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Controllability and observability of an artificial advection-diffusion problem

Pierre Cornilleau*& Sergio Guerrero[†]

March 3, 2012

Abstract

In this paper we study the controllability of an artificial advection-diffusion system through the boundary. Suitable Carleman estimates give us the observability of the adjoint system in the one dimensional case. We also study some basic properties of our problem such as backward uniqueness and we get an intuitive result on the control cost for vanishing viscosity.

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Introduction

Artificial advection-diffusion problem In the present paper we deal with an advection-diffusion problem with small viscosity truncated in one space direction. Our interest for the linear advection diffusion equation comes from the Navier-Stokes equations, but it arises also from other fields as, for example, meteorology. For a given viscosity $\varepsilon > 0$, the incompressible Navier-Stokes equations can be written as

$$\begin{cases} f_t + (f \cdot \nabla)f - \varepsilon \Delta f + \nabla p = 0\\ \operatorname{div}(f) = 0 \end{cases}$$

where f is the velocity vector field, p the pressure, ∇ the gradient and Δ the usual Laplacian. Considering the flow around a body, we have that f is almost constant far away from the body and equal to a. Then, our system can be approximated by the Oseen equation

$$\begin{cases} f_t + (a.\nabla)f - \varepsilon \Delta f + \nabla p = 0\\ \operatorname{div}(f) = 0 \end{cases}$$

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(see, for instance, [9, page 309, (1.2)]). Consequently, we get the following equation for the vorticity u = rot(f)

$$u_t + a \cdot \nabla u - \varepsilon \Delta u = 0.$$

In the sequel, we assume for simplicity that a is the nth unit vector of the canonical basis of \mathbb{R}^n .

When one computes the solution of this problem, it can only be solved numerically on a bounded domain. A good way to approximate the solution on the whole space may be given by the use of artificial boundary conditions (see [8, 9]).

In this paper, we will consider an advection-diffusion system in a strip $\Omega := \mathbb{R}^{n-1} \times (-L, 0)$ (*L* some positive constant) with particular artificial boundary conditions on both sides of the domain. The following system was considered in [8, section 6]:

$$\begin{aligned} u_t + \partial_{x_n} u - \varepsilon \Delta u &= 0 & \text{in } (0, T) \times \Omega, \\ \varepsilon(u_t + \partial_{\nu} u) &= 0 & \text{on } (0, T) \times \Gamma_0, \\ \varepsilon(u_t + \partial_{\nu} u) + u &= 0 & \text{on } (0, T) \times \Gamma_1, \\ u(0, .) &= u_0 & \text{in } \Omega, \end{aligned}$$
(1)

where T > 0, $\Gamma_0 := \mathbb{R}^{n-1} \times \{0\}$, $\Gamma_1 := \mathbb{R}^{n-1} \times \{-L\}$ and we have denoted ∂_{x_n} the partial derivative with respect to x_n and ∂_{ν} the normal derivative. In [8, Theorem 3], the author proved that the solution of (1) converges in some sense to the restriction of the solution of

$$\begin{cases} u_t + \partial_{x_n} u = 0 & \text{in } (0, T) \times \mathbb{R}^{n-1} \times \mathbb{R}_-, \\ u = 0 & \text{on } (0, T) \times \Gamma_1, \\ u(0, .) = u_0 & \text{in } \mathbb{R}^{n-1} \times \mathbb{R}_-, \end{cases}$$

to $(0,T) \times \Omega$ (see also section 6 in [8]).

Up to our knowledge, there is no controllability result concerning system (1). In this paper, we will be interested in the uniform boundary controllability of (1):

$$\left\{ \begin{array}{ll} u_t + \partial_{x_n} u - \varepsilon \Delta u = 0 & \quad \mathrm{in} \ (0,T) \times \Omega, \\ \varepsilon (u_t + \partial_{\nu} u) + u \mathbf{1}_{\Gamma_1} = v \mathbf{1}_{\Gamma_i} & \quad \mathrm{on} \ (0,T) \times \partial \Omega, \\ u(0,.) = u_0 & \quad \mathrm{in} \ \Omega, \end{array} \right.$$

where i = 0 or i = 1.

In the sequel, we shall focus on the one-dimensional problem

$$(S_v) \begin{cases} u_t + u_x - \varepsilon u_{xx} = 0 & \text{in } (0, T) \times (-L, 0), \\ \varepsilon (u_t + \partial_\nu u) = v & \text{on } (0, T) \times \{0\}, \\ \varepsilon (u_t + \partial_\nu u) + u = 0 & \text{on } (0, T) \times \{-L\}, \\ u(0, .) = u_0 & \text{in } (-L, 0). \end{cases}$$

We are interested in the so-called *null controllability* of this system

for given u_0 , find v such that the solution of (S_v) satisfies $u(T) \equiv 0$.

This property is classically equivalent to an observability inequality for the so-called *adjoint system* (see Proposition 16 below for a proof), which can be written as

$$(S') \begin{cases} \varphi_t + \varphi_x + \varepsilon \varphi_{xx} = 0 & \text{in } (0,T) \times (-L,0), \\ \varepsilon(\varphi_t - \partial_\nu \varphi) - \varphi = 0 & \text{on } (0,T) \times \{0\}, \\ \varphi_t - \partial_\nu \varphi = 0 & \text{on } (0,T) \times \{-L\}, \\ \varphi(T,.) = \varphi_T & \text{in } (-L,0). \end{cases}$$

The observability inequality corresponding to the previous controllability property is:

there exists
$$C > 0$$
 such that $\|\varphi(0, .)\|_X \le C \|\varphi(., 0)\|_{L^2(0,T)} \quad \forall \varphi_T \in X,$ (2)

where the space X will be defined below. We refer the reader unfamiliar with the links between observability and controllability to the seminal paper [4] (see also [10] on hyperbolic systems).

The phenomenon of degeneration of a parabolic problem to a hyperbolic one has been studied in several papers: see, for instance, [2] or [6] (one dimensional heat equation) and [7] (Burgers equation). Similar results of interest can also be found in [1].

$$||u||_X := \left(||u||_{L^2(\Omega)}^2 + \varepsilon ||u||_{L^2(\partial\Omega)}^2 \right)^{\frac{1}{2}}.$$

Observe that X is the largest subspace of $L^2(\Omega)$ such that its elements possess a trace in $L^2(\partial\Omega)$. In particular,

$$\forall \epsilon > 0, \ H^{1/2+\epsilon}(\Omega) \subsetneq X \subsetneq L^2(\Omega).$$

We will denote by $C_{obs}(\varepsilon)$ the cost of the null-control, which is the smallest constant C which fulfills the observability estimate (2). Our main result is the following:

Theorem 1 Assume T > 0.

• For any $\varepsilon > 0$, there exists $\overline{C} > 0$ such that the solution of problem (S') satisfies (2). Consequently, for every $u_0 \in X$, there exists a control $v \in L^2(0,T)$ with

$$\|v\|_{L^2(0,T)} \le C \|u_0\|_X$$

such that the solution of the problem (S_v) satisfies $u(T) \equiv 0$.

• Furthermore, if T/L is large enough, the cost of the null-control $C_{obs}(\varepsilon)$ tends to zero exponentially as $\varepsilon \to 0$:

 $\exists C, k > 0 \text{ such that} \quad C_{obs}(\varepsilon) \leq C e^{-k/\varepsilon} \quad \forall \varepsilon \in (0, 1).$

Remark 1 One can in fact obtain an observability result for the adjoint system at x = -L, that is

$$\|\varphi(0,.)\|_X \le C \|\varphi(.,-L)\|_{L^2(0,T)}$$

for any φ solution of (S'). This provides some controllability result for the direct system on Γ_1 : a function v can be found depending continuously on u_0 so that the solution of

$$\begin{array}{ll} u_t + u_x - \varepsilon u_{xx} = 0 & in \ (0,T) \times (-L,0), \\ u_t + \partial_{\nu} u = 0 & on \ (0,T) \times \{0\}, \\ \varepsilon (u_t + \partial_{\nu} u) + u = v & on \ (0,T) \times \{-L\}, \\ u(0,.) = u_0 & in \ (-L,0), \end{array}$$

satisfies $u(T, .) \equiv 0$. Despite the fact that it is more physical to control our system at Γ_1 , we have chosen to present here the result for the system (S_v) since its proof, less obvious, requires more computations.

Remark 2 The fact that the control cost vanishes tells intuitively that the state is almost null for T/L big enough. This is to be connected with the fact that, for $\varepsilon = 0$, the system is purely advective and then that, for T > L, its state is null.

In some context of inverse problems (be able to know the origin of a polluted river for instance), it can be interesting to know if the observation of the solution of the direct problem $(S_0)^n$ on the boundary part Γ_1 or Γ_0 can allow us to recover the initial data. The corresponding result is presented now and will be proved at the end of the first section.

Proposition 2 Let T > 0 and $\varepsilon > 0$. If the solution of (S_0) with initial data $u_0 \in X$ satisfies u = 0 on $(0,T) \times \Gamma_1$, then $u_0 \equiv 0$. However, there is no constant C > 0 such that the following estimate holds

$$\|u_0\|_X \le C \|u(., -L)\|_{L^2(0,T)} \quad \forall u_0 \in X.$$
(3)

Remark 3 Using Remark 1, we can also obtain a similar result at x = 0.

The rest of the article is organized as follows: in the first section, we show the well-posedness of the direct and the adjoint problems using some semi-groups approach. In the second section, we adapt a Carleman inequality to the case of our one-dimensional problem. The third section is intended to explain how to get observability in the one-dimensional case and the equivalence between observability and controllability.

Notations:

 $A \leq B$ means that, for some universal constant c > 0, $A \leq cB$. $A \sim B$ means that, for some universal constant c > 1, $c^{-1}B \leq A \leq cB$.

1 Well-posedness and basic properties of systems

In this section, we work in dimension n.

