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## Report

## Author(s):

Kunoth, Angela; Schwab, Christoph

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# Analytic regularity and GPC approximation for control problems constrained by linear parametric elliptic and parabolic PDEs* 

A. Kunoth ${ }^{\dagger}$ and Ch. Schwab

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# ANALYTIC REGULARITY AND GPC APPROXIMATION FOR CONTROL PROBLEMS CONSTRAINED BY LINEAR PARAMETRIC ELLIPTIC AND PARABOLIC PDES* 

ANGELA KUNOTH ${ }^{\dagger}$ AND CHRISTOPH SCHWAB ${ }^{\ddagger}$


#### Abstract

This paper deals with linear-quadratic optimal control problems constrained by a parametric or stochastic elliptic or parabolic PDE. We address the (difficult) case that the number of parameters may be countable infinite, i.e., $\sigma_{j}$ with $j \in \mathbb{N}$, and that the PDE operator may depend non-affinely on the parameters. We consider tracking-type functionals and distributed as well as boundary controls. Building on recent results in [CDS1, CDS2], we show that the state and the control are analytic as functions depending on these parameters $\sigma_{j}$. Polynomial approximations of state and control in terms of the possibly countably many stochastic coordinates $\sigma_{j}$ will be used to establish sparsity of polynomial "generalized polynomial chaos (gpc)" expansions of the state and the control with respect to the parameter sequence $\left(\sigma_{j}\right)_{j \geq 1}$. These imply, in particular, convergence rates of best $N$-term truncations of these expansions. The sparsity result allows in conjunction with adaptive wavelet Galerkin schemes as in [SG11, G] for sparse, adaptive tensor discretizations of control problems constrained by linear elliptic and parabolic PDEs developed in [DK, GK, K].


Key words. Linear-quadratic optimal control, linear parametric or stochastic PDE, distributed or boundary control, elliptic or parabolic PDE, analyticity, polynomial chaos approximation.

AMS subject classifications. 41A, 65K10, 65N99, 49N10, 65C30.

1. Introduction. Increasingly, simulation and design of complex systems requires the numerical solution of partial differential equations (PDEs) involving a large number of parameters. We mention only PDEs on high dimensional, so-called "design spaces". Also stochastic PDEs driven by noise lead to parametric PDEs when Wiener chaos expansions are employed to circumvent Monte-Carlo simulations. Of particular interest in this respect are optimal control problems of parametric systems that are governed by linear parametric or stochastic PDEs: in PDE-constrained control with a tracking-type optimization functional, the goal is to steer the solution $y$ of the PDE, called the state, towards a prescribed desired state in a least-squares sense while minimizing the effort for the control $u$. If, however, the PDE depends on (possibly countably many) parameters arising, for example, from random field inputs in models of uncertainty, this would require the solution of the control problem for each instance of the parameters. Already for a single random variable $\sigma$ in the diffusion coefficient, the computational expense would be enormous: each realization of this variable, e.g., in a Monte-Carlo simulation with $N$ draws, would require the solution of the whole control problem, resulting in necessarily $N$ solutions of the control problem.

For deterministic linear-quadratic control problems constrained by elliptic PDEs, one needs to solve as first order necessary and sufficient conditions for optimality a coupled system of linear PDEs for the state $y$ and the adjoint state $p$ each, and a

[^1]third equation coupling $p$ with the control $u$. For such systems of PDEs, in recent years solvers became available which produce optimal numerical approximations of the solution triple $(y, p, u)$, in the sense that accuracy versus work to obtain these approximations is provably proportional to those of best $N$-term approximations of the solution triple ( $y, p, u$ ) which allow to achieve accuracy $\varepsilon$ with an optimal order of arithmetic operations (when compared to wavelet-best $N$ term approximations) see $[\mathrm{DeVK}]$ and the articles therein for related concepts of nonlinear approximation applied to operator equations. These solvers are based on adaptive wavelet schemes for which convergence and optimal complexity of the scheme has been proved firstly for distributed and Neumann boundary control problems in [DK]. We wish to point out that it is not crucial in the present context to work with wavelets; any (possibly adaptive) scheme with the property that it guarantees to provide the solution triple ( $y, p, u$ ) up to accuracy $\varepsilon$ each with a provably minimal amount of degrees of freedom and complexity for a well-posed system of coupled PDEs having possibly non-smooth solutions would serve our purpose.

