

IDENTIFIABILITY FOR COMPOSITE STRING VIBRATION PROBLEM

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ABSTRACT. The paper considers the identifiability (i.e., the unique identification) of a composite string in the class of piecewise constant parameters. The 1-D string vibration is measured at finitely many observation points. The observations are processed to obtain the first eigenvalue and a constant multiple of the first eigenfunction at the observation points. It is shown that the identification by the Marching Algorithm is continuous with respect to the mean convergence in the admissible set. The result is based on the continuous dependence of eigenvalues, eigenfunctions, and the solutions on the parameters. A numerical algorithm for the identification in the presence of noise is proposed and implemented.

1. Introduction

Consider the transverse vibration of a thin string of length 1 stretched with a unit force. Let the string be made of finitely many uniform pieces. Let its density be $\rho(x)$, and the piecewise constant parameter $a(x) = 1/\rho(x)$ satisfy $0 < \nu \leq a(x) \leq \mu$, $x \in [0, 1]$ with $a(x) = a_i$ for $x \in [x_{i-1}, x_i)$, $i = 1, 2, \dots, N$. Then the string's displacement from the position of equilibrium $u(x, t)$, $0 \leq x \leq 1$, $t > 0$ can be modeled by

$$(1.1) \quad \begin{cases} u_{tt} - (a(x)u_x)_x = f(x, t), & x \neq x_i, \quad t \in (0, T), \\ u(0, t) = q_1(t), \quad u(1, t) = q_2(t), & t \in (0, T), \\ u(x_i+, t) = u(x_i-, t), \\ a(x_i+)u_x(x_i+, t) = a(x_i-)u_x(x_i-, t), \\ u(x, 0) = g(x), \quad u_t(x, 0) = h(x), & x \in (0, 1). \end{cases}$$

The parameter *identification* problem for (1.1) is to find variable parameters a, f, q_1, q_2 and g such that the solution $u(x, t)$ fits given observations z in a prescribed (e.g., the best fit to data) sense.

The *identifiability* problem for (1.1) is to establish the uniqueness of the above identification.

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One can view such a model as a special case of the remote sensing problem. When identifiability results are available, the system can be designed accordingly.

Our main results are contained in Section 3. There we consider the string vibration problem

$$(1.2) \quad \begin{cases} u_{tt} - (a(x)u_x)_x = 0, & x \neq x_i, \quad t \in (0, T), \\ u(0, t) = 0, \quad u(1, t) = 0, & t \in (0, T), \\ u(x_i+, t) = u(x_i-, t), \\ a(x_i+)u_x(x_i+, t) = a(x_i-)u_x(x_i-, t), \\ u(x, 0) = g \in H_0^1[0, 1], \quad u_t(x, 0) = h \in L_2[0, 1] \end{cases}$$

and show that the piecewise constant parameter $a(x)$ can be uniquely identified from finitely many observation functions $z_m(t) = u(p_m, t)$, $t > 0$, $m = 1, 2, \dots, M - 1$. These observations should be taken at a sufficiently dense in $(0, 1)$ set of equidistant points p_m as specified in Theorem 3.3. Also in this section we show that the identification is stable, i.e., the recovered parameters a depend continuously on the observations $z_m(t)$. These results use certain properties of solutions of (1.2) established in Section 2.

The main idea of our identification method is to use the Fourier transform of the data $z_m(t)$ to find the first eigenvalue λ_1 of the associated Sturm-Liouville problem and a constant multiple $C\psi_1(p_m)$ of the first eigenfunction at the observation points. The resulting sequence of M numbers (denoted by \mathcal{G}) provides the input to the Marching Algorithm (see [7]) which uniquely recovers the sought parameter $a(x)$.

In Section 4 the identifiability results are generalized to problems with nonzero boundary inputs $q_1(t)$ and $q_2(t)$, as well as for a nonzero external input $f(x, t)$. In Section 5 we present numerical results illustrating the identification of a piecewise constant parameter $a(x)$ from the observations $z_m(t)$ of the system.

Some identifiability results for smooth or constant parameters a were obtained previously, see [14, 15, 16]. These works show that one can identify a constant parameter a in (1.2) from the measurement $z(t)$ taken at one point $p \in (0, 1)$. These works also discuss problems more general than (1.2), including problems with a broad range of boundary conditions, non-zero forcing functions, as well as elliptic and parabolic problems.

In [12], [4] and references therein identifiability results are obtained for elliptic and parabolic equations with discontinuous parameters in a multidimensional setting. A typical assumption there is that one knows the normal derivative of the solution at the boundary of the region for every Dirichlet boundary input. Different approaches can be found in [1, 3] and in more comprehensive treatments [2, 10, 11].

In [7] and [8] we studied the conductivity identifiability problems and obtained some fundamental results (the Marching Algorithm, the continuity of the

eigenvalues and the eigenfunctions) critical for the string vibration problems studied in this paper.

2. Properties of solutions

Let the admissible set be defined by

$$A_{\text{ad}} = \{a \in L^\infty[0, 1] : 0 < \nu \leq a(x) \leq \mu\}$$

for some positive constants ν and μ . Admissible sets of piecewise constant parameters are defined as follows.

Definition 2.1.

- (1) Define $\mathcal{PC} = \cup_{N=1}^\infty \mathcal{PC}_N \subset A_{\text{ad}}$, where \mathcal{PC}_N consists of piecewise constant functions $a(x)$ of the form $a(x) = a_i$ for $x \in [x_{i-1}, x_i)$, $i = 1, 2, \dots, N$.
- (2) Let $\sigma > 0$. Define

$$\mathcal{PC}(\sigma) = \{a \in \mathcal{PC} : x_i - x_{i-1} \geq \sigma, \quad i = 1, 2, \dots, N\},$$

where x_1, x_2, \dots, x_{N-1} are the discontinuity points of a , and $x_0 = 0, x_N = 1$.

Note that $a \in \mathcal{PC}(\sigma)$ attains at most $N = \lceil 1/\sigma \rceil$ distinct values $a_i, 0 < \nu \leq a_i \leq \mu$.

Theorem 2.2. *Let $a \in \mathcal{PC}$. Then*

- (1) *The associated Sturm-Liouville problem for (1.1)*

$$(2.1) \quad \begin{cases} (a(x)\psi(x)')' = -\lambda\psi(x), & x \neq x_i, \\ \psi(0) = \psi(1) = 0, \\ \psi(x_i+) = \psi(x_i-), \\ a(x_i+)\psi_x(x_i+) = a(x_i-)\psi_x(x_i-) \end{cases}$$

has infinitely many eigenvalues

$$0 < \lambda_1 < \lambda_2 < \dots \rightarrow \infty.$$

The normalized eigenfunctions $\{\psi_k\}_{k=1}^\infty$ form a basis in $L^2[0, 1]$.

- (2) *Each eigenvalue is simple. For each eigenvalue λ_k there exists a unique continuous, piecewise smooth normalized eigenfunction $\psi_k(x)$ such that $\psi_k'(0+) > 0$, and the function $a(x)\psi_k'(x)$ is continuous on $[0, 1]$. Also $\psi_1(x) > 0$ for $x \in (0, 1)$.*
- (3) *Eigenvalues $\{\lambda_k\}_{k=1}^\infty$ satisfy the inequality*

$$\nu\pi^2k^2 \leq \lambda_k \leq \mu\pi^2k^2.$$

These and other properties of the eigenvalues and the eigenfunctions of (2.1) follow from standard arguments, see e.g. [5]. A more detailed derivation is presented in [7] and [8].

