Null-controllability of underactuated linear parabolic-transport systems with constant coefficients

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Abstract

The goal of the present article is to study controllability properties of mixed systems of linear parabolic-transport equations, with possibly non-diagonalizable diffusion matrix, on the one-dimensional torus. The equations are coupled by zero or first order coupling terms, with constant coupling matrices, without any structure assumptions on them. The distributed control acts through a constant matrix operator on the system, so that there might be notably less controls than equations, encompassing the case of indirect and simultaneous controllability. More precisely, we prove that in small time, such kind of systems are never controllable in any Sobolev spaces, whereas in large time, null-controllability holds, for sufficiently regular initial data, if and only if a spectral Kalman rank condition is verified. We also prove that initial data that are not regular enough are not controllable. Positive results are obtained by using the so-called fictitious control method together with an algebraic solvability argument, whereas the negative results are obtained by using an appropriate WKB construction of approximate solutions for the adjoint system associated to the control problem. As an application to our general results, we also investigate into details the case of 2×2 systems (*i.e.*, one pure transport equation and one parabolic equation).

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1 Introduction

1.1 Context and state of the art

Controllability properties of coupled systems of PDEs has attracted a lot of attention these last two decades, due to their link with real-life models and also the specific mathematical difficulties arising in this context. An important part of the literature is devoted to systems where all components of the equations have the same qualitative behaviour (meaning that they are for instance all parabolic, or all hyperbolic, etc.). However, the case where different dynamics are mixed has been less studied, despite its mathematical interest. Indeed, in this context, the controllability properties of each equation

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taken separately might be totally different (for instance, the heat equation with distributed control is controllable in arbitrary small time from any open subset [30, 22], whereas the wave equation with distributed control is controllable in large time and under some geometric conditions [6]), so that the controllability properties of the final coupled system might be difficult to guess. Moreover, when we are considering underactuated systems (in the sense that there are less controls than equations) as in the present article, additional mathematical difficulties are appearing, due notably to the algebraic and analytic effects of the coupling terms, that become predominant in the understanding of the controllability or observability properties of the system under study. Here, in the present article, we aim to study the indirect controllability properties of a model of coupled parabolic-transport equations as introduced in [7].

Let us mention that many realistic models already studied in the literature can be reformulated in terms of coupled parabolic-transport equations, notably the wave equation with structural damping [38, 35, 10, 24], the heat equation with memory [26, 23], the 1D-Linearized compressible Navier-Stokes equations [20, 13, 12, 8], or the Benjamin-Bona-Mahony equation [39]. For more details, we also refer to [7, §1.4]. This justifies the interest of studying a general version of coupled parabolic-transport systems as in the present article, that can be seen as an attempt to find a unified framework in order to encompass many existing results of the literature and to generalize them. Other results of interest, related to the present work, are [2], where the authors study a one-dimensional system of one transport equation and one parabolic equation, for which they prove a non-controllability result in small time by a WKB approach, and [11], where the authors prove a controllability result in large time for a one-dimensional system of one transport equation and one elliptic equation.

1.2 Presentation of the parabolic-transport system under study

Let $T > 0$ some final time , $\mathbb{T} = \mathbb{R}/(2\pi\mathbb{Z})$ the one-dimensional torus, ω an nonempty open subset of $\mathbb{T}$, $d \in \mathbb{N}^*$ (which represents the number of equations in our system) , $m \in \{1, \dots, d\}$ (which represents the number of controls in our system), $A, B, K \in \mathcal{M}_d(\mathbb{R})$ (that are some constant coupling matrices), and $M \in \mathcal{M}_{d,m}(\mathbb{R})$ (that is a constant control operator). Our goal is to study the controllability properties of the following coupled system of parabolic-transport equations:

$$\begin{cases} \partial_t f - B\partial_x^2 f + A\partial_x f + Kf = Mu\mathbb{1}_\omega & \text{in } (0,T)\times\mathbb{T}, \\ f(0,\cdot) = f_0 & \text{in } \mathbb{T}. \end{cases} \tag{Sys}$$

Here, the state is $f : [0,T]\times\mathbb{T} \to \mathbb{R}^d$, and the control is $u : [0,T]\times\mathbb{T} \to \mathbb{R}^m$. The exact regularity chosen for f and u will be made more precise later on.

We assume that

$$d = d_{\mathrm{h}} + d_{\mathrm{p}} \text{ with } 1 \le d_{\mathrm{h}} < d,\ 1 \le d_{\mathrm{p}} < d, \tag{H.1}$$

$$B = \begin{pmatrix} 0 & 0 \\ 0 & D \end{pmatrix}, \text{ with } D \in \mathcal{M}_{d_{\mathrm{p}}}(\mathbb{R}), \tag{H.2}$$

$$\Re(\mathrm{Sp}(D)) \subset (0, +\infty). \tag{H.3}$$

d_{h} represents the number of purely hyperbolic equations, whereas d_{p} represents the number of parabolic equations.

Notice that (H.3) is necessary to ensure that the matrix operator $\partial_t - D\Delta$ is parabolic is the sense of Petrovskii ([29, Chapter 7, Definition 2]). Introducing the similar block decomposition for the $d \times d$ matrix $A = \left(\begin{smallmatrix} A' & A_{12} \\ A_{21} & A_{22} \end{smallmatrix}\right)$, we make the following hypothesis on the matrix $A' \in \mathcal{M}_{d_h}(\mathbb{R})$

$$A' \text{ is diagonalizable with } \mathrm{Sp}(A') \subset \mathbb{R}. \tag{H.4}$$

Notice that it is well-known that (H.4) is necessary (and sufficient, see [7, §2.2]) to ensure the well-posedness of (Sys).

1.3 Main results

To state our results, we need to introduce the following notations:

$$\ell(\omega) := \sup\{|I|;\ I \text{ connected component of } \mathbb{T} \setminus \omega\}, \tag{1}$$

$$\mu_* := \min\{|\mu|;\ \mu \in \mathrm{Sp}(A')\},$$

and

$$T^* = T^*(\omega) := \begin{cases} \frac{\ell(\omega)}{\mu_*} & \text{if } \mu_* > 0, \\ +\infty & \text{if } \mu_* = 0. \end{cases} \tag{2}$$

For $n \in \mathbb{Z}$, we also set

$$B_n := -n^2 B - inA - K \tag{3}$$

and

$$[B_n|M] := \begin{pmatrix} M & B_n M & \dots & B_n^{d-1} M \end{pmatrix} \in \mathcal{M}_{n,nm}(\mathbb{C}). \tag{4}$$

Our main result is the following one.

Theorem 1. *Assume that the hypotheses* (H.1)–(H.4) *hold, and that* $T > T_*$.

Then, the spectral Kalman rank condition $\mathrm{rank}([B_n|M]) = d$ *holds for all* $n \in \mathbb{Z}$ *if and only if for every* $f_0 \in H^{4d(d-1)}(\mathbb{T})^d$*, there exists a control* $u \in L^2([0,T] \times \omega)^m$ *such that the solution* f *of the parabolic-transport system* (Sys) *with initial condition* f_0 *satisfies* $f(T,\cdot) = 0$.

Remark 2.

- Recall that the Kalman rank condition is necessary for the control of ODE systems [27]. Therefore, writing the parabolic-transport system in the Fourier variable, we immediately find that for every $T > 0$, the spectral Kalman-rank condition "$\forall n \in \mathbb{Z},\ \mathrm{rank}([B_n|M]) = d$" is necessary for the null-controllability of every H^k initial conditions in time T, for any $k \in \mathbb{N}$.

- Actually, we prove two slightly stronger versions of this theorem, namely theorems 12 and 15, that are useful in order to obtain some controllability results when the system has some invariant, *e.g.*, when the average of a component is invariant whatever the control is. Theorem 12 also gives better regularity assumptions. We refer to propositions 22 to 24 for examples of applications of these refined statements.

- One can refine a little bit the regularity stated in theorem 1, as follows. Assume that $T > T_*$ and that for all $n \in \mathbb{Z}$, the spectral Kalman rank condition $\mathrm{rank}([B_n|M]) = d$ holds. Then, for every $f_0 \in H^{4d(d-1)}(\mathbb{T})^{d_h} \times H^{4d(d-1)-1}(\mathbb{T})^{d_p}$, there exists a control $u \in L^2([0,T] \times \omega)^m$ such that the solution f of the parabolic-transport system (Sys) with initial condition f_0 satisfies $f(T,\cdot) = 0$.

  Indeed, by letting evolve the system freely on a short interval of time, we can show using the method of lemma 25 that the parabolic component becomes $H^{4d(d-1)}(\mathbb{T})^{d_p}$, so that theorem 1 can be applied, taking into account that the condition $T > T^*$ is open and that the system is time-invariant.

  In general, if one identifies a subspace $V \subset L^2(\mathbb{T})^d$ such that $\mathrm{e}^{-t\mathcal{L}} V \subset H^{4d(d-1)}(\mathbb{T})^d$ for all $t > 0$, then, assuming the spectral Kalman rank condition, we can steer every $f_0 \in V$ to 0 with L^2 controls.

- The spectral Kalman rank condition $\mathrm{rank}([B_n|M]) = d$ was first introduced in [5] for coupled systems of heat equations with diagonalizable diffusions (see also [34] for non-diagonalizable diffusions).

We also prove the following necessary condition on the regularities of solutions and also the minimal time of controllability for the null-controllability to hold.

Theorem 3. *Let $\mu \in \mathrm{Sp}(A')$, $N \in \mathbb{N}$ and $T > 0$. Assume that every initial condition in $L^2(\mathbb{T})^d \cap \{f_0 \in L^2(\mathbb{T})^d : \forall |n| \leq N, (f_0, e_n)_{L^2} = 0\}$ is steerable to 0 in time T with control in $L^2((0,T) \times \omega)$. Then, there exists $V_0 \in \ker(A'^* - \mu)$ such that $M^*\left(\begin{smallmatrix} V_0 \\ 0 \end{smallmatrix}\right) \neq 0$.*

Remark 4. Theorems 1, 12 and 15 only ensure null-controllability of smooth enough initial conditions. Theorem 3 proves that such a regularity condition is needed in general: even if the time is large enough and if the Kalman rank condition is satisfied for every n, it might happen that some L^2 initial condition (or even L^2 initial conditions that only have high frequencies) cannot be steered to 0 with a L^2 control.

Theorem 5. *Let $T > 0$ and assume that there exists $N \in \mathbb{N} \setminus \{0\}$ and $k \in \mathbb{N}$ such that every initial conditions in $H^k(\mathbb{T})^d \cap \{\sum_{|n|>N} X_n e^{inx}\}$ for the parabolic-transport system (Sys) can be steered to 0 in time T. Then $T \geq T_*$.*

1.4 Precise scope and organization of the article

This article can be seen as a continuation of [7], insofar as we generalize the results of the above-mentioned article, since we are able to treat any matrices A, B, K, M without any restrictions on their structure. Indeed, in [7], the authors treated the case where $M = I_d$ (where no Kalman rank condition is needed), or particular cases where only the parabolic or the hyperbolic parts are controlled, under strong restrictions on the structure of the coupling matrices A, B and K and also on the diffusion matrix B.

Let us mention that our results are sharp in terms of the controllability conditions we obtain. However, it is very likely that the initial state space (whose choice is determined by technical reasons coming from the specific strategy we use, that is consuming in terms of regularity, see Section 3.2) is almost never sharp and depends strongly on the structure of the coupling terms. Finding the exact "good" state space remains an open problem that seems difficult to solve in all generality.

The article is organized as follows. In section 2, we give some notations and we gather some existing results that will be used in our proof. Section 3 is devoted to proving that the condition $\mathrm{rank}([B_n|M]) = d$ is sufficient in order to obtain our desired controllability result in large time. The argument is based on a fictitious control argument detailed in section 3.1, where we first prove an auxiliary controllability result, in the case $M = I_d$, with regular enough controls for regular enough initial data. Then, in section 3.2, we explain how to obtain a control in the range on M by performing algebraic manipulations. Notice that the method of fictitious control plus algebraic solvability, that has been introduced in [16] in the context of the controllability of PDEs, has been successfully used for various (systems of) PDEs [4, 18, 19, 33, 15, 40, 41, 17]. One of the main novelties here is that the algebraic solvability is not directly performed on the system (or its adjoint as in [17]) but on a projected version of the system on its Fourier components. Section 4 is devoted to proving some necessary conditions of controllability. Section 4.1 is devoted to constructing WKB solutions. These solutions are used to disprove controllability in small time in section 4.2 and to prove theorem 3 in section 4.3. Section 5 aims to give an application of our results to the particular case of 2×2 systems together with some considerations about the sharpness of our regularity assumptions in this precise setting. To conclude, appendix A proves a general result about a "control up to a finite-dimensional space plus unique continuation" strategy that is used in section 3.1, in the spirit of [31, 7].

2 Some notations and preliminary results

We will rely on some basic results on the parabolic-transport system (Sys) that are already known [7]. For the reader convenience, we collect here the notations and results we will use most often, and we will recall some others along the way as they are used.

Let $\mathcal{L}$ be the unbounded operator on $L^2(\mathbb{T})^d$ with domain $H^1(\mathbb{T})^{d_\mathrm{h}} \times H^2(\mathbb{T})^{d_\mathrm{p}}$ defined by

$$\mathcal{L}f = -B\partial_x^2 f + A\partial_x f + Kf.$$

The operator $-\mathcal{L}$ generates a strongly continuous semigroup of bounded operators of $L^2(\mathbb{T})^d$ [7, Proposition 11]. Every $H^k(\mathbb{T})^d$ is stable by $\mathrm{e}^{-t\mathcal{L}}$, and the restriction of $\mathrm{e}^{-t\mathcal{L}}$ on $H^k(\mathbb{T})^d$ is a strongly continuous semigroup of bounded operators [7, Remark 13]. We denote by $S(T,\ f_0,\ u)$ the solution at time T of the parabolic-transport system (Sys) with control matrix $M = I_d$ (the identity matrix of size d, *i.e.*, we control every component with a different control), initial condition f_0 and control u.

Let $n_0 \in \mathbb{N}$ to be chosen large enough later on. For $n \in \mathbb{Z}$, we denote by $e_n : x \in \mathbb{T} \mapsto \mathrm{e}^{\mathrm{i}nx}$. We also denote by $E : \mathbb{C} \to \mathcal{M}_d(\mathbb{C})$ the following function:

$$E(z) = B + zA - z^2K.$$

Let $r > 0$ small enough. For $|z| < r$, let $P^\mathrm{h}(z)$ be the eigenprojection on the sum of generalized eigenspaces of $E(z)$ associated to the set of eigenvalues $\lambda(z) \in \mathrm{Sp}(E(z))$ such that $|\lambda(z)| < r$. According to [7, Proposition 5], $P^\mathrm{h}(z)$ satisfies:

- $P^\mathrm{h}(0) = \left(\begin{smallmatrix} I & 0 \\ 0 & 0 \end{smallmatrix}\right)$;
- $z \mapsto P^\mathrm{h}(z)$ is holomorphic;
- $P^\mathrm{h}(z)$ is a projection that commutes with $E(z)$;
- $P^\mathrm{h}(z)E(z) = O(z)$ as $z \to 0$.

We also set $P^\mathrm{p}(z) = I - P^\mathrm{h}(z)$. This projection $P^\mathrm{p}(z)$ satisfies similar properties as $P^\mathrm{h}(z)$ ([7, Proposition 6]).

Following [7, Proposition 18], we denote by F^0 the space of frequencies less than n_0 and by F^h (respectively F^p) the space of hyperbolic frequencies greater than n_0 (respectively the space of parabolic frequencies greater than n_0), *i.e.*

$$\begin{aligned} F^0 &:= \bigoplus\nolimits_{|n|\le n_0} \mathbb{C}^d e_n; \\ F^\mathrm{p} &:= \bigoplus\nolimits_{|n|>n_0} \mathrm{Range}(P^\mathrm{p}(\mathrm{i}/n))e_n; \\ F^\mathrm{h} &:= \bigoplus\nolimits_{|n|>n_0} \mathrm{Range}(P^\mathrm{h}(\mathrm{i}/n))e_n. \end{aligned} \tag{5}$$

By [7, Proposition 18], we notably have

$$L^2(\mathbb{T})^d = F^0 \oplus F^\mathrm{p} \oplus F^\mathrm{h}.$$

The space F^p is stable by the semigroup $\mathrm{e}^{-t\mathcal{L}}$ (see the definition of P^p [7, Proposition 5] and the definition of F^p [7, Proposition 18]). We denote by $\mathcal{L}^\mathrm{p}$ the restriction of $\mathcal{L}$ to F^p.

