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Mathematical analysis of an observer for ² solving inverse source wave problem

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38 Abstract

The objective of this work is to propose a method using observers to estimate a sourceterm of a wave equation from internal measurements in a subdomain. The first part of

the work consists in proving an identifiability result from classical observability

- 42 conditions for wave equations. We show that the source reconstruction is an ill-posed
- inverse problem of degree 1 or 2 depending on the measurements type. This inverse
 problem is solved using observers a sequential strategy that is proven to be
- ⁴⁴ problem is solved using observers a sequential strategy that is proven to be ⁴⁵ equivalent to a minimization of a cost functional with Tikhonov regularization.
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47 1 Introduction

46

Inverse problems applied to wave-like propagation phenomena has received a lot of 48 attention with numerous applications in geophysics [35, 16], tomography [38, 44, 32]. 49 medical imaging [37, 56] with the development of elastography [45, 52], non-destructive 50 testing [23]. In fact, several subcategories of inverse problems can be distinguished: 51 the reconstruction of initial conditions, the reconstruction of sources, and finally the 52 identification of constitutive parameters. First, the reconstruction of initial conditions 53 has been widely studied [39, 52] and is based on the proof of a so-called observability 54 inequality. This inequality, which has been studied extensively by the control community 55 [14, 27, 30] quantifies the initial condition to measurement operator invertibility which 56 in this case is linear, allowing for the analysis of reconstruction strategies [47, 55]. 57 At the other end of the spectrum, parameter reconstruction is a nonlinear inversion 58 problem [28, 53], where identifiability and stability results are hard to obtain in the 59 general cases and often rely on Carleman inequalities, see for example [43, 49, 54] and 60 references therein. Finally, source reconstruction is the intermediate step between initial 61 conditions reconstruction and parameter reconstruction. In a first step, we consider a 62 source term with a separate variable, *i.e.*, the source is the product of a time-dependent 63 function multiplying a space-dependent function to be estimated. This source term 64 can take various forms, even being point sources [24, 26]. One of the advantages of 65 estimating such source terms is that one can show identifiability by relying on initial 66 condition observability inequalities [34, 17], while more general source terms may require 67 the use of dedicated Carleman estimates [48, 21]. From a practical point of view, there 68 are two main categories of methods for solving these inverse problems: Least-squares 69 minimization approaches [38, 46, 55] - also known as variational approaches or 4D-var 70 in the data assimilation community [10, 12, 22] – or sequential approaches, where a 71 new evolution problem integrating the measurements as a feedback pursues the target 72 solution. The least squares approach is certainly the oldest strategy [1, 4] and is 73 typically solved using adjoint equations [11]. The goal is to eventually find the optimal 74 trajectory that minimizes the discrepancy between the solution and the observations. 75 The sequential approach, also called the observer approach, is less studied. The 76 principle is to introduce feedback into the dynamics based on measurement to create 77 a new system, called the observer, which converges asymptotically to the observed 78 system over time. Starting from the reconstruction of the initial conditions of wave 79 problems, algorithms were developed [39] combining the observer and time-reversal 80 strategies in a comparable spirit to the back-and-forth nudging introduced in data 81 assimilation [31]. This strategy has been extended to unbounded domains by [52]. 82 The proposed feedbacks are often easy to compute, and their efficiency is based on 83 stabilization principles inspired by the so-called Luenberger observer introduced for 84 finite-dimensional control problems [2, 6]. In the same community, it is also known that 85 the principles of dynamic programming [3] can also help to formulate observers that 86 asymptotically solve the least square minimization problem. Such a strategy is often 87 referred to as optimal filtering, since the feedback is also called a filter, and it is derived 88 directly from an optimal criterion. It has been extended to the estimation of initial 89

condition for infinite dimensional systems [5, 7] with recent works on numerical methods
and their analysis, see e.g. [57] and references therein. Observers for reconstructing
sources or identifying parameters have also been introduced in finite dimensional cases
[25] with applications to discretized wave-like problems [33, 56].

Our goal is to extend the definition of such an optimal observer to source reconstruction in the infinite dimensional context of source inversion for wave problem. Our contribution includes the formulation of the observer, namely its existence, but also its asymptotic property, here called stability, and its robustness to measurement noise, as in all inverse problems. Notably, we base our analysis on general field measurements on a subdomain whereas most observability results for wave equations are classically based on velocity measurements.

This work is organized as follows. In the first section, we present the problem 101 underlying the nature of the present observations. In the second section, we present 102 how observability inequality allows us to quantify how poor the reconstruction of source 103 terms is from the available field measurements. The next section is devoted to the 104 formulation and analysis of the observer in the infinite dimensional context, developing 105 an infinite dimensional Kalman filtering reduced to the source space. We then propose 106 a discretization strategy based on the same principles at the discrete level. Finally, we 107 illustrate the strategy with numerical experiments. 108

¹⁰⁹ 2 Problem setting

Let Ω be a bounded, connected open domain of class C^2 in \mathbb{R}^d . We consider the following wave equation

$$\begin{cases} \partial_{tt}u(x,t) - \Delta u(x,t) = \sigma(t)\theta(x), & (x,t) \in \Omega \times (0,T), \\ u(x,t) = 0, & (x,t) \in \partial\Omega \times (0,T), \\ u(x,0) = u_0(x), & x \in \Omega, \\ \partial_t u(x,0) = v_0(x), & x \in \Omega, \end{cases}$$
(1)

where (u_0, v_0) is the initial condition, $\sigma \in L^2(0, T)$ is known and $\theta \in \mathcal{P} \subseteq L^2(\Omega)$, a Hilbert space, to be recovered from available measurements.

¹¹⁴ 2.1 Semigroup approach

We rewrite (1) as a first order system in the state-space $\mathcal{Z} := H_0^1(\Omega) \times L^2(\Omega)$ equipped with the semi-norm H^1 and the L^2 standard scalar product,

$$\begin{cases} \dot{z}(t) = Az(t) + B(t)\theta, & t \in [0, T], \\ z(0) = z_0, \end{cases}$$
(2)

117 where

$$z_0 = \begin{pmatrix} u_0 \\ v_0 \end{pmatrix}, \quad B = \begin{pmatrix} 0 \\ \sigma(t) \mathrm{Id} \end{pmatrix}, \tag{3}$$

and the operator A is an unbounded skew-adjoint operator from $\mathcal{D}(A) \subset \mathcal{Z}$ into \mathcal{Z} defined by

$$A = \begin{pmatrix} 0 & \mathrm{Id} \\ \Delta_0 & 0 \end{pmatrix}, \quad \mathcal{D}(A) = \mathcal{D}(-\Delta_0) \times \mathrm{H}^1_0(\Omega),$$

where Δ_0 denotes the Laplacian operator with homogeneous Dirichlet boundary conditions [29, Part II, Section 1.2.10]. The operator A is maximal dissipative, hence it is a generator of a C^0 -semigroup $(\Phi(t))_{t\geq 0}$ on \mathcal{Z} also denoted $(e^{tA})_{t\geq 0}$. Note that $(B(t))_{t\geq 0}$ is a one parameter family of bounded operator from \mathcal{P} to \mathcal{Z} . We recall the various notions of solution and associated regularity.

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Theorem 2.1. For all T > 0, the evolution equation (2) admits the following solutions: 124 1. For all $z_0 \in \mathcal{D}(A)$, $\sigma \in \mathrm{H}^1([0,T])$ and $\theta \in \mathcal{P}$, there exists a unique strong solution 125 $z \in \mathrm{C}^0([0,T]; \mathcal{D}(A)) \cap \mathrm{C}^1([0,T]; \mathcal{Z})$ of (2) and (3).

127 2. For all $z_0 \in \mathbb{Z}$, $\sigma \in L^1(0,T)$ and $\theta \in \mathcal{P}$, there exists a unique mild-solution 128 $z \in C^0([0,T]; \mathbb{Z})$ of (2) and (3) given by Duhamel's formula

$$z(t) = e^{tA}z_0 + \int_0^t e^{(t-s)A} \begin{pmatrix} 0\\\sigma(s) \end{pmatrix} \theta \,\mathrm{d}s, \quad t \in [0,T].$$

$$\tag{4}$$

129 Moreover, the solution is a weak solution, namely $z \in L^p((0,T); \mathbb{Z})$, for all 130 $v \in D(A^*), v, z(\cdot) \in W^{1,p}(0,T)$ and

$$\begin{cases} \frac{\mathrm{d}}{\mathrm{d}t}(v, z(t)) = (A^*v, z(t)) + (B(t)\sigma, v), & t \in [0, T] \ a.e., \\ z(0) = z_0. \end{cases}$$

¹³¹ *Proof.* See [29, Part II, Chapter 1, Section 3, Propositions 3.1 and 3.2].

132 Remark 2.2. From Theorem 2.1, it follows naturally:

133 1. For all $(u_0, v_0) \in \mathcal{D}(-\Delta_0) \times \mathrm{H}^1_0(\Omega)$, $\sigma \in \mathrm{H}^1([0,T])$ and $\theta \in \mathcal{P}$, there exists a 134 unique strong solution u of (1) belonging to

$$C^{0}([0,T]; \mathcal{D}(-\Delta_{0})) \cap C^{1}([0,T]; H^{1}_{0}(\Omega)) \cap C^{2}([0,T]; L^{2}(\Omega)).$$

¹³⁵ 2. For all $(u_0, v_0) \in H^1_0(\Omega) \times L^2(\Omega)$, $\sigma \in L^1(0, T)$ and $\theta \in \mathcal{P}$, there exists a unique ¹³⁶ mild-solution u of (1) belonging to

$$C^{0}([0,T]; H^{1}_{0}(\Omega)) \cap C^{1}([0,T]; L^{2}(\Omega)).$$

We introduce the so-called *sensitivity* operators, namely the one parameter family (L(t))_{t\geq0} of bounded operator, such that $t \mapsto L(t)\theta \in C^0([0,T]; \mathcal{Z})$ is solution to

$$\begin{cases} \dot{z}(t) = Az(t) + B(t)\theta, & t \in [0, T] \\ z(0) = 0, \end{cases}$$
(5)

then, the Duhamel's formula (4) also reads $z(t) = \Phi(t)z_0 + L(t)\theta$, $t \in [0, T]$. For further use, we define \bar{L} which corresponds to L in the case of $\sigma \equiv 1$ namely the one parameter family $(\bar{L}(t))_{t\geq 0}$ of bounded operator, such that $t \mapsto \bar{L}(t)\theta \in C^0([0, T]; \mathcal{Z})$ is solution to

$$\begin{cases} \dot{z}(t) = Az(t) + \begin{pmatrix} 0\\ \theta \end{pmatrix}, \quad t \in [0, T]\\ z(0) = 0. \end{cases}$$
(6)

¹⁴² 2.2 Observation operator

We now need to describe the assumption concerning the available observations. In this work, we suppose that we have at our disposal, at every time t, measurements during the period (0,T) of the restriction of u in $\omega \subset \Omega$ an open and nonempty subset of Ω with Lipchitz boundary $\partial \omega$ and $\partial(\Omega \setminus \omega)$ also being Lipchitz. Moreover, we assume that $\omega \supseteq \omega_{gcc}$ where (ω_{gcc}, T) satisfies the following condition [51]

$$\exists T, \forall \xi \in \Omega, \exists s \in [0, T], \exists \delta > 0; \omega_{gcc} \supset B(x_{\xi}(s); \delta), \tag{GCC}$$

where $B(x_{\xi}(s); \delta)$ is a ball centered in $x_{\xi}(s)$ with a radius of δ and where a ray, starting from ξ and following the Descartes law of reflection in its most general senses [13, 51], is parametrized by $\{x_{\epsilon}(s), s \in [0, T]\}$. Note that is well known that such subdomain should contain a part of the boundary that we denote $\Gamma = \partial \Omega \cap \partial \omega$ with no cups. Otherwise, some *whispering gallery* rays may prevent the GCC to be fulfilled.

153

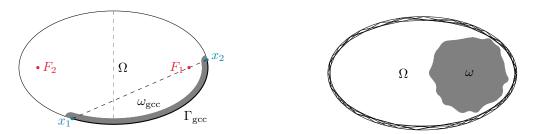


Figure 1. An interior subdomain of observation ω which does satisfy the GCC condition (left) or does not (right)

154 Let us now present the definition of the observation operator C usually introduced in control and observation theory – see for instance [36] and references therein. We 155 consider a restriction operator $I_{\omega} \in \mathcal{L}(\mathrm{H}^{1}_{0}(\Omega), \mathcal{Y})$, where \mathcal{Y} is the observation space 156 to be defined. Then formally, and in the absence of noise, at any time $t \in [0, T]$, the 157 measurement procedure could produce a measurement $y(t) = I_{\omega}u(t)$. If we assume 158 that the wave equation solution $u(t) \in H_0^1(\Omega)$ for all t, it is then natural to consider 159 that $\mathcal{Y} = \mathrm{H}^{1}_{\Gamma}(\omega)$ – the subspace of functions in $\mathrm{H}^{1}(\omega)$ null on Γ . However, in practice, 160 the measurements are always polluted by noise, hence in our deterministic context, the 161 resulting measurement should belong to a less regular space, typically $\mathcal{Y} = L^2(\omega)$. 162

If we consider regular measurements namely the observation space is $\mathcal{Y} = \mathrm{H}^{1}_{\Gamma}(\omega)$, it is rightful to introduce

$$C = \begin{pmatrix} I_{\omega} & 0 \end{pmatrix} \in \mathcal{L}(\mathcal{Z}, \mathcal{Y}), \text{ where } I_{\omega} \in \mathcal{L}(\mathrm{H}^{1}_{0}(\Omega), \mathrm{H}^{1}_{\Gamma}(\omega)).$$
(7)

As in [41], we equipped this space with the norm $\|\cdot\|_{\mathrm{H}^{1}_{\omega}} = \|E_{\omega}(\cdot)\|_{\mathrm{H}^{1}_{0}(\Omega)}$, where $E_{\omega} \in \mathcal{L}(\mathrm{H}^{1}_{\Gamma}(\omega), \mathrm{H}^{1}_{0}(\Omega))$ is defined by

$$E_{\omega}\phi := \psi \text{ the solution of } \begin{cases} \Delta\psi = 0, & \text{in } \Omega \setminus \omega, \\ \psi = \phi, & \text{in } \bar{\omega}, \\ \psi = 0, & \text{on } \partial\Omega. \end{cases}$$
(8)

- ¹⁶⁷ We recall the following lemma.
- **Lemma 2.3.** The norm $\|\cdot\|_{\mathrm{H}^{1}_{\omega}}$ is equivalent to the norm $\|\cdot\|_{\mathrm{H}^{1}(\omega)}$.

Proof. As in [41], on the one hand, we show that E_{ω} is bounded from $\mathrm{H}^{1}_{\Gamma}(\omega)$ to $\mathrm{H}^{1}_{0}(\Omega)$. Using a trace property as well as a lifting property, we obtain

$$\begin{aligned} \|\nabla(E_{\omega}\phi)\|_{\mathrm{L}^{2}(\Omega)}^{2} &\leq \|\nabla\phi\|_{\mathrm{L}^{2}(\omega)}^{2} + C^{\mathrm{st}} \|\phi_{|\partial\omega}\|_{\mathrm{H}^{1/2}(\partial\omega)}^{2} \\ &\leq C^{\mathrm{st}} \|\phi\|_{\mathrm{H}^{1}(\omega)}^{2}. \end{aligned}$$

On the other hand, we have

$$\begin{aligned} \|\phi\|_{\mathrm{H}^{1}(\omega)}^{2} &= \|\nabla(E_{\omega}\phi)\|_{\mathrm{L}^{2}(\omega)}^{2} + \|\phi\|_{\mathrm{L}^{2}(\omega)}^{2} \\ &\leq \|\nabla(E_{\omega}\phi)\|_{\mathrm{L}^{2}(\Omega)}^{2} + \|E_{\omega}\phi\|_{\mathrm{L}^{2}(\Omega)}^{2}. \end{aligned}$$

169 Using Poincaré inequality, we get

$$\|\phi\|_{\mathrm{H}^{1}(\omega)}^{2} \leq (1+C_{p})\|\nabla(E_{\omega}\phi)\|_{\mathrm{L}^{2}(\Omega)}^{2}.$$

170

For the sake of completeness and since it will be a central tool, for our results, we define the observation operator adjoint. One can show the following orthogonality property using (8)

$$\forall \phi \in \mathrm{H}^{1}_{\Gamma}(\omega), \ \forall \psi \in \mathrm{H}^{1}_{0}(\Omega) \text{ such that } \psi_{|\omega} = 0, \ (\nabla(E_{\omega}\phi), \nabla\psi)_{\mathrm{L}^{2}(\Omega)} = 0.$$

Therefore, we obtain that for all $\phi \in \mathrm{H}^{1}_{\Gamma}(\omega)$ and $\psi \in \mathrm{H}^{1}_{0}(\Omega)$,

$$(\nabla(E_{\omega}\phi),\nabla\psi)_{\mathrm{L}^{2}(\Omega)} = (\nabla(E_{\omega}\phi),\nabla(E_{\omega}\psi_{|\omega}))_{\mathrm{L}^{2}(\Omega)} - (\nabla(E_{\omega}\phi),\nabla(E_{\omega}\psi_{|\omega}-\psi))_{\mathrm{L}^{2}(\Omega)} = (\nabla(E_{\omega}\phi),\nabla(E_{\omega}\psi_{|\omega}))_{\mathrm{L}^{2}(\Omega)} = (\phi,\psi_{|\omega})_{\mathrm{H}^{1}_{*}}.$$

Then, with the Riesz representation, the adjoint of E_{ω} is I_{ω} , and we can directly infer that

$$C^* = \begin{pmatrix} E_{\omega} & 0 \end{pmatrix}^t \in \mathcal{L}(\mathrm{H}^1_{\Gamma}(\omega); \mathcal{Z}).$$

¹⁷⁶ We can see that the choice of norms for the observation operator is important, in ¹⁷⁷ particular it directly impacts the definition of its adjoint. With this choice of observation ¹⁷⁸ operator C, we restrict the space of observations to $H^1_{\Gamma}(\omega)$. In practice, for noisy ¹⁷⁹ measurements, the noise could be less regular, *i.e.* it could belong to a larger space ¹⁸⁰ such as $L^2(\omega)$. In this case, we define an observation operator where the image space ¹⁸¹ is larger than in the previous case, namely $L^2(\omega)$ instead of $H^1_{\Gamma}(\omega)$. The observation ¹⁸² operator becomes

$$C = \begin{pmatrix} I_{\omega} & 0 \end{pmatrix} : \mathcal{Z} \to L^{2}(\omega) \text{ where } I_{\omega} \in \mathcal{L}(\mathrm{H}^{1}_{0}(\Omega), \mathrm{L}^{2}(\omega)).$$
(9)

Then, the observation operator adjoint is given by $C^* = (F_{\omega} \ 0)^t$ where $F_{\omega} \in \mathcal{L}(L^2(\omega), H^1_0(\Omega))$ is defined by

$$F_{\omega}\phi := \psi \text{ the solution of } \begin{cases} -\Delta\psi = \chi_{\omega}\phi, & \text{in } \Omega, \\ \psi = 0, & \text{on } \partial\Omega, \end{cases}$$
(10)

185 where χ_{ω} is the characteristic function of ω .

