# On the optimality of the observability inequalities for parabolic and hyperbolic systems with potentials* 

Thomas Duyckaerts ${ }^{\dagger}$ Xu Zhang ${ }^{\ddagger}$ and Enrique Zuazua ${ }^{\S}$


#### Abstract

In this paper we prove the optimality of the observability inequality for parabolic systems with potentials in even space dimensions $n \geq 2$. This inequality (derived by E. Fernández-Cara and the third author in the context of the scalar heat equation with potentials in any space dimension) asserts, roughly, that for small time, the total energy of solutions can be estimated from above in terms of the energy localized in a subdomain with an observability constant of the order of $\exp \left(C\|a\|_{\infty}^{2 / 3}\right)$, $a$ being the potential involved in the system. The problem of the optimality of the observability inequality remains open for scalar equations.

The optimality is a consequence of a construction due to V. Z. Meshkov of a complexvalued bounded potential $q=q(x)$ in $\mathbb{R}^{2}$ and a non-trivial solution $u$ of $\Delta u=q(x) u$ with the decay property $|u(x)| \leq \exp \left(-|x|^{4 / 3}\right)$. Meshkov's construction may be generalized to any even dimension. We give an extension to it to odd dimensions, which gives a sharp decay rate up to some logarithmic factor and yields a weaker optimality result in odd space-dimensions.

We address the same problem for the wave equation. In this case it is well known that, in space-dimension $n=1$, observability holds with a sharp constant of the order of $\exp \left(C\|a\|_{\infty}^{1 / 2}\right)$. For systems in even space dimensions $n \geq 2$ we prove that the best constant one can expect is of the order of $\exp \left(C\|a\|_{\infty}^{2 / 3}\right)$ for any $T>0$ and any observation domain. Based on Carleman inequalities, we show that the positive counterpart is also true when $T$ is large enough and the observation is made in a neighborhood of the boundary. As in the context of the heat equation, the optimality of this estimate is open for scalar equations.

We address similar questions, for both equations, with potentials involving the first order term. We also discuss issues related with the impact of the growth rates of the nonlinearities on the controllability of semilinear equations. Some other open problems are raised.


2000 Mathematics Subject Classification. Primary: 93B07; Secondary: 93B05, 35B37.

[^0]Key Words. Optimality. Meshkov's construction. Observability inequality. Heat equation. Wave equation. Potential. Carleman inequality. Decay at infinity.

## Contents

1 Introduction and main results ..... 3
2 Sharp observability estimates for hyperbolic and parabolic equations with po- tentials ..... 8
2.1 Statement of the results ..... 8
2.1.1 The parabolic system ..... 8
2.1.2 The hyperbolic system ..... 9
2.2 Global Carleman estimate for hyperbolic operators ..... 11
2.3 Proof of sharp observability estimates for hyperbolic equations with potentials ..... 13
3 Super-exponentially decaying solutions to elliptic equations ..... 19
3.1 Construction in even dimension ..... 21
3.1.1 Construction for large $r$ ..... 22
3.1.2 Decay of $u$ at infinity ..... 30
3.1.3 Extension of $u$ to all $\mathbb{R}^{2}$ ..... 31
3.2 Construction in odd dimension ..... 31
4 Optimality of the observability constant for the heat equation with zero order potential ..... 37
5 Optimality of the observability constant for the wave equation with zero order potential ..... 41
6 Equations with convective terms ..... 44
7 Connections with the controllability of semilinear equations ..... 45
7.1 The semilinear heat equation ..... 45
7.2 The semilinear wave equation ..... 46
8 Open problems ..... 47
8.1 Meshkov's construction for scalar equations ..... 48
8.2 Meshkov's construction in dimension $n=3$ ..... 48
8.3 Sharp observability and time-dependent potentials ..... 48
$8.4 L^{r}$-potentials ..... 49
$8.51-d$ problems ..... 49
8.6 Other equations ..... 49

## 1 Introduction and main results

In this paper we discuss the optimality of some observability estimates for the heat and wave equations with potentials that may be obtained by means of Carleman inequalities and that arise naturally in the context of control theory. To better illustrate the problem under consideration, let us first analyze the case of the heat equation.

Let $n \geq 1$ and $N \geq 1$ be two integers, $T>0, \Omega$ be a simply connected, bounded domain of $\mathbb{R}^{n}$ with smooth boundary $\Gamma$. Put $Q=(0, T) \times \Omega$ and $\Sigma=(0, T) \times \Gamma$. Consider the heat equation with a potential $a=a(t, x)$ in $L^{\infty}\left(Q ; \mathbb{R}^{N \times N}\right)$ :

$$
\begin{cases}\varphi_{t}-\Delta \varphi+a \varphi=0, & \text { in } Q,  \tag{1.1}\\ \varphi=0, & \text { on } \Sigma, \\ \varphi(0, x)=\varphi^{0}(x), & \text { in } \Omega,\end{cases}
$$

where $\varphi$ takes values in $\mathbb{R}^{N}$. Denote by $\|\cdot\|_{\infty}$ and $|\cdot|$ the (usual) norms on $L^{\infty}\left(Q ; \mathbb{R}^{N \times N}\right)$ and $\mathbb{R}^{N}$, respectively. We recall the following known result:

Theorem A ([11]) Assume that $\omega$ is an open non-empty subset of $\Omega$. Then, there exists a constant $C=C(\Omega, \omega)>0$, depending on $\Omega$ and $\omega$ but independent of $T$, the potential $a=a(t, x)$ and the solution $\varphi$ of (1.1), such that

$$
\begin{equation*}
\|\varphi(T)\|_{\left(L^{2}(\Omega)\right)^{N}}^{2} \leq \exp \left(C\left(1+\frac{1}{T}+T\|a\|_{\infty}+\|a\|_{\infty}^{2 / 3}\right)\right) \int_{0}^{T} \int_{\omega}|\varphi|^{2} d x d t \tag{1.2}
\end{equation*}
$$

for every solution $\varphi$ of (1.1), potential $a \in L^{\infty}\left(Q ; \mathbb{R}^{N \times N}\right)$ and time $T>0$.
Inequality (1.2), that we shall refer to as observability inequality for system (1.1), plays a key role for solving control problems for linear and nonlinear heat processes (see [11] and [12]). This inequality was proved in [11] using the Carleman inequality approach developed by A. Fursikov and O. Imanuvilov (see, for example, [13]) by paying special attention to the dependence of the observability constants on the size of the potential $a=a(t, x)$ entering in the equation. Indeed, the inequality was only proven in the case $N=1$, but the method of proof yields the same result in the vectorial case $N \geq 2$.

The observability constant in this situation is, by definition, the best constant in the above observability inequality. We shall denote it by $C^{*}(T, a)$. According to (1.2) we can guarantee that the following upper bound holds:

$$
\begin{equation*}
C^{*}(T, a) \leq \exp \left(C\left(1+\frac{1}{T}+T\|a\|_{\infty}+\|a\|_{\infty}^{2 / 3}\right)\right) . \tag{1.3}
\end{equation*}
$$

In this paper we show that this estimate is sharp, in a sense that will be made precise. Similar questions will be addressed for the wave equation as well.

Before going further it is convenient to observe that the estimate in (1.3) includes three different terms. More precisely:

$$
\begin{equation*}
\exp \left(C\left(1+\frac{1}{T}+T\|a\|_{\infty}+\|a\|_{\infty}^{2 / 3}\right)\right)=C_{1}^{*}(T, a) C_{2}^{*}(T, a) C_{3}^{*}(T, a) \tag{1.4}
\end{equation*}
$$

where

$$
\begin{equation*}
C_{1}^{*}(T, a)=\exp \left(C\left(1+\frac{1}{T}\right)\right), \quad C_{2}^{*}(T, a)=\exp \left(C T\|a\|_{\infty}\right), \quad C_{3}^{*}(T, a)=\exp \left(C\|a\|_{\infty}^{2 / 3}\right) . \tag{1.5}
\end{equation*}
$$

The role that each of these constants plays in the observability inequality is of different nature:

- When $a \equiv 0$, i.e. in the absence of potential, the observability constant is simply $C_{1}^{*}(T, a)$. This constant blows-up exponentially as $T \downarrow 0$. This growth rate is easily seen to be optimal by inspection of the heat kernel and has been analyzed in more detail in [11] and [20] to see what the influence of the geometry of $\Omega$ and $\omega$ is. We refer also to [23] for a discussion of the optimal growth rate in one space dimension. For large $a$, the constant $C_{1}^{*}$ is the leading one in the region $T \lesssim\|a\|_{\infty}^{-2 / 3}{ }^{1}$.
- The second constant $C_{2}^{*}(T, a)$ is very natural as it arises when applying Gronwall's inequality to analyze the time evolution of the $L^{2}$-norm of solutions. For large $a$, this is the leading constant in the region $T \gtrsim\|a\|_{\infty}^{-1 / 3}$.
- The constant $C_{3}^{*}(T, a)$, which, actually, only depends on the potential $a$, is the most intriguing one. Indeed, the $2 / 3$ exponent does not seem to arise naturally in the context of the heat equation since, taking into account that the heat operator is of order one and two in the time and space variables respectively, one could rather expect terms of the form $\exp \left(c\|a\|_{\infty}\right)$ and $\exp \left(c\|a\|_{\infty}^{1 / 2}\right)$, as a simple ODE argument would indicate.

This paper is devoted to discuss the optimality of this third contribution, $C_{3}^{*}(T, a)$, in the observability constant. Note that this contribution is only predominant in the region:

$$
\begin{equation*}
\|a\|_{\infty}^{-2 / 3} \lesssim T \lesssim\|a\|_{\infty}^{-1 / 3} . \tag{1.6}
\end{equation*}
$$

Nevertheless, its study is of great importance in view of applications to non-linear problem, as it is the largest constant depending on $a$ for small $T$ and large $a$. In fact, it is the key constant in the main result in [12] that asserts that some weakly blowing-up nonlinear heat processes may be controlled to zero under suitable growth conditions on the nonlinearity. We shall return to this matter bellow.

In this paper we show that the estimate (1.3) is sharp in what concerns the dependence on the potential $a$, for times of order less than $\|a\|_{\infty}^{-1 / 3}$, in space dimension $n \geq 2$ for systems with at least two equations. More precisely, in even dimension, the following holds:

Theorem 1.1. Assume that $n \geq 2$ is even and that $N \geq 2$. Let $\omega$ be any given non-empty open subset of $\Omega$ such that $\Omega \backslash \bar{\omega} \neq \emptyset$. Then there exist two constants $c>0$ and $\mu>0$, a family of potentials $\left\{a_{R}\right\}_{R>0} \subset L^{\infty}\left(Q ; \mathbb{R}^{N \times N}\right)$ satisfying

$$
\left\|a_{R}\right\|_{\infty} \xrightarrow[R \rightarrow+\infty]{\longrightarrow}+\infty,
$$

and a family of initial data $\left\{\varphi_{R}^{0}\right\}_{R>0}$ in $\left(L^{2}(\Omega)\right)^{N}$ such that the corresponding solution $\varphi_{R}$ of (1.1) satisfies

$$
\begin{equation*}
\lim _{R \rightarrow \infty}\left\{\inf _{T \in I_{\mu}} \frac{\left\|\varphi_{R}(T)\right\|_{\left(L^{2}(\Omega)\right)^{N}}^{2}}{\exp \left(c\left\|a_{R}\right\|_{\infty}^{2 / 3}\right) \int_{0}^{T} \int_{\omega}\left|\varphi_{R}\right|^{2} d x d t}\right\}=+\infty, \tag{1.7}
\end{equation*}
$$

[^1]where $I_{\mu} \triangleq\left(0, \mu\left\|a_{R}\right\|_{\infty}^{-1 / 3}\right]$.
Remark 1.1. As we shall see, the construction of Theorem 1.1 only requires potentials $a=a(x)$ which are independent of time.

According to (1.7) the inequality (1.2) is sharp in the sense that, for large potentials $a$, when the observation time is of order smaller than $\|a\|_{\infty}^{-1 / 3}$, the observability constant may not grow slower than an exponential factor of the order of $\exp \left(c\|a\|_{\infty}^{2 / 3}\right)$ for some $c>0$. This sharpness of the constant for small time is particularly important in view of the applications to semi-linear equations (see Section 7). In Section 4 (see Theorem 4.1), we give a slightly weaker extension to odd dimensions, for $N \geq 8$, and with a logarithmic correction. Getting exactly (1.7) for $n$ odd is an open problem.

Note that Theorems 1.1 and 4.1 do not exclude a possible improvement of the estimate for $N<8$ in odd space dimension, and for scalar equations in any space dimension. It will be clear from the proof of Theorem 1.1 that this last issue is closely related to that of the optimal decay at infinity for solutions of scalar equations on $\mathbb{R}^{n}$ of the form $\Delta u=q(x) u$ with bounded potential $q$, a problem that, to our knowledge, is not completely solved.

Theorem 1.1 is a consequence of the following known result:
Theorem B (Meshkov [19]). Assume that the space dimension is $n=2$. Then, there exists a nonzero complex-valued bounded potential $q=q(x)$ and a non-trivial complex valued solution $u=u(x)$ of

$$
\begin{equation*}
\Delta u=q u \quad \text { in } \mathbb{R}^{2} \tag{1.8}
\end{equation*}
$$

with the property that

$$
\begin{equation*}
|u(x)| \leq C \exp \left(-|x|^{4 / 3}\right), \quad \forall x \in \mathbb{R}^{2}, \tag{1.9}
\end{equation*}
$$

for some positive constant $C>0$.
Taking into account that the potential $q$ and the solution $u$ are complex valued, equation (1.8) can be viewed as an elliptic system with real valued coefficients and two components ( $N=2$ ).

As we shall see, Theorem 1.1 holds from the construction by Meshkov by scaling and localization arguments. The proof of Theorem 1.1 will be given in Section 4.

By separation of variables, Theorem B holds in any even dimension. The validity of this result in odd dimensions is, to the knowledge of the authors, an open problem. Section 3 is devoted to this issue. We obtain a slightly weaker version of Theorem B in $3-d$ for $\mathbb{C}^{4}$-valued solutions and with a potential $q$ growing at infinity in a logarithmic way. This construction is the main tool to prove Theorem 4.1 on the optimality of the observability constant in odd space-dimensions.

Let us mention some variants of the preceding results. One may assume that the potential $a$ is in $L^{\infty}\left(0, T ; L^{p}\left(\Omega ; \mathbb{R}^{N \times N}\right)\right)$, where $n \leq p<+\infty$, but is not necessarily bounded. For such $a$, an analogue of the observability inequality (1.2) is shown in Theorem 2.1. In this case, the constant $C_{3}^{*}$ has to be replaced by $\exp \left(C\|a\|_{L^{\infty}\left(0, T ; L^{p}\left(\Omega ; \mathbb{R}^{N \times N}\right)\right)}^{\alpha}\right)$, for some $\alpha$ depending on $p$. Suprisingly, Meshkov's construction is not sufficient to show that this constant is sharp. We shall also consider the case of convective potentials, that is first order operators of the form $a_{1} . \nabla$, where $a_{1}$ is bounded. The analogue of the observability inequality (1.2) is also stated in Theorem 2.1. In
this case, in space-dimension $n=2$, we are able to show the optimality of this inequality, using an easy adaptation of Meshkov's construction (see Section 6).

We now consider the wave equation and, more precisely, a system of $N \geq 1$ wave equations of the form

$$
\begin{cases}w_{t t}-\Delta w+a w=0, & \text { in } Q,  \tag{1.10}\\ w=0, & \text { on } \Sigma, \\ w(0, x)=w^{0}(x), \quad w_{t}(0, x)=w^{1}(x), & \text { in } \Omega,\end{cases}
$$

where the unknown function $w$ takes value in $\mathbb{R}^{N}$ and $a=a(t, x)$ is a matrix-valued potential as in (1.1).

Let $\omega$ be a nonempty open subset of $\Omega$. As above, we shall study the observability constant $D^{*}=D^{*}(T, a)$, defined, for fixed $a \in L^{\infty}\left(Q ; \mathbb{R}^{N \times N}\right)$ and $T>0$, as the smallest (possibly infinite) positive constant such that the following observability inequality holds:

$$
\begin{equation*}
\left\|w^{0}\right\|_{\left(L^{2}(\Omega)\right)^{N}}^{2}+\left\|w^{1}\right\|_{\left(H^{-1}(\Omega)\right)^{N}}^{2} \leq D^{*}(T, a) \int_{0}^{T} \int_{\omega}|w|^{2} d x d t, \tag{1.11}
\end{equation*}
$$

for every solution of (1.10).
Concerning the dependence of $D^{*}(T, a)$ on the potential $a$ the following can be said:

- Unlike the heat case, the existence of (finite) $D^{*}(T, a)$ is not guaranteed for all triple $(\Omega, T, \omega)$. In [5], Bardos, Lebeau and Rauch have established (under strong smoothness assumptions) an essentially necessary and sufficient condition for (1.11) to hold, the geometric control condition, which asserts that all rays of geometric optics in $\Omega$ enter the subdomain $\omega$ in an uniform time $T>0$. But the micro-local techniques used in their work do not seem to give the explicit dependence of the constant $D^{*}(T, a)$ on the potential $a$. In Section 2 we shall introduce a stronger condition on $(\Omega, T, \omega)$ which yields the existence of $D^{*}(T, a)$ for the multi-dimensional case. Indeed, by means of Carleman estimates, we show in Theorem 2.2 that under this condition:

$$
\begin{equation*}
D^{*}(T, a) \leq \exp \left(C(T)\left(1+\|a\|_{\infty}^{2 / 3}\right)\right) \tag{1.12}
\end{equation*}
$$

with a positive constant $C(T)$ that depends only on $(\Omega, T, \omega)$.

- Observability constant (1.12) is similar to that in [11] for the heat equation (see (1.3) above). Note however that there is an important difference between $C^{*}(T, a)$ for the heat equation and $D^{*}(T, a)$ for the wave one. Indeed, in (1.3), there is a term $\exp \left(C T\|a\|_{\infty}\right)$, due to the use of Gronwall's inequality. Because of this term, one needs to choose $T$ to be of the order (1.6) (thus small when $a$ is large), to get an obervability constant bounded by $\exp \left(C\|a\|^{2 / 3}\right)$. The situation is different for the wave equation. Indeed, as noted in [30], a modified energy estimate gives an upper bound for the evolution of the energy of the order of $\exp \left(C T\|a\|_{\infty}^{1 / 2}\right)$. As we shall see, this point is crucial to derive (1.12).
- In one space dimension, in [30], using sidewise energy estimates it was proved that

$$
\begin{equation*}
D^{*}(T, a) \leq C \exp \left(C\|a\|_{\infty}^{1 / 2}\right) . \tag{1.13}
\end{equation*}
$$

This is true whatever the number $N$ of components of the system is and even for $1-d$ wave equations with $B V$-coefficients in the principal part. But it may fail for equations with Hölder continuous coefficients (see [8]). Estimate (1.13) is known to be sharp. It is even sharp for the scalar $1-d$ wave equations with constant potentials.

However, as we shall see, (1.13) is not satisfied in space dimensions $n \geq 2$. Indeed, our next main goal of this paper is to show that estimate (1.12) is sharp in several space dimensions. We state the result in even space-dimension. As in the case of parabolic systems a similar, but weaker result holds in odd space-dimensions (Theorem 5.1).

Theorem 1.2. Assume that $n \geq 2$ is even and $N \geq 2$. Let $\omega$ be any given open non-empty subset of $\Omega$ such that $\Omega \backslash \bar{\omega} \neq \emptyset$. Then, for all $T>0$ there exist a constant $c>0$, a family of potentials $\left\{a_{R}\right\}_{R>0} \subset L^{\infty}\left(Q ; \mathbb{R}^{N \times N}\right)$ satisfying

$$
\left\|a_{R}\right\|_{\infty} \xrightarrow[R \rightarrow+\infty]{\longrightarrow}+\infty
$$

and a family of initial data $\left\{\left(w_{R}^{0}, w_{R}^{1}\right)\right\}_{R>0}$ in $\left(L^{2}(\Omega)\right)^{N} \times\left(H^{-1}(\Omega)\right)^{N}$ such that the corresponding solution $w_{R}$ of (1.10) satisfies

$$
\begin{equation*}
\lim _{R \rightarrow \infty}\left\{\frac{\left\|w_{R}^{0}\right\|_{\left(L^{2}(\Omega)\right)^{N}}^{2}+\left\|w_{R}^{1}\right\|_{\left(H^{-1}(\Omega)\right)^{N}}^{2}}{\exp \left(c\left\|a_{R}\right\|_{\infty}^{2 / 3}\right) \int_{0}^{T} \int_{\omega}\left|w_{R}\right|^{2} d x d t}\right\}=+\infty \tag{1.14}
\end{equation*}
$$

Remark 1.2. The potentials $a_{R}$ in Theorem 1.2 will be chosen to be time-invariant. Furthermore, as one shall see in the proof of Theorem 1.2 in Section 5, we actually choose the initial velocity $w_{R}^{1}=0$.

As in the context of the heat equation, the proof of this result is based on the construction by Meshkov in Theorem B. Theorem 1.2 shows that the estimate (1.12) is sharp for systems of wave equations in even dimensions. Hence, in this case, when $T$ is fixed, the observability constant $D^{*}(T, a)$ has to grow, at least, at the order of $\exp \left(C\|a\|_{\infty}^{2 / 3}\right)$ as $\|a\|_{\infty} \rightarrow \infty$. The problem of the optimality of the estimate (1.12) for scalar equations is open. As in the parabolic case, Theorem 3.2 yields a weaker version of Theorem 1.2 when $n$ is odd, that we state in Section 5 (see Theorem 5.1).

The rest of this paper is organized as follows. In Section 2, we shall show sharp observability estimates for parabolic and hyperbolic equations with both zero and, in some cases, first order potentials. We also consider the case of boundary observability for the wave equation. In Section 3 , we construct solutions on $\mathbb{R}^{n}$ of elliptic linear equations with a maximal speed of decay at infinity. These constructions generalize that of Meshkov, as it includes first order potentials, and odd $n$. They are the main tool of the three next sections, which are devoted to prove negative results. In Sections 4 and 5 we prove respectively the optimality of the observability estimates for heat and wave equations with zero order potentials, i.e., Theorems 1.1 and 1.2. In Section 6 we address the case of heat and wave equations with convective potentials, and show the optimality of the observability estimates which depend exponentially on the square of the $L^{\infty}$-norm of the convective potential. It is well-known that, in the context of both heat and wave equations, the
observability inequalities with explicit bounds in terms of the potentials are intimately related with the optimal growth rates for the control of semilinear equations. This issue will be briefly discussed in Section 7. Finally, in Section 8 we comment some closely related issues and open problems.

## 2 Sharp observability estimates for hyperbolic and parabolic equations with potentials

This section is devoted to show sharp observability estimates for hyperbolic and parabolic equations with zero and first order potentials. We shall consider elliptic operators with variable coefficients, which does not change the proof with respect to the simpler case of Laplace's operator, and gives the same dependence of the observability constant on the potentials. As it was mentionned in the introduction, we also consider the case of zero order potentials which are in $L^{p}$-spaces in the space variable, with $n \leq p \leq \infty$.