Let $H = L_2[0, 1]$ with the norm $\|\cdot\|$ and the inner product $\langle \cdot, \cdot \rangle$. Let $V = H_0^1[0, 1]$. Then (1.2) can be restated as an abstract differential equation within the standard triple $V \subset H \subset V'$

$$(2.2) \quad u'' + Au = 0, \quad u(0) = g \in V, \quad u'(0) = h \in H.$$

According to [13, 5], equation (2.2) has a unique weak solution $u \in C([0, T]; V)$. We need a somewhat more detailed information about it.

Theorem 2.3. *Let $a \in \mathcal{PC}$, $h \in H$, $g \in V$. Then on $D = [0, 1] \times [0, \infty)$ the weak solution u of (2.2) is given by*

$$(2.3) \quad u(x, t) = \sum_{k=1}^{\infty} (A_k \cos(\sqrt{\lambda_k}t) + B_k \sin(\sqrt{\lambda_k}t))\psi_k(x)$$

with

$$(2.4) \quad A_k = \langle g, \psi_k \rangle, \quad B_k = \frac{\langle h, \psi_k \rangle}{\sqrt{\lambda_k}}, \quad k = 1, 2, \dots$$

Moreover, the series in (2.3) converges uniformly on D , and the solution $u(x, t)$ is a continuous and bounded function in D .

Proof. The partial sums in (2.3) are Galerkin approximations for the solution of (2.2). Their weak convergence to the solution $u(x, t)$ is established using the energy estimate as in [13] or [5]. Under our assumptions it takes the form of the energy conservation

$$(2.5) \quad \int_0^1 (a(x)|u_x(x, t)|^2 + |u_t(x, t)|^2)dx = \int_0^1 (a(x)|g'(x)|^2 + |h(x)|^2)dx$$

for each $t \geq 0$. We are going to give an explicit derivation of this equality as well as to show the uniform convergence of the involved series.

Recall that the eigenfunctions $\{\psi_k(x)\}_{k=1}^{\infty}$ form an orthonormal basis in H . Let the coefficients A_k and B_k be defined as in (2.4). Then

$$\sum_{k=1}^{\infty} |A_k|^2 = \|g\|_H^2, \quad \sum_{k=1}^{\infty} \lambda_k |B_k|^2 = \|h\|_H^2.$$

Also functions $\{\psi_k/\sqrt{\lambda_k}\}_{k=1}^{\infty}$ form an orthonormal basis in V_a , i.e., in V equipped with the equivalent inner product $\langle av', w' \rangle$ for $v, w \in V$, see Theorem 6.5.2 in [5]. Thus

$$\sum_{k=1}^{\infty} \lambda_k |A_k|^2 = \|g\|_{V_a}^2 = \int_0^1 a(x)|g'(x)|^2 dx.$$

Fix $t > 0$. Then

$$\|u^N(x, t)\|_{V_a}^2 = \int_0^1 a(x) \left| \sum_{k=1}^N (A_k \cos(\sqrt{\lambda_k}t) + B_k \sin(\sqrt{\lambda_k}t))\psi'_k(x) \right|^2 dx$$

$$\begin{aligned} &= \sum_{k=1}^N \lambda_k |A_k \cos(\sqrt{\lambda_k t}) + B_k \sin(\sqrt{\lambda_k t})|^2, \\ \|u_t^N(x, t)\|_H^2 &= \int_0^1 \left| \sum_{k=1}^N \sqrt{\lambda_k} (-A_k \sin(\sqrt{\lambda_k t}) + B_k \cos(\sqrt{\lambda_k t})) \psi_k(x) \right|^2 dx \\ &= \sum_{k=1}^N \lambda_k | -A_k \sin(\sqrt{\lambda_k t}) + B_k \cos(\sqrt{\lambda_k t})|^2. \end{aligned}$$

Thus

$$(2.6) \quad \int_0^1 (a(x)|u_x^N(x, t)|^2 + |u_t^N(x, t)|^2) dx = \sum_{k=1}^N \lambda_k (A_k^2 + B_k^2)$$

and the series in (2.3) converges in V_a uniformly with respect to $t > 0$. Since $\nu \|v\|_V^2 \leq \|v\|_{V_a}^2$ and $C[0, 1]$ is continuously imbedded in V it follows that the series in (2.3) converges uniformly on D , and the limit $u(x, t)$ is continuous. Also the norm $\|u(x, t)\|_V$ is uniformly bounded for $t > 0$, thus $u(x, t)$ is bounded on D . Finally, one obtains (2.5) by taking the limit as $N \rightarrow \infty$ in (2.6). \square

3. Identification map and its continuity

This section contains our main results. Here is an outline of our arguments. Let $a \in \mathcal{PC}$ and $p \in (0, 1)$. Then the string motion at the point p is given by

$$(3.1) \quad z_p(t; a) = u(p, t; a) = \sum_{k=1}^{\infty} (A_k \cos(\sqrt{\lambda_k t}) + B_k \sin(\sqrt{\lambda_k t})) \psi_k(p; a), \quad t \geq 0$$

with the coefficients A_k and B_k defined in (2.4). By Theorem 2.3, $z_p(t; a)$ is a bounded continuous function for $t \geq 0$.

Our goal is to extract the first eigenvalue $\lambda_1(a)$ and a constant multiple of the first eigenfunction $C(a)\psi_1(p; a)$ from $z_p(t; a)$, with $C(a)$ being independent of the point $p \in (0, 1)$. This is accomplished in Theorem 3.1 by applying the Fourier transform to the even or odd extension of $z_p(t; a)$. Given $M - 1$ observation points p_m this procedure defines the solution map $\mathcal{G}(a)$ by

$$(3.2) \quad \mathcal{G}(a) = (\lambda_1(a), G_1(a), \dots, G_{M-1}(a)) \in \mathbb{R}^M, \quad a \in \mathcal{PC},$$

where $G_m(a) = C(a)\psi_1(p_m; a)$. Note that the M -tuple $\mathcal{G}(a)$ is determined from the observations $z_m(t; a) = u(p_m, t; a)$ without the knowledge of the parameter a .

Suppose that one observes the system (1.2) at sufficiently many observation points p_m , $m = 1, 2, \dots, M - 1$ equidistant on the interval $(0, 1)$ as specified in Theorem 3.3. This theorem asserts that the solution map $\mathcal{G}(a)$ can be inverted on $\mathcal{PC}(\sigma)$, i.e., the piecewise constant parameter a can be uniquely identified from $\mathcal{G}(a)$.

This inversion procedure is accomplished by the Marching Algorithm described and justified in [7]. The Marching Algorithm ([7], Theorem 4.6) is applied there to the identifiability of piecewise constant conductivities in a heat conduction problem. Since the solution of the heat conduction problem is different from (3.1), a different procedure is used there to extract the M -tuple $\mathcal{G}(a)$ from the data. However, the application of the Marching Algorithm to $\mathcal{G}(a)$ is the same.