Similarly, the space F^h is stable by the semigroup $\mathrm{e}^{-t\mathcal{L}}$. We denote by $\mathcal{L}^\mathrm{h}$ the restriction of $\mathcal{L}$ to F^h, and $-\mathcal{L}^\mathrm{h}$ generates a strongly continuous *group* of bounded operators on F^h [7, Proposition 19].

Let Π^0, Π^p, Π^h and Π be the projections defined by

$$\begin{aligned} L^2(\mathbb{T})^d &= F^0 \oplus F^\mathrm{p} \oplus F^\mathrm{h}; \\ \Pi^0 &= I_{F^0} + 0 + 0; \\ \Pi^\mathrm{p} &= 0 + I_{F^\mathrm{p}} + 0; \\ \Pi^\mathrm{h} &= 0 + 0 + I_{F^\mathrm{h}}; \\ \Pi &= 0 + I_{F^\mathrm{p}} + I_{F^\mathrm{h}} = \Pi^\mathrm{p} + \Pi^\mathrm{h}. \end{aligned}$$

These projections are bounded operators on $L^2(\mathbb{T})^d$ [7, Proposition 18] (and also on every $H^k(\mathbb{T})^d$, as one can readily convince by following the proof of [7, Proposition 18]).

3 Null controllability of regular initial conditions

3.1 Regular controls for regular initial conditions

As a technical preparation for the proof of theorem 1, we need some results regarding the regularity of controls, when the control matrix is $M = I_d$.

Proposition 6. *Assume that $T > T_*$ (as defined in eq. (2)) and that $M = I_d$. Let $k, \ell \in \mathbb{N}$. For every $f_0 \in H^k(\mathbb{T})^d$, there exists $u \in H_0^k((0,T)\times\omega)^{d_\mathrm{h}} \times H_0^\ell((0,T)\times\omega)^{d_\mathrm{p}}$ such that the solution of the parabolic-transport system (Sys) with initial condition f_0 and control u satisfies $f(T,\cdot) = 0$.*

Remark 7. The regularity of the controls on the hyperbolic part is limited by the regularity of the initial conditions of the hyperbolic part, whereas the controls on the parabolic part can be in arbitrary Sobolev spaces, independently on the regularity of the initial conditions of the parabolic part.

We adapt the proof of the corresponding result when $k = 0$ [7, Theorem 2]. First, we prove the following adaptation of [7, Proposition 21].

Proposition 8. *Let $T' \in (T^*, T)$ and $k \in \mathbb{N}$. If n_0 (in the definition of F^0, see eq. (5)) is large enough, there exists a continuous operator*

$$\begin{aligned} \mathcal{U}^\mathrm{h} : H^k(\mathbb{T})^d \times H_0^k((T',T)\times\omega)^{d_\mathrm{p}} &\to H_0^k((0,T')\times\omega)^{d_\mathrm{h}} \\ (f_0, u_\mathrm{p}) &\mapsto u_\mathrm{h}, \end{aligned}$$

such that for every $(f_0, u_\mathrm{p}) \in H^k(\mathbb{T})^d \times H_0^k((T',T)\times\omega)^{d_\mathrm{p}}$ (where u_p is extended by 0 on $(0,T')$ and u_h is extended by 0 on (T',T)),

$$\Pi^\mathrm{h} S(T; f_0, (\mathcal{U}^\mathrm{h}(f_0, u_\mathrm{p}), u_\mathrm{p})) = 0.$$

Proof. As in [7, §4.3.1], the conclusion of proposition 8 is equivalent to the exact controllability of the system $\partial_t f + \mathcal{L}^\mathrm{h} f = \Pi^\mathrm{h}(u, 0)$ at time T'. Since $-\mathcal{L}^\mathrm{h}$ generates a strongly continuous *group*, the exact controllability at time T' is equivalent to the null-controllability at time T', which is what we are going to prove.

When $k = 0$, [7, Proposition 23] is the claimed result. To extend this result to $k > 0$, we use a general result of Ervedoza and Zuazua concerning the regularity of controls for regular initial data in the context of groups of operators [21, Theorem 1.4]. Let $\widetilde{\omega}$ an open subset of $\mathbb{T}$ such that $\overline{\widetilde{\omega}} \subset \omega$ and $T^*(\widetilde{\omega}) < T'$. Let $\chi \in C_c^\infty(\omega)$ such that $\chi = 1$ on $\widetilde{\omega}$. Let $\eta \in C_0^\infty(0,T')$. Let $z_0 \in H^k(\mathbb{T})^d$ be an initial condition. Let $Y_{T'}$ as defined by [21, Proposition 1.3] and define the control as

$$V(t) = \eta(t)\chi(x)M^* Y(t),$$

where Y is the solution to

$$\partial_t Y - B^* \partial_x^2 Y - A^* \partial_x Y + K^* Y = 0$$

associated to the initial condition $Y(T') = Y_{T'}$. According to [21, Proposition 1.3], $V(t)$ is a control that steers z_0 to 0 at time T'. According to [21, Theorem 1.4], $Y_{T'} \in H^k(\mathbb{T})^d$ (hence $V \in L^2(0,T'; H^k(\omega)^d)$) and $V \in H^k(0,T'; L^2(\omega)^d)$, with estimates of the form

$$\|V\|^2_{L^2(0,T';H^k(\omega)^d)} + \|V\|^2_{H^k(0,T';L^2(\omega)^d)} \le C_k \|z_0\|^2_{H^k(\mathbb{T})^d}.$$

We claim that $L^2(0,T';\ H_0^k(\omega)) \cap H_0^k(0,T';\ L^2(\omega)) \subset H^k((0,T')\times\omega)$. Indeed, for every $\tau \in \mathbb{R}$ and $\xi \in \mathbb{R}$,

$$(1+\tau^2+\xi^2)^k \le C_k\big((1+\tau^2)^k + (1+\xi^2)^k\big).$$

Hence, integrating in Fourier space, for any $f \in L^2(0,T';\ H_0^k(\omega)) \cap H_0^k(0,T';\ L^2(\omega))$, we have

$$\|f\|^2_{H^k(\mathbb{R}^2)} \le C_k\big(\|f\|^2_{L^2(\mathbb{R};H^k(\mathbb{R}))} + \|f\|^2_{H^k(\mathbb{R};L^2(\mathbb{R}))}\big),$$

so that $f \in H^k((0,T')\times\omega)$.

Recall that for $\Omega \subset \mathbb{R}^n$ convex[1], $H_0^k(\Omega)$ is the set of functions whose extension by zero outside Ω are $H^k(\mathbb{R}^n)$. Hence, $L^2(0,T';\ H_0^k(\omega)) \cap H_0^k(0,T';\ L^2(\omega)) \subset H_0^k((0,T')\times\omega)$ as claimed, so that $V \in H_0^k((0,T')\times\omega)^d$. □

For the proof of proposition 6, we will also use:

Proposition 9 ([7], proposition 22). *Let $T' \in (T^*, T)$ and $k \in \mathbb{N}$. If n_0 is large enough, there exists a continuous operator*

$$\begin{aligned} \mathcal{U}^{\mathrm{p}} :\ L^2(\mathbb{T})^d \times L^2((0,T')\times\omega)^{d_{\mathrm{h}}} &\to C_c^\infty((T',T)\times\omega)^{d_{\mathrm{p}}} \\ (f_0, u_{\mathrm{h}}) &\mapsto u_{\mathrm{p}}, \end{aligned}$$

(in the sense that for any $s \in \mathbb{N}$, $\mathcal{U}^{\mathrm{p}} : L^2(\mathbb{T})^d \times L^2((0,T')\times\omega)^{d_{\mathrm{h}}} \to H_0^s(T',T)\times\omega)^{d_{\mathrm{p}}}$ is continuous for the natural topologies associated to these spaces) such that for every $(f_0, u_{\mathrm{h}}) \in L^2(\mathbb{T})^d \times L^2((0,T')\times\omega)^{d_{\mathrm{h}}}$,

$$\Pi^{\mathrm{p}} S(T; f_0, (u_{\mathrm{h}}, \mathcal{U}^{\mathrm{p}}(f_0,\ u_{\mathrm{h}})) = 0.$$

We can now prove proposition 6 by mimicking the proof of the case $k = 0$ [7, Proposition 20 & §4.5].

Proof of proposition 6. Step 1: Control up to final dimensional space. — We claim that there exists a closed finite codimensional space $\mathcal{G}$ of $H^k(\mathbb{T})^d$ and a continuous operator $\mathcal{U} :\ \mathcal{G} \to H_0^k((0,T')\times\omega)^{d_{\mathrm{h}}} \times C_c^\infty((T',T)\times\omega)^{d_{\mathrm{p}}}$ (in the sense that for any $s \in \mathbb{N}$, $\mathcal{U} :\ \mathcal{G} \to H_0^k((0,T')\times\omega)^{d_{\mathrm{h}}} \times H_0^s(T',T)\times\omega)^{d_{\mathrm{p}}}$ is continuous for the natural topologies associated to these spaces) such that for every $f_0 \in \mathcal{G}$, $\Pi S(T,\ f_0,\ \mathcal{U} f_0) = 0$.

The property $\Pi S(T,\ f_0,\ (u_{\mathrm{h}}, u_{\mathrm{p}})) = 0$ holds if

$$\begin{cases} u_{\mathrm{h}} = \mathcal{U}^{\mathrm{h}}(f_0,\ u_{\mathrm{p}}) = \mathcal{U}_1^{\mathrm{h}}(f_0) + \mathcal{U}_2^{\mathrm{h}}(u_{\mathrm{p}}), \\ u_{\mathrm{p}} = \mathcal{U}^{\mathrm{p}}(f_0,\ u_{\mathrm{h}}) = \mathcal{U}_1^{\mathrm{p}}(f_0) + \mathcal{U}_2^{\mathrm{p}}(u_{\mathrm{h}}). \end{cases} \tag{6}$$

Set $\mathcal{C} = \mathcal{U}_1^{\mathrm{p}} + \mathcal{U}_{\mathrm{p}}^2 \mathcal{U}_1^{\mathrm{h}}$. Then, the previous relations hold if

$$\mathcal{C} f_0 = (I - \mathcal{U}_2^{\mathrm{p}} \mathcal{U}_2^{\mathrm{h}}) u_{\mathrm{p}}. \tag{7}$$

Since $\mathcal{U}_2^{\mathrm{p}}$ is continuous from $H_0^k((T',T)\times\omega)^{d_{\mathrm{p}}}$ into $C_c^\infty((T',T)\times\omega)^{d_{\mathrm{p}}}$, we deduce that the operator $\mathcal{U}_2^{\mathrm{p}}\mathcal{U}_2^{\mathrm{h}} :\ H_0^k((T',T)\times\omega)^{d_{\mathrm{p}}} \to H_0^k((T',T)\times\omega)^{d_{\mathrm{p}}}$ is compact. Thus, according to Fredholm's alternative, the relation (7) holds on a closed finite codimensional space $\mathcal{G}$.

Step 2: Conclusion. — Dealing with the finite (co)dimensional spaces F^0 and $\mathcal{G}$ is a straightforward adaptation of [7, §4.5]; more specifically, we use proposition 27 proved in Appendix A with $H = V = H^k(\mathbb{T})^d$, $U_T = H_0^k((0,T)\times\omega)^{d_{\mathrm{h}}} \times H_0^\ell((0,T)\times\omega)$, $A = -\mathcal{L}$, $B = \mathbb{1}_\omega$, $\mathcal{G} = \mathcal{G}$ and $\mathcal{F} = F^0$. The control up to a finite dimensional space hypothesis is satisfied according to the previous step. The unique

[1] More generally, satisfiying the *segment condition*, see[1, Definition 3.21 & Theorem 5.29].

continuation hypothesis is satisfied because every generalised eigenvector is a finite sum of elements of the form $X_n e^{inx}$ ($X_n \in \mathbb{C}^d$), and finite linear combinations of $X_n e^{inx}$ have the unique continuation property thanks to, *e.g.*, Jerison-Lebeau's spectral inequality (see [31, Theorem 3], or [7, Eq. (90)] for our specific case). □

For technical reasons that will be explained later on, we will need the control to be in the form $P(\partial_x)u$, where $P(\partial_x)$ is a constant coefficients differential operator to be chosen later on.

Proposition 10. *Assume that $T > T_*$ (as defined in (2)) and that $M = I_d$. Let $k, \ell \in \mathbb{N}$. Let P be a nonzero polynomial with complex coefficients. Assume that $\ell \geqslant \deg(P)$. Let $f_0 \in H^k(\mathbb{T})^d$ be such that for every $n \in \mathbb{Z}$, $P(in) = 0 \implies c_n(f_0) = 0$. Then, there exists $u \in H_0^{k+\deg(P)}((0,T)\times\omega)^{d_\mathrm{h}} \times H_0^{\ell}((0,T)\times\omega)^{d_\mathrm{p}}$ such that the solution of the parabolic-transport system (Sys) with initial condition f_0 and control $P(\partial_x)u$ satisfies $f(T,\cdot) = 0$.*

Remark 11. The condition '$n \in \mathbb{Z}$, $P(in) = 0 \implies c_n(f_0) = 0$" is necessary to ensure some kind of right invertibility for $P(\partial_x)$, as it will appear clearly during the proof.

Proof. Let $k, \ell \in \mathbb{N}$ with $\ell \geqslant \deg(P)$. Let $f_0 \in H^k(\mathbb{T})^d$ be such that for every $n \in \mathbb{Z}$, $P(in) = 0 \implies c_n(f_0) = 0$. We define $\widetilde{f_0} := P(\partial_x)^{-1} f_0$ by $c_n(\widetilde{f_0}) := P(in)^{-1} c_n(f_0)$ if $P(in) \neq 0$ and $c_n(\widetilde{f_0}) := 0$ if $P(in) = 0$. Note that $P(\partial_x)\widetilde{f_0} = f_0$ and that $\widetilde{f_0} \in H_0^{k+\deg(P)}(\omega)^d$. Then, applying proposition 6 to $\widetilde{f_0}$ leads to the fact that there exists $\tilde{u} \in H_0^{k+\deg(P)}((0,T)\times\omega)^{d_\mathrm{h}} \times H_0^{\ell}((0,T)\times\omega)^{d_\mathrm{p}}$ such that the solution $\widetilde{f}$ of the parabolic-transport system (Sys) with initial condition $\widetilde{f_0}$ and control $\tilde{u}$ satisfies $\widetilde{f}(T,\cdot) = 0$. Moreover, since $\widetilde{f_0} \in H_0^{k+\deg(P)}(\omega)^d$ and $\tilde{u} \in H_0^{k+\deg(P)}((0,T)\times\omega)^{d_\mathrm{h}} \times H_0^{\ell}((0,T)\times\omega)^{d_\mathrm{p}}$ with $\ell \geqslant \deg(P)$, we notably have $\widetilde{f} \in L^2((0,T); H^{k+\deg(P)}(\mathbb{T}))$. Hence, setting $f = P(\partial_x)\widetilde{f}$ and $u = P(\partial_x)\tilde{u}$, and using that $P(\partial_x)$ has constant coefficients (so that it commutes with the operator $\partial_t - B\partial_x^2 + A\partial_x + K$) ensures that f satisfies (Sys) with initial condition f_0 and control $P(\partial_x)u$. Moreover, since $\widetilde{f}(T,\cdot) = 0$, we also have $f(T,\cdot) = P(\partial_x)\widetilde{f}(T,\cdot) = 0$, which leads to the desired result. □

3.2 Algebraic solvability

For $k \in \mathbb{N} \setminus \{0\}$, we define

$$[B_n|M]_k := \begin{pmatrix} M & B_n M & \dots & B_n^{k-1} M \end{pmatrix}. \tag{8}$$

We prove the following variant of theorem 1.

