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# Systematic discretization of input-output maps of linear infinite dimensional <br> systems <br> (Completely revised version of Preprint 2006-06) 

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# SYSTEMATIC DISCRETIZATION OF INPUT-OUTPUT MAPS OF LINEAR INFINITE-DIMENSIONAL SYSTEMS 

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#### Abstract

Many model reduction techniques take a semi-discretization of a PDE model as starting point and aim then at an accurate approximation of its input/output map. In this contribution, we discuss the direct discretization of the i/o map of the infinite-dimensional system for a general class of linear time-invariant systems with distributed inputs and outputs.

First, the input and output signals are discretized in space and time, resulting in the matrix representation of an approximated i/o-map. In a generalized sense, the matrix contains the Markov parameters of a corresponding time-discrete multi-input-multi-output system. Second, the system dynamics is approximated in form of the underlying evolution operator, in order to calculate the matrix representation numerically. The discretization framework, corresponding error estimates, a SVD-based system reduction method and a numerical application in an optimization problem are presented, and illustrated for a heat control system.


Key words. input-output map, discretization, infinite-dimensional control system, time-discrete MIMO system, model reduction, optimization, feedback control

AMS subject classifications. 39C20, 35B37

## 1. Introduction.

1.1. Motivation. The control of complex physical systems is a big challenge in many engineering applications as well as in mathematical research. Frequently, these control systems are modeled by infinite-dimensional state space systems on the basis of (instationary and nonlinear) partial differential equations (PDEs). On the one hand, space-discretizations resolving most of the state information typically lead to very large semi-discrete systems, on the other hand, popular design techniques for real-time controllers like robust control require linear models of very moderate size [13, 39].

Numerous approaches to bridge this gap are proposed in the literature, see e.g. $[3,12,28]$ and the references therein. In some applications one is interested in loworder models capturing essential state dynamics. Then e.g. low-order modeling on the basis of physical insight $[27,36,40]$ and models reduced by means of mathematical methods like Proper Orthogonal Decomposition (POD) [14] can be very useful. In this paper we focus on the frequent situation that models merely describing accurately the system's input/output (i/o) map are sufficient to realize efficient controls. Empirical or simulation-based black-box system identification [7, 30], and mathematical model reduction techniques like balanced truncation [28], moment matching [26] and recent variants of POD [42, 48] are tools to extract appropriate models.

Empirically and physically motivated approaches usually lack analytical estimates for the accuracy of the i/o map, for some mathematical model reduction techniques like balanced truncation such estimates exist. Most mathematical methods take, however, the space-discretized PDE as starting point, some of them even a timeinvariant linearization. The preceding PDE space discretizations are often neglected by assuming that they are 'sufficiently accurate'. Thereby, the PDE discretizations rarely take the efficient approximation of the i/o map into account. On the one hand,

[^0]the state space discretization typically aims at a reduction of the global state space error and is thus still oriented at a state simulation problem. On the other hand, the i/o error due to the discretization of spatially distributed inputs and outputs is rarely considered rigorously. Aiming at a low-dimensional model of the i/o behavior in the end, starting with a space-discretization of the original state space model can be considered as a conceptual detour.
1.2. An integral approach to derive i/o models with error estimates. In this paper we investigate a new and integral approach to derive low-order models with error estimates for the i/o behavior. We focus directly on the i/o map of the original infinite-dimensional system, in the following denoted by
$$
\mathbb{G}: \mathcal{U} \rightarrow \mathcal{Y}, \quad u=u(t, \theta) \mapsto y=y(t, \xi)
$$
and suggest the following framework for its direct discretization for a general class of linear time-invariant systems (introduced in Section 2). Here $u$ and $y$ are input and output signals from Hilbert spaces $\mathcal{U}$ and $\mathcal{Y}$, respectively, which may vary in time $t$ and space $\theta \in \Theta$ and $\xi \in \Xi$, with appropriate spatial domains $\Theta$ and $\Xi$. The framework consists of two steps.

1. Approximation of signals (cf. Section 3). We choose finite-dimensional subspaces $\overline{\mathcal{U}} \subset \mathcal{U}$ and $\overline{\mathcal{Y}} \subset \mathcal{Y}$ with bases $\left\{u_{1}, \ldots, u_{\bar{p}}\right\} \subset \overline{\mathcal{U}}$ and $\left\{y_{1}, \ldots, y_{\bar{q}}\right\} \subset \overline{\mathcal{Y}}$, and denote the corresponding orthogonal projections by $\mathbb{P}_{\overline{\mathcal{U}}}$ and $\mathbb{P}_{\overline{\mathcal{Y}}}$, respectively. Then, the approximation

$$
\mathbb{G}_{S}=\mathbb{P}_{\overline{\mathcal{Y}}} \mathbb{G P}_{\bar{u}}
$$

has a matrix representation $\mathbf{G} \in \mathbb{R}^{\bar{q} \times \bar{p}}$, for instance with elements $\mathbf{G}_{i j}=$ $\left(y_{i}, \mathbb{G} u_{j}\right)_{\mathcal{Y}}$ if orthonormal bases are chosen in $\overline{\mathcal{U}}$ and $\overline{\mathcal{Y}}$.
2. Approximation of system dynamics (cf. Section 4). Frequently, $\mathbb{G}$ arises from a linear PDE state space model. Then the components $\mathbf{G}_{i j}=\left(y_{i}, \mathbb{G} u_{j}\right)_{\mathcal{y}}$ can be approximated by numerically simulating the state space model successively for inputs $u_{j}, j=1, \ldots, \bar{p}$ and by testing the resulting outputs against all $y_{1}, \ldots, y_{\bar{q}}$. The result is an approximation $\mathbb{G}_{D S}$ of $\mathbb{G}_{S}$. Considering time-invariant systems and choosing basis functions with a space-time tensor structure, like

$$
u_{i(j, l)}(t ; \theta)=\phi_{j}(t) \mu_{l}(\theta), \quad y_{j(i, k)}(t ; \xi)=\psi_{i}(t) \nu_{k}(\xi)
$$

this task reduces to the approximation of observations $\left(\nu_{k}, C z_{l}(t)\right)_{Y}$, with states $z_{l}(t)=S(t) B \mu_{l}$. Here $S(t), B$ and $C$ are the system's evolution semigroup, input and output operator, respectively. Hence, $C z_{l}(t)$ can be considered as the system's impulse response corresponding to an initial value $\mu_{l}$, and $z_{l}(t)$ can be approximated by numerically solving a homogeneous PDE. We discuss some prospects of this framework.

Error estimation (cf. Section 5). The total error $\epsilon_{D S}$ can be estimated by the signal approximation error $\epsilon_{S}$ and the dynamical approximation error $\epsilon_{D}$, i.e.

$$
\begin{equation*}
\underbrace{\left\|\mathbb{G}-\mathbb{G}_{D S}\right\|}_{=: \epsilon_{D S}} \leq \underbrace{\left\|\mathbb{G}-\mathbb{G}_{S}\right\|}_{=: \epsilon_{S}}+\underbrace{\left\|\mathbb{G}_{S}-\mathbb{G}_{D S}\right\|}_{=: \epsilon_{D}}, \tag{1.1}
\end{equation*}
$$

where the norms still have to be specified. As main result of this paper, Thm. 5.1 shows how to choose $\overline{\mathcal{U}}$ and $\overline{\mathcal{Y}}$ in the first step and the accuracy tolerances for the
numerical solutions of the underlying PDEs in the second step such that $\epsilon_{S}$ and $\epsilon_{D}$ balance and that $\epsilon_{S}+\epsilon_{D}<$ tol for a given tolerance tol.

Progressive reduction of the signal error. Choosing hierarchical bases in $\overline{\mathcal{U}}$ and $\overline{\mathcal{Y}}$, the error $\epsilon_{S}$ can be progressively reduced by adding further basis functions $u_{\bar{p}+1}, u_{\bar{p}+2}, \ldots$ and $y_{\bar{q}+1}, y_{\bar{q}+2}, \ldots$ resulting in additional columns and rows of the matrix representation.

Matrix reduction via multilinear SVDs (cf. Section 6). The matrix representation of $\mathbb{G}_{D S}$ allows for low rank approximations with error estimates on the basis of socalled higher order singular value decompositions (HOSVDs) [19], respecting the timespace tensor structure of the basis functions. The corresponding singular vectors represent the most relevant input and output signals.

Actuators and sensors for distributed inputs and outputs. Thinking of practical applications, input signals $u(t ; \theta)$ and output signals $y(t ; \xi)$ are often generated and measured by actuators and sensors with limited spatial and temporal resolutions, such that 'realizable' input and output signals naturally belong to finite dimensional subspaces $\overline{\mathcal{U}}$ and $\overline{\mathcal{Y}}$, respectively. Error estimates of the form (1.1) and the extraction of relevant input and output signals on the basis of HOSVDs may thus provide useful information for efficient sensor and actuator design, see Section 6. Note that classical approaches (where the control system is first discretized in space and then model reduction is applied) rarely take the error due to input and output space-discretizations into account.

Control Design (cf. Section 6). The matrix representation $\mathbf{G}=\left[\mathbf{G}_{i j}\right]$ may directly be used in control design, or a state realization of the i/o model $\mathbb{G}_{D S}$ can be used as basis for many classical control design algorithms.

### 1.3. Relation to numerical analysis, control theory and optimal control.

 From the point of view of numerical analysis, the presented approach is a Galerkin approximation of the i/o map, which is a Volterra integral operator arising from the semigroup representation of the evolution system. The corresponding error estimates are based on standard interpolation theory in Sobolev spaces and on error results for the numerical solution of evolution equations.From the point of view of control theory, the linear time-continuous infinitedimensional system with distributed controls and observations is first approximated by a time-discrete multi-input-multi-output system, the corresponding Markov parameters are then approximated by numerically calculating impulse responses.

Using the approximated i/o map in optimal control applications corresponds to a pronounced form of the concept 'first discretize, then optimize' since the original system is discretized in space and time. On the one hand, this entails the risk of loosing essential structural features of the original control problem, which may lead e.g. to instabilities or simply failing of the calculated controls. The analytical investigation of the behavior of the approximated i/o map in control applications is an important future task. On the other hand, the algebraic representation of the i/o-map enables the use of very fast methods for model reduction and control design, and the error results for the full discretization may help to take effects of digitizing inputs and outputs for processing by discrete controllers into account.

Finally, we aim to mention some other approaches which directly focus on the original infinite-dimensional system or control problem. Balanced Truncation and POD have been formulated in infinite-dimensional function spaces, see e.g. [14, 18, 43]. In sophisticated simulations and optimal control applications, state and control discretization errors and even modeling errors are adaptively controlled with respect
to their effect on quantities of interest, see e.g. [11, 15].
1.4. Notation. For $\Omega \subset \mathbb{R}^{d}, d \in \mathbb{N}, L^{2}(\Omega)$ denotes the usual Lebesgue space of square-integrable functions, and $H^{\alpha}(\Omega), \alpha \in \mathbb{N}_{0}$ denotes the corresponding Sobolev spaces of $\alpha$-times weakly differentiable functions. We interpret functions $v$, which vary in space and time, optionally as classical functions $v:[0, T] \times \Omega \rightarrow \mathbb{R}$ with values $v(t ; x) \in \mathbb{R}$, or as abstract functions $v:[0, T] \rightarrow \mathbb{R}$ with values in a function space $X$ such as $X=H^{\alpha}(\Omega)$. Correspondingly, $H^{\alpha}\left(0, T ; H^{\beta}(\Omega)\right)$, with $\alpha, \beta \in \mathbb{N}_{0}$, denotes the space of equivalence classes of functions $v:[0, T] \rightarrow H^{\beta}(\Omega)$ with $t \mapsto\|v\|_{H^{\beta}(\Omega)}$ being $\alpha$-times weakly differentiable, for details see e.g. [22]. We introduce Hilbert spaces

$$
\begin{aligned}
H^{\alpha, \beta}((0, T) \times \Omega) & :=H^{\alpha}\left(0, T ; L^{2}(\Omega)\right) \cap L^{2}\left(0, T ; H^{\beta}(\Omega)\right) \\
& \|v\|_{H^{\alpha, \beta}((0, T) \times \Omega)}:=\|v\|_{H^{\alpha}\left(0, T ; L^{2}(\Omega)\right)}+\|v\|_{L^{2}\left(0, T ; H^{\beta}(\Omega)\right)}
\end{aligned}
$$

see e.g. [37]. By $C([0, T] ; X)$ and $C^{\alpha}([0, T] ; X)$ we denote the space of functions $v:[0, T] \rightarrow X$ which are continuous respectively $\alpha$-times continuously differentiable. For two normed spaces $X$ and $Y, \mathscr{L}(X, Y)$ denotes the set of bounded linear operators $X \rightarrow Y$, and we abbreviate $\mathscr{L}(X):=\mathscr{L}(X, X)$. For $\alpha \in \mathbb{N}, L^{\alpha}(0, T ; \mathscr{L}(X, Y))$ denotes the space of operator-valued functions $K:[0, T] \rightarrow \mathscr{L}(X, Y)$ with $t \mapsto$ $\|K(t)\|_{\mathscr{L}(X, Y)}=\sup _{x \neq 0}\|K(t) x\|_{Y} /\|x\|_{X}$ lying in $L^{\alpha}(0, T)$. Vectors, often representing a discretization of a function $v$, are written in corresponding small bold letters $\mathbf{v}$, whereas matrices, often representing a discrete version of an operator like $\mathbb{G}$ or $G$, are written in bold capital letters $\mathbf{G} . \mathbb{R}^{\alpha \times \beta}$ stands for the set of real $\alpha \times \beta$ matrices, and $\mathbf{A} \otimes \mathbf{B}$ denotes the Kronecker tensor product of two matrices $\mathbf{A}$ and $\mathbf{B}$.
2. I/o maps of $\infty$-dimensional LTI state space systems. We consider infinite-dimensional linear time-invariant systems of first order

$$
\begin{align*}
\partial_{t} z(t) & =A z(t)+B u(t), \quad t \in(0, T]  \tag{2.1a}\\
z(0) & =z^{0},  \tag{2.1~b}\\
y(t) & =C z(t), \quad t \in[0, T] \tag{2.1c}
\end{align*}
$$

Here for every time $t \in[0, T]$, the state $z(t)$ is supposed to belong to a Hilbert space $Z$ like $Z=L^{2}(\Omega)$, where $\Omega$ is a subset of $\mathbb{R}^{d_{\Omega}}$ with $d_{\Omega} \in \mathbb{N}$. $A$ is a densely defined unbounded operator $A: Z \supset D(A) \rightarrow Z$, generating a $C^{0}$-semigroup $(S(t))_{t \geq 0}$ on $Z$. The control operator $B$ belongs to $\mathscr{L}(U, Z)$ and the observation operator $C$ to $\mathscr{L}(Z, Y)$, where $U=L^{2}(\Theta)$ and $Y=L^{2}(\Xi)$ with subsets $\Theta \subset \mathbb{R}^{d_{1}}$ and $\Xi \subset \mathbb{R}^{d_{2}}$, $d_{1}, d_{2} \in \mathbb{N}$.