Finally in this section, we show that \mathcal{G} is continuous (as a function of a) if $\mathcal{PC}(\sigma)$ is equipped with the $L_1[0, 1]$ topology. Observing that $\mathcal{PC}(\sigma)$ is compact in this topology we can conclude that the identification map \mathcal{G}^{-1} is continuous. This means that the identification procedure is stable.

Theorem 3.1. *Let $a, b \in \mathcal{PC}$, $g \in V$, $h \in H$. Let $p \in (0, 1)$ and $z_p(t; a)$ be defined by $z_p(t; a) = u(p, t; a)$, $t \geq 0$, where $u(x, t; a)$ is the weak solution of (1.2).*

- (1) *Suppose that $h = 0$. Then the Fourier transform of the even extension $z_{p,even}(t; a)$ of $z_p(t; a)$ is given by*

$$\begin{aligned}
 (3.3) \quad & \mathcal{F}(z_{p,even}(t; a)) \\
 &= \frac{1}{\sqrt{2\pi}} \int_{-\infty}^{\infty} z_{p,even}(t; a) e^{-iwt} dt \\
 &= \sum_{k=1}^{\infty} \sqrt{\frac{\pi}{2}} \left[A_k \delta(w - \sqrt{\lambda_k}) + A_k \delta(w + \sqrt{\lambda_k}) \right] \psi_k(p; a).
 \end{aligned}$$

- (2) *Suppose that $g = 0$. Then the Fourier transform of the odd extension $z_{p,odd}(t; a)$ of $z_p(t; a)$ is given by*

$$\begin{aligned}
 (3.4) \quad & \mathcal{F}(z_{p,odd}(t; a)) \\
 &= \frac{1}{\sqrt{2\pi}} \int_{-\infty}^{\infty} z_{p,odd}(t; a) e^{-iwt} dt \\
 &= \sum_{k=1}^{\infty} \sqrt{\frac{\pi}{2}} \left[-iB_k \delta(w - \sqrt{\lambda_k}) + iB_k \delta(w + \sqrt{\lambda_k}) \right] \psi_k(p; a).
 \end{aligned}$$

- (3) *Suppose that we have either a) $g(x) > 0$ on $(0, 1)$ and $h = 0$, or b) $h(x) > 0$ a.e on $(0, 1)$ and $g = 0$. Then $z_p(t; a) = z_p(t; b)$ for $t > 0$ implies $\lambda_1(b) = \lambda_1(a)$. If $g(x) > 0$ on $(0, 1)$, then $A_1^a = A_1^b$. If $h(x) > 0$ a.e on $(0, 1)$, then $B_1^a = B_1^b$.*

Proof. (1) By Theorem 2.3 the observations $z_p(t)$ are continuous and bounded functions. Here the dependency on a is dropped for convenience. Since $h = 0$, the even extension of $z_p(t)$

$$(3.5) \quad z_{p,even}(t) = \sum_{k=1}^{\infty} A_k \cos(\sqrt{\lambda_k}t) \psi_k(p), \quad t \in \mathbb{R}$$

is continuous and bounded on \mathbb{R} .

Thus $z_{p,even}(t) \in \mathcal{S}'(\mathbb{R})$, i.e., it defines a tempered distribution on \mathbb{R} . See [19] for relevant definitions. Since the series in (3.1) converges uniformly on \mathbb{R}^+ , the series in (3.5) converges in $\mathcal{S}'(\mathbb{R})$. By [19], Chapter VI, Section 2, the Fourier transform \mathcal{F} and its inverse are bijective continuous linear mappings of $\mathcal{S}'(\mathbb{R})$ onto itself. Therefore the Fourier transform of $z_{p,even}(t)$ can be found by the termwise application of \mathcal{F} to the series (3.5). Since $\mathcal{F}(e^{i\lambda t}) = \sqrt{2\pi}\delta(w - \lambda)$ the representation (3.3) follows.

(2) If $g = 0$, then the odd extension of $z_p(t)$ is given by

$$(3.6) \quad z_{p,odd}(t) = \sum_{k=1}^{\infty} B_k \sin(\sqrt{\lambda_k t}) \psi_k(p), \quad t \in \mathbb{R}.$$

Arguing as in part (1), one gets the representation (3.4).

(3) Suppose that $g(x) > 0$ on $(0, 1)$ and $h = 0$. By Theorem 2.2(2), $\psi_1(x; a) > 0$ on $(0, 1)$ for any $a \in \mathcal{PC}$. Thus $A_1^a > 0$. Therefore the first nonzero term in the expansion (3.3) is $A_1^a \delta(w - \sqrt{\lambda_1(a)}) + A_1^a \delta(w + \sqrt{\lambda_1(a)})$. Since $z_p(t; a) = z_p(t; b)$ the same argument shows that the first nonzero term in (3.3) must be $A_1^b \delta(w - \sqrt{\lambda_1(b)}) + A_1^b \delta(w + \sqrt{\lambda_1(b)})$. Thus $\lambda_1(a) = \lambda_1(b)$ and $A_1^a = A_1^b$.

If $h(x) > 0$ a.e on $(0, 1)$ and $g = 0$, then one uses the Fourier transform expansion (3.4) for the odd extension of $z_p(t)$ to conclude that $\lambda_1(a) = \lambda_1(b)$ and $B_1^a = B_1^b$. \square

Remark 3.2. Suppose that $g(x) > 0$ on $(0, 1)$ and $h = 0$. By Theorem 3.1 one can construct the M -tuple

$$\mathcal{G}(a) = (\lambda_1(a), G_1(a), \dots, G_{M-1}(a)) \in \mathbb{R}^M$$

defined in (3.2) from the data $z_m(t)$, $m = 1, 2, \dots, M - 1$ as follows. First, find the Fourier transform of $z_{m,even}(t)$. It has the form (3.3). The smallest positive value of w where $\mathcal{F}(z_{m,even}(t)) \neq 0$ gives the square root of the eigenvalue $\lambda_1(a)$. With $\lambda_1(a)$ being determined, find $G_m(a) = cA_1\psi_1(p_m; a)$, where c is some constant. Since $A_1(a) = \langle g(x), \psi_1(x; a) \rangle$, the factor $C(a) = cA_1$ is independent of the observation point p_m . In case $h(x) > 0$ a.e. on $(0, 1)$ and $g = 0$ apply the Fourier transform to the odd extension of $z_m(t)$. Use (3.4) to find the eigenvalue $\lambda_1(a)$ and $G_m(a) = cB_1\psi_1(p_m; a)$. A numerical algorithm for the construction of $\mathcal{G}(a)$ is discussed in Section 5.

Theorem 3.3. *Given $\sigma > 0$ let an integer M be such that*

$$M \geq \frac{3}{\sigma} \quad \text{and} \quad M > 2\sqrt{\frac{\mu}{\nu}}.$$

Suppose that the observations $z_m(t; a) = u(p_m, t; a)$, $t > 0$ of (1.2) for $p_m = m/M$, $m = 1, 2, \dots, M - 1$ are given. It is assumed that the initial conditions g and h in (1.2) satisfy $h \in H$, $g \in V$ and either $g(x) > 0$ and $h = 0$ on $(0, 1)$, or $h(x) > 0$ a.e. on $(0, 1)$ and $g = 0$. Then the piecewise constant parameter $a \in A_{ad}$ is identifiable in the class of piecewise constant functions $\mathcal{PC}(\sigma)$.