Theorem 12. *Assume that the hypotheses (H.1)–(H.4) hold, and that $T > T_*$. Let $k \in \mathbb{N}$. Assume that for all $|n| \in \mathbb{N}$ large enough, the Kalman rank condition* $\operatorname{rank}([B_n|M]_k) = d$ *holds. Define the following space of functions*

$$E := \{f \in L^2(\mathbb{T})^d : \forall n \in \mathbb{Z},\ c_n(f) \in \operatorname{Range}([B_n|M])\}.$$

Set, when it is defined,

$$[B_n|M]_k^+ := [B_n|M]_k^* \left([B_n|M]_k [B_n|M]_k^*\right)^{-1}.$$

Write $[B_n|M]_k^+$ by blocks as

$$[B_n|M]_k^+ = \begin{pmatrix} L_{n,1}^\mathrm{h} & L_{n,1}^\mathrm{p} \\ \vdots & \vdots \\ L_{n,k}^\mathrm{h} & L_{n,k}^\mathrm{p} \end{pmatrix},$$

where the $L_{n,j}^\mathrm{h}$ are of size $m \times d_\mathrm{h}$ and the $L_{n,j}^\mathrm{p}$ are of size $m \times d_\mathrm{p}$. Considering the $L_{n,j}^\mathrm{h}$ as rational functions of n, and denoting their degree by $\deg(L_{n,j}^\mathrm{h})$*, set*

$$p := \max_{1\leq j\leq k} \deg(n^{j-1} L_{n,j}^\mathrm{h}) = \max_{1\leq j\leq k} \left(j - 1 + \deg(L_{n,j}^\mathrm{h})\right).$$

Then, for every $f_0 \in H^p(\mathbb{T})^d \cap E$, there exists a control $u \in L^2([0,T]\times\omega)$ such that the solution f of the parabolic-transport system (Sys) with initial condition f_0 satisfies $f(T,\cdot) = 0$.

The idea of the proof is to first choose a "fictitious" control that acts on every components. Then, we look at the Fourier coefficients of f. This transforms the control system (Sys) into a family of finite-dimensional control systems. On each of these finite-dimensional system, we perform some algebraic manipulations, called algebraic solvability, that transform the fictitious control (that acted on every component) into an "actual" control (that acts only on Range(M)).

We begin with the algebraic solvability result we will use, which is essentially taken from [33, §2.1].

Lemma 13. *Let $k \in \mathbb{N}\backslash\{0\}$. Let $\widetilde{B} \in \mathcal{M}_d(\mathbb{C})$ and $\widetilde{M} \in \mathcal{M}_{m,d}(\mathbb{R})$. Let $X_0 \in \mathbb{C}^d$ and $w \in H_0^{k-1}(0,T;\mathbb{C}^{mk})$. Write w by blocks as*

$$w = \begin{pmatrix} w_1 \\ \vdots \\ w_k \end{pmatrix},$$

where $w_j \in H_0^{k-1}(\mathbb{T};\mathbb{C}^m)$, and set $u = w_1 + w_2' + \cdots + w_k^{(k-1)}$. Let $X, \widetilde{X} \in C^0(0,T;\mathbb{C}^d)$ be the solutions of

$$X' = \widetilde{B}X + [\widetilde{B}|\widetilde{M}]_k w, \qquad \widetilde{X}' = \widetilde{B}\widetilde{X} + \widetilde{M}u, \qquad X(0) = \widetilde{X}(0) = X_0,$$

where

$$[\widetilde{B}|\widetilde{M}]_k := \begin{pmatrix} \widetilde{M} & \widetilde{B}\widetilde{M} & \cdots & \widetilde{B}^{k-1}\widetilde{M} \end{pmatrix}.$$

Then $X(T) = \widetilde{X}(T)$.

Proof. Consider $\widetilde{\mathcal{M}}_k$ the operator matrix with $d+m$ rows and km columns defined by blocks as

$$\widetilde{\mathcal{M}}_k := \begin{pmatrix} 0 & -\widetilde{M} & \cdots & -\sum_{j=0}^{k-2} \partial_t^j \widetilde{B}^{k-2-j}\widetilde{M} \\ -I & -\partial_t & \cdots & -\partial_t^{k-1} \end{pmatrix} = \begin{pmatrix} \widetilde{\mathcal{M}}_{k,1} \\ \widetilde{\mathcal{M}}_{k,2} \end{pmatrix}.$$

Set also

$$\mathcal{P}:\ (X,W) \in H_0^1(0,T;\mathbb{C}^d) \times L^2(0,T;\mathbb{C}^m) \to \partial_t X - \widetilde{B}X - \widetilde{M}W \in L^2(0,T;\mathbb{C}^d).$$

We claim that

$$\mathcal{P} \circ \widetilde{\mathcal{M}}_k = [\widetilde{B}|\widetilde{M}]_k. \tag{9}$$

Indeed, we have by blocks $\mathcal{P} \circ \widetilde{\mathcal{M}}_k = \begin{pmatrix} C_0 & \cdots & C_{k-1} \end{pmatrix}$ with

$$C_\ell = -(\partial_t - \widetilde{B}) \sum_{j=0}^{\ell-1} \partial_t^j \widetilde{B}^{\ell-1-j}\widetilde{M} + \widetilde{M}\partial_t^\ell.$$

Then, remarking that this is a telescoping sum,

$$\begin{aligned} C_\ell &= -\sum_{j=1}^{\ell} \partial_t^j \widetilde{B}^{\ell-j}\widetilde{M} + \sum_{j=0}^{\ell-1} \partial_t^j \widetilde{B}^{\ell-j}\widetilde{M} + \widetilde{M}\partial_t^\ell \\ &= -\partial_t^\ell \widetilde{M} + \widetilde{B}^\ell \widetilde{M} - \widetilde{M}\partial_t^\ell, \end{aligned}$$

which proves the claimed formula (9).

Now, plug eq. (9) into the differential equation $X' = \widetilde{B}X + [\widetilde{B}|\widetilde{M}]_k w$, which gives

$$X' = \widetilde{B}X + (\partial_t - \widetilde{B})\widetilde{\mathcal{M}}_{k,1} w - \widetilde{M}\widetilde{\mathcal{M}}_{k,2} w.$$

With $Y := X - \widetilde{\mathcal{M}}_{k,1}w$, and remarking that $\widetilde{\mathcal{M}}_{k,2}w = -u$, this can be written as $Y' = \widetilde{B}Y + \widetilde{M}u$. Since $w \in H_0^{k-1}(0,T;\mathbb{C}^{mk})$, $\widetilde{\mathcal{M}}_{k,1}w(0) = \widetilde{\mathcal{M}}_{k,1}w(T) = 0$. Hence $Y(0) = X(0) = \widetilde{X}(0)$ and $Y(T) = X(T)$. Thus Y solves the same Cauchy problem as $\widetilde{X}$. This proves that $Y = \widetilde{X}$, hence $\widetilde{X}(T) = Y(T) = X(T)$. □

We can now prove theorem 12.

Proof of theorem 12. Let $f_0 \in H^p(\mathbb{T})^d$. Set $X_n(t) = c_n(f(t,\cdot))$ and $u_n(t) = c_n(u(t,\cdot))$. The desired conclusion $f(T,\cdot) = 0$ reads in Fourier as: $\forall n \in \mathbb{Z}, X_n(T) = 0$. Moreover, X_n satisfies

$$\begin{cases} X_n'(t) = B_n X_n(t) + M u_n(t), & t \in (0,T), \\ X_n(0) = c_n(f_0). \end{cases} \tag{10}$$

First, let us give the idea of the proof: if v steers f_0 to 0 when $M = I$, we want to define w_n by $c_n(v(t,\cdot)) = [B_n|M]_k w_n$ (this is possible for n large enough) and choose $u_n := w_{n1} + w_{n2}' + \cdots + w_{nk}^{(k-1)}$. Then, according to lemma 13, the function u_n steers X_n from $c_n(f_0)$ to 0. There are two problems with this crude choice of u_n: this construction only works for n large enough, and more importantly, we have no guarantee that the support of $\sum u_n e^{inx}$ is included in $[0,T] \times \omega$: a natural construction is to chose $w_n := [B_n|M]_k^+ c_n(v(t,\cdot))$, that is a rational function and does not preserve the support. We refine this construction in order to get rid of the denominator of $[B_n|M]_k^+$ in order to use that Fourier multipliers that are polynomial preserve the support. This is here that we need proposition 10.

Step 1: Control of a finite number of frequencies. — To implement the strategy outlined above, we need to deal with low frequencies first. Recall that Π is the projection on frequencies larger than n_0 and that E was defined in the statement of theorem 12. We claim that for any $n_0 \in \mathbb{N} \setminus \{0\}$, $\epsilon > 0$ and $f_0 \in E$, there exists $u \in L^2(0,\epsilon;\, H_0^p(\omega))^m$ such that $(1-\Pi)S(\epsilon, f_0, Mu) = 0$.

This property is equivalent to the null-controllability of the system (Sys) projected on frequencies less or equal than n_0. The observability inequality associated with this problem [14, Theorem 2.44] is:

$$\exists C > 0,\ \forall g_0 \in (1-\Pi)E,\ \|e^{-\epsilon\mathcal{L}^*} g_0\|^2_{H^{-p}(\mathbb{T})^d} \le C \int_0^\epsilon \|M^* e^{-t\mathcal{L}^*} g_0\|^2_{L^2(\omega)^m}\, dt.$$

Since $(1-\Pi)E$ is finite dimensional, this is equivalent to the unique continuation property

$$\forall g_0 \in (1-\Pi)E,\ \left(M^* e^{-t\mathcal{L}^*} g_0(x) = 0 \text{ for } (t,x) \in (0,\epsilon) \times \omega\right) \implies g_0 = 0.$$

Let us prove this property. Let $g_0 \in (1-\Pi)E$ be such that $M^* e^{-t\mathcal{L}^*} g_0(x) = 0$ for $(t,x) \in (0,\epsilon) \times \omega$. Since finite sums of e^{inx} have the unique continuation property, we have $M^* e^{-t\mathcal{L}^*} g_0 = 0$ for $(t,x) \in (0,\epsilon) \times \mathbb{R}$. Hence, for every $0 < t < \epsilon$ and $|n| \le n_0$,

$$c_n(M^* e^{-t\mathcal{L}^*} g_0) = 0.$$

We can rewrite this as

$$M^* e^{-tB_n^*} c_n(g_0) = 0.$$

Differentiating ℓ times in t and evaluating at $t = 0$, we get that for all $\ell \in \mathbb{N}$ and $|n| \le n_0$,

$$M^* (B_n^*)^\ell c_n(g_0) = 0.$$

Since we assumed that for $|n| > n_0$, $c_n(g_0) = 0$, this means that $c_n(g_0) \in \ker([B_n|M]^*)$. But, by definition of E, $c_n(g_0) \in \mathrm{Range}([B_n|M]) = \ker([B_n|M]^*)^\perp$. Thus, $c_n(g_0) = 0$ and $g_0 = 0$. This proves the unique continuation property, and the claim.

Step 2: Construction of u_n. — We set $T' = T_* + \epsilon = T - \epsilon$.

If we denote the adjugate matrix of a matrix C by $\mathrm{Adj}(C)$, define

$$Q(\mathrm{i}n) := [B_n|M]_k^* \,\mathrm{Adj}([B_n|M]_k[B_n|M]_k^*);$$
$$P(\mathrm{i}n) := \det([B_n|M]_k[B_n|M]_k^*).$$

Notice that Q is a polynomial with matrix coefficients, and P is a nonzero polynomial (with scalar coefficients).

Increasing n_0 if necessary, we may assume that for every $|n| > n_0$, $P(\mathrm{i}n) \neq 0$. Then, for $|n| > n_0$,

$$[B_n|M]_k^+ = Q(\mathrm{i}n)/P(\mathrm{i}n).$$

We first apply a control as in step 1: for any $f_0 \in E$, there exists $u \in L^2(0,\epsilon;\, H_0^p(\omega))^m$ such that $(1-\Pi)S(\epsilon, f_0, Mu) = 0$. Then, the resulting solution $f(\epsilon,\cdot)$ is such that $P(\mathrm{i}n) = 0 \implies c_n(f(\epsilon,\cdot)) = 0$, since $P(\mathrm{i}n) \neq 0$ for $|n| > n_0$ and $c_n(f(\epsilon,\cdot)) = 0$ for $|n| \leqslant n_0$.

We consider this $f(\epsilon,\cdot)$ as our new initial condition, that we denote by f_ϵ, and we have to steer it to 0 in time T'. Note that since $f_0 \in H^p(\mathbb{T})$ and $u \in L^2(0,\epsilon;\, H_0^p(\omega))^d$, according to Duhamel's formula and the fact that the semigroup $\mathrm{e}^{-t\mathcal{L}}$ is strongly continuous on $H^p(\mathbb{T})^d$, the state f_ϵ also belongs to $H^p(\mathbb{T})^d$.

Let $\ell \in \mathbb{N}$ large enough. According to proposition 10, there exists

$$v \in H_0^{p+\deg P}((0,T')\times\omega)^{d_\mathrm{h}} \times H_0^{\ell}((0,T')\times\omega)^{d_\mathrm{p}}$$

such that $S(T', f_\epsilon, P(\partial_x)v) = 0$. Write $Q(\mathrm{i}n)$ by blocks as:

$$Q(\mathrm{i}n) = \begin{pmatrix} Q_1(\mathrm{i}n) \\ \vdots \\ Q_k(\mathrm{i}n) \end{pmatrix} = \begin{pmatrix} Q_1^\mathrm{h}(\mathrm{i}n) & Q_1^\mathrm{p}(\mathrm{i}n) \\ \vdots & \vdots \\ Q_k^\mathrm{h}(\mathrm{i}n) & Q_k^\mathrm{p}(\mathrm{i}n) \end{pmatrix},$$

where the $Q_j(\mathrm{i}n)$ are of size $m \times d$, the $Q_j^\mathrm{h}(\mathrm{i}n)$ are of size $m \times d_\mathrm{h}$ and $Q_j^\mathrm{p}(\mathrm{i}n)$ are of size $m \times d_\mathrm{p}$. Notice that the $L_{n,j}^\mathrm{h}$ defined in the statement of theorem 12 are $L_{n,j}^\mathrm{h} = Q_j^\mathrm{h}(\mathrm{i}n)/P(\mathrm{i}n)$. Set also

$$w_n(t) := Q(\mathrm{i}n)c_n(v(t,\cdot)).$$

Write it by blocks as

$$w_n(t) = \begin{pmatrix} w_{n,1}(t) \\ \vdots \\ w_{n,k}(t) \end{pmatrix} = \begin{pmatrix} Q_1(\mathrm{i}n)c_n(v(t,\cdot)) \\ \vdots \\ Q_k(\mathrm{i}n)c_n(v(t,\cdot)) \end{pmatrix}.$$

Finally, according to lemma 13, we also set

$$u_n(t) := w_{n,1}(t) + w_{n,2}'(t) + \cdots + w_{n,k}^{(k-1)}(t).$$

Step 3: Conclusion. — Remark that for every $n \in \mathbb{Z}$,

$$\begin{aligned}
[B_n|M]_k w_n(t) &= [B_n|M]_k Q(\mathrm{i}n)c_n(v(t,\cdot)) \\
&= [B_n|M]_k[B_n|M]_k^* \,\mathrm{Adj}([B_n|M]_k[B_n|M]_k^*)c_n(v(t,\cdot)) \\
&= \det([B_n|M]_k[B_n|M]_k^*)c_n(v(t,\cdot)) \\
&= P(\mathrm{i}n)c_n(v(t,\cdot)).
\end{aligned}$$

Moreover, since $S(T', f_\epsilon, P(\partial_x)v) = 0$, the control $\tilde{v}_n(t) := P(\mathrm{i}n)c_n(v(t,\cdot))$ steers $c_n(f_\epsilon)$ to 0 for the system $X_n' = B_nX_n + \tilde{v}_n$ in time T'. That is to say, w_n steers $c_n(f_\epsilon)$ to 0 for the system $X_n' = B_nX_n +$

$[B_n|M]_k w_n$ in time T'. Thus, according to lemma 13, u_n steers $c_n(f_\epsilon)$ to 0 for the system (10) in time T'.