### 2.1 Statement of the results

In the sequel, we fix real valued functions $b^{i j} \in C^{1}(\bar{\Omega})$ satisfying:

$$
\begin{equation*}
b^{i j}(x)=b^{j i}(x), \quad \forall x \in \bar{\Omega}, i, j=1,2, \cdots, n, \tag{2.1}
\end{equation*}
$$

and

$$
\begin{equation*}
\sum_{i, j=1}^{n} b^{i j}(x) \xi_{i} \xi_{j} \geq \beta|\xi|^{2}, \quad \forall(x, \xi) \equiv\left(x, \xi_{1}, \cdots, \xi_{n}\right) \in \bar{\Omega} \times \mathbb{R}^{n}, \tag{2.2}
\end{equation*}
$$

for some constant $\beta>0$.
We also consider $\mathbb{R}^{N \times N_{-}}$-valued functions $a, a_{1}^{k}(k=1, \cdots, n)$ and $a_{2}$ on $Q$ satisfying:
(H1) $a \in L^{\infty}\left(0, T ; L^{p}\left(\Omega ; \mathbb{R}^{N \times N}\right)\right)$ for some $p \in[n, \infty]$, and $a_{1}^{1}, \cdots, a_{1}^{n}, a_{2} \in L^{\infty}\left(Q ; \mathbb{R}^{N \times N}\right)$.
Put

$$
\begin{equation*}
r_{0} \triangleq\|a\|_{L^{\infty}\left(0, T ; L^{p}\left(\Omega ; \mathbb{R}^{N \times N}\right)\right)}, \quad r_{1} \triangleq \sum_{k=1}^{n}\left\|a_{1}^{k}\right\|_{\infty}, \quad r_{2} \triangleq r_{1}+\left\|a_{2}\right\|_{\infty} \tag{2.3}
\end{equation*}
$$

### 2.1.1 The parabolic system

We consider first the following parabolic system:

$$
\begin{cases}\varphi_{t}-\sum_{i, j=1}^{n}\left(b^{i j}(x) \varphi_{x_{i}}\right)_{x_{j}}=a \varphi+\sum_{k=1}^{n} a_{1}^{k} \varphi_{x_{k}}, & \text { in } Q  \tag{2.4}\\ \varphi=0, & \text { on } \Sigma, \\ \varphi(0)=\varphi^{0}, & \text { in } \Omega\end{cases}
$$

We have the following observability estimate for system (2.4).

Theorem 2.1. Let $\omega$ be an open non-empty subset of $\Omega$ and $b^{i j}(\cdot) \in C^{1}(\bar{\Omega})$ satisfy (2.1)-(2.2). Then, there exists a constant $C=C(\Omega, \omega)>0$ depending only on $\Omega$ and $\omega$ such that for any time $T>0$, potentials a and $a_{1}^{k}(k=1, \cdots, n)$ satisfying (H1), and initial data $\varphi^{0} \in\left(L^{2}(\Omega)\right)^{N}$, the corresponding solution $\varphi \in C\left([0, T] ;\left(L^{2}(\Omega)\right)^{N}\right)$ of (2.4) satisfies

$$
\begin{equation*}
\|\varphi(T)\|_{\left(L^{2}(\Omega)\right)^{N}}^{2} \leq \exp \left\{C\left[1+\frac{1}{T}+T r_{0}+r_{0}^{\frac{1}{3 / 2-n / p}}+(1+T) r_{1}^{2}\right]\right\} \int_{0}^{T} \int_{\omega}|\varphi|^{2} d x d t \tag{2.5}
\end{equation*}
$$

By adapting an argument in Step 5 in the proof of Theorem 2.4 in Subsection 2.2 (for the hyperbolic system), the proof of Theorem 2.1 is almost the same as that in [9, Theorem 2.3] and [11, Theorem 1.2]. Hence we omit the details. Theorem 1.1 shows that, when $\left(b^{i j}\right)_{n \times n}=I$ (the identity matrix), the exponent $2 / 3$ in $r_{0}^{2 / 3}$ (for the special case $p=\infty$ in the estimate (2.5)) is sharp. As we shall see in Section 6, the quadratic dependence on $r_{1}$ is also sharp. The problem of the optimality of the term in $r_{0}^{\frac{1}{3 / 2-n / p}}$ is completely open when $p<\infty$.

### 2.1.2 The hyperbolic system

Let us now consider the following hyperbolic system:

$$
\begin{cases}w_{t t}-\sum_{i, j=1}^{n}\left(b^{i j}(x) w_{x_{i}}\right)_{x_{j}}=a w+\sum_{k=1}^{n} a_{1}^{k} w_{x_{k}}+a_{2} w_{t}, & \text { in } Q,  \tag{2.6}\\ w=0, & \text { on } \Sigma, \\ w(0)=w^{0}, w_{t}(0)=w^{1}, & \text { in } \Omega,\end{cases}
$$

where $w=\left(w_{1}, \cdots, w_{N}\right)^{\top}$ is a $\mathbb{R}^{N}$-valued unknown. Under some assumption on the observation domain, we will show boundary and interior observability inequalities for system (2.6), with an explicit dependence on the constants $r_{0}$ (and, in the case of boundary observability also $r_{2}$ ). For this, we introduce the following condition.
(H2) There exists a function $d(\cdot) \in C^{2}(\bar{\Omega})$ satisfying the following:
i) For some constant $\mu_{0}>0$,

$$
\begin{equation*}
\sum_{i, j=1}^{n}\left\{\sum_{i^{\prime}, j^{\prime}=1}^{n}\left[2 b^{i j^{\prime}}\left(b^{i^{\prime} j} d_{x_{i^{\prime}}}\right)_{x_{j^{\prime}}}-b_{x_{j^{\prime}}}^{i j} i^{i^{\prime} j^{\prime}} d_{x_{i^{\prime}}}\right]\right\} \xi_{i} \xi_{j} \geq \mu_{0} \sum_{i, j=1}^{n} b^{i j} \xi_{i} \xi_{j}, \quad \forall(x, \xi) \in \bar{\Omega} \times \mathbb{R}^{n} \tag{2.7}
\end{equation*}
$$

ii) The function $d(\cdot)$ does not have any critical point in $\bar{\Omega}$, i.e.,

$$
\begin{equation*}
\min _{x \in \bar{\Omega}}|\nabla d(x)|>0 . \tag{2.8}
\end{equation*}
$$

Denote by $\nu=\nu(x)=\left(\nu_{1}, \nu_{2}, \cdots, \nu_{n}\right)$ the unit outward normal vector of $\Omega$ at $x \in \Gamma$. For the function $d(\cdot)$ satisfying Condition (H2), we introduce the following set:

$$
\begin{equation*}
\Gamma_{0} \triangleq\left\{x \in \Gamma \mid \sum_{i, j=1}^{n} b^{i j} \nu_{i} d_{x_{j}}>0\right\} \tag{2.9}
\end{equation*}
$$

Note that for the case of $\left(b^{i j}\right)_{n \times n}=I$, by choosing $d(x)=\left|x-x_{0}\right|^{2}$ with any fixed $x_{0} \in \mathbb{R}^{n} \backslash \bar{\Omega}$, (H2) is satisfied with $\mu_{0}=4$ and (2.7) holds with an equality. In this case, $\Gamma_{0}$ in (2.9) is given by

$$
\left\{x \in \Gamma \mid\left(x-x_{0}\right) \cdot \nu(x)>0\right\}
$$

which coincides with the subset of the boundary arising usually in the context of the multiplier method ([17]).

On the other hand, it is easy to check that, if $d(\cdot) \in C^{2}(\bar{\Omega})$ satisfies (2.7), then for any given constants $\alpha \geq 1$ and $\beta \in \mathbb{R}$, the function

$$
\begin{equation*}
\hat{d}=\hat{d}(x) \triangleq \alpha d(x)+\beta \tag{2.10}
\end{equation*}
$$

still satisfies Condition (H2) with $\mu_{0}$ replaced by $\alpha \mu_{0}$, the set $\Gamma_{0}$ remaining unchanged. Hence, without loss of generality, we may assume that

$$
\left\{\begin{array}{l}
(2.7) \text { holds with } \mu_{0} \geq 4,  \tag{2.11}\\
\frac{1}{4} \sum_{i, j=1}^{n} b^{i j}(x) d_{x_{i}}(x) d_{x_{j}}(x) \geq \max _{x \in \bar{\Omega}} d(x) \geq \min _{x \in \bar{\Omega}} d(x)>0, \quad \forall x \in \bar{\Omega} .
\end{array}\right.
$$

In what follows, put

$$
\begin{equation*}
R_{1} \triangleq \max _{x \in \bar{\Omega}} \sqrt{d(x)}, \quad T_{0} \triangleq 2 \inf \left\{R_{1} \mid d(\cdot) \text { satisfies }(2.11)\right\} . \tag{2.12}
\end{equation*}
$$

For the interior observation, we introduce the following assumption:
(H3) There is a constant $\delta>0$ such that

$$
\begin{equation*}
\omega=\mathcal{O}_{\delta}\left(\Gamma_{0}\right) \bigcap \Omega \tag{2.13}
\end{equation*}
$$

where $\mathcal{O}_{\delta}\left(\Gamma_{0}\right)=\left\{x \in \mathbb{R}^{n}| | x-x^{\prime} \mid<\delta\right.$ for some $\left.x^{\prime} \in \Gamma_{0}\right\}$.
In other words, the observation subdomain $\omega$ is assumed to be a neighborhood of a subset of the boundary satisfying the condition above for boundary observability.

In the rest of this section, we will use $C=C(T, \Omega)$ to denote a generic positive constant which may vary from line to line. Our observability estimates for system (2.6) are stated as follows.
Theorem 2.2. Let $b^{i j} \in C^{1}(\bar{\Omega})$ satisfy (2.1)-(2.2). Assume that conditions (H2)-(H3) hold. Let $T>T_{0}$, where $T_{0}$ is defined in (2.12). Then the following two assertions hold:
i) boundary observability: There is a constant $C>0$ such that for any potentials a, $a_{1}^{1}, \cdots$, $a_{1}^{n}$ and $a_{2}$ satisfying (H1) for some $p \in[n,+\infty]$, and initial data $\left(w^{0}, w^{1}\right) \in\left(H_{0}^{1}(\Omega)\right)^{N} \times\left(L^{2}(\Omega)\right)^{N}$, the corresponding weak solution:

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

of (2.6) satisfies

$$
\begin{equation*}
\left\|w^{0}\right\|_{\left(H_{0}^{1}(\Omega)\right)^{N}}+\left\|w^{1}\right\|_{\left(L^{2}(\Omega)\right)^{N}} \leq \exp \left(C\left(1+r_{0}^{\frac{1}{3 / 2-n / p}}+r_{2}^{2}\right)\right)\left\|\frac{\partial w}{\partial \nu}\right\|_{\left(L^{2}\left((0, T) \times \Gamma_{0}\right)\right)^{N}} . \tag{2.14}
\end{equation*}
$$

ii) interior observability: If $a_{1}^{k} \equiv 0(k=1, \cdots, n)$ and $a_{2} \equiv 0$, then there is a constant $C>0$ such that for any potential a satisfying (H1) and initial data $\left(w^{0}, w^{1}\right) \in\left(L^{2}(\Omega)\right)^{N} \times\left(H^{-1}(\Omega)\right)^{N}$, the corresponding weak solution:

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

of (2.6) satisfies

$$
\begin{equation*}
\left\|w^{0}\right\|_{\left(L^{2}(\Omega)\right)^{N}}+\left\|w^{1}\right\|_{\left(H^{-1}(\Omega)\right)^{N}} \leq \exp \left(C\left(1+r_{0}^{\frac{1}{3 / 2-n / p}}\right)\right)\|w\|_{\left(L^{2}((0, T) \times \omega)\right)^{N}} \tag{2.15}
\end{equation*}
$$

Part $i i$ of Theorem 2.2 is the internal observability inequality for the wave equation announced in (1.11). Note that unlike the case of the heat equation (Theorem 2.1) and the boundary observability for the wave equation (part $i$ of Theorem 2.2), we do not give a generalization of this inequality to first order potentials. The proof of such a generalization would yield new technical difficulties, essentially due to the fact that in view of control-theoretic applications, we treat the case of the adjoint system (with initial data in $\left(L^{2}(\Omega) \times H^{-1}(\Omega)\right)^{N}$ ), instead of the usual energy space (with initial data in $\left.\left(H_{0}^{1}(\Omega) \times L^{2}(\Omega)\right)^{N}\right)$.

As in the parabolic case, Theorem 1.2 shows that the exponent $2 / 3$ in the estimate $r_{0}^{2 / 3}$ (in (2.15) for the special case $p=\infty$ ) is sharp. As we shall see in Section 6 , the exponent 2 in the estimate $r_{2}^{2}$ (in (2.14)) is sharp, too.

In the rest of this section, to simplify the presentation, we only consider the case $N=1$. The proof in the case $N \geq 2$ is exactly the same. The proof of Theorem 2.2 will be given in Subsection 2.3. For this, we need to show a global Carleman estimate for hyperbolic operators, which is the object of the next subsection.

### 2.2 Global Carleman estimate for hyperbolic operators

Recall (2.12) for the definitions of $R_{1}$ and $T_{0}$. Let $T>T_{0}$ be given. We may choose $d$ such that

$$
\begin{equation*}
T>2 R_{1} \tag{2.16}
\end{equation*}
$$

By (2.16), one may choose a constant $c \in(0,1)$ so that

$$
\begin{equation*}
\left(\frac{2 R_{1}}{T}\right)^{2}<c<\frac{2 R_{1}}{T} \tag{2.17}
\end{equation*}
$$

Put

$$
\begin{equation*}
\phi=\phi(t, x) \triangleq d(x)-c(t-T / 2)^{2} \tag{2.18}
\end{equation*}
$$

with $T$ and $c$ satisfying respectively (2.16) and (2.17).
Define a formal differential operator $\mathcal{P}$ by

$$
\begin{equation*}
\mathcal{P} u \triangleq u_{t t}-\sum_{i, j=1}^{n}\left(b^{i j}(x) u_{x_{i}}\right)_{x_{j}} \tag{2.19}
\end{equation*}
$$

By [14, (5.15) in the proof of Theorem 5.1], we have the following Carleman estimate.

Theorem 2.3. Let $b^{i j} \in C^{1}(\bar{\Omega})$ satisfy (2.1)-(2.2). Assume that condition (H2) holds and $\Gamma_{0}$ is given by (2.9). Then there exists a $\lambda_{0}>1$ such that for all $\lambda \geq \lambda_{0}$ and all $u \in H_{0}^{1}(Q)$ with $\mathcal{P} u \in L^{2}(Q)$, it holds:

$$
\begin{equation*}
\lambda \int_{Q} e^{2 \lambda \phi}\left(\lambda^{2} u^{2}+u_{t}^{2}+|\nabla u|^{2}\right) d x d t \leq C\left(\int_{Q} e^{2 \lambda \phi}|\mathcal{P} u|^{2} d x d t+\lambda \int_{0}^{T} \int_{\Gamma_{0}} e^{2 \lambda \phi}\left|\frac{\partial u}{\partial \nu}\right|^{2} d x d t\right) . \tag{2.20}
\end{equation*}
$$

In order to prove the interior observability result in Theorem 2.2, we also need the following:
Theorem 2.4. Let $b^{i j} \in C^{1}(\bar{\Omega})$ satisfy (2.1)-(2.2), and $a \in L^{\infty}\left(0, T ; L^{p}(\Omega)\right)$ with $p \in[n, \infty]$. Assume that conditions (H2)-(H3) hold. Then there exists a $\lambda_{0}>1$ such that for all $\lambda \geq \lambda_{0}$, any $u \in C\left([0, T] ; L^{2}(\Omega)\right)$ satisfying $u(0, x)=u(T, x)=0$ for $x \in \Omega, \mathcal{P} u \in H^{-1}(Q)$ and

$$
\begin{equation*}
(u, \mathcal{P} \eta)_{L^{2}(Q)}=\langle\mathcal{P} u, \eta\rangle_{H^{-1}(Q), H_{0}^{1}(Q)}, \quad \forall \eta \in H_{0}^{1}(Q) \text { with } \mathcal{P} \eta \in L^{2}(Q) \tag{2.21}
\end{equation*}
$$

it holds

$$
\begin{gather*}
\lambda\left\|e^{\lambda \phi} u\right\|_{L^{2}(Q)}^{2} \leq C\left(\left\|e^{\lambda \phi}(\mathcal{P} u-a u)\right\|_{H^{-1}(Q)}^{2}+\frac{1}{\lambda^{2(1-n / p)}}\left\|e^{\lambda \phi} a u\right\|_{L^{2}\left(0, T ; H^{-n / p}(\Omega)\right)}^{2}\right.  \tag{2.22}\\
\left.+\lambda^{2}\left\|e^{\lambda \phi} u\right\|_{L^{2}((0, T) \times \omega)}^{2}\right)
\end{gather*}
$$

Proof. We shall borrow some idea from [16], which consists in applying (2.21) to some special choice of $\eta$ with $\mathcal{P} \eta=\cdots+\lambda e^{2 \lambda \phi} u$, which yields the desired term $\lambda\left\|e^{\lambda \phi} u\right\|_{L^{2}(Q)}^{2}$ and reduces the estimate (2.22) to an estimate on $\|\eta\|_{H_{0}^{1}(Q)}$.

In [14, (7.20) and (7.21) in the proof of Theorem 7.1], it is shown, using an optimization argument, that there is a $\lambda_{0}$ greater than 1 so that for any $\lambda$ greater than $\lambda_{0}$ and any $u$ in $C\left([0, T] ; L^{2}(\Omega)\right)$ vanishing at $t=0$ and $t=T$, there exist $(\check{z}, \check{r})$ in $H_{0}^{1}(Q) \times L^{2}((0, T) \times \Omega)$, such that:

$$
\begin{cases}\operatorname{supp} \check{r} \subset \overline{(0, T) \times \omega}, &  \tag{2.23}\\ \mathcal{P} \check{z}=\check{r}+\lambda e^{2 \lambda \phi} u, & \text { in } Q \\ \check{z}=0, & \text { on } \partial Q\end{cases}
$$

and for some constant $C>0$, independent of $\lambda$, it holds

$$
\begin{equation*}
\int_{Q} e^{-2 \lambda \phi}\left(|\nabla \check{z}|^{2}+\check{z}_{t}^{2}+\lambda^{2} \check{z}^{2}\right) d x d t+\frac{1}{\lambda^{2}} \int_{0}^{T} \int_{\omega} e^{-2 \lambda \phi} \check{r}^{2} d x d t \leq C \lambda \int_{Q} e^{2 \lambda \phi} u^{2} d x d t . \tag{2.24}
\end{equation*}
$$

Now, by (2.21) with $\eta$ replaced by $\check{z}$ above, one gets

$$
\left(u, \check{r}+\lambda e^{2 \lambda \phi} u\right)_{L^{2}(Q)}=\langle\mathcal{P} u, \check{z}\rangle_{H^{-1}(Q), H_{0}^{1}(Q)} .
$$

Hence, noting supp $\check{r} \subset \overline{(0, T) \times \omega}$ :

$$
\begin{align*}
& \lambda\left\|e^{\lambda \phi} u\right\|_{L^{2}(Q)}^{2}=\langle\mathcal{P} u, \check{z}\rangle_{H^{-1}(Q), H_{0}^{1}(Q)}-(u, \check{r})_{L^{2}((0, T) \times \omega)} \\
&=\langle\mathcal{P} u-a u, \check{z}\rangle_{H^{-1}(Q), H_{0}^{1}(Q)}+(a u, \check{z})_{L^{2}(Q)}-(u, \check{r})_{L^{2}((0, T) \times \omega)} \\
& \leq\left\|e^{\lambda \phi}(\mathcal{P} u-a u)\right\|_{H^{-1}(Q)}\left\|e^{-\lambda \phi} \check{z}\right\|_{H_{0}^{1}(Q)} \\
&\left.+\left\|e^{\lambda \phi} a u\right\|_{L^{2}\left(0, T ; H^{-n / p}(\Omega)\right)}\right) \mid e^{-\lambda \phi} \|_{L^{2}\left(0, T ; H_{0}^{n / p}(\Omega)\right)} \\
& \quad+\left\|e^{\lambda \phi} u\right\|_{L^{2}((0, T) \times \omega)}\left\|e^{-\lambda \phi} \check{r}\right\|_{L^{2}((0, T) \times \omega)}  \tag{2.25}\\
& \leq C \sqrt{\mathcal{Q}}\left[\left\|e^{-\lambda \phi} \check{z}\right\|_{H_{0}^{1}(Q)}^{2}+\lambda^{2(1-n / p)}\left\|e^{-\lambda \phi} \check{z}\right\|_{L^{2}\left(0, T ; H_{0}^{n / p}(\Omega)\right)}^{2}\right. \\
&\left.+\frac{1}{\lambda^{2}}\left\|e^{-\lambda \phi} \check{r}\right\|_{L^{2}((0, T) \times \omega)}^{2}\right]^{1 / 2},
\end{align*}
$$

where

$$
\mathcal{Q} \triangleq\left\|e^{\lambda \phi}(\mathcal{P} u-a u)\right\|_{H^{-1}(Q)}^{2}+\frac{1}{\lambda^{2(1-n / p)}}\left\|e^{\lambda \phi} a u\right\|_{L^{2}\left(0, T ; H^{-n / p}(\Omega)\right)}^{2}+\lambda^{2}\left\|e^{\lambda \phi} u\right\|_{L^{2}((0, T) \times \omega)}^{2}
$$

is the right-hand term of (2.22). On the other hand, for any $f \in H^{1}\left(\mathbb{R}^{n}\right)$, by Hölder's inequality, one has

$$
\|f\|_{H^{n / p}\left(\mathbb{R}^{n}\right)}^{2}=\int_{\mathbb{R}^{n}}\left(1+|\xi|^{2}\right)^{n / p}|\hat{f}(\xi)|^{2 n / p}|\hat{f}(\xi)|^{2(1-n / p)} d \xi \leq\|f\|_{H^{1}\left(\mathbb{R}^{n}\right)}^{n / p}\|f\|_{L^{2}\left(\mathbb{R}^{n}\right)}^{1-n / p}
$$

This yields immediately:

$$
\begin{equation*}
\|g\|_{H_{0}^{n / p}(\Omega)}^{2} \leq C\|g\|_{H_{0}^{1}(\Omega)}^{n / p}\|g\|_{L^{2}(\Omega)}^{1-n / p}, \quad \forall g \in H_{0}^{1}(\Omega) \tag{2.26}
\end{equation*}
$$

for some constant $C>0$, independent of $g$. Hence, for any $h \in L^{2}\left(0, T ; H_{0}^{1}(\Omega)\right)$ :

$$
\begin{equation*}
\|h\|_{L^{2}\left(0, T ; H_{0}^{n / p}(\Omega)\right)} \leq C\|h\|_{L^{2}\left(0, T ; H_{0}^{1}(\Omega)\right)}^{n / p}\|h\|_{L^{2}(Q)}^{1-n / p} . \tag{2.27}
\end{equation*}
$$

Now, by (2.27) and using Young's inequality, it follows

$$
\begin{align*}
\lambda^{2(1-n / p)}\left\|e^{-\lambda \phi} \check{z}\right\|_{L^{2}\left(0, T ; H_{0}^{n / p}(\Omega)\right)}^{2} & \leq C \lambda^{2(1-n / p)}\left\|e^{-\lambda \phi} \check{z}\right\|_{L^{2}\left(0, T ; H_{0}^{1}(\Omega)\right)}^{2 n / p}\left\|e^{-\lambda \phi} \check{z}\right\|_{L^{2}(Q)}^{2(1-n / p)}  \tag{2.28}\\
& \leq C\left[\left\|e^{-\lambda \phi} \check{z}\right\|_{L^{2}\left(0, T ; H_{0}^{1}(\Omega)\right)}^{2}+\lambda^{2}\left\|e^{-\lambda \phi} \check{z}\right\|_{L^{2}(Q)}^{2}\right] .
\end{align*}
$$

Finally, using (2.24), (2.25) and (2.28), we arrive at the desired estimate (2.22). This completes the proof of Theorem 2.4.

### 2.3 Proof of sharp observability estimates for hyperbolic equations with potentials

We now prove Theorem 2.2. The main idea is to use the Carleman estimate in Theorems 2.3-2.4. The proof is divided into several steps.