1.1 Homogeneous problems

We will use some semi-group results to show existence and uniqueness of the homogeneous direct problem (that is $(S_0)^n$). This will enable to define solutions of the system $(S_0)^n$ as a semigroup value. We define $H = (X, \|.\|_X), V = H^1(\Omega)$ endowed with the usual norm $\|.\|$ and we consider the bilinear form on V defined by

$$a(u_1, u_2) = \varepsilon \int_{\Omega} \nabla u_1 \cdot \nabla u_2 + \int_{\Omega} \partial_{x_n} u_1 u_2 + \int_{\Gamma_1} u_1 u_2.$$
(4)

With the help of this bilinear form, one may now consider the space

$$\mathcal{D} := \left\{ u_1 \in X; \sup_{u_2 \in \mathcal{D}(\overline{\Omega}); \ ||u_2||_X \le 1} |a(u_1, u_2)| < +\infty \right\}$$

equipped with the natural norm

$$||u_1||_{\mathcal{D}} = ||u_1||_X + \sup_{u_2 \in \mathcal{D}(\overline{\Omega}); ||u_2||_X \le 1} |a(u_1, u_2)|.$$

Note that, using an integration by parts, one shows that $a(u_1, u_2)$ is well-defined for $u_1 \in X$ and $u_2 \in \mathcal{C}^{\infty}(\overline{\Omega})$ and that the map

$$u_2 \in X \mapsto a(u_1, u_2) \in \mathbb{R}$$

is well-defined and continuous for any $u_1 \in \mathcal{D}$.

Using the Riesz representation theorem, we can define an operator \mathcal{A} with domain $\mathcal{D}(\mathcal{A}) = \mathcal{D}$ and such that

$$\forall u_1 \in \mathcal{D}(\mathcal{A}), \ \forall u_2 \in X, \quad < -\mathcal{A}u_1, u_2 >_X = a(u_1, u_2).$$

 $(S_0)^n$ might now be written in the following abstract way

$$\begin{cases} u_t = \mathcal{A}u, \\ u(0, .) = u_0 \end{cases}$$

since, for $u \in \mathcal{C}(\mathbb{R}^+, \mathcal{D}(\mathcal{A})) \cap \mathcal{C}^1(\mathbb{R}^+, X)$ a classical solution of $(S_0)^n$, we directly get that

$$\forall v \in X, \quad \langle u_t, v \rangle_X = -a(u, v).$$

Proposition 3 Let $\varepsilon > 0$. Then \mathcal{A} generates a continuous semi-group $(e^{t\mathcal{A}})_{t\geq 0}$ on X.

Proof Using the Hille-Yoshida theorem, it is sufficient to show that \mathcal{A} is a maximal monotone operator.

• First, Green-Riemann formula gives

$$<\mathcal{A}u, u>_X = -\varepsilon \int_{\Omega} |\nabla u|^2 - \frac{1}{2} \int_{\partial \Omega} |u|^2,$$

so the monotonicity is proved.

• Given $v \in X$, we have to solve (I - A)u = v, that is to find $u \in \mathcal{D}(A)$ such that

$$\forall u' \in X, < (I - \mathcal{A})u, u' >_X = < v, u' >_X.$$

The left-hand side term of this equation is a continuous bilinear form B on the space V while the right-hand side term is a continuous linear form L on V. Moreover, using (4), we can easily compute

$$B(u,u) = \|u\|_X^2 + \varepsilon \int_{\Omega} |\nabla u|^2 + \frac{1}{2} \int_{\partial \Omega} |u|^2 \ge \min\{1,\varepsilon\} \|u\|^2,$$

which show that B is coercive on V. Consequently, Lax-Milgram theorem shows that there exists $u \in V$ such that $B(u, u') = L(u'), \forall u' \in V$. Using test functions $u' \in \mathcal{D}(\Omega)$ additionally gives $u \in \mathcal{D}(\mathcal{A})$ and our proof ends.

In particular, for every initial data $u_0 \in X$, we have existence and uniqueness of a solution $u \in \mathcal{C}(\mathbb{R}^+, X)$ to $(S_0)^n$. We will call these solutions weak solutions opposed to strong solutions i.e. such that $u_0 \in \mathcal{D}(\mathcal{A})$ and which fulfill $u \in \mathcal{C}(\mathbb{R}^+, \mathcal{D}(\mathcal{A})) \cap \mathcal{C}^1(\mathbb{R}^+, X)$. One can notice that, using the density of $\mathcal{D}(\mathcal{A})$ in X, a weak solution can always be approximated by a strong solution.

Let us introduce the adjoint system in dimension n

$$(S')^n \begin{cases} \varphi_t + \partial_{x_n} \varphi + \varepsilon \Delta \varphi = 0 & \text{ in } (0, T) \times \Omega, \\ \varepsilon (\varphi_t - \partial_{\nu} \varphi) - \varphi = 0 & \text{ on } (0, T) \times \Gamma_0, \\ \varphi_t - \partial_{\nu} \varphi = 0 & \text{ on } (0, T) \times \Gamma_1, \\ \varphi (T, .) = \varphi_T & \text{ in } \Omega. \end{cases}$$

Defining the adjoint \mathcal{A}^* of \mathcal{A} , we can show in a similar way as before that $(S')^n$ may be written in the following abstract way

$$\begin{cases} \varphi_t + \mathcal{A}^* \varphi = 0, \\ \varphi(T, .) = \varphi_T. \end{cases}$$

The adjoint operator \mathcal{A}^* is also a maximal monotone operator with domain $\mathcal{D}(\mathcal{A}^*)$. Thus, for every initial data $\varphi_T \in X$, we have existence and uniqueness of a solution $\varphi \in \mathcal{C}(\mathbb{R}^+, X)$ to $(S')^n$ given by means of the backward semigroup $(e^{(T-t)\mathcal{A}^*})_{t\geq 0}$. We will also speak of *weak solutions* or *strong solutions* in this situation.

1.2 Nonhomogeneous direct problems

We consider a slightly more general system

$$(S_{f,g_0,g_1})^n \begin{cases} u_t + \partial_{x_n} u - \varepsilon \Delta u = f & \text{in } (0,T) \times \Omega, \\ \varepsilon (u_t + \partial_{\nu} u) = g_0 & \text{on } (0,T) \times \Gamma_0, \\ \varepsilon (u_t + \partial_{\nu} u) + u = g_1 & \text{on } (0,T) \times \Gamma_1, \\ u(0,.) = u_0 & \text{in } \Omega, \end{cases}$$

with $f \in L^2((0,T) \times \Omega)$, $g_0 \in L^2((0,T) \times \Gamma_0)$ and $g_1 \in L^2((0,T) \times \Gamma_1)$. We consider a function g on $\partial \Omega$ such that $g = g_i$ on Γ_i (i = 0, 1).

Definition 4 We say that $u \in \mathcal{C}([0,T], X)$ is a transposition solution of $(S_{f,g_0,g_1})^n$ if, for every function $\varphi \in \mathcal{C}([0,T], \mathcal{D}(\mathcal{A}^*)) \cap \mathcal{C}^1([0,T], X)$, the following identity holds

$$\int_0^\tau (\langle u, \varphi_t \rangle_X + \langle u, \mathcal{A}^* \varphi \rangle_X + \langle F, \varphi \rangle_X) = [\langle u(t), \varphi(t) \rangle_X]_{t=0}^{t=\tau} \quad \forall \tau \in [0, T],$$

where we have defined, using the Riesz representation theorem, $F(t) \in X$ such that

$$\langle F(t), u' \rangle_X = \int_{\Omega} f(t)u' + \int_{\partial\Omega} g(t)u', \quad \forall u' \in X.$$
 (5)

Proposition 5 Let T > 0, $u_0 \in X$, $f \in L^2((0,T) \times \Omega)$, $g_0 \in L^2((0,T) \times \Gamma_0)$ and $g_1 \in L^2((0,T) \times \Gamma_1)$. Then $(S_{f,g_0,g_1})^n$ possesses a unique solution u.

Proof

Let us first assume that u belongs to $u \in \mathcal{C}([0,T], \mathcal{D}(\mathcal{A})) \cap \mathcal{C}^1([0,T], X)$. Then, it is easy to prove that u is a solution to $(S_{f,g_0,g_1})^n$ if and only if, for any $\phi \in \mathcal{D}(\mathcal{A}^*)$,

$$< u_t, \phi >_X = < u, \mathcal{A}^* \phi >_X + < F, \phi >_X \text{ in } [0, T].$$

Consequently, thanks to Duhamel formula, u is a solution of $(S_{f,g_0,g_1})^n$ is equivalent to

$$u(t) = e^{t\mathcal{A}}u_0 + \int_0^t e^{(t-s)\mathcal{A}}F(s)ds, \quad \forall t \in [0,T],$$

since $\mathcal{D}(\mathcal{A}^*)$ is dense in X.

In order to prove the existence of a solution $u \in \mathcal{C}([0,T],X)$, we approximate (u_0, f, g) by

$$(u_0^p, f^p, g_0^p, g_1^p) \in \mathcal{D}(\mathcal{A}) \times \mathcal{C}([0, T], L^2(\Omega)) \times \mathcal{C}([0, T], L^2(\Gamma_0)) \times \mathcal{C}([0, T], L^2(\Gamma_1))$$

in

$$X \times L^{2}((0,T), L^{2}(\Omega)) \times L^{2}((0,T), L^{2}(\Gamma_{0})) \times L^{2}((0,T), L^{2}(\Gamma_{1})).$$

Then, we know that u^p given by

$$u^p(t) = e^{t\mathcal{A}}u_0^p + \int_0^t e^{(t-s)\mathcal{A}}F^p(s)ds, \quad \forall t \in [0,T],$$

where F^p is related to (f^p, g_0^p, g_1^p) by formula (5), is the solution of $(S_{f^p, g_0^p, g_1^p})^n$. Passing to the limit in this identity, we get that $u^p \to u$ in $\mathcal{C}([0, T], X)$ with

$$u(t) := e^{t\mathcal{A}}u_0 + \int_0^t e^{(t-s)\mathcal{A}}F(s)ds, \quad \forall t \in [0,T].$$

Furthermore, passing to the limit in Definition 4 for u^p , we obtain that u is a solution of $(S_{f,g_0,g_1})^n$. Moreover, this argument also proves the uniqueness of the solution.

We now focus on some regularization effect for our general system and for some technical reasons we want to have some explicit dependence of the bounds on ε . For this result, it is essential that $\varepsilon > 0$.

 $\textbf{Lemma 6} \ Let \ \varepsilon \in (0,1), \ f \in L^2((0,T) \times \Omega), \ g_0 \in L^2((0,T) \times \Gamma_0) \ and \ g_1 \in L^2((0,T) \times \Gamma_1).$

• Let $u_0 \in X$. Then if u is the solution of $(S_{f,g_0,g_1})^n$, u belongs to $L^2((0,T), H^1(\Omega)) \cap \mathcal{C}([0,T], X)$ and

$$\varepsilon^{1/2} \|u\|_{L^2((0,T),H^1(\Omega))} + \|u\|_{L^{\infty}((0,T),X)} \le C(\|u_0\|_X + \|f\|_{L^2((0,T)\times\Omega)} + \|g\|_{L^2((0,T)\times\partial\Omega)})$$
(6)

for some C > 0 only depending on T.