Thus, the control u formally defined by $u := \sum_{n\in\mathbb{Z}} u_n e_n$ is such that $S(f_\epsilon, T', Mu) = 0$. We claim that the previous sum is well-defined in $L^2(0, T'; L^2(\mathbb{T}))$. Indeed, if we define u in the sense of distributions,

$$u = (Q_1(\partial_x) + \partial_t Q_2(\partial_x) + \cdots + \partial_t^{k-1} Q_k(\partial_x))v.$$

Since v is supported on $[0, T'] \times \omega$, so is u. Consider the differential operator $\mathcal{Q} := Q_1(\partial_x) + \partial_t Q_2(\partial_x) + \cdots + \partial_t^{k-1} Q_k(\partial_x)$. We have $u = \mathcal{Q}w$. Write this operator by blocks as $\mathcal{Q} = \begin{pmatrix} \mathcal{Q}^{\mathrm{h}} & \mathcal{Q}^{\mathrm{p}} \end{pmatrix}$. In other words,

$$\mathcal{Q}^{\mathrm{h}} := Q_1^{\mathrm{h}}(\partial_x) + \partial_t Q_2^{\mathrm{h}}(\partial_x) + \cdots + \partial_t^{k-1} Q_k^{\mathrm{h}}(\partial_x).$$

The order of the differential operator $\mathcal{Q}^{\mathrm{h}}$ is at most

$$\mathrm{Order}(\mathcal{Q}^{\mathrm{h}}) \leq \max_{1\leq j\leq k} (j - 1 + \deg(Q_j^{\mathrm{h}})).$$

Since $L_{n,j}^{\mathrm{h}} = Q_j^{\mathrm{h}}(\mathrm{i}n)/P(\mathrm{i}n)$, according to the definition of p (see theorem 12), $\mathrm{Order}(\mathcal{Q}^{\mathrm{h}}) \leq p + \deg(P)$. Moreover, recall that $v \in H_0^{p+\deg(P)}((0, T')\times\omega)^{d_{\mathrm{h}}} \times H_0^{\ell}((0, T')\times\omega)^{d_{\mathrm{p}}}$. Thus, if we choose $\ell \geq \mathrm{Order}(\mathcal{Q}^{\mathrm{p}})$, $u \in L^2((0, T') \times \omega)$. □

3.3 Upper bound on the loss of regularity

Theorem 12 requires initial condition to be H^p for some p. In this section, we provide an elementary upper bound on p.

Proposition 14. *Assume that for $|n|$ large enough, the Kalman rank condition* $\mathrm{rank}([B_n|M]) = d$ *holds. Let*

$$k(n) := \inf\{k : \ \mathrm{rank}([B_n|M]_k) = d\} \in \{\infty\} \cup \mathbb{N}.$$

Then, the sequence $(k(n))_{n\in\mathbb{Z}}$ is eventually constant when $|n| \to +\infty$. We will denote $k_0 := \lim_{|n|\to+\infty} k(n)$.

Proof. The rank condition $\mathrm{rank}([B_n|M]_k) = d$ is equivalent to $\det([B_n|M]_k[B_n|M]_k^*) \neq 0$. Let $P_k(n) = \det([B_n|M]_k[B_n|M]_k^*)$. P_k is a polynomial in n, hence if $P_k(n_0) \neq 0$ for some n_0, then $P_k(n) \neq 0$ for every large enough $|n|$. Thus, for every n_0, there exists n_1 such that $k(n) \leq k(n_0)$ whenever $|n| \geq n_1$. Since $k(n)$ is integer valued, it is eventually constant. □

Then, we have the following version of theorem 12.

Theorem 15. *Assume that the hypotheses (H.1)–(H.4) hold, that $T > T_*$ and that for all $|n| \in \mathbb{N}$ large enough, the Kalman rank condition* $\mathrm{rank}([B_n|M]) = d$ *holds. Let k_0 as in proposition 14. Let E as in theorem 12.*

Then, for every $f_0 \in H^{4d(k_0-1)}(\mathbb{T})^d \cap E$, there exists a control $u \in L^2([0, T] \times \omega)$ such that the solution f of the parabolic-transport system (Sys) with initial condition f_0 satisfies $f(T, \cdot) = 0$.

The sufficient part of theorem 1, as stated in the introduction is a special case of this theorem, since we always have $k_0 \leqslant d$. Here is the main lemma that allows us to bound the p of theorem 12 (see also [5, Theorem 2.1] for similar considerations).

Lemma 16. *Let $A \in \mathcal{M}_d(\mathbb{C})_p[X]$ a polynomial of degree at most p with $d \times d$ matrices coefficients. Assume that for some $z_0 \in \mathbb{C}$, $A(z_0)$ is invertible. Then, $A^{-1} \in \mathcal{M}_d(\mathbb{C})_{p(d-1)}(X)$, i.e., the coefficients of $(A(z))^{-1}$ are rational functions of z of degree at most $p(d-1)$.*

Proof. Write

$$A(z)^{-1} = \frac{1}{\det(A(z))} \operatorname{Adj}(A(z)),$$

where $\operatorname{Adj}(A(z))$ is the adjugate matrix of $A(z)$. $\det(A(z))$ and $\operatorname{Adj}(A(z))$ are nonzero polynomials in z. Moreover, the coefficients of $\operatorname{Adj}(A(z))$ are sums of products on $d-1$ coefficients of $A(z)$. Hence, they are polynomials of degree at most $(d-1)p$. □

The case we are interested in is:

Corollary 17. *With k_0 as in proposition 14, set, when it is defined*

$$[B_n|M]^+_{k_0} := [B_n|M]^*_{k_0}\big([B_n|M]_{k_0}[B_n|M]^*_{k_0}\big)^{-1}.$$

Then, as a function of n, $[B_n|M]^+_{k_0} \in \mathcal{M}_d(\mathbb{C})^{d\times d}_{2(k_0-1)(2d-1)}(X)$.

Proof. We have $[B_n|M]_{k_0} \in \mathcal{M}_{d,mk_0}(\mathbb{C})_{2(k_0-1)}[X]$, hence

$$[B_n|M]_{k_0}[B_n|M]^*_{k_0} \in \mathcal{M}_d(\mathbb{C})_{4(k_0-1)}(X).$$

According to the previous lemma,

$$\big([B_n|M]_{k_0}[B_n|M]^*_{k_0}\big)^{-1} \in \mathcal{M}_d(\mathbb{C})_{4(k_0-1)(d-1)}(X).$$

Hence $[B_n|M]^+_{k_0} \in \mathcal{M}_d(\mathbb{C})_k(X)$ with $k = 4(k_0-1)(d-1) + 2(k_0-1) = 2(k_0-1)(2d-1)$. □

Proof of theorem 15. According to theorem 12, every initial condition in $E \cap H^p(\mathbb{T})^d$ can be steered to 0, where $p = \deg([B_n|M]^+_{k_0}) + k_0 - 1$ (degree as a rational function of n). But according to corollary 17, $\deg([B_n|M]^+_{k_0}) \le 2(k_0-1)(2d-1)$. Thus $p \le 4d(k_0-1)$. Hence, every initial condition in $E \cap H^{4d(k_0-1)}(\mathbb{T})^d$ can be steered to 0. □

4 Necessary conditions for null-controllability

4.1 Construction of WKB solutions

We will give necessary conditions of null-controllability using so called *WKB solutions*, that we construct here. Using these kind of approximate solutions is standard for the wave equation (see, e.g., [25, pp. 426–428] or [32, Appendix B] for a more elementary presentation) or the Schrödinger equation (see, e.g., [36, pp. 16–17]). WKB solutions were also used to disprove observability of some 2×2 parabolic-transport system with variable coefficients [2, §3] (see also [3, §3] for a Navier-Stokes system with Maxwell's law). Our construction is a generalization of their construction for system of arbitrary size, which brings a few difficulties. For the sake of clarity, we construct WKB solutions only for systems with constant coefficients, which is enough for our purposes. But it is likely that such a construction could be adapted to a large class of variable-coefficients parabolic-transport systems of arbitrary sizes.

To disprove the observability inequality, these WKB solutions ought to be constructed for the adjoint system. But the parabolic-transport system (Sys) and its adjoint have the same structure, so, in order to lighten the notations, we construct the WKB solutions for the system (Sys).

Let $\phi \in C^\infty([0,T]\times\mathbb{T};\mathbb{C})$ such that $\Im(\phi)\geq 0$ and $\partial_x\phi$ never vanishes. We search approximate solutions $g_h^{\mathrm{WKB}}(t,x)$ of the parabolic-transport system (Sys) with the following ansatz, where $h>0$ is assumed to be small:

$$\begin{cases} g_h^{\mathrm{WKB}}(t,x) = X_h(t,x)\mathrm{e}^{\mathrm{i}\phi(t,x)/h}, \\ X_h(t,x) \sim \sum_{j\geq 0} h^j Y_j(t,x). \end{cases} \tag{11}$$

We have

$$\begin{aligned} \partial_x g_h^{\mathrm{WKB}} &= \left(\partial_x X_h + \frac{\mathrm{i}}{h}\partial_x\phi X_h\right)\mathrm{e}^{\mathrm{i}\phi/h}, \\ \partial_t g_h^{\mathrm{WKB}} &= \left(\partial_t X_h + \frac{\mathrm{i}}{h}\partial_t\phi X_h\right)\mathrm{e}^{\mathrm{i}\phi/h}, \\ \partial_x^2 g_h^{\mathrm{WKB}} &= \left(\partial_x^2 X_h + \frac{2\mathrm{i}}{h}\partial_x\phi\partial_x X_h - \frac{1}{h^2}(\partial_x\phi)^2 X_h + \frac{\mathrm{i}}{h}\partial_x^2\phi X_h\right)\mathrm{e}^{\mathrm{i}\phi/h}. \end{aligned}$$

Assuming that this g_h^{WKB} is solution of the parabolic-transport system (Sys), we get

$$\begin{aligned} 0 &= \left(\partial_t - B\partial_x^2 + A\partial_x + K\right)(X_h\mathrm{e}^{\mathrm{i}\phi/h}) \\ &= \left[\left(\partial_t - B\partial_x^2 + A\partial_x + K\right)X_h + \frac{1}{h}\left(\mathrm{i}\partial_t\phi + \mathrm{i}A\partial_x\phi - \mathrm{i}B\partial_x^2\phi - 2\mathrm{i}B\partial_x\phi\partial_x\right)X_h + \frac{1}{h^2}B(\partial_x\phi)^2 X_h\right]\mathrm{e}^{\mathrm{i}\phi/h}. \end{aligned}$$

Plugging in the asymptotic expansion of X_h, we get

$$0 \sim \sum_{j\geq -2}\left[(\partial_x\phi)^2 BY_{j+2} + \left(\mathrm{i}\partial_t\phi + \mathrm{i}A\partial_x\phi - \mathrm{i}B\partial_x^2\phi - 2\mathrm{i}B\partial_x\phi\partial_x\right)Y_{j+1} + \left(\partial_t - B\partial_x^2 + A\partial_x + K\right)Y_j\right]h^j,$$

where, by convention, $Y_j = 0$ for $j<0$. We want to cancel each of the terms in this sum. Thus, we are looking for $(Y_j)_{j\geq 0}$ such that for all $j\geq -2$,

$$(\partial_x\phi)^2 BY_{j+2} + \left(\mathrm{i}\partial_t\phi + \mathrm{i}A\partial_x\phi - \mathrm{i}B\partial_x^2\phi - 2\mathrm{i}B\partial_x\phi\partial_x\right)Y_{j+1} + \left(\partial_t - B\partial_x^2 + A\partial_x + K\right)Y_j = 0. \tag{12}$$

Solving this induction relation requires us to look at different projections of this equation. From now on, we will denote

$$Y_j = \begin{pmatrix} Y_j^{\mathrm{h}} \\ Y_j^{\mathrm{p}} \end{pmatrix} \quad \text{with } Y_j^{\mathrm{h}} \in \mathbb{C}^{d_{\mathrm{h}}} \text{ and } Y_j^{\mathrm{p}} \in \mathbb{C}^{d_{\mathrm{p}}}.$$

Then, recalling that $B = \left(\begin{smallmatrix} 0 & 0 \\ 0 & D \end{smallmatrix}\right)$ and taking the parabolic components of eq. (12) (*i.e.*, the d_{p} last components), we get

$$(\partial_x\phi)^2 DY_j^{\mathrm{p}} = -\begin{pmatrix} 0 & I \end{pmatrix}\left[\left(\mathrm{i}\partial_t\phi + \mathrm{i}A\partial_x\phi - \mathrm{i}B\partial_x^2\phi - 2\mathrm{i}B\partial_x\phi\partial_x\right)Y_{j-1} + \left(\partial_t - B\partial_x^2 + A\partial_x + K\right)Y_{j-2}\right]. \tag{13}$$

Since D is invertible, this formula determines Y_j^{p} as a function of Y_{j-1} and Y_{j-2}. Note that for $j=0$, this implies $Y_0^{\mathrm{p}} = 0$.

Before looking at the other projections of eq. (12), let us recall that $A = \left(\begin{smallmatrix} A' & A_{12} \\ A_{21} & A_{22} \end{smallmatrix}\right)$. We similarly write K in blocks as $\left(\begin{smallmatrix} K' & K_{12} \\ K_{21} & K_{22} \end{smallmatrix}\right)$. Then, taking the transport (*i.e.*, the first d_{h}) components of eq. (12), we get

$$0 = (\mathrm{i}\partial_t\phi + \mathrm{i}\partial_x\phi A')Y_j^{\mathrm{h}} + \mathrm{i}\partial_x\phi A_{12}Y_j^{\mathrm{p}} + \begin{pmatrix} I & 0 \end{pmatrix}(\partial_t + A\partial_x + K)Y_{j-1}. \tag{14}$$

From now on, we choose ϕ of the form[2]

$$\phi(t,x) = \psi(x-\mu t),\ \psi \in C^\infty(\mathbb{R}),\ \Im(\phi) \geq 0,\ \psi' \text{does not vanish}, \psi\ 2\pi\text{-periodic}, \tag{15}$$

where μ is an eigenvalue of A' an ψ' never vanishes. With this ϕ, eq. (14) reads

$$\begin{aligned} 0 &= \mathrm{i}\psi'(x-\mu t)(A'-\mu)Y_j^{\mathrm{h}} + \mathrm{i}\psi'(x-\mu t)A_{12}Y_j^{\mathrm{p}} + \begin{pmatrix} I & 0 \end{pmatrix}(\partial_t + A\partial_x + K)Y_{j-1} \\ &= \mathrm{i}\psi'(x-\mu t)(A'-\mu)Y_j^{\mathrm{h}} + \mathrm{i}\psi'(x-\mu t)A_{12}Y_j^{\mathrm{p}} + (\partial_t + A'\partial_x + K')Y_{j-1}^{\mathrm{h}} + (A_{12}\partial_x + K_{12})Y_{j-1}^{\mathrm{p}}. \end{aligned} \tag{16}$$

Denote by P'_μ the projection on the eigenspace of A' associated with μ along the other eigenspaces. We consider $Y_{j,\mu}^{\mathrm{h}} \in \mathrm{Range}(P'_\mu)$ defined by $Y_{j,\mu}^{\mathrm{h}} = P'_\mu Y_j^{\mathrm{h}}$. Similarly, we set $Y_{j,\neq\mu}^{\mathrm{h}} \in \ker(P'_\mu)$ as $Y_{j,\neq\mu}^{\mathrm{h}} = (I - P'_\mu)Y_j^{\mathrm{h}}$. Finally, we write in blocks A' and K' along the sum $\mathbb{R}^d = \mathrm{Range}(P'_\mu) \oplus \ker(P'_\mu)$ as

$$A' = \begin{pmatrix} \mu & 0 \\ 0 & A'_{22} \end{pmatrix}, \quad K' = \begin{pmatrix} K'_{11} & K'_{12} \\ K'_{21} & K'_{22} \end{pmatrix},$$

where $A'_{22} \in \mathcal{L}(\ker(P'_\mu))$, $K'_{11} = P'_\mu K' P'_\mu \in \mathcal{L}(\mathrm{Range}(P'_\mu))$, $K'_{12} \in \mathcal{L}(\ker(P'_\mu), \mathrm{Range}(P'_\mu))$, etc. Then, projecting eq. (16) on $\ker(P'_\mu)$ along $\mathrm{Range}(P'_\mu)$ (*i.e.*, multiplying by $(I - P'_\mu)$), we get

$$\mathrm{i}\psi'(x-\mu t)(A'_{22}-\mu)Y_{j,\neq\mu}^{\mathrm{h}} = -(I-P'_\mu)\left[\mathrm{i}\psi'(x-\mu t)A_{12}Y_j^{\mathrm{p}} + \begin{pmatrix} I & 0 \end{pmatrix}(\partial_t + A\partial_x + K)Y_{j-1}\right]. \tag{17}$$

Since P'_μ is the projection on the eigenspace of A' associated with the eigenvalue μ, $A' - \mu$ is invertible on $\ker(P'_\mu)$, *i.e.*, $A'_{22} - \mu$ is invertible. Hence, eq. (17) determines $Y_{j,\neq\mu}^{\mathrm{h}}$ as a function of Y_j^{p} and Y_{j-1}.