Proof. By Remark 3.2 one can use the observations $z_m(t)$, $m = 1, 2, \dots, M-1$ to construct the M -tuple $\mathcal{G}(a)$. It is shown in [7] that the knowledge of $\mathcal{G}(a)$ is sufficient for the unique identification of the piecewise constant parameter $a \in \mathcal{PC}(\sigma)$. The identification of a from $\mathcal{G}(a)$ is accomplished by the Marching Algorithm (see [7], Theorem 4.6). \square

The continuity properties of the mapping $\mathcal{G}(a)$ are summarized in the next theorem.

Theorem 3.4. *Let $\mathcal{PC} \subset A_{\text{ad}}$ be equipped with the $L^1[0, 1]$ topology. Assume that the conditions of Theorem 3.3 are satisfied. Then*

- (1) *The mapping $a \rightarrow \mathcal{G}(a)$ from \mathcal{PC} into \mathbb{R}^M is continuous.*
- (2) *The identification mapping \mathcal{G}^{-1} from $\mathcal{G}(\mathcal{PC}(\sigma)) \subset \mathbb{R}^M$ into $\mathcal{PC}(\sigma)$ is continuous.*

Proof. By Theorem 3.1 in [8] the mapping $a \rightarrow \lambda_1(a)$ is continuous. By Theorem 3.3 in [8] the mapping $a \rightarrow \psi_1(x; a)$ is continuous from \mathcal{PC} into $C[0, 1]$. Since $C(a) = \langle g(x), \psi_1(x; a) \rangle$ and $G_m(a) = C(a)\psi_1(p_m; a)$ the first assertion follows.

Concerning (2), the set $\mathcal{PC}(\sigma)$ is compact in A_{ad} according to Theorem 4.2 in [8]. Thus \mathcal{G} is continuous on a compact set. By Theorem 3.3 it is invertible. Therefore \mathcal{G}^{-1} is continuous on $\mathcal{G}(\mathcal{PC}(\sigma))$. \square

In the next two theorems we establish the continuity of the identification with respect to the observations $z_m(t)$.

Theorem 3.5. *Let $a \in \mathcal{PC} \subset A_{\text{ad}}$ equipped with the $L^1[0, 1]$ topology, and $u(a)$ be the solution of the string vibration problem (1.2) with the initial conditions $g \in V$, $h \in H$. Then the mapping $a \rightarrow u(a)$ from \mathcal{PC} into $C([0, 1] \times [0, T])$ is continuous for any $T > 0$.*

Proof. According to Theorem 2.3 the partial sums $u^N(x, t; a)$ converge to the solution $u(x, y; a)$ as $N \rightarrow \infty$ uniformly on $D = [0, 1] \times [0, \infty)$. Moreover, the rate of convergence does not depend on $a \in \mathcal{PC}$. Thus, it is enough to show that the map $a \rightarrow u^N(t; a)$ is continuous from \mathcal{PC} into $C([0, 1] \times [0, T])$. The continuity of the eigenvalues $\lambda_k(a)$ and the eigenfunctions $\psi_k(x, a)$ as the functions of a was established in Theorems 3.1 and 3.3 in [8]. Therefore each term in the expansion

$$(3.7) \quad u^N(x, t; a) = \sum_{k=1}^N (A_k \cos(\sqrt{\lambda_k(a)}t) + B_k \sin(\sqrt{\lambda_k(a)}t))\psi_k(x; a)$$

with

$$A_k = \langle g, \psi_k(a) \rangle, \quad B_k = \frac{\langle h, \psi_k(a) \rangle}{\sqrt{\lambda_k(a)}}, \quad k = 1, 2, \dots$$

is continuous with respect to a and the result follows. \square

Theorem 3.6. *Let $a \in \mathcal{PC}(\sigma) \subset A_{\text{ad}}$ equipped with the $L^1[0, 1]$ topology, and the initial conditions of the string vibration problem (1.1) satisfy $h \in H$, $g \in V$ and either $g(x) > 0$ and $h = 0$, or $g = 0$ and $h > 0$ a.e. on $(0, 1)$. Suppose that $\{a_n\}_{n=1}^\infty \subset \mathcal{PC}(\sigma)$ be such that $z_m(t; a_n) \rightarrow z_m(t; a)$ in $C[0, T]$ as $n \rightarrow \infty$ for any $T > 0$ and all $m = 1, 2, \dots, M - 1$. Then $a_n \rightarrow a$ in $\mathcal{PC}(\sigma)$ as $n \rightarrow \infty$.*

Proof. By Theorem 4.2 in [8], the set $\mathcal{PC}(\sigma)$ is compact. Therefore, we can assume without loss of generality (passing to a subsequence) that $a_n \rightarrow b \in \mathcal{PC}(\sigma)$ as $n \rightarrow \infty$. By Theorem 3.5, $z_m(t; a_n) \rightarrow z_m(t; b)$ in $C[0, T]$ as $n \rightarrow \infty$ for any $T > 0$. Thus $z_m(t; a) = z_m(t; b)$ for $t > 0$ and all $m = 1, 2, \dots, M - 1$. By Theorem 3.1 one gets $\mathcal{G}(a) = \mathcal{G}(b)$. Then Theorem 3.3 implies $a = b$. \square

4. Extension to boundary and external inputs

Our goal here is to extend the results of the previous section to system (1.1). First, we derive a formula for the solution $u(x, t; a)$ of (1.1). Then we show how to use observations $z_p(t; a) = u(p, t; a)$ to construct the M -tuple $\mathcal{G}(a)$. The parameter a is reconstructed from $\mathcal{G}(a)$ using the Marching Algorithm.

Theorem 4.1. *Suppose that $T > 0$, $a \in \mathcal{PC}$, $g \in V$, $h \in H$ and $f \in C([0, T]; H)$. Furthermore, assume that $q_1(t), q_2(t) \in C^2[0, T]$ and $q_1(0) = q'_1(0) = q_2(0) = q'_2(0)$.*

(1) *Let $\{\lambda_k, \psi_k\}_{k=1}^\infty$ be the eigenvalues and the eigenfunctions of (2.1). Let*

$$(4.1) \quad \Phi(x, t; a) = \frac{q_2(t) - q_1(t)}{\int_0^1 \frac{1}{a(s)} ds} \int_0^x \frac{1}{a(s)} ds + q_1(t).$$

Let $g_k = \langle g, \psi_k \rangle$, $h_k = \langle h, \psi_k \rangle$, $\Phi_k(t) = \langle \Phi(\cdot, t), \psi_k \rangle$ and $f_k(t) = \langle f(\cdot, t), \psi_k \rangle$ for $k = 1, 2, \dots$. Then the solution $u(x, t; a)$ of (1.1) is given by

$$(4.2) \quad u(x, t; a) = \Phi(x, t; a) + \sum_{k=1}^\infty \beta_k(t; a) \psi_k(x),$$

where

$$(4.3) \quad \begin{aligned} \beta_k(t; a) &= g_k \cos \sqrt{\lambda_k} t + \frac{1}{\sqrt{\lambda_k}} h_k \sin \sqrt{\lambda_k} t \\ &+ \frac{1}{\sqrt{\lambda_k}} \int_0^t [f_k(\tau) - \Phi_k''(\tau)] \sin \sqrt{\lambda_k}(t - \tau) d\tau \end{aligned}$$

for $k = 1, 2, \dots$

(2) *For each $t > 0$ and $a \in \mathcal{PC}$ the series in (4.2) converges in V . Moreover, this convergence is uniform with respect to t on $0 \leq t \leq T$ and $a \in \mathcal{PC}$. The solution $u(x, t; a)$ is continuous on $[0, 1] \times [0, T]$.*