We recall how a linear bounded i/o-map $\mathbb{G} \in \mathscr{L}(\mathcal{U}, \mathcal{Y})$ with

$$
\mathcal{U}=L^{2}(0, T ; U) \quad \text { and } \quad \mathcal{Y}=L^{2}(0, T ; Y)
$$

can be associated to (2.1), for details see e.g. [41, Ch. 4]. It is well-known that for initial values $z_{0} \in D(A)$ and controls $u \in C^{1}([0, T] ; Z)$, a unique classical solution $z \in C([0, T] ; Z) \cap C^{1}((0, T) ; Z)$ of (2.1) exists. For $z_{0} \in Z$ and $u \in \mathcal{U}$, the well-defined function

$$
\begin{equation*}
z(t)=S(t) z_{0}+\int_{0}^{t} S(t-s) B u(s) d s, \quad t \in[0, T] \tag{2.2}
\end{equation*}
$$

is called a mild solution of (2.1). A mild solution of (2.1) is unique, belongs to $C([0, T] ; Z)$ and is the uniform limit of classical solutions [41]. Hence, the output
signal $y(t)=C z(t)$ is well-defined and belongs to $\mathcal{Y} \cap C([0, T] ; Y)$. In particular, the ouput signals $y(u) \in \mathcal{Y}$ arising from input signals $u \in \mathcal{U}$ and zero initial conditions $z_{0} \equiv 0$ allow to define the linear i/o-map $\mathbb{G}: \mathcal{U} \rightarrow \mathcal{Y}$ of the system (2.1) by $u \mapsto$ $y(u)$. It is possible to represent $\mathbb{G}$ as a convolution with the kernel function $K \in$ $L^{2}(-T, T ; \mathscr{L}(U, Y))$,

$$
K(t)=\left\{\begin{array}{ll}
C S(t) B, & t \geq 0  \tag{2.3}\\
0, & t<0
\end{array} .\right.
$$

Lemma 2.1. The $i / o-m a p ~ \mathbb{G}$ of (2.1) has the representation

$$
\begin{equation*}
(\mathbb{G} u)(t)=\int_{0}^{T} K(t-s) u(s) d s, \quad t \in[0, T] \tag{2.4}
\end{equation*}
$$

belongs to $\mathscr{L}(\mathcal{U}, \mathcal{Y}) \cap \mathscr{L}(\mathcal{U}, C([0, T], \mathcal{Y}))$ and satisfies

$$
\begin{equation*}
\|\mathbb{G}\|_{\mathscr{L}(\mathcal{U}, \mathcal{Y})} \leq \sqrt{T}\|K\|_{L^{2}(0, T ; \mathscr{L}(U, Y))} \tag{2.5}
\end{equation*}
$$

Proof. Since $C$ is bounded, the representation of $y=C z$ based on (2.2) can be reformulated as in (2.4), see e.g. [22] for the theory of Bochner integrals. For general $K \in L^{2}(-T, T ; \mathscr{L}(U, Y))$, a generalized Hölder's inequality yields that for fixed $t \in[0, T]$, the function $s \rightarrow K(t-s) u(s)$ belongs to $L^{1}(0, T ; \mathscr{L}(U, Y))$ with

$$
\|(\mathbb{G} u)(t)\|_{Y} \leq\|u\|_{\mathcal{U}}\|K(t-\cdot)\|_{L^{2}(0, T ; \mathscr{L}(U, Y)},
$$

and by integrating over $[0, T]$ we obtain (2.5).
Remark 1. The i/o-map $\mathbb{G}$ is causal in the sense that $y(t)$ only depends on $u_{\mid[0, t)}$ for all $t \in[0, T]$, and $\mathbb{G}$ is time-invariant in the sense that if $y=\mathbb{G} u$ then $\sigma_{\tau} y=\mathbb{G}\left(\sigma_{\tau} u\right)$ for all $\tau \in[0, T]$. Here $\sigma_{\tau}$ is a shift operator with $\left(\sigma_{\tau} u\right)(t)=u(t-\tau)$ for $t \in[\tau, T]$ and $\left(\sigma_{\tau} u\right)(t)=0$ for $t \in[0, \tau)$.

EXAMPLE 1. As prototype for a parabolic system, we consider the heat equation with homogeneous Dirichlet boundary conditions and assume that $\Omega$ has a $C^{2}$ boundary. In this case, $Z=L^{2}(\Omega)$ and the operator $A$ in (2.1) coincides with the Laplace operator

$$
\begin{equation*}
A=\triangle: D(A)=H^{2}(\Omega) \cap H_{0}^{1}(\Omega) \subset Z \rightarrow Z \tag{2.6}
\end{equation*}
$$

Since $A$ is the infinitesimal generator of an analytic $C^{0}$-semigroup of contractions $(S(t))_{t \geq 0}$, the mild solution $z$ of (2.1) exhibits the following stability and regularity properties, see e.g. [41, Ch. 7] and [25].
(i) If $z_{0}=0$ and $u \in \mathcal{U}$, then $z \in H^{1,2}((0, T) \times \Omega)$ with

$$
\begin{equation*}
\|z\|_{H^{1,2}((0, T) \times \Omega)} \leq c\|u\|_{\mathcal{U}} \tag{2.7}
\end{equation*}
$$

(ii) Assume that $u \equiv 0$. For $z_{0} \in D(A)$ we have $z \in C^{1}([0, T] ; D(A))$, but for $z_{0} \in Z$ we only have $z \in C^{1}((0, T] ; D(A))$.
We will consider concrete choices of $\Omega, B$ and $C$ in Section 6. We note that if the observation preserves the inherent state regularity in the sense that

$$
\begin{equation*}
C_{\mid H^{2}(\Omega)} \in \mathscr{L}\left(H^{2}(\Omega), H^{2}(\Xi)\right) \tag{2.8}
\end{equation*}
$$

then $\mathbb{G} \in \mathscr{L}\left(\mathcal{U}, \mathcal{Y}_{s}\right)$ and also

$$
\begin{equation*}
\mathbb{G}_{\mid \mathcal{U}_{s}} \in \mathscr{L}\left(\mathcal{U}_{s}, \mathcal{Y}_{s}\right), \quad \text { with } \mathcal{U}_{s}=H^{1,2}((0, T) \times \Theta), \quad \mathcal{Y}_{s}=H^{1,2}((0, T) \times \Xi) \tag{2.9}
\end{equation*}
$$

In fact, for $u \in \mathcal{U}_{s}$, we have $\|u\|_{\mathcal{U}} \leq\|u\|_{\mathcal{U}_{s}}$, and for $u \in \mathcal{U}$, we have $\|\mathbb{G} u\|_{\mathcal{Y}_{s}} \leq$
 and $c$ is the constant in (2.7).

Remark 2. Many other linear time-invariant systems with distributed controls and observations admit a representation of the i/o map via (2.4) and exhibit properties similar to (2.9). This is, for instance, the case for the heat equation with homogeneous Neumann boundary conditions, and also for more general parabolic equations, see [37] and [38]. For Stokes systems, results similar to (2.4) and (2.9) are obtained by working with appropriate subspaces of divergence-free functions, see [44]. Wave equations with second order time derivatives can be represented in form of (2.1) and (2.4) by means of an order reduction. Though hyperbolic systems do not have the smoothing property of parabolic systems, they preserve the regularity of the data and results similar to (2.9) can be obtained due to the restriction to input signals of higher regularity in time, see [37, p. 95]. Note, however, that systems with boundary control or pointwise observation do not fit directly into the setting (2.1).

## 3. Discretization of signals.

3.1. Space-time discretization and matrix representation. In order to discretize the input signals $u \in \mathcal{U}$ and $y \in \mathcal{Y}$ in space and time, we choose four families $\left\{U_{h_{1}}\right\}_{h_{1}>0},\left\{Y_{h_{2}}\right\}_{h_{2}>0},\left\{\mathcal{R}_{\tau_{1}}\right\}_{\tau_{1}>0}$ and $\left\{\mathcal{S}_{\tau_{2}}\right\}_{\tau_{2}>0}$ of subspaces $U_{h_{1}} \subset U, Y_{h_{2}} \subset Y$, $\mathcal{R}_{\tau_{1}} \subset L^{2}(0, T)$ and $\mathcal{S}_{\tau_{2}} \subset L^{2}(0, T)$ of finite dimensions $p\left(h_{1}\right)=\operatorname{dim}\left(U_{h_{1}}\right), q\left(h_{2}\right)=$ $\operatorname{dim}\left(Y_{h_{2}}\right), r\left(\tau_{1}\right)=\operatorname{dim}\left(\mathcal{R}_{\tau_{1}}\right)$ and $s\left(\tau_{2}\right)=\operatorname{dim}\left(\mathcal{S}_{\tau_{2}}\right)$. We then define

$$
\begin{aligned}
& \mathcal{U}_{h_{1}, \tau_{1}}=\left\{u \in \mathcal{U}: u(t ; \cdot) \in U_{h_{1}}, u(\cdot ; \theta) \in \mathcal{R}_{\tau_{1}} \quad \text { for almost every } t \in[0, T], \theta \in \Theta\right\}, \\
& \mathcal{Y}_{h_{2}, \tau_{2}}=\left\{y \in \mathcal{Y}: y(t ; \cdot) \in Y_{h_{2}}, y(\cdot ; \xi) \in \mathcal{S}_{\tau_{2}} \quad \text { for almost every } t \in[0, T], \xi \in \Xi\right\}
\end{aligned}
$$

We denote the orthogonal projections onto these subspaces by $\mathbb{P}_{\mathcal{U}, h_{1}, \tau_{1}} \in \mathscr{L}(\mathcal{U})$ and $\mathbb{P}_{\mathcal{Y}, h_{2}, \tau_{2}} \in \mathscr{L}(\mathcal{Y})$. As first step of the approximation of $\mathbb{G}$, we define

$$
\begin{equation*}
\mathbb{G}_{S}=\mathbb{G}_{S}\left(h_{1}, \tau_{1}, h_{2}, \tau_{2}\right)=\mathbb{P}_{\mathcal{Y}, h_{2}, \tau_{2}} \mathbb{G}_{\mathbb{P}_{\mathcal{U}, h_{1}, \tau_{1}} \in \mathscr{L}(\mathcal{U}, \mathcal{Y}) . . . . . . . . ~} \tag{3.1}
\end{equation*}
$$

In order to obtain a matrix representation of $\mathbb{G}_{S}$, we introduce families of bases $\left\{\mu_{1}, \ldots, \mu_{p}\right\}$ of $U_{h_{1}},\left\{\nu_{1}, \ldots, \nu_{q}\right\}$ of $Y_{h_{2}},\left\{\phi_{1}, \ldots, \phi_{r}\right\}$ of $\mathcal{R}_{\tau_{1}}$ and $\left\{\psi_{1}, \ldots, \psi_{s}\right\}$ of $\mathcal{S}_{\tau_{2}}$ and corresponding mass matrices $\mathbf{M}_{U, h_{1}} \in \mathbb{R}^{p \times p}, \mathbf{M}_{Y, h_{2}} \in \mathbb{R}^{q \times q}, \mathbf{M}_{\mathcal{R}, \tau_{1}} \in \mathbb{R}^{r \times r}$ and $\mathbf{M}_{\mathcal{S}, \tau_{2}} \in \mathbb{R}^{s \times s}$, for instance via

$$
\left[\mathbf{M}_{U, h_{1}}\right]_{i j}=\left(\mu_{j}, \mu_{i}\right)_{U}, \quad i, j=1, \ldots, p
$$

These mass matrices induce, for instance via

$$
(\mathbf{v}, \mathbf{w})_{p ; w}=\mathbf{v}^{T} \mathbf{M}_{U, h_{1}} \mathbf{w} \quad \text { for all } \mathbf{v}, \mathbf{w} \in \mathbb{R}^{p}
$$

weighted scalar products and corresponding norms in the respective spaces, which we indicate by a subscript $w$, like $\mathbb{R}_{w}^{p}$ with $(\cdot, \cdot)_{p ; w}$ and $\|\cdot\|_{p ; w}$, in contrast to the canonical spaces like $\mathbb{R}^{p}$ with $(\cdot, \cdot)_{p}$ and $\|\cdot\|_{p}$. We represent signals $u \in \mathcal{U}_{h_{1}, \tau_{1}}$ and $y \in \mathcal{Y}_{h_{2}, \tau_{2}}$ as

$$
u(t ; \theta)=\sum_{k=1}^{p} \sum_{i=1}^{r} \mathbf{u}_{i}^{k} \phi_{i}(t) \mu_{k}(\theta), \quad y(t ; \xi)=\sum_{l=1}^{q} \sum_{j=1}^{s} \mathbf{y}_{j}^{l} \psi_{j}(t) \nu_{k}(\xi)
$$

where $\mathbf{u}_{i}^{k}$ are the elements of a block-structured vector $\mathbf{u} \in \mathbb{R}^{p r}$ with $p$ blocks $\mathbf{u}^{k} \in \mathbb{R}^{r}$, and the vector $\mathbf{y} \in \mathbb{R}^{q s}$ is defined similarly. Then

$$
\|u\|_{\mathcal{U}}=\|\mathbf{u}\|_{p r ; w}, \quad \text { and } \quad\|y\|_{\mathcal{Y}}=\|\mathbf{y}\|_{q s ; w}
$$

where $\|\cdot\|_{p r ; w}$ and $\|\cdot\|_{q s ; w}$ denote the weighted norms with respect to the mass matrices

$$
\mathbf{M}_{\mathcal{U}, h_{1}, \tau_{1}}=\mathbf{M}_{U, h_{1}} \otimes \mathbf{M}_{\mathcal{R}, \tau_{1}} \in \mathbb{R}^{p r \times p r}, \quad \mathbf{M}_{\mathcal{Y}, h_{2}, \tau_{2}}=\mathbf{M}_{Y, h_{2}} \otimes \mathbf{M}_{\mathcal{S}, \tau_{2}} \in \mathbb{R}^{q s \times q s}
$$

i.e. the corresponding coordinate isomorphisms $\kappa_{\mathcal{U}, h_{1}, \tau_{1}} \in \mathscr{L}\left(\mathcal{U}_{h_{1}, \tau_{1}}, \mathbb{R}_{w}^{p r}\right)$ and $\kappa_{\mathcal{Y}, h_{2}, \tau_{2}} \in$ $\mathscr{L}\left(\mathcal{Y}_{h_{2}, \tau_{2}}, \mathbb{R}_{w}^{q S}\right)$ are unitary.

