi$ and $\theta$ with $\alpha \geq 1$, the terms given by the equations (3.3.17)-(3.3.20) can be bounded in a suitable way by taking $s$ large enough. This enable us to control these new boundary terms.

- Step 4: Substituting (3.3.7)-(3.3.20) into (3.3.6) and gathering the terms we have

$$
\begin{align*}
& \quad-2 s^{3} \int_{0}^{T} \int_{\Omega_{L}}\left|\partial_{x} \varphi\right|^{2} \partial_{x}^{2} \varphi|z|^{2} d x d t-2 s \int_{0}^{T} \int_{\Omega_{L}} \partial_{x}^{2} \varphi\left|\partial_{x} z\right|^{2} d x d t \\
& +s^{3} \int_{0}^{T}\left|\partial_{x} \varphi(0, t)\right|^{3}|z(0, t)|^{2} d t+s \int_{0}^{T} \partial_{x} \varphi(0, t)\left|\partial_{x} z(0, t)\right|^{2} d t \\
& +\int_{0}^{T}\left|\partial_{t} z(0, t)\right|^{2} d t+\left\|M_{1} z\right\|_{L^{2}\left(\Omega_{L \times(0, T))}^{2}+\left\|M_{2} z\right\|_{L^{2}\left(\Omega_{L \times(0, T))}^{2}\right.}^{2}\right.}+\left\|N_{1}(z)(0, t)\right\|_{L^{2}(0, T)}^{2}+\left\|N_{2}(z)(0, t)\right\|_{L^{2}(0, T)}^{2}  \tag{3.3.21}\\
& =\left\|e^{-s \varphi} f\right\|_{L^{2}\left(\Omega_{L} \times(0, T)\right)}^{2}+\left\|e^{-s \varphi(0, t)} g\right\|_{L^{2}(0, T)}^{2}+s \int_{0}^{T} \partial_{x} \varphi\left(-L_{1}, t\right)\left|\partial_{x} z\left(-L_{1}, t\right)\right|^{2} d t \\
& \quad+X+Y .
\end{align*}
$$

where $X$ and $Y$ are given by

$$
X=\frac{1}{2} s \int_{0}^{T} \int_{\Omega_{L}} \partial_{t}^{2} \varphi|z|^{2} d x d t+s^{2} \int_{0}^{T} \int_{\Omega_{L}} \partial_{t} \partial_{x} \varphi \partial_{x} \varphi|z|^{2} d x d t
$$

and

$$
\begin{aligned}
Y= & -\frac{1}{2} s \int_{0}^{T} \partial_{t} \partial_{x} \varphi(0, t)|z(0, t)|^{2} d t-\frac{1}{2} s \int_{0}^{T} \partial_{t}^{2} \varphi(0, t)|z(0, t)|^{2} d t \\
& +\int_{0}^{T} e^{-s \varphi(0, t)} g \partial_{t} z(0, t) d t-s \int_{0}^{T} \partial_{x}^{2} \varphi(0, t) z(0, t) \partial_{x} z(0, t) d t \\
& -\frac{1}{2} s \int_{0}^{T} \partial_{t}^{2} \varphi(0, t)|z(0, t)|^{2} d t+s^{2} \int_{0}^{T} \partial_{t} \varphi(0, t) \partial_{x} \varphi(0, t)|z(0, t)|^{2} d t \\
& +\int_{0}^{T} \partial_{x} z(0, t) \partial_{t} z(0, t) d t+s \int_{0}^{T} \partial_{x} \varphi(0, t) z(0, t) \partial_{x} z(0, t) d t .
\end{aligned}
$$

Our next task is to eliminate the terms $X$ and $Y$. In order to do that, let us point out that the derivatives of $\theta$ can be bounded as follows:

$$
\begin{equation*}
\left|\theta^{\prime}(t)\right| \leq \alpha T(t(T-t))^{-(\alpha+1)}, \quad\left|\theta^{\prime \prime}(t)\right| \leq C(t(T-t))^{-(\alpha+2)}, \quad \forall t \in(0, T) \tag{3.3.22}
\end{equation*}
$$

for some positive constant $C=C\left(\alpha, \Omega_{L}, T\right)$. On the other hand, $\psi$ and their derivatives satisfy

$$
\begin{equation*}
\frac{3}{2} L_{1} \leq \psi(x) \leq 2 L_{1}, \quad \frac{1}{2} \leq \psi^{\prime}(x) \leq 1, \quad \psi^{\prime \prime}(x)=-\frac{1}{2 L_{1}}, \quad \forall x \in \overline{\Omega_{L}} \tag{3.3.23}
\end{equation*}
$$

Then, by using inequalities (3.3.22) and (3.3.23) in (3.3.21) we obtain

$$
\begin{align*}
& s^{3} \int_{0}^{T} \int_{\Omega_{L}}(t(T-t))^{-3 \alpha}|z|^{2} d x d t+s \int_{0}^{T} \int_{\Omega_{L}}(t(T-t))^{-\alpha}\left|\partial_{x} z\right|^{2} d x d t \\
& +s^{3} \int_{0}^{T}(t(T-t))^{-3 \alpha}|z(0, t)|^{2} d t+s \int_{0}^{T}(t(T-t))^{-\alpha}\left|\partial_{x} z(0, t)\right|^{2} d t \\
& +\int_{0}^{T}\left|\partial_{t} z(0, t)\right|^{2} d t+\left\|M_{1} z\right\|_{L^{2}\left(\Omega_{L \times(0, T))}\right.}^{2}+\left\|M_{2} z\right\|_{L^{2}\left(\Omega_{L} \times(0, T)\right)}^{2}  \tag{3.3.24}\\
& \left\|N_{1}(z)(0, t)\right\|_{L^{2}(0, T)}^{2}+\left\|N_{2}(z)(0, t)\right\|_{L^{2}(0, T)}^{2} \leq C_{1}\left\|e^{-s \varphi} f\right\|_{L^{2}\left(\Omega_{L} \times(0, T)\right)} \\
& C_{1}\left\|e^{-s \varphi(0, t)} g\right\|_{L^{2}(0, T)}^{2}+C_{1} s \int_{0}^{T}(t(T-t))^{-\alpha}\left|\partial_{x} z\left(-L_{1}, t\right)\right|^{2} d t+C_{1}|X|+C_{1}|Y|,
\end{align*}
$$

for some constant $C_{1}=C_{1}\left(\alpha, \Omega_{L}, T\right)$. Notice that

$$
C_{1}|X| \leq C_{2} s \int_{0}^{T} \int_{\Omega_{L}}(t(T-t))^{-\alpha+2}|z|^{2} d x d t+C_{2} s^{2} \int_{0}^{T} \int_{\Omega_{L}}(t(T-t))^{-(2 \alpha+1)}|z|^{2} d x d t .
$$

Moreover, since $\alpha \geq 1$ we can choose $s_{1}>0$ large enough to get

$$
C_{1}|X| \leq \frac{1}{2} s^{3} \int_{0}^{T} \int_{\Omega_{L}}(t(T-t))^{-3 \alpha}|z|^{2} d x d t, \quad \forall s \geq s_{1} .
$$

By using Young's inequality and the previous arguments we deduce the existence of a constant $s_{2}=s_{2}\left(\alpha, C_{1}, \Omega, T\right) \geq s_{1}$ such that the following estimate for $Y_{1}$ holds:

$$
\begin{align*}
C_{1}|Y| \leq & \frac{1}{2} s^{3} \int_{0}^{T}(t(T-t))^{-3 \alpha}|z(0, t)|^{2}+\frac{1}{2} s \int_{0}^{T}(t(T-t))^{-\alpha}\left|\partial_{x} z(0, t)\right|^{2} d t \\
& +\frac{1}{2} \int_{0}^{T}\left|\partial_{t} z(0, t)\right|^{2} d t+\frac{1}{2} \int_{0}^{T} e^{-s \varphi(0, t)}|g|^{2} d t, \quad \forall s \geq s_{2} . \tag{3.3.25}
\end{align*}
$$

Therefore, for each $s \geq s_{3}$ with $s_{3}=\max \left\{s_{1}, s_{2}\right\}$ we get

$$
\begin{align*}
& s^{3} \int_{0}^{T} \int_{\Omega_{L}}(t(T-t))^{-3 \alpha}|z|^{2} d x d t+s \int_{0}^{T} \int_{\Omega_{L}}(t(T-t))^{-\alpha}\left|\partial_{x} z\right|^{2} d x d t \\
& \quad+s^{3} \int_{0}^{T}(t(T-t))^{-3 \alpha}|z(0, t)|^{2} d t+s \int_{0}^{T}(t(T-t))^{-\alpha}\left|\partial_{x} z(0, t)\right|^{2} d t \\
& \quad+\int_{0}^{T}\left|\partial_{t} z(0, t)\right|^{2} d t+\left\|M_{1} z\right\|_{L^{2}\left(\Omega_{L \times(0, T))}^{2}+\left\|M_{2} z\right\|_{L^{2}\left(\Omega_{L} \times(0, T)\right)}^{2}\right.}^{\quad+\left\|N_{1}(z)(0, t)\right\|_{L^{2}(0, T)}^{2}+\left\|N_{2}(z)(0, t)\right\|_{L^{2}(0, T)}^{2}} \\
& \leq C_{1}\left\|e^{-s \varphi} f\right\|_{L^{2}\left(\Omega_{L} \times(0, T)\right)}^{2}+C_{1}\left\|e^{-s \varphi(0, t)} g\right\|_{L^{2}(0, T)}^{2}+C_{1} s \int_{0}^{T}(t(T-t))^{-\alpha}\left|\partial_{x} z\left(-L_{1}, t\right)\right|^{2} d t .
\end{align*}
$$

It remains to deduce estimates for $\partial_{t} z$ and $\partial_{x}^{2} z$. To do this, by definition of $M_{1}$ and $M_{2}$ we can assert that

$$
\begin{align*}
& \quad s^{-1} \int_{0}^{T} \int_{\Omega_{L}}(t(T-t))^{\alpha}\left|\partial_{x}^{2} z\right|^{2} d x d t \\
& \leq  \tag{3.3.27}\\
& \leq s^{-1} \int_{0}^{T} \int_{\Omega_{L}}(t(T-t))^{\alpha}\left|M_{1} z\right|^{2} d x d t+C s \int_{0}^{T} \int_{\Omega_{L}}(t(T-t))^{-(\alpha+2)}|z|^{2} d x d t \\
& \quad+C s^{3} \int_{0}^{T} \int_{\Omega_{L}}(t(T-t))^{-3 \alpha}|z|^{2} d x d t, \quad \forall s>0,
\end{align*}
$$

and

$$
\begin{align*}
& s^{-1} \int_{0}^{T} \int_{\Omega_{L}}(t(T-t))^{\alpha}\left|\partial_{t} z\right|^{2} d x d t \\
\leq & C s^{-1} \int_{0}^{T} \int_{\Omega_{L}}(t(T-t))^{\alpha}\left|M_{2} z\right|^{2} d x d t+C s \int_{0}^{T} \int_{\Omega_{L}}(t(T-t))^{-\alpha}|z|^{2} d x d t  \tag{3.3.28}\\
& +C s \int_{0}^{T} \int_{\Omega_{L}}(t(T-t))^{-\alpha}\left|\partial_{x} z\right|^{2} d x d t, \forall s>0 .
\end{align*}
$$

Thus, the global terms of $\partial_{t} z$ and $\partial_{x}^{2} z$ can be incorporated in the left-hand side of
(3.3.26). Thus, for each $s \geq s_{3}$ we have

$$
\begin{align*}
& s^{3} \int_{0}^{T} \int_{\Omega_{L}}(t(T-t))^{-3 \alpha}|z|^{2} d x d t+s \int_{0}^{T} \int_{\Omega_{L}}(t(T-t))^{-\alpha}\left|\partial_{x} z\right|^{2} d x d t \\
& \quad+s^{-1} \int_{0}^{T} \int_{\Omega_{L}}(t(T-t))^{\alpha}\left(\left|\partial_{x}^{2} z\right|^{2}+\left|\partial_{t} z\right|^{2}\right) d x d t \\
& \quad+s^{3} \int_{0}^{T}(t(T-t))^{-3 \alpha}|z(0, t)|^{2} d t+s \int_{0}^{T}(t(T-t))^{-\alpha}\left|\partial_{x} z(0, t)\right|^{2} d t \\
& \quad+\int_{0}^{T}\left|\partial_{t} z(0, t)\right|^{2} d t+\left\|M_{1} z\right\|_{L^{2}\left(\Omega_{L} \times(0, T)\right)}^{2}+\left\|M_{2} z\right\|_{L^{2}\left(\Omega_{L} \times(0, T)\right)}^{2} \\
& \quad+\left\|N_{1}(z)(0, t)\right\|_{L^{2}(0, T)}^{2}+\left\|N_{2}(z)(0, t)\right\|_{L^{2}(0, T)}^{2} \\
& \leq C_{2}\left\|e^{-s \varphi} f\right\|_{L^{2}\left(\Omega_{L} \times(0, T)\right)}^{2}+C_{2}\left\|e^{-s \varphi(0, t)} g\right\|_{L^{2}(0, T)}^{2}+C_{2} s \int_{0}^{T}(t(T-t))^{-\alpha}\left|\partial_{x} z\left(-L_{1}, t\right)\right|^{2} d t \tag{3.3.29}
\end{align*}
$$

Finally, let us come back to the original variables. By definition of $z$, we know that

$$
\begin{equation*}
s^{3} \int_{0}^{T} \int_{\Omega_{L}} e^{-2 s \varphi}(t(T-t))^{-3 \alpha}|y|^{2} d x d t=s^{3} \int_{0}^{T} \int_{\Omega_{L}}(t(T-t))^{-3 \alpha}|z|^{2} d x d t, \forall s>0 . \tag{3.3.30}
\end{equation*}
$$

Moreover, since $\partial_{x} z=-s e^{-s \varphi} \partial_{x} \varphi y+e^{-s \varphi} \partial_{x} y$, in $\Omega_{L} \times(0, T)$ we have

$$
\begin{align*}
& s \int_{0}^{T} \int_{\Omega_{L}} e^{-2 s \varphi}(t(T-t))^{-\alpha}\left|\partial_{x} y\right|^{2} d x d t \\
\leq & C s^{3} \int_{0}^{T} \int_{\Omega_{L}}(t(T-t))^{-3 \alpha}|z|^{2} d x d t+C s \int_{0}^{T} \int_{\Omega_{L}}(t(T-t))^{-\alpha}\left|\partial_{x} z\right|^{2} d x d t, \quad \forall s>0 . \tag{3.3.31}
\end{align*}
$$

In the same manner, the global terms of $\partial_{x}^{2} y$ and $\partial_{t} y$ can be estimated in the following way:

$$
\begin{align*}
& s^{-1} \int_{0}^{T} \int_{\Omega_{L}} e^{-2 s \varphi}(t(T-t))^{\alpha}\left|\partial_{x}^{2} y\right|^{2} d x d t \\
\leq & C s^{3} \int_{0}^{T} \int_{\Omega_{L}}(t(T-t))^{-3 \alpha}|z|^{2} d x d t+C s \int_{0}^{T} \int_{\Omega_{L}}(t(T-t))^{-\alpha}\left|\partial_{x} z\right|^{2} d x d t  \tag{3.3.32}\\
& +C s^{-1} \int_{0}^{T} \int_{\Omega_{L}}(t(T-t))^{\alpha}\left|\partial_{x}^{2} z\right|^{2} d x d t, \quad \forall s>0 .
\end{align*}
$$

and

$$
\begin{align*}
& s^{-1} \int_{0}^{T} \int_{\Omega_{L}} e^{-2 s \varphi}(t(T-t))^{\alpha}\left|\partial_{t} y\right|^{2} d x d t \\
\leq & C s \int_{0}^{T} \int_{\Omega_{L}}(t(T-t))^{-2 \alpha-1}|z|^{2} d x d t+C s^{-1} \int_{0}^{T} \int_{\Omega_{L}}(t(T-t))^{\alpha}\left|\partial_{t} z\right|^{2} d x d t, \quad \forall s>0 . \tag{3.3.33}
\end{align*}
$$

Then, using (3.3.30)-(3.3.33) in (3.3.32) we obtain 3.3.4. This completes the proof of Lemma 3.9 .

Remark 3.12 We point out that Lemma 3.9 can be used to prove a boundary controllability for problems in the form

$$
\begin{cases}\partial_{t} u_{L}(x, t)-\partial_{x}^{2} u_{L}(x, t)=0, & \forall(x, t) \in \Omega_{L} \times(0, T),  \tag{3.3.34}\\ \left(u(x, 0), u_{\Gamma}(0)\right)=\left(u_{0}(x), u_{\Gamma, 0}\right), & \forall x \in \Omega_{L}, \\ u_{\Gamma}(t)=u(0, t), & \forall t \in(0, T), \\ u\left(-L_{1}, t\right)=v(t), & \forall t \in(0, T), \\ u_{\Gamma}^{\prime}(t)+\partial_{x} u(0, t)=0, & \forall t \in(0, T) .\end{cases}
$$

with $\left(u_{0}, u_{\Gamma, 0}\right) \in L^{2}\left(\Omega_{L}\right) \times \mathbb{R}$ the initial data and control $v \in L^{2}(0, T)$ acts only on the flux of solutions in the left-hand side of the domain $x=-L_{1}$. In fact, the adjoint of the control problem of (3.3.34) is the same as the one in (3.3.1) but in this case we have to prove the following observability inequality:

$$
\|z(\cdot, 0)\|_{L^{2}\left(\Omega_{L}\right)}^{2}+\left|z_{\Gamma}(0)\right|^{2} \leq C \int_{0}^{T}\left|\partial_{x} z\left(-L_{1}, t\right)\right|^{2} d t
$$

for some constant $C=C\left(\Omega_{L}, T\right)$.

### 3.4 Convergence of the approximate system

In this section, the goal is to prove Theorem 2. To do this, we introduce the notions of weak solutions in the sense of distributions for the problems (3.1.2) and (3.1.4).

Let us start giving a remark on the approximate system. As we said before, sometimes parabolic equations with discontinuous diffusion coefficients can be viewed as a transmission problem. This means that (3.1.2) can be written in the following way:

$$
\begin{cases}\partial_{t} u_{L}^{K}(x, t)-\partial_{x}^{2} u_{L}^{K}(x, t)=\chi_{\omega}(x) v^{K}(x, t), & \forall(x, t) \in \Omega_{L} \times(0, T),  \tag{3.4.1}\\ \partial_{t} u_{R}^{K}(x, t)-K^{2} \partial_{x}^{2} u_{R}^{K}(x, t)=0, & \forall(x, t) \in \Omega_{R} \times(0, T), \\ u_{L}^{K}(x, 0)=u_{0, L}(x), & \forall x \in \Omega_{L}, \\ u_{R}^{K}(x, 0)=u_{0, R}(x), & \forall x \in \Omega_{R}, \\ u_{R}^{K}\left(0^{+}, t\right)=u_{L}^{K}\left(0^{-}, t\right), & \forall t \in(0, T), \\ K^{2} \partial_{x} u_{R}\left(0^{+}, t\right)=\partial_{x} u_{L}\left(0^{-}, t\right), & \forall t \in(0, T), \\ u_{L}^{K}\left(-L_{1}, t\right)=\partial_{x} u_{R}^{K}(1, t)=0, & \forall t \in(0, T) .\end{cases}
$$

We point out that equation (3.4.1 ${ }_{5}$ and $(3.4 .1)_{6}$ describes the continuity of the solution and the flux at $x=0$.

Let $K>0$ and $T>0$. We say that $u \in C^{0}\left([0, T] ; L^{2}(\Omega)\right)$ is a weak solution of (3.4.1)
if for all $\tau \in[0, T]$ the following inequality holds

$$
\begin{align*}
& -\int_{0}^{\tau} \int_{\Omega_{L}} u_{L}^{K}\left(\partial_{t} \psi_{L}+\partial_{x}^{2} \psi_{L}\right) d x d t-\int_{0}^{\tau} \int_{\Omega_{R}} u_{R}^{K}\left(\partial_{t} \psi_{R}+K^{2} \partial_{x}^{2} \psi_{R}\right) d x d t \\
& +\int_{\Omega_{R}}\left(u_{L}^{K}(x, \tau) \psi_{L}(x, \tau)-u_{0, L}(x) \psi_{L}(x, 0)\right) d x+\int_{\Omega_{R}}\left(u_{R}^{K}(x, \tau) \psi_{R}(x, \tau)-u_{0, R}(x) \psi_{R}(x, 0)\right) d x \\
& +\int_{0}^{\tau} u_{L}^{K}(0, t)\left(\partial_{x} \psi_{L}(0, t)-K^{2} \partial_{x} \psi_{R}(0, t)\right) d t=\int_{0}^{\tau} \int_{\Omega_{L}} \chi_{\omega} v^{K} \psi_{L} d x d t \tag{3.4.2}
\end{align*}
$$

for all $\psi=\psi(x, t) \in \Psi(\tau)$, where

$$
\Psi(\tau)=\left\{\psi \in C^{0}(\bar{\Omega} \times[0, \tau]) ; \psi_{L}, \psi_{R} \text { are smooth and } \psi_{L}\left(-L_{1}, t\right)=\partial_{x} \psi_{R}(1, t)=0\right\}
$$

On the other hand, we say that $y \in C^{0}\left([0, T] ; L^{2}(\Omega)\right)$ is a weak solution of (3.1.2) if

$$
\begin{align*}
& -\int_{0}^{\tau} \int_{\Omega_{L}} y_{L}\left(\partial_{t} \phi_{L}+\partial_{x}^{2} \phi_{L}\right) d x d t+\int_{\Omega_{L}}\left(y_{L}(x, \tau) \phi_{L}(x, \tau)-y_{0, L}(x) \phi_{L}(x, 0)\right) d x \\
& y_{R}(\tau) \phi_{R}(\tau)-y_{0, \Gamma} \phi(0)+\int_{0}^{\tau} y_{\Gamma}(t)\left(\partial_{x} \phi_{L}(0, t)-\phi_{R}^{\prime}(t)\right) d t  \tag{3.4.3}\\
= & \int_{0}^{\tau} \int_{\Omega_{L}} f \phi_{R} d x d t+\int_{0}^{\tau} g(t) \phi_{L}(0, t) d t,
\end{align*}
$$

for all $\phi \in \Phi(\tau)$ where
$\Phi(\tau)=\left\{\phi \in C^{0}(\bar{\Omega} \times[0, T]) ; \phi_{L}\right.$ is smooth,$\partial_{x} \phi_{R}=0$ in $\Omega_{R}$ and $\left.\phi_{L}\left(-L_{1}, t\right)=\partial_{x} \phi_{R}(1, t)=0\right\}$

Now we have all the ingredients to start the proof of Theorem 3.2.

Proof of Theorem 3.2. First, we recall that for each $K>0$, (3.4.1) admits a unique weak solution

$$
u^{K} \in C^{0}\left([0, T] ; L^{2}(\Omega)\right) \cap L^{2}\left(0, T ; H_{L}^{1}(\Omega)\right),
$$

with $u_{0} \in L^{2}(\Omega)$ and $v \in L^{2}(\omega \times(0, T))$. Moreover, the following energy estimate holds:

$$
\begin{equation*}
\left\|u^{K}\right\|_{C^{0}\left([0, T] ; L^{2}(\Omega)\right)}+\left\|\sigma^{K} \partial_{x} u^{K}\right\|_{L^{2}(\Omega \times(0, T))}^{2} \leq C\left(\left\|u_{0}\right\|_{L^{2}(\Omega)}^{2}+\left\|v^{K}\right\|_{L^{2}(\omega \times(0, T))}^{2}\right) \tag{3.4.4}
\end{equation*}
$$

for some positive constant $C=(\Omega, \omega, T)$ independent of $K$. Since $v^{K}$ converges weakly to $v$, we deduce that

$$
\left(u^{K}\right)_{K>1} \text { is uniformly bounded } L^{\infty}\left(0, T ; L^{2}(\Omega)\right) \cap L^{2}\left(0, T ; H_{L}^{1}(\Omega)\right) .
$$

Then, there exists a subsequence $\left(u^{K}\right)_{K>1}$ (which denotes by the same index for simplicity) such that (3.1.6) and (3.1.7) holds. Moreover, we can use classical compactness results (see for instance [25], 41] and [89]) to deduce that

$$
u^{K} \rightarrow u \text { strongly in } L^{2}(\Omega \times(0, T))
$$

It remains to identify the limit problem for $u$. First, from (3.4.4) it is easy to see that

$$
\left\|\partial_{x} u_{R}^{K}\right\|_{L^{2}\left(\Omega_{R} \times(0, T)\right)} \leq C K^{-2}
$$

and as $K \rightarrow+\infty$ we see that $\partial_{x} u_{R}=0$ a.e. in $\Omega_{R}$, i.e., $u_{R}=u_{R}(t)$ is a function of $t$ only. Moreover, Trace Theorem implies that $u_{L}\left(-L_{1}, t\right)=0$ and $u_{R}(t)=u_{L}(0, t), \forall t \in(0, T)$.

Now we focus on the weak solution (3.4.2). In order to avoid the explicit dependence of $K$, we choose $\psi=\phi$ with $\phi \in \Phi(\tau)$. Then, we have

$$
\begin{align*}
& -\int_{0}^{\tau} \int_{\Omega_{L}} u_{L}^{K}\left(\partial_{t} \phi_{L}+\partial_{x}^{2} \phi_{L}\right) d x d t-\int_{0}^{\tau} \int_{\Omega_{R}} u_{R}^{K} \phi^{\prime}(t) d x d t \\
& +\int_{\Omega_{L}}\left(u_{L}^{K}(x, \tau) \phi_{L}(x, \tau)-u_{0, L} \phi_{L}(x, 0)\right) d x+\phi(\tau) \int_{\Omega_{R}} u_{R}^{K} d x-\phi(0) \int_{\Omega_{R}} u_{0, R} d x  \tag{3.4.5}\\
& +\int_{0}^{\tau} u_{L}^{K}(0, t) \partial_{x} \phi_{L}(0, t) d t=\int_{0}^{\tau} \int_{\Omega_{L}} \chi_{\omega} v^{K} \phi_{L} d x d t .
\end{align*}
$$

Letting $K \rightarrow+\infty$ in (3.4.5), using the fact that $u_{R}(t)=u_{L}(0, t)$ and $\left|\Omega_{R}\right|=1$ we get

$$
\begin{aligned}
& -\int_{0}^{\tau} \int_{\Omega_{L}} u_{L}\left(\partial_{t} \phi_{L}+\partial_{x}^{2} \phi_{L}\right) d x d t+\int_{\Omega_{L}}\left(u_{L}(x, \tau) \phi_{L}(x, \tau)-u_{0, L}(x) \phi_{L}(x, 0)\right) d x \\
& +u_{R}(\tau) \phi_{R}(\tau)-\phi(0) \int_{\Omega_{R}} u_{0, R} d x+\int_{0}^{\tau} u_{L}(0, t)\left(\partial_{x} \phi_{L}(0, t)-\phi_{R}^{\prime}(t)\right) d t \\
= & \int_{0}^{\tau} \int_{\Omega_{L}} \chi_{\omega} v \phi_{L} d x d t,
\end{aligned}
$$

which is the definition of weak solution (3.4.3) for (3.1.2) with $f=\chi_{\omega} v, g=0, y_{0, L}=u_{0, L}$ and $y_{0, \Gamma}=\int_{\Omega_{R}} u_{0, R} d x$. This completes the proof of Theorem 3.2 .

### 3.5 Plugging the limit control in the approximate system

In this section we focus on the proof of Theorem 3.3. In general terms, the proof is based on the convergence result given in Theorem 3.2 together with regularity results of (3.1.9) and (3.1.10).