Proof. Since $(a\Phi_x)_x = 0$ we define the weak solution of (1.1) to be $u(x, t; a) = \Phi(x, t; a) + v(x, t; a)$, where v is the weak solution of

$$(4.4) \quad \begin{cases} v_{tt} - (av_x)_x = -\Phi_{tt} + f, & 0 < x < 1, 0 < t < T, \\ v(0, t) = 0, & 0 < t < T, \\ v(1, t) = 0, & 0 < t < T, \\ v(x, 0) = g(x), & 0 < x < 1, \\ v_t(x, 0) = h(x), & 0 < x < 1. \end{cases}$$

Let $\beta_k(t) = \langle v(\cdot, t), \psi_k \rangle$. To simplify the notation the dependency of β_k on a is suppressed. Then

$$\beta_k''(t) + \lambda_k \beta_k(t) = -\Phi_k''(t) + f_k(t), \quad \beta_k(0) = g_k, \quad \beta_k'(0) = h_k.$$

Thus $\beta_k(t)$ has the representation stated in (4.3).

Note that (4.3) implies that

$$\lambda_k \beta_k^2(t) \leq 2\lambda_k g_k^2 + 2h_k^2 + 2 \left[\int_0^t (f_k(\tau) - \Phi_k''(\tau)) \sin \sqrt{\lambda_k}(t - \tau) d\tau \right]^2,$$

$$\sum_{k=1}^\infty \lambda_k g_k^2 = \|g\|_{V_a}^2, \quad \sum_{k=1}^\infty h_k^2 = \|h\|_H^2$$

since $g \in V, h \in H$, and $q_1, q_2 \in C^2[0, T]$.

Also

$$\begin{aligned} & \sum_{k=1}^\infty \left[\int_0^t (f_k(\tau) - \Phi_k''(\tau)) \sin \sqrt{\lambda_k}(t - \tau) d\tau \right]^2 \\ & \leq T \sum_{k=1}^\infty \int_0^T |f_k(\tau) - \Phi_k''(\tau)|^2 d\tau \\ & = T \int_0^T \|f - \Phi_{tt}\|_H^2 d\tau \end{aligned}$$

since $f \in C([0, T]; H)$.

By Theorem 6.5.2 in [5], eigenfunctions $\{\psi_k(x)/\sqrt{\lambda_k}\}_{k=1}^\infty$ form an orthonormal basis in the energy space $V_a = \{w \in V : \|w\|_{V_a}^2 = \int_0^1 a(x)|w'(x)|^2 dx < \infty\}$. Since $\sum_{k=1}^\infty \lambda_k \beta_k^2(t) < \infty$, the series $\sum_{k=1}^\infty \beta_k(t)\psi_k$ converges in V_a .

Note that $\sqrt{\nu}\|w\|_V \leq \|w\|_{V_a} \leq \sqrt{\mu}\|w\|_V$. Therefore the series $\sum_{k=1}^\infty \beta_k(t)\psi_k$ converges in V for each $t \geq 0$ uniformly with respect to $a \in \mathcal{PC}$. Since the estimates for $\beta_k(t; a)$ are independent of t , the convergence is uniform with respect to t on $0 \leq t \leq T$. Each coefficient $\beta_k(t)$ is continuous on $[0, T]$. Therefore $v \in C([0, T]; V)$ is the unique weak solution of (4.4). The continuity of v on $[0, 1] \times [0, T]$ follows from the continuous imbedding of $C[0, 1]$ into V . \square

Next theorem describes some conditions under which the identifiability for (1.1) is possible.

Theorem 4.2. *Given $\sigma > 0$ let an integer M be such that*

$$M \geq \frac{3}{\sigma} \quad \text{and} \quad M > 2\sqrt{\frac{\mu}{\nu}}.$$

Suppose that the observations $z_m(t; a) = u(p_m, t; a)$ for $p_m = m/M$, $m = 1, 2, \dots, M - 1$ and $t > 0$ of the string vibration system (1.1) are given. Then the parameter $a \in A_{\text{ad}}$ is identifiable in the class of piecewise constant functions $\mathcal{PC}(\sigma)$ in each one of the following four cases.

- (1) $f = 0$, $q_1 = 0$, $q_2 = 0$, $g \in V$, $h \in H$ and either $g(x) > 0$ and $h = 0$ on $(0, 1)$, or $g = 0$ and $h(x) > 0$ a.e. on $(0, 1)$.
- (2) $g = 0$, $h = 0$, $q_1 = 0$, $q_2 = 0$, $f(x, t) = s(x)r(t)$, $s \in H$, $s > 0$ a.e. on $(0, 1)$, $r \in C[0, \infty)$.
- (3) $g = 0$, $h = 0$, $f = 0$, $q_2 = 0$, $q_1 \neq 0$, $q_1 \in C^2[0, \infty)$, $q_1(0) = q_1'(0) = 0$.
- (4) $g = 0$, $h = 0$, $f = 0$, $q_1 = 0$, $q_2 \neq 0$, $q_2 \in C^2[0, \infty)$, $q_2(0) = q_2'(0) = 0$.

Proof. In every case we show how to extract the M -tuple

$$\mathcal{G}(a) = (\lambda_1(a), G_1(a), \dots, G_{M-1}(a)) \in \mathbb{R}^M, \quad a \in \mathcal{PC}(\sigma),$$

where $G_m(a) = C(a)\psi_1(p_m; a)$, from the observations $z_m(t; a)$. Then, as in [7], Theorem 4.6, the Marching Algorithm uniquely recovers the sought coefficient $a(x)$.

- (1) This is Theorem 3.3.
- (2) Let

$$y_m(t) = \sum_{k=1}^{\infty} \frac{1}{\sqrt{\lambda_k}} \langle s, \psi_k \rangle \sin \sqrt{\lambda_k} t \psi_k(p_m).$$

Arguing as in Theorem 2.3 we conclude that this series converges uniformly on $[0, \infty)$, and $y_m(t) \in C[0, \infty)$. This fact and Theorem 4.1 imply that the observations $z_m(t) = u(p_m, t; a)$ are given by

$$\begin{aligned} z_m(t; a) &= \sum_{k=1}^{\infty} \int_0^t \frac{1}{\sqrt{\lambda_k}} \langle s, \psi_k \rangle \psi_k(p_m) \sin \sqrt{\lambda_k}(t - \tau) r(\tau) d\tau \\ &= \int_0^t \left[\sum_{k=1}^{\infty} \frac{1}{\sqrt{\lambda_k}} \langle s, \psi_k \rangle \sin \sqrt{\lambda_k}(t - \tau) \psi_k(p_m) \right] r(\tau) d\tau. \end{aligned}$$

Since $r(t) \in C[0, \infty)$ by the assumption, Titchmarsh Theorem ([18, Theorem 152, Chap. XI, p. 325] or [19, Section 6.5]) implies that the Volterra integral equation

$$z_m(t; a) = \int_0^t y_m(t - \tau) r(\tau) d\tau$$

is uniquely solvable for $y_m(t)$.