¹⁸⁶ 2.3 Inverse problem setting

We can now formulate the inverse problems. We consider an actual trajectory \check{u} modeled as the first component of a mild solution \check{z} of (2) for a given and known (u_0, v_0) and $\sigma(t)$ but an unknown $\check{\theta}$. Since the initial data are known, we can assume that they are zero thanks to the linearity of the problem. Therefore, \check{u} is a particular solution of the system

$$\begin{cases} \partial_{tt}u(x,t) - \Delta u(x,t) = \sigma(t)\theta(x), & (x,t) \in \Omega \times (0,T), \\ u(x,t) = 0, & (x,t) \in \partial\Omega \times (0,T), \\ u(x,0) = 0, & x \in \Omega, \\ \partial_t u(x,0) = 0, & x \in \Omega, \end{cases}$$
(11)

¹⁹² and its first-order formulation

$$\begin{cases} \dot{z}(t) = Az(t) + B(t)\theta, & t \in [0, T], \\ z(0) = 0. \end{cases}$$
(12)

This actual trajectory is subject to measurements modeled by the bounded observation operator $C \in \mathcal{L}(\mathcal{Z}, \mathcal{Y})$. We denote the available measurements (or observations) y^{δ} which, as they are perturbed by some measurements' error – also known as the noise in our deterministic framework – are assumed to belong to $L^2(0, T; \mathcal{Y})$ equipped with the standard L^2 norm, hence $y^{\delta}(t)$ is defined for almost all $t \in (0, T)$. In our deterministic framework, the measurement *noise* intensity is represented by saying that there exists $\delta \in \mathbb{R}^+$ such that

$$\|y^{\delta} - C\check{z}\|_{\mathrm{L}^{2}(0,T;\mathcal{Y})}^{2} \leq \delta^{2}T$$

Therefore, estimating $\check{\theta} \in \mathcal{P} \subset L^2(\Omega)$ from y^{δ} consists in inverting the linear *input*output operator

$$\Psi_T: \begin{array}{c} L^2(\Omega) \to L^2((0,T); L^2(\omega)), \\ \theta \mapsto (t \mapsto CL(t)\theta). \end{array}$$
(13)

We immediately find that Ψ_T is bounded as C is bounded from \mathcal{Z} to \mathcal{Y} and L(t) is bounded from \mathcal{P} to \mathcal{Z} for all $t \in [0,T]$. Our objective is to study the condition for which Ψ_T is injective and more importantly when Ψ_T is surjective. In fact, we are going to show that Ψ_T is not surjective in $L^2(0,T; L^2(\omega))$ but define a *mildly* ill-posed problem of degree 1 or 2 (see [18]) depending on the measurements space. We will then propose a regularization strategy to invert Ψ_T using observers' methods.

²⁰⁸ 3 From observability to identifiability

In this first part, our objective is to quantify the degree of ill-posedness of Ψ_T . In particular, we prove that we have enough information in the measurements to reconstruct θ . A fundamental tool for measuring the degree of ill-posedness are observability inequalities. These inequalities are used to control parameters from observations. The most important elements are the norms that appear in this type of inequality which quantifies the ill-posedness [9, 18].

²¹⁵ 3.1 Observability conditions for the wave equation

For the homogeneous wave equation with Dirichlet conditions in the case of *field*measurements in a subdomain, we will first develop a so-called observability inequality.
In our framework, we present the proposition which allows us to express an observability condition with *field* measurements by relying on the geometric control condition.
We adopt the observability condition shown in [41], here revisited for observation domains that potentially contain a boundary part. In our case, the observation domain must contain part of the boundary.

Proposition 3.1. Let u be an arbitrary solution of (1) with $\theta = 0$. Let $T_0 > 0$. Assume that there exists $T > T_0$ such that the GCC condition is satisfied for some ω_{gcc} strict subset of ω , in the time interval (0,T). Then the following observability condition $\forall (u_0, v_0) \in \mathrm{H}^1_0(\Omega) \times \mathrm{L}^2(\Omega)$

$$\int_{0}^{T} \|u(\cdot,t)\|_{\mathrm{H}^{1}(\omega)}^{2} \,\mathrm{d}t \ge c_{obs}(T) \left(\|u_{0}\|_{\mathrm{H}^{1}_{0}(\Omega)}^{2} + \|v_{0}\|_{\mathrm{L}^{2}(\Omega)}^{2}\right) \tag{14}$$

227 holds.

Proof. The proof extended the one of [41]. In this proof, we write $a \leq b$ if there exists c > 0 independent of T and M such that $a \leq cb$. Since the geometric control condition of [13] holds, we have the classical observability inequality, with $c'_{obs}(T) > 0$

$$\int_{0}^{T} \int_{\omega} |\partial_{t} u(x,t)|^{2} \, \mathrm{d}x \, \mathrm{d}t \ge c_{\mathrm{obs}}'(T) \left(\|u_{0}\|_{\mathrm{H}_{0}^{1}(\Omega)}^{2} + \|v_{0}\|_{\mathrm{L}^{2}(\Omega)}^{2} \right).$$
(15)

We consider several subdomains satisfying $\omega_{gcc} \subseteq \omega_0 \subset \omega_\epsilon \subset \omega$. We define a cutoff function $\psi \in C^2(\overline{\Omega})$ satisfying

$$\psi(x) = \begin{cases} 0, & \text{if } x \in \Omega \setminus \omega_{\epsilon} \\ 1, & \text{if } x \in \omega_0, \end{cases}$$

and $0 \le \psi(x) \le 1$ for every $x \in \overline{\Omega}$. Particularly, we have

$$\left. \frac{\partial \psi}{\partial n} \right|_{\partial \omega \setminus \partial \Omega} = 0.$$

We introduce $\varphi(t) = t^2(T-t)^2$. Successively, integrating by parts, we obtain

$$0 = \int_0^T \int_\omega \varphi \psi(\partial_{tt} u - \Delta u) u \, dx \, dt$$

=
$$\int_0^T \int_\omega \ddot{\varphi} \psi \frac{|u|^2}{2} \, dx \, dt - \int_0^T \int_\omega \varphi \psi |\partial_t u|^2 \, dx \, dt$$

+
$$\int_0^T \int_{\partial\omega} \varphi \frac{\partial \psi}{\partial n} \frac{|u|^2}{2} \, d\sigma \, dt - \int_0^T \int_\omega \varphi \Delta \psi \frac{|u|^2}{2} \, dx \, dt$$

+
$$\int_0^T \int_\omega \varphi \psi |\nabla u|^2 \, dx \, dt$$

Using the definition of ψ , we get

$$\int_0^T \int_\omega \varphi \psi |\partial_t u|^2 \, \mathrm{d}x \, \mathrm{d}t + \int_0^T \int_\omega \varphi \Delta \psi \frac{|u|^2}{2} \, \mathrm{d}x \, \mathrm{d}t \\ = \int_0^T \int_\omega \varphi \psi |\nabla u|^2 \, \mathrm{d}x \, \mathrm{d}t + \int_0^T \int_\omega \ddot{\varphi} \psi \frac{|u|^2}{2} \, \mathrm{d}x \, \mathrm{d}t.$$

We recall that for the system (1) with $\theta = 0$, we have classically the conservation energy property namely

$$\mathcal{E}(0) := \frac{1}{2} \|u_0\|_{H_0^1(\Omega)}^2 + \frac{1}{2} \|v_0\|_{L_0^2(\Omega)}^2 = \frac{1}{2} \|u(t)\|_{H_0^1(\Omega)}^2 + \frac{1}{2} \|\partial_t u(t)\|_{L^2(\Omega)}^2 = \mathcal{E}(t).$$
(16)

234 Therefore, for a time $T > T_0$ and ε such that $T - 2\varepsilon > T_0$, we find

$$\begin{aligned} \mathcal{E}(0) &= \mathcal{E}(\varepsilon) \lesssim (c_{\rm obs}'(T))^{-1} \int_{\varepsilon}^{T-\varepsilon} \int_{\omega_0} |\partial_t u|^2 \, \mathrm{d}x \, \mathrm{d}t \\ &\lesssim (c_{\rm obs}'(T))^{-1} \int_{\varepsilon}^{T-\varepsilon} \int_{\omega_0} \left[\varphi \psi \, |\partial_t u|^2 + |\nabla u|^2\right] \, \mathrm{d}x \, \mathrm{d}t \\ &\lesssim (c_{\rm obs}'(T))^{-1} \int_0^T \int_{\omega} \left[|\nabla u|^2 + |u|^2 \right] \, \mathrm{d}x \, \mathrm{d}t, \end{aligned}$$

 $_{235}$ hence justifying (14).

Proposition 3.2. The observability constant $c_{obs}(T)$ of Proposition 3.1 satisfies, for 237 T large enough

$$c_{obs}(T) \ge cT,$$

where c > 0 is independent of T.

239 Proof. Let t > 0. Applying Proposition 3.1 for a time $T_0 > 0$, we have

$$\int_{t}^{T_{0}+t} \|u(\cdot,t)\|_{\mathrm{H}^{1}(\omega)}^{2} \,\mathrm{d}t \ge 2c_{\mathrm{obs}}(T_{0})\mathcal{E}(t).$$
(17)

240 Let $T \in [N T_0; (N+1) T_0], N \in \mathbb{N}^*$, we have

$$\int_0^T \|u(\cdot,t)\|_{\mathrm{H}^1(\omega)}^2 \,\mathrm{d}t \ge \sum_{n=0}^{N-1} \int_{nT_0}^{(n+1)T_0} \|u(\cdot,t)\|_{\mathrm{H}^1(\omega)}^2 \,\mathrm{d}t.$$

 $_{241}$ Using (17), we get

$$\int_0^T \|u(\cdot,t)\|_{\mathrm{H}^1(\omega)}^2 \,\mathrm{d}t \ge \sum_{n=0}^{N-1} 2c_{\mathrm{obs}}(T_0)\mathcal{E}(nT_0).$$

For all weak solution of (1) with $\theta = 0$, we have classically by energy conservation (16)

$$\int_0^T \|u(\cdot, t)\|_{\mathrm{H}^1(\omega)}^2 \, \mathrm{d}t \ge 2 N T_0 \, c_{\mathrm{obs}}(T_0) \mathcal{E}(0) \ge 2 \, (T - T_0) \, c_{\mathrm{obs}}(T_0) \mathcal{E}(0).$$

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This proposition provides a rough estimate of the behavior in time T of the observation constant $c_{obs}(T)$. A more precise estimate can be obtained using microlocal techniques, as in [50].

²⁴⁶ 3.2 Regular field measurements in H^1

²⁴⁷ 3.2.1 Identifiability with regular field measurements

From the observability condition adapted to initial condition reconstruction, we are now going to deduce the parameter identifiability of our source reconstruction problem. As a first step, we start by studying (12) when $\sigma \equiv 1$.

Theorem 3.3. Let $\theta \in \mathcal{P}$, when ω satisfies the GCC condition, then for T large enough, there exists a constant $c_{idf}(T) \geq cT$ where c > 0 independent of T such that

$$\int_{0}^{T} \|C\bar{L}(t)\theta\|_{\mathbf{H}_{\omega}^{1}}^{2} \, \mathrm{d}t \ge c_{idf}(T) \|\theta\|_{\mathbf{H}^{-1}(\Omega)}^{2}, \tag{18}$$

where the observation operator is the observation operator in $\mathrm{H}^{1}(\omega)$, namely $C = (I_{\omega} \ 0) \in \mathcal{L}(\mathcal{Z}, \mathrm{H}^{1}_{\Gamma}(\omega))$ and \bar{L} is given by (6).

Proof. The idea is to introduce the Riesz representation of θ seen as an element of H⁻¹ that we denote $g \in H_0^1(\Omega) \cap \mathcal{D}(-\Delta_0)$ and is given by $g = \Delta_0^{-1} \theta$. We introduce the following dynamics which in fact is equivalent to (11) with $\sigma(t) \equiv 1$.

$$\begin{cases} \partial_{tt}u_g(x,t) - \Delta u_g(x,t) = -\Delta_0 g(x), & (x,t) \in \Omega \times (0,T), \\ u_g(x,t) = 0, & (x,t) \in \partial \Omega \times (0,T), \\ u_g(x,0) = 0, & x \in \Omega, \\ \partial_t u_g(x,0) = 0, & x \in \Omega. \end{cases}$$
(19)

Note that $u_g \in \mathcal{D}(\Delta_0)$ and $\partial_t u_g \in \mathrm{H}^1(\Omega)$. We now proceed by standard energy estimate using the following energy functional

$$\mathcal{E}_g(t) := \frac{1}{2} \|\nabla u_g(t)\|_{\mathrm{L}^2(\Omega)}^2 + \frac{1}{2} \|\partial_t u_g(t)\|_{\mathrm{L}^2(\Omega)}^2$$

²⁶⁰ Integrating in space and time the first equation of (19) multiplied by $\partial_t u_g$, we get

$$\int_0^T \int_\Omega \partial_{tt} u_g \partial_t u_g \, \mathrm{d}x \, \mathrm{d}t - \int_0^T \int_\Omega \Delta u_g \partial_t u_g \, \mathrm{d}x \, \mathrm{d}t = -\int_0^T \int_\Omega \Delta_0 g \partial_t u_g \, \mathrm{d}x \, \mathrm{d}t$$

²⁶¹ Integrating by parts while using that g is time-independent, we obtain

$$\mathcal{E}_g(T) = -\int_0^T \int_\Omega \Delta_0 g \partial_t u_g \, \mathrm{d}x \, \mathrm{d}t = -\int_\Omega \Delta_0 g u_g(T) \, \mathrm{d}x = \int_\Omega \nabla g \nabla u_g(T) \, \mathrm{d}x$$

262 By Cauchy–Schwarz inequality, we get

$$\mathcal{E}_g(T) \le \|\nabla g\|_{\mathrm{L}^2(\Omega)} \|\nabla u_g(T)\|_{\mathrm{L}^2(\Omega)} \le \sqrt{2} \|\nabla g\|_{\mathrm{L}^2(\Omega)} \mathcal{E}_g^{\frac{1}{2}}(T).$$

²⁶³ This calculation finally leads to

$$\|\partial_t u_g(T)\|_{\mathcal{L}^2(\Omega)} \le 2\|\nabla g\|_{\mathcal{L}^2(\Omega)}.$$
(20)

We now introduce the system satisfied by $\tilde{u}_g = u_g - g$ in order to obtain an initial value problem for which we can use the observability condition of Proposition 3.1.

$$\begin{cases} \partial_{tt}\tilde{u}_g(x,t) - \Delta\tilde{u}_g(x,t) = 0, & (x,t) \in \Omega \times (0,T), \\ \tilde{u}_g(x,t) = 0, & (x,t) \in \partial\Omega \times (0,T), \\ \tilde{u}_g(x,0) = -g(x), & x \in \Omega, \\ \partial_t\tilde{u}_g(x,0) = 0, & x \in \Omega. \end{cases}$$

$$(21)$$

Using the observability condition of Proposition 3.1, we get

$$\int_0^T \|u_g - g\|_{\mathrm{H}^1(\omega)}^2 \,\mathrm{d}t = \int_0^T \|\tilde{u}_g\|_{\mathrm{H}^1(\omega)}^2 \,\mathrm{d}t \ge C_{\mathrm{obs}}(T) \|g\|_{\mathrm{H}^1_0(\Omega)}^2.$$
(22)

In order to get rid of the variable g in the left-hand side, we write that

$$\int_0^T \|u_g - g\|_{\mathrm{H}^1(\omega)}^2 \,\mathrm{d}t = \int_0^T \|u_g\|_{\mathrm{H}^1(\omega)}^2 \,\mathrm{d}t + \int_0^T \|g\|_{\mathrm{H}^1(\omega)}^2 \,\mathrm{d}t - 2\int_0^T (u_g, g)_{\mathrm{H}^1(\omega)} \,\mathrm{d}t. \tag{23}$$

The cross product can be controlled. Integrating the first equation of (19) in time times Δ_0^{-1} and then multiplying by g in the sense of the scalar product, we have

$$\int_0^T (u_g, g)_{\mathrm{H}^1(\omega)} \, \mathrm{d}t = T \|g\|_{\mathrm{H}^1(\omega)}^2 + (\Delta_0^{-1} \partial_t u_g(T), g)_{\mathrm{H}^1(\omega)}.$$

 $_{269}$ Injecting this result in (23), we get

$$\int_0^T \|u_g - g\|_{\mathrm{H}^1(\omega)}^2 \,\mathrm{d}t = \int_0^T \|u_g\|_{\mathrm{H}^1(\omega)}^2 \,\mathrm{d}t - T\|g\|_{\mathrm{H}^1(\omega)}^2 - 2(\Delta_0^{-1}\partial_t u_g(T), g)_{\mathrm{H}^1(\omega)}.$$

From Cauchy–Schwarz and Poincaré inequalities, we obtain

$$\int_{0}^{T} \|u_{g} - g\|_{\mathrm{H}^{1}(\omega)}^{2} \,\mathrm{d}t \leq \int_{0}^{T} \|u_{g}\|_{\mathrm{H}^{1}(\omega)}^{2} \,\mathrm{d}t - T\|g\|_{\mathrm{H}^{1}(\omega)}^{2} \\ + 2(C_{p} + 1)\|\nabla\Delta_{0}^{-1}\partial_{t}u_{g}(T)\|_{\mathrm{L}^{2}(\Omega)}\|g\|_{\mathrm{H}^{1}(\omega)}.$$
(24)

Since $\nabla \Delta_0^{-1}$ is a bounded operator from $L^2(\Omega)$ to $L^2(\Omega)$, there exists a positive scalar C^{st} such that

$$\|\nabla \Delta_0^{-1} \partial_t u_g(T)\|_{\mathrm{L}^2(\Omega)} \le C^{\mathrm{st}} \|\partial_t u_g(T)\|_{\mathrm{L}^2(\Omega)}.$$
(25)

Injecting now (25) in (24), we have

$$\begin{split} \int_0^T \|u_g - g\|_{\mathrm{H}^1(\omega)}^2 \, \mathrm{d}t &\leq \int_0^T \|u_g\|_{\mathrm{H}^1(\omega)}^2 \, \mathrm{d}t - T \|g\|_{\mathrm{H}^1(\omega)}^2 \\ &\quad + 2C^{\mathrm{st}}(C_p + 1) \|\partial_t u_g(T)\|_{\mathrm{L}^2(\Omega)} \|g\|_{\mathrm{H}^1(\omega)}. \end{split}$$

Thanks to (20) and using Young inequality, we obtain

$$\begin{split} \int_0^T \|u_g - g\|_{\mathrm{H}^1(\omega)}^2 \, \mathrm{d}t &\leq \int_0^T \|u_g\|_{\mathrm{H}^1(\omega)}^2 \, \mathrm{d}t - T\|g\|_{\mathrm{H}^1(\omega)}^2 \\ &+ 2\epsilon C^{\mathrm{st}}(C_p + 1)\|g\|_{\mathrm{H}^1_0(\Omega)}^2 + \frac{2}{\epsilon} C^{\mathrm{st}}(C_p + 1)\|g\|_{\mathrm{H}^1(\omega)}^2. \end{split}$$

From the observability condition (22), we finally obtain

$$\left(C_{\rm obs}(T) - 2\epsilon C^{\rm st}(C_p+1)\right) \|g\|_{\mathrm{H}^1_0(\Omega)}^2 \le \int_0^T \|u_g\|_{\mathrm{H}^1(\omega)}^2 \mathrm{d}t - \left(T - \frac{2}{\epsilon}C^{\rm st}(C_p+1)\right) \|g\|_{\mathrm{H}^1(\omega)}^2$$

272 Choosing $\epsilon = \frac{C_{\rm obs}(T)}{4C^{\rm st}(1+C_p)}$ gives us

$$\frac{C_{\rm obs}(T)}{2} \|g\|_{\rm H_0^1(\Omega)}^2 \le \int_0^T \|u_g\|_{\rm H^1(\omega)}^2 \,\mathrm{d}t - \left(T - \frac{8C^{\rm st^2}(1+C_p)^2}{C_{\rm obs}(T)}\right) \|g\|_{\rm H^1(\omega)}^2$$

For T large enough such that $T \ge \frac{8C^{\text{st}^2}(1+C_p)^2}{C_{\text{obs}}(T)}$, we have

$$\frac{C_{\rm obs}(T)}{2} \|g\|_{{\rm H}^1_0(\Omega)}^2 \le \int_0^T \|u_g\|_{{\rm H}^1(\omega)}^2 \,\mathrm{d}t$$

We conclude by noting that, by definition, $\|g\|_{\mathrm{H}^{1}_{0}(\Omega)} = \|\theta\|_{\mathrm{H}^{-1}(\Omega)}$ and by Lemma 2.3, $\|u_{g}\|_{\mathrm{H}^{1}(\omega)} \leq C^{\mathrm{st}} \|u_{g}\|_{\mathrm{H}^{1}_{\omega}}$.