Finally, we project eq. (16) on $\mathrm{Range}(P'_\mu)$, and we get

$$0 = (\partial_t + \mu\partial_x + K'_{11})Y_{j,\mu}^{\mathrm{h}} + K'_{12}Y_{j,\neq\mu}^{\mathrm{h}} + P'_\mu(A_{12}\partial_x + K_{12})Y_j^{\mathrm{p}} + \mathrm{i}\psi'(x-\mu t)P'_\mu A_{12}Y_{j+1}^{\mathrm{p}}. \tag{18}$$

We then use eq. (13) to express Y_{j+1}^{p} as

$$Y_{j+1}^{\mathrm{p}} = D_1 Y_j^{\mathrm{h}} + D_2 Y_j^{\mathrm{p}} + D_3 Y_{j-1}, \ \text{ with } \ D_1 = -\frac{\mathrm{i}}{\psi'(x-\mu t)}D^{-1}A_{21},$$

and where D_2 and D_3 are matrix first or second-order differential operators. Their specific expressions do not matter for our purpose. Plugging this in eq. (18), we get

$$\begin{aligned} &(\partial_t + \mu\partial_x + K'_{11} + P'_\mu A_{12}D^{-1}A_{21}P'_\mu)Y_{j,\mu}^{\mathrm{h}} \\ &\quad = -K'_{12}Y_{j,\neq\mu}^{\mathrm{h}} - P'_\mu(A_{12}\partial_x + K_{12})Y_j^{\mathrm{p}} - \mathrm{i}\psi'(x-\mu t)P'_\mu A_{12}(D_1(I-P'_\mu)Y_{j,\neq\mu}^{\mathrm{h}} + D_2 Y_j^{\mathrm{p}} + D_3 Y_{j-1}). \end{aligned} \tag{19}$$

If we choose an initial condition $Y_{j,\mu,0}^{\mathrm{h}}$ for $Y_{j,\mu}^{\mathrm{h}}$, eq. (19) determines $Y_{j,\mu}^{\mathrm{h}}$ as a function of $Y_{j,\mu,0}^{\mathrm{h}}$, $Y_{j,\neq\mu}^{\mathrm{h}}$, Y_j^{p} and Y_{j-1}.

We have seen that if ϕ is given by eq. (15), the $(Y_j)_{j\in\mathbb{N}}$ that solve the WKB recurrence equation (12) are given by eqs. (13), (17) and (19).

To be rigorous, we have only proved that *if* $(Y_j)_{j\in\mathbb{N}}$ *solves eq.* (12), *then* Y_j^{p}, $Y_{j,\neq\mu}^{\mathrm{h}}$ *and* $Y_{j,\mu}^{\mathrm{h}}$ *solves eqs.* (13), (17) *and* (19) *respectively*, but not the reciprocal (which is what we are actually interested in). However, we easily rephrase the computations of this section as a sequence of equivalences:

- ($\forall j \geq -2$, Y_j solves eq. (12)), if and only if

[2]Equations (13) and (14) with $j = 0$ imply $(\partial_t\phi + \partial_x\phi A')Y_0^{\mathrm{h}} = 0$. If we want a non-trivial Y_0^{h}, this imposes ϕ to depend only on $x - \mu t$ for some $\mu \in \mathrm{Sp}(A')$.

- ($\forall j \geq 0$, Y_j^{p} solves eq. (13), $Y_{j,\neq\mu}^{\mathrm{h}}$ solves eq. (17) and $Y_{j,\mu}^{\mathrm{h}}$ solves eq. (18)), if and only if
- ($\forall j \geq 0$, Y_j^{p} solves eq. (13), $Y_{j,\neq\mu}^{\mathrm{h}}$ solves eq. (17) and $Y_{j,\mu}^{\mathrm{h}}$ solves eq. (19)).

We summarize the computations of this section in the following proposition:

Proposition 18. *Let $\psi \in C^\infty(\mathbb{T})$ such that ψ' never vanishes and $\Im(\psi) \geq 0$. Let $\mu \in \mathrm{Sp}(A')$ and set ϕ as in eq. (15).*

For every $j \geq 0$, let $Y_{j,\mu,0}^{\mathrm{h}} \in C^\infty(\mathbb{T}; \ker(A' - \mu))$. Define $(Y_j^{\mathrm{p}})_{j\geq -2}$, $(Y_{j,\neq\mu}^{\mathrm{h}})_{j\geq -2}$ and $(Y_{j,\mu}^{\mathrm{h}})_{j\geq -2}$ with the following recursive procedure:

- *set $Y_{-2}^{\mathrm{p}} = Y_{-1}^{\mathrm{p}} = 0$, $Y_{-2,\neq\mu}^{\mathrm{h}} = Y_{-1,\neq\mu}^{\mathrm{h}} = 0$, $Y_{-2,\mu}^{\mathrm{h}} = Y_{-1,\mu}^{\mathrm{h}} = 0$;*
- *if Y_k^{p}, $Y_{k,\neq\mu}^{\mathrm{h}}$, $Y_{k,\mu}^{\mathrm{h}}$ are defined for $-2 \leq k \leq j-1$, define Y_j^{p} with eq. (13), $Y_{j,\neq\mu}^{\mathrm{h}}$ with eq. (17) and $Y_{j,\mu}^{\mathrm{h}}$ with eq. (19) with initial condition $Y_{j,\mu,0}^{\mathrm{h}}$.*

For $j \geq 0$, set $Y_j(t,x) = \begin{pmatrix} Y_{j,\mu}^{\mathrm{h}} + Y_{j,\neq\mu}^{\mathrm{h}} \\ Y_j^{\mathrm{p}} \end{pmatrix}$. Let $q \in \mathbb{N}$. Let the function g_h^{WKB} be defined by

$$g_h^{\mathrm{WKB}}(t,x) = \sum_{j=0}^{q} h^j Y_j \mathrm{e}^{\mathrm{i}\phi(t,x)/h}. \tag{20}$$

Then, defining r_h by

$$(\partial_t - B\partial_x^2 + A\partial_x + K) g_h^{\mathrm{WKB}}(t,x) = r_h(t,x)\mathrm{e}^{\mathrm{i}\phi(t,x)/h},$$

for every $k \in \mathbb{N}$, $\ell \in \mathbb{N}$, $t \in [0,T]$ and $x \in \mathbb{T}$,

$$|\partial_t^k \partial_x^\ell r_h(t,x)| \leq C_{k,\ell} h^{q-1}.$$

Remark 19. Assume that $h^{-1} \in \mathbb{N}$. Then, replacing ϕ by $\phi + 2k\pi$ in eq. (11) does not change the WKB solution g_h^{WKB}. In fact, in the construction, we do not really need ϕ to be 2π-periodic in x, we only need

$$\phi(t, x + 2\pi) = \phi(t,x) \mod 2\pi.$$

The construction is entirely unchanged with this weaker hypothesis on ϕ. Thus, we can choose

$$\phi(t,x) = \mathrm{i}\varphi(x - \mu t) + (x - \mu t) \text{ with } \mu \in \mathrm{Sp}(A'),\ \varphi \geq 0.$$

These WKB solutions will be used to disprove observability inequalities that often feature a projections on high frequencies. To deal with these projection on high frequencies, we will use the following lemma.

Lemma 20. *Let $n \in \mathbb{Z}$. Under the assumptions of proposition 18, for every $\ell \in \mathbb{N}$, we have uniformly in $0 \leq t \leq T$, in the limit $h \to 0^+$,*

$$(g_h^{\mathrm{WKB}}(t,\cdot), e_n)_{L^2} = O(h^\ell).$$

Proof. The scalar product $(g_h^{\mathrm{WKB}}(t,\cdot), \mathrm{e}^{\mathrm{i}nx})_{L^2}$ can be written as

$$(g_h^{\mathrm{WKB}}(t,\cdot), e_n)_{L^2} = \int_{\mathbb{T}} w_{t,h,n}(x) \mathrm{e}^{\mathrm{i}\psi(x-\mu t)/h}\, \mathrm{d}x,$$

where

$$w_{t,h,n}(x) := \sum_{j=0}^{q} h^j Y_j(t,x) \mathrm{e}^{-\mathrm{i}nx}.$$

Note that $w_{t,h,n}$ and its derivative are uniformly bounded for $0 \leq t \leq T$ and $h \leq 1$. Consider the differential operator $L := (\mathrm{i}\psi'(x-\mu t))^{-1}\partial_x$. Here, we use the fact that ψ' never vanishes. This operator is such that

$$hL\mathrm{e}^{\mathrm{i}\psi(x-\mu t)/h} = \mathrm{e}^{\mathrm{i}\psi(x-\mu t)/h}.$$

Thus, denoting L^* the adjoint of L, by integration by parts,

$$(g_h^{\mathrm{WKB}}(t,\cdot), e_n)_{L^2} = h^\ell \int_{\mathbb{T}} (\overline{L^*})^\ell (w_{t,h,n})(x)\mathrm{e}^{\mathrm{i}\psi(x-\mu t)/h}\,\mathrm{d}x.$$

The operator L^* is a differential operator independent of h. Hence, by definition of $w_{t,h,n}$

$$(g_h^{\mathrm{WKB}}(t,\cdot), e_n)_{L^2} = O(h^\ell). \qquad \square$$

4.2 The parabolic-transport system is not null controllable in small time

We now prove that the time condition $T \geqslant T_*$ is necessary (remark that the equality case $T = T^*$ remains an open question). It was already proved to be necessary for the null-controllability of every L^2 initial conditions [7]. But this proof did not exclude the null-controllability of every H^k initial condition when $T < T_*$.

Proof of theorem 5. Let $\mu \in \mathrm{Sp}(A')$ with maximum modulus. By definition, $T_* = \ell(\omega)/|\mu|$. Let $T < T_*$.

We aim to disprove the observability inequality associated to the control problem of theorem 5 using the WKB solution constructed above. We claim that this observability inequality is: there exists $C > 0$ such that for every $g_0 \in L^2(\mathbb{T})^d$, the solution g of

$$(\partial_t - B^*\partial_x - A^*\partial_x + K^*\partial_x)g(t,x) = 0, \quad g(0,x) = g_0(x) \tag{21}$$

satisfies

$$\|\pi_N g(T,\cdot)\|_{H^{-k}(\mathbb{T})} \leq C\|M^* g\|_{L^2((0,T)\times\omega)}, \tag{22}$$

where $\pi_N : \sum_{n\in\mathbb{Z}} X_n \mathrm{e}^{\mathrm{i}nx} \in L^2(\mathbb{T}) \mapsto \sum_{|n|>N} X_n \mathrm{e}^{\mathrm{i}nx}$.

This is proved using a standard duality lemma, see *e.g.* [14, Lemma 2.48] with $C_2 = \mathrm{e}^{-t\mathcal{L}} \circ \pi_N^* \circ \iota_k$ and $C_1 : u \in L^2((0,T)\times\omega) \mapsto \int_0^T \mathrm{e}^{-(T-t)\mathcal{L}} Mu(t)\,\mathrm{d}t$, where ι_k is the injection $H^k(\mathbb{T}) \to L^2(\mathbb{T})$. Note that π_N^* is the injection $\{\sum_{|n|>N} X_n \mathrm{e}^{\mathrm{i}nx}\} \to L^2(\mathbb{T})$, and that ι_k^* is a bijective isometry $H^{-k}(\mathbb{T}) \to H^k(\mathbb{T})$ ([7, Lemma 33]).

Testing this observability inequality on initial conditions of the form $\partial_x^k g_0$ instead of g_0, we get

$$\|\pi_N g(T,\cdot)\|_{L^2(\mathbb{T})} \leq C\|\partial_x^k M^* g\|_{L^2((0,T)\times\omega)}. \tag{23}$$

Step 1: Construction of the counterexample. — According to remark 19, we assume implicitly that $h^{-1} \in \mathbb{N}$ in what follows. Let $T < T_*$. There exists $x_0 \notin \overline{\omega}$ such that $x_0 - \mu t \notin \overline{\omega}$ for every $0 \leq t \leq T$. Choose $\varphi \in C^\infty(\mathbb{T})$, real-valued, such that $\varphi(x_0) = 0$, $\varphi''(x_0) \neq 0$ and $\varphi(x) > 0$ for every $x \neq x_0$. Then, choose $\phi(t,x) = \mathrm{i}\varphi(x+\mu t) + (x+\mu t)n_0$, as we did in remark 19 (the change from μ to $-\mu$ is because we are considering $-A^*$ instead of A).

This choice of ϕ ensures that whatever the choices of the Y_j, the WKB solution g_h^{WKB} defined by eq. (11) stays concentrated around $x_0 + \mu t$.

Let $Y_{0,\mu,0}^{\mathrm{h}} \in C^\infty(\mathbb{T};\ \ker(A'^* + \mu))$ with $Y_{0,\mu,0}^{\mathrm{h}}(x_0) \neq 0$. For $j \geq 1$, set $Y_{j,\mu,0}^{\mathrm{h}} = 0$. Let $q > k+1$. Consider the function g_h^{WKB} defined by proposition 18 (where B, and K are replaced respectively by B^* and K^*, and where A is replaced by $-A^*$).

Set also $g_h(t,x)$ the solution of the adjoint system (21) with initial condition $g_h^{\mathrm{WKB}}(t=0,\cdot)$.

Step 2: Estimation of the difference between g_h^{WKB} and g_h. — According to proposition 18,

$$(\partial_t - B^*\partial_x^2 - A^*\partial_x + K^*)g_h^{\mathrm{WKB}} = O(h^{k+1})\mathrm{e}^{\mathrm{i}\phi(t,x)/h},$$

Hence, with $r_h := g_h^{\mathrm{WKB}} - g_h$, we have $r_h(0,x) = 0$ and

$$(\partial_t - B^*\partial_x^2 - A^*\partial_x + K^*)r_h = O(h^{k+1})\mathrm{e}^{\mathrm{i}\phi(t,x)/h}.$$

where the O has to be understood in the C^∞-topology. Since the parabolic-transport system is well-posed in $H^k(\mathbb{T})^d$, we get that for every $j \in \mathbb{N}$, uniformly in $0 < t < T$,

$$\|\partial_x^j(r_h(t,\cdot))\|_{L^2} \le C_j h^{k-j+1}. \tag{24}$$

Step 3: Upper bound on the right-hand side of the observability inequality. — According to the triangle inequality,

$$\|\partial_x^k M^* g_h\|_{L^2((0,T)\times\omega)} \le \|\partial_x^k M^* g_h^{\mathrm{WKB}}\|_{L^2((0,T)\times\omega)} + \|\partial_x^k M^* r_h\|_{L^2((0,T)\times\omega)}.$$

According to (24) applied to $j = k$, the second term of the right-hand side is $O(h)$. For the first term of the right-hand side, we recall that $g_h^{\mathrm{WKB}} = \sum_{j=0}^q h^j Y_j \mathrm{e}^{\mathrm{i}\psi(x+\mu t)}$, and that, thanks to our choice of ψ, $\mathrm{e}^{\mathrm{i}\psi(x+\mu t)}$ is exponentially small when $x + \mu t \neq x_0$ *i.e.* $x \neq x_0 - \mu t$. Therefore, since $x_0 - \mu t \notin \overline{\omega}$ for every $0 \le t \le T$, for some $c > 0$,

$$\|\partial_x^k M^* g_h^{\mathrm{WKB}}\|_{L^2((0,T)\times\omega)} = O(\mathrm{e}^{-c/h}).$$

This proves that

$$\|\partial_x^k M^* g_h\|_{L^2((0,T)\times\omega)} = O(h). \tag{25}$$

Step 4: Lower bound on the left-hand side of the observability inequality. — According to lemma 20 applied for $\ell = 1$, we have

$$\|\pi_N g_h^{\mathrm{WKB}}(T,\cdot)\|_{L^2(\mathbb{T})} = \|g_h^{\mathrm{WKB}}(T,\cdot)\|_{L^2(\mathbb{T})} + O(h). \tag{26}$$