Finally, we obtain a matrix representation $\mathbf{G}$ of $\mathbb{G}_{S}$ by setting

$$
\begin{equation*}
\mathbf{G}=\mathbf{G}\left(h_{1}, \tau_{1}, h_{2}, \tau_{2}\right)=\kappa \mathcal{Y} \mathbb{P}_{\mathcal{Y}} \mathbb{G P}_{\mathcal{U}} \kappa_{\mathcal{U}}^{-1} \in \mathbb{R}^{q s \times p r}, \tag{3.2}
\end{equation*}
$$

where the dependencies on $h_{1}, \tau_{1}, h_{2}, \tau_{2}$ have been partially omitted. Considering

$$
\mathbf{H}=\mathbf{H}\left(h_{1}, \tau_{1}, h_{2}, \tau_{2}\right):=\mathbf{M}_{\mathcal{Y}, h_{2}, \tau_{2}} \mathbf{G} \in \mathbb{R}^{q s \times p r}
$$

as a block-structured matrix with $q \times p$ blocks $\mathbf{H}^{k l} \in \mathbb{R}^{s \times r}$ and block elements $\mathbf{H}_{i j}^{k l} \in \mathbb{R}$, we obtain the element-wise representation

$$
\begin{equation*}
\mathbf{H}_{i j}^{k l}=\left[\mathbf{M}_{\mathcal{Y}} \kappa \mathcal{Y} \mathbb{P}_{\mathcal{Y}} \mathbb{G}\left(\mu_{l} \phi_{j}\right)\right]_{i}^{k}=\left(\nu_{k} \psi_{i}, \mathbb{G}\left(\mu_{l} \phi_{j}\right)\right)_{\mathcal{Y}} . \tag{3.3}
\end{equation*}
$$

Remark 3. Alternatively, $\mathbf{H}$ can be considered as a fourth-order tensor in $\mathbb{R}^{s \times r \times q \times p}$ with elements $\mathbf{H}_{i j k l}=\mathbf{H}_{i j}^{k l}$.

To have a discrete analogon of the $\mathscr{L}(\mathcal{U}, \mathcal{Y})$-norm, we introduce for given $\mathcal{U}_{h_{1}, \tau_{1}}$ and $\mathcal{Y}_{h_{2}, \tau_{2}}$ the weighted matrix norm

$$
\begin{equation*}
\left\|\mathbf{G}\left(h_{1}, \tau_{1}, h_{2}, \tau_{2}\right)\right\|_{q s \times p r ; w}:=\sup _{\mathbf{u} \in \mathbb{R}^{p r}} \frac{\|\mathbf{G} \mathbf{u}\|_{q s ; w}}{\|\mathbf{u}\|_{p r ; w}}=\left\|\mathbf{M}_{\mathcal{Y}, h_{2}, \tau_{2}}^{1 / 2} \mathbf{G M}_{\mathcal{U}, h_{1}, \tau_{1}}^{-1 / 2}\right\|_{q s \times p r}, \tag{3.4}
\end{equation*}
$$

and we write $\left(h_{1}^{\prime}, \tau_{1}^{\prime}, h_{2}^{\prime}, \tau_{2}^{\prime}\right) \leq\left(h_{1}, \tau_{1}, h_{2}, \tau_{2}\right)$ if the inequality holds component-wise.
Lemma 3.1. For all $\left(h_{1}, \tau_{1}, h_{2}, \tau_{2}\right) \in \mathbb{R}_{+}^{4}$, we have

$$
\begin{equation*}
\left\|\mathbf{G}\left(h_{1}, \tau_{1}, h_{2}, \tau_{2}\right)\right\|_{q s \times p r ; w}=\left\|\mathbb{G}_{S}\left(h_{1}, \tau_{1}, h_{2}, \tau_{2}\right)\right\|_{\mathscr{L}(\mathcal{U}, \mathcal{Y})} \leq\|\mathbb{G}\|_{\mathscr{L}(\mathcal{U}, \mathcal{Y})} \tag{3.5}
\end{equation*}
$$

If the subspaces $\left\{\mathcal{U}_{h_{1}, \tau_{1}}\right\}_{h_{1}, \tau_{1}>0}$ and $\left\{\mathcal{Y}_{h_{2}, \tau_{2}}\right\}_{h_{2}, \tau_{2}>0}$ are nested in the sense that

$$
\begin{equation*}
\mathcal{U}_{h_{1}, \tau_{1}} \subset \mathcal{U}_{h_{1}^{\prime}, \tau_{1}^{\prime}}, \quad \mathcal{Y}_{h_{2}, \tau_{2}} \subset \mathcal{Y}_{h_{2}^{\prime}, \tau_{2}^{\prime}} \quad \text { for }\left(h_{1}^{\prime}, \tau_{1}^{\prime}, h_{2}^{\prime}, \tau_{2}^{\prime}\right) \leq\left(h_{1}, \tau_{1}, h_{2}, \tau_{2}\right) \tag{3.6}
\end{equation*}
$$

then $\left\|\mathbf{G}\left(h_{1}, \tau_{1}, h_{2}, \tau_{2}\right)\right\|_{q s \times p r ; w}$ monotonically increases for decreasing $\left(h_{1}, \tau_{1}, h_{2}, \tau_{2}\right) \in$ $\mathbb{R}_{+}^{4}$, and $\left\|\mathbf{G}\left(h_{1}, \tau_{1}, h_{2}, \tau_{2}\right)\right\|_{q s \times p r ; w}$ is convergent for $\left(h_{1}, \tau_{1}, h_{2}, \tau_{2}\right) \searrow 0$.

Proof. In order to show (3.5), we calculate

$$
\left\|\mathbb{G}_{S}\right\|_{\mathscr{L}(\mathcal{U}, \mathcal{Y})}=\sup _{u \in \mathcal{U}_{h_{1}, \tau_{1}}} \frac{\left\|\mathbb{P}_{\mathcal{Y}, h_{2}, \tau_{2}} \mathbb{G} u\right\|_{\mathcal{Y}}}{\|u\|_{\mathcal{U}}} \leq \sup _{u \in \mathcal{U}_{h_{1}, \tau_{1}}} \frac{\|\mathbb{G} u\|_{\mathcal{Y}}}{\|u\|_{\mathcal{U}}} \leq\|\mathbb{G}\|_{\mathscr{L}(\mathcal{U}, \mathcal{Y})}
$$

and observe that for $u \in \mathcal{U}_{h_{1}, \tau_{1}}$ and $\mathbf{u}=\kappa_{\mathcal{U}, h_{1}, \tau_{1}} u \in \mathbb{R}^{p r}$, we have

$$
\begin{aligned}
&\left\|\mathbb{G}_{S} u\right\|_{\mathcal{Y}}=\left\|\kappa_{\mathcal{Y}, h_{2}, \tau_{2}}^{-1} \mathbf{G} \kappa_{\mathcal{U}, h_{1}, \tau_{1}} \mathbb{P}_{\mathcal{U}, h_{1}, \tau_{1}} u\right\|_{\mathcal{Y}}=\|\mathbf{G} \mathbf{u}\|_{q s ; w} \leq\|\mathbf{G}\|_{q s \times p r ; w}\|u\|_{\mathcal{U}}, \\
&\|\mathbf{G u}\|_{q s ; w} \leq \| \kappa \mathcal{Y}, h_{2}, \tau_{2} \\
& \mathbb{S}_{S} \kappa_{\mathcal{U}, h_{1}, \tau_{1}}^{-1} \mathbf{u}\left\|_{q_{s} ; w}=\right\| \mathbb{G}_{S} u\left\|_{\mathcal{Y}} \leq\right\| \mathbb{G}_{S}\left\|_{\mathscr{L}(\mathcal{U}, \mathcal{Y}}\right\| \mathbf{u} \|_{p r ; w} .
\end{aligned}
$$

Assume that (3.6) holds. Since $\left\|\mathbb{P}_{\mathcal{Y}, h_{2}, \tau_{2}} y\right\|_{\mathcal{Y}} \leq\left\|\mathbb{P}_{\mathcal{Y}, h_{2}^{\prime}, \tau_{2}^{\prime}} y\right\|_{\mathcal{Y}}$ for all $y \in \mathcal{Y}$, we have

$$
\left\|\mathbb{G}_{S}\left(h_{1}, \tau_{1}, h_{2}, \tau_{2}\right)\right\|_{q s \times p r ; w} \leq \sup _{u \in \mathcal{U}_{h_{1}^{\prime}, \tau_{1}^{\prime}}} \frac{\left\|\mathbb{P}_{\mathcal{Y}, h_{2}^{\prime}, \tau_{2}} \mathbb{G} u\right\|_{\mathcal{Y}}}{\|u\|_{\mathcal{U}}}=\left\|\mathbb{G}_{S}\left(h_{1}^{\prime}, \tau_{1}^{\prime}, h_{2}^{\prime}, \tau_{2}^{\prime}\right)\right\|_{q^{\prime} s^{\prime} \times p^{\prime} r^{\prime} ; w}
$$

Hence, (3.5) ensures the convergence of $\left\|\mathbb{G}_{S}(\mathbf{h})\right\|_{q s \times p r ; w}$.


Fig. 3.1. Hierarchical basis for $L^{2}(0,1)$-subspaces of piecewise linear functions: (a) $\mu_{1}$ and $\mu_{2}$ (b) $\mu_{3}$ (c) $\mu_{4}$ and $\mu_{5}$ (d) $\mu_{6}, \ldots, \mu_{9}$.


Fig. 3.2. Haar wavelet basis for $L^{2}(0,1)$-subspaces of piecewise constant functions: (a) $\phi_{1}$ (b) $\phi_{2}(c) \phi_{3}$ and $\phi_{4}(d) \phi_{5}, \ldots, \phi_{8}$.
3.2. An example for signal discretizations. As an example, consider the case $U=Y=L^{2}(0,1)$, and choose $U_{h_{1}}$ and $Y_{h_{2}}$ as spaces of continuous piecewise linear functions and $\mathcal{R}_{\tau_{1}}$ and $\mathcal{S}_{\tau_{2}}$ as spaces of piecewise constant functions, all with respect to equidistant grids.

For $p \in \mathbb{N}, p \geq 2$ and $h_{1}(p)=1 /(p-1)$, let $\mathscr{T}_{h_{1}}=\left\{I_{k}\right\}_{1 \leq k \leq p-1}$ be the equidistant partition of $(0,1]$ into intervals $I_{k}=\left((k-1) h_{1}, k h_{1}\right]$. The corresponding space $U_{h_{1}}$ of continuous piecewise linear functions is, for instance, spanned by the nodal basis

$$
\left\{\mu_{1}^{\left(h_{1}\right)}, \ldots, \mu_{p\left(h_{1}\right)}^{\left(h_{1}\right)}\right\} \subset U_{h_{1}}, \quad \text { with } \mu_{l}^{\left(h_{1}\right)}\left(k h_{1}\right)=\delta_{l-1}(k), \quad k=0, \ldots, p,
$$

i.e. the $\mu_{k}^{\left(h_{1}\right)}$ are the well-known hat functions. The subspaces $\left\{U_{h_{1}}\right\}$ are nested if the choice is restricted to $h_{1} \in\left\{2^{-n}\right\}_{n \in \mathbb{N}_{0}}$ and $p \in\left\{2^{n}+1\right\}_{n \in \mathbb{N}_{0}}$. Since the nodal bases of $U_{h_{1}}$ and $U_{h_{1}^{\prime}}$ do not have any common element for $h_{1} \neq h_{1}^{\prime}$, one may prefer to choose a hierarchical basis of finite element functions $\hat{\mu}_{l}$, as in Fig. 3.1, see e.g. [49], [50]. Then, $U_{h_{1}}=\operatorname{span}\left\{\hat{\mu}_{1}, \ldots, \hat{\mu}_{p\left(h_{1}\right)}\right\}$ for all $h_{1} \in\left\{2^{-n}\right\}_{n \in \mathbb{N}_{0}}$ with basis functions $\hat{\mu}_{k}$ independent of $h_{1}$.

For $r \in \mathbb{N}$ and $\tau_{1}=T / r$, let $\Gamma_{\tau_{1}}=\left\{I_{j}\right\}_{1 \leq j \leq r}$ be the equidistant partition of $(0, T]$ into intervals $I_{j}=\left((j-1) \tau_{1}, j \tau_{1}\right]$. The corresponding space $\mathcal{R}_{\tau_{1}}$ of piecewise constant functions is, for instance, spanned by the nodal and orthogonal basis

$$
\begin{equation*}
\left\{\phi_{1}^{\left(\tau_{1}\right)}(t), \ldots, \phi_{r}^{\left(\tau_{1}\right)}(t)\right\}, \quad \text { with } \phi_{j}^{\left(\tau_{1}\right)}(t)=\chi_{I_{j}}(t), \quad j=1, \ldots, r \tag{3.7}
\end{equation*}
$$

The spaces are nested by requiring $\tau_{1} \in\left\{2^{-n} T\right\}_{n \in \mathbb{N}_{0}}$. An orthonormal hierarchical basis for $\mathcal{R}_{\tau_{1}}$ is obtained by choosing $\phi_{j}$ as Haar-wavelets, cf. Fig. 3.2 and [17].

Denoting the orthogonal projections onto $U_{h_{1}}$ and $\mathcal{R}_{\tau_{1}}$ by $P_{U, h_{1}}$ and $P_{\mathcal{R}, \tau_{1}}$, respectively, the Poincaré-Friedrich's inequality shows that there exist constants $c_{U}=1 / 2$
and $c_{\mathcal{R}}=1 / \sqrt{2}$, independent of $h_{1}, \tau_{1}$ and $T$, such that

$$
\begin{array}{ll}
\left\|u-P_{U_{h_{1}}} u\right\|_{L^{2}(0,1)} \leq c_{U} h_{1}^{2}\left\|\partial_{\xi}^{2} u\right\|_{L^{2}(0,1)} & \text { for } u \in H^{2}(0,1), \\
\left\|v-P_{\mathcal{R}_{\tau_{1}}} v\right\|_{L^{2}(0, T)} \leq c_{\mathcal{R}} \tau_{1}\left\|\partial_{t} v\right\|_{L^{2}(0, T)} & \text { for } v \in H^{1}(0, T), \tag{3.8b}
\end{array}
$$

see e.g. [16, 51]. By Fubini's theorem, it follows that the corresponding projection $\mathbb{P}_{\mathcal{U}, h_{1}, \tau_{1}}$ onto $\mathcal{U}_{h_{1}, \tau_{1}}=\left\{u \in \mathcal{U}, u_{\mid I_{j}} \equiv u^{(j)}, u^{(j)} \in U_{h_{1}}, \quad j=1, \ldots, r\right\}$ satisfies

$$
\begin{equation*}
\left\|u-\mathbb{P}_{\mathcal{U}, h_{1}, \tau_{1}} u\right\|_{\mathcal{U}} \leq\left(c_{U} h_{1}^{2}+c_{\mathcal{R}} \tau_{1}\right)\|u\|_{\mathcal{U}_{s}} \quad \text { for all } u \in \mathcal{U}_{s}=H^{1,2}((0, T) \times(0,1)) . \tag{3.9}
\end{equation*}
$$

We define $Y_{h_{2}}, \mathcal{R}_{\tau_{2}}$ and $\mathcal{Y}_{h_{2}, \tau_{2}}$ accordingly and a corresponding estimate as (3.9) holds for the projection $\mathbb{P}_{\mathcal{Y}, h_{2}, \tau_{2}} y$ of elements $y \in \mathcal{Y}_{s}=\mathcal{U}_{s}$.