From the identifiability property with $\sigma \equiv 1$, we deduce an identifiability property for the source reconstruction in (12) for a class of $\sigma \in H^1(0,T)$.

Theorem 3.4. Let $\theta \in \mathcal{P}$ and u be a solution of (11). Let $T_0 > 0$. Assuming $\sigma(t) \in \mathrm{H}^1(0,T)$ and $\sigma(0) \neq 0$ and (ω, T_0) satisfies the GCC condition, for $T \geq T_0$ large enough, then there exists a positive scalar $c'_{idf}(T)$ depending only on T such that

$$\int_{0}^{T} \|CL\theta\|_{\mathcal{Y}}^{2} \,\mathrm{d}t = \int_{0}^{T} \|u\|_{\mathrm{H}_{\omega}^{1}}^{2} \,\mathrm{d}t \ge c_{idf}'(T) \|\theta\|_{\mathrm{H}^{-1}(\Omega)}^{2}, \tag{26}$$

where the observation operator is the observation operator in $\mathrm{H}^{1}_{\Gamma}(\omega)$, namely $C = (I_{\omega} \ 0) \in \mathcal{L}(\mathcal{Z}, \mathrm{H}^{1}_{\Gamma}(\omega)).$

Proof. This proof is inspired by [34]. Given the scalar nature of σ , we use the first-order formulation of the system, so that observations of the solution read, for all $t \in [0, T]$

$$\begin{split} u_{|\omega}(t) &= \int_0^t C e^{(t-s)A} \sigma(s) \begin{pmatrix} 0\\ \theta \end{pmatrix} \, \mathrm{d}s = CL(t)\theta \\ &= \int_0^t \sigma'(t-s) \left(\int_0^s C e^{\tau A} \begin{pmatrix} 0\\ \theta \end{pmatrix} \, \mathrm{d}\tau \right) \, \mathrm{d}s + \sigma(0) \left(\int_0^t C e^{\tau A} \begin{pmatrix} 0\\ \theta \end{pmatrix} \, \mathrm{d}\tau \right) \\ &= \int_0^t \sigma'(t-s) C \bar{L}(s) \theta \, \mathrm{d}s + \sigma(0) C \bar{L}(t) \theta. \end{split}$$

We define a Volterra operator $S: L^2(0,T;\mathcal{Y}) \to L^2(0,T;\mathcal{Y})$ ([15, Part 1, Chapter 2]),

$$Sy(t) = y(t) + \int_0^t k_T(t-s)y(s) \,\mathrm{d}s, \text{ where } k_T(t) = \begin{cases} \frac{\sigma'(t)}{\sigma(0)} & t \in [0,T], \\ 0 & t > T. \end{cases}$$
(27)

²⁸⁴ Using this Volterra operator, we have

$$u_{|\omega}(t) = \sigma(0) S \circ C\bar{L}(t)\theta, \ t \in [0, T].$$

$$(28)$$

Since $k_T \in L^1(\mathbb{R}^+)$, using [15, Part 1, Chapter 2, Theorem 3.1], it is well-known that Sis a linear bijective continuous operator from $L^2(0,T;\mathcal{Y})$ onto $L^2(0,T;\mathcal{Y})$. In particular, there exists a unique $r_T \in L^1(0,T)$ such that

$$y(t) = Sy(t) - \int_0^t r_T(t-s)Sy(s) \,\mathrm{d}s.$$
 (29)

288 Consequently, we obtain

$$\|y\|_{L^{2}(0,T;\mathcal{Y})} \leq \left(1 + \|r_{T}\|_{L^{1}(0,T)}\right) \|Sy\|_{L^{2}(0,T;\mathcal{Y})}.$$
(30)

Combining (28) and (30) with $y = C\bar{L}\theta$, we have

$$\|C\bar{L}\theta\|_{L^{2}(0,T;\mathcal{Y})} \leq \left(1 + \|r_{T}\|_{L^{1}(0,T)}\right) \frac{1}{|\sigma(0)|} \|u_{|\omega}\|_{L^{2}(0,T;\mathcal{Y})}.$$
(31)

Combining the above equation with the observability condition of Proposition 3.3, we conclude the proof. $\hfill \Box$

We deduce from the observability inequality (26) that the observations should belong to $H^1(\omega)$, for a stable reconstruction is only in a H⁻¹ norm. As a consequence, we face an ill-posed inverse problem of degree 1 for an expected reconstruction in L² for measurements in H¹,

Remark 3.5. In some cases, it is possible to affirm that the identifiability constant of Theorem 3.4 has the same dependence in T as that of the Theorem 3.3, for instance, when σ generates a persistent excitation. Let us take an example of such function

$$\sigma(t) = \frac{1}{1+\beta}(e^{-\alpha t} + \beta), \text{ where } \beta, \alpha \text{ are two positive constants.}$$

In fact, this corresponds to a small perturbation of the case where $\sigma(t) \equiv 1$.

Following the proof of Theorem 3.4 one has to show that $||r_T||_{L^1(0,T)}$ is independent of T where r_T is defined implicitly by σ (and k_T) in (29). On one hand, we have

$$\begin{aligned} \left\| \frac{\sigma'(t)}{\sigma(0)} \right\|_{\mathrm{L}^{1}(0,T)} &= \left\| k_{T} \right\|_{\mathrm{L}^{1}(0,T)} = \left\| -\frac{\alpha}{1+\beta} e^{-\alpha t} \right\|_{\mathrm{L}^{1}(0,T)} \leq \frac{\alpha}{1+\beta} \int_{0}^{T} e^{-\alpha t} \,\mathrm{d}t, \\ &\leq \frac{1}{1+\beta} (1-e^{-\alpha T}) \leq \frac{1}{1+\beta}. \end{aligned}$$

From [15, Part 1, Chapter 2], we have an explicit expression for r_T in terms of k_T ,

$$r_T = \sum_{j=1}^{\infty} (-1)^{j-1} (k_T)^{*j},$$

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303 where $(k_T)^{*j}$ is the (j-1)-fold convolution of k_T with itself. Hence,

$$||r_T||_{\mathrm{L}^1(0,T)} \le \sum_{j=1}^{\infty} ||k_T||_{\mathrm{L}^1(0,T)}^j$$

The common ratio $||k_T||_{L^1(0,+\infty)}$ is smaller than one, hence,

$$\|r_T\|_{\mathbf{L}^1(0,T)} \le \left(1 - \|k_T\|_{L^1(0,T)}\right)^{-1} \le \frac{1}{1 - \frac{1}{1 + \beta}} = \frac{1 + \beta}{\beta},$$

305 which does not depend on time T.

As we have seen in the previous proof, the dependency of the identifiability constant with respect time is not always explicit. When we consider the simpler case $\sigma(t) \equiv 1$, we have a constant $C^{\text{st}}(T) \approx T$ as well as when σ is defined as follows:

$$\forall T > 0, \sigma = \bar{\sigma}_{|[0,T]} \text{ with } \bar{\sigma} \in \mathrm{H}^{1}(\mathbb{R}^{+}) \text{ and } \left\| \frac{\bar{\sigma}'}{\sigma(0)} \right\|_{\mathrm{L}^{1}(\mathbb{R}^{+})} \leq 1,$$

as illustrated in the previous example in Remark 3.5. One can compute that any
compactly support function will not satisfy this condition. Later, we use the following
assumption when explicitly mentioned.

Assumption 1. The identifiability constant of Theorem (3.4) satisfies

$$c'_{idf}(T) \ge cT$$
, with $c > 0$ independent of T.

313 3.2.2 From the identifiability result to a generalized Tikhonov regularization 314 strategy

We recall that we assume that a target trajectory \check{u} produced by a target source term parameter $\check{\theta}$ has resulted in a set of measurements y^{δ} where

$$\int_{0}^{T} \|y^{\delta}(x,t) - \check{u}(x,t)\|_{\mathrm{H}^{1}(\omega)}^{2} \,\mathrm{d}t \le \delta^{2}T,\tag{32}$$

for a parameter $\delta > 0$. We introduce a cost functional to minimize $\mathscr{J}_T : \mathrm{H}^1_0(\Omega) \to \mathbb{R}$ defined by

$$\mathscr{J}_{T}(\theta) = \frac{\epsilon}{2} \|\theta\|_{\mathrm{H}_{0}^{1}(\Omega)}^{2} + \frac{1}{2} \int_{0}^{T} \|y^{\delta}(t) - u_{|\theta}\|_{\mathrm{H}^{1}(\omega)}^{2} \,\mathrm{d}t$$
$$= \frac{\epsilon}{2} \|\theta\|_{\mathrm{H}_{0}^{1}(\Omega)}^{2} + \frac{1}{2} \int_{0}^{T} \|y^{\delta}(t) - Cz_{|\theta}\|_{\mathrm{H}^{1}(\omega)}^{2} \,\mathrm{d}t, \quad (33)$$

where $u_{|_{\theta}}$ is the solution of (11) associated with the parameter θ , C is defined by (7) and $\epsilon > 0$. Since the functional is quadratic and the application $\theta \mapsto u_{|_{\omega}}$ is continuous and $\theta \mapsto \|\theta\|_{\mathrm{H}^{1}_{0}(\Omega)}^{2}$ is coercive, we have then the existence of a unique minimum of \mathscr{J}_{T} . We denote the unique minimizer in $\mathrm{H}^{1}_{0}(\Omega)$ by $\bar{\theta}_{T} := \underset{\theta \in \mathrm{H}^{1}(\Omega)}{\operatorname{argmin}} \mathscr{J}_{T}$.

Theorem 3.6. We assume that there exists M > 0 such that

$$\check{\theta} \in \mathrm{H}^{1}_{0}(\Omega) \text{ and } \|\check{\theta}\|_{\mathrm{H}^{1}_{0}(\Omega)} \leq M.$$
(34)

Furthermore, we also assume that the available measurements are noise-free, i.e. we have for all $t \in [0,T]$

$$y^{\delta}(t) = \check{u}_{|\omega}(t).$$

Let $T_0 > 0$. Assuming $\sigma(t) \in H^1(0,T)$ and $\sigma(0) \neq 0$ and (ω,T_0) satisfies the GCC condition, for $T \geq T_0$ large enough and ϵ small enough, then there exists a positive scalar $C^{st}(T)$ depending only on T such that

$$\|\bar{\theta}_T - \check{\theta}\|_{L^2(\Omega)} \le C^{st}(T) \,\epsilon^{\frac{1}{4}} \, M,\tag{35}$$

327 where $\bar{\theta}_T$ minimizes the cost functional (33).

Moreover, under Assumption 1, there exists a positive scalar C^{st} independent of T such that

$$\|\bar{\theta}_T - \check{\theta}\|_{\mathcal{L}^2(\Omega)} \le C^{st} T^{-\frac{1}{4}} \epsilon^{\frac{1}{4}} M.$$

$$(36)$$

Proof. In this proof, we write $a \leq b$ if there exists c > 0 independent of T, ϵ and Msuch that $a \leq cb$.

332 Let $\tilde{u} = \bar{u} - \check{u}$ be the solution of the following dynamics

$$\begin{cases} \partial_{tt}\tilde{u}(x,t) - \Delta\tilde{u}(x,t) = \sigma(t)(\bar{\theta}_T(x) - \check{\theta}(x)), & (x,t) \in \Omega \times (0,T), \\ \tilde{u}(x,t) = 0, & (x,t) \in \partial\Omega \times (0,T), \\ \tilde{u}(x,0) = 0, & x \in \Omega, \\ \partial_t\tilde{u}(x,0) = 0, & x \in \Omega. \end{cases}$$
(37)

Identifying L^2 to its dual, we have

$$\|\bar{\theta}_T - \check{\theta}\|_{\mathrm{L}^2(\Omega)}^2 \le \langle \bar{\theta}_T - \check{\theta}, \bar{\theta}_T - \check{\theta} \rangle_{\mathrm{H}^1_0(\Omega), \mathrm{H}^{-1}(\Omega)} \le \|\bar{\theta}_T - \check{\theta}\|_{\mathrm{H}^1_0(\Omega)} \|\bar{\theta}_T - \check{\theta}\|_{\mathrm{H}^{-1}(\Omega)}.$$

³³⁴ Using the identifiability property (26) of Theorem 3.4, we get

$$\|\bar{\theta}_T - \check{\theta}\|_{\mathrm{L}^2(\Omega)}^2 \le (c'_{\mathrm{idf}}(T))^{-\frac{1}{2}} \|\bar{\theta}_T - \check{\theta}\|_{\mathrm{H}^1_0(\Omega)} \left(\int_0^T \|\bar{u} - \check{u}\|_{\mathrm{H}^1_\omega}^2 \,\mathrm{d}t\right)^{\frac{1}{2}}.$$
 (38)

335 By triangle inequality and since we consider free-noise observations, we have

$$\|\bar{\theta}_T - \check{\theta}\|_{\mathrm{L}^2(\Omega)}^2 \le (c'_{\mathrm{idf}}(T))^{-\frac{1}{2}} \left(\|\bar{\theta}_T\|_{\mathrm{H}^1_0(\Omega)} + \|\check{\theta}\|_{\mathrm{H}^1_0(\Omega)} \right) \left(\int_0^T \|\bar{u} - y^{\delta}\|_{\mathrm{H}^1_\omega}^2 \,\mathrm{d}t \right)^{\frac{1}{2}}.$$

Using the definition of the cost-functional (33) taking into account the norm equivalence of Lemma 2.3, and the prior on the parameter (34), we get

$$\|\bar{\theta}_T - \check{\theta}\|_{\mathrm{L}^2(\Omega)}^2 \lesssim (c'_{\mathrm{idf}}(T))^{-\frac{1}{2}} \left(\left(2\epsilon^{-1} \mathcal{J}_T(\bar{\theta}_T) \right)^{\frac{1}{2}} + M \right) \left(2\mathcal{J}_T(\bar{\theta}_T) \right)^{\frac{1}{2}}$$

Remembering that $y^{\delta}(t) = \check{u}_{|\omega}(t)$ and using again (34), we can bound the cost-functional as follows

$$\mathscr{J}_T(\bar{\theta}_T) \le \mathscr{J}_T(\check{\theta}) \le \frac{\epsilon}{2} M^2.$$

340 Eventually, we obtain

$$\|\bar{\theta}_T - \check{\theta}\|_{\mathrm{L}^2(\Omega)}^2 \lesssim (c'_{\mathrm{idf}}(T))^{-\frac{1}{2}} \left((M^2)^{\frac{1}{2}} + M \right) (\epsilon M^2)^{\frac{1}{2}},$$

which corresponds to the desired result (35).

In addition, under Assumption 1, we have that there exists c > independent of T such that

$$c'_{\text{idf}}(T) \ge cT,$$

which leads to (36).

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Theorem 3.7. We assume that there exists M > 0 such that (34) holds. Let $T_0 > 0$. Assuming $\sigma(t) \in H^1(0,T)$ and $\sigma(0) \neq 0$ and (ω,T_0) satisfies the GCC condition, for $T \geq T_0$ large enough and ϵ small enough, then there exists a positive scalar $C^{st}(T)$ depending only on T such that

$$\|\bar{\theta}_T - \check{\theta}\|_{\mathcal{L}^2(\Omega)} \le C^{st}(T) \,\epsilon^{-\frac{1}{4}} \,\delta,\tag{39}$$

where $\bar{\theta}_T$ minimizes the cost functional (33).

Moreover, under Assumption 1, there exists a positive scalar C^{st} independent of T such that

$$\|\bar{\theta}_T - \check{\theta}\|_{\mathcal{L}^2(\Omega)} \le C^{st} T^{\frac{1}{4}} \epsilon^{-\frac{1}{4}} \delta.$$

$$\tag{40}$$

Proof. In this proof, we write $a \leq b$ if there exists c > 0 independent of T, ϵ, δ and M such that $a \leq cb$.

Reasoning exactly as in the proof of Theorem 3.6, we can obtain (38). Inserting $y^{\delta} - y^{\delta} = 0$ in (38) and by triangular inequality we find

$$\begin{split} \|\bar{\theta}_{T} - \check{\theta}\|_{\mathrm{L}^{2}(\Omega)}^{2} &\lesssim (c_{\mathrm{idf}}^{\prime}(T))^{-\frac{1}{2}} \left(\|\check{\theta}\|_{\mathrm{H}^{1}_{0}(\Omega)} + \|\bar{\theta}_{T}\|_{\mathrm{H}^{1}_{0}(\Omega)} \right) \\ &\times \left[\left(\int_{0}^{T} \|u_{|\bar{\theta}_{T}} - y^{\delta}\|_{\mathrm{H}^{1}_{\omega}}^{2} \,\mathrm{d}t \right)^{\frac{1}{2}} + \left(\int_{0}^{T} \|\check{u} - y^{\delta}\|_{\mathrm{H}^{1}_{\omega}}^{2} \,\mathrm{d}t \right)^{\frac{1}{2}} \right]. \end{split}$$

Using Lemma 2.3, the estimate on the noise amplitude (32) and the prior on the parameter (34), we find

$$\|\bar{\theta}_T - \check{\theta}\|_{\mathrm{L}^2(\Omega)}^2 \lesssim (c'_{\mathrm{idf}}(T))^{-\frac{1}{2}} \left(M + \sqrt{2\epsilon^{-1} \mathscr{J}_T(\bar{\theta}_T)} \right) \left(\sqrt{2 \mathscr{J}_T(\bar{\theta}_T)} + \delta \sqrt{T} \right).$$

Using again (32) and (34), we can bound the cost functional as follows

$$\mathscr{J}_T(\bar{\theta}_T) \le \mathscr{J}_T(\check{\theta}) \le \frac{\epsilon}{2}M^2 + \frac{1}{2}\delta^2 T.$$

355 Hence, we finally get

$$\|\bar{\theta}_T - \check{\theta}\|_{\mathrm{L}^2(\Omega)}^2 \lesssim (c'_{\mathrm{idf}}(T))^{-\frac{1}{2}} \left(M + \sqrt{M^2 + \epsilon^{-1}\delta^2 T} \right) \left(\sqrt{\epsilon M^2 + \delta^2 T} + \delta \sqrt{T} \right), \quad (41)$$

 $_{356}$ which implies (39).

Furthermore, under Assumption 1, we have that there exists c > independent of Tsuch that

$$c'_{\text{idf}}(T) \ge cT,$$

 $_{359}$ which proves (40).