Moreover, using the inverse triangle inequality,

$$\|\pi_N g_h(T,\cdot)\|_{L^2(\mathbb{T})} \ge \|\pi_N g_h^{\mathrm{WKB}}(T,\cdot)\|_{L^2(\mathbb{T})} - \|\pi_N r_h(T,\cdot)\|_{L^2(\mathbb{T})}.$$

Using the error estimate (24) for $j = 0$ and eq. (26), we get

$$\|\pi_N g_h(T,\cdot)\|_{L^2(\mathbb{T})} \ge \|g_h^{\mathrm{WKB}}(T,\cdot)\|_{L^2(\mathbb{T})} - Ch^k. \tag{27}$$

Thus, we only need to find a lower-bound for $\|g_h^{\mathrm{WKB}}(T,\cdot)\|_{L^2(\mathbb{T})}$. We have

$$\|g_h^{\mathrm{WKB}}(T,\cdot)\|_{L^2(\mathbb{T})}^2 = \int_{\mathbb{T}} \left|\sum_{j=0}^q h^j Y_j(t,x)\right|^2 \mathrm{e}^{-2\varphi(x+\mu T)/h}\,\mathrm{d}x = \int_{\mathbb{T}} |Y_0(t,x)|^2 \mathrm{e}^{-2\varphi(x+\mu T)/h}\,\mathrm{d}x + O(h).$$

Recall that $\varphi(x_0) = 0$, that for $x \neq x_0$, $\varphi > 0$ and that $\varphi''(x_0) \neq 0$. Then, using Laplace's method (see *e.g.* [37, §2.2] and in particular [37, eq. (2.34)]), we get

$$\|g_h^{\mathrm{WKB}}(T,\cdot)\|_{L^2(\mathbb{T})}^2 = c\sqrt{h} + O(h)$$

for some $c > 0$. Plugging this into eq. (27), we get that for h small enough,

$$\|\pi_N g_h(T,\cdot)\|_{L^2(\mathbb{T})} \ge c\sqrt{h}. \tag{28}$$

Step 5: Conclusion. — Comparing the lower bound (28) and the upper bound (25) and taking h small enough, we see that the observability inequality (22) cannot hold if $T < T_*$, hence the parabolic-transport system (Sys) with initial conditions in $H^k \cap \pi_N(L^2(\mathbb{T}))$ is not null-controllable in time $T < T_*$. □

4.3 Rough initial conditions are not null-controllable

We now give necessary conditions for every L^2 initial condition to be steerable to 0. To do this, we only need the first term of the WKB expansion of proposition 18. By analyzing higher-order terms of the WKB expansion, it is likely that we could get necessary conditions for the null-controllability of every H^k initial conditions. But doing this analysis in general seems hard, and we leave this for future work, or on a case-by-case basis. We will prove the following statement, which is a refined version of theorem 3.

Proposition 21. *Let $\mu \in \mathrm{Sp}(A')$, $N \in \mathbb{N}$ and $T > 0$. Let P'_μ be the projection on the eigenspace of A' associated to μ. Write K in blocks as $\left(\begin{smallmatrix} K' & K_{12} \\ K_{21} & K_{22} \end{smallmatrix}\right)$, with $K' \in \mathcal{M}_{d_h}(\mathbb{R})$. Set*

$$K^*_\mu := (P'_\mu)^* \left((K')^* + A^*_{21}(D^*)^{-1}A^*_{12}\right)(P'_\mu)^*$$

Assume that every initial condition $f_0 \in L^2(\mathbb{T})^d \cap \{\sum_{|n|>N} X_n \mathrm{e}^{\mathrm{i}nx}\}$ is steerable to 0 in time T with control in $L^2((0,T)\times\omega)$. Then, for every $\mu \in \mathrm{Sp}(A')$ and for every non-zero subspace $S \subset \mathrm{Range}((P'_\mu)^*)$ that is stable by K^*_μ, there exists $V_0 \in S$ such that $M^*\left(\begin{smallmatrix} V_0 \\ 0 \end{smallmatrix}\right) \neq 0$.

Proof. Step 1: Observability inequality. — Using a standard duality lemma [14, Lemma 2.48], and as in the proof of theorem 5, we get an observability inequality that is equivalent to the null-controllability of the system (Sys) with initial conditions in $L^2(\mathbb{T})^d \cap \{\sum_{|n|>N} X_n \mathrm{e}^{\mathrm{i}nx}\}$. This observability inequality is: there exists $C > 0$ such that for every $g_0 \in L^2(\mathbb{T})^d$, the solution g of

$$(\partial_t - B^*\partial_x^2 - A^*\partial_x + K^*)g(t,x) = 0, \quad g(0,x) = g_0(x) \tag{29}$$

satisfies

$$\|\pi_N g(T,\cdot)\|_{L^2(\mathbb{T})} \leq C\|M^* g\|_{L^2((0,T)\times\omega)}, \tag{30}$$

where, as in the proof of theorem 5, $\pi_N : \sum_{n\in\mathbb{Z}} X_n \mathrm{e}^{\mathrm{i}nx} \in L^2(\mathbb{T}) \mapsto \sum_{|n|>N} X_n \mathrm{e}^{\mathrm{i}nx}$.

Step 2: Construction of the counterexample. — According to remark 19, we assume implicitly that $h^{-1} \in \mathbb{N}$ in what follows. Let $V_0 \in S \setminus \{0\}$. Set $\varphi := 0$ and let $\phi(t,x) = x - \mu t$ as in remark 19. Set $Y^{\mathrm{h}}_{0,\mu,0}(x) := V_0$. For $j > 0$, set $Y^{\mathrm{h}}_{j,\mu,0} := 0$. Let g_h^{WKB} be defined by proposition 18 with B and K replaced respectively by B^* and K^* and A by $-A^*$, and with $q \geq 2$. Let g_h be the solution of the parabolic-transport system (Sys) with initial condition $g_h^{\mathrm{WKB}}(0,\cdot)$.

Remark that according to proposition 18, and in particular eq. (19),

$$(\partial_t - \mu\partial_x + K^*_\mu)Y^{\mathrm{h}}_{0,\mu} = 0.$$

Thus, $Y^{\mathrm{h}}_{0,\mu}(t,x) = \mathrm{e}^{-tK^*_\mu}V_0$. In particular, since S is stable by K^*_μ, $Y^{\mathrm{h}}_{0,\mu}(t,x) \in S$ for all t, x.

Step 3: Error estimate between g_h^{WKB} and g_h. — Set $r_h := g_h - g_h^{\mathrm{WKB}}$. Then $r_h(0,x) = 0$, and according to proposition 18,

$$(\partial_t - B^*\partial_x^2 - A^*\partial_x + K^*)r_h = O(h).$$

Since the parabolic-transport system is well-posed in $L^2(\mathbb{T})^d$, we get notably that uniformly in $0 < t < T$,

$$\|r_h(t,\cdot)\|_{L^2} \leq Ch. \tag{31}$$

Step 4: Upper bound of the right-hand side of the observability inequality. — Using the error estimate (31), the right-hand side of the observability inequality (30) satisfies

$$\begin{aligned}\|M^* g_h\|^2_{L^2((0,T)\times\omega)} &\le \|M^* g_h^{\mathrm{WKB}}\|^2_{L^2((0,T)\times\omega)} + Ch \\ &\le \|M^* Y_0^{\mathrm{h}} \mathrm{e}^{\mathrm{i}\phi/h}\|^2_{L^2((0,T)\times\omega)} + Ch \\ &= \left\| M^* \begin{pmatrix} Y_{0,\mu}^{\mathrm{h}} \\ 0 \end{pmatrix} \right\|^2_{L^2((0,T)\times\omega)} + Ch \\ &= 2\pi \int_0^T \left| M^* \begin{pmatrix} \mathrm{e}^{-tK_\mu^*} V_0 \\ 0 \end{pmatrix} \right|^2 \mathrm{d}t + Ch, \end{aligned} \tag{32}$$

where we used the definition of g_h^{WKB} for the last three inequalities.

Step 5: Lower-bound of the left-hand side of the observability inequality. — Using the error estimate (31), the left-hand side of the observability inequality (30) satisfies

$$\|\pi_N g_h(T,\cdot)\|^2_{L^2} \ge \|\pi_N g_h^{\mathrm{WKB}}(T,\cdot)\|^2_{L^2} - Ch.$$

Then, using the estimate on low frequencies of g_h^{WKB} (lemma 20 with $\ell = 1$)

$$\|\pi_N g_h(T,\cdot)\|^2_{L^2} \ge \|g_h^{\mathrm{WKB}}(T,\cdot)\|^2_{L^2} - Ch.$$

Now, using the definition of g_h^{WKB}, and the fact that $|\mathrm{e}^{\mathrm{i}\phi}| = 1$,

$$\begin{aligned}\|\pi_N g_h(T,\cdot)\|^2_{L^2} &\ge \|Y_{0,\mu}^{\mathrm{h}}(T,\cdot)\|^2_{L^2} - Ch. \\ &= 2\pi |\mathrm{e}^{-TK_\mu^*} V_0|^2 - Ch. \end{aligned} \tag{33}$$

Step 6: Conclusion. — Comparing the upper bound on the right-hand side of the observability inequality (eq. (32)) and the lower bound on the left-hand side (eq. (33)), we see that $M^* \mathrm{e}^{-tK_\mu^*} V_0$ cannot vanish for every $0 \le t \le T$. Since $\mathrm{e}^{-tK_\mu^*} V_0 \in S$ for every t, this proves the proposition. □

5 Systems of two equations

We apply the general theorems of the previous sections on 2×2 systems. Some of these results are not new (see, *e.g.*, [13]). Our goal here is only to check whether our results are optimal, at least in this setting.

5.1 Control properties of 2×2 systems: statements

Here, we consider the parabolic transport-system (Sys) with

$$B = \begin{pmatrix} 0 & 0 \\ 0 & d \end{pmatrix}, \quad A = \begin{pmatrix} a' & a_{12} \\ a_{21} & a_{22} \end{pmatrix}, \quad K = \begin{pmatrix} k_{11} & k_{12} \\ k_{21} & k_{22} \end{pmatrix}, \quad M = \begin{pmatrix} m_1 \\ m_2 \end{pmatrix}, \tag{34}$$

where all lower-case letters are real numbers, with $d > 0$ and $a' \neq 0$. Here, we assume that M has rank one. We do not need to treat the case where $\mathrm{rank}(M) = 2$, because it is already covered with the general theorem where there is a control on every component (see [7, Theorem 2] or theorem 15 with $k = 1$): every initial condition in $L^2(\mathbb{T})^d$ is null-controllable in time $T > T_*$. In the following three propositions, we detail the applications of our general theorem to eleven cases, showcasing the variety of phenomena that can appear depending on the values of every coefficients. The proofs are given in the next subsections.

Proposition 22. *Assume that B, A, K, M are given by eq. (34). Assume that $(m_1, m_2) = (1, 0)$.*

- *If $(a_{21}, k_{21}) = (0, 0)$, the parabolic-transport system (Sys) is not null-controllable, whatever the time T is.*

Let $T > \ell(\omega)/|a'|$ (where $\ell(\omega)$ is defined in eq. (1)).

- *If $k_{21} \neq 0$, every initial condition in $L^2(\mathbb{T})^2$ for the system (Sys) can be steered to 0 in time T with L^2 controls.*
- *If $a_{21} \neq 0$ and $k_{21} = 0$, every initial condition $f_0 = (f_0^{\mathrm{h}}, f_0^{\mathrm{p}})$ in $L^2(\mathbb{T})^2$ such that $\int_{\mathbb{T}} f_0^{\mathrm{p}} = 0$ for the system (Sys) can be steered to 0 in time T with L^2 controls.*

Proposition 23. *Assume that B, A, K, M are given by eq. (34). Assume that $(m_1, m_2) = (0, 1)$.*

- *If $(a_{12}, k_{12}) = (0, 0)$, the parabolic-transport system (Sys) is not null-controllable, whatever the time T is.*

Let $T > \ell(\omega)/|a'|$.

- *If $a_{12} \neq 0$ and $k_{12} \neq 0$, every initial condition in $H^1(\mathbb{T}) \times L^2(\mathbb{T})$ for the system (Sys) can be steered to 0 in time T with L^2 controls.*
- *If $a_{12} \neq 0$ and $k_{12} = 0$, every initial condition $f_0 = (f_0^{\mathrm{h}}, f_0^{\mathrm{p}})$ in $H^1(\mathbb{T}) \times L^2(\mathbb{T})$ such that $\int_{\mathbb{T}} f_0^{\mathrm{h}} = 0$ for the system (Sys) can be steered to 0 in time T with L^2 controls.*
- *If $a_{12} = 0$ and $k_{12} \neq 0$, every initial condition in $H^2(\mathbb{T}) \times L^2(\mathbb{T})$ for the system (Sys) can be steered to 0 in time T with L^2 controls.*

In every cases, there exists an initial condition f_0 in $L^2(\mathbb{T})$ such that $\int_{\mathbb{T}} f_0 = 0$ that cannot be steered to 0 in time T with L^2 controls.

In the case where $a_{21} = 0$ and $k_{21} \neq 0$, there is a gap in the regularity condition that is sufficient for the null controllability (*i.e.*, $H^2 \times L^2$), and the lack of null-controllability of $L^2 \times L^2$ initial conditions. Are every $H^1 \times L^2$ initial conditions steerable to 0? We conjecture that this is not the case, but theorem 3 is not enough to prove so. We would need to look at the second term in the WKB expansion to find out, or use another method; maybe using a refined version of regularization properties of lemma 25.

We do not detail in general the case where $m_1 \neq 0$ and $m_2 \neq 0$. Let us just mention that there is no regularity condition for null-controllability to hold. But depending on whether the solution of $\det([B_n, M]) = 0$ (which is a quadratic equation in n) are integer, there might be a condition on at most two fourier components for an initial condition to be steerable to 0. We only detail the following case that is about the simultaneous control of a transport and a parabolic equation.

Proposition 24. *Assume that B, M are given by eq. (34). Assume that $A = \left(\begin{smallmatrix} a' & 0 \\ 0 & a_{22} \end{smallmatrix}\right)$ and $K = \left(\begin{smallmatrix} k_{11} & 0 \\ 0 & k_{22} \end{smallmatrix}\right)$. Assume that $(m_1, m_2) = (1, 1)$. Let $T > \ell(\omega)/|a'|$.*

- *If $a' \neq a_{22}$ and $k_{11} = k_{22}$, every initial condition $f_0 = (f_0^{\mathrm{h}}, f_0^{\mathrm{p}}) \in L^2(\mathbb{T})^2$ such that $\int_{\mathbb{T}} f_0^{\mathrm{h}} = \int_{\mathbb{T}} f_0^{\mathrm{p}}$ can be steered to zero with controls in L^2.*
- *If $a' \neq a_{22}$ and $k_{11} \neq k_{22}$, every initial condition in $L^2(\mathbb{T})^2$ can be steered to zero with controls in L^2.*
- *If $a' = a_{22}$ and $\sqrt{(k_{22} - k_{11})/d} \notin \mathbb{N}$, every initial condition in $L^2(\mathbb{T})^2$ can be steered to zero with controls in L^2.*

- *If $a' = a_{22}$ and $n_0 := \sqrt{(k_{22} - k_{11})/d} \in \mathbb{N}$, every initial condition $f_0 = (f_0^{\mathrm{h}}, f_0^{\mathrm{p}}) \in L^2(\mathbb{T})^2$ such that $c_{\pm n_0}(f_0^{\mathrm{h}}) = c_{\pm n_0}(f_0^{\mathrm{p}})$ can be steered to zero with controls in L^2.*

The case $a' \neq a_{22}$ and $k_{11} = k_{22}$ is not new, at least in spirit: the simultaneous controllability (equivalently, additive observability) of a heat equation and a wave equation has been studied by Zuazua [42, §2.1–2.2].

5.2 Regularity of the free equation

We will use some basic regularity results.

Lemma 25. *Let $f_0 \in H^1(\mathbb{T})^{d_{\mathrm{h}}} \times L^2(\mathbb{T})^{d_{\mathrm{p}}}$. For every $t > 0$, $\mathrm{e}^{-t\mathcal{L}} f_0 \in H^1(\mathbb{T})^d$.*

Assume in addition that $A_{12} = 0$, and that $f_0 \in H^2(\mathbb{T})^{d_{\mathrm{h}}} \times L^2(\mathbb{T})^{d_{\mathrm{p}}}$. For every $t > 0$, $\mathrm{e}^{-t\mathcal{L}} f_0 \in H^2(\mathbb{T})^d$.