• Let now $u_0 \in H^1(\Omega)$. Then $u \in L^2((0,T), \mathcal{D}(\mathcal{A})) \cap \mathcal{C}([0,T], H^1(\Omega)), u_t \in L^2((0,T), X)$ and

$$\varepsilon^{1/2} \|u\|_{L^{\infty}([0,T],H^{1}(\Omega))} + \varepsilon(\|\Delta u\|_{L^{2}((0,T)\times\Omega)} + \|\partial_{\nu}u\|_{L^{2}((0,T)\times\partial\Omega)}) + \|u_{t}\|_{L^{2}([0,T],X)} \\
\leq \frac{C}{\varepsilon^{1/2}} (\|u_{0}\|_{H^{1}(\Omega)} + \|f\|_{L^{2}((0,T)\times\Omega)} + \|g\|_{L^{2}((0,T)\times\partial\Omega)}),$$
(7)

where C > 0 only depends on T.

Proof Approximating u by a regular function (in $L^2((0,T), \mathcal{D}(\overline{\Omega}))$), all computations below are justified.

• First, we multiply our main equation by u and we integrate on Ω . An application of Green-Riemann formula with use of the boundary conditions gives the identity

$$\frac{1}{2}\frac{d}{dt}\|u\|_X^2 + \frac{1}{2}\int_{\partial\Omega}|u|^2 + \varepsilon\int_{\Omega}|\nabla u|^2 = \int_{\Omega}fu + \int_{\partial\Omega}gu,$$

which immediately yields

$$\frac{d}{dt} \|u\|_X^2 + \int_{\partial\Omega} |u|^2 + 2\varepsilon \int_{\Omega} |\nabla u|^2 \le \|u\|_X^2 + \|f\|_{L^2(\Omega)}^2 + \|g\|_{L^2(\partial\Omega)}^2$$

Finally, Gronwall's lemma yields (6).

• Now, we multiply the main equation by u_t and we integrate on Ω , which, using Green-Riemann formula and the boundary conditions provide the identity

$$\|u_t\|_X^2 + \int_{\Omega} \partial_{x_n} u u_t + \frac{\varepsilon}{2} \frac{d}{dt} \left(\int_{\Omega} |\nabla u|^2 \right) + \frac{1}{2} \frac{d}{dt} \left(\int_{\Gamma_1} |u|^2 \right) = \int_{\Omega} f u_t + \int_{\partial \Omega} g u_t,$$

Now

$$-\int_{\Omega} u_t \partial_{x_n} u \leqslant \frac{1}{2} \int_{\Omega} |\nabla u|^2 + \frac{1}{2} ||u_t||_{\lambda}^2$$

and, using again Young's inequality, we have

$$\int_{\Omega} f u_t + \int_{\partial \Omega} g u_t \leq \frac{1}{4} \| u_t \|_X^2 + C(\|f\|_{L^2(\Omega)}^2 + \|g\|_{L^2(\partial \Omega)}^2).$$

We have thus proven the following inequality

$$\|u_t\|_X^2 + 2\varepsilon \frac{d}{dt} \left(\int_{\Omega} |\nabla u|^2 \right) + 2\frac{d}{dt} \int_{\Gamma_1} |u|^2 \le 2 \int_{\Omega} |\nabla u|^2 + C \left(\|g\|_{L^2(\partial\Omega)}^2 + \|f\|_{L^2(\Omega)}^2 \right).$$
(8)

We integrate here between 0 and t and estimate the first term in the right-hand side using (6) to get

$$\int_{\Omega} |\nabla u|^2 \leq \frac{C}{\varepsilon^2} \left(\int_0^T (\|f(t)\|_{L^2(\Omega)}^2 + \|g(t)\|_{L^2(\partial\Omega)}^2) dt + \|u_0\|_{H^1(\Omega)}^2 \right),$$

for $t \in (0, T)$. We now inject this into (8) to get the second part of the required result

$$\|u_t\|_{L^2((0,T);X)}^2 \leqslant \frac{C}{\varepsilon} \left(\int_0^T (\|f(t)\|_{L^2(\Omega)}^2 + \|g(t)\|_{L^2(\partial\Omega)}^2) dt + \|u_0\|_{H^1(\Omega)}^2 \right).$$

Writing the problem as

$$\left\{ \begin{array}{ll} -\Delta u = \frac{1}{\varepsilon}(f-u_t-\partial_{x_n}u) & \text{ in } (0,T)\times\Omega,\\ \partial_\nu u = \frac{g_0}{\varepsilon}-u_t & \text{ on } (0,T)\times\Gamma_0,\\ \partial_\nu u = \frac{g_1}{\varepsilon}-u_t-\frac{1}{\varepsilon}u & \text{ on } (0,T)\times\Gamma_1, \end{array} \right.$$

we get that $u \in L^2((0,T), \mathcal{D}(\mathcal{A}))$ and the existence of some constant C > 0 such that

$$\begin{split} \|\Delta u\|_{L^{2}((0,T)\times\Omega)} &+ \|\partial_{\nu}u\|_{L^{2}((0,T)\times\partial\Omega)} \\ &\leq \frac{C}{\varepsilon} (\|f\|_{L^{2}((0,T)\times\Omega)} + \|g\|_{L^{2}((0,T)\times\partial\Omega)} + \|u\|_{L^{2}((0,T),H^{1}(\Omega))} + \|u_{t}\|_{L^{2}((0,T),X)}) \end{split}$$

which gives the last part of the required result (7).

Remark 4 Working on the non-homogeneous adjoint problem and using the backward semigroup $e^{(T-t)\mathcal{A}^*}$, one is able to show that the estimates of Lemma 6 hold for the system

$$(S'_{f,g_0,g_1})^n \begin{cases} \varphi_t + \partial_{x_n} \varphi + \varepsilon \Delta \varphi = f & in \ (0,T) \times \Omega, \\ \varepsilon(\varphi_t - \partial_{\nu} \varphi) - \varphi = g_0 & on \ (0,T) \times \Gamma_0, \\ \varepsilon(\varphi_t - \partial_{\nu} \varphi) = g_1 & on \ (0,T) \times \Gamma_1, \\ \varphi(T,.) = \varphi_T & in \ \Omega. \end{cases}$$

1.3 Backward uniqueness

We will show here the backward uniqueness of systems $(S_0)^n$ and $(S')^n$ for $\varepsilon > 0$, thanks to the well-known result of Lions-Malgrange ([12]) and the regularization effect.

Lemma 7 Let $u_0 \in X$ and $\delta \in (0,T)$. Then, the solution u of $(S_0)^n$ satisfies $u \in \mathcal{C}([\delta,T],\mathcal{D}(\mathcal{A}))$.

Proof Our strategy is to show that $u \in L^2((\delta, T); H^2(\Omega))$ and $u_t \in L^2((\delta, T); H^2(\Omega))$. This easily implies that $u \in \mathcal{C}([\delta, T]; H^2(\Omega))$.

We first select some regular cut-off function θ_1 such that $\theta_1 = 1$ on $(\delta/2, T)$ and $\theta_1 = 0$ on $(0, \delta/4)$. If $u_1 = \theta_1 u$, we have that

$$\begin{array}{ll} \left(\begin{array}{cc} u_{1,t} + \partial_{x_n} u_1 - \varepsilon \Delta u_1 = \theta'_1 u & \text{ in } (0,T) \times \Omega, \\ \varepsilon(u_{1,t} + \partial_{\nu} u_1) = \varepsilon \theta'_1 u & \text{ on } (0,T) \times \Gamma_0, \\ \varepsilon(u_{1,t} + \partial_{\nu} u_1) + u_1 = \varepsilon \theta'_1 u & \text{ on } (0,T) \times \Gamma_1, \\ u_1(0,.) = 0 & \text{ in } \Omega, \end{array} \right)$$

that is u_1 satisfies $(S_{\theta'_1 u, \varepsilon \theta'_1 u, \varepsilon \theta'_1 u})^n$. An application of Lemma 6 gives that $u_1 \in L^2((0,T), \mathcal{D}(\mathcal{A}))$ and $u_{1,t} \in L^2((0,T), X)$. This implies in particular that $u \in L^2((\delta,T), \mathcal{D}(\mathcal{A}))$.

We now focus on u_t and we select another cut-off function θ_2 such that $\theta_2 = 1$ on (δ, T) and $\theta_2 = 0$ on $(0, \delta/2)$. We write the system satisfied by $\theta_2 u$ and we differentiate it with respect to time. One deduces that $u_2 = (\theta_2 u)_t$ is the solution of

$$(S_{\theta_2'u_t+\theta_2''u,\varepsilon(\theta_2'u_t+\theta_2''u),\varepsilon(\theta_2'u_t+\theta_2''u)})^n$$

Using that $|\theta'_2| \leq |\theta_1|$, Lemma 6 gives that $u_2 \in L^2((0,T), \mathcal{D}(\mathcal{A}))$ that is $u_t \in L^2((\delta,T), \mathcal{D}(\mathcal{A}))$.

Proposition 8 Assume that u is a weak solution of $(S_0)^n$ such that u(T) = 0. Then $u_0 \equiv 0$.

Proof If $0 < \delta < T$, then Lemma 7 shows that $u(\delta) \in \mathcal{D}(\mathcal{A})$. We will now apply Théorème 1.1 of [12] to u as a solution of $(S_0)^n$ in the time interval (δ, T) to show that $u(\delta, .) = 0$. The bilinear form a defined in (4) can be split into two bilinear forms on V defined by

$$a(u_1, u_2) = a_0(t, u_1, u_2) + a_1(t, u_1, u_2)$$

with

$$a_0(t,u_1,u_2) = \varepsilon \int_{\Omega} \nabla u_1 \cdot \nabla u_2, \quad a_1(t,u_1,u_2) = \int_{\Omega} \partial_{x_n} u_1 u_2 + \int_{\Gamma_1} u_1 u_2.$$

The first four hypotheses in [12] (see (1.1)-(1.4) in that reference) are satisfied since $u \in L^2(\delta, T; V) \cap H^1(\delta, T; H)$ (from Lemma 6 above), $u(t) \in \mathcal{D}(\mathcal{A})$ for almost every $t \in (\delta, T)$, $u_t = \mathcal{A}u$ and $u(T, .) \equiv 0$. On the other hand, it is clear that a_0 and a_1 are continuous bilinear forms on V and that they do not depend on time t. It is also straightforward to see that

$$\forall u \in V, a_0(t, u, u) + \|u\|_X^2 \ge \min\{1, \varepsilon\} \|u\|^2$$

and, for some constant C > 0,

$$\forall u, v \in V, |a_1(t, u_1, u_2)| \leq C ||u_1|| ||u_2||_X.$$

This means that Hypothesis I in that reference is fulfilled. We have shown that $u(\delta, .) = 0$, which finishes the proof since

$$u_0 = \lim_{\delta \to 0} u(\delta) = 0.$$

Remark 5 The result also holds for the adjoint system: if φ is a weak solution of $(S')^n$ such that $\varphi(0) \equiv 0$ then $\varphi_T = 0$. The proof is very similar to the one above and is left to the reader.