This paper is structured as follows. In the next section, it is proved that the solution of a linear operator equation involving a general parameter-dependent saddle point operator in an abstract setting is analytic, with precise bounds on the growth of the partial derivatives. This allows us in Section 2.4 to obtain rates of $N$-term generalized polynomial chaos approximations. These results are specified in Section 3 to linear-quadratic control problems constrained by an elliptic PDE with distributed, Neumann or Dirichlet boundary controls and in Section 4 to control problems constrained by linear parabolic PDEs. We conclude in Section 5 with some remarks how to realize this practically and how to combine the gpc approximations with discretizations with respect to space and time.
2. Parametric saddle point problems. We generalize the results of [CDS1] and study well-posedness, regularity and polynomial approximation of solutions for a family of abstract parametric saddle point problems. Particular attention is paid to the case of countably many parameters. The abstract results in the present section are more general than what is required in our ensuing treatment of optimal control problems and are of independent interest. We have in mind (and will discuss in detail in the following sections) optimal control problems for systems constrained by elliptic and parabolic PDEs with random coefficients.
2.1. An abstract result. Throughout, we denote by $\mathcal{X}$ and $\mathcal{Y}$ two reflexive Banach spaces over $\mathbb{R}$ (all results will hold with the obvious modifications also for spaces over $\mathbb{C}$ ) with (topological) duals $\mathcal{X}^{\prime}$ and $\mathcal{Y}^{\prime}$, respectively. By $\mathcal{L}\left(\mathcal{X}, \mathcal{Y}^{\prime}\right)$, we denote the set of bounded linear operators $G: \mathcal{X} \rightarrow \mathcal{Y}^{\prime}$. The Riesz representation theorem associates with each $G \in \mathcal{L}\left(\mathcal{X}, \mathcal{Y}^{\prime}\right)$ a unique bilinear form $\mathcal{G}(\cdot, \cdot): \mathcal{X} \times \mathcal{Y} \rightarrow \mathbb{R}$ by means of

$$
\begin{equation*}
\mathcal{G}(v, w)=\langle w, G v\rangle_{\mathcal{Y} \times \mathcal{Y}^{\prime}} \quad \text { for all } v \in \mathcal{X}, w \in \mathcal{Y} \tag{2.1}
\end{equation*}
$$

Here and in what follows, we indicate spaces in duality pairings $\langle\cdot, \cdot\rangle$ by subscripts.
We shall be interested in the solution of linear operator equations $G q=g$ and make use of the following solvability result which is a straightforward consequence of the closed graph theorem, see, e.g., $[\mathrm{BF}]$.

Proposition 1. An operator $G \in \mathcal{L}\left(\mathcal{X}, \mathcal{Y}^{\prime}\right)$ is boundedly invertible if and only if its associated bilinear form satisfies the inf-sup conditions: there exists a constant
$\gamma>0$ such that

$$
\begin{equation*}
\inf _{0 \neq v \in \mathcal{X}} \sup _{0 \neq w \in \mathcal{Y}} \frac{\mathcal{G}(v, w)}{\|v\| \mathcal{X}\|w\| \mathcal{Y}} \geq \gamma \tag{2.2}
\end{equation*}
$$

and

$$
\begin{equation*}
\inf _{0 \neq w \in \mathcal{Y}} \sup _{0 \neq v \in \mathcal{X}} \frac{\mathcal{\mathcal { X }}(v, w)}{\|v\| \boldsymbol{\mathcal { X }}\|w\| \mathcal{Y}} \geq \gamma \tag{2.3}
\end{equation*}
$$