REMARK 4. Estimates similar to (3.9) also exist for domains $\Theta \subset \mathbb{R}^{d}$ with $d \geq 1$ and are classical results from the interpolation theory in Sobolev spaces, see e.g. [16]. Note that the interpolation constants then often have to be estimated numerically. Estimates with higher approximation orders can be obtained, if ansatz functions of higher polynomial degree are used and if the input and output signals exhibit corresponding higher regularity in space and time.

### 3.3. Interpretation as discrete-time multi-input-multi-output system.

 $\mathbb{G}_{S}$ can be considered as a generalization of a classical linear discrete-time multi-input-multi-output (MIMO) system. Input signals $u \in \mathcal{U}_{h_{1}, \tau_{1}}$ and output signals $y \in \mathcal{Y}_{h_{2}, \tau_{2}}$ can be uniquely represented by finite sequences$$
\left\{\mathbf{u}_{1}, \ldots, \mathbf{u}_{r}\right\} \subset \mathbb{R}^{p}, \quad\left\{\tilde{\mathbf{y}}_{1}, \ldots, \tilde{\mathbf{y}}_{s}\right\} \subset \mathbb{R}^{q}
$$

with coefficient vectors $\mathbf{u}_{i}=\left(\mathbf{u}_{i}^{1}, \ldots, \mathbf{u}_{i}^{p}\right)^{T}$ and $\tilde{\mathbf{y}}_{i}=\left(\tilde{\mathbf{y}}_{i}^{1}, \ldots, \tilde{\mathbf{y}}_{i}^{q}\right)^{T}$, where $\mathbf{u}=$ $\kappa_{\mathcal{U}, h_{1}, \tau_{1}} u$ and $\tilde{\mathbf{y}}=\mathbf{M}_{\mathcal{Y}, h_{2}, \tau_{2}} \kappa \mathcal{Y}, h_{2}, \tau_{2} y$. Hence, $y=\mathbb{G}_{S} u$ writes as

$$
\begin{equation*}
\Sigma: \quad \tilde{\mathbf{y}}_{i}=\sum_{j \in \mathbb{Z}} \mathbf{H}_{i j} \mathbf{u}_{j}, \quad i \in \mathbb{Z} \tag{3.10}
\end{equation*}
$$

where $\mathbf{H}_{i j}=\left[\mathbf{H}_{i j}^{k l}\right]_{k l} \in \mathbb{R}^{q \times p}$ for $1 \leq i \leq s$ and $1 \leq j \leq r$, and we define $\mathbf{H}_{i j}:=$ $0 \in \mathbb{R}^{q \times p}$ for other $i, j \in \mathbb{Z}$. (3.10) is the external representation of a general linear discrete-time system $\Sigma$ with $p$ input channels and $q$ output channels, see e.g. [2]. In this context, $\mathbf{u}_{j}$ (resp. $\tilde{\mathbf{y}}_{j}$ ) usually denotes the input value (resp. output value) at the point of time $t_{j}=j \tau$ with some fixed time step size $\tau$. $\Sigma$ is called causal if $\mathbf{H}_{i j}=0$ for $j>i$ and $\Sigma$ is called time-invariant if $\mathbf{H}_{i j}=\mathbf{H}_{i-j}$. For a time-invariant system, the sequence of $q \times p$ constant matrices $h_{\Sigma}=\left(\ldots, \mathbf{H}_{-2}, \mathbf{H}_{-1}, \mathbf{H}_{0}, \mathbf{H}_{1}, \mathbf{H}_{2}, \ldots\right)$ is called the impulse response of $\Sigma$ because it is the output obtained to a unit impulse $\mathbf{u}_{j}=\delta_{0}(j)$. For a causal time-invariant system, the matrices $\mathbf{H}_{0}, \mathbf{H}_{1}, \ldots$ are often referred to as the Markov parameters of $\Sigma$.

The causality and time-invariance of the i/o-map $\mathbb{G}$ (cf. Rem. 1) transfer to the causality and time-invariance of $\Sigma$ if the time basis functions $\phi_{j}$ and $\psi_{i}$ are chosen in an appropriate way. Choosing, for instance, $\phi_{1}=\psi_{1}, \ldots, \phi_{r}=\psi_{r}$ as in (3.7) for some $r=s \in \mathbb{N}$ and $\tau=\tau_{1}=\tau_{2}=1 / r$, the matrices $\mathbf{H}^{k l} \in \mathbb{R}^{r \times r}$ are lower triangular Toeplitz matrices and $\mathbf{u}_{j}$ and $\tilde{\mathbf{y}}_{j}$ can be interpreted as signal values at time $t_{j}=j \tau$, cf. Remark 7. Note that all $\phi_{j}$ and $\psi_{i}$ satisfying $\phi_{j}=\sigma_{(j-1) \tau} \phi_{1}$ and $\psi_{i}=\sigma_{(i-1) \tau} \psi_{1}$ will also ensure this property.

The problem of finding an internal or state space representation of $\Sigma$ with minimal state space dimension is referred to as minimal partial realization problem, and in particular for causal time-invariant systems, many algorithms exist to solve it, see e.g. [13, 21, 39].
3.4. Signal approximation error estimates. We first consider the signal error $\epsilon_{D}$ with respect to the $\mathscr{L}(\mathcal{U}, \mathcal{Y})$-norm.

Lemma 3.2. The signal approximation error $\epsilon_{s}:=\left\|\mathbb{G}-\mathbb{G}_{S}\right\|_{\mathscr{L}(\mathcal{U}, \mathcal{Y})}$ decomposes into $\epsilon_{s}=\epsilon_{s, \text { inp }}+\epsilon_{s, \text { outp }}$ with

$$
\epsilon_{s, \text { inp }}:=\sup _{u \in \operatorname{ker} \mathbb{P}_{\mathcal{U}, h_{1}, \tau_{1}}} \frac{\|\mathbb{G} u\|_{\mathcal{Y}}}{\|u\|_{\mathcal{U}}}, \quad \epsilon_{s, \text { outp }}:=\max _{u \in \mathcal{U}_{h_{1}, \tau_{1}}} \frac{\left.\|\left(I-\mathbb{P}_{\mathcal{Y}, h_{2}, \tau_{2}}\right) \mathbb{G} u\right) \| \mathcal{Y}}{\|u\|_{\mathcal{U}}} .
$$

Proof. We estimate

$$
\begin{equation*}
\epsilon_{s} \leq\left\|\mathbb{G}\left(\mathbb{I}-\mathbb{P}_{\mathcal{U}, h_{1}, \tau_{1}}\right)\right\|_{\mathscr{L}(\mathcal{U}, \mathcal{Y})}+\left\|\left(\mathbb{I}-\mathbb{P}_{\mathcal{Y}, h_{2}, \tau_{2}}\right) \mathbb{G P}_{\mathcal{U}, h_{1}, \tau_{1}}\right\|_{\mathscr{L}(\mathcal{U}, \mathcal{Y})} \tag{3.11}
\end{equation*}
$$

and observe that

$$
\epsilon_{s, \text { outp }}=\left\|\left(\mathbb{I}-\mathbb{P}_{\mathcal{Y}, h_{2}, \tau_{2}}\right) \mathbb{G P}_{\mathcal{U}, h_{1}, \tau_{1}}\right\|_{\mathscr{L}(\mathcal{U}, \mathcal{Y})}, \quad \epsilon_{s, \text { inp }}=\left\|\mathbb{G}\left(\mathbb{I}-\mathbb{P}_{\mathcal{U}, h_{1}, \tau_{1}}\right)\right\|_{\mathscr{L}(\mathcal{U}, \mathcal{Y})} .
$$

We obtain the equality in (3.11) by considering a sequence $u_{j}=u^{*}+u_{j}^{\prime}$, where $u^{*}$ is the maximizer in the definition of $\epsilon_{s, \text { outp }}$ and $\left(u_{j}^{\prime}\right)_{j} \subset \operatorname{ker} \mathbb{P}_{\mathcal{U}, h_{1}, \tau_{1}}$ is a supremal sequence in the definition of $\epsilon_{s, i n p}$. $\square$

The next remarks show that we can only hope for a good approximation in $\|\cdot\|_{\mathscr{L}(\mathcal{U}, \mathcal{Y})}$ if the subspaces $\mathcal{U}_{h_{1}, \tau_{1}}$ and $\mathcal{Y}_{h_{2}, \tau_{2}}$ can be chosen specifically for $\mathbb{G}$ such that output signals from input signals $u \in \mathcal{U}_{h_{1}, \tau_{1}}$ are well approximated in $\mathcal{Y}_{h_{2}, \tau_{2}}$ and that neglected input signal components in $\operatorname{ker} \mathbb{P}_{\mathcal{U}, h_{1}, \tau_{1}}$ only lead to small output signals.

Remark 5. A usual requirement for families of approximating subspaces $\mathcal{U}_{h_{1}, \tau_{1}}$ and $\mathcal{Y}_{h_{2}, \tau_{2}}$ is that they become dense if $h_{1}, \tau_{1}, h_{2}, \tau_{2} \rightarrow 0$. We note that this condition implies that $\left\|\left(\mathbb{G}-\mathbb{G}_{S}\right) u\right\|_{\mathcal{Y}} \rightarrow 0$ for every $u \in \mathcal{U}$, but does not guarantee the uniform convergence $\left\|\mathbb{G}-\mathbb{G}_{S}\right\|_{\mathscr{L}(\mathcal{U}, \mathcal{Y})} \rightarrow 0$. Considering, for instance, the identity operator $\mathbb{G}=I d \in \mathscr{L}(\mathcal{U}, \mathcal{Y})$ in the case $\mathcal{U}=\mathcal{Y}$, $\epsilon_{s, \text { inp }}$ equals one for every finite-dimensional $\mathcal{U}_{h_{1}, \tau_{1}}$. Similar effects can be expected for operators with feedthrough components acting between infinite-dimensional subspaces.

Remark 6. If $\mathbb{G} \in \mathscr{L}(\mathcal{U}, \mathcal{Y})$ is a compact operator, then there exist orthonormal systems $\left\{\hat{u}_{1}, \hat{u}_{2} \ldots\right\}$ of $\mathcal{U}$ and $\left\{\hat{y}_{1}, \hat{y}_{2}, \ldots\right\}$ of $\mathcal{Y}$ and nonnegative numbers $\sigma_{1} \geq \sigma_{2} \geq \ldots$ with $\sigma_{k} \rightarrow 0$ such that $\mathbb{G} u=\sum_{k=1}^{\infty} \sigma_{k}\left(u, \hat{u}_{k}\right) \mathcal{U} \hat{y}_{k}$ for all $u \in \mathcal{U}$, see e.g. [47]. Choosing $\mathcal{U}_{h_{1}, \tau_{1}}$ and $\mathcal{Y}_{h_{2}, \tau_{2}}$ as the span of $\hat{u}_{1}, \ldots, \hat{u}_{r}$ and $\hat{y}_{1}, \ldots, \hat{y}_{s}$, respectively, with $s=r$ and $r \in \mathbb{N}$, we obtain an efficient approximation $\mathbb{G}_{S}$ of $\mathbb{G}$ with $\left\|\mathbb{G}-\mathbb{G}_{S}\right\|_{\mathscr{L}(\mathcal{U}, \mathcal{Y})} \leq \sigma_{r+1}$.

The case where less specific information about $\mathbb{G}$ is available and we only know that $\mathbb{G}_{\mathcal{U}_{s}} \in \mathscr{L}\left(\mathcal{U}_{s}, \mathcal{Y}_{s}\right)$ with spaces of higher regularity in space and time $\mathcal{U}_{s} \subset \mathcal{U}$ and $\mathcal{Y}_{s} \subset \mathcal{Y}$ like in Example 1 will be considered in Theorem 5.1.
4. Approximation of system dynamics. We discuss the efficient approximation of $\mathbb{G}_{S}$ respectively of its matrix representation $\mathbf{G}=\mathbf{M}_{\mathcal{Y}}^{-1} \mathbf{H}$. For time-invariant systems with distributed control and observation, this task reduces to the approximation of the convolution kernel $K \in L^{2}(0, T ; \mathscr{L}(U, Y))$.
4.1. Kernel function approximation. Inserting (2.4) in (3.3), by a change of variables we obtain

$$
\mathbf{H}_{i j}^{k l}=\int_{0}^{T} \int_{0}^{T} \psi_{i}(t) \phi_{j}(s)\left(\nu_{k}, K(t-s) \mu_{l}\right)_{Y} d s d t=\int_{0}^{T} \mathbf{W}_{i j}(t) \mathbf{K}_{k l}(t) d t
$$

with matrix-valued functions $\mathbf{W}:[0, T] \rightarrow \mathbb{R}^{s \times r}$ and $\mathbf{K}:[0, T] \rightarrow \mathbb{R}^{q \times p}$,

$$
\mathbf{W}_{i j}(t)=\int_{0}^{T-t} \psi_{i}(t+s) \phi_{j}(s) d s, \quad \mathbf{K}_{k l}(t)=\left(\nu_{k}, K(t) \mu_{l}\right)_{Y}
$$

and thus

$$
\begin{equation*}
\mathbf{H}=\mathbf{M}_{\mathcal{Y}} \mathbf{G}=\int_{0}^{T} \mathbf{K}(t) \otimes \mathbf{W}(t) d t \tag{4.2}
\end{equation*}
$$

Remark 7. $\mathbf{W}(t)$ can be exactly calculated if piecewise polynomial ansatz functions $\psi_{i}(t)$ and $\phi_{j}(t)$ are chosen. For the special choice (3.7), we see in this way that $\mathbf{W}(t) \in \mathbb{R}^{r \times r}$ is a lower triangular Toeplitz matrix for all $t \in[0, T]$, and hence the matrices $\mathbf{H}_{i j}=\int_{0}^{T} \mathbf{W}_{i j}(t) \mathbf{K}(t) d t \in \mathbb{R}^{q \times p}$ satisfy $\mathbf{H}_{i j}=\mathbf{H}_{i-j}$ for $1 \leq i, j \leq r$ and $\mathbf{H}_{i j}=0$ for $1 \leq i<j \leq r$. In view of Section 3.3, the $\mathbf{H}_{i j}$ are thus the Markov parameters of a discrete-time linear time-invariant causal MIMO system.