To prove it, we will use the following (sub)lemma:

Lemma 26. *Consider $\mathcal{L}^{\mathrm{p}}$ and F^{p} as defined in section 2 (or [7, §4.1]). For every $t > 0$, $k \in \mathbb{N}$ and $f_0 \in F^{\mathrm{p}}$, $\mathrm{e}^{-t\mathcal{L}^{\mathrm{p}}} f_0 \in H^k(\mathbb{T})^d$.*

Proof. Set $f(t) = \mathrm{e}^{-t\mathcal{L}^{\mathrm{p}}} f_0$. Denote the first d_{h} components of $f(t)$ by $f^{\mathrm{h}}(t)$ and the last d_{p} components of $f(t)$ by $f^{\mathrm{p}}(t)$ (and similarly for f_0).

We will use some simple tools from [7, §4.4.1]. For the sake of readability, we redo the proof in full here.

Step 1: Computing $f^{\mathrm{h}}(t)$ as a function of $f^{\mathrm{p}}(t)$. — Since $f(t) \in F^{\mathrm{p}}$, by definition of F^{p} (section 2), for every $|n| > n_0$,

$$P^{\mathrm{p}}(\mathrm{i}/n) c_n(f(t)) = c_n(f(t)).$$

Writing $P^{\mathrm{p}}(z)$ by blocks as $\left(\begin{smallmatrix} p_{11}(z) & p_{12}(z) \\ p_{21}(z) & p_{22}(z) \end{smallmatrix}\right)$, and taking the first d_{h} components,

$$p_{11}(\mathrm{i}/n) c_n(f^{\mathrm{h}}(t)) + p_{12}(\mathrm{i}/n) c_n(f^{\mathrm{p}}(t)) = c_n(f^{\mathrm{h}}(t)).$$

Since $P^{\mathrm{p}}(0) = \left(\begin{smallmatrix} 0 & 0 \\ 0 & I \end{smallmatrix}\right)$, $p_{11}(0) = 0$ and for z small enough, $|p_{11}(z)| < 1$. Then, increasing n_0 if necessary, for $|n| > n_0$,

$$c_n(f^{\mathrm{h}}(t)) = (I - p_{11}(\mathrm{i}/n))^{-1} p_{12}(\mathrm{i}/n) c_n(f^{\mathrm{p}}(t)).$$

For $z \in \mathbb{C}$ small enough, let $G(z) = (I - p_{11}(z))^{-1} p_{12}(z)$. Then, G depends holomorphically in z small enough, and for $|n| > n_0$ $c_n(f^{\mathrm{h}}(t)) = G(\mathrm{i}/n) c_n(f^{\mathrm{p}}(t))$.

Step 2: Conclusion. — Define $\mathcal{D}$ the unbounded operator on $L^2(\mathbb{T})^{d_{\mathrm{p}}}$ with domain $H^2(\mathbb{T})^{d_{\mathrm{p}}}$ by

$$\mathcal{D}\Big(\sum_n X_n \mathrm{e}^{\mathrm{i}nx}\Big) := \sum_n \big(n^2 D - \mathrm{i}nA_{22} - K_{22} - G(\mathrm{i}/n)(\mathrm{i}nA_{21} + K_{21})\big) X_n \mathrm{e}^{\mathrm{i}nx}.$$

Recall that

$$(\partial_t - D\partial_x^2 + A_{22}\partial_x + K_{22}) f^{\mathrm{p}}(t) + (A_{21}\partial_x + K_{21}) f^{\mathrm{h}}(t) = 0.$$

Since $c_n(f^{\mathrm{h}}(t)) = G(\mathrm{i}/n) c_n(f^{\mathrm{p}}(t))$, this can be written as $(\partial_t + \mathcal{D}) f^{\mathrm{p}}(t) = 0$. Hence,

$$f^{\mathrm{p}}(t) = \mathrm{e}^{-t\mathcal{D}} f_0^{\mathrm{p}} = \sum_{|n|>n_0} \mathrm{e}^{-t\left(n^2 D + \mathrm{i}nA_{22} + K_{22} + G(\mathrm{i}/n)(\mathrm{i}nA_{21} + K_{21})\right)} c_n(f_0^{\mathrm{p}}).$$

Since $\Re(\mathrm{Sp}(D)) \subset (0, +\infty)$, $f^{\mathrm{p}}(t)$ is in every $H^k(\mathbb{T})^{d_{\mathrm{p}}}$. Since the first d_{h} components of $f(t)$ are

$$f^{\mathrm{h}}(t) = \sum_{|n|>n_0} G(\mathrm{i}/n) c_n(f^{\mathrm{p}}(t)) e_n,$$

and since $G(\mathrm{i}/n)$ is bounded as $|n| \to +\infty$, $f^{\mathrm{h}}(t)$ also belongs in every $H^k(\mathbb{T})^{d_{\mathrm{h}}}$. □

Proof of lemma 25. The proof consists in looking at the projection on hyperbolic (respectively parabolic) components of $\mathrm{e}^{-t\mathcal{L}} f_0$, using the asymptotics for the hyperbolic projection. As in the previous proof, we denote the first d_h components of f_0 by f_0^h and the last d_p components by f_0^p.

Let us also recall that according to [7, §4.1],

$$\mathrm{e}^{-t\mathcal{L}} f_0 = \mathrm{e}^{-t\mathcal{L}^0}\Pi^0 f_0 + \mathrm{e}^{-t\mathcal{L}^\mathrm{h}}\Pi^\mathrm{h} f_0 + \mathrm{e}^{-t\mathcal{L}^\mathrm{p}}\Pi^0 f_\mathrm{p}. \tag{35}$$

Step 1: Asymptotics for the hyperbolic projection. — We use the notations $P^\mathrm{p}(z)$, $P^\mathrm{h}(z)$ defined in [7, Proposition 5–6]. Using the series for the perturbation of the total eigenprojections [28, Ch. II, eq. (2.14)], we get

$$\begin{aligned} P^\mathrm{h}(z) &= \begin{pmatrix} I & 0 \\ 0 & 0 \end{pmatrix} - z\left(\begin{pmatrix} I & 0 \\ 0 & 0 \end{pmatrix} A \begin{pmatrix} 0 & 0 \\ 0 & D^{-1} \end{pmatrix} + \begin{pmatrix} 0 & 0 \\ 0 & D^{-1} \end{pmatrix} A \begin{pmatrix} I & 0 \\ 0 & 0 \end{pmatrix}\right) + O(z^2) \\ &= \begin{pmatrix} I & 0 \\ 0 & 0 \end{pmatrix} - z\begin{pmatrix} 0 & A_{12}D^{-1} \\ D^{-1}A_{21} & 0 \end{pmatrix} + O(z^2). \end{aligned}$$

Thus,

$$\Pi^\mathrm{h} f_0 = \sum_{|n|>n_0} \left[\begin{pmatrix} c_n(f_0^\mathrm{h}) \\ 0 \end{pmatrix} - \frac{\mathrm{i}}{n}\begin{pmatrix} A_{12}D^{-1}c_n(f_0^\mathrm{p}) \\ D^{-1}A_{21}c_n(f_0^\mathrm{h}) \end{pmatrix} + O(n^{-2}c_n(f_0))\right]\mathrm{e}^{\mathrm{i}nx}. \tag{36}$$

Step 2: Case where $f_0 \in H^1 \times L^2$. — Since $\Pi^0 f_0$ is a finite sum of $\mathrm{e}^{\mathrm{i}nx}$, it is in every H^k, and so is $\mathrm{e}^{-t\mathcal{L}^0}\Pi^0 f_0$. According to the regularity of the parabolic frequencies (lemma 26), $\mathrm{e}^{-t\mathcal{L}^\mathrm{p}}\Pi^\mathrm{p} f_0$ is in every H^k.

Since $f_0^\mathrm{h} \in H^1(\mathbb{T})^{d_\mathrm{h}}$, $(c_n(f_0^\mathrm{h}))_n \in \ell^2(\mathbb{Z};\, 1+n^2)$ (the ℓ^2 space with weight $1+n^2$). Since $f_0^\mathrm{p} \in L^2(\mathbb{T})^{d_\mathrm{p}}$, $(c_n(f_0^\mathrm{p}))_n \in \ell^2(\mathbb{Z})$. Hence,

$$\left(c_n(f_0^\mathrm{h}) - \frac{\mathrm{i}}{n}A_{12}D^{-1}c_n(f_0^\mathrm{p})\right)_{|n|>n_0} \in \ell^2(|n|>n_0;\, 1+n^2),$$

and

$$\left(D^{-1}A_{21}c_n(f_0^\mathrm{h})\right)_{|n|>n_0} \in \ell^2(|n|>n_0;\, 1+n^2).$$

Hence, according to the asymptotics for Π^h of eq. (36), $\Pi^\mathrm{h} f_0 \in H^1(\mathbb{T})^d$. Since $\mathrm{e}^{-t\mathcal{L}^\mathrm{h}}$ is continuous on every H^k, $\mathrm{e}^{-t\mathcal{L}^\mathrm{h}}\Pi^\mathrm{h} f_0 \in H^1$.

Step 3: Case where $f_0 \in H^2 \times L^2$ and $A_{12} = 0$. — The asymptotics (36) reads

$$\Pi^\mathrm{h} f_0 = \sum_{|n|>n_0} \left[\begin{pmatrix} c_n(f_0^\mathrm{h}) \\ 0 \end{pmatrix} - \frac{\mathrm{i}}{n}\begin{pmatrix} 0 \\ D^{-1}A_{21}c_n(f_0^\mathrm{h}) \end{pmatrix} + O(n^{-2}c_n(f_0))\right]\mathrm{e}^{\mathrm{i}nx}. \tag{37}$$

The rest of the proof is very similar to the previous case: $\mathrm{e}^{-t\mathcal{L}^0}\Pi^0 f_0$ and $\mathrm{e}^{-t\mathcal{L}^\mathrm{p}}\Pi^\mathrm{p} f_0$ are in every H^k, while the asymptotics (37) proves that $\Pi^h f_0$ "gains" two derivatives compared to f_0^p. □

5.3 Control properties of 2×2 systems: proofs

Proof of proposition 22. In this case,

$$[B_n|M] = \begin{pmatrix} 1 & \mathrm{i}na' + k_{22} \\ 0 & \mathrm{i}na_{21} + k_{21} \end{pmatrix}.$$

In particular, $\det([B_n|M]) = \mathrm{i}na_{21} + k_{21}$. We see that if $(a_{21}, k_{21}) = (0,0)$, the Kalman rank condition never holds, whatever n is. Hence, according to remark 2, item 1, null-controllability does not hold, whatever T is.

Note that in our case, $[B_n|M]^+ = [B_n|M]^{-1}$ (when the right-hand side exists). Hence,

$$[B_n|M]^{-1} = \frac{1}{\mathrm{i}na_{21} + k_{21}} \begin{pmatrix} \mathrm{i}na_{21} + k_{21} & -\mathrm{i}na' - k_{22} \\ 0 & 1 \end{pmatrix}.$$

In particular, with the notations of theorem 12 with $k = 2$, $L^{\mathrm{h}}_{n,1} = 1$ and $L^{\mathrm{h}}_{n,2} = 0$. Thus, $p = 0$.

If $k_{21} \neq 0$, $\det([B_n|M]) = \mathrm{i}na_{21} + k_{21}$ never vanishes. In this case, E (as defined in theorem 12) is $E = L^2(\mathbb{T})^2$. Hence, according to theorem 12, every $L^2(\mathbb{T})^2$ can be steered to 0 with L^2 controls in time $T > \ell(\omega)/|a'|$.

If $a_{21} \neq 0$ and $k_{21} = 0$, the Kalman rank condition holds for every $n \neq 0$. For $n = 0$, according to the formula for $[B_n|M]$, $\mathrm{rank}([B_0|M]) = \mathbb{C} \times \{0\}$. Thus, $E = \{(f_0^{\mathrm{h}}, f_0^{\mathrm{p}}) \in L^2(\mathbb{T})^2,\ \int_{\mathbb{T}} f_0^{\mathrm{p}} = 0\}$. Therefore, according to theorem 12, every initial condition $(f_0^{\mathrm{h}}, f_0^{\mathrm{p}}) \in L^2(\mathbb{T})^2$ such that $\int_{\mathbb{T}} f_0^{\mathrm{p}} = 0$ can be steered to 0 with controls in L^2 in time $T > \ell(\omega)/|a'|$. □

Proof of proposition 23. In this case,

$$[B_n|M] = \begin{pmatrix} 0 & \mathrm{i}na_{12} + k_{12} \\ 1 & -n^2 d + \mathrm{i}na_{22} + k_{22} \end{pmatrix}.$$

In particular, $\det([B_n|M]) = -\mathrm{i}na_{12} - k_{12}$. We see that if $(a_{12}, k_{12}) = (0, 0)$, the Kalman rank condition never holds, whatever n is. Hence, according to remark 2 item 1, null-controllability does not hold, whatever T is.

As in the previous proof, $[B_n|M]^+ = [B_n|M]^{-1}$. Hence,

$$[B_n|M]^{-1} = \frac{1}{-\mathrm{i}na_{12} - k_{12}} \begin{pmatrix} -n^2 d + \mathrm{i}na_{22} + k_{22} & -\mathrm{i}na_{12} - k_{12} \\ -1 & 0 \end{pmatrix}.$$

In particular, with the notations of theorem 12 with $k = 2$, $L^{\mathrm{h}}_{n,1} = -(-n^2 d + \mathrm{i}na_{22} + k_{22})/(\mathrm{i}na_{12} + k_{12})$ and $L^{\mathrm{h}}_{n,2} = 1/(\mathrm{i}na_{12} + k_{12})$. In particular, if $a_{12} \neq 0$, $p = \max(1, 1-1) = 1$. And if $a_{12} = 0$ and $k_{12} \neq 0$, $p = \max(2, 1+0) = 2$.

Step 1: Case $a_{12} \neq 0$ and $k_{12} \neq 0$. — The Kalman rank condition holds for every n. Hence, with the notations of theorem 12, $p = 1$ and $E = L^2(\mathbb{T})^2$, and every initial condition in $H^1(\mathbb{T})^2$ can be steered to 0 with controls in L^2 in time $T > \ell(\omega)/|a'|$.

The strategy to control initial conditions in $H^1 \times L^2$ is first to let the solution evolve freely during an arbitrarily small time, which gives a $H^1(\mathbb{T})^2$ state (lemma 25), that we can steer to 0 according to the previous discussion.

Step 2: Case $a_{12} \neq 0$ and $k_{12} = 0$. — The case is almost the same as the previous one, except that the Kalman rank condition is not satisfied for $n = 0$ (and only for $n = 0$). We have $\mathrm{rank}([B_0|M]) = \{0\} \times \mathbb{C}$ and $E = \{(f_0^{\mathrm{h}}, f_0^{\mathrm{p}}) \in L^2(\mathbb{T})^2,\ \int_{\mathbb{T}} f_0^{\mathrm{h}} = 0\}$. We still have $p = 1$. Hence, we can steer every initial condition $(f_0^{\mathrm{h}}, f_0^{\mathrm{p}}) \in H^1(\mathbb{T})^2$ such that $\int_{\mathbb{T}} f_0^{\mathrm{h}} = 0$ an be steered to 0 with controls in time $T > \ell(\omega)/|a'|$.

As in the previous case, to control initial conditions in $H^1 \times L^2$, we let the solution evolve freely, which gives a $H^1(\mathbb{T})^2$ state, and preserves the property $\int_{\mathbb{T}} f_0^{\mathrm{h}} = 0$. Then, we can steer this state in time $T > \ell(\omega)/|a'|$.

Step 3: Case $a_{12} = 0$ and $k_{12} \neq 0$. — In this case, the Kalman rank condition is satisfied for every n, and $p = 2$. Hence, according to theorem 12, we can steer every $H^2(\mathbb{T})^2$ initial condition to 0 in time $T > \ell(\omega)/|a'|$ with controls in L^2.

Again, to control an initial condition in $H^2 \times L^2$, we let the solution evolve freely for a small time, which gives a $H^2(\mathbb{T})^2$ state (lemma 25), that we can steer to 0 in time $T > \ell(\omega)/|a'|$.