Step 1: Choice of a cutoff function. Note that our $w$ satisfying (2.6) does not necessarily vanish at $t=0, T$. Therefore we need to introduce a suitable cutoff function. To this end, set (recall (2.11))

$$
\begin{equation*}
T_{i} \triangleq T / 2-\varepsilon_{i} T, \quad T_{i}^{\prime} \triangleq T / 2+\varepsilon_{i} T, \quad R_{0} \triangleq \min _{x \in \bar{\Omega}} \sqrt{d(x)}(>0) \tag{2.29}
\end{equation*}
$$

where $i=0,1 ; 0<\varepsilon_{0}<\varepsilon_{1}<1 / 2$ will be given below.
¿From (2.17), (2.18) and the definition (2.12) of $R_{1}$, it is easy to see that

$$
\begin{equation*}
\phi(0, x)=\phi(T, x) \leq R_{1}^{2}-c T^{2} / 4<0, \quad \forall x \in \Omega . \tag{2.30}
\end{equation*}
$$

Therefore there exists an $\varepsilon_{1} \in(0,1 / 2)$, close to $1 / 2$, such that

$$
\begin{equation*}
\phi(t, x) \leq R_{1}^{2} / 2-c T^{2} / 8<0, \quad \forall(t, x) \in\left(\left(0, T_{1}\right) \bigcup\left(T_{1}^{\prime}, T\right)\right) \times \Omega \tag{2.31}
\end{equation*}
$$

with $T_{1}$ and $T_{1}^{\prime}$ given by (2.29). Further, by (2.18), we see that

$$
\phi(T / 2, x)=d(x) \geq R_{0}^{2}, \quad \forall x \in \Omega .
$$

Hence, one can find an $\varepsilon_{0} \in(0,1 / 2)$, close to 0 , such that

$$
\begin{equation*}
\phi(t, x) \geq R_{0}^{2} / 2, \quad \forall(t, x) \in\left(T_{0}, T_{0}^{\prime}\right) \times \Omega, \tag{2.32}
\end{equation*}
$$

with $T_{0}$ and $T_{0}^{\prime}$ given by (2.29). We now choose a nonnegative function $\xi \in C_{0}^{\infty}(0, T)$ so that

$$
\begin{equation*}
\xi(t) \equiv 1 \quad \text { in }\left(T_{1}, T_{1}^{\prime}\right) . \tag{2.33}
\end{equation*}
$$

We start to show the second assertion of Theorem 2.2, whose proof is more technical.
Step 2: An intermediate inequality. Assume that the assumptions of the second part of Theorem 2.2 hold. We first shall apply Carleman inequality (2.22) to prove the following observability inequality:

$$
\begin{align*}
& \exists \lambda_{1}>0, \quad \forall \lambda \geq\left(1+r_{0}^{\frac{1}{3 / 2-n / p}}\right) \lambda_{1},  \tag{2.34}\\
& \lambda \int_{Q} e^{2 \lambda \phi} w^{2} d x d t \leq C\left(\lambda^{2} \int_{0}^{T} \int_{\omega} e^{2 \lambda \phi} w^{2} d x d t+\|w\|_{L^{2}(J \times \Omega)}^{2}\right) .
\end{align*}
$$

where $J \triangleq\left(0, T_{1}\right) \cup\left(T_{1}^{\prime}, T\right)$. Note that in the right-hand side of (2.34) the square of the $L^{2}$ norm of $w$ on all $\Omega$ appears, but only for time smaller than $T_{1}$ or greater than $T_{1}^{\prime}$. To get read of these term, we will need to use, in the next step, a modified energy method.

Recall that $\xi w$ vanishes at $t=0, T$. Hence, by Theorem 2.4, for any $\lambda \geq \lambda_{0}$, we have

$$
\begin{align*}
\lambda \int_{Q} e^{2 \lambda \phi}(\xi w)^{2} d x d t \leq C & {\left[\left\|e^{\lambda \phi}(\mathcal{P}(\xi w)-a \xi w)\right\|_{H^{-1}(Q)}^{2}\right.}  \tag{2.35}\\
& \left.+\frac{1}{\lambda^{2(1-n / p)}}\left\|e^{\lambda \phi} a \xi w\right\|_{L^{2}\left(0, T ; H^{-n / p}(\Omega)\right)}^{2}+\lambda^{2} \int_{0}^{T} \int_{\omega} e^{2 \lambda \phi} w^{2} d x d t\right] .
\end{align*}
$$

By equation (2.6), recalling $a_{1}^{k} \equiv 0(k=1, \cdots, n)$ and $a_{2} \equiv 0$, and noting (2.31) and (2.33), we have

$$
\begin{align*}
\left\|e^{\lambda \phi}(\mathcal{P}(\xi w)-a \xi w)\right\|_{H^{-1}(Q)} & =\left\|e^{\lambda \phi}\left(2 \xi_{t} w_{t}+w \xi_{t t}\right)\right\|_{H^{-1}(Q)} \\
& =\sup _{|f|_{H_{0}^{1}(Q)}=1} \int_{Q} e^{\lambda \phi}\left(2 \xi_{t} w_{t}+w \xi_{t t}\right) f d x d t  \tag{2.36}\\
& =\sup _{|f|_{H_{0}^{1}(Q)}=1} \int_{Q} e^{\lambda \phi} w\left(-\xi_{t t} f-2 \xi_{t} f_{t}-2 \lambda \phi_{t} \xi_{t} f\right) d x d t \\
& \leq C e^{\left(R_{1}^{2} / 2-c T^{2} / 8\right) \lambda}(1+\lambda)\|w\|_{L^{2}(J \times \Omega)} .
\end{align*}
$$

Recalling the definition of $r_{0}$ in (2.3) and noting that the Sobolev embedding $H_{0}^{n / p}(\Omega) \hookrightarrow L^{\frac{2 p}{p-2}}(\Omega)$, implies by duality the embedding $L^{\frac{2 p}{p+2}}(\Omega) \hookrightarrow H^{-n / p}(\Omega)$, we get:

$$
\begin{align*}
\left\|e^{\lambda \phi} a \xi w\right\|_{L^{2}\left(0, T ; H^{-n / p}(\Omega)\right)} & \leq\left\|e^{\lambda \phi} a \xi w\right\|_{L^{2}\left(0, T ; L^{2 p /(p+2)}(\Omega)\right)}  \tag{2.37}\\
& \leq C r_{0}\left\|e^{\lambda \phi} w\right\|_{L^{2}(Q)},
\end{align*}
$$

where at the second line we simply used Hölder's inequality. Further, by (2.31) and (2.33), we have

$$
\begin{align*}
\int_{Q} e^{2 \lambda \phi}(\xi w)^{2} d x d t & =\int_{Q} e^{2 \lambda \phi} w^{2} d x d t-\int_{Q} e^{2 \lambda \phi}\left(1-\xi^{2}\right) w^{2} d x d t  \tag{2.38}\\
& \geq \int_{Q} e^{2 \lambda \phi} w^{2} d x d t-C e^{\left(R_{1}^{2}-c T^{2} / 4\right) \lambda}\|w\|_{L^{2}(J \times \Omega)}^{2}
\end{align*}
$$

Combining (2.35)-(2.38), we conclude that there is a constant $C_{1}=C_{1}(T, \Omega)$, independent of $\lambda$ and $r_{0}$, such that:

$$
\begin{align*}
& \lambda \int_{Q} e^{2 \lambda \phi} w^{2} d x d t \leq C_{1}\left[\frac{r_{0}^{2}}{\lambda^{2(1-n / p)}} \int_{Q} e^{2 \lambda \phi} w^{2} d x d t\right.  \tag{2.39}\\
&\left.\quad+\lambda^{2} \int_{0}^{T} \int_{\omega} e^{2 \lambda \phi} w^{2} d x d t+e^{\left(R_{1}^{2}-c T^{2} / 4\right) \lambda}\left(1+\lambda^{2}\right)\|w\|_{L^{2}(J \times \Omega)}^{2}\right]
\end{align*}
$$

Since $R_{1}^{2}-c T^{2} / 4<0$, one may find a $\lambda_{1}^{\prime} \geq \lambda_{0}$ such that $e^{\left(R_{1}^{2}-c T^{2} / 4\right) \lambda}\left(1+\lambda^{2}\right)<1$ for all $\lambda \geq \lambda_{1}^{\prime}$.
Take a $\lambda_{1} \geq \lambda_{1}^{\prime}$ and such that:

$$
\begin{equation*}
\lambda \geq r_{0}^{\frac{1}{3 / 2-n / p}} \lambda_{1} \Longrightarrow \lambda-C_{1} \frac{r_{0}^{2}}{\lambda^{2(1-n / p)}} \geq \frac{\lambda}{2} . \tag{2.40}
\end{equation*}
$$

For such a choice of $\lambda_{1}$, the desired inequality (2.34) follows from (2.39).
Step 3: A modified energy method. From (2.32), we see that

$$
\begin{equation*}
\int_{Q} e^{2 \lambda \phi} w^{2} d x d t \geq e^{R_{0}^{2} \lambda} \int_{T_{0}}^{T_{0}^{\prime}} \int_{\Omega} w^{2} d x d t \tag{2.41}
\end{equation*}
$$

Put

$$
\begin{equation*}
E(t) \triangleq \frac{1}{2}\left[\|w(t, \cdot)\|_{L^{2}(\Omega)}^{2}+\left\|w_{t}(t, \cdot)\right\|_{H^{-1}(\Omega)}^{2}\right] . \tag{2.42}
\end{equation*}
$$

For any $S_{0} \in\left(T_{0}, T / 2\right)$ and $S_{0}^{\prime} \in\left(T / 2, T_{0}^{\prime}\right)$, by means of the classical energy estimate, one has

$$
\begin{equation*}
\int_{S_{0}}^{S_{0}^{\prime}} E(t) d t \leq C\left(1+r_{0}\right) \int_{T_{0}}^{T_{0}^{\prime}} \int_{\Omega} w^{2} d x d t \tag{2.43}
\end{equation*}
$$

We claim that, there is a constant $C>0$ such that

$$
\begin{equation*}
E(t) \leq C e^{C r_{0}^{\frac{1}{2-n / p}}} E(s), \quad \forall t, s \in[0, T] . \tag{2.44}
\end{equation*}
$$

Note however that this does not follow from the usual energy method. Instead, we need to use the duality argument and adopt a modified energy estimate introduced in [30]. For this, for any $\left(z^{0}, z^{1}\right) \in H_{0}^{1}(\Omega) \times L^{2}(\Omega)$, we introduce the following system

$$
\begin{cases}z_{t t}-\sum_{i, j=1}^{n}\left(b^{i j}(x) z_{x_{i}}\right)_{x_{j}}=a z, & \text { in } Q,  \tag{2.45}\\ z=0, & \text { on } \Sigma, \\ z(T)=z^{0}, \quad z_{t}(T)=z^{1}, & \text { in } \Omega .\end{cases}
$$

Denote a (modified) energy of system (2.45) by

$$
\begin{equation*}
\mathcal{E}(t)=\frac{1}{2} \int_{\Omega}\left[\left|z_{t}(t, x)\right|^{2}+\sum_{i, j=1}^{n} b^{i j}(x) z_{x_{i}}(t, x) z_{x_{j}}(t, x)+r_{0}^{\frac{2}{2-n / p}}|z(t, x)|^{2}\right] d x . \tag{2.46}
\end{equation*}
$$

Then, by (2.45) and recalling the definition of $r_{0}$ in (2.3), it follows

$$
\begin{equation*}
\frac{d \mathcal{E}(t)}{d t}=\int_{\Omega} a z z_{t} d x+r_{0}^{\frac{2}{2-n / p}} \int_{\Omega} z z_{t} d x . \tag{2.47}
\end{equation*}
$$

Put $p_{1}=\frac{2 p}{n-2}$ and $p_{2}=\frac{2 p}{p-n}$. Noting that $\frac{1}{p}+\frac{1}{p_{1}}+\frac{1}{p_{2}}+\frac{1}{2}=1$ and $\frac{1}{2(n / p)^{-1}}+\frac{1}{2(1-n / p)^{-1}}+\frac{1}{2}=1$, by Hölder's inequality and Sobolev's embedding theorem, and recalling (2.46), we get

$$
\begin{align*}
\int_{\Omega} a z z_{t} d x & \leq \int_{\Omega}|a||z|^{\frac{n}{p}}|z|^{1-\frac{n}{p}}\left|z_{t}\right| d x \\
& \leq r_{0}| ||z(t, \cdot)|^{\frac{n}{p}}\left\|_{L^{p_{1}}(\Omega)}\right\||z(t, \cdot)|^{1-\frac{n}{p}}\left\|_{L^{p_{2}}(\Omega)}\right\| z_{t}(t, \cdot) \|_{L^{2}(\Omega)} \\
& =r_{0}\|z(t, \cdot)\|_{L^{\frac{2 n}{p}}}\|z(t, \cdot)\|_{L^{2}(\Omega)}^{1-\frac{n}{p}}\left\|z_{t}(t, \cdot)\right\|_{L^{2}(\Omega)}  \tag{2.48}\\
& =r_{0}^{\frac{1}{2-n / p}} \underbrace{\|z(t, \cdot)\|_{L^{\frac{n}{p}}}^{\frac{2 n}{n-2}}(\Omega)}_{\leq \mathcal{E}(t)^{\frac{n}{2 p}}} \underbrace{\left(r_{0}^{\frac{1-n / p}{2-n / p}}\|z(t, \cdot)\|_{L^{2}(\Omega)}^{1-\frac{n}{p}}\right)}_{\leq \mathcal{E}(t)^{\frac{1}{2}-\frac{n}{2 p}}} \underbrace{\left\|z_{t}(t, \cdot)\right\|_{L^{2}(\Omega)}}_{\leq \mathcal{E}(t)^{1 / 2}} \\
& \leq C r_{0}^{\frac{1}{2-n / p}} \mathcal{E}(t) .
\end{align*}
$$

Similarly,

$$
\begin{equation*}
r_{0}^{\frac{2}{2-n / p}} \int_{\Omega} z z_{t} d x \leq \frac{r_{0}^{\frac{1}{2-n / p}}}{2} \int_{\Omega}\left(r_{0}^{\frac{2}{2-n / p}} z^{2}+z_{t}^{2}\right) d x \leq C r_{0}^{\frac{1}{2-n / p}} \mathcal{E}(t) \tag{2.49}
\end{equation*}
$$

Hence, combining (2.47)-(2.49), we conclude that

$$
\frac{d \mathcal{E}(t)}{d t} \leq C r_{0}^{\frac{1}{2-n / p}} \mathcal{E}(t)
$$

By this and noting the time reversibility of system (2.45), we get

$$
\mathcal{E}(t) \leq C e^{C r_{0}^{\frac{1}{2-n / p}} \mathcal{E}(s), \quad \forall t, s \in[0, T] . . . . ~}
$$

Hence,

Now, taking the scalar product of the first equation of (2.6) by $z$, integrating it in $(t, s) \times \Omega$, recalling that by assumption $a_{1}^{k} \equiv 0(k=1, \cdots, n)$ and $a_{2} \equiv 0$, and using (2.45) and integrations by parts, we get

$$
\begin{align*}
\left(w(s), z_{t}(s)\right)_{L^{2}(\Omega)}+\langle & \left.w_{t}(s),-z(s)\right\rangle_{H^{-1}(\Omega), H_{0}^{1}(\Omega)}  \tag{2.51}\\
& =\left(w(t), z_{t}(t)\right)_{L^{2}(\Omega)}+\left\langle w_{t}(t),-z(t)\right\rangle_{H^{-1}(\Omega), H_{0}^{1}(\Omega)}, \quad \forall t, s \in[0, T]
\end{align*}
$$

Hence, by (2.42) and (2.51), and noting the last equation in (2.45), and using (2.50), we get (denoting by $S$ the unit sphere of the space $H_{0}^{1}(\Omega) \times L^{2}(\Omega)$ )

$$
\begin{aligned}
\sqrt{2 E(T)} & =\sup _{\left(z^{0}, z^{1}\right) \in S}\left[\left(w(T), z_{1}\right)_{L^{2}(\Omega)}+\left\langle w_{t}(T),-z_{0}\right\rangle_{H^{-1}(\Omega), H_{0}^{1}(\Omega)}\right] \\
& =\sup _{\left(z^{0}, z^{1}\right) \in S}\left[\left(w(t), z_{t}(t)\right)_{L^{2}(\Omega)}+\left\langle w_{t}(t),-z(t)\right\rangle_{H^{-1}(\Omega), H_{0}^{1}(\Omega)}\right] \\
& \leq C \sqrt{E(t)} \sup _{\left(z^{0}, z^{1}\right) \in S}\left\|\left(z(t), z_{t}(t)\right)\right\|_{H_{0}^{1}(\Omega) \times L^{2}(\Omega)} \\
& \leq C e^{C r_{0}^{\frac{1}{2-n / p}} \sqrt{E(t)} \sup _{\left(z^{0}, z^{1}\right) \in S}\left\|\left(z(T), z_{t}(T)\right)\right\|_{H_{0}^{1}(\Omega) \times L^{2}(\Omega)}} \\
& =C e^{C r_{0}^{\frac{1}{2-n / p}} \sqrt{E(t)}} .
\end{aligned}
$$

This fact, combined with the time reversibility of system (2.6), yields the desired estimate (2.44).
Step 4: We now return to the proof of the second assertion. By (2.44), we get

$$
\begin{equation*}
\|w\|_{L^{2}(J \times \Omega)}^{2} \leq C E(0) e^{C r_{0}^{\frac{1}{2-n / p}}} \tag{2.52}
\end{equation*}
$$

and

$$
\begin{equation*}
\int_{S_{0}}^{S_{0}^{\prime}} E(t) d t \geq \frac{1}{C} E(0) e^{-C r_{0}^{\frac{1}{2-n / p}}} \tag{2.53}
\end{equation*}
$$

Combining (2.53) with (2.41) and (2.43), we get:

$$
\begin{equation*}
\lambda \int_{Q} w^{2} d x d t \geq \frac{\lambda}{C\left(1+r_{0}\right)} e^{R_{0}^{2} \lambda-C r_{0}^{\frac{1}{2-n / p}}} E(0) \tag{2.54}
\end{equation*}
$$

Inequality (2.34) together with (2.52) and (2.54) yields a constant $C_{2}$ such that

$$
\begin{align*}
& \lambda \geq\left(1+r_{0}^{\frac{1}{3 / 2-n / p}}\right) \lambda_{1} \Longrightarrow  \tag{2.55}\\
& \underbrace{\left[\lambda e^{R_{0}^{2} \lambda-C_{2} r_{0}^{\frac{1}{2-n / p}}}-C_{2}\left(1+r_{0}\right) e^{\left.C_{2} r_{0}^{\frac{1}{2-n / p}}\right]}\right.}_{\alpha\left(\lambda, r_{0}\right)} E(0) \leq C_{2} \lambda^{2}\left(1+r_{0}\right) e^{C_{2} \lambda} \int_{0}^{T} \int_{\omega} w^{2} d x d t .
\end{align*}
$$

Assume that $\lambda \geq\left(1+r_{0}^{\frac{1}{3 / 2-n / p}}\right) \lambda_{1}$. Taking, if necessary, a greater $\lambda_{1}$ we have:

$$
\begin{equation*}
\lambda e^{\frac{R_{0}^{2} \lambda}{2}} \geq 1+C_{2}\left(1+r_{0}\right), \quad \frac{R_{0}^{2} \lambda}{2} \geq 2 C_{2} r_{0}^{\frac{1}{2-n / p}} \tag{2.56}
\end{equation*}
$$

(To obtain the second inequality we used that $\frac{1}{2-n / p}<\frac{1}{3 / 2-n / p}$ ). Thus:

$$
\begin{equation*}
\alpha\left(\lambda, r_{0}\right) \geq e^{C_{2} r_{0}^{\frac{1}{2-n / p}}} \geq 1 \tag{2.57}
\end{equation*}
$$

Then, from (2.55) and (2.57), we obtain:

$$
\begin{equation*}
\exists \lambda_{1}, \quad \lambda \geq\left(1+r_{0}^{\frac{1}{3 / 2-n / p}}\right) \lambda_{1} \Longrightarrow E(0) \leq C_{2} \lambda^{2}\left(1+r_{0}\right) e^{C_{2} \lambda}\|w\|_{L^{2}((0, T) \times \omega)}^{2} \tag{2.58}
\end{equation*}
$$

taking the preceding inequality at $\lambda=\left(1+r_{0}^{\frac{1}{3 / 2-n / p}}\right) \lambda_{1}$ gives the desired observability inequality (2.15).

Step 5: Let us now show the first assertion. Again we will show an intermediate inequality:

$$
\begin{align*}
\exists \lambda_{1}>0: \forall \lambda & >\lambda_{1}\left(1+r_{0}^{\frac{1}{3 / 2-n / p}}+r_{2}^{2}\right),  \tag{2.59}\\
& \lambda \int_{Q} e^{2 \lambda \phi}\left(w_{t}^{2}+|\nabla w|^{2}\right) d x d t \leq C\left[\|w\|_{H^{1}(J \times \Omega)}^{2}+\lambda \int_{0}^{T} \int_{\Gamma_{0}} e^{2 \lambda \phi}\left|\frac{\partial w}{\partial \nu}\right|^{2} d x d t\right] .
\end{align*}
$$

Clearly, $\xi w$ vanishes at $t=0, T$. Therefore, by Theorem 2.3, we get

$$
\begin{align*}
& \left.\lambda \int_{Q} e^{2 \lambda \phi}\left[\lambda^{2}(\xi w)^{2}+\mid(\xi w)\right)^{2}+|\nabla(\xi w)|^{2}\right] d x d t  \tag{2.60}\\
& \qquad \quad \leq C\left(\int_{Q} e^{2 \lambda \phi}|\mathcal{P}(\xi w)|^{2} d x d t+\lambda \int_{0}^{T} \int_{\Gamma_{0}} e^{2 \lambda \phi}\left|\frac{\partial w}{\partial \nu}\right|^{2} d x d t\right), \quad \forall \lambda \geq \lambda_{0}
\end{align*}
$$

By equation (2.6), we get:

$$
\begin{equation*}
\int_{Q} e^{2 \lambda \phi}|\mathcal{P}(\xi w)|^{2} d x d t=\int_{Q} e^{2 \lambda \phi}\left|\xi_{t t} w+2 \xi_{t} w_{t}+\xi\left(a w+\sum_{k=1}^{n} a_{1}^{k} w_{x_{k}}+a_{2} w_{t}\right)\right|^{2} d x d t \tag{2.61}
\end{equation*}
$$

Furthemore, recalling the definition of $r_{0}$ in (2.3), and using successively Hölder's and Sobolev's inequalities, then inequality (2.27):

$$
\begin{aligned}
\left\|a e^{\lambda \phi} \xi w\right\|_{L^{2}(Q)} & \leq r_{0}\left\|e^{\lambda \phi} \xi w\right\|_{L^{2}\left(0, T ; L^{s}(\Omega)\right)}, \quad 1 / s+1 / p=1 / 2 \\
& \leq r_{0}\left\|e^{\lambda \phi} \xi w\right\|_{L^{2}\left(0, T ; H_{0}^{n / p}(\Omega)\right)} \\
& \leq r_{0}\left\|e^{\lambda \phi} \xi w\right\|_{L^{2}\left(0, T ; H_{0}^{1}(\Omega)\right)}^{n / p}\left\|e^{\lambda \phi} \xi w\right\|_{L^{2}(Q)}^{1-n / p} .
\end{aligned}
$$

Hence, using Young's inequality we get, for any $\varepsilon>0$ :

$$
\begin{equation*}
\left\|a e^{\lambda \phi} \xi w\right\|_{L^{2}(Q)}^{2} \leq \varepsilon \lambda\left\|e^{\lambda \phi} \xi w\right\|_{L^{2}\left(0, T ; H_{0}^{1}(\Omega)\right)}^{2}+C_{\varepsilon} r_{0}^{2 p /(p-n)} \lambda^{-n /(p-n)}\left\|e^{\lambda \phi} \xi w\right\|_{L^{2}(Q)}^{2}, \tag{2.62}
\end{equation*}
$$

where $C_{\varepsilon}$ is a positive constant depending on $\varepsilon$. Further, by (2.31) and (2.33), we have

$$
\begin{align*}
& \int_{Q} e^{2 \lambda \phi}\left[\left|(\xi w)_{t}\right|^{2}+|\nabla(\xi w)|^{2}\right] d x d t  \tag{2.63}\\
& \quad \geq \int_{Q} e^{2 \lambda \phi}\left[w_{t}^{2}+|\nabla w|^{2}\right] d x d t-C e^{\left(R_{1}^{2}-c T^{2} / 4\right) \lambda}\left(\|w\|_{H^{1}\left(\left(0, T_{1}\right) \times \Omega\right)}^{2}+\|w\|_{H^{1}\left(\left(T_{1}^{\prime}, T\right) \times \Omega\right)}^{2}\right) .
\end{align*}
$$

Combining (2.60)-(2.63), and taking $\varepsilon>0$ small enough in (2.62) we get, for some large constant $C_{3}>0$ :

$$
\begin{align*}
& \lambda \int_{Q} e^{2 \lambda \phi}\left[\lambda^{2}(\xi w)^{2}+w_{t}^{2}+|\nabla w|^{2}\right] d x d t \\
& \leq C_{3}\left[e^{\left(R_{1}^{2}-c T^{2} / 4\right) \lambda}\|w\|_{H^{1}(J \times \Omega)}^{2}+r_{0}^{2 p /(p-n)} \lambda^{-n /(p-n)}| | e^{\lambda \phi} \xi w \|_{L^{2}(Q)}^{2}\right.  \tag{2.64}\\
& \left.\quad+r_{2}^{2} \int_{Q} e^{2 \lambda \phi}\left(w_{t}^{2}+|\nabla w|^{2}\right) d x d t+\lambda \int_{0}^{T} \int_{\Gamma_{0}} e^{2 \lambda \phi}\left|\frac{\partial w}{\partial \nu}\right|^{2} d x d t\right] .
\end{align*}
$$

Now, choosing $\lambda_{1}>0$ large enough such that:

$$
\lambda>\lambda_{1}\left(1+r_{0}^{\frac{1}{3 / 2-n / p}}+r_{2}^{2}\right) \Longrightarrow C_{3}\left(r_{0}^{\frac{2 p}{p-n}}+r_{2}^{2}\right) \leq \frac{\lambda}{2},
$$

and noting that by (2.17), $R_{1}^{2}-c T^{2} / 4<0$, we get (2.59).
As in Step 3, (2.59) and the modified energy method lead to the first assertion in Theorem 2.2. Since the preceding case is more complex we omit the details. This completes the proof of Theorem 2.2.