For systems of the form (2.1), the matrix-valued function $\mathbf{K}$ is given by

$$
\mathbf{K}_{k l}(t)=\left(\nu_{k}, C S(t) B \mu_{l}\right)_{Y}=\left(c_{k}^{*}, S(t) b_{l}\right)_{Z}
$$

where $c_{k}^{*}=C^{*} \nu_{k} \in Z$ and $b_{l}=B \mu_{l}$ for $k=1, \ldots, q$ and $l=1, \ldots, p$. Hence, $\mathbf{K}$ can be calculated by solving $p$ homogeneous systems

$$
\begin{align*}
& \dot{z}_{l}(t)=A z_{l}(t), \quad t \in(0, T],  \tag{4.3a}\\
& z_{l}(0)=b_{l} \tag{4.3b}
\end{align*}
$$

since (4.3) has the mild solution $z_{l}(t)=S(t) b_{l} \in C\left([0, T] ; L^{2}(\Omega)\right)$. We obtain an approximation $\tilde{\mathbf{H}}$ of $\mathbf{H}$ by replacing $z_{l}(t)$ by numerical approximations $z_{l, \text { tol }}(t)$, i.e.

$$
\begin{equation*}
\tilde{\mathbf{H}}=\int_{0}^{T} \tilde{\mathbf{K}}(t) \otimes \mathbf{W}(t) d t \tag{4.4}
\end{equation*}
$$

with $\tilde{\mathbf{K}}_{k l}(t)=\left(\nu_{k}, C z_{l, \text { tol }}(t)\right)_{Y}=\left(c_{k}^{*}, z_{l, \text { tol }}(t)\right)_{Z}$. Here the subscript tol indicates that the error $z_{l}-z_{l, \text { tol }}$ is assumed to satisfy some tolerance criterion which will be specified later. The corresponding approximation $\mathbb{G}_{D S}$ of $\mathbb{G}_{S}$ is given by

$$
\begin{equation*}
\mathbb{G}_{D S}=\kappa_{\mathcal{Y}}^{-1} \tilde{\mathbf{G}} \kappa_{\mathcal{U}} \mathbb{P}_{\mathcal{U}}, \quad \text { with } \tilde{\mathbf{G}}=\mathbf{M}_{\mathcal{Y}}^{-1} \tilde{\mathbf{H}} \tag{4.5}
\end{equation*}
$$

and depends on $h_{1}, h_{2}, \tau_{1}, \tau_{2}$ and tol.
REMARK 8. The matrix function $\mathbf{K}$ is approximated column-wise. The kernel may also be calculated row-wise by solving an adjoint autonomous system, which may be preferable if $q<p$ or if the output approximation is successively improved by adding further basis functions $\nu_{q+1}, \nu_{q+2}, \ldots$.

Remark 9. The calculation of $\tilde{\mathbf{H}}$ can be parallelized in an obvious way by calculating the $p$ solutions $z_{l, \text { tol }}$ in parallel and we note that no state trajectories have to be stored. In general, the matrix $\tilde{\mathbf{H}}$ is not sparse, such that the memory requirements become significant if a high resolution of the signals in space and time is required, and the question of a data-sparse representation arises. Recalling Remark 7, the blocks $\tilde{\mathbf{H}}^{k l}$ are lower triangular Toeplitz matrices for the special choice of time basis funtions (3.7) and thus only $q \cdot p \cdot r$ elements have to be stored. Another approach to obtain data-sparse representations uses approximate factorizations $\check{\mathbf{K}}_{k l}(t-s)=\sum_{m, n=1}^{M} \alpha_{m n} L_{m}(t) L_{n}(s)$ for $s, t \in[0, T]$ with suitable ansatz functions $L_{n}(t)$, see e.g. [2g].
4.2. Dynamics approximation error. The following proposition relates the system dynamics error $\epsilon_{D}$ to the errors made in solving the $\operatorname{PDE}(4.3)$ for $l=1, \ldots, p$.

Proposition 4.1. The system dynamics error $\epsilon_{D}:=\left\|\mathbb{G}_{S}-\mathbb{G}_{D S}\right\|_{\mathscr{L}(\mathcal{U}, \mathcal{Y})}$ satisfies

$$
\begin{equation*}
\epsilon_{D} \leq \sqrt{T}\|\mathbf{K}-\tilde{\mathbf{K}}\|_{L^{2}\left(0, T ; \mathbb{R}_{w}^{q \times p}\right)} \leq p \sqrt{T} \sqrt{\frac{\lambda_{\max }\left(\mathbf{M}_{Y, h_{2}}\right)}{\lambda_{\min }\left(\mathbf{M}_{U, h_{1}}\right)}} \max _{1 \leq l \leq p}\left\|\mathbf{K}_{:, l}-\tilde{\mathbf{K}}_{:, l}\right\|_{L^{2}\left(0, T ; \mathbb{R}^{q}\right)} . \tag{4.6}
\end{equation*}
$$

Here $\mathbf{K}_{:, l}$ and $\tilde{\mathbf{K}}_{:, l}$ denote the l'th column of $\mathbf{K}(t)$ and $\tilde{\mathbf{K}}(t)$, respectively, $\lambda_{\max }\left(\mathbf{M}_{Y, h_{2}}\right)$ is the largest eigenvalue of $\mathbf{M}_{Y, h_{2}}$ and $\lambda_{\min }\left(\mathbf{M}_{U, h_{1}}\right)$ the smallest eigenvalue of $\mathbf{M}_{U, h_{1}}$. $\mathbb{R}_{w}^{q \times p}$ denotes the space of real $q \times p$-matrices equipped with the weighted matrix norm $\|\mathbf{M}\|_{q \times p ; w}=\sup _{\mathbf{u} \neq 0}\|\mathbf{M} \mathbf{u}\|_{q ; w} /\|\mathbf{u}\|_{p ; w}$.

Proof. $\mathbf{K}$ is the matrix function representation of the space-projected kernel function $K_{m}:[-T, T] \rightarrow \mathscr{L}(U, Y)$ with $K_{m}(t)=P_{Y, h_{2}} K(t) P_{U, h_{1}}$, where $P_{Y, h_{2}}$ and $P_{U, h_{1}}$ are the orthogonal projections onto the subspaces $Y_{h_{2}}$ and $U_{h_{1}}$, respectively. Introducing the corresponding i/o-map $\mathbb{G}_{m}=\mathbb{G}_{m}\left(h_{1}, h_{2}\right)$,

$$
\begin{equation*}
\left(\mathbb{G}_{m} u\right)(t)=\int_{0}^{T} K_{m}(t-s) u(s) d s, \quad t \in[0, T] \tag{4.7}
\end{equation*}
$$

we note that $\mathbb{G}_{S}=\mathbb{P}_{\mathcal{Y}, h_{2}, \tau_{2}} \mathbb{G}_{m} \mathbb{P}_{\mathcal{U}, h_{1}, \tau_{1}}$. Similarly, we associate with $\tilde{\mathbf{K}}(t)$ the kernel function $\tilde{K}:[-T, T] \rightarrow \mathscr{L}(U, Y)$ with $\tilde{K}(t)=\kappa_{Y, h_{2}}^{-1} \tilde{\mathbf{K}}(t) \kappa_{U, h_{1}} P_{U, h_{1}}$, and with corresponding i/o-map

$$
\begin{equation*}
\left(\mathbb{G}_{D} u\right)(t)=\int_{0}^{T} \tilde{K}(t-s) u(s) d s, \quad t \in[0, T] \tag{4.8}
\end{equation*}
$$

We observe that $\mathbb{G}_{D S}$ as defined in (4.5) satisfies $\mathbb{G}_{D S}=\mathbb{P}_{\mathcal{Y}, h_{2}, \tau_{2}} \mathbb{G}_{D} \mathbb{P}_{\mathcal{U}, h_{1}, \tau_{1}}$ by showing according to (3.2)-(4.2) that the matrix representation of $\mathbb{P}_{\mathcal{Y}, h_{2}, \tau_{2}} \mathbb{G}_{D} \mathbb{P}_{\mathcal{U}, h_{1}, \tau_{1}}$ coincides with (4.4). We note that $\left\|K_{m}(t)\right\|_{\mathscr{L}(U, Y)}=\|\mathbf{K}(t)\|_{q \times p ; w}$ and $\|\tilde{K}(t)\|_{\mathscr{L}(U, Y)}=$ $\|\tilde{\mathbf{K}}(t)\|_{q \times p ; w}$ for all $t \in[0, T]$. Lemma 2.1 yields

$$
\left\|\mathbb{G}_{m}-\mathbb{G}_{D}\right\|_{\mathscr{L}(\mathcal{U}, \mathcal{Y})} \leq \sqrt{T}\left\|K_{m}-\tilde{K}\right\|_{L^{2}(0, T ; \mathscr{L}(U, Y))}=\sqrt{T}\|\mathbf{K}-\tilde{\mathbf{K}}\|_{L^{2}\left(0, T ; \mathbb{R}_{w}^{q \times p}\right)}
$$

Defining $\mathbf{E}(t)=\mathbf{K}(t)-\tilde{\mathbf{K}}(t)$, for $\mathbf{u} \in \mathbb{R}^{p}$ with $\|\mathbf{u}\|_{\mathbb{R}^{p}}=1$ and $t \in[0, T]$, by using the equivalence of the 1 -norm and 2 -vector norms in $\mathbb{R}^{p}$ we have that

$$
\|\mathbf{E}(t) \mathbf{u}\|_{\mathbb{R}^{q}} \leq \sum_{l=1}^{p} \mid \mathbf{u}_{l}\| \| \mathbf{E}_{:, l}(t) \|_{\mathbb{R}^{q}} \leq \sqrt{p}\left(\sum_{l=1}^{p}\left\|\mathbf{E}_{:, l}(t)\right\|_{\mathbb{R}^{q}}^{2}\right)^{1 / 2}
$$

and hence $\|\mathbf{E}\|_{L^{2}\left(0, T ; \mathbb{R}^{q \times p}\right)}^{2} \leq p \sum_{l=1}^{p} \int_{0}^{T}\left\|\mathbf{E}_{:, l}(t)\right\|_{\mathbb{R}^{q}}^{2} d t \leq p^{2} \max _{l=1, \ldots, p} \int_{0}^{T}\left\|\mathbf{E}_{:, l}(t)\right\|_{\mathbb{R}^{q}}^{2} d t$, which concludes the proof.

Remark 10. Calculating directly the columns of $\mathbf{K}$ and estimating $\epsilon_{D}$ via (4.6), the quotient of the eigenvalues of the mass matrices $\mathbf{M}_{U, h_{1}}$ and $\mathbf{M}_{Y, h_{2}}$ has to be compensated by the approximation accuracy of $\mathbf{K}_{:, l}$. This may be problematic if hierarchical basis functions are chosen, since the quotient grows unboundedly with decreasing $h_{1}$ and $h_{2}$. One may circumvent this problem by calculating $\mathbf{K}$ with respect to different bases. Approximating the columns of $\mathbf{K}^{w}(t)=\mathbf{M}_{\mathcal{Y}}^{1 / 2} \mathbf{K}(t) \mathbf{M}_{\mathcal{U}}^{-1 / 2}$ via an adapted problem (4.3), we have $\epsilon_{D} \leq p \sqrt{T} \max _{1 \leq l \leq p}\left\|\mathbf{K}_{:, l}^{w}-\tilde{\mathbf{K}}_{;, l}^{w}\right\|_{L^{2}\left(0, T ; \mathbb{R}^{q}\right)}$. Note that the necessary back transformations have to be carried out with sufficient accuracy.
4.3. Error estimation for the homogeneous PDE. In order to approximate the system dynamics, the homogeneous PDE (4.3) has to be solved via a fully-discrete numerical scheme for $p$ different initial values. A first goal is to choose the time and space grids (and possibly other discretization parameters) such that

$$
\begin{equation*}
\left\|\mathbf{K}_{:, l}-\tilde{\mathbf{K}}_{:, l}\right\|_{L^{2}\left(0, T ; \mathbb{R}^{q}\right)}<\text { tol } \quad \text { resp. } \quad\left\|\mathbf{K}_{:, l}^{w}-\tilde{\mathbf{K}}_{:, l}^{w}\right\|_{L^{2}\left(0, T ; \mathbb{R}^{q}\right)}<\text { tol } \tag{4.9}
\end{equation*}
$$

is guaranteed for a given tol $>0$ by means of reliable error estimators. A second goal is to achieve this accuracy in a cost-economic way. A special difficulty in solving (4.3) numerically is the handling of initial values $b_{l}$, which belong in general only to $Z$ but not necessarily to $D(A)$. Considering the example heat equation, this means that the space and time derivatives of the exact solution $z_{l} \in C^{1}\left((0, T], H^{2}(\Omega) \cap H_{0}^{1}(\Omega)\right)$ may become very large for small $t$, but decay quickly for $t>0$. In fact, in general we only have the analytical bound

$$
\left\|\partial_{t} z(t)\right\|_{L^{2}(\Omega)}=\|\Delta z(t)\|_{L^{2}(\Omega)} \leq \frac{c}{t}\left\|z^{0}\right\|_{L^{2}(\Omega)} \quad \text { for all } t \in(0, T]
$$

with some constant $c>0$ independent of $z_{0}$ and $T$, cf. [32, p. 148]. Adaptive space and time discretizations on the basis of a posteriori error estimates are the method of choice to match these requirements [23].

Discontinuous Galerkin time discretizations in combination with standard Galerkin space discretizations provide an appropriate framework to derive corresponding (a priori and a posteriori) error estimates, also for the case of adaptively refined grids which are in general no longer quasi-uniform [24, 32, 45]. We distinguish two types of error estimates.

Global state error estimates measure the error $\left(z_{l}-z_{l, \text { tol }}\right)$ in some global norm. For parabolic problems, a priori and a posteriori estimates for the error in $L^{\infty}\left(0, T ; L^{2}(\Omega)\right)$ and $L^{\infty}\left(0, T ; L^{\infty}(\Omega)\right)$ can be found in [24]. Such results permit to guarantee (4.9) in view of

$$
\begin{equation*}
\left\|\mathbf{K}_{:, l}-\tilde{\mathbf{K}}_{:, l}\right\|_{L^{2}\left(0, T ; \mathbb{R}^{q}\right)} \leq\|C\|_{\mathscr{L}(Z, Y)}\left(\sum_{i=1}^{q}\left\|\nu_{i}\right\|_{Y}^{2}\right)^{1 / 2}\left\|z-z_{\mathrm{tol}}^{(l)}\right\|_{L^{2}(0, T ; Z)} \tag{4.10}
\end{equation*}
$$

Goal-oriented error estimates can be used to measure the error $\left\|\mathbf{K}_{:, l}-\tilde{\mathbf{K}}_{:, l}\right\|_{L^{2}\left(0, T ; \mathbb{R}^{q}\right)}$ directly. This may be advantageous, since (4.10) may be very conservative: the error in the observations $\mathbf{K}_{:, l}$ can be small even if some norm of the state error is large. The core of these error estimation techniques is an exact error representation formula, which can be evaluated if one knows the residual and the solution of an auxiliary dual PDE. This so-called dual-weighted residuals (DWR) approach goes back to [4], substantial contributions have since then been made in [1], [6], [8], [9], [10], [31], [33], [34] and the references therein. Note that the numerical examples presented later are calculated using adaptive time and space grid refinements on the basis of DWR error estimation techniques.

The previous discussion justifies the following assumption.
Assumption 1. Given a tolerance tol $>0$, we can ensure (by using appropriate error estimators and mesh refinements) that the solutions $z_{l}$ of (4.3) and the solutions $z_{l, \text { tol }}$ calculated by means of an appropriate fully-discrete numerical scheme satisfy

$$
\begin{equation*}
\left\|\mathbf{K}_{:, l}-\tilde{\mathbf{K}}_{:, l}\right\|_{L^{2}\left(0, T ; \mathbb{R}^{q}\right)}<\mathrm{tol}, \quad l=1, \ldots, p \tag{4.11}
\end{equation*}
$$

5. Total error estimates. We present estimates for the total error in the approximation of $\mathbb{G}$ and of its adjoint $\mathbb{G}^{*}$. Using general-purpose ansatz spaces $\mathcal{U}_{h_{1}, \tau_{1}}$ and $\mathcal{Y}_{h_{2}, \tau_{2}}$ for the signal approximation, we only obtain error results in a weaker $\mathscr{L}\left(\mathcal{U}_{s}, \mathcal{Y}\right)$-norm respectively $\mathscr{L}\left(\mathcal{Y}_{s}, \mathcal{U}\right)$-norm.