If (2.2) and (2.3) hold, then for every $g \in \mathcal{Y}^{\prime}$ the operator equation

$$
\begin{equation*}
\text { find } q \in \mathcal{X}: \quad \mathcal{G}(q, v)=\langle g, v\rangle_{\mathcal{Y}^{\prime} \times \mathcal{Y}} \quad \forall v \in \mathcal{Y} \tag{2.4}
\end{equation*}
$$

admits a unique solution $q \in \mathcal{X}$. There holds the a-priori estimate

$$
\begin{equation*}
\|q\|_{\mathcal{X}} \leq \frac{\|g\|_{\boldsymbol{\mathcal { Y }}^{\prime}}}{\gamma} \tag{2.5}
\end{equation*}
$$

2.2. Parametric operator families. In the present paper, we shall be interested in parametric families of operators $G$. We admit both, finitely many as well as infinitely many parameters. To this end, we denote by $\sigma:=\left(\sigma_{j}\right)_{j \in \mathbb{S}} \in \mathcal{S}$ the set of parameters where $\mathbb{S} \subseteq \mathbb{N}$ is an at most countable index set. We assume the parameters to take values in $\mathcal{S} \subset \mathbb{R}^{\mathbb{S}}$. In particular, in the case $\mathbb{S}=\mathbb{N}$ it holds $\mathcal{S} \subseteq \mathbb{R}^{\mathbb{N}}$, i.e., each realization of $\sigma$ is a sequence of real numbers. We shall consider in particular the parameter domain $\mathcal{S}=[-1,1]^{\mathbb{N}}$ which we equip with the uniform probability measure

$$
\begin{equation*}
\rho(\sigma)=\bigotimes_{j \geq 1} \frac{d \sigma_{j}}{2} \tag{2.6}
\end{equation*}
$$

By $\mathbb{N}_{0}^{\mathbb{N}}$ we denote the set of all sequences of nonnegative integers, and by $\mathfrak{F}=\left\{\nu \in \mathbb{N}_{0}^{\mathbb{N}}\right.$ : $|\nu|<\infty\}$ the set of "finitely supported" such sequences, i.e., sequences of nonnegative integers which have only a finite number of nonzero entries. For $\nu \in \mathfrak{F}$, we denote by $\mathfrak{n} \subset \mathbb{N}$ the set of coordinates $j$ such that $\nu_{j} \neq 0$, with $j$ repeated $\nu_{j} \geq 1$ many times. Analogously, $\mathfrak{m} \subset \mathbb{N}$ denotes the supporting coordinate set for $\mu \in \mathfrak{F}$.

We consider parametric families of continuous, linear operators which we denote as $G(\sigma) \in \mathcal{L}\left(\mathcal{X}, \mathcal{Y}^{\prime}\right)$. We now make precise the dependence of $G(\sigma)$ on the parameter sequence $\sigma$ which is required for our regularity and approximation results.

Assumption 1. The parametric operator family $\left\{G(\sigma) \in \mathcal{L}\left(\mathcal{X}, \mathcal{Y}^{\prime}\right): \sigma \in \mathcal{S}\right\}$ is a regular $\mathfrak{p}$-analytic operator family for some $0<\mathfrak{p} \leq 1$, i.e.,
(i) $G(\sigma) \in \mathcal{L}\left(\mathcal{X}, \mathcal{Y}^{\prime}\right)$ is boundedly invertible for every $\sigma \in \mathcal{S}$ with uniformly bounded inverses $G(\sigma)^{-1} \in \mathcal{L}\left(\mathcal{Y}^{\prime}, \boldsymbol{\mathcal { X }}\right)$, i.e., there exists $C_{0}>0$ such that

$$
\begin{equation*}
\sup _{\sigma \in \mathcal{S}}\left\|G(\sigma)^{-1}\right\|_{\mathcal{L}\left(\boldsymbol{\mathcal { Y }}^{\prime}, \boldsymbol{\mathcal { X }}\right)} \leq C_{0} \tag{2.7}
\end{equation*}
$$

and
(ii) for any fixed $\sigma \in[-1,1]^{\mathbb{N}}$, the operators $G(\sigma)$ are analytic with respect to each $\sigma_{j}$ such that there exists a nonnegative sequence $b=\left(b_{j}\right)_{j \geq 1} \in \ell^{\mathfrak{p}}(\mathbb{N})$ such that

$$
\begin{equation*}
\forall \nu \in \mathfrak{F} \backslash\{0\}: \quad \sup _{\sigma \in \mathcal{S}}\left\|(G(0))^{-1}\left(\partial_{\sigma}^{\nu} G(\sigma)\right)\right\|_{\mathcal{L}(\boldsymbol{\mathcal { X }}, \boldsymbol{\mathcal { X }})} \leq C_{0} b^{\nu} \tag{2.8}
\end{equation*}
$$