Proof of Theorem [3.3. Let us define $y=u-u^{K}$ in $\bar{\Omega} \times[0, T]$. It is clear that $y$ depends on $K$ and therefore we shall write $y^{K}$ instead of $y$, however in this case we avoid this dependence for simplicity. Then, according to the equations (3.1.9) and 3.1.10 y is a solution of

$$
\begin{cases}\partial_{t} y_{L}(x, t)-\partial_{x}^{2} y_{L}(x, t)=0, & \forall(x, t) \in \Omega_{L} \times(0, T),  \tag{3.5.1}\\ \partial_{t} y_{R}(x, t)-K^{2} \partial_{x}^{2} y_{R}(x, t)=\partial_{x} u_{L}(0, t), & \forall(x, t) \in \Omega_{R} \times(0, T), \\ y_{L}(x, 0)=0, & \forall x \in \Omega_{L}, \\ y_{R}(x, 0)=u_{0, R}(x)-\int_{\Omega_{R}} u_{0, R} d x, & \forall x \in \Omega_{R}, \\ y_{R}(0, t)=y_{L}(0, t), & \forall t \in(0, T), \\ K^{2} \partial_{x} y_{R}(0, t)=\partial_{x} y_{L}(0, t)-\partial_{x} u_{L}(0, t), & \forall t \in(0, T), \\ y\left(-L_{1}, t\right)=\partial_{x} y_{R}(1, t)=0, & \forall t \in(0, T) .\end{cases}
$$

Due to the regularity of $u$ and $u^{K}$ (see for instance Propositions 3.7 and 3.8 in Section 2 and [80] or [41]), and classical arguments we can write for all $t \in(0, T)$

$$
\begin{aligned}
& \frac{1}{2} \frac{d}{d t} \int_{\Omega_{L}}\left|y_{L}(t)\right|^{2} d x+\frac{1}{2} \frac{d}{d t} \int_{\Omega_{R}}\left|y_{R}(t)\right|^{2} d x+\int_{\Omega} \sigma^{K}\left|\partial_{x} y(t)\right|^{2} d x \\
= & -\partial_{x} u_{L}(0, t) \int_{\Omega_{R}} y_{R} d x-y_{L}(0, t) \partial_{x} u_{L}(0, t) .
\end{aligned}
$$

Thus, for all $\tau \in[0, T]$ we have

$$
\begin{align*}
\int_{\Omega_{L}}|y(\tau)|^{2} d x \leq & 2 \int_{0}^{T} \int_{\Omega_{L}} y_{R} \partial_{x} u_{L}(0, t) d x d t+2 \int_{0}^{T} y_{L}(0, t) \partial_{x} u_{L}(0, t) d t  \tag{3.5.2}\\
& +\int_{\Omega_{R}}\left|u_{0, R}-\overline{u_{0, R}}\right|^{2} d x
\end{align*}
$$

Now, by (3.1.6) we choose $K_{1}>1$ such that for all $K \geq K_{1}$ we have

$$
\begin{equation*}
\left|\int_{0}^{T} \int_{\Omega} y w d x d t\right| \leq \frac{\varepsilon^{2}}{3}, \quad \forall w \in L^{2}(\Omega \times(0, T)) . \tag{3.5.3}
\end{equation*}
$$

Moreover, by Trace Theorem and the strong convergence of $(y)_{K>1}$ of we choose $K_{2}>1$ such that for all $K \geq K_{2}$ we get

$$
\begin{equation*}
\left|\int_{0}^{T} y(0, t) z d t\right| \leq \frac{\varepsilon^{2}}{3}, \quad \forall z \in L^{2}(0, T) . \tag{3.5.4}
\end{equation*}
$$

Then, applying the inequalities (3.1.11), (3.5.2) and (3.5.3), choosing $\tau=T$ and by using the fact that $u(T)=0$ in $\bar{\Omega}$ we get

$$
\left\|u^{K}(\cdot, T)\right\|_{L^{2}(\Omega \times(0, T))} \leq \varepsilon, \quad \forall K \geq K_{0}
$$

with $K_{0}=\max \left\{K_{1}, K_{2}\right\}$. This completes the proof of Theorem 3.3.

## Chapter 4

## Controllability of 1-D heat equation with discontinuous diffusion coefficients

In this chapter, we will study the null-controllability of the heat equation with Lipschitz diffusion coefficient and mixed boundary conditions. In this case, we consider a boundary control acting in the left-hand side of the domain. This will be done by using a suitable Carleman estimate, which implies the observability inequality for the adjoint system asociated to the original problem.

### 4.1 Introduction, setting and main result

Let $\Omega=\left(-L_{1}, L_{2}\right)$ be an open interval of $\mathbb{R}$. Throughout this chapter, we consider the following heat equation with mixed boundary conditions:

$$
\begin{cases}\partial_{t} u(x, t)-\partial_{x}\left(\sigma(x) \partial_{x} u(x, t)\right)=0, & \forall(x, t) \in \Omega \times(0, T),  \tag{4.1.1}\\ u(x, 0)=u^{0}(x), & \forall x \in \Omega, \\ u\left(-L_{1}, t\right)=v(t), & \forall t \in(0, T), \\ \partial_{x} u\left(L_{2}, t\right)=0, & \forall t \in(0, T) .\end{cases}
$$

Here, the piecewise diffusion coefficient $\sigma=\sigma(x)$ is defined by

$$
\sigma(x)= \begin{cases}\sigma_{1}^{2}, & \forall x \in\left(-L_{1}, 0\right), \\ \sigma_{2}^{2}, & \forall x \in\left(0, L_{2}\right)\end{cases}
$$

where $\sigma_{1}, \sigma_{2}>0$. Moreover, the initial datum $u^{0}$ belongs in $L^{2}(\Omega)$ is given and $v \in$ $L^{2}(0, T)$ is the boundary control. It is well-known that, under these assumptions, the problem (4.1.1) is well-posed in the Hadamard's sense. As we said before, we are interested in the problem of null-controllability of 4.1.1). In other words, we focus in the following question: given $u^{0} \in L^{2}(\Omega)$ in 4.1.1 and $T>0$, can we find a boundary control $v \in$ $L^{2}(0, T)$ such that the associated solution $u$ satisfy

$$
u(x, T)=0, \quad \forall x \in \Omega ?
$$

It is well-known that the question above is equivalent to prove the observability inequality for the adjoint system of (4.1.1). To be more precise, let $w$ be the solution of

$$
\begin{cases}\partial_{t} w(x, t)+\partial_{x}\left(\sigma(x) \partial_{x} w(x, t)\right)=0, & \forall(x, t) \in \Omega \times(0, T),  \tag{4.1.2}\\ w(T)=w^{T}, & \forall x \in \Omega, \\ w\left(-L_{1}, t\right)=0, & \forall t \in(0, T), \\ \partial_{x} w\left(L_{2}, t\right)=0, & \forall t \in(0, T) .\end{cases}
$$

Then, the observability property of $(4.1 .2$ is the following: can we find a constant $C_{\text {obs }}>0$ such that every solution $w$ of 4.1.2) satisfies

$$
\|w(x, 0)\|_{L^{2}(\Omega)}^{2} \leq C_{o b s} \int_{0}^{T}\left|\partial_{x} w\left(-L_{1}, t\right)\right|^{2} d t ?
$$

The main ingredient in the proof of the observability inequality is a suitable Carleman estimate where the observation is on the flux of the left-hand side of the domain. More precisely, we will get an estimate of the following form

$$
\begin{aligned}
& s^{3} \int_{0}^{T} \int_{\Omega} \rho|z|^{2} d x d t+s \int_{0}^{T} \int_{\Omega} \mu\left|\partial_{x} z\right|^{2} d x d t \\
\leq & C \int_{0}^{T} \int_{\Omega} \nu\left|\partial_{t} \pm \partial_{x}\left(\sigma \partial_{x} z\right)\right|^{2} d x d t+C s \int_{0}^{T} \mu\left(t,-L_{1}\right)\left|\partial_{x} z\left(t,-L_{1}\right)\right|^{2} d t,
\end{aligned}
$$

for some functions $\rho, \mu$ and $\nu$ depending on space and time.
Controllability issues of parabolic equations have been intensely studied by several authors since the 70 's. For this reason, we will mention some of the most important results concerning scalar parabolic equations with smooths coefficients. In [42] and [43] H.O. Fattorini and D.L. Russel obtained for first time results about null boundary controllability for the one dimensional heat equation using the so-called method of moments. In contrast, in 88 the author proved a null controllability result for N -dimensional heat equation with a boundary control supported on the whole boundary of the domain. To be more specific, he proved that the null controllability of wave equation at a positive time implies the null controllability of the heat equation at any positive time.

In 1995, the null-controllability of the heat equation for high-dimensional case was solved by G. Lebeau and L. Robbiano [78] using a spectral inequality which was proved using local Carleman estimates. On the other hand, in 1996 the same problem was solved by A. Fursikov and O. Imanuvilov in [50]. Besides, the authors consider a general parabolic operator. This result was obtained by proving Carleman estimate for a general parabolic equation and for an arbitrary internal observation region.

The literature is also rich about controllability of other types of parabolic equations like Stokes or Navier-Stokes. For a deeper discussion, see for instance the survey of E. Fernandez-Cara and S. Guerrero [48] (see also [8] and the references given there).

Now we reduce our scope to controllability issues for parabolic problems in the case of non-smooth coefficients. In 2002, E. Fernandez-Cara and E. Zuazua in 49 proved a
controllability result for $1-D$ linear parabolic equations for coefficients with bounded variations using the Russel's method (see [88]).

In [39] the authors proved a Carleman estimate and consequently a null controllability result for a semilinear heat equation in the case where the control is supported in the region where the diffusion coefficient is the 'lowest'. The key is the construction of a non-smooth weight function satisfying the same transmission condition as the solution.

In [20], the authors achieve a Carleman estimate for the operators of the form $\partial_{t} \pm$ $\partial_{x}\left(\sigma \partial_{x}\right)$ without any restriction on the observation region, but this strategy does not extend to higher-dimensional cases. In the one-dimensional case, this monotonicity assumption of [39] on the diffusion coefficient was relaxed in [20] and [19] introducing more requirements on the non-smooth weight function. Also, they achieve a Carleman estimate with boundary observation in the Dirichlet case and finite jumps on the diffusion coefficient on the domain.

In [74, J. Le Rousseau derived a Carleman estimate for the problem above where the diffusion coefficient $\sigma$ is a bounded variation function. The proofs relies in the idea of approximate the diffusion coefficient $\sigma$ of $\partial_{t} \pm \partial_{x}\left(\sigma \partial_{x}\right)$ by a sequence of piecewise functions $\sigma_{\varepsilon}$ and study the controllability properties of each problem with $\partial_{t} \pm \partial_{x}\left(\sigma_{\varepsilon} \partial_{x}\right)$ and later pass to the limit. The main issue in this limiting process is to keep both the weight functions and constants in the Carleman estimate under control. They also obtain Carleman estimates in the case of boundary observation considering Dirichlet boundary conditions for bounded variation diffusion coefficients.

On the other hand, in [76] the authors obtain a Carleman estimate for an operator of the type $\nabla \cdot(c(x) \nabla z)$ without any isotropy assumption. Specifically, in this article $c$ is a symmetric positive-definite matrix with a jump discontinuity across a smooth hypersurface. Also, they give conditions on the Carleman weight functions that are rather simple to handle, and they prove that these functions are sharp.

Recently, in [82] the authors achieve a null controllability result for one-dimensional parabolic equation with generalized Robin-Neumann conditions at both extremities, with one boundary control. They following the flatness approach. Also, they obtain some numerical results on the reconstruction of the control for the problem above.

Our next task is formulate the main result of this work. To do this, we will introduce some notation. Let $A: D(A) \subset L^{2}(\Omega) \rightarrow L^{2}(\Omega)$ be the operator formally defined by

$$
A=\partial_{x}\left(\sigma \partial_{x}\right),
$$

and its domain of $A$ is given by

$$
D(A)=\left\{u \in H^{1}(\Omega) ; \sigma \partial_{x} u \in H^{1}(\Omega), u\left(-L_{1}\right)=0\right\} .
$$

Throught this section, we consider the following system:

$$
\begin{cases}\partial_{t} z(x, t)-\partial_{x}\left(\sigma(x) \partial_{x} z(x, t)\right)=f(x, t), & \forall(x, t) \in \Omega \times(0, T),  \tag{4.1.3}\\ z(x, 0)=z^{0}(x), & \forall x \in \Omega, \\ z\left(-L_{1}, t\right)=0, & \forall t \in(0, T), \\ \partial_{x} z\left(L_{2}, t\right)=0, & \forall t \in(0, T) .\end{cases}
$$

Suppose that $f \in L^{2}(\Omega \times(0, T))$ and $z^{0} \in L^{2}(\Omega)$. It is clear that, under these assumptions, the problem $(4.1 .3)$ is well-posed. Moreover, due to the classical semigroup approach, the solution $z$ of (4.1.3) satisfies

$$
\begin{equation*}
z(\cdot, t) \in D(A), \text { for all } t \in(0, T) \tag{4.1.4}
\end{equation*}
$$

Before going further, it is convenient to write the system (4.1.3) in terms of each parts of the domain separated by the interface located in $\{0\}$. To be more specific, let us $\Omega_{1}=\left(-L_{1}, 0\right)$ and $\Omega_{2}=\left(0, L_{2}\right)$. Here and consequently, for a spatial function $h$ defined on the domain $\Omega, h_{j}$ stands its restriction to the subdomains $\Omega_{j}$, for each $j=1,2$.

Thus, with this notation, system 4.1.3) can be written as follows:

$$
\begin{cases}\partial_{t} z_{1}(x, t)-\sigma_{1}(x)^{2} \partial_{x}^{2} z_{1}(x, t)=f_{1}(x, t), & \forall(x, t) \in \Omega_{1} \times(0, T),  \tag{4.1.5}\\ \partial_{t} z_{2}(x, t)-\sigma_{2}(x)^{2} \partial_{x}^{2} z_{2}(x, t)=f_{2}(x, t), & \forall(x, t) \in \Omega_{2} \times(0, T), \\ z_{1}(x, 0)=z_{1}^{0}(x), & \forall x \in \Omega_{1}, \\ z_{2}(x, 0)=z_{2}^{0}(x), & \forall x \in \Omega_{2}, \\ z_{2}(0, t)=z_{1}(0, t), & \forall t \in(0, T), \\ \sigma_{2}^{2} \partial_{x} z_{2}(0, t)=\sigma_{1}^{2} \partial_{x} z_{1}(0, t), & \forall t \in(0, T), \\ z_{1}\left(-L_{1}, t\right)=0, & \forall t \in(0, T), \\ \partial_{x} z_{2}\left(L_{2}, t\right)=0, & \forall t \in(0, T), \\ \partial_{x} z_{1}\left(-L_{1}, t\right)=g_{N}(t), & \forall t \in(0, T) .\end{cases}
$$

The boundary conditions $\left.4 .{ }^{4.5}\right)_{3}$ and 4. A $_{4}$ are the transmission conditions of the system. Let us emphasize that these conditions arise due to the regularity of the solution $z$ in (4.1.4).

Let us introduce the weight functions that we will use for state the Carleman inequality for systems like 4.1.3). For each $j=1,2$, we consider the positive functions $\phi_{j}=\phi_{j}(x) \in$ $C^{2}\left(\overline{\Omega_{j}}, \mathbb{R}\right)$ such that $\phi_{1}(0)=\phi_{2}(0)$ and

$$
\begin{cases}m_{j} \leq \sigma_{j} \phi_{j}^{\prime}(x) \leq M_{j}, & \text { in } \Omega_{j},  \tag{4.1.6}\\ r_{j} \leq-\sigma_{j} \phi_{j}^{\prime \prime}(x) \leq R_{j}, & \text { in } \Omega_{j},\end{cases}
$$

for some positive constants $m_{j}, M_{j}, r_{j}$ and $R_{j}$ with $M_{2}<m_{1}$. Note that the assumptions in (4.1.6) imply:

$$
m_{1}=\sigma_{1} \phi_{1}^{\prime}(0), M_{1}=\sigma_{1} \phi_{1}\left(-L_{1}\right)^{\prime}, m_{2}=\sigma_{2} \phi_{2}\left(L_{2}\right)^{\prime}, M_{2}=\sigma_{2} \phi_{2}(0)^{\prime}
$$

The Figure 2.1 sketch a prototype of functions $\phi_{1}$ and $\phi_{2}$ that we will use:


Figure 4.1: Sketch of the jump of the derivative of functions $\phi_{1}$ and $\phi_{2}$

For instance, note that the following class of functions:

$$
\phi_{1}(x)=-\frac{\left(M_{1}-m_{1}\right)}{2 L_{1} \sigma_{1}} x^{2}+\frac{m_{1}}{\sigma_{1}} x+a, \quad \phi_{2}(x)=-\frac{\left(M_{2}-m_{2}\right)}{2 L_{2} \sigma_{2}} x^{2}+\frac{M_{2}}{\sigma_{2}} x+a, \quad a \in \mathbb{R},
$$

satisfy the above assumptions, for $a \in \mathbb{R}$ sufficiently large such that $\phi_{1}$ and $\phi_{2}$ are positive.
Additionally, for a parameter $\alpha>0$ and $j=1,2$, let us denote by $\varphi_{j}$ the following functions:

$$
\begin{equation*}
\varphi_{j}(x, t)=\theta^{\alpha}(t) \phi_{j}(x), \quad \forall(x, t) \in \Omega_{j} \times(0, T) \tag{4.1.7}
\end{equation*}
$$

where $\theta(t)=t^{-1}(T-t)^{-1}$. For simplicity of notation, we ignore the dependence of $\alpha$ and $\sigma_{j}$ on $\varphi_{j}$.

We can now formulate the Carleman estimate for 1-D heat equation:
Theorem 4.1 Consider the functions $\varphi_{j}$ in (4.1.7) for $\alpha \geq 2$. Moreover, assume that $3 \alpha \geq 2 \beta$, with $\beta \geq 1$. Then, there exist two positive constants $C$ and $s_{*}$ such that for all
$s \geq s_{*}$, the following inequality holds:

$$
\begin{aligned}
& s^{3} \int_{0}^{T} \int_{\Omega_{1}} e^{-2 s \varphi_{1}} \theta^{3 \alpha-\beta}\left|z_{1}\right|^{2} d x d t+s^{3} \int_{0}^{T} \int_{\Omega_{2}} e^{-2 s \varphi_{2}} \theta^{3 \alpha-\beta}\left|z_{2}\right|^{2} d x d t \\
& \quad+s \int_{0}^{T} \int_{\Omega_{1}} e^{-2 s \varphi_{1}} \theta^{\alpha-\beta}\left|\partial_{x} z_{1}\right|^{2} d x d t+s \int_{0}^{T} \int_{\Omega_{2}} e^{-2 s \varphi_{2}} \theta^{\alpha-\beta}\left|\partial_{x} z_{2}\right|^{2} d x d t \\
& \quad+s^{3} \int_{0}^{T} e^{-2 s \varphi_{1}(t, 0)} \theta^{3 \alpha-\beta}\left|z_{1}(t, 0)\right|^{2} d t+s^{3} \int_{0}^{T} e^{-2 s \varphi_{2}\left(t, L_{2}\right)} \theta^{3 \alpha-\beta}\left|z_{2}\left(t, L_{2}\right)\right|^{2} d t \\
& \leq C \int_{0}^{T} \int_{\Omega_{1}} e^{-2 s \varphi_{1}} \theta^{-\beta}\left|\partial_{t} z_{1} \pm \sigma_{1}^{2} \partial_{x}^{2} z_{1}\right|^{2} d x d t+C \int_{0}^{T} \int_{\Omega_{2}} e^{-2 s \varphi_{2}} \theta^{-\beta}\left|\partial_{t} z_{2} \pm \sigma_{2}^{2} \partial_{x}^{2} z_{2}\right|^{2} d x d t \\
& \quad+C s \int_{0}^{T} e^{-2 s \varphi_{1}\left(t,-L_{1}\right)} \theta^{\alpha-\beta}\left|\partial_{x} z_{1}\left(-L_{1}\right)\right|^{2} d t,
\end{aligned}
$$

for each function $z_{j} \in L^{2}\left(\Omega_{j} \times(0, T)\right)$ such that $\partial_{t} z_{j} \pm \sigma_{j}^{2} \partial_{x}^{2} z_{j} \in L^{2}\left(\Omega_{j} \times(0, T)\right)$, for each $j=1,2$, and $z_{1}\left(-L_{1}, t\right)=\partial_{x} z_{2}\left(L_{2}, t\right)=0$ for each $t \in(0, T)$.

As a consequence of Theorem 4.1, we have the following:
Corollary 4.2 Let $\left(u_{1}, u_{2}\right)$ be the solution of the following system

$$
\begin{cases}\partial_{t} u_{1}(x, t)-\sigma_{1}^{2} \partial_{x}^{2} u_{1}(x, t)=0, & \forall(x, t) \in \Omega_{1} \times(0, T),  \tag{4.1.8}\\ \partial_{t} u_{2}(x, t)-\sigma_{2}^{2} \partial_{x}^{2} u_{2}(x, t)=0, & \forall(x, t) \in \Omega_{2} \times(0, T), \\ u_{1}(x, 0)=u_{1}^{0}(x), & \forall x \in \Omega_{1}, \\ u_{2}(x, 0)=u_{2}^{0}(x), & \forall x \in \Omega_{2}, \\ u_{1}\left(-L_{1}, t\right)=v(t), & \forall t \in(0, T), \\ \partial_{x} u_{2}\left(L_{2}, t\right)=0, & \forall t \in(0, T),\end{cases}
$$

where $u_{j}^{0} \in L^{2}\left(\Omega_{j}\right)$, for each $j=1,2$ and $v \in L^{2}(0, T)$. Then, system 4.1.8 is nullcontrollable at time $T>0$.

We will prove the Corollary 4.2 at the end of this chapter.

### 4.2 Proof of the Theorem 4.1

### 4.2.1 Setting

For simplicity, we prove the case where the operator involve is $\partial_{t} z_{j}-\sigma_{j}^{2} \partial_{x}^{2} z_{j}$ for each $j=1,2$. The other case is completely analogous. First, we start reformulating our parabolic problem 4.1.5). This will be done into three steps:

## - Step 1: Localization in time

Let us define $\eta$ as a smooth function compactly supported in $]-1,1[$ with $\eta(0)=1$. For $\beta \geq 0, \lambda \geq 1$ and $t_{0} \in(0, T)$, we define the function $\eta_{t_{0}}$ as follows:

$$
\begin{equation*}
\eta_{t_{0}}(t)=\eta\left(\frac{\lambda\left(t-t_{0}\right)}{t_{0}^{\beta}\left(T-t_{0}\right)^{\beta}}\right), \quad \forall t \in(0, T) \tag{4.2.1}
\end{equation*}
$$

and shorten notation, we write $\eta_{t_{0}}$ instead of $\eta_{\beta, \lambda, t_{0}}$. Moreover, for each $j=1,2$, we define

$$
z_{j, t_{0}}(x, t)=\eta_{t_{0}}(t) z_{j}(x, t) \quad \forall(x, t) \in \Omega_{j} \times(0, T),
$$

where $\left(z_{1}, z_{2}\right)$ is the solution of system (4.1.5). Then, the new variables $\left(z_{1, t_{0}}, z_{2, t_{0}}\right)$ solve the following problem:

$$
\begin{cases}\partial_{t} z_{1, t_{0}}-\sigma_{1}^{2} \partial_{x}^{2} z_{1, t_{0}}=\eta_{t_{0}} f_{1}+\partial_{t} \eta_{t_{0}} z_{1}, & \forall(x, t) \in \Omega_{1} \times(0, T), \\ \partial_{t} z_{2, t_{0}}-\sigma_{2}^{2} \partial_{x}^{2} z_{2, t_{0}}=\eta_{t_{0}} f_{2}+\partial_{t} \eta_{t_{0}} z_{2}, & \forall(x, t) \in \Omega_{2} \times(0, T), \\ z_{2, t_{0}}(0, t)=z_{1, t_{0}}(0, t), & \forall t \in(0, T), \\ \sigma_{2}^{2} \partial_{x} z_{2, t_{0}}(0, t)=\sigma_{1}^{2} \partial_{x} z_{1, t_{0}}(0, t), & \forall t \in(0, T), \\ z_{1, t_{0}}\left(-L_{1}, t\right)=0, & \forall t \in(0, T), \\ \partial_{x} z_{2, t_{0}}\left(L_{2}, t\right)=0, & \forall t \in(0, T), \\ \partial_{x} z_{1, t_{0}}\left(-L_{1}, t\right)=g_{N, t_{0}}(t), & \forall t \in(0, T),\end{cases}
$$

where $g_{N, t_{0}}=\eta_{t_{0}} z\left(t,-L_{1}\right)$, for each $t \in(0, T)$ is the observation of our new system.
Notice that the form of $\eta_{t_{0}}$ in (4.2.1) does not play an important role here. Besides, we have defined in this section for indicate the dependence of $t_{0}$ of the support of $z_{1, t_{0}}$ and $z_{2, t_{0}}$.

## - Step 2: Conjugation

For $\alpha>0$, we consider the weight functions defined in 4.1.7). For each $j=1,2$, let us denote by $Z_{j}$ the following function:

$$
Z_{j}=e^{-s \varphi_{j}} z_{j, t_{0}}=e^{-s \varphi_{j}} \eta_{t_{0}} z_{j}, \quad \forall(x, t) \in \Omega_{j} \times(0, T)
$$

where $s$ is the Carleman parameter associated to the weight functions (4.1.7). Once again, for abbreviation, we write $Z_{j}$ instead of $Z_{j, t_{0}}$. Then the unknown variables $\left(Z_{1}, Z_{2}\right)$ solve the following problem:

$$
\begin{cases}\partial_{t} Z_{1}-\sigma_{1}^{2}\left(\partial_{x}+s \theta^{\alpha}(t) \partial_{x} \phi_{1}\right)^{2} Z_{1}=g_{1}, & \forall(x, t) \in \Omega_{1} \times(0, T),  \tag{4.2.2}\\ \partial_{t} Z_{2}-\sigma_{2}^{2}\left(\partial_{x}+s \theta^{\alpha}(t) \partial_{x} \phi_{2}\right)^{2} Z_{2}=g_{2}, & \forall(x, t) \in \Omega_{2} \times(0, T), \\ Z_{2}(0, t)=Z_{1}(0, t), & \forall t \in(0, T), \\ \sigma_{2}^{2} \partial_{x} Z_{2}(0, t)=\sigma_{1}^{2} \partial_{x} Z_{1}(0, t)+\left(m_{1} \sigma_{1}-M_{2} \sigma_{2}\right) s \theta^{\alpha}(t) Z_{2}(0, t), & \forall t \in(0, T), \\ Z_{1}\left(-L_{1}, t\right)=0, & \forall t \in(0, T), \\ \partial_{x} Z_{2}\left(L_{2}, t\right)=-m_{2} \sigma_{2}^{-1} \theta^{\alpha}(t) s Z_{2}\left(L_{2}, t\right), & \forall t \in(0, T), \\ \partial_{x} Z_{1}\left(-L_{1}, t\right)=G_{N}(t), & \forall t \in(0, T) .\end{cases}
$$

Here, the function $G_{N}$ is defined by $G_{N}=e^{-s \varphi_{1}\left(-L_{1}, t\right)} g_{N, t_{0}}(t)$, in $(0, T)$. Moreover, for each $j=1,2$, the function $g_{j}$ is given by

$$
\begin{equation*}
g_{j}=e^{-s \varphi_{j}} \eta_{t_{0}} f_{j}+e^{-s \varphi_{j}} \partial_{t} \eta_{t_{0}} z_{j}+\alpha \theta^{\alpha-1}(t) \partial_{t} \theta(t) \phi_{j} Z_{j}, \quad \forall(x, t) \in \Omega_{j} \times(0, T) \tag{4.2.3}
\end{equation*}
$$

## - Step 3: Freeze in time

Now, our next task is to avoid the dependence of $t$ on the coefficients of the right-hand side of 4.2 .2$)_{1}$ and 4.2 .2$)_{2}$ and also on the boundary conditions 4.2.2 4 and 4.2.2 ${ }_{6}$.

To simplify our notation, for $t_{0}$ fixed, we denote by $s_{0}$ the following expression:

$$
\begin{equation*}
s_{0}=s \theta^{\alpha}\left(t_{0}\right) \tag{4.2.4}
\end{equation*}
$$

Then, we rewrite the system (4.2.2) as follows:

$$
\begin{cases}\partial_{t} Z_{1}-\sigma_{1}^{2}\left(\partial_{x}+\partial_{x} \phi_{1} s_{0}\right)^{2} Z_{1}=F_{1}, & \forall(x, t) \in \Omega_{1} \times(0, T),  \tag{4.2.5}\\ \partial_{t} Z_{2}-\sigma_{2}^{2}\left(\partial_{x}+\partial_{x} \phi_{2} s_{0}\right)^{2} Z_{2}=F_{2}, & \forall(x, t) \in \Omega_{2} \times(0, T), \\ Z_{2}(0, t)=Z_{1}(0, t), & \forall t \in(0, T), \\ \sigma_{2}^{2} \partial_{x} Z_{2}(0, t)=\sigma_{1}^{2} \partial_{x} Z_{1}(0, t)+\left(m_{1} \sigma_{1}-M_{2} \sigma_{2}\right) s_{0} Z_{2}(0, t)+H(t), & \forall t \in(0, T), \\ Z_{1}\left(-L_{1}, t\right)=0, & \forall t \in(0, T), \\ \partial_{x} Z_{2}\left(L_{2}, t\right)=-m_{2} \sigma_{2}^{-1} s_{0} Z_{2}\left(L_{2}, t\right)+J(t), & \forall t \in(0, T), \\ \partial_{x} Z_{1}\left(-L_{1}, t\right)=G_{N}(t), & \forall t \in(0, T),\end{cases}
$$

where the residual functions $H$ and $J$ are defined by

$$
\begin{align*}
H(t) & =\left(m_{1} \sigma_{1}-M_{2} \sigma_{2}\right)\left(\theta^{\alpha}(t)-\theta^{\alpha}\left(t_{0}\right)\right) s Z_{2}(0, t), \quad t \in(0, T),  \tag{4.2.6}\\
J(t) & =\left(\theta^{\alpha}\left(t_{0}\right)-\theta^{\alpha}(t)\right) m_{2} \sigma_{2}^{-1} s Z_{2}\left(L_{2}, t\right), \quad t \in(0, T) . \tag{4.2.7}
\end{align*}
$$

Furthermore, the source term $F_{j}$ in 4.2.5) is given by

$$
\begin{align*}
F_{j} & =e^{-s \varphi_{j}} \eta_{t_{0}} f_{j}+e^{-s \varphi_{j}} \partial_{t} \eta_{t_{0}} z_{j}+\alpha \theta^{\alpha-1}(t) \partial_{t} \theta(t) \phi_{j} Z_{j}+2 \sigma_{j}^{2}\left(s \theta^{\alpha}(t)-s_{0}\right) \partial_{x} \phi_{j} \partial_{x} Z_{j}  \tag{4.2.8}\\
& +\sigma_{j}^{2}\left(s \theta^{\alpha}(t)-s_{0}\right) \partial_{x}^{2} \phi_{j} Z_{j}+\sigma_{j}^{2}\left(s^{2} \theta^{2 \alpha}(t)-s_{0}^{2}\right)\left|\partial_{x} \phi_{j}\right|^{2} Z_{j},
\end{align*}
$$

for each $(x, t) \in \Omega_{j} \times(0, T)$ and $j=1,2$.
Let us emphasize that in this step, several residual functions of different nature appears, see the definition of $H, J$ and $F_{j}$ for instance. At the moment, we will treat these ones as a source terms and later we will choose the parameters $\alpha, \beta, \lambda$ and $s$ in order to eliminate them in the classical spirit of Carleman estimates.