Theorem 5.1. Consider the i/o map $\mathbb{G} \in \mathscr{L}(\mathcal{U}, \mathcal{Y})$ of the infinite-dimensional linear time-invariant system (2.4) and assume that
(i) $\mathbb{G}_{\mid \mathcal{U}_{s}} \in \mathscr{L}\left(\mathcal{U}_{s}, \mathcal{Y}_{s}\right)$ with spaces of higher regularity in space and time

$$
\mathcal{U}_{s}=H^{\alpha_{1}, \beta_{1}}((0, T) \times \Theta), \quad \mathcal{Y}_{s}=H^{\alpha_{2}, \beta_{2}}((0, T) \times \Xi), \quad \alpha_{1}, \beta_{1}, \alpha_{2}, \beta_{2} \in \mathbb{N} .
$$

(ii) The families of subspaces $\left\{\mathcal{U}_{h_{1}, \tau_{1}}\right\}_{h_{1}, \tau_{1}}$ and $\left\{\mathcal{Y}_{h_{2}, \tau_{2}}\right\}_{h_{2}, \tau_{2}}$ satisfy

$$
\begin{aligned}
\left\|u-\mathbb{P}_{\mathcal{U}, h_{1}, \tau_{1}} u\right\|_{\mathcal{U}} \leq\left(c_{\mathcal{R}} \tau_{1}^{\alpha_{1}}+c_{U} h_{1}^{\beta_{1}}\right)\|u\|_{\mathcal{U}_{s}}, & u \in \mathcal{U}_{s}, \\
\left\|y-\mathbb{P}_{\mathcal{Y}, h_{2}, \tau_{2}} y\right\|_{\mathcal{Y}} \leq\left(c_{\mathcal{S}} \tau_{2}^{\alpha_{2}}+c_{Y} h_{2}^{\beta_{2}}\right)\|y\|_{\mathcal{Y}_{s}}, & y \in \mathcal{Y}_{s},
\end{aligned}
$$

with positive constants $c_{\mathcal{R}}, c_{\mathcal{S}}, c_{U}$ and $c_{Y}$.
(iii) The error in solving for the state dynamics can be made arbitrarily small, i.e. Assumption 1 holds.
Let $\delta>0$ be given. Then one can choose subspaces $\mathcal{U}_{h_{1}^{*}, \tau_{1}^{*}}$ and $\mathcal{Y}_{h_{2}^{*}, \tau_{2}^{*}}$ such that

$$
\begin{array}{ll}
\tau_{1}^{*}<\left(\frac{\delta}{8 c_{\mathcal{R}}\|\mathbb{G}\|_{\mathscr{L}(\mathcal{U}, \mathcal{Y})}}\right)^{1 / \alpha_{1}}, & h_{1}^{*}<\left(\frac{\delta}{8 c_{U}\|\mathbb{G}\|_{\mathscr{L}(\mathcal{U}, \mathcal{Y})}}\right)^{1 / \beta_{1}} \\
\tau_{2}^{*}<\left(\frac{\delta}{8 c_{\mathcal{S}}\|\mathbb{G}\|_{\mathscr{L}\left(\mathcal{U}_{s}, \mathcal{Y}_{s}\right)}}\right)^{1 / \alpha_{2}}, & h_{2}^{*}<\left(\frac{\delta}{8 c_{Y}\|\mathbb{G}\|_{\mathscr{L}\left(\mathcal{U}_{s}, \mathcal{Y}_{s}\right)}}\right)^{1 / \beta_{2}} \tag{5.1b}
\end{array}
$$

and one can solve numerically the PDEs (4.3) for $l=1, \ldots, p\left(h_{1}\right)$ such that one of the following conditions holds.
(i)

$$
\begin{equation*}
\left\|\mathbf{K}_{:, l}^{w}-\tilde{\mathbf{K}}_{:, l}^{w}\right\|_{L^{2}\left(0, T ; \mathbb{R}^{q}\right)}<\frac{\delta}{2 \sqrt{T} p\left(h_{1}^{*}\right)}, \tag{5.2a}
\end{equation*}
$$

(ii) $\left\|\mathbf{K}_{:, l}-\tilde{\mathbf{K}}_{:, l}\right\|_{L^{2}\left(0, T ; \mathbb{R}^{q}\right)}<\frac{\delta}{2 \sqrt{T} p\left(h_{1}^{*}\right)} \sqrt{\frac{\lambda_{\min }\left(\mathbf{M}_{U, h_{1}^{*}}\right)}{\lambda_{\max }\left(\mathbf{M}_{Y, h_{2}^{*}}\right)}}$,


In this case,

$$
\left\|\mathbb{G}-\mathbb{G}_{D S}\right\|_{\mathscr{L}\left(\mathcal{U}_{s}, \mathcal{Y}\right)}<\delta .
$$

Moreover, the signal error $\epsilon_{S}^{\prime}:=\left\|\mathbb{G}-\mathbb{G}_{S}\right\|_{\mathscr{L}\left(\mathcal{U}_{s}, \mathcal{Y}\right)}$ and the system dynamics error $\epsilon_{D}:=\left\|\mathbb{G}_{S}-\mathbb{G}_{D S}\right\|_{\mathscr{L}(\mathcal{U}, \mathcal{Y})}$ are balanced in the sense that $\epsilon_{S}^{\prime}, \epsilon_{D}<\delta / 2$.

Proof. For $u \in \mathcal{U}_{s}$, we have

$$
\begin{aligned}
\left\|\mathbb{G} u-\mathbb{G}_{S} u\right\|_{\mathcal{Y}} & \leq\left\|\mathbb{G} u-\mathbb{P}_{\mathcal{Y}, h_{2}, \tau_{2}} \mathbb{G} u\right\|_{\mathcal{Y}}+\left\|\mathbb{P}_{\mathcal{Y}, h_{2}, \tau_{2}} \mathbb{G} u-\mathbb{P}_{\mathcal{Y}, h_{2}, \tau_{2}} \mathbb{G P}_{\mathcal{U}, h_{1}, \tau_{1}} u\right\|_{\mathcal{Y}}, \\
& \leq\left(c_{\mathcal{S}} \tau_{2}^{\alpha_{2}}+c_{Y} h_{2}^{\beta_{2}}\right)\|\mathbb{G} u\|_{\mathcal{Y}_{s}}+\left(c_{\mathcal{R}} \tau_{1}^{\alpha_{1}}+c_{U} h_{1}^{\beta_{1}}\right)\left\|\mathbb{P}_{\mathcal{Y}}\right\|_{\mathscr{L}(\mathcal{Y})}\|\mathbb{G}\|_{\mathscr{L}(\mathcal{U}, \mathcal{Y})}\|u\|_{\mathcal{U}_{s}}, \\
& \leq\left\{\left(c_{\mathcal{S}} \tau_{2}^{\alpha_{2}}+c_{Y} h_{2}^{\beta_{2}}\right)\|\mathbb{G}\|_{\mathscr{L}\left(\mathcal{U}_{r s}, \mathcal{Y}_{s}\right)}+\left(c_{\mathcal{R}} \tau_{1}^{\alpha_{1}}+c_{U} h_{1}^{\beta_{1}}\right)\|\mathbb{G}\|_{\mathscr{L}(\mathcal{U}, \mathcal{Y})}\right\}\|u\|_{\mathcal{U}_{s}},
\end{aligned}
$$

and thus (5.1) ensures that $\epsilon_{S}^{\prime}=\left\|\mathbb{G}-\mathbb{G}_{S}\right\|_{\mathscr{L}\left(\mathcal{U}_{s}, \mathcal{Y}\right)}<\delta / 2$. Proposition 4.1 in combination with (5.2) and in view of (4.10) ensures that $\epsilon_{D}=\left\|\mathbb{G}_{S}-\mathbb{G}_{D S}\right\|_{\mathscr{L}(\mathcal{U}, \mathcal{Y})}<\delta / 2$, which concludes the proof.

Remark 11. Condition (i) holds for many systems of practical relevance, cf. Remark 2. Condition (ii) can be achieved by choosing $U_{h_{1}}, Y_{h_{2}}, \mathcal{R}_{\tau_{1}}$ and $\mathcal{S}_{\tau_{2}}$ for instance as spaces of piecewise polynomial functions of appropriate degrees, cf. (3.8) and refer e.g. to [16] for interpolation theory in Sobolev spaces in the case of more general settings.

Remark 12. Considering $\mathbb{G}$ in the $\mathscr{L}\left(\mathcal{U}_{s}, \mathcal{Y}\right)$-norm we restrict the input space to input signals of higher regularity in space and time. This is not a strong limitation if we think of $\mathcal{U}_{s}$ as the space wherein we search for an approximate optimal control, since exact optimal controls exhibit such higher regularity in many applications [46]. Note that we do not exclude the use of $L^{2}$-controls, in fact, we can consider the spaces $\mathcal{U}_{h_{1}, \tau_{1}}$ as spaces of controls which we are able to realize in technical applications.

In view of (2.9) and (3.9) we can apply Theorem 5.1 to Example 1.
Corollary 5.2. Consider the heat control system in Example 1 with $U=Y=$ $L^{2}(0,1)$ and let $\delta>0$ be given. We assume that $C_{\mid H^{2}(\Omega)} \in \mathscr{L}\left(H^{2}(\Omega), H^{2}(0,1)\right)$, i.e. $\mathbb{G}_{\mid \mathcal{U}_{s}} \in \mathscr{L}\left(\mathcal{U}_{s}, \mathcal{Y}_{s}\right)$ with $\left.\mathcal{U}_{s}=\mathcal{Y}_{s}=H^{1,2}((0, T) \times(0,1))\right)$. We choose $U_{h_{1}}$ and $Y_{h_{2}}$ as spaces of continuous piecewise linear functions with respect to equidistant grids on $[0,1]$, and we choose $\mathcal{R}_{\tau_{1}}$ and $\mathcal{S}_{\tau_{2}}$ as spaces of piecewise constant functions with respect to equidistant grids on $[0, T]$, with dimensions satisfying

$$
\begin{array}{ll}
p>2 \sqrt{\frac{\|\mathbb{G}\|_{\mathscr{L}(\mathcal{U}, \mathcal{Y})}}{\delta}}+1, & q>2 \sqrt{\frac{\|\mathbb{G}\|_{\mathscr{L}\left(\mathcal{U}_{s}, \mathcal{Y}_{s}\right)}}{\delta}}+1, \\
r>\frac{\sqrt{2} 4\|\mathbb{G}\|_{\mathscr{L}(\mathcal{U}, \mathcal{Y})}}{\delta}, & s>\frac{\sqrt{2} 4\|\mathbb{G}\|_{\mathscr{L}\left(\mathcal{U}_{s}, \mathcal{Y}_{s}\right)}}{\delta} . \tag{5.3b}
\end{array}
$$

If the homogeneous heat equations (4.3) are solved for $l=1, \ldots, p$ such that one of the conditions (i)-(iii) in (5.2) holds, then the i/o-maps $\mathbb{G}$ and $\mathbb{G}_{D S}$ restricted to $\mathcal{U}_{s}$ satisfy

$$
\left\|\mathbb{G}-\mathbb{G}_{D S}\right\|_{\mathscr{L}\left(\mathcal{U}_{s}, \mathcal{Y}\right)}<\delta .
$$

The next result shows that $\left(\mathbb{G}_{D S}\right)^{*} \in \mathscr{L}(\mathcal{Y}, \mathcal{U})$ automatically approximates the adjoint $\mathbb{G}^{*}$ with $\left\|\mathbb{G}^{*}-\left(\mathbb{G}_{D S}\right)^{*}\right\|_{\mathscr{L}\left(\mathcal{Y}_{s}, \mathcal{U}\right)}<c \delta$ with a $\mathbb{G}$-specific constant $c$, under the assumption that $\mathbb{G}_{\mid \mathcal{Y}_{s}}^{*} \in \mathscr{L}\left(\mathcal{Y}_{s}, \mathcal{U}_{s}\right)$. Note that $\mathbb{G}^{*} \in \mathscr{L}(\mathcal{Y}, \mathcal{U})$ is given by

$$
\left(\mathbb{G}^{*} y\right)(s)=\int_{0}^{T} K(s-t)^{*} y(t) d t
$$

Theorem 5.3. The adjoint $\left(\mathbb{G}_{D S}\right)^{*} \in \mathscr{L}(\mathcal{Y}, \mathcal{U})$ of $\mathbb{G}_{D S} \in \mathscr{L}(\mathcal{U}, \mathcal{Y})$ has the matrix representation

$$
\begin{equation*}
\tilde{\mathbf{G}}^{*}:=\mathbf{M}_{\mathcal{U}}^{-1} \tilde{\mathbf{G}}^{T} \mathbf{M}_{\mathcal{Y}}=\mathbf{M}_{\mathcal{U}}^{-1} \tilde{\mathbf{H}}^{T} \in \mathbb{R}^{p r \times q s} \tag{5.4}
\end{equation*}
$$

For a given $\delta>0$, assume that all conditions in Theorem 5.1 hold, ensuring that $\left\|\mathbb{G}-\mathbb{G}_{D S}\right\|_{\mathscr{L}\left(\mathcal{U}_{s}, \mathcal{Y}\right)}<\delta$. If, in addition, $\mathbb{G}_{\mid \mathcal{Y}_{s}}^{*} \in \mathscr{L}\left(\mathcal{Y}_{s}, \mathcal{U}_{s}\right)$, then

$$
\begin{equation*}
\left\|\mathbb{G}^{*}-\left(\mathbb{G}_{D S}\right)^{*}\right\|_{\mathscr{L}\left(\mathcal{Y}_{s}, \mathcal{U}\right)}<\delta\left(\frac{1}{2}+c_{*}\right) \tag{5.5}
\end{equation*}
$$

with $c_{*}=\frac{1}{4}\left(\left\|\mathbb{G}^{*}\right\|_{\mathscr{L}\left(\mathcal{Y}_{s}, \mathcal{U}_{s}\right)} /\|\mathbb{G}\|_{\mathscr{L}(\mathcal{Y}, \mathcal{U})}+\left\|\mathbb{G}^{*}\right\|_{\mathscr{L}(\mathcal{Y}, \mathcal{U})} /\|\mathbb{G}\|_{\mathscr{L}\left(\mathcal{Y}_{s}, \mathcal{U}_{s}\right)}\right)$.
Proof. We first observe that $\tilde{\mathbf{G}}^{*}$ is the adjoint of $\tilde{\mathbf{G}}: \mathbb{R}_{w}^{p r} \rightarrow \mathbb{R}_{w}^{q s}$, since

$$
(\tilde{\mathbf{G}} \mathbf{u}, \mathbf{y})_{q s ; w}=\mathbf{u}^{T} \tilde{\mathbf{G}}^{T} \mathbf{M}_{\mathcal{y}} \mathbf{y}=\left(\mathbf{u}, \mathbf{M}_{\mathcal{U}}^{-1} \tilde{\mathbf{G}}^{T} \mathbf{M}_{\mathcal{y}} \mathbf{y}\right)_{p r ; w}
$$