Since $s > 0$ is assumed to be in H one has $\langle s(x), \psi_1(x; a) \rangle \neq 0$. Now one can proceed as in Theorem 3.1, i.e., use the Fourier transform of the odd extension of $y_m(t)$ to find the first eigenvalue $\lambda_1(a)$ and a

constant multiple $C(a)\psi_1(p_m; a)$ of the first eigenfunction at p_m . In this case $C(a) = c\langle s, \psi_1 \rangle / \sqrt{\lambda_1}$. Thus $\mathcal{G}(a)$ is obtained.

- (3) In this case function $\Phi(x, t; a)$ defined in (4.1) has the form $\Phi(x, t; a) = \alpha(x)q_1(t)$ with $\alpha(0) = 1$, $\alpha(1) = 0$. Note that $\alpha \in H$, but $\alpha \notin V$. Theorem 4.1 gives

$$\begin{aligned} z_m(t; a) &= u(p_m, t; a) \\ &= \Phi(p_m, t; a) - \sum_{k=1}^{\infty} \int_0^t \frac{1}{\sqrt{\lambda_k}} \langle \alpha, \psi_k \rangle \psi_k(p_m) \sin \sqrt{\lambda_k}(t - \tau) q_1''(\tau) d\tau. \end{aligned}$$

Let

$$(4.5) \quad y_m(t) = \sum_{k=1}^{\infty} \frac{1}{\sqrt{\lambda_k}} \langle \alpha, \psi_k \rangle \sin \sqrt{\lambda_k} t \psi_k(p_m).$$

Arguing the uniform convergence of this series as in case (2), we obtain that $y_m(t)$ is a continuous bounded function on $([0, \infty)$ and

$$(4.6) \quad z_m(t; a) = \alpha(p_m)q_1(t) - \int_0^t y_m(t - \tau) q_1''(\tau) d\tau.$$

Applying the Laplace transform,

$$Z_m(s) = \alpha(p_m)Q(s) - Y_m(s)s^2Q(s),$$

where the capitalized functions are the Laplace transforms of the corresponding lower case functions in (4.6). Since $y_m(t)$ is bounded on $[0, \infty)$, its Laplace transform $Y_m(s)$ is analytic in $D = \{z \in \mathbb{C} : \operatorname{Re} z > 0\}$. Therefore

$$\frac{\alpha(p_m)}{s^2} - Y_m(s) = \frac{Z_m(s)}{s^2Q(s)}$$

in D . The inverse Laplace transform gives

$$\alpha(p_m)t - y_m(t) = L^{-1} \left\{ \frac{Z_m(s)}{s^2Q(s)} \right\}.$$

This means that given the data $z_m(t; a)$ and the boundary input $q_1(t)$ one can uniquely determine the continuous function $y_m^{(\alpha)}(t) = \alpha(p_m)t - y_m(t)$ for $t \geq 0$. Extend it to $t < 0$ so that it will be odd. Because of the representation (4.5), the Fourier transform of this odd extension $y_{m, \text{odd}}^{(\alpha)}(t)$ is given by

$$(4.7) \quad \begin{aligned} \mathcal{F}(y_{m, \text{odd}}^{(\alpha)}(t)) &= i\sqrt{2\pi} \delta'(w) \alpha(p_m) \\ &+ i\frac{\pi}{2} \sum_{k=1}^{\infty} \frac{\langle a, \psi_k \rangle}{\sqrt{\lambda_k}} [\delta(w - \sqrt{\lambda_k}) - \delta(w + \sqrt{\lambda_k})] \psi_k(p_m). \end{aligned}$$

Recall that $\Phi(x, t; a) = \alpha(x)q_1(t)$. Therefore $\alpha(x) > 0$ on $(0, 1)$, and the coefficient $\langle \alpha, \psi_1 \rangle$ is positive. Thus the first eigenvalue $\lambda_1(a)$ and a constant multiple of the first eigenfunction $C(a)\psi_1(p_m)$ can be uniquely determined from (4.7). That is the M -tuple $\mathcal{G}(a)$ is uniquely determined from the data.

(4) The proof is the same as in case (3). □

In each case of Theorem 4.2 one can define the solution map and its inverse. The continuity results (with respect to a) for these maps are analogous to the results stated in Theorems 3.5 and 3.6. Their proofs follow the same lines as in Section 3.

5. Numerical results

In a typical numerical experiment we used $\nu = 0.1, \mu = 1.0$ for the bounds of the admissible set, $M = 12$ for the number of observation points p_m , and $T = 80$ for the observation time interval $[0, T]$. The governing system was (1.2) with $h(x) = 0$ and $g(x) = x(1 - x)$. The number of discontinuity points in the original piecewise constant parameter \hat{a} varied from 2 to 3.

In one such experiment the parameter \hat{a} was chosen to be

$$(5.1) \quad \hat{a} = \begin{cases} 0.1, & 0 \leq x < 0.21 \\ 1.0, & 0.21 \leq x < 0.60 \\ 0.1, & 0.60 \leq x \leq 1.0. \end{cases}$$

Its first eigenvalue is $\lambda_1(\hat{a}) = 1.1528$. For details of the numerical computation of eigenvalues and eigenfunctions for (2.1) see [8]. Given such a parameter \hat{a} , the observation data $\hat{z}_m(t)$ was computed according to (3.1) with $B_k = 0$ for all k . Then the data was contaminated by noise of level η according to

$$z_m(t) = \hat{z}_m(t) + 2\eta(r(\zeta) - 0.5) \max_{0 \leq x \leq 1} |g(x)|,$$

where $r(\zeta)$ is a random variable uniformly distributed on interval $[0, 1)$, and $g(x)$ is the initial position. The identification of the coefficient \hat{a} was conducted using $J = 1025$ values $z_m(t_j)$ at time instants $t_j, j = 0, 1, \dots, J$ equidistant on the interval $[0, T]$.

According to the algorithm developed in the previous sections, we proceed in two steps. First, we determine the first eigenvalue λ_1 and a constant multiple of the first eigenfunction, i.e., the M -tuple $\mathcal{G}(a)$. In the second step the coefficient a is recovered from $\mathcal{G}(a)$.

Identification of $\mathcal{G}(a)$

Fix an observation point $p_m = m/M, m = 1, 2, \dots, M - 1$. To find λ_1 from the data collected at this point we apply the Fourier cosine transform to the

observations $z_m(t_j)$, $j = 0, 1, \dots, J$. Let

$$F_C(n, m) = \frac{2}{T} \int_0^T z_m(t) \cos\left(\frac{\pi n t}{T}\right) dt,$$

and let \tilde{n} correspond to the first maximum in the sequence $|F_C(n, m)|$, $n = 1, 2, \dots$. The eigenvalue λ_1 is assigned to be $(\pi\tilde{n}/T)^2$. The value of $G_m(a)$ in $\mathcal{G}(a)$ is chosen to be $F_C(\tilde{n}, m)$. The process is repeated for each $m = 1, \dots, M-1$. It implements numerically finding $\mathcal{G}(a)$ from the data $z_{p,even}(t)$ in (3.5).