As we said before, the proof of the main result is is deduced by the following:
Lemma 4.3 There exists a positive constant $C_{1}=C_{1}\left(\alpha, \beta, m_{1}, M_{1}, m_{2}, M_{2}, r_{1}, r_{2}, \sigma_{1}, \sigma_{2}\right)$ independent of $s$ such that each solution $\left(Z_{1}, Z_{2}\right)$ of (4.2.5) with $F_{j} \in L^{2}\left(\Omega_{j} \times(0, T)\right)$,
$j=1,2, G_{N} \in L^{2}(0, T)$ and $H, J \in L^{2}(0, T)$ satisfies

$$
\begin{align*}
& s_{0}^{3} \int_{0}^{T} \int_{\Omega_{1}}\left|Z_{1}\right|^{2} d x d t+s_{0}^{3} \int_{0}^{T} \int_{\Omega_{2}}\left|Z_{2}\right|^{2} d x d t+s_{0} \int_{0}^{T} \int_{\Omega_{1}}\left|\partial_{x} Z_{1}\right|^{2} d x d t \\
& \quad+s_{0} \int_{0}^{T} \int_{\Omega_{2}}\left|\partial_{x} Z_{2}\right|^{2} d x d t+s_{0}^{3} \int_{0}^{T}\left|Z_{2}(0, t)\right|^{2} d t+s_{0}^{3} \int_{0}^{T}\left|Z_{2}\left(L_{2}, t\right)\right|^{2} d t \\
& \leq C_{1} \int_{0}^{T} \int_{\Omega_{1}}\left|F_{1}\right|^{2} d x d t+C_{1} \int_{0}^{T} \int_{\Omega_{2}}^{T}\left|F_{2}\right|^{2} d x d t+C_{1} s_{0} \int_{0}^{T}\left|G_{N}\right|^{2} d t  \tag{4.2.9}\\
&+C_{1} s_{0} \int_{0}^{T}|H|^{2} d t+C_{1} s_{0} \int_{0}^{T}|J|^{2} d t,
\end{align*}
$$

for all $s \geq s_{*}$.
We will give the proof of the Lemma 4.3 later. Now, with this Lemma at hand, we can conclude the proof of Theorem 4.1.

### 4.2.2 End of the proof of Theorem 4.1

Our main goal is eliminate the residual terms appeared in Lemma 1 for each solution of the system 4.2.5). For an easier comprehension, the rest of the proof of Theorem 4.1 falls naturally into three parts:

## - Step 1: First residual terms: $H$ and $J$

Proposition 4.4 We assume the same hypotheses as Lemma 4.3. Additionally, suppose that $2 \beta+\alpha \geq 2$. Then, there exist a positive constant $\lambda_{*}$ such that for every $\lambda \geq \lambda_{*}$, the following inequality holds:

$$
\begin{equation*}
C_{1} s_{0} \int_{0}^{T}|H|^{2} d t+C_{1} s_{0} \int_{0}^{T}|J|^{2} d t \leq \frac{1}{2} s_{0}^{3} \int_{0}^{T}\left|Z_{2}(0, t)\right|^{2} d t+\frac{1}{2} s_{0}^{3} \int_{0}^{T}\left|Z_{2}\left(L_{2}, t\right)\right|^{2} d t, \tag{4.2.10}
\end{equation*}
$$

where $C_{1}$ is the constant associated to Lemma 4.3.

Proof. Our proof starts with the observation that for each $t \in \operatorname{supp}\left(\eta_{t_{0}}\right)$ satisfies the following:

$$
\begin{equation*}
\left|t-t_{0}\right| \leq \lambda^{-1}\left(t_{0}\left(T-t_{0}\right)\right)^{\beta} . \tag{4.2.11}
\end{equation*}
$$

Then, for each $t \in \operatorname{supp}\left(\eta_{t_{0}}\right)$ we have

$$
\begin{equation*}
\left|\theta^{\alpha}(t)-\theta^{\alpha}\left(t_{0}\right)\right| \leq C\left(\theta^{\alpha}\left(t_{0}\right)\right)^{\prime}\left|t-t_{0}\right| \leq C \lambda^{-1}\left|t_{0}\left(T-t_{0}\right)\right|^{\beta-\alpha-1}, \tag{4.2.12}
\end{equation*}
$$

for some constant $C>0$. Then, the $L^{2}(0, T)$-norm of $H$ can be bounded as follows:

$$
\begin{equation*}
C_{1} s_{0} \int_{0}^{T}|H|^{2} d t \leq C\left|t_{0}\left(T-t_{0}\right)\right|^{2 \beta-\alpha-2} \lambda^{-2} s^{3} \int_{0}^{T}\left|Z_{2}(0, t)\right|^{2} d t \tag{4.2.13}
\end{equation*}
$$

Furthermore, notice that if we choose $\lambda_{*}$ such that

$$
\begin{equation*}
2 C\left|\frac{T^{2}}{4}\right|^{2 \beta+2 \alpha-2} \leq \lambda_{*}^{2} \tag{4.2.14}
\end{equation*}
$$

where $C$ is the constant appear in 4.2.13, it is evident that

$$
\begin{equation*}
C\left|t_{0}\left(T-t_{0}\right)\right|^{2 \beta-\alpha-2} \lambda^{-2} \leq \frac{1}{2}\left|t_{0}\left(T-t_{0}\right)\right|^{-3 \alpha}, \quad \forall \lambda \geq \lambda_{*} \tag{4.2.15}
\end{equation*}
$$

Thus, combining (4.2.13) with 4.2.15, we obtain

$$
\begin{equation*}
C_{1} s_{0} \int_{0}^{T}|H|^{2} d t \leq \frac{1}{2} s_{0}^{3} \int_{0}^{T}\left|Z_{2}(0, t)\right|^{2} d t . \tag{4.2.16}
\end{equation*}
$$

Analogously, straightforward computations show that

$$
\begin{equation*}
C_{1} s_{0} \int_{0}^{T}|J|^{2} d t \leq C\left|t_{0}\left(T-t_{0}\right)\right|^{2 \beta-\alpha-2} \lambda^{-2} s^{3} \int_{0}^{T}\left|Z_{2}\left(L_{2}, t\right)\right|^{2} d t \tag{4.2.17}
\end{equation*}
$$

Then, if we choose $\lambda_{*}$ such that

$$
\begin{equation*}
2 C\left|\frac{T^{2}}{4}\right|^{2 \beta+2 \alpha-2} \leq \lambda_{*}^{2}, \tag{4.2.18}
\end{equation*}
$$

where the constant $C$ is as the right-hand side of (4.2.17), we deduce that

$$
\begin{equation*}
C\left|t_{0}\left(T-t_{0}\right)\right|^{2 \beta-\alpha-2} \lambda^{-2} \leq \frac{1}{2}\left|t_{0}\left(T-t_{0}\right)\right|^{-3 \alpha}, \tag{4.2.19}
\end{equation*}
$$

for each $\lambda \geq \lambda_{*}$. Hence, we have the following upper-bound for $L^{2}(0, T)-$ norm of $J$ :

$$
\begin{equation*}
C_{1} s_{0} \int_{0}^{T}|J|^{2} d t \leq \frac{1}{2} s_{0}^{3} \int_{0}^{T}\left|Z_{2}(0, t)\right|^{2} d t \tag{4.2.20}
\end{equation*}
$$

Consequently, we add the inequalities (4.2.16) and (4.2.20), with $\lambda_{*}$ the maximum value which satisfies (4.2.14) and (4.2.18) and the proof of Proposition 4.4 is complete.

Applying the Proposition 1 into 4.2.9, we see that

$$
\begin{align*}
& \quad s_{0}^{3} \int_{0}^{T} \int_{\Omega_{1}}\left|Z_{1}\right|^{2} d x d t+s_{0}^{3} \int_{0}^{T} \int_{\Omega_{2}}\left|Z_{2}\right|^{2} d x d t+s_{0} \int_{0}^{T} \int_{\Omega_{1}}\left|\partial_{x} Z_{1}\right|^{2} d x d t \\
& \quad+s_{0} \int_{0}^{T} \int_{\Omega_{2}}\left|\partial_{x} Z_{2}\right|^{2} d x d t+s_{0}^{3} \int_{0}^{T}\left|Z_{2}(0, t)\right|^{2} d t+s_{0}^{3} \int_{0}^{T}\left|Z_{2}\left(L_{2}, t\right)\right|^{2} d t  \tag{4.2.21}\\
& \leq \\
& \leq C_{2} \int_{0}^{T} \int_{\Omega_{1}}\left|F_{1}\right|^{2} d x d t+C_{2} \int_{0}^{T} \int_{\Omega_{2}}\left|F_{2}\right|^{2} d x d t+C_{2} s_{0} \int_{0}^{T}\left|G_{N}\right|^{2} d t,
\end{align*}
$$

for some constant $C_{2}>0$ independent of $s$ and for every $\lambda \geq \lambda_{*}$.

## - Step 2: Residual terms of $F_{j}$

Now, we are interested in eliminate the residual terms of $F_{1}$ (and $F_{2}$ respectively) in (4.2.8) depending on $Z_{1}$ and $\partial_{x} Z_{1}$ (and $Z_{2}$ and $\partial_{x} Z_{2}$ respectively). First, we rewrite the terms of $F_{j}$ as follows:

$$
F_{j}=\tilde{f}_{j}+K_{j}
$$

where $\tilde{f}_{j}$ and $K_{j}$ is defined by

$$
\tilde{f}_{j}=e^{-s \varphi_{j}} \eta_{t_{0}} f_{j}+e^{-s \varphi_{j}} \partial_{t} \eta_{t_{0}} z_{j}
$$

and

$$
\begin{aligned}
K_{j}= & \alpha s \theta^{\alpha-1} \partial_{t} \theta \phi_{j} Z_{j}+\sigma_{j}^{2}\left(\theta^{2 \alpha}(t)-\theta^{2 \alpha}\left(t_{0}\right)\right)\left|\partial_{x} \phi_{j}\right|^{2} s^{2} Z_{j}+\sigma_{j}^{2}\left(\theta^{\alpha}(t)-\theta^{\alpha}\left(t_{0}\right)\right) \partial_{x}^{2} \phi_{j} s Z_{j} \\
& -\sigma_{j}^{2}\left(\theta^{\alpha}(t)-\theta^{\alpha}\left(t_{0}\right)\right) s \partial_{x} Z_{j}
\end{aligned}
$$

in $(0, T) \times \Omega_{j}$. We recall that $\eta_{t_{0}}$ is defined as follows:

$$
\eta_{t_{0}}(t)=\eta\left(\frac{\lambda\left(t-t_{0}\right)}{\left(t_{0}\left(T-t_{0}\right)\right)^{\beta}}\right),
$$

where the support of the function $\eta$ belongs in $]-1,1[$ and $\eta(0)=1$.
Proposition 4.5 Let $t_{0} \in(0, T)$ fixed and suppose that $\beta \geq 1$ in the definition of $\eta_{t_{0}}$ in (4.2.1). Then, there exist three positive constants $c_{1}, c_{2}$ and $\lambda_{*}$ independent of $t_{0}$ such that

$$
c_{1} \leq \frac{t(T-t)}{t_{0}\left(T-t_{0}\right)} \leq c_{2}, \quad \forall \lambda \geq \lambda_{*}
$$

Proof. By definition of $\eta_{t_{0}}$ we can write each term of its support as follows:

$$
t=t_{0}+a \lambda^{-1} t_{0}^{\beta}\left(T-t_{0}\right)^{\beta}, \quad \text { with }-1<a<1 .
$$

Then, we have

$$
\begin{align*}
\psi_{t_{0}}(t) & :=\frac{t(T-t)}{t_{0}\left(T-t_{0}\right)} \\
& =1+a \lambda^{-1} t_{0}^{\beta-1}\left(T-t_{0}\right)^{\beta-1}\left(T-2 t_{0}-a \lambda^{-1} t_{0}^{\beta}\left(T-t_{0}\right)^{\beta}\right) \tag{4.2.22}
\end{align*}
$$

for each $t$ lying in the support of $\psi_{t_{0}}$. We divide the proof into four cases:

- Case 1: Suppose that $0<t_{0} \leq T / 2$ and $a \geq 0$. In this case, notice that

$$
\begin{aligned}
\psi_{t_{0}}(t) & \geq 1-s^{2} \lambda^{-2}\left(t_{0}\left(T-t_{0}\right)\right)^{2 \beta-1} \\
& \geq 1-\lambda^{-2} \frac{T^{2(2 \beta-1)}}{4^{2 \beta-1}} .
\end{aligned}
$$

It is clear that if we choose $\lambda_{*}$ such that

$$
\sqrt{2} \frac{T^{2 \beta-1}}{2^{\beta-1}} \leq \lambda_{*},
$$

it is clear that

$$
1-\lambda^{-2} \frac{T^{2(2 \beta-1)}}{4^{2 \beta-1}} \geq \frac{1}{2}, \quad \lambda \geq \lambda_{*}
$$

Therefore, from the estimates above, we can assert that

$$
\psi_{t_{0}}(t) \geq \frac{1}{2}, \quad \forall \lambda \geq \lambda_{*}
$$

On the other hand,

$$
\begin{aligned}
\psi_{t_{0}}(t) & \leq 1+a \lambda^{-1} t_{0}^{\beta-1}\left(T-t_{0}\right)^{\beta-1} \\
& \leq 1+\lambda_{*}^{-1}\left(\frac{T^{2}}{4}\right)^{\beta-1} .
\end{aligned}
$$

Therefore, we obtain the following bounds for $\psi_{t_{0}}$ :

$$
\begin{equation*}
\frac{1}{2} \leq \psi_{t_{0}}(t) \leq 1+a \lambda_{*}\left(\frac{T}{4}\right)^{\beta-1} \tag{4.2.23}
\end{equation*}
$$

where $\lambda_{*}$ does not depend on $t_{0}$.

- Case 2: Suppose that $0<t_{0} \leq T / 2$ and $a<0$. For simplicity, we will use the temporary notation $b=-a$ that is, $b>0$. Then, 4.2.22 can be rewritten as follows:

$$
\psi_{t_{0}}(t)=1-b \lambda^{-1} t_{0}^{\beta-1}\left(T-t_{0}\right)^{\beta-1}\left(T-2 t_{0}+b \lambda^{-1} t_{0}^{\beta}\left(T-t_{0}\right)^{\beta}\right) .
$$

It is easy to see that

$$
\psi_{t_{0}}(t) \leq 1
$$

On the other hand,

$$
\psi_{t_{0}}(t) \geq 1-b\left(T-2 t_{0}\right) \frac{T^{2(\beta-1)}}{2^{2(\beta-1)}} \lambda^{-1}-b^{2} \frac{T^{2(2 \beta-1)}}{2^{2(2 \beta-1)}} \lambda^{-2}
$$

Straightforward computations show that if we choose $\lambda_{*}$ such that

$$
\lambda_{*} \geq\left(2+\sqrt{\frac{3}{2}}\right) \frac{T^{2 \beta-1}}{2^{\beta-2}}
$$

we can assert that

$$
1-b\left(T-2 t_{0}\right) \frac{T^{2(\beta-1)}}{2^{2(\beta-1)}} \lambda^{-1}-b^{2} \frac{T^{2(2 \beta-1)}}{2^{2(2 \beta-1)}} \lambda^{-2} \geq \frac{1}{2}, \quad \forall \lambda \geq \lambda_{*} .
$$

Hence, we have

$$
\begin{equation*}
\frac{1}{2} \leq \psi_{t_{0}}(t) \leq 1, \quad \forall \lambda \geq \lambda_{*} \tag{4.2.24}
\end{equation*}
$$

- Case 3: Suppose that $T / 2<t_{0}<T$ and $a \geq 0$. Then, it is easy to see that

$$
\psi_{t_{0}}(t) \leq 1 .
$$

On the other hand, a simple computation shows that if we choose $\lambda_{*}$ fullfilling

$$
\lambda_{*} \geq \frac{5}{2} \frac{T^{2 \beta-1}}{2^{2 \beta-1}}
$$

we obtain

$$
\psi_{t_{0}}(t) \geq 1-\left(2 t_{0}-T\right) a \lambda^{-1}\left|\frac{T^{2}}{4}\right|^{\beta-1}-a^{2} \lambda^{-2}\left|\frac{T^{2}}{4}\right|^{2 \beta-1} \geq \frac{1}{2}, \quad \lambda \geq \lambda_{*}
$$

Thus, we have

$$
\begin{equation*}
\frac{1}{2} \leq \psi_{t_{0}}(t) \leq 1, \quad \lambda \geq \lambda_{*} \tag{4.2.25}
\end{equation*}
$$

- Case 4: Suppose that $T / 2<t_{0}<T$ and $a<0$. Once again, we set $b=-a$. In this case, we note that

$$
\begin{aligned}
\psi_{t_{0}}(t) & =1-b \lambda^{-1} t_{0}^{\beta-1}\left(T-t_{0}\right)^{\beta-1}\left(T-2 t_{0}+b \lambda^{-1} t_{0}^{\beta}\left(T-t_{0}\right)^{\beta}\right) \\
& \leq 1+b\left(2 t_{0}-T\right) \lambda^{-1} t_{0}^{\beta-1}\left(T-t_{0}\right)^{\beta-1} \\
& \leq 1+2^{2-2 \beta} b \lambda_{*}^{-1} T^{2 \beta-1} .
\end{aligned}
$$

Furthermore, it is clear that if we choose $\lambda_{*}$ such that

$$
\lambda \geq 2^{2 \beta-1} \sqrt{2} b T^{2 \beta-1},
$$

we obtain

$$
1-b^{2} \lambda^{-2} 2^{2-4 \beta} T^{4 \beta-2} \geq \frac{1}{2}, \quad \lambda \geq \lambda_{*} .
$$

Therefore,

$$
\begin{equation*}
\frac{1}{2} \leq \psi_{t_{0}}(t) \leq 1+b 2^{2-2 \beta} \lambda_{*}^{-1} T^{2 \beta-1}, \quad \forall \lambda \geq \lambda_{*} \tag{4.2.26}
\end{equation*}
$$

From (4.2.23), 4.2.24, (4.2.25) and (4.2.26), we conclude the proof of the Proposition 4.5

Proposition 4.6 We Suppose the same hypotheses of Lemma 1. Additionally, let $\alpha$ and $\beta$ such that $\alpha \geq 2$ and $2 \beta-\alpha \geq 2$. Then, for each $j=1,2$, there exist two positive constants $\lambda_{*}$ and $s_{*}$ such that for every $\lambda \geq \lambda_{*}$ and $s \geq s_{*}$ we have

$$
C_{2} \int_{0}^{T} \int_{\Omega_{j}}\left|K_{j}\right|^{2} d x d t \leq \frac{1}{2} s_{0}^{3} \int_{0}^{T} \int_{\Omega_{j}}\left|Z_{j}\right|^{2} d x d t+\frac{1}{2} s_{0} \int_{0}^{T} \int_{\Omega_{j}}\left|\partial_{x} Z_{j}\right|^{2} d x d t, \quad j=1,2,
$$

where $C_{2}$ is the constant previously defined in 4.2.21).
Proof. By definition of $K_{j}$, we have

$$
\begin{align*}
& C_{2} \int_{0}^{T} \int_{\Omega_{j}}\left|K_{j}\right|^{2} d x d t \\
& \leq \\
& \leq\left(t_{0}\left(T-t_{0}\right)\right)^{-2(\alpha+1)} s^{2} \int_{0}^{T} \int_{\Omega_{j}}\left|Z_{j}\right|^{2} d x d t+C s^{4} \int_{0}^{T} \int_{\Omega_{j}}\left|\theta^{2 \alpha}(t)-\theta^{2 \alpha}\left(t_{0}\right)\right|^{2}\left|Z_{j}\right|^{2} d x d t  \tag{4.2.27}\\
& \quad+C s^{2} \int_{0}^{T} \int_{\Omega_{j}}\left|\theta^{\alpha}(t)-\theta^{\alpha}\left(t_{0}\right)\right|^{2}\left|Z_{j}\right|^{2} d x d t+C s^{2} \int_{0}^{T} \int_{\Omega_{j}}\left|\theta^{\alpha}(t)-\theta^{\alpha}\left(t_{0}\right)\right|^{2}\left|\partial_{x} Z_{j}\right|^{2} d x d t .
\end{align*}
$$

Let us estimate each term of the right-hand side of the inequality above. First, notice that if we choose $s_{*}$ such that

$$
\begin{equation*}
4 C_{2}\left|\frac{T^{2}}{4}\right|^{\alpha-2} \leq s_{*} \tag{4.2.28}
\end{equation*}
$$

the following inequality holds:

$$
\begin{equation*}
C_{2}\left|t_{0}\left(T-t_{0}\right)\right|^{-2(\alpha+1)} \leq \frac{1}{4}\left|t_{0}\left(T-t_{0}\right)\right|^{-3 \alpha} s, \quad s \geq s_{*} . \tag{4.2.29}
\end{equation*}
$$

This implies,

$$
\begin{equation*}
C_{2}\left|t_{0}\left(T-t_{0}\right)\right|^{-2(\alpha+1)} s^{2} \int_{0}^{T} \int_{\Omega_{j}}\left|Z_{j}\right|^{2} d x d t \leq \frac{1}{4} s_{0}^{3} \int_{0}^{T} \int_{\Omega_{j}}\left|Z_{j}\right|^{2} d x d t, \quad s \geq s_{*} \tag{4.2.30}
\end{equation*}
$$

On the other hand, we note that for each $t \in \operatorname{supp}\left(\eta_{t_{0}}\right)$, we have

$$
\begin{equation*}
\left|\theta^{2 \alpha}(t)-\theta^{2 \alpha}\left(t_{0}\right)\right| \leq C \theta^{2 \alpha-1}\left(t_{0}\right) \partial_{t} \theta\left(t_{0}\right)\left|t-t_{0}\right| \leq C\left|t_{0}\left(T-t_{0}\right)\right|^{\beta-2 \alpha-1} \lambda^{-1} \tag{4.2.31}
\end{equation*}
$$

for some constant $C>0$. Then, notice that for each $t \in \operatorname{supp}\left(\eta_{t_{0}}\right)$, we can assert that

$$
\begin{align*}
& C_{2} s^{4} \int_{0}^{T} \int_{\Omega_{j}}\left|\theta^{2 \alpha}(t)-\theta^{2 \alpha}\left(t_{0}\right)\right|^{2}\left|Z_{j}\right|^{2} d x d t \\
\leq & C\left|t_{0}\left(T-t_{0}\right)\right|^{2 \beta-4 \alpha-2} \lambda^{-2} s^{4} \int_{0}^{T} \int_{\Omega_{j}}\left|Z_{j}\right|^{2} d x d t \tag{4.2.32}
\end{align*}
$$

If we choose the parameters $s_{*}$ and $\lambda_{*}$ such that

$$
\begin{equation*}
4 C\left|\frac{T^{2}}{4}\right|^{2 \beta-\alpha-2} \leq \lambda_{*}^{2} s_{*}^{-1} \tag{4.2.33}
\end{equation*}
$$

where the constant $C>0$ is given in 4.2.32, it is easy to check that

$$
\begin{equation*}
C\left|t_{0}\left(T-t_{0}\right)\right|^{2 \beta-4 \alpha-2} \lambda^{-2} s^{4} \leq\left.\frac{1}{4} s^{3}\left(t_{0}\right)\left(T-t_{0}\right)\right|^{-3 \alpha}, \quad \lambda \geq \lambda_{*}, s \geq s_{*} \tag{4.2.34}
\end{equation*}
$$

Thus, combining (4.2.32) and (4.2.34), the first term of the right-hand side of 4.2.27) can be bounded as follows:

$$
\begin{equation*}
C_{2} s^{4} \int_{0}^{T} \int_{\Omega_{j}}\left|\theta^{2 \alpha}(t)-\theta^{2 \alpha}\left(t_{0}\right)\right|^{2}\left|Z_{j}\right|^{2} d x d t \leq \frac{1}{4} s_{0}^{3} \int_{0}^{T} \int_{\Omega_{j}}\left|Z_{j}\right|^{2} d x d t \tag{4.2.35}
\end{equation*}
$$

Now we deal with the second term of the right-hand side of 4.2.27). Using the estimate of $\theta^{\alpha}$ in (4.2.12), we see that

$$
\begin{equation*}
C_{2} s^{2} \int_{0}^{T} \int_{\Omega_{j}}\left|\theta^{\alpha}(t)-\theta^{\alpha}\left(t_{0}\right)\right|^{2}\left|Z_{j}\right|^{2} d x d t \leq C\left|t_{0}\left(T-t_{0}\right)\right|^{2 \beta-2 \alpha-2} \lambda^{-2} s^{2} \int_{0}^{T} \int_{\Omega_{j}}\left|Z_{j}\right|^{2} d x d t \tag{4.2.36}
\end{equation*}
$$

If we choose the parameters $\lambda_{*}$ and $s_{*}$ such that

$$
\begin{equation*}
4 C\left|\frac{T^{2}}{4}\right|^{2 \beta+\alpha-2} \leq \lambda_{*}^{2} s_{*}, \tag{4.2.37}
\end{equation*}
$$

where $C$ is the constant appeared in 4.2.36, it is easy to seen that

$$
\begin{equation*}
C\left|t_{0}\left(T-t_{0}\right)\right|^{2 \beta-\alpha-2} \leq \frac{1}{4}\left|t_{0}\left(T-t_{0}\right)\right|^{-3 \alpha} s^{3} \int_{0}^{T} \int_{\Omega_{j}}\left|Z_{j}\right|^{2} d x d t, \quad \lambda \geq \lambda_{*}, s \geq s_{*} . \tag{4.2.38}
\end{equation*}
$$

Substituting 4.2.38) into 4.2.36 yields

$$
\begin{equation*}
C s^{2} \int_{0}^{T} \int_{\Omega_{j}}\left|\theta^{\alpha}(t)-\theta^{\alpha}\left(t_{0}\right)\right|^{2}\left|Z_{j}\right|^{2} d x d t \leq \frac{1}{4} s_{0}^{3} \int_{0}^{T} \int_{\Omega_{j}}\left|Z_{j}\right|^{2} d x d t \tag{4.2.39}
\end{equation*}
$$

Once again, applying the estimate 4.2.12, the third term of the right-hand side of (4.2.27) can be bounded as follows:
$C s^{2} \int_{0}^{T} \int_{\Omega_{j}}\left|\theta^{\alpha}(t)-\theta^{\alpha}\left(t_{0}\right)\right|^{2}\left|\partial_{x} Z_{j}\right|^{2} d x d t \leq C\left|t_{0}\left(T-t_{0}\right)\right|^{2 \beta-2 \alpha-2} \lambda^{-2} s^{2} \int_{0}^{T} \int_{\Omega_{j}}\left|\partial_{x} Z_{j}\right|^{2} d x d t$.

It is inmediate that if we choose $\lambda_{*}$ and $s_{*}$ satisfying

$$
\begin{equation*}
2 C\left|\frac{T^{2}}{4}\right|^{2 \beta-\alpha-2} \leq \lambda_{*}^{2} s_{*}^{-1} \tag{4.2.41}
\end{equation*}
$$

it follows that

$$
\begin{equation*}
C\left|t_{0}\left(T-t_{0}\right)\right|^{2 \beta-2 \alpha-2} \leq \frac{1}{2}\left|t_{0}\left(T-t_{0}\right)\right|^{-\alpha}, \quad \forall \lambda \geq \lambda_{*}, \forall s \geq s_{*} . \tag{4.2.42}
\end{equation*}
$$

Substituting (4.2.42) into (4.2.40), we get

$$
\begin{equation*}
C s^{2} \int_{0}^{T} \int_{\Omega_{j}}\left|\theta^{\alpha}(t)-\theta^{\alpha}\left(t_{0}\right)\right|^{2}\left|\partial_{x} Z_{j}\right|^{2} d x d t \leq \frac{1}{2} s_{0} \int_{0}^{T} \int_{\Omega_{j}}\left|\partial_{x} Z_{j}\right|^{2} d x d t \tag{4.2.43}
\end{equation*}
$$

Finally, we add the inequalities 4.2.35, 4.2.39) and 4.2.43), and the proof of the Proposition 4.6 is complete.

Thus, applying the Proposition 4.4 and 4.6 into (4.2.9), we conclude that

$$
\begin{align*}
& s_{0}^{3} \int_{0}^{T} \int_{\Omega_{1}}\left|Z_{1}\right|^{2} d x d t+s_{0}^{3} \int_{0}^{T} \int_{\Omega_{2}}\left|Z_{2}\right|^{2} d x d t+s_{0} \int_{0}^{T} \int_{\Omega_{1}}\left|\partial_{x} Z_{1}\right|^{2} d x d t \\
& \quad+s_{0} \int_{0}^{T} \int_{\Omega_{2}}\left|\partial_{x} Z_{2}\right|^{2} d x d t  \tag{4.2.44}\\
& \leq C \int_{0}^{T} \int_{\Omega_{1}}\left|\tilde{f}_{1}\right|^{2} d x d t+C \int_{0}^{T} \int_{\Omega_{2}}\left|\tilde{f}_{2}\right|^{2} d x d t+C s_{0} \int_{0}^{T}\left|G_{N}\right|^{2} d t,
\end{align*}
$$

for some constant $C>0$ independent of $\lambda$ and $s$.