For $u \in \mathcal{U}$ and $y \in \mathcal{Y}$ and omitting the dependencies on $h_{1}, h_{2}, \tau_{1}, \tau_{2}$ and tol, we have

$$
\begin{aligned}
\left(\mathbb{G}_{D S} u, y\right)_{\mathcal{Y}} & =\left(\mathbb{P}_{\mathcal{Y}} \kappa_{\mathcal{Y}}^{-1} \tilde{\mathbf{G}} \kappa_{\mathcal{U}} \mathbb{P}_{\mathcal{U}} u, y\right)_{\mathcal{Y}}=\left(\tilde{\mathbf{G}} \kappa_{\mathcal{U}} \mathbb{P}_{\mathcal{U} u}, \kappa_{\mathcal{Y}} \mathbb{P}_{\mathcal{Y} y}\right)_{q s ; w}, \\
& =\left(\kappa_{\mathcal{U}} \mathbb{P}_{\mathcal{U}} u, \tilde{\mathbf{G}}^{*} \kappa_{\mathcal{y}} \mathbb{P}_{\mathcal{Y}} y\right)_{q s ; w}=\left(u, \mathbb{P}_{\mathcal{U}} \kappa_{\mathcal{U}}^{-1} \tilde{\mathbf{G}}^{*} \kappa \mathcal{P _ { \mathcal { Y } }} \mathbb{P}_{\mathcal{Y}},\right. \\
& =\left(u,\left(\mathbb{G}_{D S}\right)^{*} y\right)_{\mathcal{U}},
\end{aligned}
$$

where we have used that $\mathbb{P}_{\mathcal{U}}=\mathbb{P}_{\mathcal{U}}^{*}, \mathbb{P}_{\mathcal{Y}}=\mathbb{P}_{\mathcal{Y}}^{*}, \kappa_{\mathcal{U}}^{*}=\kappa_{\mathcal{U}}^{-1}$ and $\kappa_{\mathcal{Y}}^{*}=\kappa_{\mathcal{y}}^{-1}$. To show (5.5), we estimate

$$
\left\|\mathbb{G}^{*}-\left(\mathbb{G}_{D S}\right)^{*}\right\|_{\mathscr{L}\left(\mathcal{y}_{s}, \mathcal{u}\right)} \leq\left\|\mathbb{G}^{*}-\left(\mathbb{G}_{S}\right)^{*}\right\|_{\mathscr{L}\left(\mathcal{y}_{s}, \mathcal{u}\right)}+\left\|\mathbb{G}_{S}^{*}-\left(\mathbb{G}_{D S}\right)^{*}\right\|_{\mathscr{L}(y, \mathcal{U})},
$$

where $\left(\mathbb{G}_{S}\right)^{*}=\mathbb{P}_{\mathcal{U}} \mathbb{G}^{*} \mathbb{P}_{\mathcal{Y}}$ is the adjoint of $\mathbb{G}_{S} \in \mathscr{L}(\mathcal{U}, \mathcal{Y})$. In analogy to Thm. 5.1, one shows that

$$
\epsilon_{S}^{*}:=\left\|\mathbb{G}^{*}-\left(\mathbb{G}_{S}\right)^{*}\right\|_{\mathscr{L}\left(\mathcal{y}_{s}, \mathcal{U}\right)} \leq c_{\mathcal{R}}^{\prime \prime} \tau_{1}^{\alpha_{1}}+c_{U}^{\prime \prime} h_{1}^{\beta_{1}}+c_{S}^{\prime \prime} \tau_{2}^{\alpha_{2}}+c_{Y}^{\prime \prime} h_{2}^{\beta_{2}},
$$

with $c_{U}^{\prime \prime}=\left\|\mathbb{G}^{*}\right\|_{\mathscr{L}\left(\mathcal{Y}_{S}, \mathcal{u}_{S}\right)} c_{U}, c_{Y}^{\prime \prime}=\left\|\mathbb{G}^{*}\right\|_{\mathscr{L}(\mathcal{Y}, \mathcal{U})} c_{Y}, c_{\mathcal{R}}^{\prime \prime}=\left\|\mathbb{G}^{*}\right\|_{\mathscr{L}\left(\mathcal{Y}_{s}, \mathcal{U}_{s}\right)} c_{\mathcal{R}}$ and $c_{\mathcal{S}}^{\prime \prime}=$ $\left\|\mathbb{G}^{*}\right\| \mathscr{L}(y, \mathcal{U}) c_{\mathcal{S}}$. Hence, (5.1) implies

In order to estimate $\epsilon_{D}^{*}:=\left\|\mathbb{G}_{S}^{*}-\left(\mathbb{G}_{D S}\right)^{*}\right\|_{\mathscr{L}(\mathcal{Y}, \mathcal{U})}$, we recall the definition of $\mathbb{G}_{m}$ in (4.7) and of $\mathbb{G}_{D}$ in (4.8) and obtain

$$
\begin{align*}
\epsilon_{D}^{*} & \leq\left\|\mathbb{P}_{\mathcal{U}}\right\|_{\mathscr{L}(\mathcal{u})}\left\|\left(\mathbb{G}_{m}\right)^{*}-\left(\mathbb{G}_{D}\right)^{*}\right\|_{\mathscr{L}(\mathcal{Y}, \mathcal{U})}\left\|_{\mathcal{P}}^{\mathcal{V}}\right\|_{\mathscr{L}(\mathcal{Y})} \\
& \leq \sqrt{T}\left(\int_{0}^{T}\left\|K_{m}(t)^{*}-\tilde{K}(t)^{*}\right\|_{\mathscr{L}(Y, U)}^{2} d t\right)^{1 / 2} \tag{5.6}
\end{align*}
$$

We observe that

$$
\begin{equation*}
K_{m}(t)=P_{U} \kappa_{U}^{-1} \mathbf{K}(t)^{*} \kappa_{Y} P_{Y} \quad \text { and } \quad \tilde{K}(t)^{*} P_{U} \kappa_{U}^{-1} \tilde{\mathbf{K}} \kappa_{Y} P_{Y} \tag{5.7}
\end{equation*}
$$

with $\mathbf{K}(t)^{*}=\mathbf{M}_{U}^{-1} \mathbf{K}(t)^{T} \mathbf{M}_{Y}$ and $\tilde{\mathbf{K}}(t)^{*}=\mathbf{M}_{U}^{-1} \tilde{\mathbf{K}}(t)^{T} \mathbf{M}_{Y}$, and that

$$
\begin{equation*}
\left\|\mathbf{K}(t)^{*}-\tilde{\mathbf{K}}(t)^{*}\right\|_{\mathbb{R}_{w}^{p, q}}=\|\mathbf{K}(t)-\tilde{\mathbf{K}}(t)\|_{\mathbb{R}_{w}^{q \times p}} \quad \text { for } t \in[0, T] . \tag{5.8}
\end{equation*}
$$

Since each of the conditions in (5.2) ensures $\|\mathbf{K}-\tilde{\mathbf{K}}\|_{\mathbb{R}_{w}^{q \times p}}<\delta / 2 \sqrt{T}$, we have by means of (5.6) - (5.8) that $\epsilon_{D}^{*}<\delta / 2$, which concludes the proof. $\square$

Remark 13. It remains to investigate the accuracy of the respective approximation of the (possibly regularized) pseudo-inverses $\left(\mathbb{G}^{*} \mathbb{G}+\alpha \mathbb{I}\right)^{-1} \mathbb{G}^{*}$ with $\alpha \geq 0$, which play an important role e.g. in optimal control problems.

## 6. Applications and numerical results.

6.1. Test problems. As test cases, we consider two heat equations on domains $\Omega \subset \mathbb{R}^{2}$ as depicted in Fig. 6.1 and with control and observation operators of the following form. Let $\Omega_{c}=\left(a_{c, 1}, a_{c, 2}\right) \times\left(b_{c, 1}, b_{c, 2}\right)$ and $\Omega_{m}=\left(a_{m, 1}, a_{m, 2}\right) \times\left(b_{m, 1}, b_{m, 2}\right)$ be rectangular subsets of $\Omega$ where the control is active and the observation takes place, respectively, with 4 appropriate points $a_{c}, b_{c}, a_{m}, b_{m} \in \bar{\Omega}$. Setting $U=Y=L^{2}(0,1)$, we define $C \in \mathscr{L}\left(L^{2}(\Omega), Y\right)$ and $B \in \mathscr{L}\left(U, L^{2}(\Omega)\right)$ for test case 1 by
$(C z)(\xi)=\int_{a_{m, 1}}^{b_{m, 1}} \frac{z\left(x_{1}, x_{2}(\xi)\right)}{b_{m, 1}-a_{m, 1}} d x_{1}, \quad(B u)\left(x_{1}, x_{2}\right)=\left\{\begin{array}{ll}u\left(\theta\left(x_{1}\right)\right) \omega_{c}\left(x_{2}\right) & ,\left(x_{1}, x_{2}\right) \in \Omega_{c} \\ 0 & ,\left(x_{1}, x_{2}\right) \notin \Omega_{c}\end{array}\right.$,
where $\omega_{c} \in L^{2}\left(a_{c, 2}, b_{c, 2}\right)$ is a weight function and $\theta:\left[a_{c, 1}, b_{c, 1}\right] \rightarrow[0,1]$ and $x_{1}$ : $[0,1] \rightarrow\left[a_{m, 1}, b_{m, 1}\right]$ are affine-linear transformations. For test case 2, we just invert the roles of $x_{1}$ and $x_{2}$ in the definition of $C$. Note that $C$ preserves an inherent spatial state regularity, i.e. $C_{\mid H^{2}(\Omega)} \in \mathscr{L}\left(H^{2}(\Omega), H^{2}(0,1)\right)$.

Test case 1. We consider a heat equation with homogeneous Dirichlet boundary conditions on $(0, T] \times \Omega$ with $T=1$ and $\Omega=(0,1)^{2}$. We choose $\Omega_{c}=\Omega, \Omega_{m}=$ $(0.1,0.2) \times(0.1,0.9)$ and $\omega_{c}\left(x_{2}\right)=\sin \left(\pi x_{2}\right)$. In this case, the output to inputs of the special form $u(t ; \theta)=\sin \left(\omega_{T} \pi t\right) \sin (m \pi \theta)$ with $\omega_{T}, m \in \mathbb{N}$ can be explicitely formulated in terms of the eigenfunctions of the Laplace operator.

As test case 2 , we consider two infinitely long plates of width 5 and height 0.2 which are connected by two rectangular bars as shown in the cross section in Fig. 6.1. We assume that the plates are surrounded by an insulating material and that we can heat the bottom plate and measure the temperature distribution in the upper plate.

Test case 2. We consider a heat equation with homogeneous Neumann boundary conditions on $(0, T] \times \Omega$ with $T=1$ and $\Omega$ as in Fig. 6.1, and choose $\Omega_{c}=$ $(0.05,4.95) \times(0.05,0.15), \Omega_{m}=(0.05,4.95) \times(0.85,0.95)$ and $\omega_{c}\left(x_{2}\right)=\sin \left(\pi\left(x_{2}-\right.\right.$ 0.05)/0.1).

The matrix approximations $\tilde{\mathbf{G}}$ of the i/o-maps $\mathbb{G}$ corresponding to the test cases have been calculated by means of a heat equation solver, which is based on the $\mathrm{C}++$ FEM software library DEAL.II[5]. It realizes a discontinuous Galerkin scheme with adaptive space and time grids and applies goal-oriented DWR-based error control to ensure (4.9), see [35] for details.


Fig. 6.1. Test cases heat equation: (a) with homogeneous Dirichlet boundary conditions, (b) with homogeneous Neumann boundary conditions.
6.2. Tests of convergence. The following numerical convergence tests have all been carried out with approximations $\mathbb{G}_{D S}\left(h_{1}, \tau_{1}, h_{2}, \tau_{2}\right.$, tol) of the i/o-map $\mathbb{G}$ corresponding to Test case 1. Hierarchical linear finite elements in $U_{h_{1}}$ and $Y_{h_{2}}$ and Haar wavelets in $\mathcal{R}_{\tau_{1}}$ and $\mathcal{S}_{\tau_{2}}$ have been chosen. The tolerance tol refers to the estimate (4.11).

Convergence of single outputs. Considering Test case 1 with inputs $u(t ; \theta)=$ $\sin \left(\omega_{T} \pi t\right) \sin (m \pi \theta)$, and exactly known outputs $y=\mathbb{G} u$, we investigate the relative error $\|y-\tilde{y}\|_{\mathcal{Y}} /\|u\|_{\mathcal{U}_{s}}$, with $\tilde{y}=\mathbb{G}_{D S}\left(h_{1}, \tau_{1}, h_{2}, \tau_{2}\right.$, tol) $u$, for varying discretization parameters $h_{1}, \tau_{1}, h_{2}, \tau_{2}$ and tol. Choosing e.g. $m=5$ and $\omega_{T}=10$, we observe a quadratic convergence in $h_{1}=h_{2}$ (cf. Fig. 6.2-a) and a linear convergence in $\tau_{1}=\tau_{2}$ (cf. Fig. 6.2-b) in correspondence to Thm. 5.1. However, the error does not converge to zero but to a positive plateau value, which is due to the system dynamics error and which becomes smaller for lower tolerances tol. For input signals with $m>5$ and $\omega_{T}>10$ the convergence order can only be observed for smaller discretization parameters $h_{1}, h_{2}, \tau_{1}$ and $\tau_{2}$.


Fig. 6.2. (a) Relative output errors for input $u(t ; \theta)=\sin (10 \pi t) \sin (5 \pi \theta)$, varying $h_{1}=h_{2}$ and fixed $\tau_{1}=\tau_{2}=1 / 64$. (b) Relative output errors for input $u(t ; \theta)=\sin (10 \pi t) \sin (5 \pi \theta)$, varying $\tau_{1}=\tau_{2}$ and fixed $h_{1}=h_{2}=1 / 17$. (c) Norm $\left\|\mathbb{G}_{D S}(\mathbf{h})\right\|_{\mathscr{L}(\mathcal{U}, \mathcal{Y})}$ for synchronously increasing approximation space dimensions $p=q=r+1=s+1$ and fixed tolerance tol $=4.0 e-5$.