Note, that the above algorithm assumes the choice of T such that $\lambda_1 \geq \left(\frac{\pi}{T}\right)^2$. According to Theorem 2.2(3) it is sufficient to take $T \geq \frac{1}{\sqrt{\nu}}$. However, the numerical identification of λ_1 is better for much larger T .

Since each observation point $p_m \in (0, 1)$ identifies its own value for λ_1 one has to make a choice as to how to use these values. In [8] the average of the first eigenvalues over the middle third of the observation points was used. Here we fixed the observation point p_m to be near the middle of the observation interval $(0, 1)$ by choosing $m = M/2$, since the influence of noise at such a point would be less pronounced. Then the relative errors

$$E_\lambda = \frac{|\lambda_1 - \lambda_1(\hat{a})|}{\lambda_1(\hat{a})}$$

were computed in several independent runs of the program. It turns out that the first eigenvalue identified by this method is practically insensitive to low noise levels. In our experiments we obtained the same value $\lambda_1 = 1.1242$ for all the noise levels η in the data $z_m(t)$ as indicated in Table 1.

The values of G_m in $\mathcal{G}(a)$ are supposed to be a constant multiple of the first eigenfunction $\psi_1(p_m; \hat{a})$ at the observation points p_m . The quality of this identification can be judged by the deviation of the ratio $G_m/\psi_1(p_m; \hat{a})$ from its average over the observation points p_m . Ideally, this deviation should be equal to zero, since the ratio is expected to be a constant. To quantify this deviation let

$$G_{av} = \frac{1}{M-1} \sum_{m=1}^{M-1} \frac{G_m}{\psi_1(p_m; \hat{a})}$$

and

$$E_G = \frac{1}{G_{av}} \max_m \left| \frac{G_m}{\psi_1(p_m; \hat{a})} - G_{av} \right|.$$

The third column in Table 1 shows the values of the relative deviation E_G for various noise levels η .

Having $\mathcal{G}(a)$ determined in the first step of the algorithm, the second step consists of finding the parameter \bar{a} . While this goal may be attempted to be accomplished by the Marching Algorithm applied to $\mathcal{G}(a)$, the numerical evidence shows that such an approach is unsatisfactory. The numerical performance of the Marching Algorithm is excellent in accordance with the theoretical justification only for the data $\mathcal{G}(a)$ very close to $\mathcal{G}(\hat{a})$. This would mean the relative

errors E_λ and E_G being several orders of magnitude smaller, than the ones shown in Table 1. In practice such a close identification cannot be achieved, and a different method such as the one described below is needed.

Identification of piecewise constant parameter \bar{a}

The data is the M -tuple $\mathcal{G}(a) = \{\lambda_1, G_1, \dots, G_{M-1}\}$.

- (1) Fix $N > 0$. Form the objective function $\Pi(a)$ by

$$(5.2) \quad \Pi(a) = \min_{c \in \mathbb{R}} \sum_{m=1}^M (cG_m - \psi_1(p_m; a))^2$$

for the parameters $a \in A_N \subset A_{\text{ad}}$ having at most $N - 1$ discontinuity points on the interval $[0, 1]$.

- (2) Use Powell's minimization method in $K = 2N - 1$ variables ($N - 1$ discontinuity points and N parameter a values) to find

$$\Pi(\bar{a}) = \min_{a \in A_N} \Pi(a).$$

The minimizer \bar{a} is the sought piecewise constant parameter.

We employ the *Powell's minimization method* since it does not require gradient computations, but still has a quadratic convergence near the points of minima. The modification used here is from [9]. We used Brent's one-dimensional minimization method BA, see [17, 9], in all the minimization steps below.

Powell's minimization method

The method iteratively minimizes a function $\Pi(q)$, $q \in \mathbb{R}^K$ of K variables. Steps 1-7 describe one iteration of the method.

- (1) Initialize the set of directions $u_i \in \mathbb{R}^K$ to the standard basis vectors in \mathbb{R}^K

$$u_i = e_i, \quad i = 1, \dots, K.$$

- (2) Save your starting position as $q_0 \in \mathbb{R}^K$.
- (3) For $i = 1, \dots, K$ move from q_0 along the direction u_i and find the point of minimum p_i .
- (4) Re-index the directions u_i , so that (for the new indices) $\Pi(p_1) \leq \Pi(p_2) \leq \dots \leq \Pi(p_K) \leq \Pi(q_0)$.
- (5) Move from q_0 along the new direction u_1 and find the point of minimum r_1 . Move from r_1 along the direction u_2 and find the point of minimum r_2 , etc. Move from r_{K-1} along the direction u_K and find the point of minimum r_K .
- (6) Set $v = r_K - q_0$.
- (7) Move from q_0 along the direction v and find the minimum. Call it q_0 . It replaces q_0 from step 2.
- (8) Repeat the above steps until a stopping criterion is satisfied.

TABLE 1. Relative errors E_λ , E_G and E_a in the identification of the piecewise constant parameter \hat{a} for various noise levels η .

η	E_λ	E_G	E_a
0.00	0.02478	0.0018	0.2299
0.02	0.02478	0.0043	0.2169
0.05	0.02478	0.0051	0.2360
0.10	0.02478	0.0077	0.2157
0.20	0.02478	0.0181	0.2067

If the minimization is to be restricted to a subset $A_{\text{ad}} \subset \mathbb{R}^K$, then the moves in the one-dimensional minimization steps above are restricted so that the trial points would not leave A_{ad} .

Thus, for $N = 4$, the minimization of $\Pi(a)$, $a \in A_N$ defined in (5.2) is in $K = 2N - 1 = 7$ variables representing 3 possible discontinuity points and 4 values of the parameter a on the intervals where it is constant.

The quality of the identification was measured as the relative L^1 error E_a between the original $\hat{a}(x)$ and the identified parameter $\bar{a}(x)$:

$$(5.3) \quad E_a = \frac{\int_0^1 |\bar{a}(x) - \hat{a}(x)| dx}{\int_0^1 \hat{a}(x) dx}.$$

The relative errors E_a in the identification of the piecewise constant parameter \hat{a} for various noise levels η are shown in the fourth column of Table 1. Various parameters used in the stopping criteria in the iterative processes in the above algorithms were determined experimentally.

For example, for noise level $\eta = 0.20$, the identified parameter \bar{a} was

$$(5.4) \quad \bar{a} = \begin{cases} 0.1000, & 0 \leq x < 0.2081 \\ 0.9943, & 0.2081 \leq x < 0.4992 \\ 0.1112, & 0.4992 \leq x < 0.6325 \\ 0.1000, & 0.6325 \leq x \leq 1.0, \end{cases}$$

giving the relative identification error $E_a = 0.2067$.

Figure 1 shows both the original parameter \hat{a} defined in (5.1) (dashed line) and the identified parameter \bar{a} (solid line).

The modification of the objective function $\Pi(a)$ in (5.2) to

$$\Pi(a) = \min_{c \in \mathbb{R}} \sum_{m=1}^M (cG_m - \psi_1(p_m; a))^2 + \beta(\lambda_1 - \lambda_1(a))^2,$$

where $\beta > 0$, did not produce an improvement in the identification.