## 3 Super-exponentially decaying solutions to elliptic equations

In this section, we construct solutions $u$ of equations on $\mathbb{R}^{n}, n \geq 2$ of one of the following forms:

$$
\begin{equation*}
\Delta u+q u=0 \tag{3.1}
\end{equation*}
$$

or:

$$
\begin{equation*}
\Delta u+q_{1} \cdot \nabla u=0, \tag{3.2}
\end{equation*}
$$

where $q$ and $q_{1}$ admit some bounds at infinity, and such that $u$ decays at infinity as fast as possible. When $n=1$, if $q$ is bounded, the solution $u$ of (3.1) may not decay faster than exponentially. In higher dimensions, an elementary but optimal Carleman estimate due to Meshkov [19], shows that one cannot hope a decay at infinity faster than $e^{-C|x|^{4 / 3}}$ for all $C>0$. In the same work is also given a surprising example on $\mathbb{R}^{2}$ of complex-valued fonctions $q$ and $u$ satisfying (3.1), with $q$ bounded and $u$ decaying like $e^{-|x|^{4 / 3}}$ (see Theorem B). The example of Meshkov gives by separation of variables examples of solutions with the optimal rate of decay in any even dimension. In Theorem
3.1 we state this result and its counterpart in dimension 2 in the case of equation (3.2), where the convective potential $q_{1}$ decays like $|x|^{-1 / 3}$ (which is sharp according to the Carleman estimate in [19]). For the sake of completness we give a proof of both results.

The case of odd dimension is more complex and to our knowledge no results exist in this direction. In Theorem 3.2, we give, for odd $n$, an example of a solution $u$ of (3.1) on $\mathbb{R}^{n}$ decaying at the same super-exponential speed, but which is $\mathbb{C}^{4}$-valued, and with some $q$ growing logarithmically at infinity. Once again, according to the Carleman inequality, this construction is quasi-optimal. The existence of a solution $u$ of (3.1) taking values in $\mathbb{C}$ and/or with a bounded $q$ remains open, as are similar questions concerning real-valued functions in all dimensions $n \geq 2$.

For $x$ in $\mathbb{R}^{n}$, we shall write $r=|x|$.
Theorem 3.1. Let $n \geq 2$ be even and $c_{*}>0$. There exists nontrivial functions:

$$
u \in C^{\infty}\left(\mathbb{R}^{n} ; \mathbb{C}\right), \quad q \in C^{\infty}\left(\mathbb{R}^{n} ; \mathbb{C}\right) \cap L^{\infty}\left(\mathbb{R}^{n} ; \mathbb{C}\right)
$$

such that (3.1) is satisfied on $\mathbb{R}^{n}$, and, for some constant $C$ :

$$
\begin{equation*}
|u(x)| \leq C e^{-c_{*} r^{4 / 3}} . \tag{3.3}
\end{equation*}
$$

Furthermore, when $n=2$, and for the same function $u$, there exists:

$$
\begin{equation*}
q_{1} \in C^{\infty}\left(\mathbb{R}^{2} ; \mathbb{C}^{2}\right), \quad \text { with }(r+1)^{1 / 3} q_{1} \in L^{\infty}\left(\mathbb{R}^{2} ; \mathbb{C}^{2}\right) \tag{3.4}
\end{equation*}
$$

such that equation (3.2) holds.
Theorem 3.2. Let $n \geq 3$ be odd and $c_{*}>0$. There exist nontrivial functions:

$$
u \in C^{\infty}\left(\mathbb{R}^{n} ; \mathbb{C}^{4}\right), \quad q \in C^{\infty}\left(\mathbb{R}^{n} ; \mathbb{C}^{4 \times 4}\right)
$$

fullfilling (3.1) and such that, for a constant $C>0$ :

$$
\begin{gather*}
(\log (2+r))^{-3} q \in L^{\infty}\left(\mathbb{R}^{n}\right)  \tag{3.5}\\
|u(x)| \leq C e^{-c_{*} r^{4 / 3}} \tag{3.6}
\end{gather*}
$$

Remark 3.1. Theorem 3.1 is optimal in the following sense. Fix a bounded function $q$ on $\mathbb{R}^{n}$. Then the only solution of equation (3.1) on $\mathbb{R}^{n}$ satisfying:

$$
\begin{equation*}
\forall c>0, \quad \int_{\mathbb{R}^{n}}|u(x)| e^{c|x|^{4 / 3}} d x<+\infty, \tag{3.7}
\end{equation*}
$$

is $u=0$. Indeed this is a consequence of the following Carleman inequality by Meshkov [19, Lemma 1 ], which holds for large $\tau$, in any space dimension $n$ :
(3.8) $\forall v \in C_{0}^{\infty}(\{r>1\})$,

$$
\tau^{3} \int|v|^{2} \exp \left(2 \tau r^{4 / 3}\right) r^{1-n} d x+\tau \int|\nabla v|^{2} \exp \left(2 \tau r^{4 / 3}\right) r^{1 / 3-n} d x \leq C \int|\Delta v|^{2} \exp \left(2 \tau r^{4 / 3}\right) r^{1-n} d x
$$

The same argument shows that if $q_{1}$ satisfies the decay property (3.4), then the only solution $u$ of (3.2) on $\mathbb{R}^{n}$ satisfying (3.7) is $u=0$. These results remain valid with vector-valued functions.

If $q$ (scalar or matrix-valued) satisfies a logarithmic bound such as:

$$
\exists N>0, \quad \frac{q}{\log ^{N}(2+r)} \in L^{\infty}\left(\mathbb{R}^{n}\right)
$$

one may write weaker uniqueness results for (3.2). For example:

$$
\int_{\mathbb{R}^{n}}|u(x)| e^{\varepsilon|x|^{\varepsilon+4 / 3}} d x<+\infty \quad \text { for some } \varepsilon>0 \Longrightarrow u=0
$$

In this sense Theorem 3.2 is quasi-optimal.
Remark 3.2. The construction of Theorem 3.1 may be adapted to potentials with polynomial bounds at infinity. Precisely, if $2 / 3<\alpha \leq 2$ and $\gamma \triangleq(4-3 \alpha) / 2$, there exist a potential $q_{\alpha}$ bounded by $C|x|^{-\gamma}$ and a solution $u$ decaying like $e^{-|x|^{\alpha}}$ such that equation (3.1) holds. By a variant of Carleman inequality (3.8), this decay is also optimal.

One may think this result would be of some help to test the optimality of the constant of the observability inequality for potentials $a \in L^{\infty}\left(0, T ; L^{p}(\Omega)\right), n \leq p<+\infty$ (see Theorems 2.1 and 2.2). Unfortunately this is not the case. Indeed the arguments of Sections 4 and 5 applied on the potentials $q_{\alpha}$ would only show that the constant of observation may not be better than $C e^{C\|a\|_{L^{\infty}\left(0, T ; L^{p}(\Omega)\right)}^{2 / 3}}$, which is only interesting in the case $p=\infty$, that is in the case of Theorem 3.1.

All the polynomially decaying potentials $q_{\alpha}$ we are refering to are locally bounded. It seems likely that, to prove the optimality of the observability constants given by Theorems 2.1 and 2.2 for potentials $a$ in $L^{p}$ spaces, one will have to construct variants of the Meshkov function $u$ with potentials $q$ that are in $L^{p}$, but not locally bounded (see also open problem 8.4).

Remark 3.3. The uniqueness result of Meshkov, and the constructions of this section are closely related to the issue of uniqueness of solutions of equations such as (3.1) or (3.2) vanishing to some specified (finite or infinite) order at a point or at a general submanifold of $\mathbb{R}^{n}$. See the book of Zuily [32] and the constructions of counterexamples in [2], [3], [15], [27].

### 3.1 Construction in even dimension

In this part we prove Theorem 3.1. We first remark that in the case of equation (3.1), we may assume that $n=2$. Indeed, if there exists a function $u$ and a potential $q$ fullfilling the first part of Theorem 3.1 for $n=2$, and if $n=2 m$ is even, we can define

$$
\tilde{u}\left(x_{1}, x_{2}, \ldots, x_{2 m}\right) \triangleq u\left(x_{1}, x_{2}\right) u\left(x_{3}, x_{4}\right) \cdots u\left(x_{2 m-1}, x_{2 m}\right),
$$

It then follows that

$$
\Delta \tilde{u}+\underbrace{\left(q\left(x_{1}, x_{2}\right)+q\left(x_{3}, x_{4}\right)+\ldots q\left(x_{2 m-1}, x_{2 m}\right)\right)}_{\tilde{q}} \tilde{u}=0 .
$$

The potential $\tilde{q}$ is bounded on $\mathbb{R}^{n}$ and $\tilde{u}$ satisfies the decay property:

$$
|\tilde{u}(x)| \leq C \exp \left(-c_{*} r^{4 / 3}\right), \quad x \in \mathbb{R}^{n} .
$$

Note that this argument does not work in the case of convective equation (3.2), as it gives only a bounded $q_{1}$.

We now assume that $n=2$. Denote by $\theta=\frac{x}{|x|}$ the polar variable. We will construct a function $u$ which is harmonic in the neighborhood of any of its zero. Thus the existence of a function $q$, bounded and $C^{\infty}$ on $\mathbb{R}^{2}$, and such that equation (3.1) holds is equivalent to the existence of a constant $C_{0}$ such that, on $\mathbb{R}^{2}$ :

$$
\begin{equation*}
|\Delta u| \leq C_{0}|u| . \tag{3.9}
\end{equation*}
$$

Indeed, if $u$ satisfies (3.9), it suffices to take $q$ to be $-\Delta u / u$ where $u$ does not vanish, and 0 elsewhere, which implies trivially equation (3.1). Likewise, the existence of a function $w$ satisfying (3.4) and such that equation (3.2) holds is equivalent to the existence of a constant $C_{0}^{\prime}$ such that, on $\mathbb{R}^{2}$ :

$$
\begin{equation*}
|\Delta u| \leq C_{0}^{\prime}(r+1)^{-1 / 3}|\nabla u| . \tag{3.10}
\end{equation*}
$$

Furthermore, it suffices to show Theorem 3.1 for some $c_{*}$. One can then obtain the general case by dilatation of $u$.

### 3.1.1 Construction for large $r$

We first construct $u$ on $\{r \geq \rho\}$, where $\rho$ is large. As in the article of Meshkov [19], we shall construct $u$ on well-chosen rings $\left\{\rho_{k} \leq r \leq \rho_{k+1}\right\}, k \in \mathbb{N}$.
Notation 3.1. In all the following, we shall write, for sequences of real numbers $\left(A_{k}\right)$ and $\left(B_{k}\right)$ :

$$
A_{k}=O\left(B_{k}\right),
$$

when there exist constants $C$ and $k_{0}$ such that:

$$
\forall k \geq k_{0}, \quad\left|A_{k}\right| \leq C\left|B_{k}\right|
$$

When $A_{k}$ and $B_{k}$ also depend on $r \in I, I$ being an interval, the estimate is also assumed to be uniform with respect to $r$ in $I$.

We will also use the notation $\approx$ in the following sense:

$$
A_{k} \approx B_{k} \Longleftrightarrow\left(A_{k}=O\left(B_{k}\right) \text { and } B_{k}=O\left(A_{k}\right)\right)
$$

Let $\rho_{0}$ be a large enough real number, and define the sequence $\rho_{k}$ by:

$$
\begin{equation*}
\rho_{k+1} \triangleq \rho_{k}+6 \rho_{k}^{1 / 3} \tag{3.11}
\end{equation*}
$$

Denote by $\rho_{k l} \triangleq \rho_{k}+l \rho_{k}^{1 / 3}$, for $l=0, \cdots, 6$. This divides the interval $\left[\rho_{k}, \rho_{k+1}\right]$ in 6 sub-intervals. Consider the harmonic function:

$$
\begin{equation*}
u_{k} \triangleq a_{k} r^{-n_{k}} e^{i n_{k} \theta}, \quad a_{k}>0, n_{k} \in \mathbb{N} \tag{3.12}
\end{equation*}
$$

The crucial point of the proof is a lemma which allows to pass from $u_{k}$ to $-\bar{u}_{k+1}$ within the interval ( $\rho_{k}, \rho_{k+1}$ ) with a function solving (3.9). We first choose the values of $n_{k}$ and $a_{k}$. Let:

$$
\begin{gather*}
n_{k} \triangleq 2\left[\frac{\rho_{k}^{4 / 3}}{2}\right], \quad d_{k} \triangleq \frac{n_{k+1}-n_{k}}{2} \in \mathbb{N}  \tag{3.13}\\
a_{0} \triangleq 1, \quad a_{k+1} \triangleq a_{k} \rho_{k 3}^{2 d_{k}}
\end{gather*}
$$

where $[y]$ stands for the integer part of $y$. The $a_{k}$ 's have been chosen so that $\left|u_{k}\right|$ and $\left|u_{k+1}\right|$ coincide when $r=\rho_{k 3}$. An easy calculation shows that for some positive constant $\delta$, independent of $k$ :

$$
\begin{equation*}
d_{k}=\delta \rho_{k}^{2 / 3}+O(1) \tag{3.14}
\end{equation*}
$$

Lemma 3.1. There exist a large integer $k_{0}$, a constant $C_{0}$ independent of $k \geq k_{0}$, and:

$$
u \in C^{\infty}\left(\left\{x \in \mathbb{R}^{2}\left|\rho_{k} \leq|x| \leq \rho_{k+1}\right\}\right)\right.
$$

such that:

$$
\begin{gather*}
u(r, \theta)=u_{k}(r, \theta), \quad \rho_{k}=\rho_{k 0} \leq r \leq \rho_{k 1}  \tag{3.15}\\
u(r, \theta)=-\bar{u}_{k+1}(r, \theta), \quad \rho_{k 5} \leq r \leq \rho_{k 6}=\rho_{k+1} \\
|u(r, \theta)|=O\left(a_{k} r^{-n_{k}}\right), \quad \rho_{k} \leq r \leq \rho_{k+1} \tag{3.16}
\end{gather*}
$$

and, for $\rho_{k} \leq r \leq \rho_{k+1}$ :

$$
\begin{gather*}
|\Delta u| \leq C_{0}|u|,  \tag{3.17}\\
|\Delta u| \leq C_{0} r^{-1 / 3}|\nabla u| . \tag{3.18}
\end{gather*}
$$

The result remains valid when replacing (3.15) by:

$$
\begin{align*}
& u(r, \theta)=-\bar{u}_{k}(r, \theta), \quad \rho_{k 0} \leq r \leq \rho_{k 1} \\
& u(r, \theta)=u_{k+1}(r, \theta), \quad \rho_{k 5} \leq r \leq \rho_{k 6}
\end{align*}
$$

Proof. We will just do the proof in the first case (when (3.15) holds), the proof of the other case being the same almost word by word.

A simple idea would be to let $u=\xi_{k} u_{k}-\tilde{\xi}_{k} \bar{u}_{k+1}$, where $\xi_{k}$ and $\tilde{\xi}_{k}$ are smooth functions of $r$, $\xi_{k}$ (resp. $\tilde{\xi}_{k}$ ) is equal to 1 (resp. 0) near $\rho_{k}, 0$ (resp. 1) near $\rho_{k+1}$. Indeed, with a suitable choice of the functions $\xi_{k}$ and $\tilde{\xi}_{k}$, such a $u$ solves (3.17) except in the neighborhood of its zeros. This is a non-negligible difficulty, taking into account that, for homotopy reasons, a continuous function $u$ satisfying (3.15) has to vanish: when $r$ increases, $u$ passes continuously from a continuous function of $\theta$ in $S^{1}$ winding clockwise $n_{k}$ times around the origin to one winding counterclockwise $n_{k+1}$ times, which is impossible without vanishing. Note that $n_{k} \neq n_{k+1}$, so that the same argument would hold with $u_{k+1}$ (which winds clockwise $n_{k+1}$ times around the origin) instead of $-\bar{u}_{k+1}$ in (3.15).

One way to avoid this problem is to consider a $\mathbb{C}^{2}$-valued function (see Remark 3.4). To treat the harder case of a complex-valued function, we need to use a trick due to Meshkov [19] consisting in introducing an intermediate function $v_{k}$, close to $-\bar{u}_{k+1}$, which is non-harmonic, but nevertheless solution of an inequality of the form (3.17) on $\rho_{k} \leq r \leq \rho_{k+1}$.

Choice of an intermediate state between $u_{k}$ and $-\bar{u}_{k+1}$. Consider a $2 \pi$-periodic function $\varphi_{k}(\theta)$ (which we shall make explicit later) such that:

$$
\begin{equation*}
\left|\varphi_{k}(\theta)\right|=O\left(\rho_{k}^{-2 / 3}\right), \quad\left|\varphi_{k}^{\prime}(\theta)\right|=O\left(\rho_{k}^{2 / 3}\right), \quad\left|\varphi_{k}^{\prime \prime}(\theta)\right|=O\left(\rho_{k}^{2}\right) \tag{3.19}
\end{equation*}
$$

Let:

$$
\begin{equation*}
v_{k}(r, \theta) \triangleq-r^{4 d_{k}} \rho_{k 3}^{-4 d_{k}} e^{i \varphi_{k}(\theta)} \bar{u}_{k+1}(r, \theta) \tag{3.20}
\end{equation*}
$$

As $-\bar{u}_{k+1}(r, \cdot)$, the function $v_{k}(r, \cdot)$, when $r$ is fixed, winds $n_{k+1}$ times, counterclockwise, around the origin. But unlike $-\bar{u}_{k+1}$, it decreases slower than $u_{k}$, so that it is more natural to replace $u_{k}$ by $v_{k}$ than by $-\bar{u}_{k+1}$ as $r$ is increasing. For these reasons, $v_{k}$ is an appropriate intermediate state between $u_{k}$ and $-\bar{u}_{k+1}$. The constant in (3.20) has been chosen so that:

$$
\begin{equation*}
\forall \theta, \quad\left|u_{k}\left(\rho_{k 3}, \theta\right)\right|=\left|u_{k+1}\left(\rho_{k 3}, \theta\right)\right|=\left|v_{k}\left(\rho_{k 3}, \theta\right)\right| . \tag{3.21}
\end{equation*}
$$

Let $g(r) \triangleq \log \left(r^{d_{k}} / \rho_{k}^{d_{k}}\right)$. By (3.14), we have:

$$
g^{\prime}(r)=\frac{d_{k}}{r}=O\left(\rho_{k}^{-1 / 3}\right) .
$$

Noting that $g\left(\rho_{k}\right)=0$ and $\rho_{k+1}-\rho_{k}=O\left(\rho_{k}^{1 / 3}\right)$, we get that $g(r)=O(1)$ for $r$ in the interval [ $\rho_{k}, \rho_{k+1}$ ]. Thus:

$$
\begin{equation*}
r^{d_{k}} \approx \rho_{k}^{d_{k}}, \quad \rho_{k} \leq r \leq \rho_{k+1} . \tag{3.22}
\end{equation*}
$$

The following lemma gathers some estimates on $u_{k}, v_{k}$ and $u_{k+1}$.

## Lemma 3.2. We have:

$$
\begin{equation*}
u_{k} \approx v_{k} \approx u_{k+1}, \quad \rho_{k} \leq r \leq \rho_{k+1} . \tag{3.23}
\end{equation*}
$$

Furthermore, any of the three sequences $w_{k}=u_{k}$, $v_{k}$ or $u_{k+1}$, satisfies, for $\rho_{k} \leq r \leq \rho_{k+1}$ :

$$
\begin{equation*}
\left|\partial_{r} w_{k}\right|=\left(r^{1 / 3}+O\left(r^{-1 / 3}\right)\right)\left|w_{k}\right|, \quad\left|\frac{1}{r} \partial_{\theta} w_{k}\right|=\left(r^{1 / 3}+O\left(r^{-1 / 3}\right)\right)\left|w_{k}\right| \tag{3.24}
\end{equation*}
$$

Finally there is a constant $C_{1}$ such that:

$$
\begin{gather*}
\left|\Delta v_{k}\right| \leq C_{1}\left|v_{k}\right|, \quad \rho_{k} \leq r \leq \rho_{k+1}  \tag{3.25}\\
\left|\Delta v_{k}\right| \leq C_{1} r^{-1 / 3}\left|\nabla v_{k}\right|, \quad \rho_{k} \leq r \leq \rho_{k+1} \tag{3.26}
\end{gather*}
$$

Proof. The estimates (3.23) are a direct consequence of (3.22) and the definitions of $u_{k}$ and $v_{k}$. By simple computations:

$$
\begin{gathered}
\partial_{r} u_{k}=-\frac{n_{k}}{r} u_{k}, \quad \partial_{r} v_{k}=\frac{2 d_{k}-n_{k}}{r} v_{k}, \\
\frac{1}{r} \partial_{\theta} u_{k}=\frac{i n_{k}}{r} u_{k}, \quad \frac{1}{r} \partial_{\theta} v_{k}=\frac{-i n_{k+1}+i \varphi_{k}^{\prime}(\theta)}{r} v_{k} .
\end{gathered}
$$

With the definition (3.13) of $n_{k}$, and the bounds (3.19) on $\varphi_{k}$, we easily get (3.24). Furthermore:

$$
\begin{equation*}
\Delta v_{k}=-2 \rho_{k 3}^{-4 d_{k}} \nabla\left(r^{4 d_{k}} e^{i \varphi_{k}}\right) \cdot \nabla \bar{u}_{k+1}-\rho_{k 3}^{-4 d_{k}} \Delta\left(r^{4 d_{k}} e^{i \varphi_{k}}\right) \bar{u}_{k+1} . \tag{3.27}
\end{equation*}
$$

Keeping in mind the estimates (3.14) and (3.22), we get:

$$
\nabla\left(r^{4 d_{k}}\right) \approx r^{4 d_{k}} \rho_{k}^{-1 / 3}, \quad \Delta\left(r^{4 d_{k}}\right)=O\left(\rho_{k}^{-2 / 3} r^{4 d_{k}}\right)
$$

The bounds (3.19) on $\varphi_{k}$ and its derivatives imply easily:

$$
\nabla e^{i \varphi_{k}}=O\left(\rho_{k}^{-1 / 3}\right), \quad \Delta\left(e^{i \varphi_{k}}\right)=O(1)
$$

Combining these estimates, one gets:

$$
\left|\nabla\left(r^{4 d_{k}} e^{i \varphi_{k}}\right)\right|=O\left(\rho_{k}^{-1 / 3} r^{4 d_{k}}\right), \quad \Delta\left(r^{4 d_{k}} e^{i \varphi_{k}}\right)=O\left(r^{4 d_{k}}\right) .
$$

Together with the estimates (3.23), (3.24), equation (3.27) yields (3.25) and (3.26).