Convergence of the norm $\left\|\mathbb{G}_{S}\left(h_{1}, \tau_{1}, h_{2}, \tau_{2}\right)\right\|_{\mathscr{L}(\mathcal{U}, \mathcal{Y})}$ for nested subspaces. Successively improving the signal approximation by adding additional basis functions, the norm $\left\|\mathbb{G}_{S}\left(h_{1}, \tau_{1}, h_{2}, \tau_{2}\right)\right\|_{\mathscr{L}(\mathcal{U}, \mathcal{Y})}$ converges, cf. Lemma 3.1. We approximate $\left\|\mathbb{G}_{S}\right\|_{\mathscr{L}(\mathcal{U}, \mathcal{Y})}$ by $\left\|\mathbb{G}_{D S}\right\|_{\mathscr{L}(\mathcal{U}, \mathcal{Y})}$, where $\mathbb{G}_{D S}$ has been calculated with tol $=4.0 e-5$. In Fig. 6.2-c, the approximations $\left\|\mathbb{G}_{S}\left(h_{1}, \tau_{1}, h_{2}, \tau_{2}\right)\right\|_{\mathscr{L}(\mathcal{U}, \mathcal{Y})}=\left\|\mathbb{G}_{S}\left(\frac{1}{p-1}, \frac{1}{r}, \frac{1}{q-1}, \frac{1}{s}\right)\right\|_{\mathscr{L}(\mathcal{U}, \mathcal{Y})}$ are plotted for increasing subspace dimensions $p=q=r+1=s+1=2,3, \ldots, 65$.
6.3. Matrix reduction on the basis of SVDs. In order to resolve the input and output signal spaces accurately by means of general purpose basis functions, a large number of basis functions is needed in general. In order to reduce the large size of the resulting i/o-matrices $\tilde{\mathbf{G}}$, we apply a reduction method known as Tucker decomposition or higher order singular value decomposition (HOSVD) [19]. It is based on singular value decompositions (SVDs) and preserves the space-time tensor structure of the input and output signal bases.

Considering $\tilde{\mathbf{G}} \in \mathbb{R}^{q s \times p r}$ as a fourth-order tensor $\tilde{\mathbf{G}} \in \mathbb{R}^{s \times r \times q \times p}$ with $\tilde{\mathbf{G}}_{i j k l}=\tilde{\mathbf{G}}_{i j}^{k l}$,
it is shown in [19] that there exists a HOSVD

$$
\begin{equation*}
\tilde{\mathbf{G}}=\mathbf{S} \times{ }_{1} \mathbf{U}^{(\psi)} \times_{2} \mathbf{U}^{(\phi)} \times_{3} \mathbf{U}^{(\nu)} \times_{4} \mathbf{U}^{(\mu)} \tag{6.1}
\end{equation*}
$$

Here $\mathbf{S} \in \mathbb{R}^{s \times r \times q \times p}$ is a so-called core tensor, satisfying some orthogonality properties, $\mathbf{U}^{(\psi)} \in \mathbb{R}^{s \times s}, \mathbf{U}^{(\phi)} \in \mathbb{R}^{r \times r}, \mathbf{U}^{(\nu)} \in \mathbb{R}^{q \times q}, \mathbf{U}^{(\mu)} \in \mathbb{R}^{p \times p}$ are unitary matrices and $\times_{1}, \ldots, \times_{4}$ denote tensor-matrix multiplications. We define a so-called matrix unfolding $\tilde{\mathbf{G}}^{(\psi)} \in \mathbb{R}^{s \times r q p}$ of the tensor $\tilde{\mathbf{G}}$ by

$$
\tilde{\mathbf{G}}_{i m}^{(\psi)}=\mathbf{G}_{i j k l}, \quad m=(k-1) p s+(l-1) s+i
$$

i.e. we put all elements belonging to $\psi_{1}, \psi_{2}, \ldots, \psi_{s}$ into one respective row, and we define the unfoldings $\tilde{\mathbf{G}}^{(\phi)} \in \mathbb{R}^{r \times q p s}, \tilde{\mathbf{G}}^{(\nu)} \in \mathbb{R}^{q \times p s r}$ and $\tilde{\mathbf{G}}^{(\mu)} \in \mathbb{R}^{p \times s r q}$ in a similar cyclic way. Then, $\mathbf{U}^{(\psi)}, \mathbf{U}^{(\phi)}, \mathbf{U}^{(\nu)}$ and $\mathbf{U}^{(\mu)}$ in (6.1) can be calculated by means of four SVDs of the respective form

$$
\tilde{\mathbf{G}}^{(\psi)}=\mathbf{U}^{(\psi)} \Sigma^{(\psi)}\left(\mathbf{V}^{(\psi)}\right)^{T}
$$

where $\Sigma^{(\psi)}$ is diagonal with entries $\sigma_{1}^{(\psi)} \geq \sigma_{2}^{(\psi)} \geq \ldots \sigma_{s}^{(\psi)} \geq 0$ and $\mathbf{V}^{(\psi)}$ is columnwise orthonormal. The $\sigma_{i}^{(\psi)}$ are so-called n-mode singular values (or in our case $\psi$-mode singular values) of the tensor $\tilde{\mathbf{G}}$ and correspond to the Frobenius norms of certain subtensors of the core tensor $\mathbf{S}$.

On the basis of (6.1) we can define an approximation $\hat{\mathbf{G}} \in \mathbb{R}^{s \times r \times q \times p}$ of $\tilde{\mathbf{G}}$ by discarding the smallest $n$-mode singular values $\left\{\sigma_{\hat{s}+1}^{(\psi)}, \ldots, \sigma_{s}^{(\psi)}\right\},\left\{\sigma_{\hat{r}+1}^{(\phi)}, \ldots, \sigma_{r}^{(\psi)}\right\}$, $\left\{\sigma_{\hat{q}+1}^{(\nu)}, \ldots, \sigma_{q}^{(\nu)}\right\}$ and $\left\{\sigma_{\hat{p}+1}^{(\mu)}, \ldots, \sigma_{p}^{(\mu)}\right\}$, i.e. we set the corresponding parts of $\mathbf{S}$ to zero. Then we have

$$
\|\tilde{\mathbf{G}}-\hat{\mathbf{G}}\|_{F}^{2} \leq \sum_{i=\hat{s}+1}^{s} \sigma_{i}^{(\psi)}+\sum_{j=\hat{r}+1}^{r} \sigma_{j}^{(\phi)}+\sum_{k=\hat{q}+1}^{q} \sigma_{k}^{(\nu)}+\sum_{l=\hat{p}+1}^{p} \sigma_{l}^{(\mu)}
$$

see [19]. Note that, in contrast to matrix SVDs, this approximation is not necessarily optimal in a least square sense. For best-rank approximations, see e.g. [20].

The truncation of $\hat{\mathbf{G}} \in \mathbb{R}^{q r \times p s}$ after a basis transformation corresponding to $\mathbf{U}^{(\psi)}$, $\mathbf{U}^{(\phi)}, \mathbf{U}^{(\nu)}$ and $\mathbf{U}^{(\mu)}$ yields a low-dimensional representation $\overline{\mathbf{G}} \in \mathbb{R}^{\hat{q} \hat{r} \times \hat{p} \hat{s}}$.

In Figure 6.3 the HOSVD has been applied to a matrix $\tilde{\mathbf{G}} \in \mathbb{R}^{q s \times p r}$ for the Test case 2 with $p=17, q=65$ and $r=s=64$. The first row shows the respective $n$-mode singular values. Underneath the first and most relevant two transformed/new basis functions $\hat{\mu}_{i}, \hat{\nu}_{i}, \hat{\phi}_{i}$ and $\hat{\psi}_{i}$, are plotted. It is not surprising that the positions of the connections between the plates can be recovered as large values of the corresponding spatial input and output basis functions.

Remark 14. The application of a HOSVD can be useful in two ways. First, in order to obtain a low-dimensional matrix-representation of the system, which is small enough to be used for real-time feedback control design. Second, in order to identify relevant input and output signals, which may be instructive for actuator and sensor design, i.e. they might help to answer where actuators and sensors have to be placed and which resolution in time and space they should have.
6.4. Application in optimization problems. We investigate the use of the i/o-map approximation in optimization problems

$$
\begin{equation*}
\min J(u, y) \quad \text { subject to } y=\mathbb{G} u, \quad u \in \mathcal{U}_{a d} \tag{6.2}
\end{equation*}
$$



Fig. 6.3. HOSVD applied to the i/o map of Test case 2. First row: n-mode singular values in semilogarithmic scales. 2nd and 3rd row: Respective two most relevant basis functions.

Here $\mathcal{U}_{a d} \subset \mathcal{U}$ is the subset of admissible controls, $J: \mathcal{U} \times \mathcal{Y} \rightarrow \mathbb{R}$ is the quadratic cost functional $J(u, y)=\frac{1}{2}\left\|y-y_{D}\right\|_{\mathcal{Y}}^{2}+\alpha\|u\|_{\mathcal{U}}^{2}, y_{D} \in \mathcal{Y}$ is an aspired system's output signal, and $\alpha>0$ is a regularization parameter. We define the discrete cost functional

$$
\begin{equation*}
\bar{J}_{\mathbf{h}}: \mathbb{R}^{p r} \times \mathbb{R}^{q s} \rightarrow \mathbb{R}, \quad \bar{J}_{\mathbf{h}}(\mathbf{u}, \mathbf{y})=\frac{1}{2}\left\|\mathbf{y}-\mathbf{y}_{D}\right\|_{q s ; w}^{2}+\alpha\|\mathbf{u}\|_{p r ; w}^{2} \tag{6.3}
\end{equation*}
$$

with $\mathbf{y}_{D}=\kappa \mathcal{Y}, h_{2}, \tau_{2} \mathbb{P}_{\mathcal{Y}, h_{2}, \tau_{2}} y_{D}$, and instead of (6.2) we solve

$$
\begin{equation*}
\min \bar{J}_{\mathbf{h}}(\mathbf{u}, \mathbf{y}) \quad \text { subject to } \mathbf{y}=\tilde{\mathbf{G}} \mathbf{u}, \quad \mathbf{u} \in \bar{U}_{a d} \tag{6.4}
\end{equation*}
$$

with $\bar{U}_{a d}=\left\{\mathbf{u} \in \mathbb{R}^{p r}: \mathbf{u}=\kappa_{\mathcal{U}, h_{1}, \tau_{1}} \mathbb{P}_{\mathcal{U}, h_{1}, \tau_{1}} u, u \in \mathcal{U}_{a d}\right\}$. Considering optimization problems without control constraints, i.e. $\mathcal{U}_{a d}=\mathcal{U}$ and $\bar{U}_{a d}=\mathbb{R}^{p r}$, the solution $\overline{\mathbf{u}}$ of (6.4) is characterized by

$$
\begin{equation*}
\left(\tilde{\mathbf{G}}^{T} \mathbf{M}_{\mathcal{Y}} \tilde{\mathbf{G}}+\alpha \mathbf{M}_{\mathcal{U}}\right) \overline{\mathbf{u}}=\tilde{\mathbf{G}}^{T} \mathbf{M}_{\mathcal{Y} \mathbf{y}_{D} .} \tag{6.5}
\end{equation*}
$$

As concrete example, we consider Test case 2 and choose $y_{D}=\mathbb{G} u_{0}$ to be the output for an input $u_{0} \equiv 1$ which is equal to 1 on all of $[0, T] \times(0,1)$. We then try to find an optimized input $u_{*}$ of less energy, such that $\mathbb{G} u_{*} \approx y_{D}$, or more exactly, $u_{*}$ that minimizes the cost functional (6.2).

First we solve (6.5) with an approximated i/o map $\tilde{\mathbf{G}} \in \mathbb{R}^{17.64 \times 65.64}$ and $\alpha=10^{-4}$, yielding an approximation $\bar{u} \approx u_{*}$. The solution takes 0.33 seconds on a normal desktop PC. The $u$-norm is reduced by $27.9 \%$ and the relative deviation of $\mathbb{G} \bar{u}$ from $y_{D}$ is $9.4 \%$. In Fig. 6.4 the same calculations have been carried out with $\hat{\mathbf{G}} \in \mathbb{R}^{3.5 \times 3.5}$, where $\hat{\mathbf{G}}$ arises from a HOSVD-based matrix reduction of $\tilde{\mathbf{G}} \in \mathbb{R}^{17.64 \times 65 \cdot 64}$, where all but the 3 most relevant spatial and the 5 most relevant temporal input and output basis functions have been truncated. Using this approximation, the norm of $u$ is reduced by $27.4 \%$, whereas the relative deviation of $\mathbb{G} \bar{u}$ from $y_{D}$ is $9.5 \%$. The cost functional has been reduced by $44.5 \%$, and the calculation of $\bar{u}$ took less than 0.0004 seconds. The outputs resulting from $u_{0}$ and $\bar{u}$ have been calculated in simulations independent from the calculation of the i/o-matrix.


Fig. 6.4. Application of the SVD-reduced approximated $i / o$ map $\hat{\mathbf{G}} \in \mathbb{R}^{3.5 \times 3.5}$ in an optimization problem. From left to right: optimized control $\bar{u}$, original output $y_{D}=\mathbb{G} u_{0}$, optimized output $\mathbb{G} \bar{u}$ and their difference.

Remark 15. Aiming to realize the optimal control via a feedback control, we may look for feedback operators $\mathbb{F} \in \mathscr{L}(\mathcal{Y}, \mathcal{U})$ satisfying $u_{*}=\mathbb{F} y_{*}$ for $y_{*}=\mathbb{G} u_{*}$. On the discrete level, we search for a matrix $\mathbf{F} \in \mathbb{R}^{p r \times q \text { s }}$ satisfying

$$
\begin{equation*}
\overline{\mathbf{u}}=\mathbf{F} \overline{\mathbf{y}}, \quad \text { where } \overline{\mathbf{y}}=\tilde{\mathbf{G}} \overline{\mathbf{u}} \tag{6.6}
\end{equation*}
$$

and $\overline{\mathbf{u}}$ is the solution of (6.4). We note that for $\overline{\mathbf{y}} \neq 0$, such a matrix $\mathbf{F}$ always exists. Requiring however, that $\mathbf{F}$ has extra structure like time-invariance and causality, (6.6) may only be solvable in a least square sense.
7. Final remarks and outlook. We have presented a systematic framework for the discretization of i/o-maps of linear infinite-dimensional control systems with spatially distributed inputs and outputs. Global error estimates have been provided, which allow to choose the involved discretization parameters in such a way that a desired overall accuracy is achieved and that the signal and the system dynamics approximation errors are balanced. Moreover, the error results are capable to take many practical and technical restrictions in sensor and actuator design like limited spatial and temporal resolutions or the use of piece-wise constant controls and observations due to digital devices directly into account.

The numerical costs of the approach are primarily governed by the numerical calculation of $p$ underlying homogeneous PDEs, where $p$ is the number of input basis functions in space, which can become problematic when the spatial resolution of the input signal space has to be accurate. In this case, code-optimization, e.g. due to parallelization and appropriate updating of mass and stiffness matrices from prior calculations, promises to have a large potential for speed-up which has not yet been investigated.

The SVD-based dimension reduction for the matrix representation can be considered as an alternative model reduction approach, and the resulting reduced i/o-models prove to be useful in first numerical optimization applications. Moreover, the SVDbased reduction may be able to provide useful insight for efficient actuator and sensor design by filtering out relevant input and output signals.

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