1.4 Proof of Proposition 2

The fact that u = 0 on $(0, T) \times \Gamma_1$ implies $u_0 \equiv 0$ is a straightforward consequence of the backward uniqueness of system $(S_0)^n$. Indeed, similarly as (2), we can prove the observability inequality

$$||u(T,.)||_X \lesssim ||u||_{L^2((0,T) \times \Gamma_1)},$$

(see Remark 6 below), which combined with Proposition 8 gives $u_0 \equiv 0$.

To show that (3) is false, we first need to prove some well-posedness result of $(S_0)^n$ with u_0 in a less regular space. For $s \in (0, 1/2)$, we define X^s as the closure of $\mathcal{D}(\bar{\Omega})$ for the norm

$$||u||_{X^s} := \left(||u||^2_{H^s(\Omega)} + \varepsilon ||u||^2_{H^s(\partial\Omega)} \right)^{\frac{1}{2}},$$

where $H^s(\Omega)$ (resp. $H^s(\partial\Omega)$) stands for the usual Sobolev space on Ω (resp. $\partial\Omega$). One easily shows that X^s is a Hilbert space. We now denote X^{-s} the set of functions u_2 such that the linear form

$$u_1 \in \mathcal{D}(\Omega) \longmapsto \langle u_1, u_2 \rangle_X$$

can be extended in a continuous way on X^s . If $u_1 \in X^{-s}$, $u_2 \in X^s$ we denote this extension by $\langle u_1, u_2 \rangle_{-s,s}$.

If $u_0 \in X^{-s}$, we say that u is a solution by transposition of $(S_0)^n$ if, for every $f \in L^2((0,T) \times \Omega)$

$$\int_{0}^{T} \int_{\Omega} uf + \langle u_{0}, \varphi(0, .) \rangle_{-s,s} = 0,$$

where φ is the weak solution (see Remark 4) of

$$(S'_{f,0,0})^n \begin{cases} \varphi_t + \partial_{x_n} \varphi + \varepsilon \Delta \varphi = f & \text{in } (0,T) \times \Omega, \\ \varepsilon (\varphi_t - \partial_\nu \varphi) - \varphi = 0 & \text{on } (0,T) \times \Gamma_0, \\ \varphi_t - \partial_\nu \varphi = 0 & \text{on } (0,T) \times \Gamma_1, \\ \varphi(T,.) = 0 & \text{in } \Omega. \end{cases}$$

Using the Riesz representation theorem and continuity (see Remark 4) of

$$f \in L^2((0,T) \times \Omega) \longmapsto \varphi(0,.) \in H^1(\Omega),$$

it is obvious that, for any $u_0 \in X^{-s}$, there exists a solution by transposition of $(S_0)^n$, $u \in L^2((0,T) \times \Omega)$, such that

$$||u||_{L^2((0,T)\times\Omega)} \lesssim ||u_0||_{X^{-s}}$$

Moreover, if $u_0 \in X = X^0$, Lemma 6 shows that the solution u satisfies

$$||u||_{L^2((0,T),H^1(\Omega))} \lesssim ||u_0||_{X^0}.$$

If one uses classical interpolation results (see [11]), we can deduce that for every $\theta \in (0, 1/2)$, for any $u_0 \in X^{-s\theta} = [X, X^{-s}]_{1-\theta}$, there exists a solution $u \in L^2((0,T), H^{1-\theta}(\Omega))$ such that

 $||u||_{L^2((0,T),H^{1-\theta}(\Omega))} \lesssim ||u_0||_{X^{-s\theta}}.$

Using classical trace result, we have that

$$\|u\|_{L^2((0,T)\times\partial\Omega)} \lesssim \|u_0\|_{X^{-s\theta}}.$$
(9)

On the other hand, estimate (3) implies in particular that

$$\|u_0\|_X \lesssim \|u\|_{L^2((0,T) \times \Gamma_1)} \quad \forall u_0 \in \mathcal{D}(\overline{\Omega}).$$
(10)

Finally, (9) and (10) yields the contradictory inclusion $X^{-s\theta} \hookrightarrow X$. \Box

2 Carleman inequality in dimension 1

In this paragraph, we will establish a Carleman-type inequality keeping track of the explicit dependence of all the constants with respect to T and ε . As in [5], we introduce the following weight functions:

$$\forall x \in [-L, 0], \quad \eta(x) := 2L + x, \quad \alpha(t, x) := \frac{\lambda - e^{\eta(x)}}{\varepsilon^2 t(T - t)}, \quad \phi(t, x) := \frac{e^{\eta(x)}}{\varepsilon^2 t(T - t)}$$

where $\lambda > e^{2L}$.

The rest of this paragraph will be dedicated to the proof of the following inequality:

Theorem 9 There exists C > 0 and $s_0 > 0$ such that for every $\varepsilon \in (0,1)$ and every $s \ge s_0(\varepsilon T + \varepsilon^2 T^2)$ the following inequality is satisfied for every $\varphi_T \in X$:

$$s^{3} \int_{(0,T)\times(-L,0)} \phi^{3} e^{-2s\alpha} |\varphi|^{2} + s \int_{(0,T)\times(-L,0)} \phi e^{-2s\alpha} |\varphi_{x}|^{2} + s^{3} \int_{(0,T)\times\{0,-L\}} \phi^{3} e^{-2s\alpha} |\varphi|^{2}$$

$$\leq Cs^{7} \int_{(0,T)\times\{0\}} e^{-4s\alpha+2s\alpha(\cdot,-L)} \phi^{7} |\varphi|^{2}.$$

$$(11)$$

Here, φ stands for the solution of (S') associated to φ_T .

Remark 6 One can in fact obtain the following Carleman estimate with control term in Γ_1

$$s^{3} \int_{(0,T)\times(-L,0)} \phi^{3} e^{-2s\alpha} |\varphi|^{2} + s \int_{(0,T)\times(-L,0)} \phi e^{-2s\alpha} |\varphi_{x}|^{2} + s^{3} \int_{(0,T)\times\{0,-L\}} \phi^{3} e^{-2s\alpha} |\varphi|^{2} \\ \leqslant Cs^{7} \int_{(0,T)\times\{-L\}} e^{-4s\alpha+2s\alpha(.,0)} \phi^{7} |\varphi|^{2}.$$

simply by choosing the weight function $\eta(x)$ equal to $x \mapsto -x + L$. The proof is very similar and the sequel will show how to deal with this case too.

In order to perform a Carleman inequality for system (S'), we first do a scaling in time. Introducing $\tilde{T} := \varepsilon T$, $\tilde{\varphi}(t,x) := \varphi(t/\varepsilon,x)$ and the weights $\tilde{\alpha}(t,x) := \alpha(t/\varepsilon,x)$, $\tilde{\phi}(t,x) := \phi(t/\varepsilon,x)$; we have the following system

$$\begin{cases}
\tilde{\varphi}_t + \varepsilon^{-1}\tilde{\varphi}_x + \tilde{\varphi}_{xx} = 0 & \text{in } q, \\
\varepsilon^2(\tilde{\varphi}_t - \varepsilon^{-1}\partial_\nu\tilde{\varphi}) - \tilde{\varphi} = 0 & \text{on } \sigma_0, \\
\tilde{\varphi}_t - \varepsilon^{-1}\partial_\nu\tilde{\varphi} = 0 & \text{on } \sigma_1, \\
\tilde{\varphi}_{|t=\tilde{T}} = \varphi_T & \text{in } (-L, 0).
\end{cases}$$
(12)

if $q := (0, \tilde{T}) \times (-L, 0)$, $\sigma := (0, \tilde{T}) \times \{-L, 0\}$, $\sigma_0 := (0, \tilde{T}) \times \{0\}$ and $\sigma_1 := (0, \tilde{T}) \times \{-L\}$. We will now explain how to get the following result.

Proposition 10 There exists C > 0, such that for every $\varepsilon \in (0,1)$ and every $s \ge C(\tilde{T} + \varepsilon^{-1}\tilde{T}^2 + \varepsilon^{1/3}\tilde{T}^{2/3})$, the following inequality is satisfied for every solution of (12) associated to $\tilde{\varphi}_T \in X$:

$$s^{3} \int_{q} \tilde{\phi}^{3} e^{-2s\tilde{\alpha}} |\tilde{\varphi}|^{2} + s \int_{q} \tilde{\phi} e^{-2s\tilde{\alpha}} |\tilde{\varphi}_{x}|^{2} + s^{3} \int_{\sigma} \tilde{\phi}^{3} e^{-2s\tilde{\alpha}} |\tilde{\varphi}|^{2} \leqslant Cs^{7} \int_{\sigma_{0}} e^{-4s\tilde{\alpha} + 2s\tilde{\alpha}(.,-L)} \tilde{\phi}^{7} |\tilde{\varphi}|^{2}.$$
(13)

Observe that Proposition 10 directly implies Theorem 9.

All along the proof we will need several properties of the weight functions:

Lemma 11

• $|\tilde{\alpha}_t| \lesssim \tilde{T}\tilde{\phi}^2$, $|\tilde{\alpha}_{xt}| \lesssim \tilde{T}\tilde{\phi}^2$, $|\tilde{\alpha}_{tt}| \lesssim \tilde{T}^2\tilde{\phi}^3$,

•
$$\tilde{\alpha}_x = -\tilde{\phi}, \ \tilde{\alpha}_{xx} = -\tilde{\phi}.$$

There are essentially two steps in this proof:

- The first one consists in doing a Carleman estimate similar to that of [5]. The idea is to compute the L^2 product of some operators applied to $e^{-s\tilde{\alpha}}\tilde{\varphi}$, do integrations by parts and make a convenient choice of the parameter s. This will give an estimate where we still have a boundary term on σ_0 in the right-hand side.
- In the second one, we will study the boundary terms appearing in the right hand side due to the boundary conditions. We will perform the proof of this theorem for smooth solutions, so that the general proof follows from a density argument.

The big difference with general parabolic equations is the special boundary condition. Up to our knowledge, no Carleman inequalities have been performed when the partial time derivative appears on the boundary condition. The fact that we still have a Carleman inequality for our system comes from the fact that our system remains somehow parabolic.