- Step 3: Last residual terms depending on $z_{j}, j=1,2$

Our next task is to deal with the residual terms depending on $z_{j}$ from $\tilde{f}_{j}, j=1,2$. In order to do that, we introduce the following functions:

$$
\omega_{0}(t)=\int_{0}^{T}\left|\eta_{t_{0}}(t)\right|^{2} d t_{0}, \quad \omega_{1}(t)=\int_{0}^{T}\left|\partial_{t} \eta_{t_{0}}(t)\right| d t_{0} .
$$

In the reminder of this step, we need the following technical result:
Proposition 4.7 There exist positive constants $c_{3}, c_{4}, c_{5}, c_{6}$ and $\lambda_{*}$ such that, for all $t \in$ $\operatorname{supp}\left(\eta_{t_{0}}\right)$,

$$
\begin{equation*}
c_{3} \lambda^{-1} \theta^{-\beta}(t) \leq \omega_{0}(t) \leq c_{4} \lambda^{-1} \theta^{-\beta}(t), \quad \forall \lambda \geq \lambda_{*}, \tag{4.2.45}
\end{equation*}
$$

and

$$
\begin{equation*}
c_{5} \lambda \theta^{\beta}(t) \leq \omega_{1}(t) \leq c_{6} \lambda \theta^{\beta}(t), \quad \forall \lambda \geq \lambda_{*} . \tag{4.2.46}
\end{equation*}
$$

Proof. By definition,

$$
\begin{align*}
\omega_{0}(t) & =\int_{0}^{T}\left|\eta\left(\frac{\lambda\left(t-t_{0}\right)}{\left(t_{0}\left(T-t_{0}\right)\right)^{\beta}}\right)\right|^{2} d t_{0} \\
& =\int_{0}^{T}|\eta(\tilde{t})|^{2}\left|\frac{d \tilde{t}}{d t_{0}}\right|^{-1} d t_{0}, \tag{4.2.47}
\end{align*}
$$

where $\tilde{t}$ is defined by

$$
\tilde{t}=\frac{\lambda\left(t-t_{0}\right)}{t_{0}^{\beta}\left(T-t_{0}\right)^{\beta}} .
$$

A direct compute show that

$$
\frac{d \tilde{t}}{d t_{0}}=-\lambda\left(\frac{1}{\left(t_{0}\left(T-t_{0}\right)\right)^{\beta}}+\frac{\beta\left(T-2 t_{0}\right)\left(t-t_{0}\right)}{\left(t_{0}\left(T-t_{0}\right)\right)^{\beta+1}}\right) .
$$

Since every $t \in \operatorname{supp}\left(\eta_{t_{0}}\right)$ satisfy

$$
\left|t-t_{0}\right| \leq \lambda^{-1}\left(t_{0}\left(T-t_{0}\right)\right)^{\beta},
$$

we can assert that

$$
\frac{\lambda}{\left(t_{0}\left(T-t_{0}\right)\right)^{\beta}}-\beta T \leq\left|\frac{d \tilde{t}}{d t_{0}}\right| \leq \frac{\lambda}{\left(t_{0}\left(T-t_{0}\right)\right)^{\beta}}+\beta T .
$$

Taking $\lambda_{*}$ large enough and using the Proposition 4.5, we have for each $t \in \operatorname{supp}\left(\eta_{t_{0}}\right)$

$$
c_{1} \lambda\left(t\left(T-t_{0}\right)\right)^{-\beta} \leq\left|\frac{d \tilde{t}}{d t_{0}}\right| \leq c_{2} \lambda(t(T-t))^{-\beta}, \quad \lambda \geq \lambda_{*} .
$$

Consequently, we have for all $t \in \operatorname{supp}\left(\eta_{t_{0}}\right)$

$$
c_{3} \lambda^{-1}(t(T-t))^{\beta} \leq \omega_{0}(t) \leq c_{4} \lambda^{-1}(t(T-t))^{\beta}, \quad \lambda \geq \lambda_{*}, t \in \operatorname{supp}\left(\eta_{t_{0}}\right) .
$$

This completes the proof of the first inequality. To proof the other one, notice that

$$
\partial_{t} \eta_{t_{0}}(t)=\frac{\lambda}{\left(t_{0}\left(T-t_{0}\right)\right)^{\beta}} \partial_{\tilde{t}} \eta(\tilde{t}),
$$

where $\tilde{t}$ is defined as before. Then,

$$
\omega_{1}(t)=\lambda^{2} \int_{0}^{T} t_{0}^{-2 \beta}(T-t)^{-2 \beta}\left|\partial_{\tilde{t}} \eta(\tilde{t})\right|^{2}\left|\frac{d \tilde{t}}{d t_{0}}\right|^{-1} d t_{0}
$$

The rest of the proof runs as before. Thus, the proof of the Proposition 4.7 is complete.

Now, we integrate in $t_{0}$ on $(0, T)$ in (4.2.44) and by definition of $Z_{1}, Z_{2}$ and $G_{N}$, we obtain

$$
\begin{aligned}
& \int_{0}^{T} \int_{0}^{T} \int_{\Omega_{1}} s_{0}^{3} e^{-2 s \varphi_{1}}\left|\eta_{t_{0}} z_{1}\right|^{2} d x d t d t_{0}+\int_{0}^{T} \int_{0}^{T} \int_{\Omega_{2}} s_{0}^{3} e^{-2 s \varphi_{2}}\left|\eta_{t_{0}} z_{2}\right|^{2} d x d t d t_{0} \\
+ & \int_{0}^{T} \int_{0}^{T} \int_{\Omega_{1}} s_{0} e^{-2 s \varphi_{1}}\left|\eta_{t_{0}} \partial_{x} z_{1}\right|^{2} d x d t d t_{0}+\int_{0}^{T} \int_{0}^{T} \int_{\Omega_{2}} s_{0} e^{-2 s \varphi_{2}}\left|\eta_{t_{0}} \partial_{x} z_{2}\right|^{2} d x d t d t_{0} \\
+ & \int_{0}^{T} \int_{0}^{T} s_{0}^{3} e^{-2 s \varphi(0, t)}\left|\eta_{t_{0}} z_{1}(0, t)\right|^{2} d t d t_{0}+\int_{0}^{T} \int_{0}^{T} s_{0}^{3} e^{-2 s \varphi\left(L_{2}, t\right)}\left|\eta_{t_{0}} z_{2}\left(L_{2}, t\right)\right|^{2} d t d t_{0} \\
\leq & C \int_{0}^{T} \int_{0}^{T} \int_{\Omega_{1}} e^{-2 s \varphi_{1}}\left|\tilde{f}_{1}\right|^{2} d x d t d t_{0}+C \int_{0}^{T} \int_{0}^{T} \int_{\Omega_{2}} e^{-2 s \varphi_{2}}\left|\tilde{f}_{2}\right|^{2} d x d t d t_{0} \\
& +C \int_{0}^{T} \int_{0}^{T} s_{0} e^{-2 s \varphi_{1}\left(-L_{1}, t\right)}\left|\eta_{t_{0}} z_{1}\left(-L_{1}, t\right)\right|^{2} d t d t_{0}
\end{aligned}
$$

where

$$
\begin{aligned}
& \int_{0}^{T} \int_{0}^{T} \int_{\Omega_{j}}\left|\tilde{f}_{j}\right|^{2} d x d t \\
\leq & C \int_{0}^{T} \int_{0}^{T} \int_{\Omega_{j}} e^{-2 s \varphi_{j}}\left|\eta_{t_{0}}\right|^{2}\left|f_{j}\right|^{2} d x d t d t_{0}+C \int_{0}^{T} \int_{0}^{T} \int_{\Omega_{2}} e^{-2 s \varphi_{j}}\left|\partial_{t} \eta_{t_{0}}\right|^{2}\left|z_{j}\right|^{2} d x d t d t_{0}
\end{aligned}
$$

Applying Fubini's Theorem and the results of Proposition 4.7, we deduce that

$$
\begin{aligned}
& \lambda^{-1} s^{3} \int_{0}^{T} \int_{\Omega_{1}} e^{-2 s \varphi_{1}} \theta^{3 \alpha-\beta}\left|z_{1}\right|^{2} d x d t+\lambda^{-1} s^{3} \int_{0}^{T} \int_{\Omega_{2}} e^{-2 s \varphi_{2}} \theta^{3 \alpha-\beta}\left|z_{2}\right|^{2} d x d t \\
& \quad+\lambda^{-1} s \int_{0}^{T} \int_{\Omega_{1}} e^{-2 s \varphi_{1}} \theta^{\alpha-\beta}\left|\partial_{x} z_{1}\right|^{2} d x d t+\lambda^{-1} s \int_{0}^{T} \int_{\Omega_{2}} e^{-2 s \varphi_{2}} \theta^{\alpha-\beta}\left|\partial_{x} z_{2}\right|^{2} d x d t \\
& \quad+\lambda^{-1} s^{3} \int_{0}^{T} e^{-2 s \varphi_{1}(t, 0)} \theta^{3 \alpha-\beta}\left|z_{1}(t, 0)\right|^{2} d t+\lambda^{-1} s^{3} \int_{0}^{T} e^{-2 s \varphi_{2}\left(t, L_{2}\right)} \theta^{3 \alpha-\beta}\left|z_{2}\left(t, L_{2}\right)\right|^{2} d t \\
& \leq C \lambda^{-1} \int_{0}^{T} \int_{\Omega_{1}} e^{-2 s \varphi_{1}} \theta^{-\beta}\left|f_{1}\right|^{2} d x d t+\lambda^{-1} \int_{0}^{T} \int_{\Omega_{2}} e^{-2 s \varphi_{2}} \theta^{-\beta}\left|f_{2}\right|^{2} d x d t \\
& \quad+C \lambda \int_{0}^{T} \int_{\Omega_{1}} e^{-2 s \varphi_{1}} \theta^{\beta}\left|z_{1}\right|^{2} d x d t+C \lambda \int_{0}^{T} \int_{\Omega_{2}} e^{-2 s \varphi_{2}} \theta^{\beta}\left|z_{2}\right|^{2} d x d t \\
& \quad+C \lambda^{-1} s \int_{0}^{T} e^{-2 s \varphi_{1}\left(-L_{1}\right)} \theta^{\alpha-\beta}\left|g_{N}\right|^{2} d t,
\end{aligned}
$$

or equivalently

$$
\begin{aligned}
& s^{3} \int_{0}^{T} \int_{\Omega_{1}} e^{-2 s \varphi_{1}} \theta^{3 \alpha-\beta}\left|z_{1}\right|^{2} d x d t+s^{3} \int_{0}^{T} \int_{\Omega_{2}} e^{-2 s \varphi_{2}} \theta^{3 \alpha-\beta}\left|z_{2}\right|^{2} d x d t \\
& \quad+s \int_{0}^{T} \int_{\Omega_{1}} e^{-2 s \varphi_{1}} \theta^{\alpha-\beta}\left|\partial_{x} z_{1}\right|^{2} d x d t+s \int_{0}^{T} \int_{\Omega_{2}} e^{-2 s \varphi_{2}} \theta^{\alpha-\beta}\left|\partial_{x} z_{2}\right|^{2} d x d t \\
&+s^{3} \int_{0}^{T} e^{-2 s \varphi_{1}(0, t)} \theta^{3 \alpha-\beta}\left|z_{1}(0, t)\right|^{2} d t+s^{3} \int_{0}^{T} e^{-2 s \varphi_{2}\left(L_{2}, t\right)} \theta^{3 \alpha-\beta}\left|z_{2}\left(L_{2}, t\right)\right|^{2} d t \\
& \leq C_{3} \int_{0}^{T} \int_{\Omega_{1}} e^{-2 s \varphi_{1}} \theta^{-\beta}\left|f_{1}\right|^{2} d x d t+C_{3} \int_{0}^{T} \int_{\Omega_{2}} e^{-2 s \varphi_{2}} \theta^{-\beta}\left|f_{2}\right|^{2} d x d t \\
&+C_{3} \lambda^{2} \int_{0}^{T} \int_{\Omega_{1}} e^{-2 s \varphi_{1}} \theta^{\beta}\left|z_{1}\right|^{2} d x d t+C_{3} \lambda^{2} \int_{0}^{T} \int_{\Omega_{2}} e^{-2 s \varphi_{2}} \theta^{\beta}\left|z_{2}\right|^{2} d x d t \\
&+C_{3} s \int_{0}^{T} e^{-2 s \varphi_{1}\left(-L_{1}\right)} \theta^{\alpha-\beta}\left|g_{N}\right|^{2} d t,
\end{aligned}
$$

To eliminate the terms of $z_{1}$ and $z_{2}$ on the right-hand side of the inequality above, we will use the following result:

Proposition 4.8 Suppose that $\alpha \geq 2$ and $3 \alpha \geq 2 \beta$. Then, there exists two positive constants $s_{*}$ and $\lambda_{*}$ such that for all $s \geq s_{*}$ and $\lambda \geq \lambda_{*}$, we have

$$
\begin{equation*}
C_{3} \lambda^{2} \int_{0}^{T} \int_{\Omega_{j}} e^{-2 s \varphi_{j}} \theta^{\beta}\left|z_{j}\right|^{2} d x d t \leq \frac{1}{2} s^{3} \int_{0}^{T} \int_{\Omega_{j}} e^{-2 s \varphi_{j}} \theta^{3 \alpha-\beta}\left|z_{j}\right|^{2} d x d t \tag{4.2.48}
\end{equation*}
$$

for each $j=1,2$.
Finally, using the estimates (4.2.48) we obtain

$$
\begin{aligned}
& s^{3} \int_{0}^{T} \int_{\Omega_{1}} e^{-2 s \varphi_{1}} \theta^{3 \alpha-\beta}\left|z_{1}\right|^{2} d x d t+s^{3} \int_{0}^{T} \int_{\Omega_{2}} e^{-2 s \varphi_{1}} \theta^{3 \alpha-\beta}\left|z_{2}\right|^{2} d x d t \\
& \quad+s \int_{0}^{T} \int_{\Omega_{1}} e^{-2 s \varphi_{1}} \theta^{\alpha-\beta}\left|\partial_{x} z_{1}\right|^{2} d x d t+s \int_{0}^{T} \int_{\Omega_{2}} e^{-2 s \varphi_{2}} \theta^{\alpha-\beta}\left|\partial_{x} z_{2}\right|^{2} d x d t \\
& \quad+s^{3} \int_{0}^{T} e^{-2 s \varphi_{1}(0, t)} \theta^{3 \alpha-\beta}\left|z_{1}(0, t)\right|^{2} d t+s^{3} \int_{0}^{T} e^{-2 s \varphi_{2}\left(L_{2}, t\right)} \theta^{3 \alpha-\beta}\left|z_{2}\left(L_{2}, t\right)\right|^{2} d t \\
& \leq C \int_{0}^{T} \int_{\Omega_{1}} e^{-2 s \varphi_{1}} \theta^{-\beta}\left|f_{1}\right|^{2} d x d t+C \int_{0}^{T} \int_{\Omega_{2}} e^{-2 s \varphi_{2}} \theta^{-\beta}\left|f_{2}\right|^{2} d x d t \\
& \quad+C s \int_{0}^{T} e^{-2 s \varphi_{1}\left(L_{1}, t\right)} \theta^{\alpha-\beta}\left|z_{1}\left(-L_{1}, t\right)\right|^{2} d t .
\end{aligned}
$$

This completes the proof of Theorem 4.1.

### 4.3 Proof of Lemma 4.3

In this section, we devote to prove the Lemma 4.3. The proof of this one falls naturally into two parts:

- The first one concerns in the $L^{2}$ estimates for the global terms of $Z_{1}$ and $Z_{2}$, the local term of $Z_{1}$ at the interface and the local term of $Z_{2}$ at the right-hand side of the domain:

$$
\begin{align*}
& s_{0}^{3} \int_{0}^{T} \int_{\Omega_{1}}\left|Z_{1}\right|^{2} d x d t+s_{0}^{3} \int_{0}^{T} \int_{\Omega_{2}}\left|Z_{2}\right|^{2} d x d t+s_{0}^{3} \int_{0}^{T}\left|Z_{1}(0, t)\right|^{2} d t+s_{0}^{3} \int_{0}^{T}\left|Z_{2}\left(L_{2}, t\right)\right|^{2} d t \\
\leq & C \int_{0}^{T} \int_{\Omega_{1}}\left|F_{1}\right|^{2} d x d t+C \int_{0}^{T} \int_{\Omega_{2}}\left|F_{2}\right|^{2} d x d t+C s_{0} \int_{0}^{T}\left|G_{N}\right|^{2} d t \\
& +C s_{0} \int_{0}^{T}|H|^{2} d t+C s_{0} \int_{0}^{T}|J|^{2} d t, \tag{4.3.1}
\end{align*}
$$

where $C$ is a positive constant independent of $t_{0}$ and $s \geq s^{*}$.

- The second part consists in prove a similar estimate for the $L^{2}\left((0, T) \times \Omega_{j}\right)$-norm of the spatial derivatives of $Z_{j}$ for $j=1,2$ :

$$
\begin{align*}
& s_{0} \int_{0}^{T} \int_{\Omega_{1}}\left|\partial_{x} Z_{1}\right|^{2} d x d t+s_{0} \int_{0}^{T} \int_{\Omega_{2}}\left|\partial_{x} Z_{2}\right|^{2} d x d t \\
\leq & C \int_{0}^{T} \int_{\Omega_{1}}\left|F_{1}\right|^{2} d x d t+C \int_{0}^{T} \int_{\Omega_{2}}\left|F_{2}\right|^{2} d x d t+C s_{0} \int_{0}^{T}\left|G_{N}\right|^{2} d t+C s_{0} \int_{0}^{T}|H|^{2} d t  \tag{4.3.2}\\
& +C s_{0} \int_{0}^{T}|J|^{2} d t
\end{align*}
$$

Clearly, if we add the estimates (4.3.1) and 4.3.2 the proof of the Lemma 4.3 is complete.
Then, we start proving the inequality 4.3.1. Before going further, let us bring a brief orientation of the proof. First, we will use the Fourier transform in time in order to get good estimates in frequency domain. However, it is not evident that these estimates are uniform on the frequency parameter in the Fourier domain. In order to prove that, we divide in some frequency ranges to analyze our estimates saying, Low, intermediate and high frequencies. That is, this phenomenon depends where the information comes from it.

For a function $h \in L^{2}(\mathbb{R})$, we introduce the partial Fourier transform in time defined by

$$
\begin{equation*}
\hat{h}(\tau)=\mathcal{F}(h)(\tau)=\int_{\mathbb{R}} h(t) e^{-i t \tau} d t \tag{4.3.3}
\end{equation*}
$$

where $i$ is the imaginary unit. Now, we extend the variables $Z_{1}$ and $Z_{2}$ of system 4.2.5) by zero, and we do the same for the functions $H, J$ and $F_{1}$ and $F_{2}$.

Thus, applying the Fourier transform (4.3.3) to the system 4.2.5), we obtain

$$
\begin{cases}\left(\sigma_{1} \partial_{x}+\sigma_{1} \partial_{x} \phi_{1} s_{0}+\gamma\right)\left(\sigma_{1} \partial_{x}+\sigma_{1} \partial_{x} \phi_{1} s_{0}-\gamma\right) \hat{Z}_{1}=-\hat{F}_{1}, & \forall(x, \tau) \in \Omega_{1} \times \mathbb{R},  \tag{4.3.4}\\ \left(\sigma_{2} \partial_{x}+\sigma_{2} \partial_{x} \phi_{2} s_{0}+\gamma\right)\left(\sigma_{2} \partial_{x}+\sigma_{2} \partial_{x} \phi_{2} s_{0}-\gamma\right) \hat{Z}_{2}=-\hat{F}_{2}, & \forall(x, \tau) \in \Omega_{2} \times \mathbb{R}, \\ \hat{Z}_{2}(0, \tau)=\hat{Z}_{1}(0, \tau), & \forall t \in \mathbb{R}, \\ \sigma_{2}^{2} \partial_{x} \hat{Z}_{2}(0, \tau)=\sigma_{1}^{2} \partial_{x} \hat{Z}_{1}(0, \tau)+\left(m_{1}-M_{2} \sigma_{2}\right) s_{0} \hat{Z}_{1}(0, \tau)+\hat{H}(\tau), & \forall \tau \in \mathbb{R}, \\ \hat{Z}_{1}\left(-L_{1}, \tau\right)=0, & \forall \tau \in \mathbb{R}, \\ \partial_{x} \hat{Z}_{2}\left(L_{2}, \tau\right)=-m_{2} \sigma_{2}^{-1} s_{0} \hat{Z}_{2}\left(L_{2}, \tau\right)+\hat{J}(\tau), & \forall \tau \in \mathbb{R}, \\ \partial_{x} \hat{Z}_{1}\left(-L_{1}, \tau\right)=\hat{G}_{N}(\tau), & \forall \tau \in \mathbb{R} .\end{cases}
$$

Here, $\gamma=\gamma(\tau)$ is defined by

$$
\gamma= \begin{cases}\sqrt{|\tau|} e^{i \frac{\pi}{4}}, & \text { if } \tau \geq 0  \tag{4.3.5}\\ \sqrt{|\tau|} e^{-i \frac{\pi}{4}}, & \text { if } \tau<0\end{cases}
$$

It is clear that $\Re(\gamma(\tau))=\sqrt{|\tau| / 2} \geq 0$ and $\gamma^{2}=i \tau$, for all $\tau \in \mathbb{R}$. Now, inspired in the structure of the operator in $(4.3 .4)_{1}$ and $(4.3 .4)_{2}$, we define the auxiliar variables:

$$
\hat{W}_{j}=\left(\sigma_{j} \partial_{x}+\sigma_{j} \partial_{x} \phi_{j} s_{0}-\gamma\right) \hat{Z}_{j}, \quad j=1,2 .
$$

### 4.3.1 First estimates

Proposition 4.9 Let $\hat{W}_{1}$ be the solution of the following system:

$$
\begin{cases}\left(\sigma_{1} \partial_{x}+\sigma_{1} \partial_{x} \phi_{1} s_{0}+\gamma\right) \hat{W}_{1}(x, \tau)=-\hat{F}_{1}(x, \tau), & \forall(x, \tau) \in \Omega_{1} \times \mathbb{R}  \tag{4.3.6}\\ \hat{W}_{1}\left(-L_{1}, \tau\right)=\sigma_{1} \hat{G}_{N}(\tau), & \forall \tau \in \mathbb{R}\end{cases}
$$

Then, each solution $\hat{W}_{1}$ of 4.3.6 satisfies

$$
\begin{equation*}
\left(s_{0}^{2}+\sqrt{|\tau|} s_{0}\right) \int_{\Omega_{1}}\left|\hat{W}_{1}(x, \tau)\right|^{2} d x+s_{0}\left|\hat{W}_{1}(0, \tau)\right|^{2} \leq C \int_{\Omega_{1}}\left|\hat{F}_{1}(x, \tau)\right|^{2} d x+C s_{0}\left|\hat{G}_{N}(\tau)\right|^{2} \tag{4.3.7}
\end{equation*}
$$

for all $\tau \in \mathbb{R}$ for some constant $C=C\left(\sigma_{1}, m_{1}\right)$ independent of $s_{0}$.

Proof. Let $\tau \in \mathbb{R}$ fixed. We multiply the equation $4.3 .6{ }_{1}$ by $s_{0} \overline{\hat{W}}_{1}(\tau)$ and we integrate on $\Omega_{1}$ :

$$
\begin{align*}
& \sigma_{1} s_{0} \int_{\Omega_{1}} \partial_{x} \hat{W}_{1}(\tau) \overline{\hat{W}_{1}}(\tau) d x+\sigma_{1} s_{0}^{2} \int_{\Omega_{1}} \partial_{x} \phi_{1}\left|\hat{W}_{1}(\tau)\right|^{2} d x+\gamma(\tau) s_{0} \int_{\Omega_{1}}\left|\hat{W}_{1}(\tau)\right|^{2} d x \\
= & -s_{0} \int_{\Omega_{1}} \hat{F}_{1}(\tau) \overline{\hat{W}_{1}}(\tau) d x . \tag{4.3.8}
\end{align*}
$$

Integration by parts yields

$$
\begin{equation*}
\sigma_{1} s_{0} \Re \int_{\Omega_{1}} \partial_{x} \hat{W}_{1}(\tau) \overline{\hat{W}_{1}}(\tau) d x=-\frac{\sigma_{1} s_{0}}{2}\left|\hat{W}_{1}(0, \tau)\right|^{2}-\frac{\sigma_{1}^{2} s_{0}}{2}\left|\hat{G}_{N}(\tau)\right|^{2} \tag{4.3.9}
\end{equation*}
$$

Taking the real part in 4.3.8) and using (4.3.9), we see that

$$
\begin{aligned}
& \frac{\sigma_{1} s_{0}}{2}\left|\hat{W}_{1}(0, \tau)\right|^{2}+\sigma_{1} s_{0}^{2} \int_{\Omega_{1}} \partial_{x} \phi_{1}\left|\hat{W}_{1}(x, \tau)\right|^{2} d x+\frac{\sqrt{|\tau|}}{2} s_{0} \int_{\Omega_{1}}\left|\hat{W}_{1}(x, \tau)\right|^{2} d x \\
= & -s_{0} \Re \int_{\Omega_{1}} \hat{F}_{1}(x, \tau) \overline{\hat{W}_{1}}(x, \tau) d x+\frac{\sigma_{1}^{2} s_{0}}{2}\left|\hat{G}_{N}(\tau)\right|^{2} .
\end{aligned}
$$

Applying Young's inequality, we have

$$
\begin{aligned}
& \sigma_{1} s_{0}^{2} \int_{\Omega_{1}} \partial_{x} \phi_{1}\left|\hat{W}_{1}(x, \tau)\right|^{2} d x+\Re(\gamma(\tau)) s_{0} \int_{\Omega_{1}}\left|\hat{W}_{1}(x, \tau)\right|^{2} d x+\frac{\sigma_{1} s_{0}}{2}\left|\hat{W}_{1}(0, \tau)\right|^{2} \\
\leq & \frac{m_{1}}{2} s_{0}^{2} \int_{\Omega_{1}}\left|\hat{W}_{1}(x, \tau)\right|^{2} d x+\frac{1}{2 m_{1}} \int_{\Omega_{1}}\left|\hat{F}_{1}(x, \tau)\right|^{2} d x+\frac{\sigma_{1}^{2} s_{0}}{2}\left|\hat{G}_{N}(\tau)\right|^{2}
\end{aligned}
$$

The proof of Proposition 4.9 is complete since $m_{1} \sigma^{-1}$ is a lower bound of $\partial_{x} \phi_{1}$.
Proposition 4.10 Let $\hat{W}_{2}$ be the solution of the system

$$
\begin{cases}\left(\sigma_{2} \partial_{x}+\sigma_{2} \partial_{x} \phi_{2} s_{0}+\gamma\right) \hat{W}_{2}(x, \tau)=-\hat{F}_{2}(x, \tau), & \forall(x, \tau) \in \Omega_{2} \times \mathbb{R}  \tag{4.3.10}\\ \sigma_{2} \hat{W}_{2}(0, \tau)=\sigma_{1} \hat{W}_{1}(0, \tau)+\left(\sigma_{1}-\sigma_{2}\right) \gamma \hat{Z}_{2}(0, \tau)+\hat{H}(\tau), & \forall \tau \in \mathbb{R} \\ \hat{W}_{2}\left(L_{2}, \tau\right)=\gamma(\tau) \hat{Z}_{2}\left(L_{2}, \tau\right)+\hat{J}(\tau), & \forall \tau \in \mathbb{R}\end{cases}
$$

Then, there exists a constant $C=C\left(m_{1}, m_{2}, \sigma_{1}, \sigma_{2}\right)$ such that each solution of 4.3.10) satisfies

$$
\begin{align*}
& \left(s_{0}^{2}+\sqrt{|\tau|} s_{0}\right) \int_{\Omega_{2}}\left|\hat{W}_{2}(\tau)\right|^{2} d x+s_{0}\left|\hat{W}_{2}\left(\tau, L_{2}\right)\right|^{2} \\
\leq & C \int_{\Omega_{1}}\left|\hat{F}_{1}(\tau)\right|^{2} d x+C \int_{\Omega_{2}}\left|\hat{F}_{2}(\tau)\right|^{2} d x+C s_{0}\left|\hat{G}_{N}(\tau)\right|^{2}+C|\tau| s_{0}\left|\hat{Z}_{2}(\tau, 0)\right|^{2}  \tag{4.3.11}\\
& +C s_{0}|\hat{H}(\tau)|^{2},
\end{align*}
$$

for each $\tau \in \mathbb{R}$.
Proof. Using the same ideas of Proposition 4.9, it is easy to check that

$$
\begin{equation*}
\left(s_{0}^{2}+\sqrt{|\tau|} s_{0}\right) \int_{\Omega_{2}}\left|\hat{W}_{2}(\tau)\right|^{2} d x+s_{0}\left|\hat{W}_{2}\left(\tau, L_{2}\right)\right|^{2} \leq C \int_{\Omega_{2}}\left|\hat{F}_{2}(\tau)\right|^{2} d x+C s_{0}\left|\hat{W}_{2}(\tau, 0)\right|^{2} \tag{4.3.12}
\end{equation*}
$$

for all $\tau \in \mathbb{R}$, where the constant $C$ in the inequality above depends only on $m_{2}$ and $\sigma_{2}$. Our next task is to estimate the local term of $\hat{W}_{2}$ at the interface. In order to do that, we use the boundary condition 4.3 .10$)_{2}$ to get the following estimate:

$$
\begin{equation*}
s_{0}\left|\hat{W}_{2}(\tau, 0)\right|^{2} \leq C s_{0}\left|\hat{W}_{1}(\tau, 0)\right|^{2}+C|\tau| s_{0}\left|\hat{Z}_{2}(\tau, 0)\right|^{2}+C s_{0}|\hat{H}(\tau)|^{2} \tag{4.3.13}
\end{equation*}
$$

Substituting (4.3.13) into (4.3.12) and applying the Proposition 4.9, we can assert that

$$
\begin{align*}
& \left(s_{0}^{2}+\sqrt{|\tau|} s_{0}\right) \int_{\Omega_{2}}\left|\hat{W}_{2}(\tau)\right|^{2} d x+s_{0}\left|\hat{W}_{2}\left(L_{2}, \tau\right)\right|^{2} \\
\leq & C \int_{\Omega_{1}}\left|\hat{F}_{1}(\tau)\right|^{2} d x+C \int_{\Omega_{2}}\left|\hat{F}_{2}(\tau)\right|^{2} d x+C s_{0}\left|\hat{G}_{N}(\tau)\right|^{2}+C|\tau| s_{0}\left|\hat{Z}_{2}(0, \tau)\right|^{2}  \tag{4.3.14}\\
& +C s_{0}|\hat{H}(\tau)|^{2}
\end{align*}
$$

which is the desired conclusion.