Here $\partial_{\sigma}^{\nu} G(\sigma):=\frac{\partial^{\nu_{1}}}{\partial \sigma_{1}} \frac{\partial^{\nu_{2}}}{\partial \sigma_{2}} \cdots G(\sigma)$; the notation $b^{\nu}$ signifies the (finite due to $\nu \in \mathfrak{F}$ ) product $b_{1}^{\nu_{1}} b_{2}^{\nu_{2}} \ldots$ and we use the convention $0^{0}:=1$.

Affine Parameter Dependence. The special case of affine parameter dependence arises, for example, in diffusion problems where the diffusion coefficients are given in terms of a Karhunen-Loève expansion (see, e.g. [ST] for such KarhunenLoève expansions and their numerical analysis, in the context of elliptic PDEs with random coefficients). Then, there exists a family $\left\{G_{j}\right\}_{j \geq 0} \subset \mathcal{L}\left(\mathcal{X}, \mathcal{Y}^{\prime}\right)$ such that $G(\sigma)$ can be written in the form

$$
\begin{equation*}
\forall \sigma \in \mathcal{S}: \quad G(\sigma)=G_{0}+\sum_{j \geq 1} \sigma_{j} G_{j} \tag{2.9}
\end{equation*}
$$

We shall refer to $G_{0}=G(0)$ as "nominal", or 'mean-field" operator, and to $G_{j}, j \geq 1$ as "fluctuation" operators. In order for the sum in (2.9) to converge, we impose the following assumptions on $\left\{G_{j}\right\}_{j \geq 0}$. In doing so, we associate with the operator $G_{j}$ the bilinear forms $\mathcal{G}_{j}(\cdot, \cdot): \mathcal{X} \times \mathcal{Y} \rightarrow \mathbb{R}$.

ASSUMPTION 2. The family of operators $\left\{G_{j}\right\}_{j \geq 0}$ in (2.9) satisfies the following conditions:

1. The "mean field" operator $G_{0} \in \mathcal{L}\left(\mathcal{X}, \mathcal{Y}^{\prime}\right)$ is boundedly invertible, i.e. (cf. Proposition 1) there exists $\gamma_{0}>0$ such that

$$
\begin{equation*}
\inf _{0 \neq v \in \mathcal{X}} \sup _{0 \neq w \in \mathcal{Y}} \frac{\mathcal{G}_{0}(v, w)}{\|v\|_{\mathcal{X}}\|w\|_{\mathcal{Y}}} \geq \gamma_{0} \tag{2.10}
\end{equation*}
$$

and that

$$
\begin{equation*}
\inf _{0 \neq w \in \mathcal{Y}} \sup _{0 \neq v \in \mathcal{X}} \frac{\mathcal{G}_{0}(v, w)}{\|v\| \boldsymbol{\mathcal { X }}\|w\| \mathcal{Y}} \geq \gamma_{0} \tag{2.11}
\end{equation*}
$$

2. The "fluctuation" operators $\left\{G_{j}\right\}_{j \geq 1}$ are small with respect to $G_{0}$ in the following sense: there exists a constant $0<\kappa<1$ such that

$$
\begin{equation*}
\sum_{j \geq 1}\left\|G_{j}\right\|_{\mathcal{X} \rightarrow \mathcal{Y}^{\prime}} \leq \kappa \gamma_{0} \tag{2.12}
\end{equation*}
$$

We remark that with (2.10), (2.11), condition (2.12) follows from

$$
\begin{equation*}
\sum_{j \geq 1}\left\|G_{0}^{-1} G_{j}\right\|_{\mathcal{X} \rightarrow \mathcal{X}} \leq \kappa \tag{2.13}
\end{equation*}
$$

We show next that, under Assumption 2, the parametric family $G(\sigma)$ is boundedly invertible uniformly with respect to the parameter vector $\sigma$ belonging to the parameter domain $\mathcal{S}=[-1,1]^{\mathbb{N}}$.