Finally, choosing $\epsilon = \delta^2 M^{-2}$ in (41), we obtain the following corollary

Corollary 3.8. Under the hypothesis of Theorem 3.7, for $T \ge T_0$ large enough, there exists a positive scalar $C^{st}(T)$ depending only on T such that

$$\|\bar{\theta}_T - \check{\theta}\|_{L^2(\Omega)} \le C^{st}(T) \sqrt{M} \sqrt{\delta}, \tag{42}$$

- where $\bar{\theta}_T$ minimizes the cost functional (33) with $\epsilon = \delta^2 M^{-2}$.
- Moreover, under Assumption 1, there exists a positive scalar C^{st} independent of T such that

$$\|\bar{\theta}_T - \check{\theta}\|_{L^2(\Omega)} \le C^{st} T^{\frac{1}{4}} \sqrt{M} \sqrt{\delta}.$$
(43)

366 3.3 Field measurements with L^2 -noise

³⁶⁷ In the rest of this work, we assume – in order to model a more realistic case in practice –

that the measurements are less regular. Namely, we suppose that the target trajectory

 \check{u} generated by the target source term parameter $\check{\theta}$ has given rise to a collection of

370 measurements y^{δ} such that

$$\int_{0}^{T} \|y^{\delta}(x,t) - \check{u}(x,t)\|_{\mathrm{L}^{2}(\omega)}^{2} \,\mathrm{d}t \le \delta^{2}T,\tag{44}$$

371 for a parameter $\delta > 0$.

We present a method to deal with this additional difficulty of regularity loss. This approach consists in modifying the definition of the observation operator with second observation operator C introduced in (9).

- 375 3.3.1 Identifiability from observability
- The proof of the desired identifiability result for the identification estimation problem is based on the classical following initial condition observability inequality – for a homogeneous wave equation in weakened norms.

Proposition 3.9. Let u be an arbitrary solution of (1) with $\theta = 0$. Assume that ω satisfies the GCC condition, then there exists $cT \leq c_{obs}(T)$ with c > 0 independent of T such that

$$\int_{0}^{T} \|u(\cdot,t)\|_{\mathrm{L}^{2}(\omega)}^{2} \,\mathrm{d}t \ge c_{obs}(T) \left(\|u_{0}\|_{\mathrm{L}^{2}(\Omega)}^{2} + \|v_{0}\|_{\mathrm{H}^{-1}(\Omega)}^{2}\right). \tag{45}$$

Theorem 3.10. Let $\theta \in \mathcal{P}$. Let u be a solution of (11). Assuming $\sigma(t) \in \mathrm{H}^1(0,T)$ and $\sigma(0) \neq 0$, and that ω satisfies the GCC condition, for T large enough, there exists $c'_{idf}(T) > 0$, such that

$$\|CL\theta\|_{L^{2}(0,T;L^{2}(\omega))}^{2} = \|u\|_{L^{2}(0,T;L^{2}(\omega))}^{2} \ge c_{idf}'(T)\|\Delta_{0}^{-1}\theta\|_{L^{2}(\Omega)}^{2}.$$

³⁸⁵ *Proof.* One can show, as for Theorem 3.4, that we have the inequality

$$\|u_{|\omega}\|_{L^{2}(0,T;L^{2}(\omega))} \ge C^{\mathrm{st}}(T) \|C\bar{L}\theta\|_{L^{2}(0,T;L^{2}(\omega))}.$$
(46)

As in the proof of Theorem 3.3, we introduce $g \in H_0^1(\Omega) \cap \mathcal{D}(-\Delta_0)$ the Riesz representation of θ (seen as an element of H^{-1}) given by $g = -\Delta_0^{-1}\theta$, and u_g solution of the dynamics (19). The difference $\tilde{u}_g = u_g - g$ satisfies the system (21) in the same way. We apply the observability condition of Proposition 3.9 on this system, and we obtain

$$\|\tilde{u}_g\|_{\mathrm{L}^2(0,T;\mathrm{L}^2(\omega))}^2 = \|u_g - g\|_{\mathrm{L}^2(0,T;\mathrm{L}^2(\omega))}^2 \ge c_{\mathrm{obs}}(T) \|g\|_{\mathrm{L}^2(\Omega)}^2.$$
(47)

Now we want to discard the variable g in the left-hand of (47). That is why, we develop the following scalar product

$$\int_0^T \|u_g - g\|_{\mathrm{L}^2(\omega)}^2 \,\mathrm{d}t = \int_0^T \|u_g\|_{\mathrm{L}^2(\omega)}^2 \,\mathrm{d}t + T\|g\|_{\mathrm{L}^2(\omega)}^2 - 2\int_0^T (u_g, g)_{\mathrm{L}^2(\omega)} \,\mathrm{d}t.$$
(48)

Applying Δ_0^{-1} to the first equation of (19), then integrate in time and finally the resulting equation given by the scalar product with g is

$$\int_0^T (u_g, g)_{\mathrm{L}^2(\omega)} \, \mathrm{d}t = (\Delta_0^{-1} \partial_t u_g(T), g)_{\mathrm{L}^2(\omega)} + T \|g\|_{\mathrm{L}^2(\omega)}^2.$$

Injecting this result in (48) and using Cauchy-Schwarz and Poincaré inequalities, we obtain

$$\int_{0}^{T} \|u_{g} - g\|_{L^{2}(\omega)}^{2} dt \leq \int_{0}^{T} \|u_{g}\|_{L^{2}(\omega)}^{2} dt - T\|g\|_{L^{2}(\omega)}^{2} + 2C_{p}\|\nabla\Delta_{0}^{-1}\partial_{t}u_{g}(T)\|_{L^{2}(\Omega)}\|g\|_{L^{2}(\omega)}.$$
 (49)

We expect to bound $\|\nabla \Delta_0^{-1} \partial_t u_g(T)\|_{L^2(\Omega)}$ by $\|g\|_{L^2(\Omega)}$. Therefore, we apply Δ_0^{-1} to the first equation of (19). Integrating in time, the scalar product with $-\partial_t u_g$ leads to

$$-\int_{\Omega} (\Delta_0^{-1} \partial_{tt} u_g) \partial_t u_g \, \mathrm{d}x + \int_{\Omega} u_g \partial_t u_g \, \mathrm{d}x = \int_{\Omega} g \partial_t u_g \, \mathrm{d}x.$$

Inserting $\Delta_0 \Delta_0^{-1} = \text{Id}$ in the first term we get

$$-\int_{\Omega} (\Delta_0^{-1} \partial_{tt} u_g) \Delta_0 \Delta_0^{-1} \partial_t u_g \, \mathrm{d}x + \frac{1}{2} \frac{\mathrm{d}}{\mathrm{d}t} \int_{\Omega} u_g^2 \, \mathrm{d}x = \int_{\Omega} g \partial_t u_g \, \mathrm{d}x.$$

397 Integrating by parts, we obtain

$$\frac{1}{2}\frac{\mathrm{d}}{\mathrm{d}t}\int_{\Omega} (\nabla\Delta_0^{-1}\partial_t u_g)^2 \,\mathrm{d}x + \frac{1}{2}\frac{\mathrm{d}}{\mathrm{d}t}\int_{\Omega} u_g^2 \,\mathrm{d}x = \int_{\Omega} g\partial_t u_g \,\mathrm{d}x.$$

Integrating in time and using Cauchy-Schwarz and Young inequalities, we get

$$\begin{aligned} \frac{1}{2} \|\nabla \Delta_0^{-1} \partial_t u_g(T)\|_{L^2(\Omega)}^2 + \frac{1}{2} \|u_g(T)\|_{L^2(\Omega)}^2 &= \int_\Omega g u_g(T) \, dx \\ &\leq \frac{1}{2} \|g\|_{L^2(\Omega)}^2 + \frac{1}{2} \|u_g(T)\|_{L^2(\Omega)}^2, \end{aligned}$$

398 thus,

$$\|\nabla \Delta_0^{-1} \partial_t u_g(T)\|_{L^2(\Omega)}^2 \le \|g\|_{L^2(\Omega)}^2$$

399 Injecting this result in (49), we have

$$\int_0^T \|u_g - g\|_{\mathrm{L}^2(\omega)}^2 \,\mathrm{d}t \le \int_0^T \|u_g\|_{\mathrm{L}^2(\omega)}^2 \,\mathrm{d}t - T\|g\|_{\mathrm{L}^2(\omega)}^2 + 2C_p\|g\|_{\mathrm{L}^2(\Omega)}\|g\|_{\mathrm{L}^2(\omega)}.$$

400 Using Young inequality, we finally have

$$\int_{0}^{T} \|u_{g} - g\|_{\mathrm{L}^{2}(\omega)}^{2} \,\mathrm{d}t \leq \int_{0}^{T} \|u_{g}\|_{\mathrm{L}^{2}(\omega)}^{2} \,\mathrm{d}t - T\|g\|_{\mathrm{L}^{2}(\omega)}^{2} + \epsilon C_{p}\|g\|_{\mathrm{L}^{2}(\Omega)}^{2} + \frac{C_{p}}{\epsilon}\|g\|_{\mathrm{L}^{2}(\omega)}^{2}.$$
 (50)

401 Gathering (47) and (50), we get

$$\int_{0}^{T} \|u_{g}\|_{L^{2}(\omega)}^{2} dt \geq \left(T - \frac{C_{p}}{\epsilon}\right) \|g\|_{L^{2}(\omega)}^{2} + (c_{\text{obs}}(T) - \epsilon C_{p}) \|g\|_{L^{2}(\Omega)}^{2}.$$

402 We choose $\epsilon = \frac{c_{\text{obs}}}{2C_p}$, then we have

$$\int_0^T \|u_g\|_{\mathrm{L}^2(\omega)}^2 \,\mathrm{d}t \ge \left(T - \frac{2C_p^2}{c_{\mathrm{obs}}(T)}\right) \|g\|_{\mathrm{L}^2(\omega)}^2 + \frac{c_{\mathrm{obs}}(T)}{2} \|g\|_{\mathrm{L}^2(\Omega)}^2$$

403 For $T > \frac{2C_p^2}{c_{\rm obs}(T)}$, and setting $-\Delta_0 g = \theta$, we obtain

$$\|C\bar{L}\theta\|_{L^{2}((0,T);L^{2}(\omega))} = \int_{0}^{T} \|u_{g}\|_{L^{2}(\omega)}^{2} dt \ge \frac{c_{\text{obs}}(T)}{2} \|\Delta_{0}^{-1}\theta\|_{L^{2}(\Omega)}^{2}.$$
 (51)

404 Gathering (46) and (51), we conclude the proof.

As in the previous section, we assume that the dependence of the observability constant is in O(T).

407 Assumption 2. The identifiability constant of Theorem (3.10) satisfies

 $c'_{idf}(T) \ge cT$, with c > 0 independent of T.

408 3.3.2 Updated convergence rate for the Tikhonov regularization

We infer from Theorem 3.10 that we now face an ill-posed problem of degree 2. This will modify the convergence rate of our Tikhonov regularization without further assumption

⁴¹⁰ modify the convergence rate of our Tikhonov regularization without further assumption ⁴¹¹ on the prior regularity. Let us first recall the following classical interpolation lemma.

Lemma 3.11. Let (μ_i, e_i) be the orthonormal basis for the L^2 scalar product such that the Laplacian eigenvalues $\mu_i > 0$ satisfies $\mu_{i+1} > \mu_i$, and $\mu_i \to +\infty$. Let $\theta \in H_0^1(\Omega)$, we have the interpolation property

$$\|\theta\|_{L^{2}(\Omega)}^{2} \leq \left(\|\Delta_{0}^{-1}\theta\|_{L^{2}(\Omega)}^{2}\right)^{\frac{1}{3}} \left(\|\theta\|_{H_{0}^{1}(\Omega)}^{2}\right)^{\frac{2}{3}},$$

415 where $\|\Delta_0^{-1}\theta\|_{L^2(\Omega)} = \sqrt{\sum_i |(e_i,\theta)|^2 \mu_i^{-2}}.$

416 We recall here the proof for the sake of completeness.

417

⁴¹⁸ *Proof.* We express the L^2 -norm of θ in the orthonormal basis.

$$\|\theta\|_{L^2(\Omega)}^2 = \sum_{i=0}^{+\infty} |(e_i,\theta)|^2 = \sum_{i=1}^{+\infty} |(e_i,\theta)|^2 \mu_i^{2q\frac{s}{s-q}} \mu_i^{2s\frac{-q}{s-q}}.$$

419 Using Hölder inequality, we have

$$\|\theta\|_{\mathrm{L}^{2}(\Omega)}^{2} \leq \left(\sum_{i} |(e_{i},\theta)|^{2} \mu_{i}^{2q}\right)^{\frac{s}{s-q}} \left(\sum_{i} |(e_{i},\theta)|^{2} \mu_{i}^{2s}\right)^{\frac{-q}{s-q}}$$

420 Choosing $s = \frac{1}{2}$ and q = -1, we obtain

$$\|\theta\|_{\mathbf{L}^{2}(\Omega)}^{2} \leq \left(\sum_{i} |(e_{i},\theta)|^{2} \mu_{i}^{-2}\right)^{\frac{1}{3}} \left(\sum_{i} |(e_{i},\theta)|^{2} \mu_{i}\right)^{\frac{2}{3}} = \left(\|\Delta_{0}^{-1}\theta\|_{\mathbf{L}^{2}(\Omega)}^{2}\right)^{\frac{1}{3}} \left(\|\theta\|_{\mathbf{H}_{0}^{1}(\Omega)}^{2}\right)^{\frac{2}{3}}.$$

421

Keeping the same regularization in our criterion, we introduce the cost functional to minimize

$$\mathcal{J}_{T}(\theta) = \frac{\epsilon}{2} \|\theta\|_{\mathrm{H}^{1}_{0}(\Omega)}^{2} + \frac{1}{2} \|y^{\delta} - Cz_{|\theta}\|_{L^{2}(0,T;\mathrm{L}^{2}(\omega))}^{2},$$
(52)

where C is defined as operator in L², namely $C = (I_{\omega} \ 0) \in \mathcal{L}(\mathcal{Z} \to L^2(\omega))$. We can then estimate the reconstruction error between $\bar{\theta}_T$ the functional minimizer and the target parameter $\check{\theta}$.

427 Theorem 3.12. We assume that there exists M > 0 such that (34) holds.

Furthermore, we also assume that the available measurements are noise-free, i.e. we have for all $t \in [0, T]$

$$y^{\delta}(t) = \check{u}_{|\omega}(t)$$

19 of 40

430 Let $T_0 > 0$. Assuming $\sigma(t) \in H^1(0,T)$ and $\sigma(0) \neq 0$ and (ω, T_0) satisfies the GCC 431 condition, for $T \geq T_0$ large enough and ϵ small enough, then there exists a positive 432 scalar $C^{st}(T)$ depending only on T such that

$$\|\bar{\theta}_T - \check{\theta}\|_{\mathrm{L}^2(\Omega)} \le C^{st}(T) \,\epsilon^{\frac{1}{6}} \, M,\tag{53}$$

- 433 where $\bar{\theta}_T$ minimizes the cost functional (52).
- ⁴³⁴ Moreover, under Assumption 2, there exists a positive scalar C^{st} independent of T such ⁴³⁵ that

$$\|\bar{\theta}_T - \check{\theta}\|_{\mathcal{L}^2(\Omega)} \le C^{st} T^{-\frac{1}{6}} \epsilon^{\frac{1}{6}} M.$$

$$\tag{54}$$

- ⁴³⁶ Proof. In this proof, we write $a \leq b$ if there exists c > 0 independent of T, ϵ and M⁴³⁷ such that $a \leq cb$.
- ⁴³⁸ We use the interpolation inequality of Lemma 3.11 to obtain

$$\|\bar{\theta}_T - \check{\theta}\|_{L^2(\Omega)}^2 \le \left(\|\bar{\theta}_T - \check{\theta}\|_{H^1_0(\Omega)}^2\right)^{\frac{2}{3}} \left(\|\Delta_0^{-1}(\bar{\theta}_T - \check{\theta})\|_{L^2(\Omega)}^2\right)^{\frac{1}{3}}.$$

439 Using the identifiability result given by Theorem 3.10, we have

$$\|\bar{\theta}_T - \check{\theta}\|_{L^2(\Omega)}^2 \le (c'_{\text{idf}}(T))^{-\frac{1}{3}} \left(\|\bar{\theta}_T\|_{H^1_0(\Omega)}^2 + \|\check{\theta}\|_{H^1_0(\Omega)}^2 \right)^{\frac{2}{3}} \left(\int_0^T \|u_{|\bar{\theta}_T} - \check{u}\|_{L^2(\omega)}^2 \, \mathrm{d}t \right)^{\frac{1}{3}}.$$
(55)

Since we consider noise-free observations, we have $y^{\delta} = \check{u}_{|\omega}$. Therefore, we obtain

$$\|\bar{\theta}_T - \check{\theta}\|_{\mathrm{L}^2(\Omega)}^2 \lesssim (c'_{\mathrm{idf}}(T))^{-\frac{1}{3}} \left(\|\bar{\theta}_T\|_{\mathrm{H}^1_0(\Omega)}^2 + \|\check{\theta}\|_{\mathrm{H}^1_0(\Omega)}^2 \right)^{\frac{2}{3}} \left(\int_0^T \|u_{|\bar{\theta}_T} - y^{\delta}\|_{\mathrm{L}^2(\omega)}^2 \,\mathrm{d}t \right)^{\frac{1}{3}}$$

⁴⁴¹ Using the definition of the cost-functional (52) and our prior on the parameter (34),⁴⁴² we obtain

$$\|\bar{\theta}_T - \check{\theta}\|_{\mathrm{L}^2(\Omega)}^2 \lesssim (c'_{\mathrm{idf}}(T))^{-\frac{1}{3}} \left(2\epsilon^{-1}\mathcal{J}_T(\bar{\theta}_T) + M^2\right)^{\frac{2}{3}} \left(2\mathcal{J}_T(\bar{\theta}_T)\right)^{\frac{1}{3}}$$

443 We can bound the cost-functional as follows

$$\mathscr{J}_T(\bar{\theta}_T) \le \mathscr{J}_T(\check{\theta}) \le \frac{\epsilon}{2} M^2.$$

444 Finally, we have

$$\|\bar{\theta}_T - \check{\theta}\|_{\mathrm{L}^2(\Omega)}^2 \lesssim (c'_{\mathrm{idf}}(T))^{-\frac{1}{3}} (M^2 + M^2)^{\frac{2}{3}} (\epsilon M^2)^{\frac{1}{3}},$$

which leads to (53). Furthermore, under Assumption 2, we have that there exists c >independent of T such that

$$c'_{\text{idf}}(T) \ge cT,$$

- 447 which proves (54).
- **Theorem 3.13.** We assume that there exists M > 0 such that (34) holds.

Let $T_0 > 0$. Assuming $\sigma(t) \in H^1(0,T)$ and $\sigma(0) \neq 0$ and (ω, T_0) satisfies the GCC condition, for $T \geq T_0$ large enough and ϵ small enough, then there exists a positive scalar $C^{st}(T)$ depending only on T such that

$$\|\bar{\theta}_T - \check{\theta}\|_{\mathrm{L}^2(\Omega)} \le C^{st}(T) \,\epsilon^{-\frac{1}{3}}\delta,\tag{56}$$

452 where $\bar{\theta}_T$ minimizes the cost functional (52).