Proposition 4.11 Let $Z_{1}$ be the solution of

$$
\begin{cases}\left(\sigma_{1} \partial_{x}+\sigma_{1} \partial_{x} \phi_{1} s_{0}-\gamma\right) \hat{Z}_{1}(x, \tau)=\hat{W}_{1}(x, \tau), & \forall(x, \tau) \in \Omega_{1} \times \mathbb{R},  \tag{4.3.15}\\ \hat{Z}_{1}\left(-L_{1}, t\right)=0, & \forall \tau \in \mathbb{R}\end{cases}
$$

Then, there exists a constant $C=C\left(m_{1}, \sigma_{1}\right)$ such that each solution $\hat{Z}_{1}$ of 4.3.15) with source term $\hat{W}_{2}(\tau) \in L^{2}(\Omega)$ for all $\tau \in \mathbb{R}$ satisfies

$$
\begin{align*}
& r_{1} \sigma_{1} s_{0}^{3} \int_{\Omega_{1}}\left|\hat{Z}_{1}(x, \tau)\right|^{2} d x+s_{0}^{2} \int_{\Omega_{1}}\left(\sigma_{1} \partial_{x} \phi_{1} s_{0}-\Re(\gamma(\tau))\right)^{2}\left|\hat{Z}_{1}(x, \tau)\right|^{2} d x \\
& +\sigma_{1}\left(m_{1} s_{0}-\Re(\gamma(\tau))\right) s_{0}^{2}\left|\hat{Z}_{1}(0, \tau)\right|^{2} \leq C \int_{\Omega_{1}}\left|\hat{F}_{1}(x, \tau)\right|^{2} d x+C s_{0}\left|\hat{G}_{N}(\tau)\right|^{2}, \quad \forall \tau \in \mathbb{R} \tag{4.3.16}
\end{align*}
$$

Proof. Let $\tau \in \mathbb{R}$. We multiply the equation 4.3.15) by $\left(\sigma_{1} \partial_{x} \phi_{1} s_{0}-\Re(\gamma)\right) \overline{\hat{Z}}_{1}$ and we integrate on $\Omega_{1}$ :

$$
\begin{align*}
& \sigma_{1} \int_{\Omega_{1}}\left(\sigma_{1} \partial_{x} \phi_{1} s_{0}-\Re(\gamma)\right) \partial_{x} \hat{Z}_{1}(x, \tau) \overline{\hat{Z}_{1}}(\tau) d x \\
& \quad+\int_{\Omega_{1}}\left(\sigma_{1} \partial_{x} \phi_{1} s_{0}-\gamma\right)\left(\sigma_{1} \partial_{x} \phi_{1} s_{0}-\Re(\gamma)\right)\left|\hat{Z}_{1}(x, \tau)\right|^{2} d x  \tag{4.3.17}\\
& =\int_{\Omega_{1}}\left(\sigma_{1} \partial_{x} \phi_{1} s_{0}-\Re(\gamma)\right) \hat{W}_{1}(x, \tau) \overline{\hat{Z}}_{1}(x, \tau) d x .
\end{align*}
$$

Integration by parts yields,

$$
\begin{align*}
& \sigma_{1} \Re \int_{\Omega_{1}}\left(\sigma_{1} \partial_{x} \phi_{1} s_{0}-\Re(\gamma)\right) \partial_{x} \hat{Z}_{1}(x, \tau) \overline{\hat{Z}_{1}}(x, \tau) d x  \tag{4.3.18}\\
= & -\frac{\sigma_{1}^{2} s_{0}}{2} \int_{\Omega_{1}} \partial_{x}^{2} \phi_{1}\left|\hat{Z}_{1}(x, \tau)\right|^{2} d x-\frac{\sigma_{1}}{2}\left(m_{1} s_{0}-\Re(\gamma)\right)\left|\hat{Z}_{1}(0, \tau)\right|^{2} .
\end{align*}
$$

Taking the real part in (4.3.17) and using (4.3.18), we have

$$
\begin{aligned}
& -\frac{\sigma_{1}^{2}}{2} s_{0} \int_{\Omega_{1}} \partial_{x} \phi_{1}\left|\hat{Z}_{1}(x, \tau)\right|^{2} d x+\int_{\Omega_{1}}\left(\sigma_{1} \partial_{x} \phi_{1} s_{0}-\Re(\gamma)\right)^{2}\left|\hat{Z}_{1}(x, \tau)\right|^{2} d x \\
& -\sigma_{1}\left(m_{1} s_{0}-\Re(\gamma)\right)\left|\hat{Z}_{1}(0, \tau)\right|^{2} \\
= & \int_{\Omega_{1}}\left(\sigma_{1} \partial_{x} \phi_{1} s_{0}-\Re(\gamma)\right) \hat{W}_{1}(x, \tau) \bar{Z}_{1}(x, \tau) d x .
\end{aligned}
$$

By the Young inequality and the assumptions on $\phi_{1}$, we see that

$$
\begin{align*}
& r_{1} \sigma_{1} s_{0} \int_{\Omega_{1}}\left|\hat{Z}_{1}(x, \tau)\right|^{2} d x+\int_{\Omega_{1}}\left(\sigma_{1} \partial_{x} \phi_{1} s_{0}-\Re(\gamma)\right)^{2}\left|\hat{Z}_{1}(x, \tau)\right|^{2} d x \\
& +\sigma_{1}\left(m_{1} s_{0}-\Re(\gamma)\right)\left|\hat{Z}_{1}(0, \tau)\right|^{2} \leq \int_{\Omega_{1}}\left|\hat{W}_{1}(x, \tau)\right|^{2} d x \tag{4.3.19}
\end{align*}
$$

Finally, we multiply 4.3.19 by $s_{0}^{2}$ and apply the Proposition 4.9, and the proof is complete.

Proposition 4.12 Let $\hat{Z}_{2}$ be the solution of

$$
\begin{cases}\left(\sigma_{2} \partial_{x}+\sigma_{2} \partial_{x} \phi_{2} s_{0}-\gamma\right) \hat{Z}_{2}(x, \tau)=\hat{W}_{2}(x, \tau), & \forall(x, \tau) \in \Omega_{2} \times \mathbb{R}  \tag{4.3.20}\\ \hat{Z}_{2}(0, \tau)=\hat{Z}_{1}(0, \tau), & \forall \tau \in \mathbb{R}\end{cases}
$$

Then, there exists a constant $C>0$ independent of $s_{0}$ such that for every solution $\hat{Z}_{2}$ of (4.12), we have

$$
\begin{align*}
& r_{2} \sigma_{2} s_{0}^{3} \int_{\Omega_{2}}\left|\hat{Z}_{2}(x, \tau)\right|^{2} d x+s_{0} \int_{\Omega_{2}}\left(\sigma_{2} \partial_{x} \phi_{2} s_{0}-\Re(\gamma(\tau))\right)\left|\hat{Z}_{2}(x, \tau)\right|^{2} d x \\
& \quad+\sigma_{2}\left(m_{2} s_{0}-\Re(\gamma(\tau))\right) s_{0}^{2}\left|\hat{Z}_{2}\left(L_{2}, \tau\right)\right|^{2}  \tag{4.3.21}\\
& \leq C \int_{\Omega_{1}}\left|\hat{F}_{1}(x, \tau)\right|^{2} d x+C \int_{\Omega_{2}}\left|\hat{F}_{2}(x, \tau)\right|^{2} d x \\
& \quad+C s_{0}\left|\hat{G}_{N}(\tau)\right|^{2}+C|\tau| s_{0}\left|\hat{Z}_{2}(0, \tau)\right|^{2}+C s_{0}|\hat{H}(\tau)|^{2},
\end{align*}
$$

for each $\tau \in \mathbb{R}$.
Proof. Firstly, as in the proof of Proposition 4.12 each solution $Z_{2}$ of 4.3.20 satisfies

$$
\begin{align*}
& \quad r_{2} \sigma_{2} s_{0}^{3} \int_{\Omega_{2}}\left|\hat{Z}_{2}(x, \tau)\right|^{2} d x+s_{0}^{2} \int_{\Omega_{2}}\left(\sigma_{2} \partial_{x} \phi_{2}-\Re(\gamma(\tau))\right)\left|\hat{Z}_{2}(x, \tau)\right|^{2} d x \\
& \quad+\sigma_{2}\left(m_{2} s_{0}-\Re(\gamma(\tau))\right) s_{0}^{2}\left|\hat{Z}_{2}\left(L_{2}, \tau\right)\right|^{2}  \tag{4.3.22}\\
& \leq s_{0}^{2} \int_{\Omega_{2}}\left|\hat{W}_{2}(x, \tau)\right|^{2} d x+\sigma_{2}\left(M_{2} s_{0}-\Re(\gamma(\tau))\right) s_{0}\left|\hat{Z}_{2}(0, \tau)\right|^{2}, \quad \forall \tau \in \mathbb{R} .
\end{align*}
$$

Then, it remains to estimate the global term of $\hat{W}_{2}$ and the local term of $\hat{Z}_{2}$ at the interface. To do this, notice that from the Proposition 4.12, the global term of $\hat{W}_{2}$ can be bounded as follows:

$$
\begin{gather*}
s_{0}^{2} \int_{\Omega_{2}}\left|\hat{W}_{2}(x, \tau)\right|^{2} d x \leq C \int_{\Omega_{1}}\left|\hat{F}_{1}(x, \tau)\right|^{2} d x+C \int_{\Omega_{2}}\left|\hat{F}_{2}(x, \tau)\right|^{2} d x+C s_{0}\left|\hat{G}_{N}(\tau)\right|^{2}  \tag{4.3.23}\\
\\
+C|\tau| s_{0}\left|\hat{Z}_{2}(0, \tau)\right|^{2}+C s_{0}|\hat{H}(\tau)|^{2}
\end{gather*}
$$

On the other hand, since $M_{2}<m_{1}$ and $Z_{1}=Z_{2}$ at the interface for all $\tau \in \mathbb{R}$, we deduce that

$$
\begin{align*}
\left(M_{2} s_{0}-\Re(\gamma(\tau))\right) s_{0}^{2}\left|\hat{Z}_{2}(0, \tau)\right|^{2} & \leq\left(m_{1} s_{0}-\Re(\gamma(\tau))\right) s_{0}^{2}\left|\hat{Z}_{1}(0, \tau)\right|^{2} \\
& \leq C \int_{\Omega_{2}}\left|\hat{F}_{1}(x, \tau)\right|^{2} d x+C s_{0}\left|\hat{G}_{N}(\tau)\right|^{2}, \quad \forall \tau \in \mathbb{R} \tag{4.3.24}
\end{align*}
$$

Finally, we substitute (4.3.23) and (4.3.24) into (4.3.24) and the proof is complete.

### 4.3.2 Global estimates in the Fourier domain

Our next task is prove the following inequality

$$
\begin{align*}
& s_{0}^{3} \int_{\Omega_{1}}\left|\hat{Z}_{1}(x, \tau)\right|^{2} d x+s_{0}^{3} \int_{\Omega_{2}}\left|\hat{Z}_{2}(x, \tau)\right|^{2} d x+s_{0}^{3}\left|Z_{1}(0, \tau)\right|^{2}+s_{0}^{3}\left|\hat{Z}_{2}\left(L_{2}, \tau\right)\right|^{2} \\
\leq & C \int_{\Omega_{1}}\left|\hat{F}_{1}(x, \tau)\right|^{2} d x+C \int_{\Omega_{2}}\left|\hat{F}_{2}(x, \tau)\right|^{2} d x+C s_{0}\left|\hat{G}_{N}(\tau)\right|^{2}+C s_{0}|\hat{H}(\tau)|^{2}+C s_{0}|\hat{J}(\tau)|^{2}, \tag{4.3.25}
\end{align*}
$$

for all $\tau \in \mathbb{R}$ by using the estimates of Propositions 4.9, 4.10, 4.11 and 4.12, In order to do that, let $\delta$ and $\delta^{\prime}$ be two positive numbers such that

$$
\delta<\sqrt{2} m_{2}<\sqrt{2} M_{2}<\delta^{\prime}<\sqrt{2} m_{1}
$$

Then, we divide the real line into three subsets, namely Low frequencies, Intermediate frequencies and High frequencies. These intervals are ilustrated in the following Figure:

|  |  |  |  |
| :--- | :---: | :---: | :---: |
| Low frequencies | Intermediate frequencies |  | High frequencies |
| $\delta s_{0}$ | $\sqrt{2} m_{2} s_{0}$ | $\sqrt{2} M_{2} s_{0}$ | $\delta^{\prime} s_{0}$ |
|  |  | $\sqrt{2} m_{1} s_{0}$ |  |
|  | Case 1 |  |  |
|  |  | Case 2 |  |

Figure 4.2: Sketch of different ranges in frequency domain

- Case 1: Low frequencies Suppose that $\tau \in \mathbb{R}$ is such that

$$
\begin{equation*}
0 \leq \sqrt{|\tau|} \leq \delta s_{0} \tag{4.3.26}
\end{equation*}
$$

We note that the condition above implies

$$
\frac{1}{\sqrt{2}}\left(\sqrt{2} m_{1}-\delta^{\prime}\right) s_{0}<m_{1} s_{0}-\sqrt{\frac{|\tau|}{2}}
$$

By Proposition 4.11, we see that

$$
\begin{equation*}
\frac{\sigma_{1}}{\sqrt{2}}\left(\sqrt{2} m_{1}-\delta^{\prime}\right) s_{0}^{3}\left|\hat{Z}_{1}(0, \tau)\right|^{2} \leq C \int_{\Omega_{1}}\left|\hat{F}_{1}(x, \tau)\right|^{2} d x+C s_{0}\left|\hat{G}_{N}(\tau)\right|^{2} \tag{4.3.27}
\end{equation*}
$$

Then, we can estimate the global term of $\hat{Z}_{1}$ and its local term at the interface as follows:

$$
\begin{equation*}
s_{0}^{3} \int_{\Omega_{1}}\left|\hat{Z}_{1}(x, \tau)\right|^{2} d x+s_{0}^{3}\left|\hat{Z}_{2}(0, \tau)\right|^{2} \leq C \int_{\Omega_{1}}\left|\hat{F}_{1}(x, \tau)\right|^{2} d x+C s_{0}\left|\hat{G}_{N}(\tau)\right|^{2} \tag{4.3.28}
\end{equation*}
$$

On the other hand, the condition 4.3.26 also implies

$$
\frac{1}{\sqrt{2}}\left(\sqrt{2} m_{2}-\delta\right) s_{0}<m_{2} s_{0}-\sqrt{\frac{|\tau|}{2}}
$$

Thus, similarly to the estimate (4.3.27), we can assert that

$$
\begin{align*}
s_{0}^{3} \int_{\Omega_{2}}\left|\hat{Z}_{2}(x, \tau)\right|^{2} d x+s_{0}^{3}\left|\hat{Z}_{2}\left(L_{2}, \tau\right)\right|^{2} \leq & C \int_{\Omega_{1}}\left|\hat{F}_{1}(x, \tau)\right|^{2} d x+C \int_{\Omega_{2}}\left|\hat{F}_{2}(\tau)\right|^{2} d x+C s_{0}\left|\hat{G}_{N}(\tau)\right|^{2} \\
& +C s_{0}|\tilde{H}(\tau)|^{2} \tag{4.3.29}
\end{align*}
$$

where we applied the Proposition 4.12. Finally, we add the inequalities (4.3.27) and (4.3.29), and the proof of 4.3.25) in the case of Low frequencies is complete.

- Case 2: Intermediate frequencies Suppose that $\tau \in \mathbb{R}$ is chosen such that

$$
\delta s_{0}<\sqrt{|\tau|} \leq \sqrt{2} m_{1} s_{0}
$$

To do this, we consider two cases (see figure 1.2)

- We assume that $\tau \in \mathbb{R}$ is such that

$$
\delta s_{0}<\sqrt{|\tau|} \leq \delta^{\prime} s_{0}
$$

In this case, we already have the estimate 4.3.27). Besides, the estimate 4.3.29) does not hold. To deal with this issue, by the boundary condition 4.3 .10$)_{3}$ we can assert that

$$
\begin{equation*}
|\tau| s_{0}\left|\hat{Z}_{2}\left(L_{2}, \tau\right)\right|^{2} \leq 2 s_{0}\left|\hat{W}_{2}\left(L_{2}, \tau\right)\right|^{2}+2 s_{0}|\tilde{J}(\tau)|^{2} \tag{4.3.30}
\end{equation*}
$$

From Proposition 4.10, we have the following upper bound of $W_{2}\left(\tau, L_{2}\right)$ :

$$
\begin{equation*}
s_{0}\left|\hat{W}_{2}\left(L_{2}, \tau\right)\right|^{2} \leq C \int_{\Omega_{1}}\left|\hat{F}_{1}(x, \tau)\right|^{2} d x+C \int_{\Omega_{2}}\left|\hat{F}_{2}(x, \tau)\right|^{2} d x+C|\tau| s_{0}\left|\hat{Z}_{2}(0, \tau)\right|^{2}+C s_{0}|\hat{H}(\tau)|^{2} \tag{4.3.31}
\end{equation*}
$$

Combining 4.3.30 and 4.3.31 with 4.3.27) we see that

$$
\begin{equation*}
s_{0}^{3}\left|\hat{Z}_{2}\left(L_{2}, \tau\right)\right|^{2} \leq C \int_{\Omega_{1}}\left|F_{1}(x, \tau)\right|^{2} d x+C \int_{\Omega_{2}}\left|F_{2}(x, \tau)\right|^{2} d x+C s_{0}\left|\hat{G}_{N}(\tau)\right|^{2}+C s_{0}|\hat{H}(\tau)|^{2}+C s_{0}|\hat{J}(\tau)|^{2} . \tag{4.3.32}
\end{equation*}
$$

Hence, by Proposition 4.11 and 4.12, we get

$$
\begin{align*}
s_{0}^{3} \int_{\Omega_{2}}\left|\hat{Z}_{2}(x, \tau)\right|^{2} d x \leq & C \int_{\Omega_{1}}\left|\hat{F}_{1}(x, \tau)\right|^{2} d x+C \int_{\Omega_{2}}\left|\hat{F}_{2}(x, \tau)\right|^{2} d x+C s_{0}\left|\hat{G}_{N}(\tau)\right|^{2}+C s_{0}|\hat{H}(\tau)|^{2} \\
& +C s_{0}|\hat{J}(\tau)|^{2} \tag{4.3.33}
\end{align*}
$$

Thus, we add 4.3.28, 4.3.32 and 4.3.33), this is precisely the claim in the case 1 of Intermediate frequencies.

- We consider $\tau \in \mathbb{R}$ such that

$$
\delta^{\prime} s_{0}<\sqrt{|\tau|}<\sqrt{2} m_{1} s_{0}
$$

In this case, we adopt a different strategy in order to estimate the local terms of $\hat{Z}_{1}$ and $\hat{Z}_{2}$ at the interface. Roughly speaking, we will descompose the solution $\hat{Z}_{2}$ into two components: the first one is unknown in the sense that this one depends of the local term $\hat{Z}_{2}(\tau, 0)$, and the second one can be estimated using the same machinery introduced in the Propositions above.
Before to start with the proof in this case, let us state some definitions. We consider the functions $\hat{Z}_{u}$ and $\hat{Z}_{k}$, defined as the solution of the following problems:

$$
\begin{cases}\left(\sigma_{2} \partial_{x}+\sigma_{2} \partial_{x} \phi_{2} s_{0}-\gamma\right) \hat{Z}_{u}(x, \tau)=u(x, \tau), & \forall(x, \tau) \Omega_{2} \times \mathbb{R},  \tag{4.3.34}\\ \hat{Z}_{u}\left(L_{2}, \tau\right)=\frac{1}{\gamma} u\left(L_{2}, \tau\right), & \forall \tau \in \mathbb{R},\end{cases}
$$

and

$$
\begin{cases}\left(\sigma_{2} \partial_{x}+\sigma_{2} \partial_{x} \phi_{2} s_{0}-\gamma\right) \hat{Z}_{k}(x, \tau)=k(x, \tau), & \forall(x, \tau) \in \Omega_{2} \times \mathbb{R}  \tag{4.3.35}\\ \hat{Z}_{k}\left(L_{2}, \tau\right)=\frac{1}{\gamma} k\left(L_{2}, \tau\right)-\frac{1}{\gamma} \hat{J}(\tau), & \forall \tau \in \mathbb{R}\end{cases}
$$

respectively. Here, the functions $u$ and $k$ solve

$$
\begin{cases}\left(\sigma_{2} \partial_{x}+\sigma_{2} \partial_{x} \phi_{2} s_{0}+\gamma\right) u(x, \tau)=0, & \forall(x, \tau) \in \Omega_{2} \times \mathbb{R}  \tag{4.3.36}\\ u(0, \tau)=\frac{\left(\sigma_{1}-\sigma_{2}\right)}{\sigma_{2}} \gamma \hat{Z}_{2}(0, \tau), & \forall \tau \in \mathbb{R}\end{cases}
$$

and

$$
\begin{cases}\left(\sigma_{2} \partial_{x}+\sigma_{2} \partial_{x} \phi_{2} s_{0}+\gamma\right) k(x, \tau)=-F_{2}(x, \tau), & \forall(x, \tau) \in \Omega_{2} \times \mathbb{R},  \tag{4.3.37}\\ k(0, \tau)=\frac{\sigma_{1}}{\sigma_{2}} \hat{W}_{1}(0, \tau)+\frac{1}{\sigma_{2}} \hat{H}(\tau), & \forall \tau \in \mathbb{R}\end{cases}
$$

respectively. It is clear that $\hat{Z}_{2}=\hat{Z}_{u}+\hat{Z}_{k}$ and $\hat{W}_{2}=u+k$ in $\mathbb{R} \times \Omega_{2}$. Let us compute the explicit solution of $\hat{Z}_{u}$. Firstly, by Duhamel's formula, the expression of $u$ in 4.3.36 is given by

$$
\begin{equation*}
u(x, \tau)=\left(\frac{\sigma_{1}-\sigma_{2}}{\sigma_{2}}\right) \gamma(\tau) \hat{Z}_{2}(0, \tau) \exp \left(\left(\phi_{2}(0)-\phi_{2}(x)\right) s_{0}-\frac{\gamma(\tau)}{\sigma_{2}} x\right) \tag{4.3.38}
\end{equation*}
$$

Then, using Duhamel's formula once again, the explicit solution of $\hat{Z}_{u}$ in terms of $u$ is given by

$$
\begin{align*}
\hat{Z}_{u}(x, \tau)= & \frac{1}{\gamma(\tau)} u\left(\tau, L_{2}\right) \exp \left(\left(\phi_{2}\left(L_{2}\right)-\phi_{2}(x)\right) s_{0}-\frac{\gamma(\tau)}{\sigma_{2}}\left(L_{2}-x\right)\right) \\
& -\frac{1}{\sigma_{2}} \exp \left(\frac{\gamma(\tau)}{\sigma_{2}} x-\phi_{2}(x) s_{0}\right) \int_{x}^{L_{2}} \exp \left(\phi_{2}(\tilde{x}) s_{0}-\frac{\gamma(\tau)}{\sigma_{2}} \tilde{x}\right) u(\tilde{x}, \tau) d \tilde{x} \tag{4.3.39}
\end{align*}
$$

Substituting 4.3.38 by 4.3.39) and evaluating at $x=0$, we have

$$
\hat{Z}_{u}(0, \tau)=\frac{\sigma_{1}-\sigma_{2}}{2 \sigma_{2}} \hat{Z}_{2}(0, \tau)\left(3 \exp \left(-\frac{2 \gamma(\tau) L_{2}}{\sigma_{2}}\right)-1\right)
$$

Therefore, since $\hat{Z}_{2}(0, \tau)=\hat{Z}_{u}(0, \tau)+\hat{Z}_{k}(0, \tau)$, we obtain

$$
\begin{equation*}
\Lambda(\tau) \hat{Z}_{2}(0, \tau)=\hat{Z}_{k}(0, \tau) \tag{4.3.40}
\end{equation*}
$$

where $\Lambda=\Lambda(\tau)$ is defined by

$$
\Lambda(\tau):=\frac{\sigma_{1}+\sigma_{2}}{2 \sigma_{2}}+3 \frac{\left(\sigma_{2}-\sigma_{1}\right)}{2 \sigma_{2}} \exp \left(-\frac{2 \gamma(\tau) L_{2}}{\sigma_{2}}\right)
$$

Now we will show that there exists a constant $C>0$ independent of $s_{0}$ such that

$$
\left|\hat{Z}_{2}(\tau, 0)\right|^{2} \leq C\left|\hat{Z}_{k}(\tau, 0)\right|^{2} .
$$

Indeed, suppose that $0<\sigma_{1}<\sigma_{2}$. Then, it is easy to check that

$$
\frac{1}{2} \leq \frac{\sigma_{1}+\sigma_{2}}{2 \sigma_{2}} \leq \Re(\Lambda(\tau)),
$$

and the assertion follows directly. Now, we want to show the same inequality when $0<\sigma_{2}<\sigma_{1}$. In that case, we consider the following assumptions: we suppose that

$$
s \geq \frac{T^{2 \alpha}}{2^{2 \alpha}}
$$

and

$$
M_{2} \geq \frac{\log (3) \sigma_{2}}{2 L_{2}}
$$

In particular these conditions implies that $s_{0} \geq 1$ and

$$
\sqrt{|\tau|} \geq \frac{\log (3) \sigma_{2}}{\sqrt{2} L_{2}}
$$

or equivalently

$$
\exp \left(-\frac{\sqrt{2|\tau| L_{2}}}{\sigma_{2}}\right) \leq \frac{1}{3}
$$

Hence, notice that under these conditions, we can assert that

$$
\begin{equation*}
\Re(\Lambda(\tau)) \geq 1 \tag{4.3.41}
\end{equation*}
$$

This implies the desired claim. It remains to prove a estimate for $\hat{Z}_{k}$ at the interface. In order to do that, using the same ideas of Proposition 4.10, we can assert that

$$
\begin{align*}
& \left(s_{0}^{2}+\sqrt{|\tau|}\right) s_{0} \int_{\Omega_{2}}|k(x, \tau)|^{2} d x+s_{0}\left|k\left(L_{2}, \tau\right)\right|^{2} \\
\leq & C \int_{\Omega_{2}}\left|\hat{F}_{2}(x, \tau)\right|^{2} d x+C s_{0}\left|\hat{W}_{1}(0, \tau)\right|^{2}+C s_{0}|\hat{H}(\tau)|^{2} . \tag{4.3.42}
\end{align*}
$$

Furthermore, we can prove the following estimate for $\hat{Z}_{2}$ :

$$
\begin{align*}
& r_{2} \sigma_{2} s_{0} \int_{\Omega_{2}}\left|\hat{Z}_{k}(0, \tau)\right|^{2} d x+\int_{\Omega_{2}}\left(\sigma_{2} \partial_{x} \phi_{2} s_{0}-\Re(\gamma(\tau))\right)^{2}\left|\hat{Z}_{k}(x, \tau)\right|^{2} d x \\
& \quad+\sigma_{2}\left(\Re(\gamma(\tau))-M_{2} s_{0}\right)\left|\hat{Z}_{k}(0, \tau)\right|^{2}  \tag{4.3.43}\\
& \leq C \int_{\Omega_{2}}|k(x, \tau)|^{2} d x+C \frac{1}{\sqrt{|\tau|}}\left(\left|k\left(L_{2}, \tau\right)\right|^{2}+|\hat{J}(\tau)|^{2}\right) .
\end{align*}
$$

Now, since we are in the second case of Intermediate frequencies, we have

$$
\Re(\gamma(\tau))-M_{2} s_{0}>\left(\delta^{\prime}-M_{2}\right) s_{0}
$$

Then, from (4.3.43), we deduce that

$$
\begin{equation*}
\sqrt{|\tau|} s_{0}\left|\hat{Z}_{k}(0, \tau)\right|^{2} \leq C \sqrt{|\tau|} \int_{\Omega_{2}}|k(x, \tau)|^{2} d x+C s_{0}\left|k\left(L_{2}, \tau\right)\right|^{2}+C s_{0}|\hat{J}(\tau)|^{2} \tag{4.3.44}
\end{equation*}
$$

Multiplying by $s_{0}$ the inequality above and using the lower and upper bound of $\sqrt{|\tau|}$, we see that

$$
s_{0}^{3}\left|\hat{Z}_{k}(0, \tau)\right|^{2} \leq C \int_{\Omega_{2}}\left|\hat{F}_{2}(x, \tau)\right|^{2} d x+C s_{0}\left|\hat{W}_{1}(0, \tau)\right|^{2}+C s_{0}|H(\tau)|^{2}+C s_{0}|\hat{J}(\tau)|^{2}
$$

Moreover, applying the Proposition 4.9, we obtain

$$
\begin{aligned}
s_{0}^{3}\left|\hat{Z}_{k}(0, \tau)\right|^{2} \leq & C \int_{\Omega_{1}}\left|\hat{F}_{1}(x, \tau)\right|^{2} d x+C \int_{\Omega_{2}}\left|\hat{F}_{2}(x, \tau)\right|^{2} d x+C s_{0}\left|\hat{G}_{N}(\tau)\right|^{2} \\
& +C s_{0}|\hat{H}(\tau)|^{2}+C s_{0}|\hat{J}(\tau)|^{2}
\end{aligned}
$$

Therefore, using the relation 4.3.40, we obtain

$$
\begin{align*}
s_{0}^{3}\left|\hat{Z}_{2}(0, \tau)\right|^{2} \leq & C \int_{\Omega_{1}}\left|\hat{F}_{1}(x, \tau)\right|^{2} d x+C \int_{\Omega_{2}}\left|\hat{F}_{2}(x, \tau)\right|^{2} d x+C s_{0}|\hat{G}(\tau)|^{2}  \tag{4.3.45}\\
& +C s_{0}|H(\tau)|^{2}+C s_{0}|J(\tau)|^{2}
\end{align*}
$$

The rest of the proof runs as before. Thus, we proved the desired inequality in the second case of Intermediate frequencies.