We will first estimate the left hand side terms of (13) like in the classical Carleman estimate (see [5]). We obtain:

Proposition 12 There exists C > 0, such that for every $\varepsilon \in (0,1)$ and every $s \ge C(\tilde{T} + \varepsilon^{-1}\tilde{T}^2 + \varepsilon^{1/3}\tilde{T}^{2/3})$, the following inequality is satisfied for every solution of (12) associated to $\tilde{\varphi}_T \in X$:

$$s^{3} \int_{q} \tilde{\phi}^{3} e^{-2s\tilde{\alpha}} |\tilde{\varphi}|^{2} + s^{2} \int_{q} \tilde{\phi} e^{-2s\tilde{\alpha}} |\tilde{\varphi}_{x}|^{2} + s^{-1} \int_{q} \tilde{\phi}^{-1} e^{-2s\tilde{\alpha}} (|\tilde{\varphi}_{xx}|^{2} + |\tilde{\varphi}_{t}|^{2})$$

$$+ s^{3} \int_{\sigma_{1}} \tilde{\phi}^{3} e^{-2s\tilde{\alpha}} |\tilde{\varphi}|^{2} \lesssim s^{5} \int_{\sigma_{0}} \tilde{\phi}^{5} e^{-2s\tilde{\alpha}} |\tilde{\varphi}|^{2} + \varepsilon^{2} s \int_{\sigma_{0}} \tilde{\phi} e^{-2s\tilde{\alpha}} |\tilde{\varphi}_{t}|^{2}.$$

$$(14)$$

Proof

Let us introduce $\psi := \tilde{\varphi} e^{-s\tilde{\alpha}}$. We state the equations satisfied by ψ . In $(0, \tilde{T}) \times (-L, 0)$, we have the identity

$$P_1\psi + P_2\psi = P_3\psi$$

where

$$P_1\psi = \psi_t + 2s\tilde{\alpha}_x\psi_x + \varepsilon^{-1}\psi_x,\tag{15}$$

$$P_2\psi = \psi_{xx} + s^2 \tilde{\alpha}_x^2 \psi + s \tilde{\alpha}_t \psi + \varepsilon^{-1} s \tilde{\alpha}_x \psi, \qquad (16)$$

and

$$P_3\psi = -s\tilde{\alpha}_{xx}\psi.$$

On the other hand, the boundary conditions are:

$$\varepsilon^2(\psi_t + s\tilde{\alpha}_t\psi - \varepsilon^{-1}(\psi_x + s\tilde{\alpha}_x\psi)) - \psi = 0 \quad \text{on } x = 0,$$
(17)

$$\varepsilon(\psi_t + s\tilde{\alpha}_t\psi) + \psi_x + s\tilde{\alpha}_x\psi = 0 \quad \text{on } x = -L.$$
(18)

We take the L^2 norm in both sides of the identity in q:

$$\|P_1\psi\|_{L^2(q)}^2 + \|P_2\psi\|_{L^2(q)}^2 + 2(P_1\psi, P_2\psi)_{L^2(q)} = \|P_3\psi\|_{L^2(q)}^2.$$
(19)

Using Lemma 11, we directly obtain

$$||P_3\psi||^2_{L^2(q)} \lesssim s^2 \int_q \tilde{\phi}^2 |\psi|^2.$$

We focus on the expression of the double product $(P_1\psi, P_2\psi)_{L^2(q)}$. This product contains 12 terms which will be denoted by $T_{ij}(\psi)$ for $1 \leq i \leq 3, 1 \leq j \leq 4$. We study them successively.

• An integration by parts in space and then in time shows that, since $\psi(0,.) = \psi(\tilde{T},.) = 0$, we have

$$T_{11}(\psi) = \int_{q} \psi_t \psi_{xx} = -\frac{1}{2} \int_{q} (|\psi_x|^2)_t + \int_{\sigma_0} \psi_x \psi_t - \int_{\sigma_1} \psi_x \psi_t$$
$$= \int_{\sigma_0} \psi_x \psi_t - \int_{\sigma_1} \psi_x \psi_t.$$

Now, we use the boundary condition (17). We have

$$\int_{\sigma_0} \psi_x \psi_t \gtrsim \varepsilon \int_{\sigma_0} |\psi_t|^2 - \varepsilon \tilde{T}^2 s \int_{\sigma_0} \tilde{\phi}^3 |\psi|^2 - \tilde{T} s \int_{\sigma_0} \tilde{\phi}^2 |\psi|^2.$$

Thanks to the fact that $\psi(0, .) = \psi(\tilde{T}, .) = 0$, the same computations can be done on x = -L using (18) so we obtain the following for this term:

$$T_{11}(\psi) \gtrsim -\varepsilon \tilde{T}^2 s \int_{\sigma} \tilde{\phi}^3 |\psi|^2 - \tilde{T} s \int_{\sigma} \tilde{\phi}^2 |\psi|^2.$$

• Integrating by parts in time, we find

$$T_{12}(\psi) = \frac{s^2}{2} \int_q \tilde{\alpha}_x^2 (|\psi|^2)_t = -s^2 \int_q \tilde{\alpha}_x \tilde{\alpha}_{xt} |\psi|^2,$$

$$T_{13}(\psi) = \frac{s}{2} \int_q \tilde{\alpha}_t (|\psi|^2)_t = -\frac{s}{2} \int_q \tilde{\alpha}_{tt} |\psi|^2,$$

$$T_{14}(\psi) = \varepsilon^{-1} \frac{s}{2} \int_q \tilde{\alpha}_x (|\psi|^2)_t = -\varepsilon^{-1} \frac{s}{2} \int_q \tilde{\alpha}_{tx} |\psi|^2,$$

and using Lemma 11, we get

$$\begin{split} T_{12}(\psi) \gtrsim &-\tilde{T}s^2 \int_q \tilde{\phi}^3 |\psi|^2, \\ T_{13}(\psi) \gtrsim &-\tilde{T}^2s \int_q \tilde{\phi}^3 |\psi|^2, \\ T_{14}(\psi) \gtrsim &-\varepsilon^{-1}\tilde{T}s \int_q \tilde{\phi}^2 |\psi|^2. \end{split}$$

• Now, integrating by parts in space, we have

$$T_{21}(\psi) = s \int_{q} \tilde{\alpha}_{x} (|\psi_{x}|^{2})_{x} = -s \int_{q} \tilde{\alpha}_{xx} |\psi_{x}|^{2} + s \int_{\sigma_{0}} \tilde{\alpha}_{x} |\psi_{x}|^{2} - s \int_{\sigma_{1}} \tilde{\alpha}_{x} |\psi_{x}|^{2}.$$

Thanks to Lemma 11 (the choice of η is important here), the last term is positive. Using the boundary condition (17) and Lemma 11, we finally get

$$T_{21}(\psi) \gtrsim s \int_{q} \tilde{\phi} |\psi_{x}|^{2} - s^{3} \int_{\sigma_{0}} \tilde{\phi}^{3} |\psi|^{2} - \varepsilon^{2} s^{3} \tilde{T}^{2} \int_{\sigma_{0}} \tilde{\phi}^{5} |\psi|^{2} + s \int_{\sigma_{1}} \tilde{\phi} |\psi_{x}|^{2} - \varepsilon^{2} s \int_{\sigma_{0}} \tilde{\phi} |\psi_{t}|^{2}.$$

• An integration by parts in space provides

$$T_{22}(\psi) = s^3 \int_q \tilde{\alpha}_x^3 (|\psi|^2)_x = -3s^3 \int_q \tilde{\alpha}_x^2 \tilde{\alpha}_{xx} |\psi|^2 + s^3 \int_{\sigma_0} \tilde{\alpha}_x^3 |\psi|^2 - s^3 \int_{\sigma_1} \tilde{\alpha}_x^3 |\psi|^2.$$

This readily yields

$$T_{22}(\psi) \gtrsim s^3 \int_q \tilde{\phi}^3 |\psi|^2 + s^3 \int_{\sigma_1} \tilde{\phi}^3 |\psi|^2 - s^3 \int_{\sigma_0} \tilde{\phi}^3 |\psi|^2.$$

• Again an integration by parts in space gives

$$T_{23}(\psi) = s^2 \int_q \tilde{\alpha}_t \tilde{\alpha}_x (|\psi|^2)_x = -s^2 \int_q (\tilde{\alpha}_t \tilde{\alpha}_x)_x |\psi|^2 + s^2 \int_{\sigma_0} \tilde{\alpha}_t \tilde{\alpha}_x |\psi|^2 - s^2 \int_{\sigma_1} \tilde{\alpha}_t \tilde{\alpha}_x |\psi|^2.$$

Using Lemma 11, we obtain

$$T_{23}(\psi) \gtrsim -s^2 \tilde{T} \int_q \tilde{\phi}^3 |\psi|^2 - s^2 \tilde{T} \int_{\sigma_0} \tilde{\phi}^3 |\psi|^2 - s^2 \tilde{T} \int_{\sigma_1} \tilde{\phi}^3 |\psi|^2.$$

• The last integral concerning the second term in the expression of $P_1\psi$ is

$$T_{24}(\psi) = \varepsilon^{-1} s^2 \int_q \tilde{\alpha}_x^2 (|\psi|^2)_x.$$

After an integration by parts in space, we get

$$T_{24}(\psi) \ge -2\varepsilon^{-1}s^2 \int_q \tilde{\phi}^2 |\psi|^2 - \varepsilon^{-1}s^2 \int_{\sigma_1} \tilde{\phi}^2 |\psi|^2.$$

• We consider now the third term in the expression of $P_1\psi$. We have

$$T_{31}(\psi) = \frac{\varepsilon^{-1}}{2} \int_{q} (|\psi_x|^2)_x \ge -\frac{\varepsilon^{-1}}{2} \int_{\sigma_1} |\psi_x|^2 dx$$

• Now, we integrate by parts with respect to x and we have

$$\begin{split} T_{32}(\psi) &= \frac{\varepsilon^{-1}}{2} s^2 \int_q \tilde{\alpha}_x^2 (|\psi|^2)_x \ge -\varepsilon^{-1} s^2 \int_q \tilde{\alpha}_x \tilde{\alpha}_{xx} |\psi|^2 - \frac{\varepsilon^{-1}}{2} s^2 \int_{\sigma_1} \tilde{\alpha}_x^2 |\psi|^2 \\ &\ge -\varepsilon^{-1} s^2 \int_q \tilde{\phi}^2 |\psi|^2 - \frac{\varepsilon^{-1}}{2} s^2 \int_{\sigma_1} \tilde{\phi}^2 |\psi|^2. \end{split}$$