Figure 1: Steps of the construction


The construction of $u$ takes three steps (see Figure 1).
Step 1. Construction of $u$ on $\left[\rho_{k 0}, \rho_{k 4}\right]$.
Lemma 3.3. There exists $u, C^{\infty}$ on $\left\{\rho_{k 0} \leq r \leq \rho_{k 4}\right\}$, satisfying (3.17) and (3.18) and such that:

$$
\begin{gather*}
u(r, \theta)=u_{k}(r, \theta), \quad \rho_{k 0} \leq r \leq \rho_{k 1},  \tag{3.28}\\
u(r, \theta)=u_{k}(r, \theta)+v_{k}(r, \theta), \quad \rho_{k 2} \leq r \leq \rho_{k 4} . \tag{3.29}
\end{gather*}
$$

Proof. Let $\chi$ be a non-decreasing, $C^{\infty}$ function such that:

$$
\begin{equation*}
\chi(s)=0 \text { if } s \leq 1 \text { and } \chi(s)=1 \text { if } s \geq 2 . \tag{3.30}
\end{equation*}
$$

Let:

$$
\chi_{k} \triangleq \chi\left(\frac{r-\rho_{k}}{\rho_{k}^{1 / 3}}\right),
$$

so that $\chi_{k}(r)$ is 0 if $r \leq \rho_{k 1}$ and 1 if $r \geq \rho_{k 2}$, and that the derivatives of $\chi_{k}$ satisfy the estimates:

$$
\begin{equation*}
\left|\chi_{k}^{(p)}\right|=O\left(\rho_{k}^{-p / 3}\right) \tag{3.31}
\end{equation*}
$$

Let:

$$
\begin{equation*}
u \triangleq u_{k}+\chi_{k} v_{k} \text { for } \rho_{k 0} \leq r \leq \rho_{k 4} . \tag{3.32}
\end{equation*}
$$

Obviously (3.28) and (3.29) are satisfied. To show that $u$ satisfies (3.17) and (3.18) with $C_{0}$ independent of $k$, we divide $\left[\rho_{k 0}, \rho_{k 4}\right]$ into two subintervals. When $\rho_{k 0} \leq r \leq \rho_{k 2}$, we will simply use that $u_{k}$ satisfy (3.17) and (3.18), that $v_{k}$ satisfy (3.25) and (3.26), and that $\left|u_{k}\right|$ is larger than $c\left|v_{k}\right|, c>1$, so that the sum $u$ of the two is of the order of $u_{k}$. When $\rho_{k 2} \leq r \leq \rho_{k 4}$, the two absolute values may coincide, and we will have to build an adequate phase $\varphi_{k}$ so that the function $u$ is harmonic around its zeros. This is the most tricky and nontrivial part of Meshkov's construction.

The region $\rho_{k 0} \leq r \leq \rho_{k 2}$. Let:

$$
\tilde{g}(r) \triangleq \log \frac{\left|v_{k}\right|}{\left|u_{k}\right|}=\log \frac{a_{k+1} r^{-n_{k+1}+4 d_{k}} \rho_{k 3}^{-4 d_{k}}}{a_{k} r^{-n_{k}}}
$$

Then:

$$
\tilde{g}(r)=2 d_{k} \log r+C(k),
$$

where $C(k)$ is a constant which depends only on $k$. By the choice of the constants $a_{k}$ and $a_{k+1}$ (see (3.21)), $\tilde{g}\left(\rho_{k 3}\right)=0$. Furthermore $\tilde{g}^{\prime}(r)=\frac{2 d_{k}}{r}$. Using (3.14) we get that if $k$ is large enough, $\tilde{g}^{\prime}(r)$ is greater than $\delta \rho_{k}^{-1 / 3}$, from which we deduce the two following crucial comparison estimates:

$$
\begin{align*}
& \left|u_{k}\right| \geq e^{\delta}\left|v_{k}\right|, \quad r \leq \rho_{k 2},  \tag{3.33}\\
& \left|v_{k}\right| \geq e^{\delta}\left|u_{k}\right|, \quad r \geq \rho_{k 4} . \tag{3.34}
\end{align*}
$$

The inequality (3.33) implies that when $\rho_{k 0} \leq r \leq \rho_{k 2}$ :

$$
\begin{equation*}
2\left|u_{k}\right| \geq|u| \geq\left(1-e^{-\delta}\right)\left|u_{k}\right| \geq c\left|u_{k}\right| \geq c\left|v_{k}\right|, \quad \text { where } c>0 \tag{3.35}
\end{equation*}
$$

Furthermore:

$$
\begin{equation*}
\Delta u=\chi_{k} \Delta v_{k}+2 \nabla \chi_{k} \cdot \nabla v_{k}+\left(\Delta \chi_{k}\right) v_{k} . \tag{3.36}
\end{equation*}
$$

According to (3.31), (3.35), (3.36), and the estimates of Lemma 3.2, the function $u$ satisfies (3.17) when $\rho_{k 0} \leq r \leq \rho_{k 2}$.

Furthermore, using again (3.35), and Lemma 3.2, we get:

$$
\begin{aligned}
\partial_{r} u & =\partial_{r} u_{k}+\chi_{k}^{\prime} v_{k}+\chi_{k} \partial_{r} v_{k}, \\
\left|\partial_{r} u\right| & \geq r^{1 / 3}\left|u_{k}\right|+O\left(r^{-1 / 3}\right)\left(\left|u_{k}\right|+\left|v_{k}\right|\right), \\
\left|\partial_{r} u\right| & \geq\left\{\frac{1}{2} r^{1 / 3}+O\left(r^{-1 / 3}\right)\right\}|u|, \quad \rho_{k 0} \leq r \leq \rho_{k 2}
\end{aligned}
$$

which, with inequality (3.17), yields inequality (3.18).
The region $\rho_{k 2} \leq r \leq \rho_{k 4}$. Notice that in this region, $\chi_{k}$ is equal to 1 , so that:

$$
\begin{equation*}
u=u_{k}+v_{k}=a_{k} r^{-n_{k}} e^{i n_{k} \theta}\{1-\underbrace{\rho_{k 3}^{-2 d_{k}} r^{2 d_{k}} e^{-i\left(2 n_{k}+2 d_{k}\right) \theta+i \varphi_{k}(\theta)}}_{w_{k}}\} . \tag{3.37}
\end{equation*}
$$

Let:

$$
\begin{equation*}
T_{k} \triangleq \frac{\pi}{n_{k}+d_{k}}, \quad \theta_{j k} \triangleq j T_{k}, 0 \leq j \leq 2 n_{k}+2 d_{k}-1 \tag{3.38}
\end{equation*}
$$

The $\theta_{j k}$ 's are the solutions of the equation $e^{-i\left(2 n_{k}+2 d_{k}\right) \theta}=1$, so that according to (3.37) ( $\varphi_{k}$ being small), the function $u$ vanishes near each $\theta_{j k}$. In order to satisfy (3.17) and (3.18), $u$ (thus $v_{k}$ ) has to be harmonic near each $\theta_{j k}$. The function $v_{k}$ being equal, up to a multiplicative constant, to

$$
v_{k}=-\rho_{k 3}^{-4 d_{k}} r^{-n_{k}+2 d_{k}} e^{-i n_{k} \theta-2 i d_{k} \theta+i \varphi_{k}(\theta)},
$$

it suffices to choose $\varphi_{k}$ satisfying the following lemma:
Lemma 3.4. There exists a real-valued $\varphi_{k} \in C^{\infty}(\mathbb{R}), 2 \pi$-periodic and satisfying (3.19), such that for all $j$, there is a constant $c_{j k}$ with:

$$
\begin{equation*}
\varphi_{k}(\theta)=4 d_{k} \theta+c_{j k}, \quad \theta_{j k}-\frac{T_{k}}{4} \leq \theta \leq \theta_{j k}+\frac{T_{k}}{4} . \tag{3.39}
\end{equation*}
$$

Figure 2: the function $f_{k}$


Proof. Consider a function $f_{k}$ on $\left[0, T_{k}\right]$ (see Figure 2) so that:

$$
\begin{gather*}
\int_{0}^{T_{k}} f_{k}(s) d s=0  \tag{3.40}\\
f_{k}(s)=4 d_{k}, \quad s \in\left[0, T_{k} / 4\right] \cup\left[3 T_{k} / 4, T_{k}\right],  \tag{3.41}\\
\left|f_{k}(s)\right| \leq C \rho_{k}^{2 / 3}, \quad\left|f_{k}^{\prime}\right| \leq C \rho_{k}^{2}, C \text { independent of } k \tag{3.42}
\end{gather*}
$$

Noting that $d_{k}$ is of the order $\rho_{k}^{2 / 3}$, and $1 / T_{k}$ of the order $\rho_{k}^{4 / 3}$, such a function exists. We extend $f_{k}$ to $\mathbb{R}$ into a $T_{k}$-periodic function, still denoted by $f_{k}$. Let:

$$
\begin{equation*}
\varphi_{k}(\theta) \triangleq \int_{0}^{\theta} f_{k}(s) d s \tag{3.43}
\end{equation*}
$$

which defines, taking into account (3.40), (3.42), and the fact that $T_{k}$ is of the order $\rho_{k}^{-4 / 3}$, a $T_{k}$-periodic function which satisfies the desired bounds (3.19). In particular $\varphi_{k}$ is $2 \pi$-periodic. Furthermore:

$$
\theta_{j k}-T_{k} / 4 \leq \theta \leq \theta_{j k}+T_{k} / 4 \Longrightarrow \varphi_{k}(\theta)=\int_{0}^{\theta_{j k}} f_{k}(s) d s+\int_{\theta_{j k}}^{\theta} f_{k}(s) d s
$$

so that according to (3.40) and (3.41), equality (3.39) holds.
We go back to the proof of Lemma 3.3. To show (3.17), we distinguish two cases. Let $\theta$ be in $[0,2 \pi)$, and choose $j$ so that $\theta=\theta_{j k}+\tau,|\tau| \leq T_{k} / 2$.

- First assume that $|\tau| \geq T_{k} / 4$. Note that $\left(2 n_{k}+2 d_{k}\right) \theta_{j k} \in 2 \pi \mathbb{Z}$. Thus the phase of the second term $w_{k}$ in (3.37) is, by estimates (3.19):

$$
\underbrace{-\left(2 n_{k}+2 d_{k}\right) \theta+\varphi_{k}(\theta)}_{\tilde{\varphi}_{k}(\theta)} \equiv-\left(2 n_{k}+2 d_{k}\right) \tau+O\left(\rho_{k}^{-2 / 3}\right)+2 \pi l_{j k}, \quad l_{j k} \in \mathbb{Z}
$$

Furthermore, depending on the sign of $\tau$ :

$$
-\pi \leq-\left(2 n_{k}+2 d_{k}\right) \tau \leq-\frac{\pi}{2} \quad \text { or } \quad \frac{\pi}{2} \leq-\left(2 n_{k}+2 d_{k}\right) \tau \leq \pi
$$

This implies, for some constant $C$ independent of $k$ and $\theta$ :

$$
\operatorname{Re}\left(e^{i \tilde{\varphi}_{k}}\right) \leq C \rho_{k}^{-2 / 3}
$$

Thus, for large $k$ :

$$
\operatorname{Re}\left(1-\rho_{k 3}^{-2 d_{k}} r^{2 d_{k}} e^{i \tilde{\varphi}_{k}}\right) \geq 1 / 2
$$

Using formula (3.37), we get that if $k$ is large enough:

$$
\begin{equation*}
2|u| \geq a_{k} r^{-n_{k}}=\left|u_{k}\right| . \tag{3.44}
\end{equation*}
$$

With lemma (3.2) and the fact that $\Delta u=\Delta v_{k}$, we get inequality (3.17) with a $C_{0}$ independent of $k$.

By a simple calculation:

$$
\partial_{r} u=-\frac{n_{k}}{r} u+\frac{2 d_{k}}{r} v_{k},
$$

so that, using successively (3.44) and (3.23):

$$
\left|\partial_{r} u\right| \geq \frac{n_{k}}{r}|u|+O\left(r^{-1 / 3}\right)\left|v_{k}\right| \geq\left\{r^{1 / 3}+O\left(r^{-1 / 3}\right)\right\}|u|,
$$

which yields, together with (3.17), inequality (3.18).

- Now assume that $|\tau| \leq T_{k} / 4$. With (3.39), we have:

$$
v_{k}=-\rho_{k 3}^{-4 d_{k}} a_{k} r^{-n_{k}+2 d_{k}} e^{-i\left(n_{k}-2 d_{k}\right) \theta} e^{i c_{j k}},
$$

so that $u$ is harmonic and inequalities (3.17) and (3.18) are trivially satisfied.

Step 2. Construction of $u$ on $\left[\rho_{k 4}, \rho_{k 5}\right]$.
Lemma 3.5. There exists $u$ in $C^{\infty}\left(\left\{\rho_{k 4} \leq r \leq \rho_{k 5}\right\}\right)$ satisfying (3.17) and (3.18), and so that:

$$
\begin{gather*}
u(r, \theta)=u_{k}(r, \theta)+v_{k}(r, \theta), \quad \text { near } \rho_{k 4},  \tag{3.45}\\
u(r, \theta)=v_{k}(r, \theta), \quad \text { near } \rho_{k 5} . \tag{3.46}
\end{gather*}
$$

Proof. Let:

$$
\psi_{k} \triangleq 1-\chi\left(2 \rho_{k}^{-1}\left(r-\rho_{k 4}\right)\right)
$$

where $\chi$ is the function defined is step 1 , and satisfying (3.30). We have:

$$
\begin{gather*}
\psi_{k}(r)=1 \text { near } \rho_{k 4}, \quad \psi_{k}(r)=0 \text { near } \rho_{k 5}, \\
\left|\psi_{k}^{(p)}\right| \lesssim \rho_{k}^{-p / 3} \tag{3.47}
\end{gather*}
$$

Let:

$$
u \triangleq \psi_{k} u_{k}+v_{k}
$$

so that (3.45) and (3.46) are satisfied. Note also that the comparison estimate (3.34) implies that for some $c>0$ independent of $k$ :

$$
\begin{equation*}
2\left|v_{k}\right| \geq|u| \geq c\left|v_{k}\right| \geq c\left|u_{k}\right|, \quad \rho_{k 4} \leq r \leq \rho_{k 5} . \tag{3.48}
\end{equation*}
$$

We have:

$$
\begin{equation*}
\Delta u=\left(\Delta \psi_{k}\right) u_{k}+2 \nabla \psi_{k} \cdot \nabla u_{k}+\Delta v_{k} . \tag{3.49}
\end{equation*}
$$

Using the estimates of Lemma 3.2 together with estimates (3.47), (3.48) and equation (3.49) one gets (3.17).

Inequality (3.18), as in the first case of step 1 , comes easily from the explicit computation of $\partial_{r} u$, inequality (3.17) and estimates (3.24), (3.47) and (3.48)

Step 3. Construction of $u$ on $\left[\rho_{k 5}, \rho_{k 6}\right]$.
Lemma 3.6. There exists $u \in C^{\infty}\left(\left\{\rho_{k 5} \leq r \leq \rho_{k 6}\right\}\right)$, satisfying (3.17) and (3.18) and so that:

$$
\begin{gather*}
u(r, \theta)=v_{k}(r, \theta), \quad \text { near } \rho_{k 5},  \tag{3.50}\\
u(r, \theta)=-\bar{u}_{k+1}(r, \theta), \quad \text { near } \rho_{k 6} . \tag{3.51}
\end{gather*}
$$

Proof. Consider the function $\widetilde{\psi}_{k}$ defined by:

$$
\widetilde{\psi}_{k}(r) \triangleq \psi_{k}\left(r-\rho_{k}^{1 / 3}\right),
$$

where $\psi_{k}$ is the function of Lemma 3.5. The function $\widetilde{\psi}_{k}$ is 1 near $\rho_{k 5}$ and 0 near $\rho_{k 6}$, and satisfies estimates (3.47). Recall the definition (3.20) of $v_{k}$ and let:

$$
u \triangleq \widetilde{\psi}_{k} v_{k}-\left(1-\widetilde{\psi}_{k}\right) \bar{u}_{k+1}=-\bar{u}_{k+1}\left\{1-\widetilde{\psi}_{k}+\widetilde{\psi}_{k} r^{4 d_{k}} \rho_{k 3}^{-4 d_{k}} e^{i \varphi_{k}(\theta)}\right\}
$$

According to (3.19), there is a constant $C>0$ independent of $k$ and $\theta$ such that:

$$
\operatorname{Re} e^{i \varphi_{k}(\theta)} \geq 1-C \rho_{k}^{-4 / 3},
$$

so that for $k$ large enough, and using (3.22):

$$
\operatorname{Re}\left(1-\widetilde{\psi}_{k}+\widetilde{\psi}_{k} r^{4 d_{k}} \rho_{k 3}^{-4 d_{k}} e^{i \tilde{\varphi}_{k}(\theta)}\right) \geq 1-\widetilde{\psi}_{k}+c_{1} \widetilde{\psi}_{k} \geq c_{1}
$$

for some positive constant $c_{1}$. Thus:

$$
\begin{equation*}
|u| \geq c_{1}\left|u_{k+1}\right|, \quad \rho_{k 5} \leq r \leq \rho_{k 6} . \tag{3.52}
\end{equation*}
$$

Furthermore:

$$
\begin{equation*}
\Delta u=\widetilde{\psi}_{k} \Delta v_{k}+2 \nabla \widetilde{\psi}_{k} \cdot \nabla v_{k}+\left(\Delta \widetilde{\psi}_{k}\right) v_{k}-2 \nabla \widetilde{\psi}_{k} \cdot \nabla u_{k+1}-\left(\Delta \widetilde{\psi}_{k}\right) u_{k+1} . \tag{3.53}
\end{equation*}
$$

Using Lemma 3.2 and the estimates (3.47) on $\widetilde{\psi}_{k}$ together with equation (3.53), we get:

$$
\Delta u=O\left(u_{k+1}\right), \quad \rho_{k 5} \leq r \leq \rho_{k 6} .
$$

We conclude with (3.52) that $u$ satisfies (3.17) on $\left[\rho_{k 5}, \rho_{k 6}\right]$.
To finish step 3 , we have to show inequality (3.18). As in the preceding steps, it comes easily from inequality (3.17), estimates (3.24) and the explicit computation of $\partial_{r} u$.

The construction of $u$ on $\left[\rho_{k}, \rho_{k+1}\right]$ is complete. According to Lemmas 3.3, 3.5 and 3.6, $u$ satisfies (3.17) and (3.18).

It remains to check that on $\left[\rho_{k}, \rho_{k+1}\right]$, $u$ satisfies the bound (3.16). Indeed, by the definition of $u$ at each step, it is easy to deduce (3.16) from the same bound on functions $u_{k}, v_{k}$ and $u_{k+1}$. This bound is trivial for $u_{k}$. But $u_{k}, v_{k}$ and $u_{k+1}$ are of the same order (see (3.23)), hence (3.16). This concludes the proof of Lemma 3.1.

If $k$ is large enough for the preceding lemma to hold and $r \in\left[\rho_{k}, \rho_{k+1}\right]$, we take $u(r)$ to be the function constructed in the lemma, satisfying (3.15) if $k$ is odd and (3.15') is $k$ is even. In this way, the pieces of $u$ stick up well together at each $\rho_{k}$, and this defines a $C^{\infty}$ function $u$ for $r \geq \rho$, where $\rho=\rho_{K}$ is a large positive real number. According to the uniform inequalities (3.17) and (3.18) satisfied by $u$ on each $\left[\rho_{k}, \rho_{k+1}\right]$, the function $u$ is solution of (3.9) and (3.10). It remains to check the decay of $u$ at infinity, and to extend $u$ to all $\mathbb{R}^{2}$.

### 3.1.2 Decay of $u$ at infinity

Take a point of $\mathbb{R}^{2}$ with coordinates $(r, \theta)$ such that:

$$
\begin{equation*}
\rho_{k} \leq r \leq \rho_{k+1} . \tag{3.54}
\end{equation*}
$$

Let:

$$
h \triangleq \frac{r-\rho_{k}}{\rho_{k}}=O\left(\rho_{k}^{-2 / 3}\right) .
$$

Estimate (3.16) yields a constant $C$, independent of $k$ and $r$ satisfying (3.54), such that:

$$
|u(r, \theta)| \leq C a_{k} r^{-n_{k}}
$$

Thus:

$$
\begin{aligned}
\log |u(r, \theta)|-\log \left|u\left(\rho_{k}, \theta\right)\right| & \leq-n_{k} \log r+n_{k} \log \rho_{k}+O(1) \\
& \leq-n_{k} \log (1+h)+O(1) \leq-n_{k} h+O(1)
\end{aligned}
$$

using the fact that $n_{k} h^{2}$ is bounded independently of $r$ and $k$. On the other hand, if $m(r) \triangleq e^{-\frac{3}{4} r^{4 / 3}}$ :

$$
\begin{aligned}
& \log m(r)-\log m\left(\rho_{k}\right)=-\frac{3}{4} r^{4 / 3}+\frac{3}{4} \rho_{k}^{4 / 3} \\
& \quad=-\frac{3}{4}\left\{\rho_{k}^{4 / 3}\left((1+h)^{4 / 3}-1\right)\right\}=-\rho_{k}^{4 / 3} h+O(1)
\end{aligned}
$$

Thus, recalling that $n_{k}=2\left[\rho_{k}^{4 / 3} / 2\right]$ :

$$
\begin{equation*}
\log |u(r, \theta)|-\log \left|u\left(\rho_{k}, \theta\right)\right| \leq \log m(r)-\log m\left(\rho_{k}\right)+O(1) . \tag{3.55}
\end{equation*}
$$

The same argument yields, if $K \leq j \leq k$ :

$$
\begin{equation*}
\log \left|u\left(\rho_{j}, \theta\right)\right|-\log \left|u\left(\rho_{j-1}, \theta\right)\right| \leq \log m\left(\rho_{j}\right)-\log m\left(\rho_{j-1}\right)+O(1) \tag{3.56}
\end{equation*}
$$

Adding inequality (3.55) and all inequalities (3.56), $K \leq j \leq k$, we get:

$$
\log |u(r, \theta)| \leq \log m(r)+O(k)
$$

It is classical that a sequence $\rho_{k}$ defined by the induction relation (3.11) is of order $k^{3 / 2}$. Hence:

$$
|u(r, \theta)| \leq e^{-3 / 4 r^{4 / 3}+C r^{2 / 3}}
$$

Which gives (3.3) for any $c_{*}<3 / 4$.

### 3.1.3 Extension of $u$ to all $\mathbb{R}^{2}$

So far, we have constructed $u$ on $r \geq \rho$, equal to $a r^{-n} e^{i n \theta}$ near $\rho$ for some integer $n$ and real $a$. Let $\psi$ be a smooth, nondecreasing, function equal to 1 for $r \geq 2 \rho / 3$ and 0 for $r \leq \rho / 3$. Let:

$$
u(r, \theta) \triangleq\left(\psi(r) r^{-n}+(1-\psi(r)) r^{n}\right) a e^{i n \theta}, \quad r \leq \rho .
$$

This extends $u$ to a $C^{\infty}$ function on $\mathbb{R}^{2}$, harmonic in a neighborhood of 0 , and who does not vanish, which implies trivially (3.9) for $r \leq \rho$. Similarly, $\nabla u$ does not vanish for $r>0$ (because $\partial_{\theta} u$ does not) which gives (3.10) for $r \leq \rho$. The construction is complete.

### 3.2 Construction in odd dimension

In this part we prove Theorem 3.2. We first remark that we only need to do the construction for $n=3$. Indeed if Theorem 3.2 holds for $n=3$, and $n=m+3$ is an odd number larger than 3 , one can define the function:

$$
\tilde{u}\left(x_{1}, x_{2}, \ldots, x_{m+3}\right) \triangleq v\left(x_{1}, \ldots, x_{m}\right) u\left(x_{m+1}, x_{m+2}, x_{m+3}\right),
$$

where $v$ is the complex-valued function $u$ defined on $\mathbb{R}^{m}$ given by Theorem 3.1, and $u$ is the $\mathbb{C}^{4}$-valued function defined on $\mathbb{R}^{3}$ given by Theorem 3.2. Note that $\tilde{u}$ takes values in $\mathbb{C}^{4}$. A straightforward computation shows that the function $\tilde{u}$ and potential $\tilde{q}$ are solutions of the equation (3.1), where $\tilde{q}$ satisfies the bound (3.5) and $\tilde{u}$ decays at the desired speed (3.6).