## - High frequencies:

Now, we consider the case of high frequencies, that is, we take $\tau \in \mathbb{R}$ such that

$$
\begin{equation*}
\sqrt{2} m_{1} s_{0} \leq \sqrt{|\tau|} . \tag{4.3.46}
\end{equation*}
$$

In this case, we will apply the same strategy as the case before. However, we can not estimate directly the term $|\tau| s_{0}\left|\hat{Z}_{2}(\tau, 0)\right|^{2}$ in (4.3.44). To avoid this difficulty, we note that the condition 4.3.46) implies

$$
\frac{\sqrt{|\tau|}}{2 \sqrt{2}} \leq \Re(\gamma(\tau))-M_{2} s_{0}
$$

Therefore, from the estimate of $\hat{Z}_{k}$ in 4.3.43, we have

$$
\begin{equation*}
\left.\sqrt{|\tau|} \hat{Z}_{k}(\tau, 0)\right|^{2} \leq C \int_{\Omega_{2}}|k(\tau)|^{2} d x+C \frac{\sigma_{2}}{\sqrt{|\tau|}}\left(\left|k\left(\tau, L_{2}\right)\right|^{2}+|\hat{J}(\tau)|^{2}\right) . \tag{4.3.47}
\end{equation*}
$$

Equivalently,

$$
|\tau| s_{0}\left|\hat{Z}_{k}(\tau, 0)\right|^{2} \leq C \sqrt{|\tau|} \int_{\Omega_{2}}|k(\tau)|^{2} d x+C \sigma_{2} s_{0}\left|k\left(\tau, L_{2}\right)\right|^{2}+C \sigma_{2} s_{0}\left|\hat{J}\left(\tau, L_{2}\right)\right|^{2}
$$

Using the estimate of $k$ in 4.3.42, we obtain

$$
\begin{aligned}
& |\tau| s_{0}\left|\hat{Z}_{2}(0, \tau)\right|^{2} \\
\leq & C|\tau| s_{0}\left|\hat{Z}_{k}(0, \tau)\right|^{2} \\
\leq & C \int_{\Omega_{1}}\left|\hat{F}_{1}(x, \tau)\right|^{2} d x+C \int_{\Omega_{2}}\left|\hat{F}_{2}(x, \tau)\right|^{2} d x+C s_{0}\left|\hat{G}_{N}(\tau)\right|^{2}+C s_{0}|\hat{H}(\tau)|^{2}+C s_{0}|\hat{J}(\tau)|^{2} .
\end{aligned}
$$

We proceed for estimate the term of $\hat{Z}_{2}\left(\tau, L_{2}\right)$ as before. This conclude the proof of (4.3.25) in the case of high frequencies. Hence, we have proved the estimate 4.3.25) for all $\tau \in \mathbb{R}$. Now, integrating on $\mathbb{R}$ in $\tau$, we deduce that

$$
\begin{align*}
& \quad s_{0}^{3} \int_{\mathbb{R}} \int_{\Omega_{1}}\left|\hat{Z}_{1}(x, \tau)\right|^{2} d x d t+s_{0}^{3} \int_{\mathbb{R}} \int_{\Omega_{2}}\left|\hat{Z}_{2}(x, \tau)\right|^{2} d x d t \\
& \leq  \tag{4.3.48}\\
& \leq \int_{\mathbb{R}} \int_{\Omega_{1}}\left|\hat{F}_{1}(x, \tau)\right|^{2} d x d \tau+\int_{\mathbb{R}} \int_{\Omega_{2}}\left|\hat{F}_{2}(x, \tau)\right|^{2} d x d \tau+C s_{0} \int_{\mathbb{R}}\left|\hat{G}_{N}\right|^{2} d \tau \\
& \quad+C s_{0} \int_{\mathbb{R}}|\hat{H}(\tau)|^{2} d \tau+C s_{0} \int_{\mathbb{R}}|\hat{J}(\tau)|^{2} d \tau .
\end{align*}
$$

Finally, we use the Parseval's identity and the fact that all functions are supported in an open subset of $(0, T)$. This concludes the proof of the inequality 4.3.1.

### 4.3.3 Estimates of the spatial derivatives

The goal of this section is show that the spatial derivatives $Z_{1}$ and $Z_{2}$ of 4.2.5 with source terms $F_{j} \in L^{2}\left(\Omega_{j} \times(0, T)\right) j=1,2$, Neumann data $G_{N} \in L^{2}(0, T)$ and residual terms $H, J \in L^{2}(0, T)$ satisfies the following:

$$
\begin{align*}
& s_{0} \int_{0}^{T} \int_{\Omega_{1}}\left|\partial_{x} Z_{1}(x, \tau)\right|^{2} d x d t+s_{0} \int_{0}^{T} \int_{\Omega_{2}}\left|\partial_{x} Z_{2}(x, \tau)\right|^{2} d x d t \\
\leq & C \int_{0}^{T} \int_{\Omega_{1}}\left|\hat{F}_{1}(x, \tau)\right|^{2} d x d t+C \int_{0}^{T} \int_{\Omega_{2}}\left|\hat{F}_{2}(x, \tau)\right|^{2} d x d t+C s_{0} \int_{0}^{T}\left|\hat{G}_{N}\right|^{2} d t+C s_{0} \int_{0}^{T}|\hat{H}|^{2} d t \\
& +C s_{0} \int_{0}^{T}|\hat{J}|^{2} d t \tag{4.3.49}
\end{align*}
$$

where the constant $C=C\left(m_{1}, m_{2}, r_{1}, r_{2}, \sigma_{1}, \sigma_{2}\right)$ is independent of $s_{0}$. We start multiplying the first equation by $Z_{1}$ and we integrate on $\Omega_{1}$ :

$$
\frac{d}{d t} \int_{\Omega_{1}}\left|Z_{1}(x, t)\right|^{2} d x-\sigma_{1} \int_{\Omega_{1}}\left(\partial_{x}+s_{0} \partial_{x} \phi_{1}\right)^{2} Z_{1}(x, t) \cdot Z_{1}(x, t) d x=\int_{\Omega_{1}} F_{1}(x, t) Z_{1}(x, t) d x
$$

for all $t \in(0, T)$ where by definition

$$
\left(\partial_{x}+s \partial_{x} \phi_{1}\right)^{2} Z_{1}=\partial_{x}^{2} Z_{1}+s_{0} \partial_{x} \phi_{1} \partial_{x} Z_{1}+s_{0} \partial_{x}^{2} \phi_{1} Z_{1}+\left|s_{0} \partial_{x} \phi_{1}\right|^{2} Z_{1}, \quad \forall(x, t) \in \Omega_{1} \times(0, T) .
$$

Then,

$$
\begin{align*}
& -\sigma_{1}^{2} \int_{\Omega_{1}}\left(\partial_{x}+s_{0} \partial_{x} \phi_{1}\right)^{2} Z_{1}(x, t) \cdot Z_{1}(x, t) d x \\
= & -\sigma_{1}^{2} \int_{\Omega_{1}} \partial_{x} Z_{1}(x, t) Z_{1}(x, t) d x-2 \sigma_{1}^{2} s_{0} \int_{\Omega_{1}} \partial_{x} \phi_{1} \partial_{x} Z_{1}(x, t) Z_{1}(x, t) d x  \tag{4.3.50}\\
& -\sigma_{1}^{2} s_{0} \int_{\Omega_{1}} \partial_{x}^{2} \phi_{1}\left|Z_{1}(x, t)\right|^{2} d x-s_{0}^{2} \int_{\Omega_{1}}\left|\partial_{x} \phi_{1}\right|^{2}\left|Z_{1}(x, t)\right|^{2} d x .
\end{align*}
$$

Integration by parts yields

$$
-\int_{\Omega_{1}} \partial_{x}^{2} Z_{1}(x, t) Z_{1}(x, t) d x=\int_{\Omega_{1}}\left|\partial_{x} Z_{1}(x, t)\right|^{2} d x-\partial_{x} Z_{1}(0, t) Z_{1}(0, t)
$$

where we used the condition $Z_{1}\left(L_{2}, t\right)=0$ for each $t \in(0, T)$ and

$$
-2 s_{0} \int_{\Omega_{1}} \partial_{x} \phi_{1} Z_{1}(x, t) Z_{1}(x, t) d x=s_{0} \int_{\Omega_{1}} \partial_{x}^{2} \phi_{1}\left|Z_{1}(x, t)\right|^{2} d x-s_{0} \partial_{x} \phi_{1}(0)\left|Z_{1}(0, t)\right|^{2} .
$$

Therefore, we can rewrite the equation 4.3.50) as follows:

$$
\begin{aligned}
& -\sigma_{1}^{2} \int_{\Omega_{1}}\left(\partial_{x}+s_{0} \partial_{x} \phi_{1}\right)^{2} Z_{1}(x, t) Z_{1}(x, t) d x \\
= & \sigma_{1}^{2} \int_{\Omega_{1}}\left|\partial_{x} Z_{1}(x, t)\right|^{2} d x-\sigma_{1}^{2} s_{0}^{2} \int_{\Omega_{1}}\left|\partial_{x} \phi_{1}\right|^{2}\left|Z_{1}(x, t)\right|^{2} d x-m_{1} \sigma_{1} s_{0}\left|Z_{1}(0, t)\right|^{2} \\
& -\sigma_{1}^{2} Z_{1}(0, t) \partial_{x} Z_{1}(0, t) .
\end{aligned}
$$

Thus,

$$
\begin{align*}
& \quad \frac{d}{d t} \int_{\Omega_{1}}\left|Z_{1}(x, t)\right|^{2} d x+\sigma_{1} \int_{\Omega_{1}}\left|\partial_{x} Z_{1}(x, t)\right|^{2} d x \\
& =  \tag{4.3.51}\\
& =\sigma_{1} s_{0}^{2} \int_{\Omega_{1}}\left|\partial_{x} \phi_{1}\right|^{2}\left|Z_{1}(x, t)\right|^{2} d x+\int_{\Omega_{1}} F_{1}(x, t) Z_{1}(x, t) d x+m_{1} \sigma_{1} s_{0}\left|Z_{1}(0, t)\right|^{2} \\
& \\
& \quad+\sigma_{1}^{2} Z_{1}(0, t) \partial_{x} Z_{1}(0, t) .
\end{align*}
$$

On the other hand, multiplying by $\hat{Z}_{2}$ the second equation of 4.2.5 and integrating on $\Omega_{2}$, we have:

$$
\begin{equation*}
\frac{d}{d t} \int_{\Omega_{2}}\left|Z_{2}(t)\right| d x-\sigma_{2}^{2} \int_{\Omega_{2}}\left(\partial_{x}+s_{0} \partial_{x} \phi_{2}\right)^{2} Z_{2}(t) Z_{2}(t) d x=\int_{\Omega_{2}} F_{2}(t) Z_{2}(t) d x \tag{4.3.52}
\end{equation*}
$$

Integration by parts yields

$$
\begin{align*}
& -\int_{\Omega_{2}}\left(\partial_{x}+s_{0} \partial_{x} \phi_{2}\right)^{2} Z_{2}(x, t) \cdot Z_{2}(x, t) d x \\
= & \int_{\Omega_{2}}\left|\partial_{x} Z_{2}(x, t)\right|^{2} d x-s_{0}^{2} \int_{\Omega_{2}}\left(\partial_{x} \phi_{2}\right)^{2}\left|Z_{2}(x, t)\right|^{2} d x+\partial_{x} Z_{2}(0, t) Z_{2}(0, t)  \tag{4.3.53}\\
& -s\left(t_{0}\right) \partial_{x} \phi_{2}\left(L_{2}\right)\left|Z_{2}\left(L_{2}, t\right)\right|^{2}+\partial_{x} \phi_{2}\left|Z_{2}(0, t)\right|^{2}-\partial_{x} Z_{2}\left(L_{2}, t\right) .
\end{align*}
$$

Moreover, from the boundary condition 4.2 .5$)_{6}$, we can assert that

$$
\begin{equation*}
-\partial_{x} Z_{2}\left(L_{2}, t\right) Z_{2}\left(L_{2}, t\right)=s_{0} \partial_{x} \phi_{2}\left(L_{2}\right)\left|Z_{2}\left(L_{2}, t\right)\right|^{2}-J\left(L_{2}, t\right) Z_{2}\left(L_{2}, t\right), \quad \forall t \in(0, T) . \tag{4.3.54}
\end{equation*}
$$

Furthermore, applying the boundary conditions 4.2.5 ${ }_{3}$ and 4.2.5 4

$$
\begin{align*}
& -\sigma_{2}^{2} \partial_{x} Z_{2}(0, t) Z_{2}(0, t) \\
= & -\sigma_{1}^{2} \partial_{x} Z_{1}(0, t) Z_{1}(0, t)-\sigma_{2}^{2}\left(\sigma_{1}^{2} \partial_{x} \phi_{1}(0)-\sigma_{2} \partial_{x} \phi_{2}(0)\right) s_{0}\left|Z_{1}(0, t)\right|^{2}  \tag{4.3.55}\\
& -\sigma_{2}^{2} H(0, t) Z_{2}(0, t) .
\end{align*}
$$

Combining (4.3.52) with (4.3.53), (4.3.54) and (4.3.55), gives:

$$
\begin{align*}
& \frac{d}{d t} \int_{\Omega_{2}}\left|Z_{2}(x, t)\right|^{2} d x+\sigma_{2}^{2} \int_{\Omega_{2}}\left|\partial_{x} Z_{2}(x, t)\right|^{2} d x \\
= & \sigma_{2}^{2} s_{0}^{2} \int_{\Omega_{2}}\left|\partial_{x} \phi_{2}\right|^{2}\left|Z_{2}(x, t)\right|^{2} d x+\int_{\Omega_{2}} F(x, t) Z_{2}(x, t) d x-\sigma_{1}^{2} \partial_{x} Z_{1}(0, t) Z_{1}(0, t)  \tag{4.3.56}\\
& -\left(m_{1} \sigma_{1}-M_{2} \sigma_{2}\right) s_{0}\left|Z_{2}(0, t)\right|^{2}-\sigma_{2}^{2} J(t) Z_{2}\left(L_{2}, t\right) .
\end{align*}
$$

Combining the inequality (4.3.51) with 4.3.56) and integrating on $(0, T)$, we obtain:

$$
\begin{aligned}
& \sigma_{1}^{2} s_{0} \int_{0}^{T} \int_{\Omega_{1}}\left|\partial_{x} Z_{1}(x, t)\right|^{2} d x d t+\sigma_{2}^{2} s_{0} \int_{0}^{T} \int_{\Omega_{1}}\left|\partial_{x} Z_{2}(x, t)\right|^{2} d x d t \\
\leq & C s_{0}^{3} \int_{0}^{T} \int_{\Omega_{1}}\left|Z_{1}(x, t)\right|^{2} d x+C s_{0}^{3} \int_{0}^{T} \int_{\Omega_{2}}\left|Z_{2}(x, t)\right|^{2} d x d t+C \int_{0}^{T} \int_{\Omega_{1}}\left|F_{1}(x, t)\right|^{2} d x d t \\
& C s_{0} \int_{0}^{T}\left|Z_{1}(0, t)\right|^{2} d x+C s_{0} \int_{0}^{T}\left|Z_{2}\left(L_{2}, t\right)\right|^{2} d t+C s_{0} \int_{0}^{T}|H|^{2} d t \\
& +C \int_{0}^{T}|J|^{2} d t .
\end{aligned}
$$

Consequently, combining the inequality above with (4.3.48), we have

$$
\begin{align*}
& s_{0} \int_{0}^{T} \int_{\Omega_{1}}\left|\partial_{x} Z_{1}(x, t)\right|^{2} d x d t+s_{0} \int_{0}^{T} \int_{\Omega_{2}}\left|\partial_{x} Z_{2}(x, t)\right|^{2} d x d t \\
\leq & C \int_{0}^{T} \int_{\Omega_{1}}\left|F_{1}(x, t)\right|^{2} d x d t+C \int_{0}^{T} \int_{\Omega_{2}}\left|F_{2}(x, t)\right|^{2} d x d t+C s_{0} \int_{0}^{T}\left|G_{N}\right|^{2} d t+C s_{0} \int_{0}^{T}|H|^{2} d t \\
& +C s_{0} \int_{0}^{T}|J|^{2} d t \tag{4.3.57}
\end{align*}
$$

which is the desired conclusion.
Finally, adding the inequalities (4.3.48) and (4.3.57) yields

$$
\begin{aligned}
& s_{0}^{3} \int_{0}^{T} \int_{\Omega_{1}}\left|Z_{1}(x, t)\right|^{2} d x d t+s_{0}^{3} \int_{0}^{T} \int_{\Omega_{2}}\left|Z_{2}(x, t)\right|^{2} d x d t+s_{0} \int_{0}^{T}\left|Z_{1}\left(-L_{1}, t\right)\right|^{2} d t \\
& \quad+s_{0} \int_{0}^{T}\left|Z_{2}\left(L_{2}, t\right)\right|^{2} d t+s_{0} \int_{0}^{T} \int_{\Omega_{1}}\left|\partial_{x} Z_{1}(x, t)\right|^{2} d x d t+s_{0} \int_{0}^{T} \int_{\Omega_{2}}\left|\partial_{x} Z_{2}(x, t)\right|^{2} d x d t \\
& \leq C \int_{0}^{T} \int_{\Omega_{1}}\left|F_{1}(x, t)\right|^{2} d x d t+C \int_{0}^{T} \int_{\Omega_{2}}\left|F_{2}(x, t)\right|^{2} d x d t+C s_{0} \int_{0}^{T}\left|G_{N}\right|^{2} d t+C s_{0} \int_{0}^{T}|H|^{2} d t \\
& \quad+C s_{0} \int_{0}^{T}|J|^{2} d t
\end{aligned}
$$

and the proof of the Lemma 4.3 is complete.

### 4.4 Proof of the Corollary 4.2

This section is devoted to proof the Corollary 4.2. In order to do that, let $\left(y_{1}, y_{2}\right)$ be the solution of

$$
\begin{cases}\partial_{t} y_{1}-\sigma_{1}^{2} \partial_{x}^{2} y_{1}=0, & \forall(x, t) \in \Omega_{1} \times(0, T),  \tag{4.4.1}\\ \partial_{t} y_{2}-\sigma_{1}^{2} \partial_{x}^{2} y_{2}=0, & \forall(x, t) \in \Omega_{2} \times(0, T), \\ y_{1}(x, 0)=y_{1}^{0}(x), & \forall x \in \Omega_{1}, \\ y_{2}(x, 0)=y_{2}^{0}(x), & \forall x \in \Omega_{2}, \\ y_{1}\left(-L_{1}, t\right)=v(t), & \forall t \in(0, T), \\ \partial_{x} y_{2}\left(L_{2}, t\right)=0, & \forall t \in(0, T),\end{cases}
$$

where $\sigma_{1}, \sigma_{2}>0, y_{j}^{0} \in L^{2}\left(\Omega_{j}\right), j=1,2$ and $v \in L^{2}(0, T)$.
For an easier comprehesion, we divide the proof into three steps:

## - Step 1: Duality

Let us consider the following adjoint system:

$$
\begin{cases}-\partial_{t} w_{1}-\sigma_{1}^{2} \partial_{x}^{2} w_{1}=0, & \forall(x, t) \in \Omega_{1} \times(0, T),  \tag{4.4.2}\\ -\partial_{t} w_{2}-\partial_{x}^{2} \partial_{x}^{2} w_{2}=0, & \forall(x, t) \in \Omega_{2} \times(0, T), \\ w_{1}(x, T)=w_{1}^{T}(x), & \forall x \in \Omega_{1}, \\ w_{1}(x, T)=w_{2}^{T}(x), & \forall x \in \Omega_{2}, \\ w_{1}\left(-L_{1}, t\right)=0, & \forall t \in(0, T), \\ \partial_{x} w_{2}\left(L_{2}, t\right)=0, & \forall t \in(0, T)\end{cases}
$$

It is clear that the null controllability of the system (4.4.1) is equivalent to proof the so-called observability inequality of the adjoint system (4.4.2): There exists a constant $C>0$ such that each solution of (4.4.2) satisfy

$$
\begin{equation*}
\int_{\Omega_{1}}\left|w_{1}(x, 0)\right|^{2} d x+\int_{\Omega_{2}}\left|w_{2}(x, 0)\right|^{2} d x \leq C \int_{0}^{T}\left|\partial_{x} w_{1}\left(-L_{1}, t\right)\right|^{2} d t \tag{4.4.3}
\end{equation*}
$$

Thus, we restrict to attention to prove the inequality (4.4.3). In order to do that, we will use the Carleman estimate of the Theorem 4.1.

## - Step 2: Applying the Carleman estimate

We apply the Theorem 4.1 to the system 4.4.2, i.e., there exist constants $C>0$ and $s_{*}>0$ such that for each $s \geq s_{*}$, each solution of 4.4.2 satisfy

$$
\begin{align*}
& s^{3} \int_{0}^{T} \int_{\Omega_{1}} e^{-2 s \varphi_{1}} \theta^{3 \alpha-2 \beta}\left|w_{1}(x, t)\right|^{2} d x d t+s^{3} \int_{0}^{T} \int_{\Omega_{2}} e^{-2 s \varphi_{2}} \theta^{3 \alpha-2 \beta}\left|w_{2}(x, t)\right|^{2} d x d t \\
\leq & C s \int_{0}^{T} e^{-2 s \varphi_{1}\left(-L_{1}, t\right)} \theta^{\alpha-\beta}\left|\partial_{x} w_{1}\left(-L_{1}, t,\right)\right|^{2} d t, \tag{4.4.4}
\end{align*}
$$

where we use $\partial_{t} w_{j}+\sigma_{j}^{2} \partial_{x}^{2} w_{j}=0$, for each $j=1,2$. Since the Carleman weights $\varphi_{1}$ and $\varphi_{2}$ are bounded, it is easy to check that

$$
\begin{align*}
& C \int_{T / 4}^{3 T / 4} \int_{\Omega_{1}}\left|w_{1}(x, t)\right|^{2} d x d t+C \int_{T / 4}^{3 T / 4} \int_{\Omega_{2}}\left|w_{2}(x, t)\right|^{2} d x d t  \tag{4.4.5}\\
\leq & s_{*}^{3} \int_{\Omega_{1}} e^{-2 s_{*} \varphi_{1}} \theta^{3 \alpha-2 \beta}\left|w_{1}(x, t)\right|^{2} d x+s_{*}^{3} \int_{\Omega_{1}} e^{-2 s_{*} \varphi_{1}}\left|w_{2}(x, t)\right|^{2} d x d x
\end{align*}
$$

and

$$
\begin{equation*}
s_{*} \int_{0}^{T} e^{-2 s \varphi_{1}\left(-L_{1}\right)} \theta^{\alpha-\beta}\left|\partial_{x} w_{1}\left(-L_{1}\right), t\right|^{2} \leq C \int_{0}^{T}\left|\partial_{x} w_{1}\left(-L_{1}, t\right)\right|^{2} d t \tag{4.4.6}
\end{equation*}
$$

Combining (4.4.5) and (4.4.6) with (4.4.6) we obtain

$$
\begin{equation*}
C \int_{T / 4}^{3 T / 4} \int_{\Omega_{1}}\left|w_{1}(x, t)\right|^{2} d x d t+C \int_{T / 4}^{3 T / 4} \int_{\Omega_{2}}\left|w_{2}(x, t)\right| d x d t \leq C \int_{0}^{T}\left|\partial_{x} w_{1}\left(-L_{1}, t\right)\right|^{2} d t \tag{4.4.7}
\end{equation*}
$$

## - Step 3: Observability inequality

Multiplying the equation (4.4.2 ${ }_{j}$ by $w_{j}$ for $j=1,2$, we obtain

$$
\begin{align*}
& -\frac{1}{2} \frac{d}{d t} \int_{\Omega_{1}}\left|w_{1}(x, t)\right|^{2} d x-\frac{1}{2} \frac{d}{d t} \int_{\Omega_{2}}\left|w_{2}(x, t)\right|^{2} d x+\sigma_{1}^{2} \int_{\Omega_{1}}\left|\partial_{x} w_{1}(x, t)\right|^{2} d x  \tag{4.4.8}\\
& +\sigma_{2}^{2} \int_{\Omega_{2}}\left|\partial_{x} w_{2}(x, t)\right|^{2} d x=0
\end{align*}
$$

where we used integration by parts. Now, integrating on $(0, \tilde{t}), \tilde{t} \in(0, T)$, we have

$$
\begin{aligned}
\int_{\Omega_{1}}\left|w_{1}(x, 0)\right|^{2} d x+\int_{\Omega_{2}}\left|w_{2}(x, 0)\right|^{2} d x= & \int_{\Omega_{1}}\left|w_{1}(\tilde{t})\right|^{2} d x+\int_{\Omega_{2}}\left|w_{2}(\tilde{t})\right|^{2} d x \\
& -\sigma_{1}^{2} \int_{0}^{\tilde{t}} \int_{\Omega_{1}}\left|\partial_{x} w_{1}\right|^{2} d x d t-\sigma_{2}^{2} \int_{0}^{\tilde{t}} \int_{\Omega_{2}}\left|\partial_{x} w_{2}\right|^{2} d x d t .
\end{aligned}
$$

Integrating on $(T / 4,3 T / 4)$ on $\tilde{t}$, we see that

$$
\begin{align*}
& \int_{\Omega_{1}}\left|w_{1}(x, 0)\right|^{2} d x+\int_{\Omega_{2}}\left|w_{2}(x, 0)\right|^{2} d x  \tag{4.4.9}\\
\leq & C \int_{T / 4}^{3 T / 4} \int_{\Omega_{1}}\left|w_{1}(x, t)\right|^{2} d x d t+C \int_{T / 4}^{3 T / 4} \int_{\Omega_{1}}\left|w_{2}(x, t)\right|^{2} d x d t \tag{4.4.10}
\end{align*}
$$

Thus, combining 4.4.7 with 4.4.9, we obtain

$$
\begin{equation*}
\int_{\Omega_{1}}\left|w_{1}(x, 0)\right|^{2} d x+\int_{\Omega_{2}}\left|w_{2}(x, 0)\right|^{2} d x \leq C \int_{0}^{T}\left|\partial_{x} w_{1}\left(-L_{1}, t\right)\right|^{2} d t \tag{4.4.11}
\end{equation*}
$$

and the proof of observability inequality is complete. Hence, system (4.4.1) is nullcontrollable.

## Conclusions

In this thesis we have obtained theoretical results about inverse and control problems on some hyperbolic and parabolic problems. In particular, we have focused in wave systems with potential in cascade and heat equation with dynamic boundary conditions. We conclude this thesis with some final remarks and perspectives related to these subjects.

In Chapter 2 we studied the simultaneous potential reconstruction for a hyperbolic system in cascade when some components of the system are not accesible. Specifically, we analyzed this inverse problem where we cannot get any measurements on the last component. Our results are based on a suitable Carleman estimate on a hyperbolic system with measurements of all components except the last one. Then, we have adapted the Bukhgeim-Klibanov method to get a Lipschitz stability result for this inverse problem.

First of all, concerning the special structure of the cascade system we considered in this study, notice that in (2.3.7), the source terms $f_{j}, \ldots, f_{n}$ arise in the estimate of $F_{j}$, for each $j=1, \ldots, n$, because of the cascade structure of system (2.1.1) and the Carleman estimate of Proposition 2.6, see also Remark 2.7. This is the main difficulty to recover the potentials $\left(q_{1}, \ldots, q_{n}\right)$ with less components of (2.1.1). Then, the stability of the inverse problem treated in this thesis with two or more inaccessible components is open.