⁴⁵³ Moreover, under Assumption 2, there exists a positive scalar C^{st} independent of T such ⁴⁵⁴ that

$$\|\bar{\theta}_T - \check{\theta}\|_{\mathcal{L}^2(\Omega)} \le C^{st} T^{\frac{1}{3}} \epsilon^{-\frac{1}{3}} \delta.$$

$$\tag{57}$$

- ⁴⁵⁷ On one hand, using the estimate on the amplitude noise (44) and the prior (34), we
- 458 intend to bound the cost functional.

$$\mathcal{J}_T(\check{\theta}) \le \frac{\epsilon}{2} M^2 + \frac{1}{2} \delta^2 T.$$
(58)

On the other hand, as in the proof of Theorem 3.12, one can show that (55) holds. Inserting $y^{\delta} - y^{\delta} = 0$ into (55) to reintroduce the cost functional and using the prior on the parameter (34), we get

$$\begin{aligned} \|\bar{\theta}_{T} - \check{\theta}\|_{L^{2}(\Omega)}^{2} \lesssim (c_{\mathrm{idf}}^{\prime}(T))^{-\frac{1}{3}} \left(2\epsilon^{-1}\mathcal{J}_{T}(\check{\theta}) + M^{2}\right)^{\frac{2}{3}} \\ \times \left(\int_{0}^{T} \|u_{|\bar{\theta}_{T}} - y^{\delta}\|_{L^{2}(\omega)}^{2} \,\mathrm{d}t + \int_{0}^{T} \|\check{u} - y^{\delta}\|_{L^{2}(\omega)}^{2} \,\mathrm{d}t\right)^{\frac{1}{3}}. \end{aligned}$$

Again, using again (32) and the cost functional bound (58), we have

$$\|\bar{\theta}_T - \check{\theta}\|_{\mathrm{L}^2(\Omega)}^2 \lesssim (c'_{\mathrm{idf}}(T))^{-\frac{1}{3}} \left(2M^2 + \epsilon^{-1}\delta^2 T\right)^{\frac{2}{3}} \left(2\delta^2 T + \epsilon M^2\right)^{\frac{1}{3}},\tag{59}$$

460 which corresponds to the desired result (56).

Furthermore, under Assumption 2, we have that there exists c > independent of T_{462} such that

$$c'_{\text{idf}}(T) \ge cT,$$

463 which leads to (57).

Eventually, choosing $\epsilon = \delta^2 M^{-2}$ in (59), we obtain the following corollary.

Corollary 3.14. Under the hypothesis of Theorem 3.13, for $T \ge T_0$ large enough, then there exists a positive scalar $C^{st}(T)$ depending only on T such that

$$\|\bar{\theta}_T - \check{\theta}\|_{L^2(\Omega)} \le C^{st}(T) M^{\frac{2}{3}} \delta^{\frac{1}{3}},\tag{60}$$

467 where $\bar{\theta}_T$ minimizes the cost functional (52) with $\epsilon = \delta^2 M^{-2}$.

⁴⁶⁸ Moreover, under Assumption 2, there exists a positive scalar C^{st} independent of T such ⁴⁶⁹ that

$$\|\bar{\theta}_T - \check{\theta}\|_{L^2(\Omega)} \le C^{st} T^{\frac{1}{3}} M^{\frac{4}{3}} \delta^{\frac{1}{3}}.$$
(61)

470 4 Observer design

Our goal in this section is to propose a sequential method for computing $\bar{\theta}_T$ and the associated wave equation trajectory \bar{z}_T . By sequential, we mean that we want to update the estimator when measurements are available, without recomputing the estimator each time a new measurement is available. This leads to the definition of an observer in the sense of [20], which is formally defined in a general framework of an observed system of a dynamics of the form

$$\begin{cases} \dot{\boldsymbol{z}} = F(\boldsymbol{z}, t) \\ \boldsymbol{z}(0) = \boldsymbol{z}_0 \end{cases}$$
(62)

where we observed a target trajectory \check{z} with unknown initial state \check{z}_0 by measuring $(y^{\delta}(t))_{t\geq 0}$. We here propose to slightly simplify this definition.

Definition 4.1. An asymptotic observer of (62) is a function $\mathbb{R}_+ \ni t \mapsto \hat{z}(t; (y^{\delta}(s))_{s \leq t})$ such that

481 1. Robustness to Noise: for all $0 \le \tau$, and for all $\varepsilon > 0$, there exist $\delta, \alpha > 0$, 482 such that

$$\frac{\|\hat{\boldsymbol{z}}(0) - \check{\boldsymbol{z}}(0)\|_{\boldsymbol{\mathcal{Z}}} \le \alpha}{\|\check{\boldsymbol{y}} - \boldsymbol{y}^{\delta}\|_{L^{2}((0,\tau);\boldsymbol{\mathcal{Y}})} \le \delta} \right\} \Rightarrow \|\check{\boldsymbol{z}}(t) - \hat{\boldsymbol{z}}(t)\|_{\boldsymbol{\mathcal{Z}}} \le \varepsilon, \quad t \ge \tau$$

483 484

2. Convergence: for all $0 \leq \tau$, and for all $\varepsilon, \alpha, \delta > 0$, there exists $T_0 > 0$ such that

$$\begin{aligned} \|\hat{\boldsymbol{z}}(0) - \check{\boldsymbol{z}}(0)\|_{\boldsymbol{\mathcal{Z}}} &\leq \alpha \\ \|\check{\boldsymbol{y}} - \boldsymbol{y}^{\delta}\|_{\mathrm{L}^{2}((0,\tau);\boldsymbol{\mathcal{Y}})} &\leq \delta \\ \mathrm{supp}(\check{\boldsymbol{y}} - \boldsymbol{y}^{\delta}) \subset (0,\tau) \end{aligned} \right\} \Rightarrow \|\check{\boldsymbol{z}}(t) - \hat{\boldsymbol{z}}(t)\|_{\boldsymbol{\mathcal{Z}}} \leq \varepsilon, \quad t \geq T_{0} + \tau. \end{aligned}$$

485

Note that the two above properties correspond respectively to an asymptotic convergence property in the absence of noise and property of robustness to noise. As the observer can be computed at time t from measurements available at time $s \leq t$, we say that it is a causal function of the measurements. Note also that, as a consequence of (i), an asymptotic observer satisfies for all $s \geq 0$

$$\hat{\boldsymbol{z}}(s) = \check{\boldsymbol{z}}(s), \quad (\boldsymbol{y}^{\delta}(t) = \check{\boldsymbol{y}}(t), t \ge s) \Rightarrow (\hat{\boldsymbol{z}}(t) = \check{\boldsymbol{z}}(t), t \ge s).$$
 (63)

In our case, we try to define such an observer for the joint dynamics of $z = (z, \theta)$, where the solution z of the wave equation is extended by the zero dynamics of the parameter θ . Our method is based on a dynamic programming strategy that relate functional minimization to Kalman observer, which we will apply to the dynamics of the state and parameters. But let us first recall some results about Kalman filtering [5, adapted here to our infinite-dimensional system.

497 4.1 Reduced-order Kalman filters for quasi-evolution operators

To present our strategy, we first need to review some results in the analysis of evolution problems with non-homogeneous generators.

- 500 4.1.1 Quasi-evolution equation
- 501 Let us consider the following abstract dynamics defined in the state-space $\boldsymbol{\mathcal{Z}}$

$$\begin{cases} \dot{\boldsymbol{z}}(t) = \boldsymbol{A}(t)\boldsymbol{z}(t) + \boldsymbol{\beta}(t), & \text{in } (0,T), \\ \boldsymbol{z}(0) = \boldsymbol{z}_0. \end{cases}$$
(64)

from given $\boldsymbol{z}_0 \in \boldsymbol{\mathcal{Z}}, \boldsymbol{\beta} \in \mathrm{L}^2((0,T); \boldsymbol{\mathcal{Z}})$. Here, the operator

$$\boldsymbol{A}(t) = \boldsymbol{A} + \boldsymbol{G}(t)$$

where $(\boldsymbol{A}, \mathcal{D}(\boldsymbol{A}))$ is maximal dissipative and for all t > 0, $\boldsymbol{G}(t)$ is a strongly continuous bounded operator, namely for all $\boldsymbol{z} \in \boldsymbol{Z}$, $\boldsymbol{G}(\cdot)\boldsymbol{z} \in C^0([0,T],\boldsymbol{Z})$. Note that by uniform boundedness theorem, there exists $c_{st} > 0$ such that $\|\boldsymbol{G}(t)\|_{\mathcal{L}(\boldsymbol{Z})} \leq c_{st}$. Classical results for evolution equation give the existence of a solution of (64) knowing \boldsymbol{z}_0 and $\boldsymbol{\beta}(t)$ and denoted $\boldsymbol{z}_{|\boldsymbol{z}_0,\boldsymbol{\beta}}$. Indeed, from [29, II-1 Proposition 3.4], if $\boldsymbol{z}_0 \in \boldsymbol{Z}$ and $\boldsymbol{\beta} \in L^2((0,T);\boldsymbol{Z})$ then $\boldsymbol{z}_{|\boldsymbol{z}_0,\boldsymbol{\beta}} \in C^0([0,T];\boldsymbol{Z})$ is a mild solution

$$\boldsymbol{z}(t) = \boldsymbol{\Phi}(t,0)\boldsymbol{z}_0 + \int_0^t \boldsymbol{\Phi}(t,s)\boldsymbol{\beta}(s) \,\mathrm{d}s, \qquad (65)$$

where $\boldsymbol{\Phi}: \Delta_T \to \mathcal{L}(\boldsymbol{Z})$ with $\Delta_T = \{(t,s): s \in [0,T], t \in [s,T]\}$ is a quasi-evolution operator in the sense of [8, Definition 1.2]. 511 Remark 4.1. Note that it is possible to give an implicit definition of Φ . Indeed, 512 consider the system,

$$\begin{cases} \dot{\boldsymbol{z}}(t) = \boldsymbol{A}\boldsymbol{z}(t) + \boldsymbol{G}(t)\boldsymbol{z}(t), & \text{ in } (s,T), \\ \boldsymbol{z}(s) = \boldsymbol{z}_s \in \boldsymbol{\mathcal{Z}}. \end{cases}$$

513 By definition of the quasi-evolution operator $\boldsymbol{\Phi}$, we have, using Duhamel's formula

$$\boldsymbol{z}(t) = \boldsymbol{\Phi}(t,s)\boldsymbol{z}_s = e^{(t-s)\boldsymbol{A}}\boldsymbol{z}_s + \int_s^t e^{(t-\tau)\boldsymbol{A}}\boldsymbol{G}(\tau)\boldsymbol{z}(\tau) \,\mathrm{d}\tau, \tag{66}$$

⁵¹⁴ which, using (66), finally becomes

$$\boldsymbol{\Phi}(t,s)\boldsymbol{z}_s = e^{(t-s)\boldsymbol{A}}\boldsymbol{z}_s + \int_s^t e^{(t-\tau)\boldsymbol{A}}\boldsymbol{G}(\tau)\boldsymbol{\Phi}(\tau,s)\boldsymbol{z}_s \,\mathrm{d}\tau.$$

The mild solution (65) is also a weak solution in the sense that $\boldsymbol{z}_{|\boldsymbol{z}_0,\boldsymbol{\beta}} \in L^2((0,T);\boldsymbol{Z}) \cap$ H¹((0,T); $\mathcal{D}(\boldsymbol{A}^*)$), namely for all $q \in \mathcal{D}(\boldsymbol{A}^*)$, $(q, \boldsymbol{z}(\cdot))_{\boldsymbol{Z}}$ belongs to H¹(0,T) and for almost all $t \in (0,T)$,

$$\forall \boldsymbol{q} \in \mathcal{D}(\boldsymbol{A}^*), \quad \frac{\mathrm{d}}{\mathrm{dt}}(\boldsymbol{q}, \boldsymbol{z}(t))_{\boldsymbol{\mathcal{Z}}} + (\boldsymbol{A}(t)^* \boldsymbol{q}, \boldsymbol{z}(t))_{\boldsymbol{\mathcal{Z}}} = (\boldsymbol{q}, \boldsymbol{\beta}(t))_{\boldsymbol{\mathcal{Z}}}.$$
(67)

⁵¹⁸ Moreover, from [29, II-1 Propostion 3.5], if we further assume that for all $\boldsymbol{z} \in \mathcal{D}(\boldsymbol{A}), \, \boldsymbol{G}(\cdot)\boldsymbol{z} \in \mathrm{C}^{0}([0,T], \mathcal{D}(\boldsymbol{A})), \, \boldsymbol{z}_{0} \in \mathcal{D}(\boldsymbol{A}) \text{ and } \boldsymbol{\beta} \in \mathrm{L}^{2}((0,T); \mathcal{D}(\boldsymbol{A})) \text{ then } \boldsymbol{z}_{|\boldsymbol{z}_{0},\boldsymbol{\beta}} \in \mathrm{C}^{0}([0,T]; \mathcal{D}(\boldsymbol{A})) \cap \mathrm{H}^{1}((0,T); \boldsymbol{\mathcal{Z}}).$ Finally, when for all $\boldsymbol{z} \in \mathcal{D}(\boldsymbol{A}), \, \boldsymbol{G}(\cdot)\boldsymbol{z} \in \mathrm{C}^{1}([0,T]; \boldsymbol{\mathcal{Z}}),$ ⁵²¹ $\boldsymbol{z}_{0} \in \mathcal{D}(\boldsymbol{A}) \text{ and } \boldsymbol{\beta} \in \mathrm{H}^{1}((0,T); \mathcal{D}(\boldsymbol{A})) \text{ then } \boldsymbol{z}_{|\boldsymbol{z}_{0},\boldsymbol{\beta}} \in \mathrm{C}^{0}([0,T], \mathcal{D}(\boldsymbol{A})) \cap \mathrm{C}^{1}([0,T]; \boldsymbol{\mathcal{Z}}).$

522 4.1.2 Least-square estimation

We introduce a space $\mathcal{V} \subset \mathcal{Z}$ with continuous injection and a symmetric and coercive bilinear form $m: \mathcal{V} \times \mathcal{V} \to \mathbb{R}$. We want to minimize a criterion $\mathscr{J}_T: \mathcal{V} \to \mathbb{R}$ defined by

$$\mathscr{J}_{T}(\zeta) = \frac{1}{2}m(\zeta,\zeta) + \frac{1}{2}\int_{0}^{T} \|y^{\delta}(t) - \boldsymbol{C}(t)\boldsymbol{z}_{|\zeta}(t)\|_{\mathcal{Y}}^{2} \,\mathrm{d}t,\tag{68}$$

where $\boldsymbol{z}_{|\zeta}$ is the solution of (64) for an initial condition $\boldsymbol{z}_0 = \hat{\boldsymbol{z}}_0 + \zeta$, where $\hat{\boldsymbol{z}}_0$ is a known prior. In this section, the observation operator is a time-dependent operator $\boldsymbol{C} \in C^0([0,T], \mathcal{L}(\boldsymbol{Z}, \mathcal{Y}))$. Note that depending on the choice of $\boldsymbol{\mathcal{V}}$, the bilinear form mcan be seen as Tikhonov regularization or a generalized Tikhonov regularization [18].

The functional \mathscr{J}_T is continuous, differentiable, and strictly convex, so it admits one and only one minimizer denoted $\overline{\zeta}_{|T}$ where we underline the dependency with respect to T. Computing the derivative of \mathscr{J}_T , we have for all $\xi \in \mathcal{V}$,

$$\langle \mathcal{D} \mathscr{J}_T(\zeta), \xi \rangle_{\boldsymbol{\mathcal{V}}', \boldsymbol{\mathcal{V}}} = m(\zeta, \xi) + \int_0^T (\xi, \boldsymbol{\varPhi}(t, 0)^* \boldsymbol{C}(t)^* (\boldsymbol{C}(t) \boldsymbol{z}_{|\zeta} - y^{\delta}))_{\mathcal{Y}} \, \mathrm{d}t.$$

Therefore, introducing the adjoint dynamics for any trajectory $\boldsymbol{z}_{|\zeta} \in L^2((0,T);\boldsymbol{Z})$ and $y^{\delta} \in L^2((0,T);\boldsymbol{Y}),$

$$\begin{cases} \dot{\boldsymbol{q}}(t) + \boldsymbol{A}(t)^* \boldsymbol{q}(t) = -\boldsymbol{C}(t)^* (y^{\delta}(t) - \boldsymbol{C}(t) \boldsymbol{z}_{|\zeta}(t)), & t \in (0, T) \\ \boldsymbol{q}(T) = 0, \end{cases}$$
(69)

whose mild solution in $C^0([0,T], \mathbf{Z})$ is given by the Duhamel formulae – adapted to this backward formulation –

$$\boldsymbol{q}(t) = \int_{t}^{T} \boldsymbol{\Phi}(t,s)^{*} \boldsymbol{C}(t)^{*} (y^{\delta}(s) - \boldsymbol{C}(t)\boldsymbol{z}_{|\zeta}(s)) \, \mathrm{d}s, \quad t \in [0,T),$$
(70)

533 we find

$$\forall \xi \in \boldsymbol{\mathcal{V}}, \quad \langle \mathbf{D} \mathscr{J}_T(\zeta), \xi \rangle_{\boldsymbol{\mathcal{V}}', \boldsymbol{\mathcal{V}}} = m(\zeta, \xi) - (\boldsymbol{q}(0), \xi)_{\boldsymbol{\mathcal{Z}}}.$$
(71)

If we now introduce the operator $\Pi_0 \in \mathcal{L}(\mathcal{Z}, \mathcal{V})$ such that for all $q \in \mathcal{Z}$, $\Pi_0 q$ is the solution of $m(\Pi_0 q, \xi) = (q, \xi)_{\mathcal{Z}}$ for all $\xi \in \mathcal{V}$, we deduce from (71) the optimality condition

 $\bar{\zeta}_{|T} = \boldsymbol{\Pi}_0 \bar{\boldsymbol{q}}_T(0),$

where \bar{q}_T is associated to the optimal trajectory $\bar{z}_T = z_{|\bar{\zeta}|T}$ through the two-ends optimality system – as called in [14, 5]

$$\begin{cases} \dot{\bar{\boldsymbol{z}}}_{T}(t) = \boldsymbol{A}(t)\bar{\boldsymbol{z}}_{T} + \boldsymbol{\beta}(t), & t \in (0,T), \\ \dot{\bar{\boldsymbol{q}}}_{T}(t) + \boldsymbol{A}(t)^{*}\bar{\boldsymbol{q}}_{T} = -\boldsymbol{C}(t)^{*}(y^{\delta}(t) - \boldsymbol{C}(t)\bar{\boldsymbol{z}}_{T}(t)), & t \in (0,T), \\ \bar{\boldsymbol{z}}_{T}(0) = \hat{\boldsymbol{z}}_{0} + \boldsymbol{\Pi}_{0}\bar{\boldsymbol{q}}_{T}(0), & \\ \bar{\boldsymbol{q}}_{T}(T) = 0. \end{cases}$$
(72)

⁵³⁹ Note finally that if $\mathcal{V} \subset \mathcal{D}(\mathbf{A})$, then $\bar{\mathbf{z}}$ and $\bar{\mathbf{q}}$ belong to $C^1([0,T], \mathbf{Z}) \cap C^0([0,T], \mathcal{D}(\mathbf{A}))$.

540 4.2 Kalman estimator definition

541 The objective of this section is to present briefly the Kalman estimator strategy to

⁵⁴² solve sequentially the two-ends problem (72). First, we introduce the Riccati dynamics

$$\begin{cases} \dot{\boldsymbol{\Pi}}(t) = \boldsymbol{A}(t)\boldsymbol{\Pi}(t) + \boldsymbol{\Pi}(t)\boldsymbol{A}(t)^* - \boldsymbol{\Pi}(t)\boldsymbol{C}(t)^*\boldsymbol{C}(t)\boldsymbol{\Pi}(t), & \text{ in } (0,T), \\ \boldsymbol{\Pi}(0) = \boldsymbol{\Pi}_0, \end{cases}$$
(73)

where $C \in C^0([0,T], \mathcal{L}(\mathcal{Z}, \mathcal{Y}))$. The initial condition $\Pi_0 \in S^+(\mathcal{Z})$ the space of bounded symmetric operators. The following theorem gives the existence of solution of (73) in $S^{+}(\mathcal{Z})$. The resulting operator $\Pi \in S^+(\mathcal{Z})$ is called Riccati operator or *covariance operator* for its interpretation in the context of stochastic filtering, see for instance [5, 7].

⁵⁴⁸ Theorem 4.2 (Existence of Riccati solution).