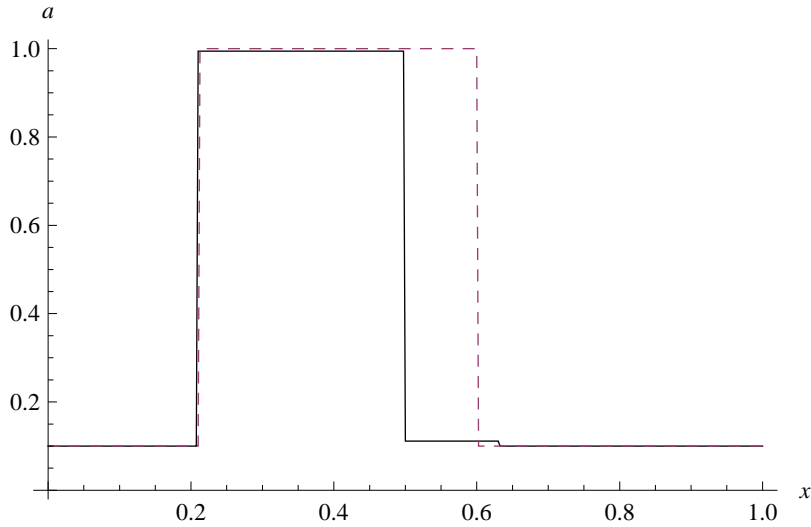


FIGURE 1. Original \hat{a} (dashed line) and identified piecewise constant parameter \bar{a} (solid line) for $\eta = 0.20$.

6. Conclusions

While in most parameter estimation problems one can hope only to achieve a best fit to data solution, sometimes it can be shown that such an identification is unique. In such case it is said that the sought parameter is identifiable within a certain class. In our recent work [9] we have shown that piecewise constant conductivities $a \in \mathcal{PC}(\sigma)$ are identifiable from observation $z_m(t; a)$ of the heat conduction process taken at finitely many points p_m . Conditions for the conductivity identifiability for nonzero boundary and external inputs are specified in [8].

In this paper we show that the piecewise constant parameter $a(x)$ associated with the variable density of a composite string can also be identified from finitely many observations $z_m(t; a)$, $m = 1, \dots, M - 1$. The extension of this result to boundary and external inputs is described in Theorem 4.2.

The identification is achieved in two steps. First, one constructs the M -tuple $\mathcal{G}(a)$ defined in (3.2). Here a is the piecewise constant parameter we are seeking to identify. The finite sequence $\mathcal{G}(a)$ consists of the first eigenvalue and a constant multiple of the first eigenfunction at the observation points p_m . In the second step the data $\mathcal{G}(a)$ is used to identify $a(x)$ by the Marching Algorithm, see [9].

It is shown in Theorems 3.4-3.6 that the Marching Algorithm not only provides the unique identification of the conductivity a , but that the identification is also continuous (stable). This result is based on the continuity of eigenvalues,

eigenfunctions, and the solutions with respect to the $L^1[0, 1]$ topology in the set of admissible parameters A_{ad} , see Section 3.

Numerical experiments based on the Marching Algorithm algorithm show the perfect identification for noiseless data, if the M -tuple $\mathcal{G}(a)$ is identified with no errors. As expected, the identification deteriorates significantly even for small levels of noise in the data. Of course, in practically interesting situations noise levels are large, so the Marching Algorithm by itself would not be useful in such cases.

The identification algorithm for noisy data is presented in Section 5. Its main novel point is, in agreement with the theoretical developments, the separation of the identification process into two separate steps. In step one the first eigenvalue and a multiple of the first eigenfunction (the M -tuple $\mathcal{G}(a)$) are extracted from the observations. In the second step a general minimization method is used to find the piecewise constant parameter $a(x)$ from $\mathcal{G}(a)$. The first eigenvalue and the eigenfunction are found using the Fourier transform of the observation $z_m(t)$. The second step is accomplished by Powell's minimization algorithm.

Numerical results in Section 5 show that this algorithm achieves very good results in the reconstruction of $\mathcal{G}(a)$. Even for high level of noise the identification of the first eigenvalue is very stable. The second step of the identification achieves satisfactory results within the 25% relative error range, but its performance is somewhat inferior to the precision achieved in step 1 of the method.

References

- [1] S. Aihara, *Maximum likelihood estimate for discontinuous parameter in stochastic hyperbolic systems*, Acta Appl. Math. **35** (1994), no. 1-2, 131–151.
- [2] H. T. Banks and K. Kunisch, *Estimation techniques for distributed parameter systems*, Birkhauser Boston, Inc., Boston, MA, 1989.
- [3] M. Burger and W. Mühlhuber, *Iterative regularization of parameter identification problems by sequential quadratic programming methods*, Inverse Problems **18** (2002), no. 4, 943–969.
- [4] A. Elayyan and V. Isakov, *On uniqueness of recovery of the discontinuous conductivity coefficient of a parabolic equation*, SIAM J. Math. Anal. **28** (1997), no. 1, 49–59.
- [5] L. C. Evans, *Partial Differential Equations*, Graduate Studies in Mathematics, 19. American Mathematical Society, Providence, RI, 1998.
- [6] S. Gutman, *Identification of discontinuous parameters in flow equations*, SIAM J. Control Optim. **28** (1990), no. 5, 1049–1060.
- [7] S. Gutman and J. Ha, *Identifiability of piecewise constant conductivity in a heat conduction process*, SIAM J. Control Optim. **46** (2007), no. 2, 694–713.
- [8] ———, *Parameter identifiability for heat conduction with a boundary input*, Math. Comput. Simulation **79** (2009), no. 7, 2192–2210.
- [9] J. Ha and S. Gutman, *Parameter estimation problem for a damped sine-Gordon equation*, International Journal of Appl. Math. and Mech. **2** (2006), no. 1, 11–23.
- [10] V. Isakov, *Inverse Problems for Partial Differential Equations*, Second edition. Applied Mathematical Sciences, 127. Springer, New York, 2006.
- [11] S. I. Kabanikhin, A. D. Satybaev, and M. A. Shishlenin, *Direct Methods of Solving Multidimensional Inverse Hyperbolic Problems*, VSP, Utrecht, Boston, 2005.

- [12] R. Kohn and M. Vogelius, *Determining conductivity by boundary measurements. II. Interior results*, Comm. Pure Appl. Math. **38** (1985), no. 5, 643–667.
- [13] J. L. Lions, *Optimal Control of Systems Governed by Partial Differential Equations*, Springer-Verlag, New York-Berlin 1971.
- [14] S. Nakagiri, *Review of Japanese work of the last ten years on identifiability in distributed parameter systems*, Inverse Problems **9** (1993), no. 2, 143–191.
- [15] Y. Orlov and J. Bentman, *Adaptive distributed parameter systems identification with enforceable identifiability conditions and reduced-order spatial differentiation*, IEEE Trans. Automat. Control **45** (2000), no. 2, 203–216.
- [16] A. Pierce, *Unique identification of eigenvalues and coefficients in a parabolic problem*, SIAM J. Control Optim. **17** (1979), no. 4, 494–499.
- [17] W. H. Press, S. A. Teukolsky, W. T. Vetterling, and B. P. Flannery, *Numerical Recipes in FORTRAN (2nd Ed.)*, Cambridge University Press, Cambridge, 1992.
- [18] E. C. Titchmarsh, *Introduction to the theory of Fourier integrals*, Oxford University Press, 1962.
- [19] E. Yosida, *Functional Analysis*, 6th Ed., Springer-Verlag, 1980.

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