Regarding relationships of the present work with controllability, let us notice that in the particular case of $h_{j}=0$ for each $j=1, \ldots, n$ in (2.2.15) and under strong assumptions on the regularity of the solutions of (2.2.14), one can obtain a Carleman inequality of (2.2.14) with internal measurements of the first component of the system. To be more precise, for each $j=1, \ldots, n$, we define $\alpha_{j}$ such that

$$
\left\{\begin{array}{l}
\alpha_{j+1}+1<\alpha_{j}<\alpha_{j+1}+2, \quad j=1, \ldots, n-1 \\
\alpha_{n}=0
\end{array}\right.
$$

Then, there exist two constants $C=C\left(\Omega, \omega, T, x_{0}\right)>0$ and $s_{0} \geq 1$ such that for all $s \geq s_{0}$, the following inequality holds:

$$
\sum_{j=1}^{n} I\left(\alpha_{j}, v_{j}, \Omega\right) \leq C \sum_{j=0}^{2^{n-1}} s^{\beta_{j}} \int_{-T}^{T} \int_{\omega_{\omega}} e^{2 s \varphi}\left|\partial_{t}^{j} v_{1}\right|^{2} d x d t
$$

for each solution of system (2.2.14) and for some positive constants $\beta_{j}, j=0, \ldots, 2^{n-1}$. In principle, this would allow to construct a control that would require stronger regularity assumptions.

Finally, let us remark that a slight change in the proof of Proposition 2.6 shows that

$$
\begin{equation*}
I(0, v, \Omega) \leq C \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \varphi}|\square v+p v|^{2} d x d t+C s \int_{-T}^{T} \int_{\omega_{2}} e^{2 s \varphi}|\nabla v|^{2} d x d t \tag{4.4.12}
\end{equation*}
$$

for all $v \in L^{2}\left(-T, T ; H_{0}^{1}(\Omega)\right)$ such that $\square v+p v \in L^{2}(\Omega \times(-T, T)), \partial_{\nu} v \in L^{2}(\partial \Omega \times(-T, T))$ and $v( \pm T)=0$ in $\Omega$. The main ingredient of the proof are the part b) of Lemma 2.4 and the weighted Poincaré inequality (see [11]). Under that form, estimate 4.4.12) can be used in the study of wave systems with first order coupling terms.

In Chapter 3, we studied the null controllability for a suitable class of parabolic equations with dynamic boundary conditions. The main result is based on the proof of the observability inequality for the associated adjoint system. In order to get it, we used a suitable Carleman estimate for a heat equation with dynamic boundary conditions.

Moreover, we present other results based on the fact that parabolic equations with this kind of boundary conditions can be viewed as a limit of heat equations with discontinuous diffusion coefficients.

The results presented in this chapter can be extended naturally to higher dimensions. Indeed, let $d \geq 1$ and set $\Omega \subset \mathbb{R}^{d}$ be an open set with smooth boundary. In addition, let us consider $\Gamma \subset \partial \Omega$ be a nonempty open subset. Let $\left(u, u_{\Gamma}\right) \in L^{2}(\Omega \times(0, T)) \times L^{2}(\Gamma \times(0, T))$ be a solution of

$$
\begin{cases}\partial_{t} u-\Delta u=\chi_{\omega} v & \text { in } \Omega \times(0, T),  \tag{4.4.13}\\ \left(u(\cdot, 0), u_{\Gamma}(\cdot, 0)\right)=\left(u_{0}, u_{\Gamma, 0}\right), & \text { in } \Omega \times \Gamma, \\ u_{\Gamma}=u, & \text { on } \Gamma \times(0, T), \\ u=0, & \text { on }(\partial \Omega \backslash \bar{\Gamma}) \times(0, T), \\ \partial_{t} u+\partial_{\nu} u=0, & \text { on } \Gamma \times(0, T),\end{cases}
$$

where $\partial_{\nu}$ denotes the outward normal derivative and $\omega \subset \Omega$. Then, one can formulate the problem of null controllability for system (4.4.13) for any time $T>0$, i.e., given $T>0$ and $\left(u_{0}, u_{0, \Gamma}\right) \in L^{2}(\Omega) \times L^{2}(\Gamma)$, there exists a control $v \in L^{2}(\omega \times(0, T))$ such that the associated solution of 4.4.13) fulfills

$$
u(T)=0, \text { in } \bar{\Omega}
$$

Then, following the approach given in Section 3.3 we have to prove the observability inequality associated to the adjoint system of (4.4.13) by using a suitable Carleman estimate. In this context, one can use weight functions which satisfy similar estimates as in (3.3.22) and (3.3.23). In particular, when $\Omega$ has radial symmetry we shall consider explicit weight functions based on $\psi$ and $\theta$ used in Section 3 .

As we mentioned in Section 3.1, parabolic equations with discontinuous diffusion coefficients can be used to approximate parabolic equations with dynamic boundary conditions. Additionally, it is well known that for each $K>0, u_{0} \in L^{2}(\Omega)$ and $T>0$, 3.1.4 is nullcontrollable at time $T>0$, see for example [38],[19],[20] and [21]. However, the constant $C>0$ appeared in the observability inequality depends (for instance) on the diffusivity parameter, and therefore if we adapt these settings to problem (3.1.4), the observability
constant may depends on $K \geq 1$. This means that the sequence of controls $\left(v^{K}\right)_{K>0}$ in (3.1.4) may not be uniformly bounded in $L^{2}(\omega \times(0, T))$.

In order to avoid this difficulty, one can build a Carleman estimate where the weight functions depends on $K>0$. In the following, we present a result in this direction. From now on, for $\alpha \geq 1$ and $K \geq 1, \psi: \bar{\Omega} \subset \mathbb{R} \rightarrow \mathbb{R}$ and $\theta:(0, T) \subset \mathbb{R} \rightarrow \mathbb{R}$ denotes the functions given by

$$
\begin{aligned}
\psi_{L}(x) & =-\frac{1}{4 L_{1}} x^{2}+x+2 L_{1}, \quad \forall x \in \overline{\Omega_{L}} \\
\psi_{R}(x) & =-\frac{1}{4 K^{2}} x^{2}+\frac{1}{K^{2}} x+2 L_{1}, \quad \forall x \in \overline{\Omega_{R}} \\
\theta(t) & =(t(T-t))^{-\alpha}, \quad \forall t \in(0, T)
\end{aligned}
$$

We point out that $\psi \geq 2 L_{1}$ in $\Omega$ and

$$
\psi_{R}(0)=\psi_{L}(0), \quad K^{2} \psi_{R}^{\prime}(0)=\psi_{L}^{\prime}(0)
$$

i.e., $\psi$ satisfies the same transmission conditions of $u^{K}$ in (3.4.1 $)_{5}$ and (3.4.1) $)_{6}$ across the interface $x=0$. Moreover, notice that $\psi_{L}$ stands for the same weight function used in Section 3.3.

Then, we have the following result:
Lemma 4.13 Let $\alpha \geq 1,0<K_{0} \leq K$, $T>0$, define $\varphi=\psi \theta$, with $\psi$ and $\theta$ defined as above and $\sigma^{K}$ given by (3.1.3). Then, there exists two positive constants $C=C(\alpha, \Omega, T)$ and $s^{*}=s^{*}(\alpha, \Omega, T)$ independent of $K$ such that for all $s \geq s^{*}$ we have

$$
\begin{align*}
& s^{3} \int_{0}^{T} \int_{\Omega} e^{-2 s \varphi}\left(\sigma^{K}\right)^{-1}(t(T-t))^{-3 \alpha}|y|^{2} d x d t+s \int_{0}^{T} \int_{\Omega} e^{-2 s \varphi} \sigma^{K}(t(T-t))^{-\alpha}\left|\partial_{x} y\right|^{2} d x d t \\
& +s^{3} \int_{0}^{T} e^{-2 s \varphi(0, t)}(t(T-t))^{-3 \alpha}|y(0, t)|^{2} d t+s \int_{0}^{T} e^{-2 s \varphi(0, t)}(t(T-t))^{-\alpha}\left|\partial_{x} y(0, t)\right|^{2} d t \\
& +K^{-2} s^{3} \int_{0}^{T} e^{-2 s \varphi\left(L_{1}, t\right)}(t(T-t))^{-3 \alpha}\left|y\left(-L_{1}, t\right)\right|^{2} d t \leq C \int_{0}^{T} \int_{\Omega} e^{-2 s \varphi}\left|\partial_{t} y+A^{K} y\right|^{2} d x d t \\
& +C s \int_{0}^{T} e^{-2 s \varphi\left(-L_{1}, t\right)}\left|\partial_{x} y\left(-L_{1}, t\right)\right|^{2} d t+C s^{2} \int_{0}^{T} e^{-2 s \varphi\left(L_{2}, t\right)}(t(T-t))^{-2 \alpha-1}|y(1, t)|^{2} d t, \tag{4.4.14}
\end{align*}
$$

for all $y \in H^{1}\left(0, T ; L^{2}(\Omega)\right) \cap L^{2}\left(0, T ; D\left(A^{K}\right)\right)$, where $A^{K}$ is defined by

$$
A^{K} y=\partial_{x}\left(\sigma^{K} \partial_{x} y\right)
$$

with domain

$$
D\left(A^{K}\right)=\left\{y \in H_{L}^{1}(\Omega) ; \sigma^{K} y^{\prime} \in H^{1}(\Omega)\right\}
$$

The proof is based on the classical approach of Carleman estimates introduced in the context of parabolic equations by O. Imanuvilov et al. We emphasize that, for our purposes, the main difficulty here is to track the dependence on $K \geq 1$ of the constant
$C$ in (4.4.14). In fact, in this context we cannot absorb the last term of the right-hand side of (4.4.14). Thus, the question of uniform controllability for parabolic problems in the form (3.1.4) remains open.

Inspired in the ideas of [9] (see also [67] and the references therein) problems in the form (3.1.1) can be viewed as limit of another class of parabolic problems. In order to get an idea, for simplicity we set $I=(0,1)$ and denote $x^{*}=1-K^{-1}$ with $K>0$. Then, we define the subsets

$$
I_{L}=\left(0, x^{*}\right) \text { and } I_{R}=\left(x^{*}, 1\right) .
$$

For $K>0$, let us consider the following problem

$$
\begin{cases}\left(1+(K-1) \chi_{I_{L}}\right) \partial_{t} u^{K}-\partial_{x}^{2} u^{K}=f^{K}, & \forall(x, t) \in I \times(0, T),  \tag{4.4.15}\\ u^{K}(x, 0)=u_{0}, & \forall x \in I, \\ u^{K}(0, t)=\partial_{x} u^{K}(1, t)=0, & \forall t \in(0, T) .\end{cases}
$$

On the other hand, we introduce the problem

$$
\begin{cases}\partial_{t} y-\partial_{x}^{2} y=g, & \forall(x, t) \in I \times(0, T),  \tag{4.4.16}\\ y(x, 0)=y_{0}(x), & \forall x \in I, \\ y(0, t)=0, & \forall t \in(0, T), \\ \partial_{t} y(1, t)+\partial_{x} y(1, t)=0, & \forall t \in(0, T),\end{cases}
$$

with $g \in L^{2}(I \times(0, T))$ and $y_{0} \in H^{1}(I)$. Then, we have the following result
Lemma 4.14 Let $0<K_{0} \leq K$ and $u_{0} \in H^{1}(\Omega)$. Suppose that

$$
f^{K} \rightharpoonup f \text { weakly in } L^{2}(I \times(0, T)) .
$$

Then, there exists a subsequence $\left(u^{K}\right)_{K>0}$ of solutions for the problem 4.4.15 which converges to $u$ in the following way

$$
u^{K} \rightharpoonup u \text { weakly in } L^{2}\left(0, T ; H^{2}(I)\right) \cap H^{1}\left(0, T ; L^{2}(I)\right) .
$$

Moreover, $u$ is a strong solution of (4.4.16) with $g=f$ and $y_{0}=u_{0}$.
Of course, all the above questions in the context of controllability can be considered for 4.4.15.

In Chapter 4, the null controllability of heat equation with discontinuous diffusion coefficients was studied. Following the arguments presented above, the idea is to prove the observability inequality for the associated adjoint system. This was done for a suitable Carleman estimate for this kind of problems. The novelty is based on the combination of microlocal analysis ideas with localization in time functions.

In this sense, we believe that the proof of this Carleman estimate allow us to deduce some insights about these systems. The original idea was to use this kind of estimates to
prove uniform observability results for systems like (3.2). However, some difficulties on the proof suggest that we might use localization in time functions which also depends on space. This is a work in progress.

## Appendix A

## Carleman estimate for heat equation with dynamic boundary conditions by using Classical weights

## A. 1 Introduction and main result

The goal is to prove a Carleman estimate which allow us to prove the following Observability inequality

$$
\begin{equation*}
\|z(\cdot, 0)\|_{L^{2}\left(\Omega_{L}\right)}^{2}+\left|z_{\Gamma}(0)\right|^{2} \leq C \int_{0}^{T} \int_{\omega}|z|^{2} d x d t \tag{A.1.1}
\end{equation*}
$$

for each $\left(z_{T}, z_{T, \Gamma}\right) \in L^{2}\left(\Omega_{L}\right) \times \mathbb{R}$, where $\left(z, z_{\Gamma}\right) \in L^{2}\left(\Omega_{L} \times(0, T)\right) \times L^{2}(0, T)$ is a solution of the adjoint system

$$
\begin{cases}\partial_{t} z(x, t)+\partial_{x}^{2} z(x, t)=0, & \forall(x, t) \in \Omega_{L} \times(0, T),  \tag{A.1.2}\\ \left(z(x, T), z_{\Gamma}(T)\right)=\left(z_{T}\left(x, z_{T, \Gamma}\right)\right), & x \in \Omega_{L}, \\ z\left(-L_{1}, t\right)=0, & \forall t \in(0, T), \\ z_{\Gamma}^{\prime}(t)-\partial_{x} z(0, t)=0, & \forall t \in(0, T) .\end{cases}
$$

In order to do that, we will consider the Classical weight functions introduced by A. Fursikov and O. Imanuvilov. We recall that these ones are based on an auxiliary function whose existence is given by the following result:

Lemma A. 1 Given nonempty open set $\omega \subset \subset \Omega_{L}$, there is a function $\eta_{0} \in C^{2}\left(\overline{\Omega_{L}}\right)$ such that

$$
\eta_{0}>0 \text { in } \Omega_{L}, \quad \eta_{0}\left(-L_{1}\right)=\eta_{0}(0)=0, \quad\left|\eta_{0}^{\prime}\right|>0 \text { in } \overline{\Omega_{L} \backslash \omega} .
$$

We notice that the functions given by the above lemma fulfills

$$
\eta_{0}^{\prime}\left(-L_{1}\right)>0, \quad \eta_{0}^{\prime}(0)<0 .
$$

From now on, we fix $\omega^{\prime} \subset \subset \Omega_{L}, \lambda, m>1$ and $\eta_{0}$ as in the previous lemma. We define the weight functions $\alpha$ and $\eta$ by

$$
\begin{align*}
& \alpha(x, t)=(t(T-t))^{-1}\left(e^{2 \lambda m\left\|\eta_{0}\right\|_{\infty}}-e^{\lambda\left(m\left\|\eta_{0}\right\|_{\infty}+\eta_{0}(x)\right)}\right),  \tag{A.1.3}\\
& \eta(x, t)=(t(T-t))^{-1} e^{\lambda\left(m\left\|\eta_{0}\right\|_{\infty}+\eta_{0}(x)\right)}, \tag{A.1.4}
\end{align*}
$$

for each $(x, t) \in \overline{\Omega_{L}} \times(0, T)$. Now we have all the ingredients to state the Carleman estimate for heat equation with dynamic boundary conditions:

Theorem A. 2 Let $T>0, \omega \subset \subset \Omega_{L}$ be a nonempty and open interval. In addition, we choose $\omega^{\prime} \subset \subset \omega$. Define $\alpha, \eta_{0}, \xi$ as above with respect to $\omega^{\prime}$. Then, there exists constants $C>0, \lambda_{1} \geq 1$ and $s_{1} \geq 1$ such that the following inequality holds

$$
\begin{aligned}
& s^{3} \lambda^{4} \int_{0}^{T} \int_{\Omega_{L}} e^{2 s \alpha} \xi^{3}|\varphi|^{2} d x d t+s \lambda \int_{0}^{T} \int_{\Omega_{L}} e^{2 s \varphi} \xi\left|\partial_{x} \varphi\right|^{2} d x d t \\
& +s^{-1} \int_{0}^{T} \int_{\Omega_{L}} e^{2 s \alpha} \xi^{-1}\left|\partial_{t} \varphi\right|^{2} d x d t+s^{-1} \int_{0}^{T} \int_{\Omega_{L}} e^{2 s \alpha} \xi^{-1}\left|\partial_{x}^{2} \varphi\right|^{2} d x d t \\
& +s^{3} \lambda^{3} \int_{0}^{T} e^{2 s \alpha(0, t)} \xi^{3}(0, t)|\varphi(0, t)|^{2} d t+s \lambda \int_{0}^{T} e^{2 s \alpha(0, t)}\left|\partial_{x} \varphi(0, t)\right|^{2} d t \\
& +\int_{0}^{T} e^{2 s \alpha(0, t)}\left|\partial_{t} \varphi(0, t)\right|^{2} d t+s \lambda \int_{0}^{T} e^{2 s \alpha\left(-L_{1}, t\right)} \xi\left(-L_{1}, t\right)\left|\partial_{x} \varphi\left(-L_{1}, t\right)\right|^{2} d t \\
& \leq C \int_{0}^{T} \int_{\Omega_{L}} e^{2 s \alpha}\left|\partial_{t} \varphi+\partial_{x}^{2} \varphi\right|^{2} d x d t+C \int_{0}^{T} e^{2 s \alpha(0, t)}\left|\partial_{t} \varphi(0, t)-\partial_{x} \varphi(0, t)\right|^{2} d t \\
& +C s^{3} \lambda^{4} \int_{0}^{T} \int_{\omega} e^{2 s \alpha} \xi^{3}|\varphi|^{2} d x d t,
\end{aligned}
$$

for all $\lambda \geq \lambda_{1}$ and $s \geq s_{1}$ and for all $\varphi \in C^{2}(\bar{\Omega} \times[0, T])$.
The rest of this appendix is devoted to prove the above Theorem.

## A. 2 Proof of the Carleman estimate

Let $\varphi \in C^{\infty}\left(\overline{\Omega_{L}} \times[0, T]\right), \lambda \geq 1$ and $s \geq s \geq 1$ be given. Define

$$
\psi=e^{-s \varphi} \varphi, \quad f=e^{-s \alpha}\left(\partial_{t} \varphi+\partial_{x}^{2} \varphi\right), \quad g=e^{-s \alpha}\left(\partial_{t} \varphi-\partial_{x} \varphi\right), \quad \forall(x, t) \in \overline{\Omega_{L}} \times(0, T) .
$$

Direct computations show that

$$
\begin{array}{r}
e^{-s \alpha} \partial_{t} \varphi=\partial_{t} \varphi+s \partial_{t} \alpha \psi, \quad e^{-s \alpha} \partial_{x} \varphi=\partial_{x} \psi+s \partial_{x} \alpha \psi, \\
e^{-s \alpha} \partial_{x}^{2} \varphi=\partial_{x}^{2} \psi+2 s \partial_{x} \alpha \partial_{x} \psi+s^{2}\left|\partial_{x} \alpha\right|^{2} \psi+s \partial_{x}^{2} \alpha \psi .
\end{array}
$$

In the following, we shall use the abbreviations

$$
\begin{aligned}
M_{1} \psi=s^{2}\left|\partial_{x} \alpha\right|^{2} \psi+\partial_{x}^{2} \psi+s \partial_{t} \alpha \psi, & M_{2} \psi=s \partial_{x}^{2} \alpha+\partial_{t} \psi+2 s \partial_{x} \alpha \partial_{x} \psi \\
N_{1} \psi=s \partial_{t} \alpha \psi-\partial_{x} \psi, & N_{2} \psi=\partial_{t} \psi-s \partial_{x} \alpha \psi .
\end{aligned}
$$

Then, according to this notation, it is clear that $\psi$ satisfies the following equations:

$$
\begin{equation*}
M_{1} \psi+M_{2} \psi=f, \quad \text { and } N_{1} \psi+N_{2} \psi=g, \text { in } \overline{\Omega_{L}} \times(0, T) . \tag{A.2.1}
\end{equation*}
$$

Applying $\|\cdot\|_{L^{2}\left(\Omega_{L} \times(0, T)\right)}$ and $\|\cdot\|_{L^{2}(0, T)}$ to the equations A.2.1) we get

$$
\begin{array}{r}
\left\|M_{1} \psi\right\|_{L^{2}\left(\Omega_{L} \times(0, T)\right)}^{2}+\left\|M_{2} \psi\right\|_{L^{2}\left(\Omega_{L} \times(0, T)\right)}^{2}+\left\|N_{1} \psi(0, \cdot)\right\|_{L^{2}(0, T)}^{2}+\left\|N_{2} \psi(0, \cdot)\right\|_{L^{2}(0, T)}^{2} \\
\left\langle M_{1} \psi, M_{2} \psi\right\rangle_{L^{2}\left(\Omega_{L} \times(0, T)\right)}+\left\langle N_{1} \psi(0, \cdot), N_{2} \psi(0, \cdot)\right\rangle_{L^{2}(0, T)}=\|f\|_{L^{2}\left(\Omega_{L} \times(0, T)\right)}^{2}+\|g(0, \cdot)\|_{L^{2}(0, T)}^{2} \tag{A.2.2}
\end{array}
$$

Our next task is to compute the inner products in the right-hand side of A.2.2. In order to do that, we shall use the notation

$$
\left\langle M_{1} \psi, M_{2} \psi\right\rangle_{L^{2}\left(\Omega_{L} \times(0, T)\right)}=\sum_{j, k=1}^{3} I_{j, k},
$$

where $I_{j k}$ stands for the scalar product in $L^{2}\left(\Omega_{L} \times(0, T)\right)$ between the $j^{\text {th }}$ term of $M_{1} \psi$ and the $k^{\text {th }}$ term of $M_{2} \psi$. Then, $I_{11}$ reads as follows

$$
I_{11}=s^{3} \int_{0}^{T} \int_{\Omega_{L}}\left|\partial_{x} \alpha\right|^{2} \partial_{x}^{2} \alpha|\psi|^{2} d x d t
$$

On the other hand, using the identity $\frac{1}{2} \psi \partial_{t} \psi=\partial_{t}\left(|\psi|^{2}\right)$ in $\Omega_{L} \times(0, T)$, we have

$$
I_{12}=s^{2} \int_{0}^{T} \int_{\Omega_{L}}\left|\partial_{x} \alpha\right|^{2} \psi \partial_{t} \psi d x d t=-\int_{0}^{T} \int_{\Omega_{L}} \partial_{x} \alpha \partial_{t} \partial_{x} \alpha|\psi|^{2} d x d t
$$

where we are used the fact that $\alpha$ blows up as $t \rightarrow 0^{+}$and $t \rightarrow T-$. Moreover, integration by parts shows that

$$
\begin{aligned}
I_{13} & =2 s^{3} \int_{0}^{T} \int_{\Omega_{L}}\left|\partial_{x} \alpha\right|^{3} \psi \partial_{x} \psi d x d t \\
& =-3 s^{3} \int_{0}^{T} \int_{\Omega_{L}}\left|\partial_{x} \alpha\right|^{2} \partial_{x}^{2} \alpha|\psi|^{2} d x d t+s^{3} \int_{0}^{T}\left|\partial_{x} \alpha(0, t)\right|^{3}|\psi(0, t)|^{2} d t
\end{aligned}
$$

where we used $\psi\left(-L_{1}, t\right)=0$ for all $t \in(0, T)$. In the same way, the term $I_{21}$ can be estimated as follows

$$
\begin{aligned}
I_{21} & =-s \int_{0}^{T} \int_{\Omega_{L}} \partial_{x}^{2} \alpha \psi \partial_{x}^{2} \psi d x d t \\
& =-s \int_{0}^{T} \int_{\Omega_{L}} \partial_{x}^{3} \alpha \psi \partial_{x} \psi d x d t-s \int_{0}^{T} \int_{\Omega_{L}} \partial_{x}^{2} \alpha\left|\partial_{x} \psi\right|^{2} d x d t+s \int_{0}^{T} \partial_{x}^{2} \alpha(0, t) \psi(0, t) \partial_{x} \psi(0, t) d t
\end{aligned}
$$

In addition, the term $I_{22}$ reads as follows:

$$
\begin{aligned}
I_{22} & =\int_{0}^{T} \int_{\Omega_{L}} \partial_{x}^{2} \psi \partial_{t} \psi d x d t \\
& =-\int_{0}^{T} \int_{\Omega_{L}} \partial_{t}\left(\left|\partial_{x} \psi\right|^{2}\right) d x d t+\int_{0}^{T} \partial_{x} \psi(0, t) \partial_{t} \psi(0, t) d t
\end{aligned}
$$

Notice that the first term in the above equality is given by

$$
-\frac{1}{2} \int_{0}^{T} \int_{\Omega_{L}} \partial_{t}\left(\left|\partial_{x} \psi\right|^{2}\right) d x d t=0
$$

On the other hand, using the equations (A.2.1) for $g$ we have

$$
\begin{aligned}
I_{22}= & \int_{0}^{T} \partial_{x} \psi(0, t) \partial_{t} \psi(0, t) d t \\
= & \int_{0}^{T}\left|\partial_{t} \psi(0, t)\right|^{2} d t+s \int_{0}^{T} \partial_{t} \alpha(0, t) \psi(0, t) \partial_{x} \psi(0, t) d t \\
& -s \int_{0}^{T} \partial_{x} \alpha(0, t) \psi(0, t) \partial_{x} \psi(0, t) d t-\int_{0}^{T} \partial_{x} \psi(0, t) g(0, t) d t .
\end{aligned}
$$

Moreover, $I_{23}$ can be estimated in the following way

$$
\begin{aligned}
I_{23}= & 2 s \int_{0}^{T} \int_{\Omega_{L}} \partial_{x} \alpha \partial_{x} \psi \partial_{x}^{2} \psi d x d t \\
= & -s \int_{0}^{T} \int_{\Omega_{L}} \partial_{x}^{2} \alpha\left|\partial_{x} \psi\right|^{2} d x d t+s \int_{0}^{T} \partial_{x} \alpha(0, t)\left|\partial_{x} \psi(0, t)\right|^{2} d t \\
& -s \int_{0}^{T} \partial_{x} \alpha\left(-L_{1}, t\right)\left|\partial_{x} \psi\left(-L_{1}, t\right)\right|^{2} d t .
\end{aligned}
$$

By definition, $I_{31}$ reads as follows

$$
I_{31}=s^{2} \int_{0}^{T} \int_{\Omega_{L}} \partial_{t} \alpha \partial_{x}^{2} \alpha|\psi|^{2} d x d t
$$

Once again, since $\alpha$ blows up as $t \rightarrow 0^{+}$and $t \rightarrow T^{-}$we get

$$
I_{32}=s \int_{0}^{T} \int_{\Omega_{L}} \partial_{t} \alpha \psi \partial_{t} \psi d x d t=-\frac{1}{2} s \int_{0}^{T} \int_{\Omega_{L}} \partial_{t}^{2} \alpha|\psi|^{2} d x d t
$$

Finally, $I_{33}$ is given by

$$
\begin{aligned}
I_{33} & =2 s^{2} \int_{0}^{T} \int_{\Omega_{L}} \partial_{x} \alpha \partial_{t} \alpha \psi \partial_{x} \psi d x d t \\
& =-s^{2} \int_{0}^{T} \int_{\Omega_{L}} \partial_{x}\left(\partial_{x} \alpha \partial_{t} \alpha\right)|\psi|^{2} d x d t+s^{2} \int_{0}^{T} \partial_{x} \alpha(0, t) \partial_{t} \alpha(0, t)|\psi(0, t)|^{2} d t .
\end{aligned}
$$

Gathering all the terms we have

$$
\begin{align*}
& \left\langle M_{1} \psi, M_{2} \psi\right\rangle_{L^{2}\left(\Omega_{L} \times(0, T)\right)} \\
= & -2 s^{3} \int_{0}^{T} \int_{\Omega_{L}}\left|\partial_{x} \alpha\right|^{2} \partial_{x}^{2} \alpha|\psi|^{2} d x d t-2 s \int_{0}^{T} \int_{\Omega_{L}} \partial_{x}^{2} \alpha\left|\partial_{x} \psi\right|^{2} d x d t \\
& +s^{3} \int_{0}^{T}\left|\partial_{x} \alpha(0, t)\right|^{3}|\psi(0, t)|^{2} d t+s \int_{0}^{T} \partial_{x} \alpha(0, t)\left|\partial_{x} \psi(0, t)\right|^{2} d t \\
& +\int_{0}^{T}\left|\partial_{t} \psi(0, t)\right|^{2} d t-s \int_{0}^{T} \partial_{x} \alpha\left(-L_{1}, t\right)\left|\partial_{x} \psi\left(-L_{1}, t\right)\right|^{2} d t \\
& -s \int_{0}^{T} \int_{\Omega_{L}} \partial_{x}^{3} \alpha \psi \partial_{x} \psi d x d t+s \int_{0}^{T} \partial_{x}^{2} \alpha(0, t) \psi(0, t) \partial_{x} \psi(0, t) d t  \tag{A.2.3}\\
& +s \int_{0}^{T} \partial_{t} \alpha(0, t) \psi(0, t) \partial_{x} \psi(0, t) d t-s \int_{0}^{T} \partial_{x} \alpha(0, t) \psi(0, t) \partial_{x} \psi(0, t) d t \\
& -\int_{0}^{T} g(0, t) \partial_{x} \psi(0, t) d t+s^{2} \int_{0}^{T} \int_{\Omega_{L}} \partial_{t} \alpha \partial_{x}^{2} \alpha|\psi|^{2} d x d t \\
& -\frac{1}{2} s \int_{0}^{T} \int_{\Omega_{L}} \partial_{t}^{2} \alpha|\psi|^{2} d x d t+s^{2} \int_{0}^{T} \partial_{x} \alpha(0, t) \partial_{t} \alpha(0, t)|\psi(0, t)|^{2} d t .
\end{align*}
$$