We now turn to the proof of the case $n=3$. One of the main ingredients of the preceding construction was the sequence of eigenfunctions $\left(e^{i n_{k} \theta}\right)_{k}$ of the Laplace operator on $S^{1}$, which trivially satisfies the estimate (in the sense given by Notation 3.1):

$$
\begin{equation*}
e^{i n_{k} \theta} \approx e^{i n_{k+1} \theta} \tag{3.57}
\end{equation*}
$$

The construction is difficult to adapt in dimension 3, since there is no sequence of spherical harmonics on $S^{2}$ satisfying (3.57). To show Theorem 3.2, we write an abstract theorem showing that an estimate of the form (3.57), but with polynomial loss in $n_{k}$, is sufficient to construct a vectorvalued, superexponentially decaying solution of an equation of the form (3.1), with a potential $q$ which only grows logarithmically.

Consider a smooth manifold $M$ without boundary, and an operator:

$$
R: C^{\infty}(M) \longrightarrow C^{\infty}(M)
$$

We define, for $\rho \geq 0$ :

$$
\widetilde{M}_{\rho} \triangleq(\rho,+\infty) \times M, \quad P \triangleq \frac{\partial^{2}}{\partial_{r}^{2}}+\frac{1}{r^{2}} R
$$

The operator $P$ acts on $C^{\infty}\left(\widetilde{M}_{0}\right)$. Up to the conjugation by a power of $r$, and the addition of a zero-order potential, this framework includes the Laplace operator on $\mathbb{R}^{n}, n \geq 2$. Let $p \geq 1$. Assume that $R$ admits a sequence of bounded eigenfunctions $\left(\Phi_{k}\right)_{k \geq 0}$ :

$$
\begin{gather*}
\Phi_{k}: M \longrightarrow \mathbb{R}^{p}, \quad R \Phi_{k} \triangleq-\lambda_{k} \Phi_{k}, \quad \lambda_{k}>0,  \tag{3.58}\\
\left\|\Phi_{k}\right\|_{L^{\infty}(M)}=1 . \tag{3.59}
\end{gather*}
$$

where the sequence $\left(\lambda_{k}\right)_{k}$ is increasing and tends to infinity. Define $n_{k}$ and $\rho_{k}$ by:

$$
\begin{gather*}
n_{k}\left(n_{k}+1\right)=\lambda_{k}, \quad n_{k} \geq 0 \\
\rho_{k} \triangleq n_{k}^{3 / 4}, \quad d_{k} \triangleq \frac{n_{k+1}-n_{k}}{2} \tag{3.60}
\end{gather*}
$$

Denote by $|\cdot|$ the euclidian norm on $\mathbb{R}^{p}$. Then the following holds:
Theorem 3.3. Assume (3.58), (3.59) and that there exist positive constants $\delta, C, N$ such that:

$$
\begin{gather*}
\forall \omega \in M, \quad \frac{1}{C n_{k}^{N}} \leq \frac{\left|\Phi_{k}(\omega)\right|}{\left|\Phi_{k+1}(\omega)\right|} \leq C n_{k}^{N}  \tag{3.61}\\
d_{k}=\delta n_{k}^{1 / 2}+O(1) \tag{3.62}
\end{gather*}
$$

Let $c_{*}>0$. Then, if $\rho$ is large enough:

$$
\begin{gather*}
\exists u \in C^{\infty}\left(\widetilde{M}_{\rho} ; \mathbb{R}^{2 p}\right), \quad \exists C>0, \quad P u=q u  \tag{3.63}\\
q \in C^{\infty}\left(\widetilde{M}_{\rho} ; \mathbb{R}^{2 p \times 2 p}\right), \quad(\log (r+2))^{-3} q \in L^{\infty},  \tag{3.64}\\
|u(r, \omega)| \leq C e^{-c_{*} r^{4 / 3}} \tag{3.65}
\end{gather*}
$$

Remark 3.4. As will appear clearly in the proof, when the power $N$ of $n_{k}$ is 0 in (3.61), the same result remains valid with a bounded $q$. This would yield Theorem 3.1, with an easier proof, but a $\mathbb{C}^{2}$-valued solution $u$.

Proof of Theorem 3.3. This construction is very similar, although much simpler because of the vectorial setting, than the preceding one. Denote by $\rho_{k j} \triangleq \rho_{k}+j \frac{\rho_{k+1}-\rho_{k}}{4}$, which divides $\left(\rho_{k}, \rho_{k+1}\right)$ in 4 subintervals. Note that according to (3.62):

$$
\begin{equation*}
\rho_{k+1}-\rho_{k}=\frac{3}{2} \delta \rho_{k}^{1 / 3}+O\left(\rho_{k}^{-1 / 3}\right) \tag{3.66}
\end{equation*}
$$

Consider the following solutions of the equation $P E=0$ :

$$
\begin{equation*}
E_{k}(r, \omega) \triangleq a_{k} r^{-n_{k}} \Phi_{k}(\omega) \tag{3.67}
\end{equation*}
$$

where the sequence $a_{k}$ is defined by:

$$
a_{0} \triangleq 1, \quad a_{k+1} \triangleq \rho_{k}^{2 d_{k}} a_{k}
$$

so that $a_{k} r^{-n_{k}}$ and $a_{k+1} r^{-n_{k+1}}$ coincide when $r=\rho_{k}$. Consider the $\mathbb{R}^{2 p}$-valued functions $\mathbf{E}_{k}$ :

$$
\mathbf{E}_{k}(r, \omega) \triangleq\binom{E_{k}(r, \omega)}{0} \text { if } k \text { is even, } \quad \mathbf{E}_{k}(r, \omega) \triangleq\binom{0}{E_{k}(r, \omega)} \text { if } k \text { is odd. }
$$

Then we have the following lemma, analogous to Lemma 3.1:
Lemma 3.7. Let $k$ be a large enough integer. There exist a constant $C_{0}$ independent of $k$, and:

$$
u \in C^{\infty}\left(\left\{\rho_{k} \leq r \leq \rho_{k+1}\right\} ; \mathbb{R}^{2 p}\right)
$$

Figure 3: Step functions for Theorem 3.3

such that:

$$
\begin{gather*}
u(r, \omega)=\mathbf{E}_{k}(r, \omega), \quad \rho_{k 0} \leq r \leq \rho_{k 1}  \tag{3.68}\\
u(r, \omega)=\mathbf{E}_{k+1}(r, \omega), \quad \rho_{k 3} \leq r \leq \rho_{k 4} \\
|u(r, \omega)|=O\left(a_{k} r^{-n_{k}}\right), \quad \rho_{k} \leq r \leq \rho_{k+1} \tag{3.69}
\end{gather*}
$$

and satisfying the inequality:

$$
\begin{equation*}
|P u| \leq C_{0}(\log r)^{3}|u| \tag{3.70}
\end{equation*}
$$

Proof. Using that the logarithmic derivative of $r^{d_{k}} / \rho_{k}^{d_{k}}$ is bounded by $\rho_{k}^{-1 / 3}$, one gets, as in the proof of Theorem 3.1, that for $\rho_{k} \leq r \leq \rho_{k+1}$ :

$$
\begin{gather*}
r^{d_{k}} \approx \rho_{k}^{d_{k}}  \tag{3.71}\\
a_{k} r^{n_{k}} \approx a_{k+1} r^{n_{k+1}} \tag{3.72}
\end{gather*}
$$

We divide the construction into several steps.
Step 1: Definition of $u$. Let $\chi$ be a $C^{\infty}$ non-increasing function on $\mathbb{R}$ such that:

$$
\begin{gather*}
s \leq 0 \Longrightarrow \chi(s)=0, \quad s \geq 1 \Longrightarrow \chi(s)=1 \\
0<s \leq 1 / 2 \Longrightarrow \chi(s)=e^{-1 / s} \tag{3.73}
\end{gather*}
$$

Near $s=0, \chi, \chi^{\prime}$ and $\chi^{\prime \prime}$ are increasing functions of $s$. Let (see Figure 3):

$$
\widetilde{\chi}_{k}(r) \triangleq \chi\left(\rho_{k}^{-1 / 3}\left(\rho_{k 3}-r\right)\right), \quad \chi_{k}(r) \triangleq \chi\left(\rho_{k}^{-1 / 3}\left(r-\rho_{k 1}\right)\right)
$$

We have:

$$
\begin{equation*}
\left|\chi_{k}^{(p)}\right|=O\left(\rho_{k}^{-p / 3}\right), \quad\left|\tilde{\chi}_{k}^{(p)}\right|=O\left(\rho_{k}^{-p / 3}\right) \tag{3.74}
\end{equation*}
$$

Assume for example that $k$ is odd. Let:

$$
u(r, \omega) \triangleq\binom{\chi_{k}(r) E_{k}(r, \omega)}{\tilde{\chi}_{k}(r) E_{k+1}(r, \omega)}
$$

so that (3.68) holds. The bound (3.69) is immediate from (3.59), (3.67) and (3.72). We have:

$$
\begin{equation*}
P u=\binom{\left(\chi_{k}^{\prime \prime}+2 \frac{n_{k}}{r} \chi_{k}^{\prime}\right) E_{k}}{\left(\widetilde{\chi}_{k}+2 \frac{n_{k+}}{r} \widetilde{\chi}_{k}^{\prime}\right) E_{k+1}} . \tag{3.75}
\end{equation*}
$$

Consider the first $p$ components of $P u$ :

$$
\begin{equation*}
v \triangleq\left(\chi_{k}^{\prime \prime}+2 \frac{n_{k}}{r} \chi_{k}^{\prime}\right) E_{k} \tag{3.76}
\end{equation*}
$$

We will show that if $k$ is large enough:

$$
\begin{equation*}
|v(r, \omega)|=O\left((\log r)^{3}|u(r, \omega)|\right), \quad \rho_{k} \leq r \leq \rho_{k+1}, \omega \in M \tag{3.77}
\end{equation*}
$$

Let $s \triangleq \rho_{k}^{-1 / 3}\left(r-\rho_{k 1}\right)$. We distinguish two regions.
Step 2: Pointwise bound on $v$ for $s \notin\left(0,\left(\log \rho_{k}\right)^{-3 / 2}\right)$. Using the explicit form (3.73) of $\chi$ near 0 , a straightforward computation shows that, for $0<s \leq 1 / 2$ :

$$
\begin{gather*}
\chi_{k}^{\prime}(r)=\rho_{k}^{-1 / 3} \chi^{\prime}(s)=\rho_{k}^{-1 / 3} s^{-2} \chi_{k}(r)  \tag{3.78}\\
\chi_{k}^{\prime \prime}(r)=\rho_{k}^{-2 / 3} \chi^{\prime \prime}(s)=\rho_{k}^{-2 / 3}\left(-2 s^{-3}+s^{-4}\right) \chi_{k}(r) \tag{3.79}
\end{gather*}
$$

This shows that if $\left(\log \rho_{k}\right)^{-3 / 2} \leq s \leq \frac{1}{2}$ :

$$
\frac{n_{k}}{r}\left|\chi_{k}^{\prime}(r)\right|=O\left(\left(\log \rho_{k}\right)^{3} \chi_{k}(r)\right), \quad\left|\chi_{k}^{\prime \prime}(r)\right|=O\left(\left(\log \rho_{k}\right)^{3} \chi_{k}(r)\right)
$$

Furthermore, these inequalities are trivial for $s<0$ (where $\chi_{k}$ is identically). When $s \geq 1 / 2$ they are a direct consequence of the estimates (3.74) on the derivatives of $\chi_{k}$. Going back to the definition of $v$, we have:

$$
|v(r, \omega)| \leq(\log r)^{3} \chi_{k}(r)\left|E_{k}(r, \omega)\right|, \quad s \leq 0 \text { or } s \geq\left(\log \rho_{k}\right)^{-3 / 2}
$$

This shows inequality (3.77), outside of the region $\left\{s \in\left(0,\left(\log \rho_{k}\right)^{-3 / 2}\right)\right\}$.
Step 3: Pointwise bound on $v$ for $s \in\left(0,\left(\log \rho_{k}\right)^{-3 / 2}\right)$. We now assume that $s \in\left(0, \log \left(\rho_{k}\right)^{-3 / 2}\right)$. Then, using formulas (3.78), (3.79) and the fact that $\chi^{\prime}$ is increasing we get, if $k$ is large:

$$
\begin{align*}
& \frac{n_{k}}{r}\left|\chi_{k}^{\prime}(s)\right|=\frac{n_{k}}{r} \rho_{k}^{-1 / 3}\left|\chi^{\prime}(s)\right| \leq C\left|\chi^{\prime}\left(\log ^{-3 / 2} \rho_{k}\right)\right| \\
& \frac{n_{k}}{r}\left|\chi_{k}^{\prime}(s)\right| \leq C\left(\log \rho_{k}\right)^{3} e^{-\left(\log \rho_{k}\right)^{3 / 2}} \tag{3.80}
\end{align*}
$$

Similarly:

$$
\begin{align*}
\left|\chi_{k}^{\prime \prime}(s)\right| & =\rho_{k}^{-2 / 3}\left|\chi^{\prime \prime}(s)\right| \leq \rho_{k}^{-2 / 3}\left|\chi^{\prime \prime}\left(\log ^{-3 / 2} \rho_{k}\right)\right| \\
& \leq \rho_{k}^{-2 / 3}\left(\log \rho_{k}\right)^{6} e^{-\left(\log \rho_{k}\right)^{3 / 2}} . \tag{3.81}
\end{align*}
$$

By the definition (3.76) of $v$, together with (3.80), (3.81), we get:

$$
|v(r, \omega)| \leq\left(\log \rho_{k}\right)^{3} e^{-\left(\log \rho_{k}\right)^{3 / 2}}\left|E_{k}(\omega)\right| \leq C\left(\log \rho_{k}\right)^{3} e^{-\left(\log \rho_{k}\right)^{3 / 2}} n_{k}^{N}\left|E_{k+1}(\omega)\right| .
$$

For the second inequality, we used the assumption (3.61) on the sequence ( $\Phi_{k}$ ) together with (3.72). Noting that $\widetilde{\chi}_{k}$ takes value 1 near $\rho_{k 1}$, we get that (3.77) holds for large $k$.

End of the proof. By the same argument, one may show the property analogous to (3.77) for the last $p$ components of $P u$, namely:

$$
\left(\widetilde{\chi}_{k}+2 \frac{n_{k+1}}{r} \widetilde{\chi}_{k}^{\prime}\right) E_{k+1}=O\left((\log r)^{3} u(r, \omega)\right), \quad \rho_{k} \leq r \leq \rho_{k+1}, \omega \in M
$$

Thus:

$$
P u=O\left((\log u)^{3} u\right), \quad \rho_{k} \leq r \leq \rho_{k+1}, \omega \in M,
$$

which completes the proof of the lemma.
The end of the proof of Theorem 3.3, which consists in sticking up the pieces of $u$ defined by Lemma 3.7, and checking the decay of $u$ at infinity, is exactly the same as the one of Theorem 3.1, and therefore we omit it.

Proof of Theorem 3.2. We shall use Theorem 3.3 with $M=S^{2}$. For this we need to choose suitable spherical harmonics. Let $\theta$ and $\phi$ be the spherical coordinates on the sphere $S^{2}, \theta \in[0, \pi]$ being the polar coordinate and $\phi \in[0,2 \pi)$ the azimutal one. Let $l=2 j$ be an even integer and $F_{l}$ be the $\mathbb{C}^{2}$-valued spherical harmonic:

$$
\begin{equation*}
F_{l}(\phi, \theta) \triangleq\binom{P_{l}^{0}(\cos \theta)}{e^{i \phi} P_{l}^{1}(\cos \theta)} . \tag{3.82}
\end{equation*}
$$

Here and in the sequel, the $P_{l}^{m}$ are the associated Legendre polynomials:

$$
\begin{equation*}
P_{l}^{m}(x)=\frac{(-1)^{m}}{2^{l} l!}\left(1-x^{2}\right)^{m / 2} \frac{d^{l+m}}{d x^{l+m}}\left(x^{2}-1\right)^{l} . \tag{3.83}
\end{equation*}
$$

The functions $P_{l}^{m}$ are solutions to the equation:

$$
\begin{equation*}
\left(1-x^{2}\right) \frac{d^{2} P}{d x^{2}}-2 x \frac{d P}{d x}+\left[l(l+1)-\frac{m^{2}}{1-x^{2}}\right] P=0, \quad x \in(-1,1) . \tag{3.84}
\end{equation*}
$$

It is standard (c.f. [26], [24]), that $\Delta_{S^{2}} F_{l}=-l(l+1) F_{l}$. To use Theorem 3.3 we need to give a pointwise estimate on $F_{l}$ :

Lemma 3.8. The eigenfunctions $F_{l}$ satisfy, for large $l \in 2 \mathbb{N}$ :

$$
\begin{equation*}
\frac{1}{C l^{7 / 4}} \leq\left|F_{l}(\phi, \theta)\right| \leq C l . \tag{3.85}
\end{equation*}
$$

Proof. Let:

$$
g(x) \triangleq l(l+1)\left|P_{l}^{0}(x)\right|^{2}+\left|P_{l}^{1}(x)\right|^{2}
$$

Note that according to (3.83),

$$
g(x)=l(l+1)\left|P_{l}^{0}(x)\right|^{2}+\left(1-x^{2}\right)\left|\frac{d P_{l}^{0}}{d x}(x)\right|^{2} .
$$

Using equation (3.84) with $m=0$ we get:

$$
g^{\prime}(x)=2 x\left|\frac{d P_{l}^{0}}{d x}(x)\right|^{2}
$$

so that the minimum of $g$ is in 0 , and its maxima are in 1 and -1 . Note that $l$ being even, $P_{l}^{0}$ is even and $P_{l}^{1}$ is odd. In particular:

$$
\begin{equation*}
g(0)=l(l+1)\left|P_{l}^{0}(0)\right|^{2} \tag{3.86}
\end{equation*}
$$

Let $\alpha_{j} \triangleq\left|P_{2 j}^{0}(0)\right|=\left|P_{l}^{0}(0)\right|$. According to formula (3.83),

$$
P_{2 j}^{0}(x)=\frac{1}{2^{2 j}(2 j)!} \frac{d^{2 j}}{d x^{2 j}}\left(x^{2}-1\right)^{2 j}
$$

The coefficient of $x^{2 j}$ in $\left(x^{2}-1\right)^{2 j}$ is $(-1)^{j}\binom{2 j}{j}$, so that we have:

$$
\alpha_{j}=\frac{1}{2^{2 j}(2 j)!}(2 j)!\binom{2 j}{j}=\frac{(2 j)!}{2^{2 j}(j!)^{2}},
$$

Stirling formula, yields:

$$
\begin{equation*}
\left|P_{l}^{0}(0)\right|=\left|P_{2 j}^{0}(0)\right|=\alpha_{j} \sim \frac{1}{2 \sqrt{\pi j}}, \quad \text { as } j \rightarrow+\infty \tag{3.87}
\end{equation*}
$$

Furthermore, going back to the definition (3.83) of $P_{m}^{l}$, we get $\left|P_{l}^{0}(1)\right|=\left|P_{l}^{0}(-1)\right|=1$, so that:

$$
\begin{equation*}
g(1)=g(-1)=l(l+1) . \tag{3.88}
\end{equation*}
$$

Using (3.86), (3.87) and (3.88) we get:

$$
\frac{1}{C l^{3 / 2}} \leq l(l+1)\left|P_{l}^{0}(x)\right|^{2}+\left|P_{l}^{1}(x)\right|^{2} \leq C l(l+1),
$$

which shows (3.85) according to the definition (3.82) of $F_{l}$.
Choose a large odd number $n_{0}$ and define the sequence $\left(n_{k}\right)_{k}$ of odd integers, and the sequence $\left(\Phi_{k}\right)_{k}$ of eigenfunctions of $\Delta_{S^{2}}$ by:

$$
n_{k+1} \triangleq n_{k}+2\left[n_{k}^{1 / 2}\right], \quad \Phi_{k}(\phi, \theta) \triangleq c_{k} F_{n_{k}-1}(\phi, \theta)
$$

where $c_{k}$ is a normalizing constant such that (3.59) holds. Noting that (3.62) is fulfilled by the choice of $\left(n_{k}\right)_{k}$, and that (3.85) implies (3.61), one can use Theorem 3.3 with $R=\Delta_{S^{2}}$. This yields a function $\tilde{u} \in C^{\infty}\left(\left\{x \in \mathbb{R}^{3}| | x \mid \geq \rho\right\} ; \mathbb{C}^{4}\right)$, solution of:

$$
P \tilde{u}=\left(\frac{\partial^{2}}{\partial_{r}^{2}}+\frac{1}{r^{2}} \Delta_{S^{2}}\right) \tilde{u}=\tilde{q} \tilde{u},
$$

and so that $\tilde{q}$ and $\tilde{u}$ decrease at the desired speed at infinity. Taking into account that $\Delta_{\mathbb{R}^{3}} f=$ $r^{-1} P(r f)$, the function $u \triangleq r^{-1} \tilde{u}$ and potential $q \triangleq \tilde{q}$ satisfy all the assertions of Theorem 3.2. It remains to extend $u$ to $r \leq \rho$, which is left to the reader.

## 4 Optimality of the observability constant for the heat equation with zero order potential

This section is addressed to the proof of Theorem 1.1 and its weaker counterpart in odd space dimensions:

Theorem 4.1. Assume that $n \geq 3$ is odd and that $N \geq 8$. Let $\omega$ be a non-empty open subset of $\Omega$ such that $\Omega \backslash \bar{\omega} \neq \emptyset$. Then there exist two constants $c>0$ and $\mu>0$, a family of potentials $\left\{a_{R}\right\}_{R>0} \subset L^{\infty}\left(Q ; \mathbb{R}^{N \times N}\right)$ satisfying

$$
\left\|a_{R}\right\|_{\infty} \xrightarrow{\longrightarrow}+\infty
$$

and a family of initial data $\left\{\varphi_{R}^{0}\right\}_{R>0}$ in $\left(L^{2}(\Omega)\right)^{N}$ such that the corresponding solution $\varphi_{R}$ of (1.1) satisfies

$$
\begin{equation*}
\lim _{R \rightarrow \infty}\left\{\inf _{T \in J_{\mu}} \frac{\left\|\varphi_{R}(T)\right\|_{\left(L^{2}(\Omega)\right)^{N}}^{2}}{\exp \left(c\left(\log \left\|a_{R}\right\|_{\infty}\right)^{-2}\left\|a_{R}\right\|_{\infty}^{2 / 3}\right) \int_{0}^{T} \int_{\omega}\left|\varphi_{R}\right|^{2} d x d t}\right\}=+\infty \tag{4.1}
\end{equation*}
$$

where $J_{\mu} \triangleq\left(0, \mu\left(\log \left\|a_{R}\right\|_{\infty}\right)^{-2}\left\|a_{R}\right\|_{\infty}^{-1 / 3}\right]$.
Theorem 4.1 would show the optimality of an observability constant a little smaller than $C_{3}^{*}(T, a)$ defined in (1.5). This is due to the logarithmic loss in the construction of Theorem 3.2.

The proofs of Theorems 1.1 and 4.1 are very similar so that we do not need to distinguish between the two cases in the main part of this section. First observe that we only need to show Theorem 1.1 in the case $N=2$ and Theorem 4.1 in the case $N=8$. Indeed, to get the other cases, it suffices to consider the solutions $\varphi_{R}$ in which the first 2 components (respectively 8 components in the case of odd dimensions) are as in the case $N=2$ (respectively $N=8$ ), the other being identically zero. The matrix $a_{R}$ of the corresponding system can be built in a similar way by adding zero entries to the $2 \times 2$ (respectively $8 \times 8$ ) matrix $a_{R}$.