• Then, using another integration by parts in x we obtain

$$T_{33}(\psi) = \frac{\varepsilon^{-1}}{2} s \int_{q} \tilde{\alpha}_{t}(|\psi|^{2})_{x} \gtrsim -\varepsilon^{-1} s \tilde{T} \int_{q} \tilde{\phi}^{2} |\psi|^{2} - \varepsilon^{-1} \tilde{T} s \int_{\sigma_{1}} \tilde{\phi}^{2} |\psi|^{2} - \varepsilon^{-1} \tilde{T} s \int_{\sigma_{0}} \tilde{\phi}^{2} |\psi|^{2}.$$

• Finally, arguing as before, we find

$$T_{34}(\psi) = \frac{\varepsilon^{-2}}{2} s \int_{q} \tilde{\alpha}_{x}(|\psi|^{2})_{x} \ge \frac{\varepsilon^{-2}}{2} s^{2} \int_{q} \tilde{\phi}|\psi|^{2} - C\varepsilon^{-2} s \int_{\sigma_{0}} \tilde{\phi}|\psi|^{2} \gtrsim -\varepsilon^{-2} s \int_{\sigma_{0}} \tilde{\phi}|\psi|^{2}.$$

Putting together all the terms and combining the resulting inequality with (19), we obtain

$$\|P_1\psi\|_{L^2(q)}^2 + \|P_2\psi\|_{L^2(q)}^2 + I_1(\psi) + I_2(\psi_x) \lesssim J_1(\psi) + J_2(\psi_t) + J_3(\psi_x) + L(\psi)$$
(20)

where the main terms are

$$I_1(\psi) = s^3 \int_q \tilde{\phi}^3 |\psi|^2 + s^3 \int_{\sigma_1} \tilde{\phi}^3 |\psi|^2, \quad I_2(\psi_x) = s \int_q \tilde{\phi} |\psi_x|^2 + s \int_{\sigma_1} \tilde{\phi} |\psi_x|^2$$

the right hand side terms are

$$J_1(\psi) = (\varepsilon^{-1}\tilde{T}s + \varepsilon^{-1}s^2) \int_q \tilde{\phi}^2 |\psi|^2 + (s^2\tilde{T} + s\tilde{T}^2) \int_q \tilde{\phi}^3 |\psi|^2 + (s\tilde{T} + s^2\varepsilon^{-1} + \varepsilon^{-1}\tilde{T}s) \int_{\sigma_1} \tilde{\phi}^2 |\psi|^2 + (s^2\tilde{T} + \varepsilon s\tilde{T}^2) \int_{\sigma_1} \tilde{\phi}^3 |\psi|^2,$$

and

$$J_2(\psi_t) = \varepsilon^2 s \int_{\sigma_0} \tilde{\phi} |\psi_t|^2, \quad J_3(\psi_x) = \varepsilon^{-1} \int_{\sigma_1} |\psi_x|^2$$

and the control terms are

$$L(\psi) = \varepsilon^{-2}s \int_{\sigma_0} \tilde{\phi} |\psi|^2 + (s\tilde{T} + \varepsilon^{-1}\tilde{T}s) \int_{\sigma_0} \tilde{\phi}^2 |\psi|^2 + (s^3 + s^2\tilde{T} + \varepsilon\tilde{T}^2s) \int_{\sigma_0} \tilde{\phi}^3 |\psi|^2 + \varepsilon^2 s^3\tilde{T}^2 \int_{\sigma_0} \tilde{\phi}^5 |\psi|^2.$$

Let us now see that we can absorb some right hand side terms with the help of the parameter s.

• First, we see that the distributed terms in $J_1(\psi)$ can be absorbed by the first term in the definition of $I_1(\psi)$ for a choice of $s \gtrsim \varepsilon^{-1/2} \tilde{T}^{3/2} + \varepsilon^{-1} \tilde{T}^2$.

• Second, we use the second term of $I_1(\psi)$ in order to absorb the integrals in the second line of the definition of $J_1(\psi)$. We find that this can be done as long as $s \gtrsim \varepsilon^{-1}\tilde{T}^2 + \varepsilon^{-1/2}\tilde{T}^{3/2}$ since $\varepsilon < 1$.

• Next, we observe that, provided $s \gtrsim \varepsilon^{-1} \tilde{T}^2$, the second term in $J_2(\psi_x)$ absorbs $J_3(\psi_x)$.

$$|L(\psi)| \lesssim s^5 \int_{\sigma_0} \tilde{\phi}^5 |\psi|^2$$

as long as $s \gtrsim \tilde{T}(1 + \varepsilon^{-1}\tilde{T})$, using that $\varepsilon < 1$.

Next, we use the expression of $P_1\psi$ and $P_2\psi$ (see (15)-(16)) in order to obtain some estimates for the terms ψ_t and ψ_{xx} respectively:

$$s^{-1} \int_{q} \tilde{\phi}^{-1} |\psi_{t}|^{2} \lesssim s^{2} \int_{q} \tilde{\phi} |\psi_{x}|^{2} + s^{-1} \varepsilon^{-2} \int_{q} \tilde{\phi}^{-1} |\psi_{x}|^{2} + s^{-1} \int_{q} \tilde{\phi}^{-1} |P_{1}\psi|^{2} \lesssim s^{2} \int_{q} \tilde{\phi} |\psi_{x}|^{2} + \int_{q} |P_{1}\psi|^{2} = s^{-1} \int_{q} \tilde{\phi}^{-1} |\psi_{x}|^{2} + s^{-1} |\psi_{x}|^{2} + s^{-1} |$$

for any $s \gtrsim \varepsilon^{-1} \tilde{T}^2$ and

$$\begin{split} s^{-1} \int_{q} \tilde{\phi}^{-1} |\psi_{xx}|^{2} &\lesssim s^{3} \int_{0}^{\tilde{T}} \int_{-L}^{0} \tilde{\phi}^{3} |\psi|^{2} + s \tilde{T}^{2} \int_{q} \tilde{\phi}^{3} |\psi_{x}|^{2} + \varepsilon^{-2} s^{2} \int_{q} \tilde{\phi} |\psi|^{2} \\ + s^{-1} \int_{q} \tilde{\phi}^{-1} |P_{2}\psi|^{2} &\lesssim s^{3} \int_{q} \tilde{\phi}^{3} |\psi|^{2} + s^{2} \int_{q} \tilde{\phi} |\psi_{x}|^{2} + \int_{q} |P_{2}\psi|^{2}, \end{split}$$

for $s\gtrsim \tilde{T}+\varepsilon^{-1}\tilde{T}^2$. Combining all this with (20), we obtain

$$s \int_{q} \tilde{\phi}(s^{2} \tilde{\phi}^{2} |\psi|^{2} + |\psi_{x}|^{2}) + s^{-1} \int_{q} \tilde{\phi}^{-1}(|\psi_{xx}|^{2} + |\psi_{t}|^{2}) + s^{3} \int_{\sigma_{1}} \tilde{\phi}^{3} |\psi|^{2}$$

$$\lesssim s^{5} \int_{\sigma_{0}} \tilde{\phi}^{5} |\psi|^{2} + \varepsilon^{2} s \int_{\sigma_{0}} \tilde{\phi} |\psi_{t}|^{2},$$

$$(21)$$

for $s \gtrsim \tilde{T}(1 + \varepsilon^{-1}\tilde{T})$.

Finally, we come back to our variable $\tilde{\varphi}$. We first remark that $\psi_x = e^{-s\tilde{\alpha}}(\tilde{\varphi}_x + s\tilde{\phi}\tilde{\varphi})$ and so

$$s \int_{q} \tilde{\phi} e^{-2s\tilde{\alpha}} |\tilde{\varphi}_{x}|^{2} \lesssim s \int_{q} \tilde{\phi} |\psi_{x}|^{2} + s^{3} \int_{q} \tilde{\phi}^{3} |\psi|^{2}.$$

Then, we have that $\psi_t = e^{-s\tilde{\alpha}} (\tilde{\varphi}_t - s\tilde{\alpha}_t \tilde{\varphi})$, hence

$$s^{-1} \int_{q} \tilde{\phi}^{-1} e^{-2s\tilde{\alpha}} |\tilde{\varphi}_{t}|^{2} \lesssim s^{-1} \int_{q} \tilde{\phi}^{-1} |\psi_{t}|^{2} + s\tilde{T}^{2} \int_{q} \tilde{\phi}^{3} |\psi|^{2} \lesssim s^{-1} \int_{q} \tilde{\phi}^{-1} |\psi_{t}|^{2} + s^{3} \int_{q} \tilde{\phi}^{3} |\psi|^{2} \lesssim s^{-1} \int_{q} \tilde{\phi}^{-1} |\psi_{t}|^{2} + s^{3} \int_{q} \tilde{\phi}^{3} |\psi|^{2} \lesssim s^{-1} \int_{q} \tilde{\phi}^{-1} |\psi_{t}|^{2} + s^{3} \int_{q} \tilde{\phi}^{3} |\psi|^{2} \lesssim s^{-1} \int_{q} \tilde{\phi}^{-1} |\psi_{t}|^{2} + s^{3} \int_{q} \tilde{\phi}^{-1} |$$

for $s \gtrsim \tilde{T}$. Analogously, we can prove that

$$s^{-1} \int_0^T \int_q \tilde{\phi}^{-1} e^{-2s\tilde{\alpha}} |\tilde{\varphi}_{xx}|^2 \lesssim s^{-1} \int_q \tilde{\phi}^{-1} |\psi_{xx}|^2 + s^3 \int_q \tilde{\phi}^3 |\psi|^2 + s^2 \int_q \tilde{\phi} |\psi_x|^2,$$

for $s \gtrsim \tilde{T}^2$.