Step 4: Lack of null-controllability of L^2 initial conditions. — We have $M^*\left(\begin{smallmatrix}1\\0\end{smallmatrix}\right) = 0$. Hence, according to theorem 3, (recall that A' has size 1×1), there exists a $L^2(\mathbb{T})^2$ initial condition with zero average that cannot be steered to 0. □

Proof of proposition 24. We have

$$[B_n|M] = \begin{pmatrix} 1 & \mathrm{i}na' + k_{11} \\ 1 & -dn^2 + \mathrm{i}na_{22} + k_{22} \end{pmatrix}.$$

In particular, $\det([B_n|M]) = -dn^2 + \mathrm{i}n(a_{22} - a') + k_{22} - k_{11}$. We see that for n large enough, this determinant is non zero. In fact, taking the real and imaginary parts,

$$\det([B_n|M]) = 0 \Leftrightarrow \begin{cases} -dn^2 + k_{22} - k_{11} & = 0, \\ n(a_{22} - a') & = 0. \end{cases} \tag{38}$$

Moreover,

$$[B_n|M]^+ = [B_n|M]^{-1} = \frac{1}{\det([B_n|M])} \begin{pmatrix} -dn^2 + \mathrm{i}na_{22} + k_{22} & -\mathrm{i}na' - k_{11} \\ -1 & 1 \end{pmatrix}.$$

Thus,

$$L^{\mathrm{h}}_{n,1} = \frac{-dn^2 + O(n)}{-dn^2 + O(n)}, \quad \text{and} \quad L^{\mathrm{h}}_{n,2} = \frac{-1}{-dn^2 + O(n)}.$$

Thus, $p = \max(0, 1 - 2) = 0$.

Step 1: Case $a' \neq a_{22}$ and $k_{11} = k_{22}$. — According to eq. (38), the Kalman condition is satisfied for $n \neq 0$. Moreover, for $n = 0$, $\mathrm{Range}([B_0|M]) = \mathbb{C}M$, thus $E = \{(f_0^{\mathrm{h}}, f_0^{\mathrm{p}}) \in L^2(\mathbb{T})^2,\ \int_{\mathbb{T}} f_0^{\mathrm{h}} = \int_{\mathbb{T}} f_0^{\mathrm{p}}\}$. The theorem 12 gives the claimed controllability result.

Step 2: Case $a' \neq a_{22}$ and $k_{11} \neq k_{22}$. — According to eq. (38), the Kalman condition is satisfied for every $n \in \mathbb{Z}$. The theorem 12 gives the claimed controllability result.

Step 3: Case $a' = a_{22}$ and $\sqrt{(k_{22} - k_{11})/d} \notin \mathbb{N}$. — As in the previous case, according to eq. (38), the Kalman condition is satisfied for every $n \in \mathbb{Z}$. The theorem 12 gives the claimed controllability result.

Step 4: Case $a' = a_{22}$ and $n_0 := \sqrt{(k_{22} - k_{11})/d} \in \mathbb{N}$. — According to eq. (38), the Kalman condition is satisfied for $n \neq \pm n_0$. For $n = \pm n_0$, $\mathrm{Range}([B_{\pm n_0}|M]) = \mathbb{C}M$. The theorem 12 gives the claimed controllability result. □

A A finite dimension-uniqueness principle for the null-controllability

In the null controllability of parabolic-transport systems, we sometimes prove null-controllability "up to a finite dimensional space", and then use functional analysis arguments to deal with the finite-dimensional spaces that are left [31, 7]. In the previous articles, this was not stated as a general result. This is the purpose of this appendix.

Proposition 27. *Let $T_0 > 0$. Let H be a complex Hilbert space. Let A be an unbounded operator on H that generates a strongly continuous semigroup on H. Let U be another Hilbert space and let $B : U \to H$ a bounded control operator. For every $T > 0$, let U_T be a Hilbert space that is a subspace of $L^2(0, T; U)$ with continuous and dense injection that satisfies the following "extension by 0 property":*[3] *if $u \in U_T$,*

[3] In the application we use here, $U = L^2(\omega)$ and $U_T = H_0^k((0, T) \times \omega)$. The hypotheses of proposition 27 are tailored to allow this situation.

$a, b > 0$, then the function $\tilde{u}$ defined by $\tilde{u}(t) = 0$ for $0 < t < a$, $\tilde{u}(t) = u(t-a)$ for $a < t < T + a$, and $\tilde{u}(t) = 0$ for $T + a < t < T + a + b$ is in U_{T+a+b}.

Assume that there exists a finite dimensional space $\mathcal{F}$ of H that is stable by the semigroup e^{tA} and a closed finite codimensional space[4] $\mathcal{G}$ of H such that:

- *(control up to finite dimension) for every $f_0 \in \mathcal{G}$, there exists $u \in U_{T_0}$ such that the solution f of $f' = Af + Bu$ satisfies $f(T_0) \in \mathcal{F}$,*
- *(unique continuation) for every $\epsilon > 0$ and for every finite linear combination of generalized eigenfunctions $g_0 \in H$ of A^*, we have $B^*(\mathrm{e}^{tA^*}g_0) = 0$ on $t \in (0, \epsilon) \implies g = 0$.*

Then, for every $T > T_0$ and every $f_0 \in H$, there exists $u \in U_T$ such that the solution f of $f' = Af + Bu$, $f(0) = f_0$ satisfies $f(T) = 0$.

Remark 28.
- In this proposition, we can weaken the hypothesis "B bounded" into "B admissible" (see [14, §2.3]), but in this article, B is always bounded.
- If the assertion "($g_0 \in H$ is a finite linear combination of generalized eigenfunctions of A^* and $B^*g_0 = 0$) $\implies g_0 = 0$" holds, the unique continuation hypothesis is satisfied by well-posedness.

Proof. Step 1: We may assume that $\mathcal{F} \subset \mathcal{G}$. — We prove that if we replace $\mathcal{G}$ by $\mathcal{F} + \mathcal{G}$, the hypotheses are still satisfied. Let $f_0 \in \mathcal{F} + \mathcal{G}$. We write $f_0 = f_{\mathcal{F}} + f_{\mathcal{G}}$. According to the hypotheses, there exists $u \in U_{T_0}$ such that the solution f of $f' = Af + Bu$, $f(0) = f_{\mathcal{G}}$ is such that $f(T_0) \in \mathcal{F}$. Then, the solution $\widetilde{f}$ of $\widetilde{f}' = A\widetilde{f} + Bu$, $\widetilde{f}(0) = f_0$ is such that

$$\widetilde{f}(T_0) = \underbrace{\mathrm{e}^{T_0 A} f_{\mathcal{F}}}_{\in \mathcal{F}} + \underbrace{f(T_0)}_{\in \mathcal{F}}.$$

Note that if we replace T_0 by any $T_1 > T_0$, the hypotheses are still satisfied.

Step 2: For $T > T_0$, the control $u \in U_T$ such that $f(T) \in \mathcal{F}$ may be chosen linearly and continuously in $f_0 \in \mathcal{G}$. — This is a standard proof of control theory. For $f_0 \in \mathcal{G}$, set

$$V(f_0) := \{u \in U_T : \ f(T) \in \mathcal{F}, \ f \text{ solves } f' = Af + Bu, \ f(0) = f_0\}.$$

Since A generates a strongly continuous semigroup, $V(f_0)$ is a closed affine subspace of U_T. Then, we can define $\mathcal{U}(f_0)$ as the orthogonal projection of 0 onto $V(f_0)$ for the U_T-norm. Using the characterization of orthogonal projection on closed convex set, we see that $\mathcal{U}$ is linear. Using the fact that A generates a strongly continuous semigroup, the characterization of the projection on closed convex subsets and the closed graph theorem, we see that $\mathcal{U}$ is bounded.

For the rest of the proof we set $\mathcal{U}_T : \ \mathcal{G} \to U_T$ such a map. We also set

$$\mathcal{N}_T := \{f_0 \in H : \ \exists u \in U_T, \ f(T) = 0, \ f \text{ solves } f' = Af + Bu, \ f(0) = f_0\}. \tag{39}$$

Step 3: For $T \geq T_0$, $\mathcal{N}_T$ is a closed finite codimensional subspace of H. — Set $S_0(t)$ the semigroup e^{tA} restricted to $\mathcal{F}$. Since $\mathcal{F}$ is finite dimensional, $S_0(t)$ can be written as e^{tA_0}, where A_0 is a bounded operator of $\mathcal{F}$. Moreover, $A_0 = A_{|\mathcal{F}}$. In particular, S_0 is actually a *group* of bounded operators.

For $f_0 \in \mathcal{G}$, and $f' = Af + B\mathcal{U}_T f_0$, $f(0) = f_0$, we have $f(T) \in \mathcal{F}$, which allows us to define

$$\mathcal{K} : \ f_0 \in \mathcal{G} \mapsto -S_0(-T)f(T) \in \mathcal{F}$$

[4] We do not require $\mathcal{G}$ to be stable by e^{tA}.

The range of this operator $\mathcal{K}$ satisfies $\mathrm{Range}(\mathcal{K}) \subset \mathcal{F}$. Hence, $\mathcal{K}$ has finite rank and is compact. Thus, according to Fredholm's alternative, $(I+\mathcal{K})\mathcal{G}$ is a closed subspace of $\mathcal{G}$ of finite codimension.

Moreover, for every $f_0 \in \mathcal{G}$, the solution $\widetilde{f}$ of $\widetilde{f}' = A\widetilde{f} + B\mathcal{U}_T f_0$, $\widetilde{f}(0) = f_0 + \mathcal{K} f_0$ satisfies

$$\widetilde{f}(T) = f(T) + \mathrm{e}^{TA}\mathcal{K} f_0 = f(T) - S_0(T)S_0(-T)f(T) = 0.$$

Thus, $(I+\mathcal{K})\mathcal{G} \subset \mathcal{F}_T$. According to [9, Proposition 11.5], this proves that $\mathcal{N}_T$ is closed and has finite codimension in H.

Step 4: There exists $\delta > 0$ such that for every $T, T' \in (T_0, T_0+\delta)$, $\mathcal{N}_T = \mathcal{N}_{T'}$. — Assume $T_0 < T < T'$. If $u \in \mathcal{N}_T$, and if we extend u by 0 on (T, T'), we have have $u \in \mathcal{N}_{T'}$. Thus $\mathrm{codim}(\mathcal{N}_{T'}) \leq \mathrm{codim}(\mathcal{N}_T)$. Since $\mathrm{codim}(\mathcal{N}_T)$ is an integer, the discontinuities of $T \mapsto \mathrm{codim}(\mathcal{N}_T)$ are isolated, which proves the claim.

From now on, we choose $\epsilon \in (0, \delta/2)$ arbitrarily small and we set $T_1 = T_0 + \epsilon$.

Step 5: For $t \in (0, \epsilon)$, $(\mathrm{e}^{tA^*}\mathcal{N}_{T_1}^\perp)^\perp \subset \mathcal{N}_{T_1}$. — Let $0 < t < \epsilon$ and $f_0 \in (\mathrm{e}^{tA^*}\mathcal{N}_{T_1}^\perp)^\perp$. For every $g_0 \in \mathcal{N}_{T_1}^\perp$, we have

$$0 = \langle \mathrm{e}^{tA^*} g_0, f_0 \rangle = \langle g_0, \mathrm{e}^{tA} f_0 \rangle.$$

Thus, $\mathrm{e}^{tA} f_0 \in (\mathcal{N}_{T_1}^\perp)^\perp$. Since $\mathcal{N}_{T_1}$ is closed (step 3), $\mathrm{e}^{tA} f_0 \in \mathcal{N}_{T_1}$. By definition of $\mathcal{N}_{T_1}$ and the "extension by 0" property of U_{T_1}, this proves that $f_0 \in \mathcal{N}_{T_1+t}$. According to the previous step, $\mathcal{N}_{T_1+t} = \mathcal{N}_{T_1}$, which proves the claim.

Step 6: $\mathcal{N}_{T_1}^\perp$ is left-invariant by e^{tA^*}. — First, consider $0 < t < \epsilon$. According to the previous step, $\mathcal{N}_{T_1}^\perp \subset \left((\mathrm{e}^{tA^*}\mathcal{N}_{T_1}^\perp)^\perp\right)^\perp$. Since $\mathcal{N}_{T_1}^\perp$ is finite dimensional hence closed, $\mathcal{N}_{T_1}^\perp \subset \mathrm{e}^{tA^*}\mathcal{N}_{T_1}^\perp$. Moreover, $\dim(\mathrm{e}^{tA^*}\mathcal{N}_{T_1}^\perp) \leq \dim(\mathcal{N}_{T_1}^\perp)$. Thus, for $0 < t < \epsilon$, $\mathrm{e}^{tA^*}\mathcal{N}_{T_1}^\perp = \mathcal{N}_{T_1}^\perp$. Thanks to the semigroup property, this is true for all $t > 0$.

Step 7: Unique continuation property associated to the control problem "steer every $f_0 \in H$ into $\mathcal{N}_{T_1}$ in time ϵ with a control in U_ϵ". — The control problem is, in mathematical form, the following:

$$\forall f_0 \in H,\ \exists u \in U_\epsilon,\ f(T) \in \mathcal{N}_{T_1},\ \text{where } f' = Af + Bu,\ f(0) = f_0. \tag{40}$$

Let $\Pi : H \to H$ the orthogonal projection on $\mathcal{N}_{T_1}^\perp$. Set also $R_T : L^2(0, T; U) \to H$ the input-to-output map defined by

$$R_T u := f(T),\ \text{where } f' = Af + Bu,\ f(0) = 0.$$

Then, the control problem (40) is equivalent to

$$\forall f_0 \in H,\ \exists u \in U_\epsilon,\ \Pi \mathrm{e}^{\epsilon A} f_0 + \Pi R_\epsilon u = 0.$$

We denote by ι_ϵ the injection map $U_\epsilon \to L^2(0, T; U)$. Then, the previous assertion is equivalent to

$$\mathrm{Range}\left(\Pi \circ \mathrm{e}^{\epsilon A}\right) \subset \mathrm{Range}\left(\Pi \circ R_\epsilon \circ \iota_\epsilon\right).$$

The observability inequality associated to this control problem is (see [14, Lemma 2.48]):

$$\forall g_0 \in H,\ \|\mathrm{e}^{\epsilon A^*} \circ \Pi^* g_0\| \leq C \|\iota_\epsilon^* \circ R_\epsilon^* \circ \Pi^* g_0\|.$$

Since $\mathrm{Range}(\Pi^*) = \mathcal{N}_{T_1}^\perp$ is finite-dimensional, and since $\ker(\iota^*) = \mathrm{Range}(\iota)^\perp = \{0\}$, this is equivalent to

$$\forall g_0 \in \mathcal{N}_{T_1}^\perp,\ R_\epsilon^* g_0 = 0 \implies \mathrm{e}^{\epsilon A^*} g_0 = 0. \tag{41}$$

To conclude, since $\mathcal{N}_{T_1}^{\perp}$ is finite dimensional and stable by e^{tA^*}, the semigroup e^{tA^*} is in fact a *group*, and in particular $e^{\epsilon A^*}$ is invertible on $\mathcal{N}_{T_1}^{\perp}$. Moreover, $R_\epsilon^* g_0(t) = B^* e^{(\epsilon - t)A^*} g_0$ (see [14, Lemma 2.47]). Thus, the assertion (41) is equivalent to

$$\forall g_0 \in \mathcal{N}_{T_1}^{\perp}, \left(B^* e^{tA^*} g_0 = 0 \text{ for } 0 < t < \epsilon\right) \implies g_0 = 0. \tag{42}$$

Step 8: Conclusion. — The unique continuation property (42) of the previous step is exactly the unique continuation property we assumed. Thus, according to the previous step, we can steer every $f_0 \in H$ into $\mathcal{N}_{T_1}$ in time ϵ with a control in U_ϵ. According to the definition of $\mathcal{N}_{T_1}$, we can steer every $f_0 \in \mathcal{N}_{T_1}$ to 0 in time $T_1 = T + \epsilon$ with a control in U_{T_1}. Hence, we can steer every $f_0 \in H$ to 0 in time $T_1 + \epsilon = T + 2\epsilon$ with a control in $U_{T+2\epsilon}$. Since ϵ can be chosen arbitrarily small, this proves the proposition. □

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