The proof is divided into several steps.
Step 1: Construction on $\mathbb{R}^{n}$.
Consider the solution $u$ and potential $q$ given by Theorem 3.1 if $n$ is even, and by Theorem 3.2 if $n$ is odd. Recalling that both $u$ and $q$ are complex-valued, by setting

$$
u_{R}(x)=\binom{\operatorname{Re} u(R x)}{\operatorname{Im} u(R x)}, \quad a_{R}(x)=-R^{2}\left(\begin{array}{cc}
\operatorname{Re} q(R x) & -\operatorname{Im} q(R x)  \tag{4.2}\\
\operatorname{Im} q(R x) & \operatorname{Re} q(R x)
\end{array}\right)
$$

we obtain a one-parameter family of potentials $\left\{a_{R}\right\}_{R>0}$ and solutions $\left\{u_{R}\right\}_{R>0}$ satisfying

$$
\begin{equation*}
\Delta u_{R}=a_{R}(x) u_{R}, \quad \text { in } \mathbb{R}^{n}, \tag{4.3}
\end{equation*}
$$

and (using Theorems 3.1-3.2 with $c_{*}=1$ )

$$
\begin{equation*}
\left|u_{R}(x)\right| \leq C \exp \left(-R^{4 / 3}|x|^{4 / 3}\right), \quad \text { in } \mathbb{R}^{n} \tag{4.4}
\end{equation*}
$$

Furthermore, for some constant $C>0$, the potential $a_{R}$ is such that

$$
\begin{align*}
& C^{-1} R^{2} \leq\left\|a_{R}\right\|_{L^{\infty}\left(\mathbb{R}^{n} ; \mathbb{R}^{N \times N}\right)} \leq C R^{2}, \quad \text { if } n \text { is even },  \tag{4.5}\\
& C^{-1} R^{2} \leq\left\|a_{R}\right\|_{L^{\infty}\left(\mathbb{R}^{n} ; \mathbb{R}^{N \times N}\right)} \leq C R^{2}(\log R)^{3}, \quad \text { if } n \text { is odd. } \tag{4.6}
\end{align*}
$$

The functions $\left\{u_{R}\right\}_{R>0}$ may also be viewed as stationary solutions of the corresponding parabolic systems. Indeed, set

$$
\begin{equation*}
\psi_{R}(t, x)=u_{R}(x), \quad x \in \mathbb{R}^{n}, t>0 . \tag{4.7}
\end{equation*}
$$

Then, it satisfies

$$
\begin{equation*}
\psi_{R, t}-\Delta \psi_{R}+a_{R} \psi_{R}=0, \quad x \in \mathbb{R}^{n}, t>0 \tag{4.8}
\end{equation*}
$$

and

$$
\begin{equation*}
\left|\psi_{R}(x, t)\right| \leq C \exp \left(-R^{4 / 3}|x|^{4 / 3}\right), \quad x \in \mathbb{R}^{n}, t>0 \tag{4.9}
\end{equation*}
$$

Assume now that $\omega$ is an open bounded subset of $\mathbb{R}^{n} \backslash B, B$ being the unit ball in $\mathbb{R}^{n}$. Then

$$
\begin{equation*}
\int_{0}^{T} \int_{\omega}\left|\psi_{R}\right|^{2} d x d t \leq T C \int_{\omega} \exp \left(-2 R^{4 / 3}|x|^{4 / 3}\right) d x \tag{4.10}
\end{equation*}
$$

and, taking into account that $|x| \geq 1$ on $\omega$,

$$
\begin{equation*}
\int_{\omega} \exp \left(-2 R^{4 / 3}|x|^{4 / 3}\right) d x=O\left(\exp \left(-2 R^{4 / 3}\right)\right) \tag{4.11}
\end{equation*}
$$

On the other hand, for some constant $c>0$,

$$
\begin{equation*}
\left\|\psi_{R}(T)\right\|_{\left(L^{2}\left(\mathbb{R}^{n}\right)\right)^{N}}^{2}=\left\|u_{R}\right\|_{\left(L^{2}\left(\mathbb{R}^{n}\right)\right)^{N}}^{2}=\frac{1}{R^{n}}\|u\|_{\left(L^{2}\left(\mathbb{R}^{n}\right)\right)^{N}}^{2}=\frac{c}{R^{n}} . \tag{4.12}
\end{equation*}
$$

In view of (4.5) and (4.10)-(4.12) it is easy to see that, if $n$ is even, an estimate of the form (1.2) would be sharp in the whole $\mathbb{R}^{n}$ in what concerns the dependence of the observability constant $C_{3}^{*}$ on the potential. The same conclusion would hold for odd $n$, up to a logarithmic factor.

## Step 2: Restriction to $\Omega$.

Let us now consider the case of a bounded domain $\Omega$ and $\omega$ to be a non-empty open subset $\Omega$ such that $\Omega \backslash \bar{\omega} \neq \emptyset$. Without loss of generality (by translation and scaling) we can assume that $B \subset \Omega \backslash \bar{\omega}$.

We can then view the functions $\left\{\psi_{R}\right\}_{R>0}$ above as a family of solutions of the Dirichlet problem in $\Omega$ with non-homogeneous Dirichlet boundary conditions:

$$
\begin{cases}\psi_{R, t}-\Delta \psi_{R}+a_{R} \psi_{R}=0, & \text { in } Q,  \tag{4.13}\\ \psi_{R}=\varepsilon_{R}, & \text { on } \Sigma,\end{cases}
$$

where

$$
\begin{equation*}
\varepsilon_{R}=\left.\psi_{R}\right|_{\Gamma}=\left.u_{R}\right|_{\Gamma} . \tag{4.14}
\end{equation*}
$$

Taking into account that both $\omega$ and $\Gamma$ are contained in the complement of $B$, we deduce that, for a suitable $C$ :

$$
\begin{array}{r}
\left|\psi_{R}(t, x)\right| \leq C \exp \left(-R^{4 / 3}\right), \quad x \in \omega, \quad 0<t<T \\
\left|\varepsilon_{R}(t, x)\right| \leq C \exp \left(-R^{4 / 3}\right), \quad x \in \Gamma, \quad 0<t<T . \tag{4.16}
\end{array}
$$

We can then correct these solutions to fulfill the Dirichlet homogeneous boundary condition. For this purpose, we introduce the correcting terms

$$
\begin{cases}\rho_{R, t}-\Delta \rho_{R}+a_{R} \rho_{R}=0, & \text { in } Q,  \tag{4.17}\\ \rho_{R}=\varepsilon_{R}, & \text { on } \Sigma, \\ \rho_{R}(0, x)=0, & \text { in } \Omega,\end{cases}
$$

and then set

$$
\begin{equation*}
\varphi_{R}=\psi_{R}-\rho_{R} . \tag{4.18}
\end{equation*}
$$

Clearly $\left\{\varphi_{R}\right\}_{R>0}$ is a family of solutions to parabolic systems of the form (1.1) with potentials $a=a_{R}(x)$.

Let us show that $\varphi_{R}$ is the family of solutions that fulfills (1.7). By the bounds (4.5) and (4.6) on $a_{R}$, we have, for some $C>0$ :

$$
\begin{array}{r}
C^{-1} R^{2} \leq\left\|a_{R}\right\|_{\infty} \leq C R^{2}, \quad \text { if } n \text { is even, } \\
C^{-1} R^{2} \leq\left\|a_{R}\right\|_{\infty} \leq C R^{2}(\log R)^{3}, \quad \text { if } n \text { is odd. } \tag{4.20}
\end{array}
$$

Furthermore, according to (4.10)-(4.11) and (4.12):

$$
\begin{gather*}
\int_{0}^{T} \int_{\omega}\left|\psi_{R}\right|^{2} d x d t \leq T \int_{\omega}\left|u_{R}\right|^{2} d x \leq C T \exp \left(-2 R^{4 / 3}\right), \quad \text { as } R \rightarrow \infty  \tag{4.21}\\
\left\|\psi_{R}(T)\right\|_{\left(L^{2}(\Omega)\right)^{N}}^{2}=\frac{1}{R^{n}}\|u\|_{\left(L^{2}(R \Omega)\right)^{N}}^{2} \geq \frac{c_{2}}{R^{n}}, \quad \text { for some constant } c_{2}>0 . \tag{4.22}
\end{gather*}
$$

Let us now analyze the corrector term $\rho_{R}$. We decompose it as

$$
\begin{equation*}
\rho_{R}=\sigma_{R}+\xi_{R} \tag{4.23}
\end{equation*}
$$

where

$$
\begin{cases}\sigma_{R, t}-\Delta \sigma_{R}=0, & \text { in } Q,  \tag{4.24}\\ \sigma_{R}=\varepsilon_{R}, & \text { on } \Sigma, \\ \sigma_{R}(0)=0, & \text { in } \Omega\end{cases}
$$

By the maximum principle and (4.16) we know that

$$
\begin{equation*}
\left\|\sigma_{R}\right\|_{\left(L^{\infty}(Q)\right)^{N}} \leq\left\|\varepsilon_{R}\right\|_{\left(L^{\infty}(\Sigma)\right)^{N}} \leq C \exp \left(-R^{4 / 3}\right) \tag{4.25}
\end{equation*}
$$

On the other hand, the reminder $\xi_{R}$ fulfills

$$
\begin{cases}\xi_{R, t}-\Delta \xi_{R}+a_{R} \xi_{R}=-a_{R} \sigma_{R}, & \text { in } Q,  \tag{4.26}\\ \xi_{R}=0, & \text { on } \Sigma, \\ \xi_{R}(0)=0, & \text { in } \Omega\end{cases}
$$

A standard energy estimate shows, together with (4.25), that

$$
\begin{equation*}
\left\|\xi_{R}(t)\right\|_{\left(L^{2}(\Omega)\right)^{N}}^{2} \leq C \exp \left(-2 R^{4 / 3}\right) \exp \left(2 t\left\|a_{R}\right\|_{\infty}\right), \quad \forall t \geq 0 \tag{4.27}
\end{equation*}
$$

We now distinguish between the two cases: $n$ even or $n$ odd.
End of the proof when $n$ is even.
Let $T \leq \mu\left\|a_{R}\right\|_{\infty}^{-1 / 3}$. According to (4.27), we have, for time $t \leq T$ :

$$
\begin{equation*}
\left\|\xi_{R}(t)\right\|_{\left(L^{2}(\Omega)\right)^{N}}^{2} \leq C \exp \left(-2 R^{4 / 3}+2 \mu\left\|a_{R}\right\|_{\infty}^{2 / 3}\right) \tag{4.28}
\end{equation*}
$$

Therefore, by choosing $\mu>0$ small enough we deduce, using estimate (4.19) on $a_{R}$ :

$$
\begin{equation*}
\left\|\xi_{R}(t)\right\|_{\left(L^{2}(\Omega)\right)^{N}} \leq C \exp \left(-\frac{1}{2} R^{4 / 3}\right) \tag{4.29}
\end{equation*}
$$

Combining (4.25) and (4.29) we deduce that

$$
\begin{equation*}
\left\|\rho_{R}(T)\right\|_{\left(L^{2}(\Omega)\right)^{N}} \leq C \exp \left(-\frac{1}{2} R^{4 / 3}\right) \tag{4.30}
\end{equation*}
$$

In view of (4.22) and (4.30) we conclude that

$$
\begin{equation*}
\left\|\varphi_{R}(T)\right\|_{\left(L^{2}(\Omega)\right)^{N}}^{2} \geq \frac{c_{2}}{2 R^{n}} \tag{4.31}
\end{equation*}
$$

as $R \rightarrow \infty$.
On the other hand, integrating (4.30) with respect to time we deduce that

$$
\int_{0}^{T} \int_{\Omega}\left|\rho_{R}\right|^{2} d x d t \leq C T \exp \left(-\frac{1}{2} R^{4 / 3}\right)
$$

Obviously this implies, in particular, that

$$
\begin{equation*}
\int_{0}^{T} \int_{\omega}\left|\rho_{R}\right|^{2} d x d t \leq C \exp \left(-\frac{1}{2} R^{4 / 3}\right) \tag{4.32}
\end{equation*}
$$

Combining (4.21) and (4.32) we get

$$
\begin{equation*}
\int_{0}^{T} \int_{\omega}\left|\varphi_{R}\right|^{2} d x d t \leq C \exp \left(-\frac{1}{2} R^{4 / 3}\right) \tag{4.33}
\end{equation*}
$$

Estimate (4.31) together with (4.33) and the bound (4.19) on $a_{R}$, guarantees (1.7) for small $c>0$, which concludes the proof of Theorem 1.1.

End of the proof when $n$ is odd.
Let $T$ be such that:

$$
T \leq \mu\left(\log \left\|a_{R}\right\|_{\infty}\right)^{-2}\left\|a_{R}\right\|_{\infty}^{-1 / 3}
$$

Note that this implies, with (4.6):

$$
\begin{gather*}
T \leq C \mu(\log R)^{-2}\left\|a_{R}\right\|_{\infty}^{-1 / 3} .  \tag{4.34}\\
40
\end{gather*}
$$

Consider a time $t<T$. By estimates (4.27) and (4.34) we have:

$$
\begin{equation*}
\left\|\xi_{R}(t)\right\|_{\left(L^{2}(\Omega)\right)^{N}}^{2} \leq C \exp \left(-2 R^{4 / 3}+2 C \mu(\log R)^{-2}\left\|a_{R}\right\|_{\infty}^{2 / 3}\right) . \tag{4.35}
\end{equation*}
$$

Therefore, by choosing $\mu>0$ small enough we deduce, using estimate (4.20) on $a_{R}$ :

$$
\begin{equation*}
\left\|\xi_{R}(t)\right\|_{\left(L^{2}(\Omega)\right)^{N}} \leq C \exp \left(-\frac{1}{2} R^{4 / 3}\right) \tag{4.36}
\end{equation*}
$$

Combining (4.25) and (4.36) we deduce that

$$
\begin{equation*}
\left\|\rho_{R}(T)\right\|_{\left(L^{2}(\Omega)\right)^{N}} \leq C \exp \left(-\frac{1}{2} R^{4 / 3}\right) \tag{4.37}
\end{equation*}
$$

As in the even-dimensional case we have, with (4.22) and (4.37):

$$
\begin{equation*}
\left\|\varphi_{R}(T)\right\|_{\left(L^{2}(\Omega)\right)^{N}}^{2} \geq \frac{c_{2}}{2 R^{n}}, \tag{4.38}
\end{equation*}
$$

as $R \rightarrow \infty$.
On the other hand, integrating (4.37) with respect to time we get:

$$
\begin{equation*}
\int_{0}^{T} \int_{\omega}\left|\rho_{R}\right|^{2} d x d t \leq C \exp \left(-\frac{1}{2} R^{4 / 3}\right) \tag{4.39}
\end{equation*}
$$

Combining with (4.21) we get

$$
\begin{equation*}
\int_{0}^{T} \int_{\omega}\left|\varphi_{R}\right|^{2} d x d t \leq C \exp \left(-\frac{1}{2} R^{4 / 3}\right) . \tag{4.40}
\end{equation*}
$$

Note that the bounds (4.20) on $a_{R}$ imply, for large $R$ :

$$
\begin{equation*}
\left(\log \left\|a_{R}\right\|_{\infty}\right)^{-2}\left\|a_{R}\right\|^{2 / 3} \leq C R^{4 / 3} . \tag{4.41}
\end{equation*}
$$

Estimate (4.31) together with (4.40) and on (4.41), guarantees (4.1) for small $c>0$. The proof of Theorem 4.1 is completed.

## 5 Optimality of the observability constant for the wave equation with zero order potential

This section is devoted to prove Theorem 1.2. We start to write a slightly weaker counterpart of Theorem 1.2 in odd dimension.

Theorem 5.1. Assume that $n \geq 3$ is odd and $N \geq 8$. Let $\omega$ be a given open non-empty subset of $\Omega$ such that $\Omega \backslash \bar{\omega} \neq \emptyset$. Then, for all $T>0$ there exist a constant $c>0$, a family of potentials $\left\{a_{R}\right\}_{R>0} \subset L^{\infty}\left(Q ; \mathbb{R}^{N \times N}\right)$ satisfying

$$
\left\|a_{R}\right\|_{\infty} \underset{R \rightarrow+\infty}{ }+\infty
$$

a family of initial data $\left\{\left(w_{R}^{0}, w_{R}^{1}\right)\right\}_{R>0} \subset\left(L^{2}(\Omega)\right)^{N} \times\left(H^{-1}(\Omega)\right)^{N}$ such that the corresponding solution $w_{R}$ of (1.10) satisfies

$$
\begin{equation*}
\lim _{R \rightarrow \infty}\left\{\frac{\left\|w_{R}^{0}\right\|_{\left(L^{2}(\Omega)\right)^{N}}^{2}+\left\|w_{R}^{1}\right\|_{\left(H^{-1}(\Omega)\right)^{N}}^{2}}{\exp \left(c\left(\log \left\|a_{R}\right\|_{\infty}\right)^{-2}\left\|a_{R}\right\|_{\infty}^{2 / 3}\right) \int_{0}^{T} \int_{\omega}\left|w_{R}\right|^{2} d x d t}\right\}=+\infty . \tag{5.1}
\end{equation*}
$$

The proof of Theorem 5.1 is very similar to the one of Theorem 1.2 and we thus leave it to the reader. To prove Theorem 1.2, we argue as in the previous section. As it was observed in that section, one may assume that $N=2$. Consider the solution $u$ and potential $q$ on $\mathbb{R}^{n}$ given by Theorem 3.1. The family $\left\{\psi_{R}\right\}_{R>0}$ as in (4.7) and (4.2) can be viewed as a family of stationary solutions of the Cauchy problem

$$
\begin{cases}\psi_{R, t t}-\Delta \psi_{R}+a_{R} \psi_{R}=0, & \text { in } \mathbb{R}^{n} \times(0, T),  \tag{5.2}\\ \psi_{R}(0)=u_{R}, \quad \psi_{R, t}(0)=0, & \text { in } \mathbb{R}^{n},\end{cases}
$$

with potentials $a_{R}=a_{R}(x)$ as in (4.2). They can also be viewed as solutions of the wave system in the domain $\Omega$ with non-homogeneous boundary conditions:

$$
\begin{cases}\psi_{R, t t}-\Delta \psi_{R}+a_{R} \psi_{R}=0, & \text { in } Q,  \tag{5.3}\\ \psi_{R}=\varepsilon_{R}, & \text { on } \Sigma, \\ \psi_{R}(0)=u_{R}, \quad \psi_{R, t}(0)=0, & \text { in } \Omega,\end{cases}
$$

where

$$
\begin{equation*}
\varepsilon_{R}(x)=u_{R}(x), \quad \forall x \in \Gamma . \tag{5.4}
\end{equation*}
$$

We can assume, without loss of generality, that both $\omega$ and $\Gamma$ are contained in the exterior of the unit ball. Then, (4.15) and (4.16) hold.

We correct the boundary conditions by introducing the weak solutions $\rho_{R}=\rho_{R}(t, x)$ of

$$
\begin{cases}\rho_{R, t t}-\Delta \rho_{R}+a_{R} \rho_{R}=0, & \text { in } Q,  \tag{5.5}\\ \rho_{R}=\varepsilon_{R}, & \text { on } \Sigma, \\ \rho_{R}(0)=\rho_{R, t}(0)=0, & \text { in } \Omega,\end{cases}
$$

and then setting

$$
\begin{equation*}
w_{R}=\psi_{R}-\rho_{R} . \tag{5.6}
\end{equation*}
$$

Clearly $\left\{w_{R}\right\}_{R>0}$ is a family of solutions to the hyperbolic systems of the form (1.10) with potentials $a=a_{R}$ of size $\frac{1}{c_{1}} R^{2} \leq\left\|a_{R}\right\|_{\infty} \leq c_{1} R^{2}$ for some constant $c_{1}>0$.

Let us show that $w_{R}$ satisfies (1.14). We have

$$
\begin{equation*}
w_{R}(0)=u_{R}, \quad w_{R, t}(0)=0, \quad \text { in } \Omega . \tag{5.7}
\end{equation*}
$$

Hence, for some constant $c_{2}>0$

$$
\begin{equation*}
\left\|w_{R}(0)\right\|_{\left(L^{2}(\Omega)\right)^{N}}^{2}+\left\|w_{R, t}(0)\right\|_{\left(H^{-1}(\Omega)\right)^{N}}^{2}=\left\|u_{R}\right\|_{\left(L^{2}(\Omega)\right)^{N}}^{2}=\frac{1}{R^{n}}\|u\|_{\left(L^{2}(R \Omega)\right)^{N}}^{2} \geq \frac{c_{2}}{R^{n}}, \tag{5.8}
\end{equation*}
$$

as $R \rightarrow \infty$.
On the other hand, taking into account that $\omega \subset \mathbb{R}^{n} \backslash B$ and by (4.21), we deduce that

$$
\begin{align*}
\int_{0}^{T} \int_{\omega}\left|w_{R}\right|^{2} d x d t & \leq 2 \int_{0}^{T} \int_{\omega}\left|\psi_{R}\right|^{2} d x d t+2 \int_{0}^{T} \int_{\omega}\left|\rho_{R}\right|^{2} d x d t  \tag{5.9}\\
& \leq 2 C \exp \left(-2 R^{4 / 3}\right)+2 \int_{0}^{T} \int_{\omega}\left|\rho_{R}\right|^{2} d x d t
\end{align*}
$$

Consequently, in order to conclude that (1.14) holds, it is sufficient to get an upper bound of the form

$$
\begin{equation*}
\int_{0}^{T} \int_{\omega}\left|\rho_{R}\right|^{2} d x d t \leq C \exp \left(-c_{0} R^{4 / 3}\right) \tag{5.10}
\end{equation*}
$$

for suitable constants $C>0$ and $c_{0}>0$.
The rest of the proof is devoted to proving (5.10).
The solution $\rho_{R}$ of (5.5) is defined by transposition. More precisely, consider the adjoint problem

$$
\begin{cases}\theta_{t t}-\Delta \theta+a_{R} \theta=f, & \text { in } Q,  \tag{5.11}\\ \theta=0, & \text { on } \Sigma, \\ \theta(T)=\theta_{t}(T)=0, & \text { in } \Omega .\end{cases}
$$

Then, multiplying the first equation in (5.5) by $\theta$, integrating it in $Q$ and using formal integration by parts we get

$$
\begin{equation*}
\int_{Q} f \cdot \rho_{R} d x d t=-\int_{\Sigma} \frac{\partial \theta}{\partial \nu} \cdot \varepsilon_{R} d \sigma d t \tag{5.12}
\end{equation*}
$$

Here, $\cdot$ means the usual scalar product in $\mathbb{R}^{N}$. We adopt (5.12) as definition of solution $\rho_{R}$ of (5.5) in the sense of transposition.

It is well known that there exists $C=C(R)$ such that

$$
\begin{equation*}
\left\|\rho_{R}\right\|_{\left(L^{\infty}\left(0, T ; L^{2}(\Omega)\right)\right)^{N}} \leq C(R)\left\|\varepsilon_{R}\right\|_{\left(L^{2}(\Sigma)\right)^{N}} . \tag{5.13}
\end{equation*}
$$

This is a consequence, by duality, of the following hidden regularity property for solutions of (5.11):

$$
\begin{equation*}
\left\|\frac{\partial \theta}{\partial \nu}\right\|_{\left(L^{2}(\Sigma)\right)^{N}} \leq C(R)\|f\|_{\left(L^{1}\left(0, T ; L^{2}(\Omega)\right)\right)^{N}} \tag{5.14}
\end{equation*}
$$

It remains to show that (5.14) holds with a constant $C(R)$ such that

$$
\begin{equation*}
C(R) \leq C \exp \left(C R^{4 / 3}\right) \tag{5.15}
\end{equation*}
$$

First of all, by an energy estimate, we observe that

$$
\begin{equation*}
\|\theta\|_{\left(L^{\infty}\left(0, T ; H_{0}^{1}(\Omega)\right)\right)^{N}}+\left\|\theta_{t}\right\|_{\left(L^{\infty}\left(0, T ; L^{2}(\Omega)\right)\right)^{N}} \leq C \exp (C R)\|f\|_{\left(L^{1}\left(0, T ; L^{2}(\Omega)\right)\right)^{N}} \tag{5.16}
\end{equation*}
$$

However, to prove this, we can not simply apply the Gronwall inequality to the energy

$$
\begin{equation*}
E(t)=\frac{1}{2}\left[\|\theta(t)\|_{\left(H_{0}^{1}(\Omega)\right)^{N}}^{2}+\left\|\theta_{t}(t)\right\|_{\left(L^{2}(\Omega)\right)^{N}}^{2}\right] . \tag{5.17}
\end{equation*}
$$

This would yield a constant in the inequality of the order of $\exp \left(C R^{2}\right)$ since the norm $\left\|a_{R}\right\|_{\infty}$ of the potential is of the order of $R^{2}$. To improve this estimate and get a constant of the order of $\exp (C R)$ we have to work, rather, with the modified energy

$$
\begin{equation*}
E_{R}(t)=\frac{1}{2}\left[\|\theta(t)\|_{\left(H_{0}^{1}(\Omega)\right)^{N}}^{2}+\left\|\theta_{t}(t)\right\|_{\left(L^{2}(\Omega)\right)^{N}}^{2}+R^{2}\|\theta(t)\|_{\left(L^{2}(\Omega)\right)^{N}}^{2}\right], \tag{5.18}
\end{equation*}
$$

as in [30]. We have:

$$
\begin{align*}
\left|\frac{d E_{R}(t)}{d t}\right| & =\left|\int_{\Omega} f \cdot \theta_{t} d x-\int_{\Omega}\left(a_{R} \theta\right) \cdot \theta_{t} d x\right|  \tag{5.19}\\
& \leq\|f(t)\|_{\left(L^{2}(\Omega)\right)^{N}}\left\|\theta_{t}(t)\right\|_{\left(L^{2}(\Omega)\right)^{N}}+C R^{2}\|\theta(t)\|_{\left(L^{2}(\Omega)\right)^{N}}\left\|\theta_{t}(t)\right\|_{\left(L^{2}(\Omega)\right)^{N}} \\
& \leq C R E_{R}(t)+\|f(t)\|_{\left(L^{2}(\Omega)\right)^{N}} \sqrt{E_{R}(t)} .
\end{align*}
$$

Solving this differential inequality and taking into account $E_{R}(T)=0$ (because the data of $\theta$ at $t=T$ vanish) we easily get (5.16).