The Riccati dynamics (73) admits one and only one mild solution $\boldsymbol{\Pi} \in C^0([0,T], \mathcal{S}^+(\boldsymbol{Z}))$ given by

$$\boldsymbol{\Pi}(t) = \boldsymbol{\Phi}(t,0)\boldsymbol{\Pi}_{0}\boldsymbol{\Phi}(t,0)^{*} - \int_{0}^{t} \boldsymbol{\Phi}(t,s)\boldsymbol{\Pi}(s)\boldsymbol{C}(t)^{*}\boldsymbol{C}(t)\boldsymbol{\Pi}(s)\boldsymbol{\Phi}(t,s)^{*} \,\mathrm{d}s.$$
(74)

Moreover, for all $(z_1, z_2) \in \mathcal{D}(\mathbf{A}^*)$, $(\mathbf{\Pi}(\cdot)z_1, z_2)_{\mathbf{Z}}$ is differentiable and satisfies the weak formulation

$$\frac{\mathrm{d}}{\mathrm{dt}} (\boldsymbol{\Pi}(t)\boldsymbol{z}_1, \boldsymbol{z}_2)_{\boldsymbol{\mathcal{Z}}} = (\boldsymbol{\Pi}(t)\boldsymbol{z}_1, \boldsymbol{A}(t)^*\boldsymbol{z}_2)_{\boldsymbol{\mathcal{Z}}} + (\boldsymbol{\Pi}(t)\boldsymbol{A}(t)^*\boldsymbol{z}_1, \boldsymbol{z}_2)_{\boldsymbol{\mathcal{Z}}} - (\boldsymbol{C}(t)\boldsymbol{\Pi}(t)\boldsymbol{z}_1, \boldsymbol{C}(t)\boldsymbol{\Pi}(t)\boldsymbol{z}_2)_{\boldsymbol{\mathcal{Y}}} = 0, \quad t \in [0, T].$$
(75)

⁵⁵¹ Proof. We refer to [8, Corrolary 2.2], for a proof of existence of a mild solution of ⁵⁵² a Riccati equation where A(t) = A + G(t) with $(A, \mathcal{D}(A))$ unbounded and G(t) a ⁵⁵³ bounded perturbation. The existence of a weak solution is a direct consequence, see ⁵⁵⁴ for instance [29, IV-1 Proposition 2.1].

From the Riccati operator definition, we can introduce the so-called *Kalman observer* - also called *Kalman estimator* in the context of stochastic filtering. Theorem 4.3. There exists one and only one mild solution in $C^0([0,T]; \mathcal{Z})$ of the dynamics

$$\begin{cases} \dot{\hat{\boldsymbol{z}}} = \boldsymbol{A}(t)\hat{\boldsymbol{z}} + \boldsymbol{\Pi}(t)\boldsymbol{C}(t)^*(\boldsymbol{y}^{\delta}(t) - \boldsymbol{C}(t)\hat{\boldsymbol{z}}(t)), & t \in (0,T), \\ \hat{\boldsymbol{z}}(0) = \hat{\boldsymbol{z}}_0, \end{cases}$$
(76)

559 in the sense that

$$\hat{\boldsymbol{z}}(t) = \boldsymbol{\Phi}(t,0)\hat{\boldsymbol{z}}_0 + \int_0^t \boldsymbol{\Phi}(t,s)\boldsymbol{\Pi}(s)\boldsymbol{C}(s)^*[\boldsymbol{y}^{\delta}(s) - \boldsymbol{C}(s)\hat{\boldsymbol{z}}(s)] \,\mathrm{d}s.$$
(77)

Moreover, this solution is the unique weak solution in the sense that, (1.) $z \in L^2((0,T); \mathbf{Z})$, (2.) for all $q \in \mathcal{D}(\mathbf{A}^*)$, $(q, \hat{z}(\cdot))_{\mathbf{Z}}$ belongs to $H^1(0,T)$ and (3.) for almost all $t \in (0,T)$,

$$\forall q \in \mathcal{D}(\boldsymbol{A}^*), \quad \frac{\mathrm{d}}{\mathrm{dt}}(\hat{\boldsymbol{z}}(t), q)_{\boldsymbol{\mathcal{Z}}} + (\hat{\boldsymbol{z}}, \boldsymbol{A}(t)^* q)_{\boldsymbol{\mathcal{Z}}} = (y^{\delta}(t) - \boldsymbol{C}(t)\hat{\boldsymbol{z}}, \boldsymbol{C}(t)\boldsymbol{\varPi}(t)q)_{\boldsymbol{\mathcal{Y}}}.$$
 (78)

From Proof. We know that $\mathbf{\Pi} \in C^0([0,T]; \mathcal{S}^*_+(\mathbf{Z}))$. Therefore, for all $t \in [0,T]$, $\mathbf{A}(t) - \mathbf{\Pi}(t)\mathbf{C}(t)^*\mathbf{C}(t)$ is a uniform bounded perturbation of \mathbf{A} , the existence enters exactly the framework of quasi-evolution operator.

566 4.2.1 The optimal filtering property

567 We can now move to the understanding on how the Riccati operator and the associated

568 Kalman observer decouples the two ends problems (72). The following proof extend

the one proposed in [57] to quasi evolution operator.

Theorem 4.4. The Kalman estimator defined by the dynamics (76) satisfies the fundamental identity

$$\forall t \in [0, T], \quad \bar{\boldsymbol{z}}_T(t) = \hat{\boldsymbol{z}}(t) + \boldsymbol{\varPi}(t)\bar{\boldsymbol{q}}_T(t).$$
(79)

Proof. Let us introduce $\eta = \hat{z} - \bar{z}_T + \Pi \bar{q}_T$ and $v \in \mathcal{D}(A^*)$, and compute

$$\frac{\mathrm{d}}{\mathrm{d}t}(\eta(t),v)_{\boldsymbol{\mathcal{Z}}} = \left(\hat{\boldsymbol{z}}(t),\boldsymbol{A}(t)^{*}v\right)_{\boldsymbol{\mathcal{Z}}} + \left(y^{\delta}(t) - \boldsymbol{C}(t)\hat{\boldsymbol{z}}(t),\boldsymbol{C}(t)\boldsymbol{\varPi}(t)v\right)_{\boldsymbol{\mathcal{Y}}} - \left(\bar{\boldsymbol{z}}_{T}(t),\boldsymbol{A}(t)^{*}v\right)_{\boldsymbol{\mathcal{Z}}} + \frac{\mathrm{d}}{\mathrm{d}t}(\boldsymbol{\varPi}(t)\bar{\boldsymbol{q}}_{T}(t),v)_{\boldsymbol{\mathcal{Z}}}, \quad (80)$$

where we need to properly define the last-term time derivative. In fact, as \bar{q} can be defined from a Duhamel formula of the form (70), and $\boldsymbol{\Pi}$ is mild solution of the Riccati equation (73), hence is absolutely continuous, we have from [8, Lemme 2.5] that

$$\frac{\mathrm{d}}{\mathrm{d}t}(\boldsymbol{\Pi}(t)\bar{\boldsymbol{q}}_{T}(t),v)_{\boldsymbol{\mathcal{Z}}} = \frac{\mathrm{d}}{\mathrm{d}t} \left(\boldsymbol{\Pi}(t)\bar{\boldsymbol{q}}_{T}(s),v\right)_{\boldsymbol{\mathcal{Z}}}\Big|_{s=t} + \frac{\mathrm{d}}{\mathrm{d}t} \left(\bar{\boldsymbol{q}}_{T}(t),\boldsymbol{\Pi}(s)v\right)_{\boldsymbol{\mathcal{Z}}}\Big|_{s=t}.$$

Moreover, as $\boldsymbol{\Pi}$ is a weak solution of (73), we have from (75),

$$\frac{\mathrm{d}}{\mathrm{d}t} (\boldsymbol{\Pi}(t)\bar{\boldsymbol{q}}_{T}(s), v)\Big|_{s=t} = (\boldsymbol{\Pi}(t)\bar{\boldsymbol{q}}_{T}(t), \boldsymbol{A}(t)^{*}v)_{\boldsymbol{\mathcal{Z}}} + (\bar{\boldsymbol{q}}_{T}(t), \boldsymbol{A}(t)^{*}\boldsymbol{\Pi}(t)v)_{\boldsymbol{\mathcal{Z}}} - (\boldsymbol{C}(t)\boldsymbol{\Pi}(t)\bar{\boldsymbol{q}}_{T}(t), \boldsymbol{C}(t)\boldsymbol{\Pi}(t)v)_{\boldsymbol{\mathcal{Y}}}, \quad (81)$$

and, as $\bar{\boldsymbol{q}}_T$ is a weak solution of (72),

$$\frac{\mathrm{d}}{\mathrm{d}t} (\bar{\boldsymbol{q}}_T(t), \boldsymbol{\Pi}(s)v) \Big|_{s=t} = -(\bar{\boldsymbol{q}}_T(t), \boldsymbol{A}(t)^* \boldsymbol{\Pi}(t)v)_{\boldsymbol{\mathcal{Z}}} - (y^{\delta}(t) - \boldsymbol{C}(t)\bar{\boldsymbol{z}}_T(t), \boldsymbol{C}(t)\boldsymbol{\Pi}(t)v)_{\boldsymbol{\mathcal{Y}}}.$$
 (82)

572 Gathering (80), (81) and (82), we finally obtain

$$\begin{cases} \frac{\mathrm{d}}{\mathrm{d}t}(\eta(t), v)_{\boldsymbol{\mathcal{Z}}} = (\eta(t), \boldsymbol{A}(t)^* - \boldsymbol{C}(t)^* \boldsymbol{C}(t) \boldsymbol{\varPi}(t) v)_{\boldsymbol{\mathcal{Z}}}, & t \in [0, T] \\ \eta(0) = 0 \end{cases}$$

whence, by Theorem 4.3, $\eta = 0$ in [0, T], which concludes the proof.

In particular, we see that by taking t = T that for all t > 0, $\hat{z}(t) = \bar{z}_t(t)$. This why the Kalman observer is also called optimal filtering estimator as it consists of sequentially computing the trajectory associated with least-square estimator at final time.

578 4.2.2 The reduced order case

As we saw in the previous section, the inverse problem can be solved using the Kalman 579 filter, but this requires us to calculate a Riccati operator in the augmented space 580 $\boldsymbol{\mathcal{Z}}$. In order to reduce the computational cost, we now show that we can instead 581 compute another Riccati operator on a reduced space \mathcal{Z}_r that takes into account a 582 priori informations on the inverse problem: in our case, described specifically later, 583 only the source terms is not known. The reduced space \mathcal{Z}_r is assumed to be a closed 584 subspace of $\boldsymbol{\mathcal{Z}}$ (equipped with the same scalar product) hence we can assume given a 585 projection operator $P_r \in \mathcal{L}(\mathbf{Z}, \mathbf{Z}_r)$ in this space. Moreover, we also assumed given a 586 compact subspace \mathcal{V}_r of \mathcal{Z}_r that is used to define a regularized version of the functionnal 587 to minimize. The minimization problem reads 588

$$\min_{\zeta_r \in \boldsymbol{\mathcal{V}}_r} \Big\{ \mathscr{J}_T(\zeta_r) = \frac{1}{2} m_r(\zeta_r, \zeta_r) + \frac{1}{2} \int_0^T \|y^{\delta} - \boldsymbol{C}\boldsymbol{z}_{|P_r^*\zeta_r}\|_{\boldsymbol{\mathcal{Y}}}^2 \,\mathrm{d}t \Big\},\$$

where $\mathbf{z}_{|P_r^*\zeta_r}$ is solution of (64) for an initial condition $\mathbf{z}_0 = \hat{\mathbf{z}}_0 + P_r^*\zeta_r$ and m_r is a continuous and coercive bilinear form on \mathcal{V}_r . The associated two-ends problem reads

$$\begin{cases} \dot{\boldsymbol{z}}_{T} = \boldsymbol{A}(t)\boldsymbol{\bar{z}}_{T} + \boldsymbol{\beta}, & t \in (0,T), \\ \dot{\boldsymbol{q}}_{T} + \boldsymbol{A}(t)^{*}\boldsymbol{\bar{q}}_{T} = -\boldsymbol{C}^{*}(y^{\delta}(t) - \boldsymbol{C}\boldsymbol{\bar{z}}_{T}), & t \in (0,T), \\ \boldsymbol{\bar{z}}_{T}(0) = \boldsymbol{\hat{z}}_{0} + P_{r}^{*}\Lambda_{0}P_{r}\boldsymbol{\bar{q}}_{T}(0), \\ \boldsymbol{\bar{q}}_{T}(T) = 0, \end{cases}$$
(83)

⁵⁹¹ with $\Lambda_0 \in \mathcal{L}(\mathcal{V}_r)$ defined such that

$$\forall (\zeta_r, \xi_r) \in \boldsymbol{\mathcal{V}}_r \times \boldsymbol{\mathcal{V}}_r, \quad m_r(\Lambda_0 \zeta_r, \xi_r) = (\zeta_r, \xi_r)_{\boldsymbol{\mathcal{Z}}}.$$

⁵⁹² The resulting Kalman observer can still be defined from the Riccati dynamics

$$\begin{cases} \dot{\boldsymbol{\Pi}}(t) = \boldsymbol{A}(t)\boldsymbol{\Pi}(t) + \boldsymbol{\Pi}(t)\boldsymbol{A}(t)^* - \boldsymbol{\Pi}(t)\boldsymbol{C}^*\boldsymbol{C}\boldsymbol{\Pi}(t), \quad t > 0, \\ \boldsymbol{\Pi}(0) = P_r^*\Lambda_0 P_r, \end{cases}$$
(84)

so that we continue to have $\bar{\boldsymbol{z}}_T(t) = \hat{\boldsymbol{z}}(t) + \boldsymbol{\Pi}(t)\bar{\boldsymbol{q}}_T(t)$ for all $t \in [0, T]$. But moreover, the Riccati solution of (84) can be deduced from a reduced covariance operator Λ , solution of a Riccati dynamics in the reduced space $\boldsymbol{\mathcal{Z}}_r$, namely

$$\begin{cases} \dot{\Lambda}(t) = -\Lambda(t)P_r\Phi(t,0)^* \boldsymbol{C}^* \boldsymbol{C}\Phi(t,0)P_r^*\Lambda(t), & t > 0, \\ \Lambda(0) = \Lambda_0. \end{cases}$$
(85)

This dynamics satisfies the condition of existence since for all $t \in [0, T]$, $C\Phi(t, 0)P_r^*$ can be considered as a time-dependent bounded observation operator. Moreover, we have the following decomposition.

- **Theorem 4.5.** Let $\Pi(t)$ the mild solution of the Riccati dynamic (84) and $\Lambda(t)$ the mild solution of the Riccati dynamic (85), then we have
 - mild solution of the Riccall aynamic (85), then we have

$$\boldsymbol{\Pi}(t) = \boldsymbol{\Phi}(t,0) P_r^* \Lambda(t) P_r \boldsymbol{\Phi}(t,0)^*.$$
(86)

⁶⁰¹ *Proof.* The mild solution of the Riccati dynamics (85) reads

$$\Lambda(t) = \Lambda(0) - \int_0^t \Lambda(s) P_r \boldsymbol{\Phi}(s,0)^* \boldsymbol{C}^* \boldsymbol{C} \boldsymbol{\Phi}(s,0) P_r^* \Lambda(s) \, \mathrm{d}s.$$

Using this mild solution, we compute the following quantity $\boldsymbol{\Phi}(t,0)P_r^*\Lambda(t)P_r\boldsymbol{\Phi}(t,0)^*$.

$$\boldsymbol{\Phi}(t,0)P_r^*\Lambda(t)P_r\boldsymbol{\Phi}(t,0)^* = \boldsymbol{\Phi}(t,0)P_r^*\Lambda(0)P_r\boldsymbol{\Phi}(t,0)^* -\int_0^t \boldsymbol{\Phi}(t,0)P_r^*\Lambda(s)P_r\boldsymbol{\Phi}(s,0)^*\boldsymbol{C}^*\boldsymbol{C}\boldsymbol{\Phi}(s,0)P_r^*\Lambda(s)P_r\boldsymbol{\Phi}(t,0)^* \,\mathrm{d}s.$$

Note that $\boldsymbol{\Phi}(t,0) = \boldsymbol{\Phi}(t,s)\boldsymbol{\Phi}(s,0)$, therefore we obtain

$$\boldsymbol{\Phi}(t,0)P_r^*\Lambda(t)P_r\boldsymbol{\Phi}(t,0)^* = \boldsymbol{\Phi}(t,0)P_r^*\Lambda(0)P_r\boldsymbol{\Phi}(t,0)^* -\int_0^t \boldsymbol{\Phi}(t,s)\boldsymbol{\Phi}(s,0)P_r^*\Lambda(s)P_r\boldsymbol{\Phi}(s,0)^*\boldsymbol{C}^*\boldsymbol{C}\boldsymbol{\Phi}(s,0)P_r^*\Lambda(s)P_r\boldsymbol{\Phi}(s,0)^*\boldsymbol{\Phi}(t,s)^* \,\mathrm{d}s.$$

602 Setting $\tilde{\boldsymbol{H}}(t) = \boldsymbol{\Phi}(t,0)P_r^* \Lambda(t)P_r \boldsymbol{\Phi}(t,0)^*$, we get

$$\tilde{\boldsymbol{\Pi}}(t) = \tilde{\boldsymbol{\Pi}}(0) - \int_0^t \boldsymbol{\Phi}(t,s) \tilde{\boldsymbol{\Pi}}(s) \boldsymbol{C}^* \boldsymbol{C} \tilde{\boldsymbol{\Pi}}(s) \boldsymbol{\Phi}(t,s)^* \, \mathrm{d}s.$$

Hence, $\tilde{\boldsymbol{\mu}}$ is a mild solution of (84) and by uniqueness of the solution of the Riccati dynamics, we have

$$\tilde{\boldsymbol{\Pi}}(t) = \boldsymbol{\Pi}(t) = \boldsymbol{\Phi}(t,0)P_r^* \Lambda(t)P_r \boldsymbol{\Phi}(t,0)^*,$$

⁶⁰⁵ which concludes the proof.

Therefore, a computational strategy can be based only on computations restricted to the reduced-space \mathbf{Z}_r .