We combine this with (21) and we obtain the required result

$$s^{3} \int_{q} \tilde{\phi}^{3} e^{-2s\tilde{\alpha}} |\tilde{\varphi}|^{2} + s \int_{q} \tilde{\phi} e^{-2s\tilde{\alpha}} |\tilde{\varphi}_{x}|^{2} + s^{-1} \int_{q} \tilde{\phi}^{-1} e^{-2s\tilde{\alpha}} (|\tilde{\varphi}_{xx}|^{2} + |\tilde{\varphi}_{t}|^{2}) + s^{3} \int_{\sigma_{1}} \tilde{\phi}^{3} e^{-2s\tilde{\alpha}} |\tilde{\varphi}|^{2}$$

$$\lesssim s^{5} \int_{\sigma_{0}} \tilde{\phi}^{5} e^{-2s\tilde{\alpha}} |\tilde{\varphi}|^{2} + \varepsilon^{2} s \int_{\sigma_{0}} \tilde{\phi} e^{-2s\tilde{\alpha}} |\tilde{\varphi}_{t}|^{2}, \qquad (22)$$

for $s\gtrsim \tilde{T}(1+\varepsilon^{-1}\tilde{T})+\varepsilon^{1/3}\tilde{T}^{2/3}$ and using that

$$\varepsilon^2 s^3 \tilde{T}^2 \int_{\sigma_0} \tilde{\phi}^5 e^{-2s\tilde{\alpha}} |\tilde{\varphi}|^2 \lesssim s^5 \int_{\sigma_0} \tilde{\phi}^5 e^{-2s\tilde{\alpha}} |\tilde{\varphi}|^2$$

which is true for $s \gtrsim \tilde{T}$ (recall that $\varepsilon < 1$).

With this result, we will now finish the proof of Proposition 10.

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Estimate of the boundary term In this paragraph we will estimate the boundary term

$$\varepsilon^2 s \int_{\sigma_0} \tilde{\phi} e^{-2s\tilde{\alpha}} |\tilde{\varphi}_t|^2.$$

After an integration by parts in time, we get

$$\varepsilon^{2}s \int_{\sigma_{0}} \tilde{\phi}e^{-2s\tilde{\alpha}} |\tilde{\varphi}_{t}|^{2} = \frac{\varepsilon^{2}}{2}s \int_{\sigma_{0}} (\tilde{\phi}e^{-2s\tilde{\alpha}})_{tt} |\tilde{\varphi}|^{2} - \varepsilon^{2}s \int_{\sigma_{0}} \tilde{\phi}e^{-2s\tilde{\alpha}} \tilde{\varphi} \tilde{\varphi}_{tt}$$

$$\lesssim \tilde{T}^{2}\varepsilon^{2}s^{3} \int_{\sigma_{0}} \tilde{\phi}^{5}e^{-2s\tilde{\alpha}} |\tilde{\varphi}|^{2} + \varepsilon^{2}s \int_{\sigma_{0}} \tilde{\phi}e^{-2s\tilde{\alpha}} |\tilde{\varphi}| |\tilde{\varphi}_{tt}|,$$
(23)

for $s \gtrsim \tilde{T} + \tilde{T}^2$. In order to estimate the second time derivative at x = 0, we will apply some a priori estimates for the adjoint system. Indeed, let us consider the following function

$$\zeta(t,x) := \theta(t)\tilde{\varphi}_t := e^{-s\tilde{\alpha}(t,-L)}\tilde{\phi}^{-5/2}(t,-L)\tilde{\varphi}_t(t,x).$$

Then, this function fulfills the following system

$$\begin{cases} \zeta_t + \varepsilon^{-1}\zeta_x + \zeta_{xx} = \theta_t \tilde{\varphi}_t & \text{in } (0, \tilde{T}) \times (-L, 0), \\ \varepsilon^2 (\zeta_t - \varepsilon^{-1} \partial_\nu \zeta) - \zeta = \varepsilon^2 \theta_t \tilde{\varphi}_t & \text{on } (0, \tilde{T}) \times \{0\}, \\ \zeta_t - \varepsilon^{-1} \partial_\nu \zeta = \theta_t \tilde{\varphi}_t & \text{on } (0, \tilde{T}) \times \{-L\}, \\ \zeta_{|t=T} = 0 & \text{in } (-L, 0). \end{cases}$$

Using Remark 4 for $(t, x) \mapsto \zeta(\varepsilon t, x)$ we find in particular

$$\varepsilon \int_{\sigma} |\zeta_t|^2 \lesssim \int_{q} |\theta_t|^2 |\tilde{\varphi}_t|^2 + \varepsilon^2 \int_{\sigma} (\theta_t)^2 |\tilde{\varphi}_t|^2.$$

This directly implies that

$$\varepsilon \int_{\sigma} \theta^2 |\tilde{\varphi}_{tt}|^2 \lesssim \int_{q} |\theta_t|^2 |\tilde{\varphi}_t|^2 + \varepsilon^2 \int_{\sigma} (\theta_t)^2 |\tilde{\varphi}_t|^2.$$

Integrating by parts in time in the last integral, we have

$$\varepsilon \int_{\sigma} \theta^2 |\tilde{\varphi}_{tt}|^2 \lesssim \left(\int_{q} |\theta_t|^2 |\tilde{\varphi}_t|^2 + \varepsilon^2 \int_{\sigma} ((\theta_t)^2)_{tt} |\tilde{\varphi}|^2 + \varepsilon^3 \int_{\sigma} (\theta_t)^4 \theta^{-2} |\tilde{\varphi}|^2 \right) + \frac{\varepsilon}{2} \int_{\sigma} \theta^2 |\tilde{\varphi}_{tt}|^2.$$

From the definition of $\theta(t)$ and multiplying the previous inequality by s^{-3} , we find that (since $\varepsilon < 1$)

$$\varepsilon s^{-3} \int_{\sigma} (e^{-2s\tilde{\alpha}} \tilde{\phi}^{-5})(.,-L) |\tilde{\varphi}_{tt}|^2 \lesssim \tilde{T}^2 \left(s^{-1} \int_{q} \tilde{\phi}^{-1} e^{-2s\tilde{\alpha}} |\tilde{\varphi}_t|^2 + \varepsilon^2 \tilde{T}^2 s \int_{\sigma} (\tilde{\phi}^3 e^{-2s\tilde{\alpha}})(.,-L) |\tilde{\varphi}|^2 \right)$$

$$\lesssim \tilde{T}^2 \left(s^{-1} \int_{q} \tilde{\phi}^{-1} e^{-2s\tilde{\alpha}} |\tilde{\varphi}_t|^2 + s^3 \int_{\sigma} (\tilde{\phi}^3 e^{-2s\tilde{\alpha}})(.,-L) |\tilde{\varphi}|^2 \right)$$

$$(24)$$

for $s \gtrsim \varepsilon \tilde{T}$. Using Cauchy-Schwarz inequality in the last term of the right hand side of (23), we obtain

Combining this with (24) and (22) yields the desired inequality (13). \Box

3 Observability and control

In this section, we prove Theorem 1.

3.1 Dissipation and observability result

Our first goal, as in [2], will be to get some dissipation result. Even if we will only use this result in dimension one, we present it in dimension n for the sake of completeness (see also [3]).

Proposition 13 For every $\varepsilon \in (0,1)$, for every time $t_1, t_2 > 0$ such that $t_2 - t_1 > L$ and for every weak solution φ of $(S')^n$, the following estimate holds

$$\|\varphi(t_1)\|_X \leq \exp\left\{-\frac{(t_2-t_1-L)^2}{4\varepsilon(t_2-t_1)}\right\} \|\varphi(t_2)\|_X.$$

Proof We first consider a weight function $\rho(t, x) = \exp(\frac{r}{\varepsilon}x_n)$ for some constant $r \in (0, 1)$ which will be fixed later. We will first treat the strong solutions case and, using a density argument, we will get the weak solutions case.

We multiply the equation satisfied by φ by $\rho\varphi$ and we integrate on Ω . We get the following identity:

$$\frac{1}{2}\frac{d}{dt}\left(\int_{\Omega}\rho|\varphi|^{2}\right) = -\frac{1}{2}\int_{\Omega}\rho\partial_{x_{n}}(|\varphi|^{2}) - \varepsilon\int_{\Omega}\rho\varphi\Delta\varphi.$$

We then integrate by parts in space, which due to $\nabla \rho = \frac{r}{\varepsilon} \rho e_n$, provides

$$-\frac{1}{2}\int_{\Omega}\rho\partial_{x_n}(|\varphi|^2) = \frac{r}{2\varepsilon}\int_{\Omega}\rho|\varphi|^2 - \frac{1}{2}\left(\int_{\Gamma_0}\rho|\varphi|^2 - \int_{\Gamma_1}\rho|\varphi|^2\right),$$

and

$$-\varepsilon \int_{\Omega} \rho \varphi \Delta \varphi = \varepsilon \int_{\Omega} \rho |\nabla \varphi|^{2} + \frac{r}{2} \int_{\Omega} \rho \partial_{x_{n}} (|\varphi|^{2}) - \varepsilon \left(\int_{\Gamma_{0}} \rho \varphi \partial_{x_{n}} \varphi - \int_{\Gamma_{1}} \rho \varphi \partial_{x_{n}} \varphi \right)$$
$$= \varepsilon \int_{\Omega} \rho |\nabla \varphi|^{2} - \frac{r^{2}}{2\varepsilon} \int_{\Omega} \rho |\varphi|^{2} + \frac{r}{2} \left(\int_{\Gamma_{0}} \rho |\varphi|^{2} - \int_{\Gamma_{1}} \rho |\varphi|^{2} \right) - \varepsilon \left(\int_{\Gamma_{0}} \rho \varphi \partial_{x_{n}} \varphi - \int_{\Gamma_{1}} \rho \varphi \partial_{x_{n}} \varphi \right).$$

Using now the boundary conditions for φ and summing up these identities, we finally get

$$\frac{d}{dt}\left(\int_{\Omega}\rho|\varphi|^{2}\right) \geq \frac{r(1-r)}{\varepsilon}\int_{\Omega}\rho|\varphi|^{2} + (1-r)\int_{\Gamma}\rho|\varphi|^{2} - 2\varepsilon\int_{\Gamma}\rho\varphi_{t}\varphi.$$

On the other hand, it is straightforward that

$$\frac{d}{dt}\left(\varepsilon\int_{\Gamma}\rho|\varphi|^{2}\right)=2\varepsilon\int_{\Gamma}\rho\varphi_{t}\varphi,$$

and, consequently, using that $r \in (0, 1)$, we have obtained

$$\frac{d}{dt}\left(\|\sqrt{\rho(.)}\varphi(t)\|_X^2\right) \ge \frac{r(1-r)}{\varepsilon}\|\sqrt{\rho(.)}\varphi(t)\|_X^2$$

Gronwall's lemma combined with $\exp(-\frac{r}{\varepsilon}L) \leq \rho(\cdot) \leq 1$ successively gives

$$\|\sqrt{\rho(\cdot)}\varphi(\cdot)\|_X^2 \leqslant \exp\left(-\frac{r(1-r)}{\varepsilon}(t_2-t_1)\right)\|\sqrt{\rho(\cdot)}\varphi(t_2)\|_X^2$$

and

$$\|\varphi(t_1)\|_X^2 \leq \exp\left(-\frac{1}{\varepsilon}\left(r(1-r)(t_2-t_1)-rL\right)\right)\|\varphi(t_2)\|_X^2$$

We finally choose

$$r := \frac{t_2 - t_1 - L}{2(t_2 - t_1)} \in (0, 1),$$

which gives the result.