Similar computations shows that the second inner product of A.2.2) is given by

$$
\begin{align*}
& \left\langle N_{1} \psi(0, t), N_{2} \psi(0, t)\right\rangle_{L^{2}(0, T)} \\
= & s \int_{0}^{T} \partial_{t} \alpha(0, t) \psi(0, t) \partial_{t} \psi(0, t) d t-s^{2} \int_{0}^{T} \partial_{t} \alpha(0, t) \partial_{x} \alpha(0, t)|\psi(0, t)|^{2} d t  \tag{A.2.4}\\
& -\int_{0}^{T} \partial_{x} \psi(0, t) \partial_{t} \psi(0, t) d t+s \int_{0}^{T} \partial_{x} \alpha(0, t) \psi(0, t) \partial_{x} \psi(0, t) d t .
\end{align*}
$$

Now we focus on some estimates on weight functions. According to the definitions of $\alpha$ and $\eta$, we get

$$
\begin{equation*}
\left|\partial_{t} \alpha(x, t)\right| \leq C(t(T-t))^{-1} \xi(x, t), \quad\left|\partial_{t}^{2} \alpha(x, t)\right| \leq C(t(T-t))^{-2} \xi(x, t) \tag{A.2.5}
\end{equation*}
$$

for each $(x, t) \in \overline{\Omega_{L}} \times(0, T)$ and for some constant $C$ dependent of $T$ but independent of $\lambda, m$ and $s$. On the other hand, a direct computations on spatial derivatives of $\alpha$ gives

$$
\begin{equation*}
\partial_{x} \alpha(x, t)=-\lambda \eta_{0}^{\prime}(x) \xi(x, t), \quad \partial_{x}^{2} \alpha(x, t)=-\lambda\left(\eta_{0}^{\prime \prime}(x)+\lambda\left|\eta_{0}^{\prime}(x)\right|^{2}\right) \xi(x, t), \tag{A.2.6}
\end{equation*}
$$

for each $(x, t) \in \overline{\Omega_{L}} \times(0, T)$. We point out that the second derivative of $\alpha$ can be bounded by below in the following way

$$
\begin{equation*}
\partial_{x}^{2} \alpha(x, t) \geq-\lambda^{2}\left|\eta_{0}^{\prime}(x)\right|^{2} \xi(x, t), \quad \forall(x, t) \in \overline{\Omega_{L}} \times(0, T) \tag{A.2.7}
\end{equation*}
$$

Then, substituting (A.2.3) and A.2.4 into A.2.2) and by using the estimates A.2.5,
(A.2.6) and A.2.7) we obtain

$$
\begin{align*}
& \left\|M_{1} \psi\right\|_{L^{2}\left(\Omega_{L} \times(0, T)\right)}^{2}+\left\|M_{2} \psi\right\|_{L^{2}\left(\Omega_{L} \times(0, T)\right)}^{2}+\left\|N_{1} \psi(0, t)\right\|_{L^{2}(0, T)}^{2}+\left\|N_{2} \psi(0, T)\right\|_{L^{2}(0, T)}^{2} \\
& \quad+s^{3} \lambda^{4} \int_{0}^{T} \int_{\Omega_{L}} \xi^{3}|\psi|^{2} d x d t+s \lambda \int_{0}^{T} \int_{\Omega_{L}} \xi\left|\partial_{x} \psi\right|^{2} d x d t \\
& \quad+s^{3} \lambda^{3} \int_{0}^{T} \int_{\Omega_{L}} \xi^{3}(0, t)|\psi(0, t)|^{2} d x d t+s \lambda \int_{0}^{T} \xi(0, t)\left|\partial_{x} \psi(0, t)\right|^{2} d t \\
& \quad+\int_{0}^{T}\left|\partial_{t} \psi(0, t)\right|^{2} d t+s \lambda \int_{0}^{T} \xi\left(-L_{1}, t\right)\left|\partial_{x} \psi\left(-L_{1}, t\right)\right|^{2} d t \\
& \leq C_{1} \int_{0}^{T} \int_{\Omega_{L}}|f|^{2} d x d t+C_{1} \int_{0}^{T}|g(0, t)|^{2} d t+C_{1} s^{3} \lambda^{4} \int_{0}^{T} \int_{\omega^{\prime}} \xi|\psi|^{2} d x d t \\
& \quad+C_{1} s \lambda \int_{0}^{T} \int_{\omega^{\prime}} \xi\left|\partial_{x} \psi\right|^{2} d x d t+X+Y \tag{A.2.8}
\end{align*}
$$

where $X$ and $Y$ are defined by

$$
\begin{aligned}
X= & C_{1} s^{2} \lambda^{2} \int_{0}^{T} \int_{\Omega_{L}}(t(T-t))^{-1} \xi^{2}|\psi|^{2} d x d t+C_{1} s^{2} \lambda^{2} \int_{0}^{T} \int_{\Omega_{L}}(t(T-t))^{-2} \xi^{2}|\psi|^{2} d x d t \\
& +C_{1} s \lambda^{3} \int_{0}^{T} \int_{\Omega_{L}} \xi\left|\psi \| \partial_{x} \psi\right| d x d t
\end{aligned}
$$

and

$$
\begin{aligned}
Y= & C_{1} s \lambda^{2} \int_{0}^{T} \xi(0, t)|\psi(0, t)|\left|\partial_{x} \psi(0, t)\right| d t+C_{1} s \int_{0}^{T}(t(T-t))^{-1} \xi(0, t)\left|\psi(0, t) \| \partial_{x} \psi(0, t)\right| d t \\
& +C_{1} s \lambda \int_{0}^{T} \xi(0, t)\left|\psi(0, t)\left\|\partial_{x} \psi(0, t)\left|d t+C_{1} \int_{0}^{T}\right| g(0, t)\right\| \partial_{x} \psi(0, t)\right| d t \\
& +C_{1} s^{2} \lambda \int_{0}^{T}(t(T-t))^{-1} \xi^{2}(0, t)|\psi(0, t)|^{2} d t+C_{1} \int_{0}^{T}\left|\partial_{x} \psi(0, t) \| \partial_{t} \psi(0, t)\right| d t \\
& +C_{1} s \int_{0}^{T}(t(T-t))^{-1} \xi(0, t)\left|\psi(0, t)\left\|\partial_{t} \psi(0, t)\left|d t+C_{1} s \lambda \int_{0}^{T} \xi(0, t)\right| \psi(0, t)\right\| \partial_{x} \psi(0, t)\right| d t
\end{aligned}
$$

Now, it is clear that there exists $\lambda_{1} \geq 1, s_{1} \geq 1$ such that for all $\lambda \geq \lambda_{1}$ and $s \geq s_{1}$ we have the following estimates:

$$
\begin{equation*}
X \leq \frac{1}{2} s^{3} \lambda^{4} \int_{0}^{T} \int_{\Omega_{L}} \xi^{3}|\psi|^{2} d x d t+\frac{1}{2} s \lambda \int_{0}^{T} \int_{\Omega_{L}} \xi\left|\partial_{x} \psi\right|^{2} d x d t \tag{A.2.9}
\end{equation*}
$$

and

$$
\begin{align*}
Y \leq & \frac{1}{2} s^{3} \lambda^{3} \int_{0}^{T} \int_{\Omega_{L}} \xi^{3}(0, t)|\psi(0, t)|^{2} d x d t+\frac{1}{2} s \lambda \int_{0}^{T} \xi(0, t)\left|\partial_{x} \psi(0, t)\right|^{2} d t  \tag{A.2.10}\\
& +\frac{1}{2} \int_{0}^{T}\left|\partial_{t} \psi(0, t)\right|^{2} d t+\frac{1}{2} C_{1}^{2} \int_{0}^{T}|g(0, t)|^{2} d t
\end{align*}
$$

Then, using (A.2.9) and A.2.10) in A.2.8) we get

$$
\begin{aligned}
& \left\|M_{1} \psi\right\|_{L^{2}\left(\Omega_{L} \times(0, T)\right)}^{2}+\left\|M_{2} \psi\right\|_{L^{2}\left(\Omega_{L} \times(0, T)\right)}^{2}+\left\|N_{1} \psi(0, \cdot)\right\|_{L^{2}(0, T)}^{2}+\left\|N_{2} \psi(0, \cdot)\right\|_{L^{2}(0, T)}^{2} \\
& \quad+s^{3} \lambda^{4} \int_{0}^{T} \int_{\Omega_{L}} \xi^{3}|\psi|^{2} d x d t+s \lambda \int_{0}^{T} \int_{\Omega_{L}} \xi\left|\partial_{x} \psi\right|^{2} d x d t \\
& \quad+s^{3} \lambda^{3} \int_{0}^{T} \int_{\Omega_{L}} \xi^{3}(0, t)|\psi(0, t)|^{2} d x d t+s \lambda \int_{0}^{T} \xi(0, t)\left|\partial_{x} \psi(0, t)\right|^{2} d t \\
& \quad+\int_{0}^{T}\left|\partial_{t} \psi(0, t)\right|^{2} d t+s \lambda \int_{0}^{T} \xi\left(-L_{1}, t\right)\left|\partial_{x} \psi\left(-L_{1}, t\right)\right|^{2} d t \\
& \leq C_{2} \int_{0}^{T} \int_{\Omega_{L}}|f|^{2} d x d t+C_{2} \int_{0}^{T}|g(0, t)|^{2} d t+C_{2} s^{3} \lambda^{4} \int_{0}^{T} \int_{\omega^{\prime}} \xi^{3}|\psi|^{2} d x d t,
\end{aligned}
$$

where the local term of $\partial_{x} \psi$ in $\omega^{\prime}$ can be absorbed as in [48]. Moreover, global terms of $\partial_{t} \psi$ and $\partial_{x}^{2} \psi$ can be obtained by using the equations A.2.1 and Young's inequality. Finally, we come back on the original variables and the proof is complete.

## Bibliography

[1] Fatiha Alabau-Boussouira. A two-level energy method for indirect boundary observability and controllability of weakly coupled hyperbolic systems. SIAM J. Control Optim., 42(3):871-906, 2003.
[2] Fatiha Alabau-Boussouira. Controllability of cascade coupled systems of multidimensional evolution PDEs by a reduced number of controls. C. R. Math. Acad. Sci. Paris, 350(11-12):577-582, 2012.
[3] Fatiha Alabau-Boussouira. A hierarchic multi-level energy method for the control of bidiagonal and mixed $n$-coupled cascade systems of PDE's by a reduced number of controls. Adv. Differential Equations, 18(11-12):1005-1072, 2013.
[4] Fatiha Alabau-Boussouira. Insensitizing exact controls for the scalar wave equation and exact controllability of 2 -coupled cascade systems of PDE's by a single control. Math. Control Signals Systems, 26(1):1-46, 2014.
[5] Fatiha Alabau-Boussouira, Roger Brockett, Olivier Glass, Jérôme Le Rousseau, and Enrique Zuazua. Control of Partial Differential Equations: Cetraro, Italy 2010, Editors: Piermarco Cannarsa, Jean-Michel Coron, volume 2048. Springer, 2012.
[6] Fatiha Alabau-Boussouira, Piermarco Cannarsa, and Masahiro Yamamoto. Source reconstruction by partial measurements for a class of hyperbolic systems in cascade. In Mathematical paradigms of climate science, volume 15 of Springer INdAM Ser., pages 35-50. Springer, [Cham], 2016.
[7] Fatiha Alabau-Boussouira and Matthieu Léautaud. Indirect controllability of locally coupled wave-type systems and applications. J. Math. Pures Appl. (9), 99(5):544576, 2013.
[8] Farid Ammar-Khodja, Assia Benabdallah, Manuel González-Burgos, and Luz De Teresa. Recent results on the controllability of linear coupled parabolic problems: a survey. Math. Control Relat. Fields, 1(3):267-306, 2011.
[9] José M Arrieta, Pavol Quittner, Aníbal Rodríguez-Bernal, et al. Parabolic problems with nonlinear dynamical boundary conditions and singular initial data. DIFFERENTIAL AND INTEGRAL EQUATIONS-ATHENS-, 14(12):1487-1510, 2001.
[10] Claude Bardos, Gilles Lebeau, and Jeffrey Rauch. Sharp sufficient conditions for the observation, control, and stabilization of waves from the boundary. SIAM journal on
control and optimization, 30(5):1024-1065, 1992.
[11] Lucie Baudouin, Maya De Buhan, and Sylvain Ervedoza. Global Carleman estimates for waves and applications. Comm. Partial Differential Equations, 38(5):823-859, 2013.
[12] Larisa Beilina, Michel Cristofol, Shumin Li, and Masahiro Yamamoto. Lipschitz stability for an inverse hyperbolic problem of determining two coefficients by a finite number of observations. Inverse Problems, 2017.
[13] Ioan Bejenaru, Jesus Ildefonso Diaz, and Ioan I Vrabie. An abstract approximate controllability result and applications to elliptic and parabolic systems with dynamic boundary conditions. Electronic Journal of Differential Equations (EJDE)[electronic only], 2001:Paper-No, 2001.
[14] Mourad Bellassoued. Global logarithmic stability in inverse hyperbolic problem by arbitrary boundary observation. Inverse Problems, 20(4):1033-1052, 2004.
[15] Mourad Bellassoued and Masahiro Yamamoto. Logarithmic stability in determination of a coefficient in an acoustic equation by arbitrary boundary observation. J. Math. Pures Appl. (9), 85(2):193-224, 2006.
[16] Mourad Bellassoued and Masahiro Yamamoto. Carleman estimates and applications to inverse problems for hyperbolic systems. Springer, 2017.
[17] Assia Benabdallah, Michel Cristofol, Patricia Gaitan, and Luz de Teresa. A new Carleman inequality for parabolic systems with a single observation and applications. C. R. Math. Acad. Sci. Paris, 348(1-2):25-29, 2010.
[18] Assia Benabdallah, Michel Cristofol, Patricia Gaitan, and Masahiro Yamamoto. Inverse problem for a parabolic system with two components by measurements of one component. Appl. Anal., 88(5):683-709, 2009.
[19] Assia Benabdallah, Yves Dermenjian, and Jérôme Le Rousseau. Carleman estimates for the one-dimensional heat equation with a discontinuous coefficient, and applications. Comptes Rendus Mécanique, 334(10):582-586, 2006.
[20] Assia Benabdallah, Yves Dermenjian, and Jérôme Le Rousseau. Carleman estimates for the one-dimensional heat equation with a discontinuous coefficient and applications to controllability and an inverse problem. Journal of Mathematical Analysis and applications, 336(2):865-887, 2007.
[21] Assia Benabdallah, Yves Dermenjian, and Jérôme Le Rousseau. Carleman estimates for stratified media. Journal of Functional Analysis, 260(12):3645-3677, 2011.
[22] Assia Benabdallah, Yves Dermenjian, and Laetitia Thevenet. Carleman estimates for some non-smooth anisotropic media. Communications in Partial Differential Equations, 38(10):1763-1790, 2013.
[23] Olivier Bodart and Caroline Fabre. Controls insensitizing the norm of the solution
of a semilinear heat equation. J. Math. Anal. Appl., 195(3):658-683, 1995.
[24] Muriel Boulakia and Axel Osses. Local null controllability of a two-dimensional fluid-structure interaction problem. ESAIM: Control, Optimisation and Calculus of Variations, 14(1):1-42, 2008.
[25] Haim Brezis. Functional analysis, Sobolev spaces and partial differential equations. Springer Science \& Business Media, 2010.
[26] Aleksandr L'vovich Bukhgeim and Mikhail Viktorovich Klibanov. Global uniqueness of a class of multidimensional inverse problems. In Doklady Akademii Nauk, volume 260, pages 269-272. Russian Academy of Sciences, 1981.
[27] N Burq. Contrôle de l'équation des ondes dans des ouverts peu réguliers. Asymptotic Analysis, 14(2):157-91, 1997.
[28] N. Carreño, S. Guerrero, and M. Gueye. Insensitizing controls with two vanishing components for the three-dimensional Boussinesq system. ESAIM Control Optim. Calc. Var., 21(1):73-100, 2015.
[29] Nicolás Carreño, Roberto Morales, and Axel Osses. Potential reconstruction for a class of hyperbolic systems from incomplete measurements. Inverse Problems, 34(8):085005, 2018.
[30] Thierry Cazenave and Alain Haraux. An introduction to semilinear evolution equations, volume 13. Oxford University Press on Demand, 1998.
[31] Giuseppe M Coclite, Angelo Favini, Ciprian G Gal, Gisèle Ruiz Goldstein, Jerome A Goldstein, Enrico Obrecht, and Silvia Romanelli. The role of wentzell boundary conditions in linear and nonlinear analysis. Tübinger Berichte, 132, 2008.
[32] Jean-Michel Coron. Control and nonlinearity. Number 136. American Mathematical Soc., 2007.
[33] Michel Cristofol, Patricia Gaitan, Kati Niinimäki, and Olivier Poisson. Inverse problem for a coupled parabolic system with discontinuous conductivities: onedimensional case. Inverse Probl. Imaging, 7(1):159-182, 2013.
[34] Michel Cristofol, Patricia Gaitan, and Hichem Ramoul. Inverse problems for a $2 \times 2$ reaction-diffusion system using a Carleman estimate with one observation. Inverse Problems, 22(5):1561-1573, 2006.
[35] Michel Cristofol, Patricia Gaitan, Hichem Ramoul, and Masahiro Yamamoto. Identification of two coefficients with data of one component for a nonlinear parabolic system. Appl. Anal., 91(11):2073-2081, 2012.
[36] René Dáger. Insensitizing controls for the 1-D wave equation. SIAM J. Control Optim., 45(5):1758-1768, 2006.
[37] Luz de Teresa. Insensitizing controls for a semilinear heat equation. Comm. Partial Differential Equations, 25(1-2):39-72, 2000.
[38] A Doubova, E Fernández-Cara, and M González-Burgos. On the controllability of the heat equation with nonlinear boundary fourier conditions. Journal of Differential Equations, 196(2):385-417, 2004.
[39] Anna Doubova, A Osses, and J-P Puel. Exact controllability to trajectories for semilinear heat equations with discontinuous diffusion coefficients. ESAIM: Control, Optimisation and Calculus of Variations, 8:621-661, 2002.
[40] Klaus-J Engel. Positivity and stability for one-sided coupled operator matrices. Positivity, 1(2):103-124, 1997.
[41] Lawrence C Evans. Partial differential equations. 2010.
[42] Hector O Fattorini and David L Russell. Exact controllability theorems for linear parabolic equations in one space dimension. Archive for Rational Mechanics and Analysis, 43(4):272-292, 1971.
[43] Hector O Fattorini and David L Russell. Uniform bounds on biorthogonal functions for real exponentials with an application to the control theory of parabolic equations. Quarterly of Applied Mathematics, 32(1):45-69, 1974.
[44] Angelo Favini, Gisele Ruiz Goldstein, Jerome A Goldstein, and Silvia Romanelli. The heat equation with generalized wentzell boundary condition. Journal of evolution equations, 2(1):1-19, 2002.
[45] Enrique Fernández-Cara, Manuel González-Burgos, and Luz de Teresa. Controllability of linear and semilinear non-diagonalizable parabolic systems. ESAIM Control Optim. Calc. Var., 21(4):1178-1204, 2015.
[46] Enrique Fernández-Cara, Manuel González-Burgos, Sergio Guerrero, and Jean-Pierre Puel. Exact controllability to the trajectories of the heat equation with fourier boundary conditions: the semilinear case. ESAIM: Control, Optimisation and Calculus of Variations, 12(3):466-483, 2006.
[47] Enrique Fernández-Cara, Manuel González-Burgos, Sergio Guerrero, and Jean-Pierre Puel. Null controllability of the heat equation with boundary fourier conditions: the linear case. ESAIM: Control, Optimisation and Calculus of Variations, 12(3):442465, 2006.
[48] Enrique Fernández-Cara and Sergio Guerrero. Global carleman inequalities for parabolic systems and applications to controllability. SIAM journal on control and optimization, 45(4):1395-1446, 2006.
[49] Enrique Fernández-Cara and Enrique Zuazua. On the null controllability of the onedimensional heat equation with bv coefficients. Computational and Applied Mathematics, 12:167-190, 2002.
[50] Andrej Vladimirovič Fursikov and O Yu Imanuvilov. Controllability of evolution equations. Number 34. Seoul National University, 1996.
[51] Patricia Gaitan and Hadjer Ouzzane. Stability result for two coefficients in a coupled hyperbolic-parabolic system. J. Inverse Ill-Posed Probl., 25(3):265-286, 2017.
[52] Ciprian G Gal. On a class of degenerate parabolic equations with dynamic boundary conditions. Journal of Differential Equations, 253(1):126-166, 2012.
[53] Ciprian G Gal. The role of surface diffusion in dynamic boundary conditions: Where do we stand? Milan Journal of Mathematics, 83(2):237-278, 2015.
[54] Ciprian G Gal and Maurizio Grasselli. The non-isothermal allen-cahn equation with dynamic boundary conditions. Discrete Contin. Dyn. Syst, 22(4):1009-1040, 2008.
[55] Ciprian G Gal, Mahamadi Warma, et al. Well posedness and the global attractor of some quasi-linear parabolic equations with nonlinear dynamic boundary conditions. Differential and Integral Equations, 23(3/4):327-358, 2010.
[56] Galina C García, Cristhian Montoya, and Axel Osses. A source reconstruction algorithm for the stokes system from incomplete velocity measurements. Inverse Problems, 33(10):105003, 2017.
[57] Gisele Ruiz Goldstein et al. Derivation and physical interpretation of general boundary conditions. Advances in Differential Equations, 11(4):457-480, 2006.
[58] Lop Fat Ho. Observabilité frontière de l'équation des ondes. Comptes rendus de l'Académie des sciences. Série 1, Mathématique, 302(12):443-446, 1986.
[59] D Hömberg, K Krumbiegel, and J Rehberg. Optimal control of a parabolic equation with dynamic boundary condition. Applied Mathematics \& Optimization, 67(1):3-31, 2013.
[60] Oleg Imanuvilov and Takéo Takahashi. Exact controllability of a fluid-rigid body system. Journal de mathématiques pures et appliquées, 87(4):408-437, 2007.
[61] Oleg Yu. Imanuvilov. On Carleman estimates for hyperbolic equations. Asymptot. Anal., 32(3-4):185-220, 2002.
[62] Oleg Yu. Imanuvilov and Masahiro Yamamoto. Global Lipschitz stability in an inverse hyperbolic problem by interior observations. Inverse Problems, 17(4):717728, 2001. Special issue to celebrate Pierre Sabatier's 65th birthday (Montpellier, 2000).
[63] Oleg Yu. Imanuvilov and Masahiro Yamamoto. Global uniqueness and stability in determining coefficients of wave equations. Comm. Partial Differential Equations, 26(7-8):1409-1425, 2001.
[64] Albert Edward Ingham. Some trigonometrical inequalities with applications to the theory of series. Mathematische Zeitschrift, 41(1):367-379, 1936.
[65] Victor Isakov. Inverse problems for partial differential equations, volume 127 of Applied Mathematical Sciences. Springer, New York, second edition, 2006.
[66] Daijun Jiang, Yikan Liu, and Masahiro Yamamoto. Inverse source problem for a wave equation with final observation data. In Mathematical Analysis of Continuum Mechanics and Industrial Applications: Proceedings of the International Conference CoMFoS15, pages 153-164. Springer, 2017.
[67] Ángela Jiménez-Casas and Anibal Rodriguez-Bernal. Dynamic boundary conditions as limit of singularly perturbed parabolic problems. Conference Publications, 2011.
[68] Sergey I Kabanikhin, Abdigany D Satybaev, and Maxim A Shishlenin. Direct methods of solving multidimensional inverse hyperbolic problems, volume 48. Walter de Gruyter, 2013.
[69] Michael V. Klibanov. Inverse problems and Carleman estimates. Inverse Problems, 8(4):575-596, 1992.
[70] Michael V. Klibanov and Masahiro Yamamoto. Lipschitz stability of an inverse problem for an acoustic equation. Appl. Anal., 85(5):515-538, 2006.
[71] Vilmos Komornik. Exact controllability and stabilization: the multiplier method, volume 36. Masson, 1994.
[72] Vilmos Komornik and Paola Loreti. Fourier series in control theory. Springer Science \& Business Media, 2005.
[73] Michael Kumpf and Gregor Nickel. Dynamic boundary conditions and boundary control for the one-dimensional heat equation. Journal of dynamical and control systems, 10(2):213-225, 2004.
[74] Jérôme Le Rousseau. Carleman estimates and controllability results for the onedimensional heat equation with bv coefficients. Journal of Differential Equations, 233(2):417-447, 2007.
[75] Jérôme Le Rousseau, Matthieu Léautaud, and Luc Robbiano. Controllability of a parabolic system with a diffusive interface. Journal of the European Mathematical Society, 15(4):1485-1574, 2013.
[76] Jérôme Le Rousseau and Nicolas Lerner. Carleman estimates for anisotropic elliptic operators with jumps at an interface. Analysis छ PDE, 6(7):1601-1648, 2013.
[77] GILLES Lebeau. Control for distributed systems: The microlocal approach. Current trends in applied mathematics. MA Herrero and E. Zuazua eds., Editorial Complutense, Madrid, pages 181-218, 1996.
[78] Gilles Lebeau and Luc Robbiano. Contrôle exact de l'équation de la chaleur. Communications in Partial Differential Equations, 20(1-2):335-356, 1995.
[79] J.-L. Lions. Contrôlabilité exacte, perturbations et stabilisation de systèmes distribués. Tome 1, volume 8 of Recherches en Mathématiques Appliquées [Research in Applied Mathematics]. Masson, Paris, 1988. Contrôlabilité exacte. [Exact controllability], With appendices by E. Zuazua, C. Bardos, G. Lebeau and J. Rauch.
[80] J.-L. Lions and E. Magenes. Problèmes aux limites non homogènes. II. Ann. Inst. Fourier (Grenoble), 11:137-178, 1961.
[81] Lahcen Maniar, Martin Meyries, and Roland Schnaubelt. Null controllability for parabolic equations with dynamic boundary conditions. Evolution Equations \& Control Theory, 6(3):381-407, 2017.
[82] Philippe Martin, Lionel Rosier, and Pierre Rouchon. Null controllability using flatness: A case study of a 1-d heat equation with discontinuous coefficients. In Control Conference (ECC), 2015 European, pages 55-60. IEEE, 2015.
[83] Sorin Micu and Enrique Zuazua. An introduction to the controllability of partial differential equations. Quelques questions de théorie du contrôle. Sari, T., ed., Collection Travaux en Cours Hermann, to appear, 2004.
[84] Axel Osses. Une nouvelle famille de multiplicateurs et applications à la contrôlabilité exacte de l'équation d'ondes. Comptes Rendus de l'Académie des Sciences-Series I-Mathematics, 326(9):1099-1104, 1998.
[85] Pierre-Arnaud Raviart, Jean-Marie Thomas, Philippe G Ciarlet, and Jacques Louis Lions. Introduction à l'analyse numérique des équations aux dérivées partielles, volume 2. Dunod Paris, 1998.
[86] Luc Robbiano. Théorème d'unicité adapté au contrôle des solutions des problèmes hyperboliques. Communications in partial differential equations, 16(4-5):789-800, 1991.
[87] Luc Robbiano. Fonction de coût et contrôle des solutions des équations hyperboliques. Asymptotic Analysis, 10(2):95-115, 1995.
[88] David L Russell. A unified boundary controllability theory for hyperbolic and parabolic partial differential equations. Studies in Applied Mathematics, 52(3):189211, 1973.
[89] Jacques Simon. Compact sets in the spacel p (o, t; b). Annali di Matematica pura ed applicata, 146(1):65-96, 1986.
[90] Louis Tebou. Locally distributed desensitizing controls for the wave equation. C. R. Math. Acad. Sci. Paris, 346(7-8):407-412, 2008.
[91] Emmanuel Trélat. Contrôle optimal: théorie 83 applications. Vuibert Paris, 2005.
[92] Marius Tucsnak and George Weiss. Observation and control for operator semigroups. Springer Science \& Business Media, 2009.
[93] Gunther Uhlmann. Electrical impedance tomography and calderón's problem. Inverse problems, 25(12):123011, 2009.
[94] Juan Luis Vázquez and Enzo Vitillaro. Heat equation with dynamical boundary conditions of reactive type. Communications in Partial Differential Equations, 33(4):561-612, 2008.
[95] Juan Luis Vázquez and Enzo Vitillaro. Heat equation with dynamical boundary conditions of reactive-diffusive type. Journal of Differential Equations, 250(4):21432161, 2011.
[96] Masahiro Yamamoto. Stability, reconstruction formula and regularization for an inverse source hyperbolic problem by a control method. Inverse Problems, 11(2):481496, 1995.
[97] Masahiro Yamamoto. Uniqueness and stability in multidimensional hyperbolic inverse problems. J. Math. Pures Appl. (9), 78(1):65-98, 1999.