4.2.3 Application to source estimation

We will now define an observer, which is equivalent to the calculation of the optimum $\bar{\theta}_T$ with increasing time. To that purpose, we rely on the Kalman estimation theory presented above and more precisely the Reduced-Order Kalman filtering approach of Section 4.2.2, where the uncertainty space is, here, the parameter space, a subspace of the augmented space $\mathbf{Z} = \mathbf{Z} \times \mathbf{P}$. In particular, we set $\mathbf{Z}_r = \{0\} \times L^2(\Omega)$ and $\mathbf{V}_r = \{0\} \times H_0^1(\Omega)$. The dynamics in the augmented space reads

$$\underbrace{\frac{\mathrm{d}}{\mathrm{d}t} \begin{pmatrix} z\\\theta \end{pmatrix}}_{\dot{z}} = \underbrace{\begin{pmatrix} A & B(t)\\0 & 0 \end{pmatrix} \begin{pmatrix} z\\\theta \end{pmatrix}}_{A(t)z} \quad \text{with} \quad \underbrace{\begin{pmatrix} z(0)\\\theta(0) \end{pmatrix}}_{z(0)} = \underbrace{\begin{pmatrix} 0\\\theta \end{pmatrix}}_{\zeta}. \tag{87}$$

The least square minimization (33) rewritten with respect to the augmented variable reads

$$\mathscr{J}_T(\theta) = \frac{\epsilon}{2} \|\theta\|_{\mathrm{H}^1_0(\Omega)}^2 + \frac{1}{2} \int_0^T \|y^{\delta}(t) - C\boldsymbol{z}_{|\theta}\|_{\mathcal{Y}}^2 \,\mathrm{d}t.$$
(88)

$$P_{\theta} = \begin{pmatrix} 0 & 0 \\ 0 & \mathrm{Id}_{\mathcal{P}} \end{pmatrix}$$

from the augmented space $\boldsymbol{\mathcal{Z}}$ to the reduced space $\boldsymbol{\mathcal{Z}}_r$, we have, using Theorem 4.5,

$$\begin{cases} \boldsymbol{\Pi}(t) = \boldsymbol{\Phi}(t,0)P_{\theta}^{*}\Lambda(t)P_{\theta}\boldsymbol{\Phi}(t,0)^{*}, \\ \boldsymbol{\Pi}(0) = P_{\theta}^{*}\Lambda(0)P_{\theta}, \end{cases}$$
(89)

where $\Lambda(t)$ is given by (85) with $P_r = P_{\theta}$. Then, we can specify the evolution operator associated with the following system

$$\boldsymbol{A}(t) = \begin{pmatrix} A & B(t) \\ 0 & 0 \end{pmatrix}, \qquad \begin{cases} \dot{\boldsymbol{z}}(t) = \boldsymbol{A}(t)\boldsymbol{z}(t), \\ \boldsymbol{z}(0) = \boldsymbol{z}_0, \end{cases} \Rightarrow \quad \begin{cases} \dot{\boldsymbol{z}}(t) = A\boldsymbol{z}(t) + B(t)\boldsymbol{\theta}, \\ \dot{\boldsymbol{\theta}}(t) = 0, \\ (\boldsymbol{z}, \boldsymbol{\theta})^t(0) = (\boldsymbol{z}_0, \boldsymbol{\theta}_0)^t, \end{cases}$$

⁶²¹ Therefore, for initial data of the form $\boldsymbol{z}_0 = P^*_{\theta}(0,\theta_0)^t = (0,\theta_0)^t$ the solution reads

$$\begin{pmatrix} z(t)\\ \theta(t) \end{pmatrix} = \boldsymbol{\varPhi}(t,0) \begin{pmatrix} 0\\ \theta_0 \end{pmatrix},$$

and, using the definition of L given by (5), $z(t) = L(t)\theta_0$ is the mild solution of

$$\begin{cases} \dot{z} = Az + B\theta_0, \quad t > 0\\ z(0) = 0, \end{cases}$$

623 Hence, we obtain that

$$\begin{pmatrix} z(t)\\ \theta(t) \end{pmatrix} = \begin{pmatrix} L(t)\\ \mathrm{Id}_{\mathcal{P}} \end{pmatrix} \theta_0$$

Eventually, we have shown that for all θ_0 , $\boldsymbol{\Phi}(t,0)P^*_{\theta}(z_0,\theta_0)^t$ is given by

$$\begin{pmatrix} z(t)\\ \theta(t) \end{pmatrix} = \boldsymbol{\Phi}(t,0) P_{\theta}^* \begin{pmatrix} z_0\\ \theta_0 \end{pmatrix} = \boldsymbol{\Phi}(t,0) \begin{pmatrix} 0\\ \theta_0 \end{pmatrix} = \begin{pmatrix} L(t)\\ \mathrm{Id}_{\mathcal{P}} \end{pmatrix} \theta_0.$$

⁶²⁵ Therefore, we can simplify the Riccati dynamics (85) into

$$\begin{cases} \dot{\Lambda} = -\Lambda L^* C^* C L \Lambda, \quad t > 0, \\ \Lambda(0) = \Lambda_0, \end{cases}$$
(90)

626 and (89) into

$$\boldsymbol{\Pi}(t) = \begin{pmatrix} L(t) \\ \mathrm{Id}_{\mathcal{P}} \end{pmatrix} \Lambda(t) \begin{pmatrix} L^*(t) & \mathrm{Id}_{\mathcal{P}} \end{pmatrix}.$$
(91)

⁶²⁷ Note that, in our case $m_r(\cdot, \cdot) = \frac{\epsilon}{2}(\cdot, \cdot)_{\mathrm{H}^1_0(\Omega)}$, therefore we can obtain an explicit ⁶²⁸ expression of Λ_0 given by $\Lambda_0 = -\epsilon^{-1}\Delta_0^{-1}$.

Finally, we obtain the equivalent sequential estimator with the following result that isa straightforward application of Theorem 4.3.

631 Corollary 4.6. The minimizer $\hat{\theta} = \bar{\theta}_t$ is solution of the following dynamics

$$\begin{cases} \hat{\theta} = \Lambda L^* C^* (y^{\delta} - C\hat{z}), & t > 0, \\ \hat{\theta}(0) = \hat{\theta}_0, & \\ \dot{\hat{z}} = A\hat{z} + B\hat{\theta} + L\dot{\hat{\theta}}, & t > 0, \\ z(0) = \hat{z}_0. & \end{cases}$$
(92)

Note that the present estimator can be defined for observations in H¹ or L² with an adapted definition of the observation operator C and its adjoint C^* . Using the results obtained in Sections 3.2.2 and 3.3.2, we show that $(\hat{z}, \hat{\theta})$ cannot be entirely defined as observers in the sense of Definition 4.1, but that they satisfy the second property (robustness to noise) and a necessary condition of the first (convergence).

Proposition 4.7. Assuming that we have at our disposal the perfect measurements $y^{\delta} = \check{y}$, then, under Assumption 1, the estimators $(\hat{z}, \hat{\theta})$ asymptotically converges to the target parameter $(\check{z}, \check{\theta})$, namely

$$\lim_{t \to +\infty} \|\hat{\theta}(t) - \check{\theta}(t)\|_{\mathrm{L}^{2}(\Omega)} = 0 \text{ and } \lim_{t \to +\infty} \|\hat{z}(t) - \check{z}(t)\|_{\mathcal{Z}} = 0.$$

640 *Proof.* We here use the link between the observer and the functional minimization. 641 Under Assumptions 1 and 2, we know that (Theorem 3.6 and 3.12), for both types of 642 observation operator, there exists a > 0 and C^{st} such that

$$\|\bar{\theta}_T - \check{\theta}\|_{\mathrm{L}^2(\Omega)} \le C^{\mathrm{st}} T^{-a},$$

which, combined with $\hat{\theta} = \bar{\theta}_t$ gives the convergence of the parameter estimator

$$\lim_{t \to +\infty} \|\hat{\theta}(t) - \check{\theta}\|_{\mathcal{L}^2(\Omega)} = 0.$$

644 Moreover, we have that

$$\hat{z} = \bar{z}_t(t) = z_{|\bar{\theta}_t}.$$

645 Therefore,

$$\lim_{t \to +\infty} \|\hat{z}(t) - \check{z}(t)\|_{\mathcal{Z}} = \lim_{t \to +\infty} \|z_{|\bar{\theta}_t}(t) - z_{|\check{\theta}_t}\|_{\mathcal{Z}} = 0$$

646 which concludes the proof.

Proposition 4.8. Assuming that we have at our disposal the noisy measurements satisfying (32) or (44), the estimator $(\hat{z}, \hat{\theta})$ is stable with respect to noise, namely there exist two positive scalars C_2^{st} and C_2^{st} independent of δ and dependent of t such that

$$\|\hat{\theta}(t) - \check{\theta}\|_{L^2(\Omega)} \le C_1^{st}(t) \,\delta \text{ and } \|\hat{z} - \check{z}\|_{\mathcal{Z}} \le C_2^{st}(t) \,\delta.$$

Proof. We take advantage of the link between the observer and the functional minimization. From Theorem 3.7 and 3.13, for both types of observation operator, there exist C^{st} such that

$$\|\theta_t - \dot{\theta}\|_{\mathrm{L}^2(\Omega)} \le C^{\mathrm{st}} \,\delta$$

which, combined with $\hat{\theta}(t) = \bar{\theta}_t$ gives the stability of the parameter estimator

$$\|\hat{\theta}(t) - \check{\theta}\|_{\mathrm{L}^2(\Omega)} \le C^{\mathrm{st}} \,\delta.$$

654 Furthermore, we have that

$$\hat{z}(t) = \bar{z}_t(t) = z_{|\bar{\theta}_t}$$

655 Therefore, we obtain

$$\|\hat{z}(t) - \check{z}\|_{\mathcal{Z}} = \|z_{|\bar{\theta}_t}(t) - z_{|\check{\theta}|}\|_{\mathcal{Z}}$$

⁶⁵⁶ Classically using the Grönwall lemma, one can show that there exists a positive scalar ⁶⁵⁷ C^{st} depending only on t and σ such that

$$\|z_{|\bar{\theta}_t}(t) - z_{|\check{\theta}|}\|_{\mathcal{Z}} \le C^{\mathrm{st}} \|\bar{\theta}_t(t) - \check{\theta}\|_{\mathrm{L}^2(\Omega)},$$

⁶⁵⁸ which concludes the proof.

Remark 4.9. Proposition 4.8 allows us to demonstrate the robustness to noise property of Definition 4.1. On the other hand, Proposition 4.7 is a necessary condition for the convergence property of Definition 4.1. Indeed, the second property of the definition is obtained only in this case $\tau = 0$. We would have to take into account the possibility of uncertainties in the initial condition for dynamics (92) to fall within the asymptotic observer framework of definition 4.1.

665 4.3 Observer discretization

We can now study the observer discretization, where we face two possible options. First we could directly propose a time and space discretization of (92), or we could first discretize the wave equation and then reapply the Tikhonov regularization strategy and observer design to the discretized system. Following [27], we are going to follow this second strategy.

671 4.3.1 Reduced-order Kalman discretization

672 As in the time-continuous case, let us first recall the Kalman discretization for a linear

discrete-time dynamics. Let $\mathcal{Z}^h \subset \mathcal{V} \subset \mathcal{Z}$ a finite dimensional space equipped with the scalar product of \mathcal{Z} . Assuming a dynamics of the following form, with a given source

scalar product of \mathbf{Z} . Assuming a dynamics of the following form, with a given source term $\boldsymbol{\beta}^h \in \mathbf{Z}^h$

$$\begin{cases} \boldsymbol{z}_{n+1|\zeta}^{h} = \boldsymbol{\varPhi}^{h,\tau} \boldsymbol{z}_{n|\zeta}^{h} + \boldsymbol{\beta}_{n+1}^{h}, & n \ge 0, \\ \boldsymbol{z}_{0|\zeta}^{h} = \hat{\boldsymbol{z}}_{0}^{h} + \zeta \end{cases}$$

associated with measurements $y_n^{\delta,h}$ and an observation operator $C^h \in \mathcal{L}(\mathbf{Z}^h, \mathcal{Y}^h)$ where \mathcal{Y}^h is the finite dimensional space in which the observations belong. We first consider the minimization of a discrete version of the criterion (68) (see [57])

$$\min_{\boldsymbol{\zeta}\in\boldsymbol{\mathcal{Z}}^{h}} \Big\{ \mathscr{J}_{N}(\boldsymbol{\zeta}^{h}) = \frac{1}{2}m(\boldsymbol{\zeta}^{h},\boldsymbol{\zeta}^{h}) + \sum_{n=1}^{N} \frac{\tau}{2} \|\boldsymbol{y}_{n}^{\boldsymbol{\delta},h} - \boldsymbol{C}^{h}\boldsymbol{z}_{n|\boldsymbol{\zeta}}^{h}\|_{\mathcal{Y}^{h}}^{2} \Big\},$$
(93)

where the time integral is replaced by a simple quadrature formula. We easily find the corresponding two-ends problem associated with the minimization of this criterion,

$$\begin{cases} \bar{\boldsymbol{z}}_{n+1|N}^{h} = \boldsymbol{\varPhi}^{h,\tau} \bar{\boldsymbol{z}}_{n|N}^{h} + \boldsymbol{\beta}_{n+1}^{h}, & 0 \le n \le N-1, \\ \bar{\boldsymbol{z}}_{0|N}^{h} = \hat{\boldsymbol{z}}_{0}^{h} + \boldsymbol{\varPi}_{0}^{h} \bar{\boldsymbol{q}}_{0|N}^{h}, \\ \bar{\boldsymbol{q}}_{n|N}^{h} = \boldsymbol{\varPhi}^{h,\tau*} \bar{\boldsymbol{q}}_{n+1|N}^{h} + \tau \boldsymbol{C}^{h*} (\boldsymbol{y}_{n}^{\delta,h} - \boldsymbol{C}^{h} \bar{\boldsymbol{z}}_{n|N}^{h}), & 0 \le n \le N-1, \\ \bar{\boldsymbol{q}}_{N|N}^{h} = 0, \end{cases}$$
(94)

where the notation |N| stresses the fact that the criterion is using the N first observations and $\Pi_0^h \in \mathcal{L}(\mathcal{Z}^h)$ is the operator defined by

$$\forall (\boldsymbol{q}_h, \xi_h) \in \boldsymbol{\mathcal{Z}}^h \times \boldsymbol{\mathcal{Z}}^h, \quad m(\boldsymbol{\Pi}_0^h \boldsymbol{q}_h, \xi_h) = (\boldsymbol{q}_h, \xi_h) \boldsymbol{\mathcal{z}}.$$

⁶⁸³ Then, we can uncouple these systems by solving the following Riccati equation,

$$\begin{cases} \boldsymbol{\Pi}_{n+1^{-}}^{h} = \boldsymbol{\Phi}^{h,\tau} \boldsymbol{\Pi}_{n}^{h} \boldsymbol{\Phi}^{h,\tau*}, & n \in \mathbb{N}, \\ \boldsymbol{\Pi}_{n+1}^{h} = \boldsymbol{\Pi}_{n+1^{-}}^{h} - \tau \boldsymbol{\Pi}_{n+1^{-}}^{h} \boldsymbol{C}^{h*} [\operatorname{Id}_{\mathcal{Y}^{h}} + \tau \boldsymbol{C}^{h} \boldsymbol{\Pi}_{n+1^{-}}^{h} \boldsymbol{C}^{h*}]^{-1} \boldsymbol{C}^{h} \boldsymbol{\Pi}_{n+1^{-}}^{h}, & n \in \mathbb{N}, \end{cases}$$
(95)

which defines a sequence of operators $\boldsymbol{\Pi}_{n^-}^h$ and $\boldsymbol{\Pi}_n^h$. The sequence of $\boldsymbol{\Pi}_{n^-}^h$ and $\boldsymbol{\Pi}_n^h$ can be interpreted as a splitting time-scheme for the Riccati dynamics (73) [57]. Then, we introduce the estimator

$$\begin{cases} \hat{\boldsymbol{z}}_{n+1^{-}}^{h} = \boldsymbol{\varPhi}^{h,\tau} \hat{\boldsymbol{z}}_{n}^{h} + \boldsymbol{\beta}_{n+1}^{h}, & n \in \mathbb{N}, \\ \hat{\boldsymbol{z}}_{n+1}^{h} = \hat{\boldsymbol{z}}_{n+1^{-}}^{h} + \tau \boldsymbol{\varPi}_{n+1}^{h} \boldsymbol{C}^{h*} \left(\boldsymbol{y}_{n+1}^{\delta,h} - \boldsymbol{C}^{h} \hat{\boldsymbol{z}}_{n+1^{-}}^{h} \right), & n \in \mathbb{N}. \end{cases}$$
(96)

⁶⁸⁷ Identically, (96) can be interpreted as a splitting time-scheme for the estimator dynamics ⁶⁸⁸ (76), see [57]. Then, we have the following equivalence results proved in [57].

689 Theorem 4.10. We have the following decoupling principle

$$\bar{\boldsymbol{z}}_{n|N}^{h} = \hat{\boldsymbol{z}}_{n^{-}}^{h} + \boldsymbol{\varPi}_{n^{-}}^{h} \bar{\boldsymbol{q}}_{n|N}^{h}, \quad \forall 0 \le n \le N+1,$$
(97)

690 and in particular for all $n \in \mathbb{N}$

$$\hat{\boldsymbol{z}}_{n-}^h = \bar{\boldsymbol{z}}_{n|n}^h. \tag{98}$$

Let us now apply such equivalence to a reduced space configuration. For this purpose, we define $\mathcal{Z}_r^h \subset \mathcal{Z}^h$ a finite dimensional space such that $\mathcal{Z}_r^h \subset \mathcal{V}_r \subset \mathcal{Z}_r$. To simplify the notation, we assume that $P_r(\mathcal{Z}^h) \subset \mathcal{Z}_r^h$ – which is satisfied for our source identification problem – hence P_r can be seen as an operator in $\mathcal{L}(\mathcal{Z}^h, \mathcal{Z}_r^h)$. The discrete criterion here reads

$$\min_{\boldsymbol{\zeta}_r^h \in \boldsymbol{\mathcal{Z}}_r^h} \Big\{ \mathscr{J}_N(\boldsymbol{\zeta}_r^h) = \frac{1}{2} m_r(\boldsymbol{\zeta}_r^h, \boldsymbol{\zeta}_r^h) + \sum_{k=1}^N \frac{\tau}{2} \|\boldsymbol{y}_n^{\delta,h} - \boldsymbol{C}\boldsymbol{z}_{n|P_r^*\boldsymbol{\zeta}_r^h}^h\|_{\mathcal{Y}^h}^2 \Big\}.$$
(99)

The two ends problem remains as in (94) except for the initial condition that is changed into

$$\bar{\boldsymbol{z}}_{0|N}^{h} = \hat{\boldsymbol{z}}_{0^{-}}^{h} + P_{r}^{*} (U_{0}^{h})^{-1} P_{r} \bar{\boldsymbol{q}}_{0|N}^{h}, \qquad (100)$$

 $_{698}$ Moreover, we have the discrete counterpart of Theorem 4.5.

Theorem 4.11. The solution of (95) with initial condition defined by

$$\boldsymbol{\Pi}_0^h = P_r^* \Lambda_0^h P_r$$

700 satisfies

$$\boldsymbol{\varPi}_{n}^{h} = \left(\boldsymbol{\varPhi}^{h,\tau}\right)^{n} P_{r}^{*} \Lambda_{n}^{h} P_{r} \left(\boldsymbol{\varPhi}^{h,\tau*}\right)^{n}$$
(101)

701 where Λ^h_n satisfies $\Lambda^h_n = (U^h_n)^{-1}$ with

$$U_{n+1}^{h} = U_{n}^{h} + \tau P_{r} \left(\boldsymbol{\varPhi}^{h,\tau*}\right)^{n+1} \boldsymbol{C}^{h*} \boldsymbol{C}^{h} \left(\boldsymbol{\varPhi}^{h,\tau}\right)^{n+1} P_{r}^{*}, \quad n \ge 0,$$
(102)

 $_{702}$ Proof. We proceed by induction. We have at iteration 0

$$\boldsymbol{\Pi}_{0}^{h} = P_{r}^{*}(U_{0}^{h})^{-1}P_{r}.$$

Let us assume (101) is satisfied at iteration n. First, we directly have

$$\boldsymbol{\varPi}_{n+1^{-}}^{h} = \boldsymbol{\varPhi}^{h,\tau} \boldsymbol{\varPi}_{n}^{h} \boldsymbol{\varPhi}^{h,\tau*} = \left(\boldsymbol{\varPhi}^{h,\tau}\right)^{n+1} P_{r}^{*} \Lambda_{n}^{h} P_{r} \left(\boldsymbol{\varPhi}^{h,\tau*}\right)^{n+1},$$

Then, using (95) and the expression above of $I\!\!I_{n+1^{-}}^{h}$, we have

$$\boldsymbol{\Pi}_{n+1}^{h} = \left(\boldsymbol{\Phi}^{h,\tau}\right)^{n+1} P_{r}^{*} \left[\Lambda_{n}^{h} - \tau \Lambda_{n}^{h} P_{r} \left(\boldsymbol{\Phi}^{h,\tau*}\right)^{n+1} \boldsymbol{C}^{h*} \left[\operatorname{Id}_{\mathcal{Y}^{h}} + \tau \boldsymbol{C}^{h} \left(\boldsymbol{\Phi}^{h,\tau}\right)^{n+1} \right. \\ \left. \times P_{r}^{*} \Lambda_{n}^{h} P_{r} \left(\boldsymbol{\Phi}^{h,\tau*}\right)^{n+1} \boldsymbol{C}^{h*} \right]^{-1} \boldsymbol{C}^{h} \left(\boldsymbol{\Phi}^{h,\tau}\right)^{n+1} P_{r}^{*} \Lambda_{n}^{h} \right] P_{r} \left(\boldsymbol{\Phi}^{h,\tau}\right)^{n+1}.$$

704 We therefore have

$$\boldsymbol{\Pi}_{n+1}^{h} = \left(\boldsymbol{\varPhi}^{h,\tau}\right)^{n+1} P_{r}^{*} \Lambda_{n+1}^{h} P_{r} \left(\boldsymbol{\varPhi}^{h,\tau}\right)^{n+1},$$

as soon as

$$\Lambda_{n+1}^{h} = \Lambda_{n}^{h} - \Lambda_{n}^{h} D_{n}^{h,*} \left(\mathrm{Id}_{\mathcal{Y}^{h}} + D_{n}^{h} \Lambda_{n}^{h} D_{n}^{h,*} \right)^{-1} D_{n}^{h} \Lambda_{n}^{h}$$

with $D_n^h = \sqrt{\tau} C^h \left(\boldsymbol{\Phi}^{h,\tau} \right)^{n+1} P_r^*$. By using Sherman–Morrison–Woodbury formula for operators – see for instance [40, Theorem 1.1] – we obtain

$$(\Lambda_{n+1}^{h})^{-1} = (\Lambda_{n}^{h})^{-1} + D_{n}^{h,*}D_{n}^{h}$$

which ends the proof since it is precisely (102).