We will now use this dissipation estimate with our Carleman inequality to get the desired result.

Proposition 14 If n = 1, $\frac{T}{L}$ is sufficiently large and ε sufficiently small, then the observability constant $C_{obs}(\varepsilon)$ is bounded by

$$C\exp\left(-\frac{k}{\varepsilon}\right)$$

where C, k are some positive constants.

Proof

• We begin with estimating both sides of the Carleman inequality obtained above. We use the same notations as above and we define $m = \lambda - e^{2L}$ and $M = \lambda - e^{L}$. We first get

On the other hand, using that $\phi \gtrsim \frac{1}{(\varepsilon T)^2}$ on $[\frac{T}{4}, \frac{3T}{4}]$, we have the following estimate from below for the left hand-side of the Carleman inequality (13)

$$\frac{s^3}{(\varepsilon T)^6} \exp\left(-\frac{2sM}{(\varepsilon T)^2}\right) \left(\int_{\frac{T}{4}}^{\frac{3T}{4}} \int_{-L}^{0} |\varphi|^2 + \int_{\frac{T}{4}}^{\frac{3T}{4}} \int_{\{-L,0\}} |\varphi|^2\right)$$

Consequently we get that

$$\|\varphi\|_{L^2((T/4,3T/4);X)}^2 \lesssim s^4(\varepsilon T)^{-8} \exp\left(\frac{4s(M-m)}{(\varepsilon T)^2}\right) \int_0^T |\varphi|^2(.,0) := C \int_0^T |\varphi|^2(.,0).$$

We now choose $s \sim (\varepsilon T)^2 + (\varepsilon T)$. The above constant C is consequently estimated by

$$\varepsilon^{-4}e^{c/\varepsilon} \lesssim e^{c'/\varepsilon}$$

for c' > c and c well-chosen.

• We now deduce the result using dissipation estimates. We have just proven

$$\|\varphi\|_{L^2((T/4,3T/4);X)}^2 \lesssim e^{\frac{c'}{\varepsilon}} \int_0^T \int_{\Gamma_0} |\varphi|^2.$$

We now use the dissipation property with $t_1 = 0$ and $t_2 = t \in \left[\frac{T}{4}, \frac{3T}{4}\right]$. We easily get, provided T > 4L,

$$\frac{T}{2} \exp\left(\frac{(T-4L)^2}{8\varepsilon T}\right) \|\varphi(0)\|_X^2 \leqslant \|\varphi\|_{L^2((T/4,3T/4);X)}^2,$$

which gives the result with $k = \frac{1}{8} \left(1 - \frac{4L}{T} \right) (T - 4L) - c' > 0$ provided that $\frac{T}{L} > 8 + 32c'$.

3.2 Proof of Theorem 1

To show the controllability result for (S_v) , we will adopt some minimization strategy inspired by the classical heat equation.

Proposition 15 A necessary and sufficient condition for the solution of problem (S_v) to satisfy u(T) = 0 is given by:

$$\forall \varphi_T \in X, < \varphi(0, .), u_0 >_X = \int_0^T \varphi(., 0) v.$$

where φ is the solution of problem (S') with final value φ_T .

Proof We apply the Definition 4 to u against φ a strong solution of (S'), which gives

$$\int_0^T \varphi(.,0)v = \langle \varphi(0,.), u_0 \rangle_X - \langle \varphi_T, u(T,.) \rangle_X$$

and we get the desired equivalence by approximation of weak solutions by strong solutions.

Proposition 16 The following properties are equivalent

- $\exists C_1 > 0, \forall \varphi_T \in X; \|\varphi(0,.)\|_X \leq C_1 \|\varphi(.,0)\|_{L^2(0,T)}$ where φ is the solution of problem (S'),
- $\exists C_2 > 0, \forall u_0 \in X, \exists v \in L^2(0,T) \text{ such that } \|v\|_{L^2(0,T)} \leq C_2 \|u_0\|_X \text{ and the solution } u \text{ of problem } (S_v) \text{ satisfies } u(T) = 0.$

Moreover, $C_1 = C_2$.

Proof (\Rightarrow) Let $u_0 \in X$. We define H as the closure of X for the norm defined by

$$\|\varphi_T\|_H = \|\varphi(.,0)\|_{L^2(0,T)},$$

where φ is the corresponding solution of (S'). Using the observability assumption and backward uniqueness (Proposition 8), one sees that it is indeed a norm on X.

We define a functional J in the following way

$$J(\varphi_T) = \frac{1}{2} \int_0^T \varphi^2(.,0) - \langle \varphi(0,.), u_0 \rangle_X$$

J is clearly convex and our assumption imply that J is continuous on H. Moreover, thanks to our observability assumption, J is coercive. Indeed, one has:

$$J(\varphi_T) \ge \frac{1}{2} \|\varphi_T\|_H^2 - C \|\varphi_T\|_H$$

for $\varphi_T \in H$.

Thus J possesses a global minimum $\widehat{\varphi}_T \in H$, which gives, writing Euler-Lagrange equations,

$$\forall \varphi_T \in H, \int_0^T \varphi(.,0)\hat{\varphi}(.,0) = \langle \varphi(0,.), u_0 \rangle_X.$$

$$(25)$$

According to Proposition 15, we have shown the existence of an admissible control defined by $v = \hat{\varphi}(., 0)$. Moreover, choosing $\varphi_T = \hat{\varphi}_T$ in (25), we obtain the following estimate

$$||v||_{L^2(0,T)} \leq ||u_0||_X ||\hat{\varphi}(0,.)||_X.$$

Using our hypothesis, we are done.

(\Leftarrow) If v is an admissible control with continuous dependence on u_0 , Proposition 15 gives us, for every $\varphi_T \in X$,

$$\langle \varphi(0,.), u_0 \rangle_X = \int_0^T \varphi(.,0) v$$

Choosing now $u_0 = \varphi(0)$ gives us the estimate

$$\|\varphi(0,.)\|_{X}^{2} \leq C_{2} \|\varphi_{T}\|_{H} \|u_{0}\|_{X}$$

that is

$$\|\varphi(0,.)\|_X \leqslant C_2 \|\varphi_T\|_H.$$

Conclusion and open problems

The question of controllability in higher dimension is open and seems to be hard to obtain using Carleman estimates. Another approach to obtain the observability estimate in the *n*-dimensional case could be similar to the one of Miller (see [13, Theorem 1.5]) but it seems that our problem is not adapted to that framework (we can check that the link between the one-dimensional problem and the general one is not as simple as it may seem, see Appendix A below).

If n = 1, we can wonder what is the minimal time to get a vanishing control cost when the viscosity goes to zero. The intuitive result would be L, but it seems that Carleman estimates cannot give such a result.

A On Miller's trick

In this paragraph, we detail why we cannot apply Miller's method to deduce observability in dimension n from the one in dimension 1. We refer the reader to the abstract setting of [13, Section 2] (see more specifically Lemma 2.2 in that reference).

We denote by Ω^1 or Ω^n the domains (resp. X^1 or X^n the function spaces, \mathcal{A}^1 or \mathcal{A}^n the evolution operators) in dimension 1 or n. The observability operator in dimension one is defined by

$$\mathcal{O}^1: \begin{array}{c} \varphi \longmapsto \varphi \mathbf{1}_{\{x=0\}} \\ X^1 \to L^2(\partial \Omega^1) \end{array}.$$

where $1_{\{x=0\}}$ is the characteristic function of $\{x=0\}$. With similar notations, it is obvious that the observability operator on the part of the boundary $\Gamma_0 = \mathbb{R}^{n-1} \times \{0\}$ is given by

$$\mathcal{O}^n: \begin{array}{c} \varphi \longmapsto \varphi \mathbf{1}_{\Gamma_0} \\ X^n \to L^2(\partial \Omega^n) \end{array}$$

and coincides with the tensorial product $I \otimes \mathcal{O}^1$ between the identity of $L^2(\mathbb{R}^{n-1})$ and \mathcal{O}^1 . We have now the following result.

Proposition 17 We define the natural "difference" between \mathcal{A}^n and \mathcal{A}^1 by

$$\mathcal{B}: w \in L^2(\mathbb{R}^{n-1}) \mapsto \varepsilon \Delta' w \in L^2(\mathbb{R}^{n-1})$$

 $(\Delta' \text{ denotes the standard Laplacian on } \mathbb{R}^{n-1})$ with domain

$$\mathcal{D}(\mathcal{B}) = H^2(\mathbb{R}^{n-1}).$$

Then \mathcal{A}^n does not extend $\mathcal{A}^1 \otimes I + I \otimes \mathcal{B}$.

Proof We compute for $u_1 \in \mathcal{D}(\mathcal{A}^1)$, $w \in \mathcal{D}(\mathcal{B})$ and $\varphi_1 \in X^1$, $\phi \in L^2(\mathbb{R}^{n-1})$ the scalar product $< (\mathcal{A}^1 \otimes I + I \otimes \mathcal{B})(u_1 \otimes w), \varphi_1 \otimes \phi >_{X^n}$. This value is successively

$$<\mathcal{A}^{1}u_{1},\varphi_{1}>_{X^{1}}< w,\phi>_{L^{2}(\mathbb{R}^{n-1})}+< u_{1},\varphi_{1}>_{X^{1}}< Bw,\phi>_{L^{2}(\mathbb{R}^{n-1})}$$
$$=\left(\int_{\Omega^{1}}((\varepsilon\partial_{x_{n}}^{2}-\partial_{x_{n}})u_{1})\varphi_{1}-\int_{\partial\Omega^{1}}((\partial_{\nu}+1_{\{x=-L\}})u_{1})\varphi_{1}\right)\int_{\mathbb{R}^{n-1}}w\phi$$
$$+\left(\int_{\Omega^{1}}u_{1}\varphi_{1}+\varepsilon\int_{\partial\Omega^{1}}u_{1}\varphi_{1}\right)\varepsilon\int_{\mathbb{R}^{n-1}}(\Delta'\phi)w$$
$$=<\mathcal{A}^{n}(u_{1}\otimes w),\varphi_{1}\otimes\phi>_{X^{n}}+\varepsilon^{2}\int_{\partial\Omega^{n}}(\Delta'(u_{1}\otimes w))\varphi_{1}\otimes\phi$$

and our claim is proved.

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