Once (5.16) is proved we can show (5.14) using the Rellich multiplier as in [17]. It follows that (5.14) holds with a constant $C(R)$ of the order of $C \exp (C R)$ which, clearly, fulfills (5.15).

This concludes the proof of Theorem 1.2.

## 6 Equations with convective terms

In Section 4 and 5 we have discussed equations with potentials entering in the zero order term. In Theorem 3.1 we have also considered an equation with a convective potential, of the form:

$$
\Delta u+q_{1} \cdot \nabla u=0, \quad x \in \mathbb{R}^{2} .
$$

with:

$$
|u| \leq C \exp \left(-|x|^{4 / 3}\right), \quad\left|q_{1}\right| \leq C(1+|x|)^{-1 / 3} .
$$

This allows extending the optimality construction of the previous sections to heat equations or systems with convective potentials of the form:

$$
\begin{cases}\varphi_{t}-\Delta \varphi+a_{1} \cdot \nabla \varphi=0, & \text { in } Q  \tag{6.1}\\ \varphi=0, & \text { on } \Sigma .\end{cases}
$$

In this case the observability inequality is known to hold with an observability constant of the order of $\exp \left(C\left\|a_{1}\right\|_{\infty}^{2}\right)$ (see Theorem 2.1 of this paper, or Theorem 2.3 in [9]).

The construction in Theorem 3.1 allows arguing as above, by scaling and cut-off, and to show that this observability estimate is sharp in what concerns the exponentially quadratic growth of the observability constant on the potential.

Let us briefly check why this is the case: We set,

$$
u_{R}(x)=\binom{\operatorname{Re} u(R x)}{\operatorname{Im} u(R x)}, \quad b_{R}(x)=-R\left(\begin{array}{cc}
\operatorname{Re} q_{1}(R x) & -\operatorname{Im} q_{1}(R x)  \tag{6.2}\\
\operatorname{Im} q_{1}(R x) & \operatorname{Re} q_{1}(R x)
\end{array}\right)
$$

Then,

$$
\begin{equation*}
\Delta u_{R}=b_{R}(x) \cdot \nabla u_{R}, \quad \text { in } \mathbb{R}^{2}, \tag{6.3}
\end{equation*}
$$

and

$$
\begin{equation*}
\left|u_{R}(x)\right| \leq C \exp \left(-c_{*} R^{4 / 3}|x|^{4 / 3}\right), \quad \text { in } \mathbb{R}^{2} \tag{6.4}
\end{equation*}
$$

Moreover,

$$
\begin{equation*}
\left|b_{R}(x)\right| \leq C R(1+R|x|)^{-1 / 3} \leq C R^{2 / 3}, \quad \forall x:|x| \leq 1 . \tag{6.5}
\end{equation*}
$$

Combining (6.4) and (6.5) one easily observes that the observability constant must necessarily grow exponentially on the square of the $L^{\infty}$-norm of the potential. The functions $u_{R}$ are defined on all $\mathbb{R}^{2}$. By restriction arguments, as in Step 2 of Section 4, one can prove the optimality of the known observability inequality for heat equation with convective potentials on a bounded open subset $\Omega$ of $\mathbb{R}^{2}$.

Note however that, in this case, due to the fact that the construction only holds in dimension 2, very little is known about the sharpness of the estimates in other dimensions. Unfortunately, the argument of separation of variables that allows to deduce, from any even $n$, $n$-dimensional examples from 2-dimensional examples does not work in the convective case. Indeed, in this argument, we would lose the decay at infinity of $q_{1}$ at the speed $|x|^{-1 / 3}$ that is necessary to show the sharpness of the estimate. Thus we do not know anything about the optimality in even dimension $n \geq 4$. As for odd space-dimension, adapting the proof of Theorem 3.2 to a convective potential one may give a weaker, vectorial form of the optimality result in dimension $n=3$. Again, we do not know anything about the case of superior odd dimensions.

## 7 Connections with the controllability of semilinear equations

As we said in the introduction the problem of the explicit dependence of the observability constant on the size of the potentials entering in the system is also closely related with the problem of controllability of semilinear equations. Indeed, the controllability of semilinear equations is usually obtained by a fixed point argument. The growth condition one needs to impose at infinity to the nonlinearity depends very much on the cost of controlling the linearized equation perturbed by a potential, which is precisely given by the observability constant. Consequently the growth of the observability constant on the potential and the growth condition on the nonlinearity for controllability to hold are closely connected. In this section we briefly describe the problem of controllability of semilinear equations, and some open problems related to the optimality results of the previous sections. We first address the semilinear heat equation to later consider the wave equation.

### 7.1 The semilinear heat equation

Consider the semilinear heat equation

$$
\begin{cases}y_{t}-\Delta y+f(y)=v \chi_{\omega}, & \text { in } Q,  \tag{7.1}\\ y=0, & \text { on } \Sigma, \\ y(0, x)=y_{0}(x), & \text { in } \Omega .\end{cases}
$$

Here $\omega$ is an open non-empty subset $\omega$ of $\Omega$, and we denote by $\chi_{\omega}$ the characteristic function of $\omega$. To make the problem non-trivial we also assume that $\Omega \backslash \bar{\omega} \neq \emptyset$ throughout this section.

We discuss the problem of null-controllability and more precisely whether for all $T>0$ and $y_{0} \in L^{2}(\Omega)$ there exists a control $v \in L^{2}(Q)$ such that the solution $y=y(t, x)$ of (7.1) satisfies

$$
\begin{equation*}
y(x, T) \equiv 0, \quad \text { in } \Omega \tag{7.2}
\end{equation*}
$$

In [12] the following result was proved:
Theorem C ([12]). Assume that the nonlinearity $f \in C^{1}(\mathbb{R})$ is such that $f(0)=0$ and

$$
\begin{equation*}
\limsup _{|s| \rightarrow \infty} \frac{\left|f^{\prime}(s)\right|}{\log ^{3 / 2}(s)}=0 . \tag{7.3}
\end{equation*}
$$

Then, whatever the open non-empty subset $\omega$ of $\Omega$ is, and for all $T>0$, system (7.2) is nullcontrollable.

In particular, Theorem C guarantees the possibility of controlling some blowing-up equations. Indeed, when $f$ is of the form

$$
\begin{equation*}
f(s)=-s \log ^{r}(1+|s|) \tag{7.4}
\end{equation*}
$$

with $r>1$ solutions of (7.1), in the absence of control, i.e. with $v \equiv 0$, blow-up in finite time occurs. According to Theorem C the process can be controlled, and, in particular, the blow-up be avoided when $1<r \leq 3 / 2$.

The growth condition in (7.3) is intimately related with the observability inequality (1.2). Indeed, the logarithmic function in (7.3) is precisely the inverse of the exponential one in (1.2).

The proof in [12] relies on the by now classical argument of fixed point consisting on linearizing the equation and estimating the cost of the control in terms of the size of the potential entering in the system.

Theorem C was proved for scalar semilinear heat equations, but the same technique applies with the same conclusions for semilinear systems.

According to the results in Section 4 the estimate (1.2) can not be improved, and the fixed point argument in [12] may not lead to any improvement of the growth condition (7.3) for controllability.

By the contrary, in [12] it was proved that there are some nonlinearities $f$ satisfying

$$
\begin{equation*}
\limsup _{|s| \rightarrow \infty} \frac{|f(s)|}{s \log ^{r}(1+|s|)}=0 \tag{7.5}
\end{equation*}
$$

with a growth rate of the order for $r>2$ for which controllability fails because the control may not avoid blow-up to occur.

Whether controllability occurs for nonlinearities with a growth rate of the order of (7.5) with $3 / 2 \leq r \leq 2$ is an open problem. However, in the light of the optimality of the explicit observability estimate we can guarantee that the method in [12] does not allow to cover the case $r>3 / 2$.

Similar questions arise for nonlinearities depending on the gradient of the state considered in [9].

### 7.2 The semilinear wave equation

The problem of controllability of the semilinear wave equation and the existing results were recently discussed in [29]. Let us recall the main known results.

Consider the semilinear wave equation:

$$
\begin{cases}y_{t t}-\Delta y+f(y)=v \chi_{\omega}, & \text { in } Q,  \tag{7.6}\\ y=0, & \text { on } \Sigma, \\ y(0, x)=y_{0}(x), \quad y_{t}(0, x)=y_{0}(x), & \text { in } \Omega .\end{cases}
$$

In one space dimension (i.e., $n=1$ ), roughly speaking, controllability is known to hold provided the nonlinearity grows as

$$
\begin{equation*}
|f(s)| \leq C|s| \log ^{2}|s|, \quad \text { as } \quad|s| \rightarrow \infty \tag{7.7}
\end{equation*}
$$

(see [29], [30], [7]). This result is sharp in the sense that blow-up is known to occur for some nonlinearities growing as (7.4) with $r>2$ and, due to the finite speed of propagation, controllability may not hold if blow-up occurs. According to this, to some extent, the picture is complete in one space dimension.

In the multi-dimensional case, i.e., $n>1$, less is known. Once again, due to blow-up, one can not expect controllability to hold for nonlinearities of the form (7.4) with $r>2$. But, in principle, one could expect the system to be controllable for $r \leq 2$. However, in view of the optimality result of Theorem 1.2 on the growth of the observability constant, one may not expect the fixed point methods in [30] to apply for $r>3 / 2$. Thus the problem is completely open for $3 / 2<r \leq 2$.

As we mentioned above, Theorem 2.2 suggests that controllability should hold for nonlinearities satisfying the growth condition (7.3). However, even this is an open problem. Indeed, controllability is only known under the more restrictive condition (see [14])

$$
\begin{equation*}
\lim _{|s| \rightarrow \infty} \frac{|f(s)|}{s \log ^{1 / 2}|s|}=0 \tag{7.8}
\end{equation*}
$$

This is so because of the lack of smoothing effect of the wave equation. Finite energy solutions belong to $y \in C\left([0, T] ; H_{0}^{1}(\Omega)\right)$ but fail to be bounded in dimensions $n \geq 2$ and this an obstacle to apply the fixed point argument for nonlinearities (7.5) with $r=3 / 2$. Note also that, by using instead the sharp observability internal observability in part $i$ if of Theorem 2.2, similar to [30] and [14], it is easy to show the exact controllability of system (7.6) with nonlinearities (7.5) up to the exponent $r<3 / 2$.

In view of this, in the application of the fixed point argument for control, one is not allowed to assume the potential to be bounded and consequently one has to work with potentials in some $L^{\infty}\left(0, T ; L^{s}(\Omega)\right)$ with $s<\infty$ in which case the observability constant has a faster growth, as we have seen in Theorem 2.2. Consequently, the controllability of the semilinear hyperbolic equations requires a stronger growth condition (7.5) (with $r<3 / 2$, compared with the parabolic one in which $r=3 / 2$ is allowed). Accordingly, the controllability for nonlinearities of the form (7.4) with $3 / 2 \leq r \leq 2$ is an open problem in the multi-dimensional hyperbolic case.

## 8 Open problems

A lot of problems remain open in this field. Some of them could need important new ideas and developments.

### 8.1 Meshkov's construction for scalar equations

One of the very first open problems in this frame is whether for scalar multidimensional equations one may construct solutions of (1.8) decaying super-exponentially. The results in [20], based on Carleman inequalities, show that the decay rate $\exp \left(-C|x|^{4 / 3}\right)$ is critical in the sense that all solutions that decay faster as $|x| \rightarrow \infty$ do necessarily vanish. But the construction we have recalled in Theorem A showing that there is a non-trivial solution decaying as $\exp \left(-C|x|^{4 / 3}\right)$ is only available for a system of two real equations. Whether such a construction can be adapted to scalar equations is an interesting and very likely difficult open problem.

Obviously this is connected with the possibility of extending the optimality results in Theorems 1.1 and 1.2 to scalar equations.

The construction in [27] of a nonzero real function vanishing to infinite order at 0 and satisfying a critical differential inequality may be a first step in this direction, at least in the case of convective potentials.

### 8.2 Meshkov's construction in dimension $n=3$

As we mentioned above the construction of Meshkov is valid for space dimension $n=2$ and, by separation of variables, can be easily extended to any even dimension. In Section 3, we gave a weaker version of this construction in dimension 3 , with a $\mathbb{C}^{4}$-valued solution $u$ and a potential $q$ that is not bounded, but grows only logarithmically. This could be adapted to get a bounded potential $q$ on $\mathbb{R}^{3}$, and a solution $u$ decaying almost like $e^{-|x|^{4 / 3}}$ (that is, faster than any $e^{-|x|^{\alpha}}$ with $\alpha<4 / 3)$. The questions whether in odd dimensions, one can exactly get the optimal decay $e^{-|x|^{4 / 3}}$, or replace the $\mathbb{C}^{4}$-valued solution $u$ by a complex-valued function (as in the original construction of Meshkov) remains open. Indeed it is not clear whether the distinction between even and odd space dimensions is purely technical or not. Note that there is no distinction in what concerns the observability inequality.

The general question of the fastest speed of decay at infinity for eigenfunctions of the Laplace operator with a potential, which includes open problems 8.1 and 8.2, does not seem to have been studied intensively. In [25], the author uses Carleman-type inequalities to give some lower bound on the eigenfunctions. In [4] and [6], some lower bounds are given, but with a strong positivity assumption on the potential. We refer to the book by Agmon [1] for a study of the exponential decay of the eigenfunctions.

### 8.3 Sharp observability and time-dependent potentials

The constructions of Sections 4 and 5 are based on time-independent functions, so that both the heat and wave equation are dealt with as perturbations of an elliptic equation. One possible way to improve the optimality results would be to consider time-dependent variants of the Meshkov's construction. Note that in [19], similar constructions are done for evolution equations, but with the purpose of getting fast-decaying solutions of these equations for large times, which is not what we need in our context.

## 8.4 $\quad L^{r}$-potentials

In Section 2 an observabillity inequality is shown for potentials that are in $L^{\infty}\left(0, T ; L^{p}(\Omega)\right), n \leq$ $p \leq \infty$. In Sections 4 and 5 , we have shown that these estimates are sharp only in the case $p=\infty$. The question remains open for finite $p$. As it was already noticed (see Remark 3.2), it is not sufficient to construct optimally decaying solutions $u$ of the elliptic equation (3.1) on $\mathbb{R}^{n}$ with a potential $q$ in $L^{p}\left(\mathbb{R}^{n}\right)$ to solve this problem.

## $8.51-d$ problems

The construction by Meshkov is impossible in one space dimension. In that case solutions of equations of the form (1.8) may not decay in a super-exponential way.

On the other hand, for heat equations, the observability estimates one obtains in $1-d$ and several space dimensions are the same. Accordingly whether the observability estimate (1.2) is sharp is an open problem.

As we mentioned in Section 7, in the context of the wave equation, using sidewise energy estimates one can obtain a much better estimate with an observability constant of the order of $\exp \left(C\|a\|_{\infty}^{1 / 2}\right)$. This estimate is sharp even for constant potentials as can be easily seen from the Fourier representation of solutions and the analysis of how the spectral gap behaves as the constant potential tends to infinity (see [10]).

When the potentials under consideration only depend on $x$ one can use well known transformations (see [20] and [22]) from wave into heat processes to obtain estimates on the observability constant of the order of $\exp \left(C\|a\|_{\infty}^{1 / 2}\right)$.

Whether the estimate (1.2) is sharp for $1-d$ heat equations with bounded potentials depending both on $x$ and $t$ is an open problem.

### 8.6 Other equations

In this article we have addressed the heat and wave equations. But the same problems arise for other equations like Schrödinger, plate and KdV equations. We refer to [31], [28] and [21] for a discussion of observability estimates for these models.

## References

[1] S. Agmon, Lectures on exponential decay of solutions of second-order elliptic equations: bounds on eigenfunctions of $N$-body Schrödinger operators, Princeton Univ. Press, Princeton, NJ, 1982.
[2] S. Alinhac, M. S. Baouendi, A nonuniqueness result for operators of principal type, Math. Z., 220 (1995), 561-568.
[3] S. Alinhac, N. Lerner, Unicité forte à partir d'une variété de dimension quelconque pour des inégalités différentielles elliptiques, Duke Math. J., 48 (1981), 49-68.
[4] B. Alziary, P. Takáč, A pointwise lower bound for positive solutions of a Schrödinger equation in $\mathbb{R}^{N}$, J. Differential Equations, 133 (1997), 280-295.
[5] C. Bardos, G. Lebeau, J. Rauch, Sharp sufficient conditions for the observation, control and stabilization of waves from the boundary, SIAM J. Control Optim., 30 (1992), 10241065.
[6] C. Bardos, M. Merigot, Asymptotic decay of the solution of a second-order elliptic equation in an unbounded domain. Applications to the spectral properties of a Hamiltonian, Proc. Roy. Soc. Edinburgh Sect. A, 76 (1977), 323-344.
[7] P. Cannarsa, V. Komornik, P. Loreti, One-sided and internal controllability of semilinear wave equations with infinitely iterated logarithms, Discrete Contin. Dyn. Syst. B, 8 (2002), 745-756.
[8] C. Castro, E. Zuazua, Concentration and lack of observability of waves in highly heterogeneous media, Arch. Ration. Mech. Anal., 164 (2002), 39-72.
[9] A. Doubova, E. Fernández-Cara, M. González-Burgos, E. Zuazua, On the controllability of parabolic systems with a nonlinear term involving the state and the gradient, SIAM J. Control Optim., 41 (2002), 798-819.
[10] H. O. Fattorini, Estimates for sequences biorthogonal to certain complex exponentials and boundary control of the wave equation, in New trends in systems analysis (Proc. Internat. Sympos., Versailles, 1976), Lecture Notes in Control and Informat. Sci., Vol. 2, Springer, Berlin, 1977, 111-124.
[11] E. Fernández-Cara, E. Zuazua, The cost of approximate controllability for heat equations: the linear case, Adv. Differential Equations, 5 (2000), 465-514.
[12] E. Fernández-Cara, E. Zuazua, Null and approximate controllability for weakly blowing-up semilinear heat equations, Ann. Inst. H. Poincaré Anal. Non Linéaire, 17 (2000), 583-616.
[13] A. V. Fursikov, O. Yu. Imanuvilov, Controllability of Evolution Equations, Lecture Notes Series, 34, Research Institute of Mathematics, Seoul National University, Seoul, Korea, 1994.
[14] X. Fu, J. Yong, X. Zhang, Exact controllability for multidimensional semilinear equations, Preprint, 2004.
[15] L. Hörmander. Non-uniqueness for the Cauchy problem, in Fourier integral operators and partial differential equations. Colloque International, Nice, edited by J. Chazarain. Lecture Notes in Mathematics, 459, Springer-Verlag, Berlin-New York, 1974, 36-72.
[16] O. Yu. Imanuvilov, On Carlerman estimates for hyperbolic equations, Asymptot. Anal., 32 (2002), 185-220.
[17] J. L. Lions, Contrôlabilité exacte, stabilisation et perturbations de systèmes distribués, Tomes $1 \& 2$, RMA $8 \& \mathbf{9}$, Paris, 1988.
[18] A. López, X. Zhang, E. Zuazua, Null controllability of the heat equation as singular limit of the exact controllability of dissipative wave equations, J. Math. Pures Appl., 79 (2000), 741-808.
[19] V. Z. Meshkov, On the possible rate of decrease at infinity of the solutions of second-order partial differential equations, (Russian), Mat. Sb., 182 (1991), 364-383; translation in Math. USSR-Sb., 72 (1992), 343-361.
[20] L. Miller, Geometric bounds on the growth rate of null-controllability cost of the heat equation in small time, J. Differential Equations, 204 (2004), 202-226.
[21] L. Rosier, Exact boundary controllability for the Korteweg-de Vries equation on a bounded domain, ESAIM: Control Optim. Calc. Var., 2 (1997), 33-55.
[22] D. L. Russell, A unified boundary controllability theory for hyperbolic and parabolic partial differential equations, Studies in Appl. Math., 52 (1973), 189-221.
[23] T. Seidman, S. Avdonin, S. A. Ivanonv, The "window problem" for series of complex exponentials, J. Fourier Anal. Appl., 6 (2000), 233-254.
[24] W. A. Strauss, Partial differential equations, Wiley, New York, 1992.
[25] J. Uchiyama, Lower bounds of decay order of eigenfunctions of second-order elliptic operators, Publ. Res. Inst. Math. Sci., 21 (1985), 1281-1297.
[26] E. T. Whittaker, G. N. Watson, A course of modern analysis, Reprint of the fourth edition (1927), Cambridge Univ. Press, Cambridge, 1996.
[27] T. Wolff, A counterexample in a unique continuation problem, Comm. Anal. Geom. 2 (1994), 79-102.
[28] X. Zhang, Exact controllability of semilinear plate equations, Asymptot. Anal., 27 (2001), 95-125.
[29] X. Zhang, E. Zuazua, Exact controllability of the semi-linear wave equation, in Sixty Open Problems in the Mathematics of Systems and Control, edited by V. D. Blondel and A. Megretski, Princeton University Press, 2004, 173-178.
[30] E. Zuazua, Exact controllability for semilinear wave equations in one space dimension, Ann. Inst. H. Poincaré Anal. Non Linéaire, 10 (1993), 109-129.
[31] E. Zuazua, Remarks on the controllability of the Schrödinger Equation, in Quantum Control: mathematical and numerical challenges, edited by A. Bandrauk, M. C. Delfour, and C. Le Bris, CRM Proc. Lect. Notes Ser., 33, AMS Publications, Providence, R.I., 2003, 181-199.
[32] C. Zuily, Uniqueness and nonuniqueness in the Cauchy problem. Progress in Mathematics, 33. Birkhäuser Boston, Inc., Boston, MA, 1983.


[^0]:    *This work is supported by the Grant MTM2005-00714 of the Spanish MEC, the DOMINO Project CIT-370200-2005-10 in the PROFIT program of the MEC (Spain), the SIMUMAT projet of the CAM (Spain) and the NSF of China under grants 10371084 and 10525105. EU TMR Projects "Smart Systems" and "HYKE" financed the first author's visit to the Universidad Autónoma de Madrid.
    ${ }^{\dagger}$ Département de mathématiques, Université de Cergy-Pontoise, Site de Saint Martin, 2 avenue Adolphe Chauvin 95302 Cergy-Pontoise, Cedex, France. e-mail: thomas.duyckaerts@u-cergy.fr.
    ${ }^{\ddagger}$ Key Laboratory of Systems and Control, Academy of Mathematics and Systems Sciences, Academia Sinica, Beijing 100080, China; and Departamento de Matemáticas, Facultad de Ciencias, Universidad Autónoma de Madrid, 28049 Madrid, Spain. e-mail: xuzhang@amss.ac.cn.
    ${ }^{\S}$ Departamento de Matemáticas, Facultad de Ciencias, Universidad Autónoma de Madrid, 28049 Madrid, Spain. $e$-mail: enrique.zuazua@uam.es.

[^1]:    ${ }^{1}$ This means that $T \leq C\|a\|_{\infty}^{-1 / 3}$ for some generic constant $C>0$. We shall use a similar notation later, for example for $T \gtrsim\|a\|_{\infty}^{-1 / 3}$.