708 4.3.2 Application to source identification

⁷⁰⁹ In this work, we discretize the wave equation using Lagrange finite element. Let $\theta^h \in \mathcal{P}^h$

the finite dimensional space corresponding to \mathcal{P} . We use a conservative mid-point scheme, that ultimately reads

$$\begin{cases} \frac{u_{n+1}^{h} - u_{n}^{h}}{\tau} = \frac{v_{n+1}^{h} + v_{n}^{h}}{2} \\ \frac{v_{n+1}^{h} - v_{n}^{h}}{\tau} - \Delta_{0}^{h} \frac{u_{n+1}^{h} + u_{n}^{h}}{2} = \sigma \left(\left(n + \frac{1}{2} \right) \tau \right) \theta^{h}, \end{cases}$$
(103)

where Δ_0^h is the bounded linear operator in $\mathcal{V}^h \subset \mathrm{H}^1(\Omega)$ associated with the conform finite element discretization of the unbounded operator Δ_0 of $\mathrm{L}^2(\Omega)$. Note that this system can be rewritten in a first order form by considering $z_n^h \in \mathcal{Z}^h = \mathcal{V}^h \times \mathcal{V}^h$ solution of

$$z_{n+1}^h = \Phi^{h,\tau} z_n^h + B_{n+1}^{h,\tau} \theta^h$$

716 where

$$\Phi^{h,\tau} = \begin{pmatrix} \mathrm{Id}^h & -\frac{\tau}{2}\mathrm{Id}^h \\ -\frac{\tau}{2}\Delta_0^h & \mathrm{Id}^h \end{pmatrix}^{-1} \begin{pmatrix} \mathrm{Id}^h & \frac{\tau}{2}\mathrm{Id}^h \\ \frac{\tau}{2}\Delta_0^h & \mathrm{Id}^h \end{pmatrix}$$

717 and

$$B_{n+1}^{h,\tau} = \tau \begin{pmatrix} \mathrm{Id}^h & -\frac{\tau}{2} \mathrm{Id}^h \\ -\frac{\tau}{2} \Delta_0^h & \mathrm{Id}^h \end{pmatrix}^{-1} \begin{pmatrix} 0 \\ \sigma \left(\left(n + \frac{1}{2} \right) \tau \right) \end{pmatrix}$$

⁷¹⁸ Note that \mathcal{Z}^h is equipped with the scalar product on \mathcal{Z} that is the natural scalar ⁷¹⁹ product on $\mathrm{H}^1_0(\Omega) \times \mathrm{L}^2(\Omega)$.

We now apply this result to our source identification problem. The joint stateparameter discrete dynamics in $\mathbf{Z}^h = \mathbf{Z}^h \times \mathbf{\mathcal{P}}^h$ reads

$$\underbrace{\begin{pmatrix} z_{n+1}^h\\ \theta_{n+1}^h\\ z_{n+1}^h \end{pmatrix}}_{\boldsymbol{z}_{n+1}^h} = \underbrace{\begin{pmatrix} \Phi^{h,\tau} & B_{n+1}^{h,\tau}\\ 0 & \mathrm{Id}^h \end{pmatrix} \begin{pmatrix} z_n^h\\ \theta_n^h \end{pmatrix}}_{\boldsymbol{\Phi}^{h,\tau} \boldsymbol{z}_n^h}, \quad \text{with} \quad \underbrace{\begin{pmatrix} z_0^h\\ \theta_0^h \end{pmatrix}}_{\boldsymbol{z}_0^h} = \underbrace{\begin{pmatrix} 0\\ \theta^h \end{pmatrix}}_{\boldsymbol{\zeta}}. \tag{104}$$

⁷²² From Theorem 4.5, we know that

$$\begin{cases} \boldsymbol{\varPi}_{n}^{h} = \left(\boldsymbol{\varPhi}^{h,\tau}\right)^{n} P_{r}^{*} \Lambda_{n}^{h} P_{r} \left(\boldsymbol{\varPhi}^{h,\tau*}\right)^{n} \text{ and } \boldsymbol{\varPi}_{n+1^{-}}^{h} = \left(\boldsymbol{\varPhi}^{h,\tau}\right)^{n+1} P_{r}^{*} \Lambda_{n}^{h} P_{r} \left(\boldsymbol{\varPhi}^{h,\tau*}\right)^{n+1}, \\ \boldsymbol{\varPi}_{0} = P_{\theta}^{*} (U_{0}^{h})^{-1} P_{\theta}, \end{cases}$$

$$(105)$$

⁷²³ where $U_n^h = (\Lambda_n^h)^{-1}$ is solution of

$$U_{n+1}^{h} = U_{n}^{h} + \tau P_{\theta} \left(\boldsymbol{\Phi}^{h,\tau*} \right)^{n+1} C^{h*} C^{h} \left(\boldsymbol{\Phi}^{h,\tau} \right)^{n+1} P_{\theta}^{*}, \quad n \ge 0.$$
(106)

724 Moreover, we have

$$\begin{pmatrix} z_n^h \\ \theta_n^h \end{pmatrix} = \left(\boldsymbol{\varPhi}^{h,\tau}\right)^n P_{\theta}^* \begin{pmatrix} 0 \\ \theta_0^h \end{pmatrix} = \left(\boldsymbol{\varPhi}^{h,\tau}\right)^n \begin{pmatrix} 0 \\ \theta_0^h \end{pmatrix} = \begin{pmatrix} L_n^h \\ \mathrm{Id}^h \end{pmatrix} \theta_0^h$$

where the operator $(L_n^h)_{n \in \mathbb{N}}$ is defined by

$$L_{n}^{h}\theta = z_{n|\theta}^{h} \quad \text{with} \quad \begin{cases} z_{n+1|\theta}^{h} = \Phi^{h,\tau} z_{n|\theta}^{h,\tau} + B_{n+1}^{h,\tau}\theta, & n \in \mathbb{N}, \\ z_{0|\theta}^{h} = 0. \end{cases}$$
(107)

Therefore, we can simplify (106) into

$$U_{n+1}^{h} = U_{n}^{h} + \tau L_{n+1}^{h*} C^{h*} C^{h} L_{n+1}^{h}, \quad n \ge 0,$$
(108)

727 and (105) into

$$\boldsymbol{\varPi}_{n}^{h} = \begin{pmatrix} L_{n}^{h} \\ \mathrm{Id} \end{pmatrix} \Lambda_{n}^{h} \begin{pmatrix} L_{n}^{h*} & \mathrm{Id} \end{pmatrix}, \text{ and } \boldsymbol{\varPi}_{n+1^{-}}^{h} = \begin{pmatrix} L_{n+1}^{h} \\ \mathrm{Id}^{h} \end{pmatrix} \Lambda_{n}^{h} \begin{pmatrix} L_{n+1}^{h*} & \mathrm{Id}^{h} \end{pmatrix}.$$
(109)

⁷²⁸ Finally, we obtain the equivalent sequential estimator.

Theorem 4.12. The minimizer $\hat{\theta}_n^h = \bar{\theta}_n^h$ is solution from $(\hat{z}_0^h, \hat{\theta}_0^h)$ of the recursive dynamics

$$\begin{cases} \hat{z}_{n+1^{-}}^{h} = \Phi^{h,\tau} z_{n}^{h} + B_{n+1}^{h,\tau} \hat{\theta}_{n}^{h}, \\ \hat{\theta}_{n+1}^{h} = \hat{\theta}_{n}^{h} + \tau \left(U_{n+1}^{h} \right)^{-1} L_{n+1}^{h*} C^{h*} (y_{n+1}^{\delta,h} - C^{h} \hat{z}_{n+1^{-}}^{h}), \\ \hat{z}_{n+1}^{h} = \hat{z}_{n+1^{-}}^{h} + L_{n+1}^{h} (\hat{\theta}_{n+1}^{h} - \hat{\theta}_{n}^{h}). \end{cases}$$
(110)

We already know that (110) is consistent. A direct convergence study of this splitting scheme could be intricate. Here, by relying on the fundamental relation $\hat{\theta}_n^h = \bar{\theta}_n^h$ we can justify the stability from the minimization. This should lead us to the convergence of the time-scheme (110).

735 **5** Numerical illustrations

⁷³⁶ We consider a domain Ω , and assume we have access to observations in the observation ⁷³⁷ subdomain ω illustrated in Figure 2. As an example, we consider the following wave ⁷³⁸ equation

$$\begin{cases} \partial_{tt}\check{u}(x,t) - \Delta\check{u}(x,t) = \sigma(t)\check{\theta}(x), & (x,t) \in \Omega \times (0,T), \\ \check{u}(x,t) = 0, & (x,t) \in \partial\Omega \times (0,T), \\ \check{u}(x,0) = \check{u}_0(x), & x \in \Omega, \\ \partial_t\check{u}(x,0) = 0, & x \in \Omega, \end{cases}$$
(111)

$$\sigma(t) = \frac{1}{2} + \frac{1}{2}e^{-0.5t},$$

740 and \check{u}_0 is solution of

$$\begin{cases} -\Delta \check{u}_0(x) = 10, & x \in \Omega, \\ \check{u}_0(x) = 0, & x \in \partial \Omega, \end{cases}$$
(112)

⁷⁴¹ and $\check{\theta}$ is solution of

$$\begin{aligned}
-\Delta\theta(x) &= 0 & x \in \Omega \setminus \Omega_{\text{hole}} \\
-\Delta\check{\theta}(x) &= 10, & x \in \Omega_{\text{hole}}, \\
\check{\theta}(x) &= 0, & x \in \partial\Omega,
\end{aligned}$$
(113)

with Ω_{hole} represented in Figure 2.

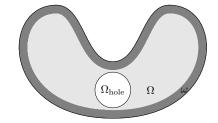


Figure 2. Representation of the domain Ω and the subdomain ω

742

The initial conditions and the target source term defined respectively by (112) and the (113) are illustrated in Figure 3 (left) and 3 (right). Note that the choice of σ is the

same as that of Remark 3.5 with $\beta = 1$ and $\alpha = 0.5$.

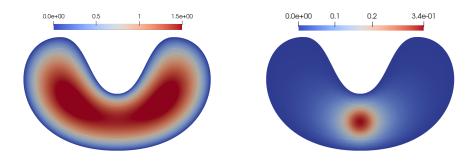


Figure 3. Representation of \check{u}_0 (left) and θ (right)

747 We decompose the target source term

$$\check{\theta}^h := \sum_{i=0}^{N_{\text{mod}}} \check{\theta}^h_i \phi^h_i, \qquad (114)$$

where $(\phi_i^h)_{i \in [0, N_{\text{mod}}]}$ are the eigenvectors of Δ_0^h . In the following numerical illustrations, we choose $N_{\text{mod}} = 27$ so that the relative error between $\check{\theta}$ and $\check{\theta}^h$ is less than 10%, see

• we choose $N_{\text{mod}} = 27$ so that the relative error between $\check{\theta}$ and $\check{\theta}^h$ is less than 10%, see Figure 4.

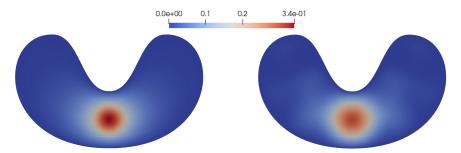


Figure 4. Source term $\check{\theta}$ (left) and its decomposition $\check{\theta}^h$ on the Laplacian basis with 27 modes (right)

750

In the following, we will compare the two methods presented in this work. The method presented is denoted H^1 -method (the noisy measurements belong to $H^1(\omega)$). Similarly, the method presented in Section 3.3 is denoted L^2 -method : the noisy measurements belong to $L^2(\omega)$.

755 5.1 Numerical illustrations for noise-free data

First, to generate the measurements, we implement the direct problem on a time window from 0 to 1 with a time step of $\tau = 0.01$ and with \mathbb{P}_1 finite elements and where the space function $\check{\theta}^h$ to reconstruct in the source term is defined (114). We give an illustration of a measurement in Figure 5.

760

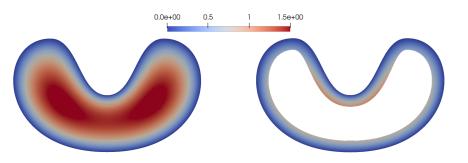


Figure 5. Solution \check{u}_n^h of the direct problem (left) and its corresponding measurements $(\check{u}_n^h)_{|\omega}$ (right) at time $t^n = 0.1$

We study the inverse problem for the H^1 -method in a noise-free case. We implement 761 the observers scheme (110) (where C is defined by (7)) in Freefem++ [42] over a time 762 window [0,1] with a time step of 0.01 and with \mathbb{P}_1 finite elements. Moreover, we start 763 from a null initial guess $\theta_0^h = 0$ but with the correct initial conditions \check{u}_0 . Furthermore, 764 we remove the Tikhonov regularization term since we suppose that we have the perfect 765 measurements at our disposal. The same simulation for the L^2 -method is performed. 766 In other words, scheme 110 is implemented, but this time the observation operator C767 is defined by (9). 768

Figure 6 shows the relative error between the target source term decomposed on the Laplacian modes and the observer defined by

$$\|\check{\theta}^h - \hat{\theta}^h_n\|_{\mathrm{L}^2(\Omega)} / \|\check{\theta}^h\|_{\mathrm{L}^2(\Omega)}$$

771 over the time.

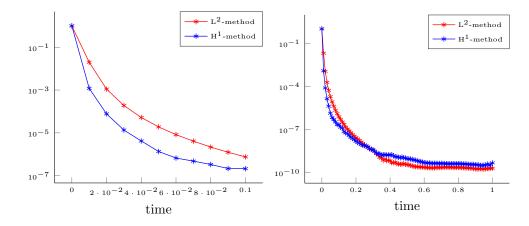


Figure 6. Evolution of the relative error with $\delta = 1$ and without Tikhonov term for T = 0.1 (left) and T = 1 (right)

For the H^1 -method, we can see that the error falls below $10^{-5}\%$ from $t^n = 0.1$. At the final time $t^n = 1$, the relative error is about $10^{-8}\%$. For the L^2 -method, the results are very similar. We can see that both methods converge for long time with the same final error level. However, in the short time range (Figure 6 left), the H^1 -method is more efficient than the L^2 -method. This indicates that the H^1 -method benefits from improved observability.

778

We will now deal with noisy measurements belonging to $L^2(\omega)$.

780 5.2 Numerical illustrations with noisy measurements

We represent the mesh in Figure 7 (left) used to solve the inverse problem. For visualization purposes, we show a coarse mesh (Figure 7 right), the used mesh is a refined version of the coarse one (Figure 7 left).



Figure 7. Representation of the mesh of Ω with 56193 degrees of freedom (left) and with 14169 degrees of freedom

- We corrupt the synthetic data $(y^{\delta,h})_{0 \le n \le N}$ of sampling rate T/N with a random artificial noise $\chi_n^h \in L^2(\omega)$. To do this, we perturb the measurements at each interpolation point with a random value ξ that follows the Gaussian law $\mathcal{N}(0,1)$ and scale with
- 788 the noise level δ , *i.e.*, at each interpolation point we have

$$(y_n^{\delta,h})_i = (\check{y}_n^h)_i + \delta\xi_i / |\omega|, \qquad (115)$$

where the scaling by $|\omega|$ ensures that the noise is normalized. We represent in Figure 8 an example of a noisy measurement and the noise associated.

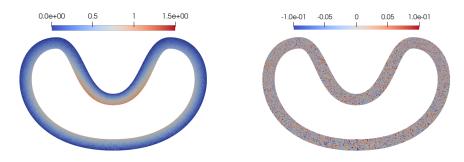


Figure 8. Noisy measurements $y_n^{\delta,h}$ (left) and the corresponding noise $\chi_n^h = y_n^{\delta,h} - \check{y}_n^h$ (right) at time $t^n = 0.1$

790

We solve the inverse problem with the *prediction-correction* scheme (110) with a time step of $\tau = 0.01$ and $\delta = 1$. Figure 9 shows the relative error

$$\frac{\|\check{\theta}^h - \hat{\theta}^h_n\|_{\mathrm{L}^2(\Omega)}}{\|\check{\theta}^h\|_{\mathrm{L}^2(\Omega)}}$$

between the target source term decomposed on the basis and the observer at the final time $t^n = 1$ as a function of the noise level δ in (115) for different regularization levels M.

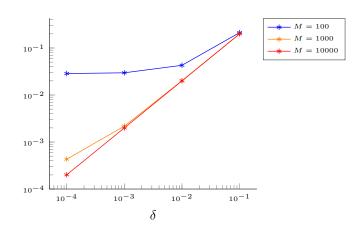


Figure 9. Representation of the relative error according the values of δ in loglog scale for different values of M

796

First, we note that the larger the value of M, the lower the weight of the regularization. In particular, we can see that the regularization with $M = 10^4$ (in red in Figure 9) offers better results. In particular, for the different noise levels considered – from $\delta = 0.01\%$ up to $\delta = 10\%$ in (115) – we reach a reconstruction of the order of the noise level.

802 6 Conclusion

We have proposed a strategy to estimate a source term in a wave equation using 803 observers from internal *field* measurements. We have divided the study into two parts 804 depending on whether we have access to regular or less regular measurements. The 805 regularity of the measurements leads to major changes in the definition of the inverse 806 problem through the definition of the observation operator. The proposed approach 807 follows three steps. First we prove identifiability results which rely on observability 808 inequalities. These identifiability results give us information about the degree of the 809 ill-posedness of the inverse problem. We then propose the minimization problem 810 associated with the reconstruction of the source term with a well-adapted Tikhonov 811 regularization. This inverse problem is solved using a sequential strategy that is proven 812 to be equivalent to the minimization of the cost functional. We have assumed that we 813 have not only internal observations on the system, but also perfect knowledge of the 814 initial conditions, which can be a strong assumption in practice. A possible perspective 815 would be to assume that there are uncertainties on the initial conditions. In this case, 816 it would be necessary to further analyze the joint state-source estimation initially 817 proposed in [19]. 818

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