Theorem 2. Under Assumption 2, for every realization $\sigma \in \mathcal{S}=[-1,1]^{\mathbb{N}}$ of the parameter vector, the parametric operator $G(\sigma)$ is boundedly invertible. Specifically, for the bilinear form $\mathcal{G}(\sigma ; \cdot, \cdot): \mathcal{X} \times \mathcal{Y} \rightarrow \mathbb{R}$ associated with $G(\sigma) \in \mathcal{L}\left(\mathcal{X}, \mathcal{Y}^{\prime}\right)$ there hold the uniform inf-sup conditions with $\gamma=(1-\kappa) \gamma_{0}>0$

$$
\begin{equation*}
\forall \sigma \in \mathcal{S}: \quad \inf _{0 \neq v \in \mathcal{X}} \sup _{0 \neq w \in \mathcal{Y}} \frac{\mathcal{G}(\sigma ; v, w)}{\|v\| \boldsymbol{\mathcal { X }}\|w\| \mathcal{Y}} \geq \gamma \tag{2.14}
\end{equation*}
$$

and

$$
\begin{equation*}
\forall \sigma \in \mathcal{S}: \quad \inf _{0 \neq w \in \mathcal{Y}} \sup _{0 \neq v \in \mathcal{X}} \frac{\mathcal{G}(\sigma ; v, w)}{\|v\| \mathcal{X}\|w\| \mathcal{Y}} \geq \gamma \tag{2.15}
\end{equation*}
$$

In particular, for every $g \in \mathcal{Y}^{\prime}$ and for every $\sigma \in \mathcal{S}$, the parametric operator equation

$$
\begin{equation*}
\text { find } q(\sigma) \in \mathcal{X}: \quad \mathcal{G}(\sigma ; q(\sigma), v)=\langle g, v\rangle_{\mathcal{Y} \times \mathcal{Y}^{\prime}} \quad \forall v \in \mathcal{Y} \tag{2.16}
\end{equation*}
$$

admits a unique solution $q(\sigma)$ which satisfies the a-priori estimate

$$
\begin{equation*}
\sup _{\sigma \in S}\|q(\sigma)\|_{\mathcal{X}} \leq \frac{\|g\|_{\boldsymbol{\mathcal { Y }}^{\prime}}}{(1-\kappa) \gamma_{0}} \tag{2.17}
\end{equation*}
$$

Proof. As the result is essentially a perturbation result, there are several ways to prove it. One approach, which was used for example in [G], is based on a Neumann Series argument. We give an alternative proof by verifying the inf-sup conditions directly. The inf-sup condition (2.2) is equivalent to the following assertion: given $v \in \mathcal{X}$, there exists $w_{v} \in \mathcal{Y}$ such that i) $\left\|w_{v}\right\| \mathcal{Y} \leq c_{1}\|v\|_{\mathcal{X}}$ and ii) $\mathcal{G}\left(v, w_{v}\right) \geq c_{2}\|v\|_{\mathcal{X}}^{2}$. Then (2.2) holds with $\gamma=c_{2} / c_{1}$.

By Assumption 2, in particular by (2.10), i) and ii) are satisfied for the bilinear form $\mathcal{G}_{0}(\cdot, \cdot)$ with constants $c_{1,0}$ and $c_{2,0}$, i.e., $\gamma_{0}=c_{2,0} / c_{1,0}$.

With $v \in \mathcal{X}$ arbitrary and with $w_{v} \in \mathcal{Y}$ as in i) and ii) for the bilinear form $\mathcal{G}_{0}(\cdot, \cdot)$ (in particular, independent of $\sigma$ ), we obtain for every $\sigma \in \mathcal{S}=[-1,1]^{\mathbb{N}}$

$$
\begin{aligned}
\mathcal{G}\left(\sigma ; v, w_{v}\right) & =\mathcal{G}_{0}\left(v, w_{v}\right)+\sum_{j \geq 1} \sigma_{j} \mathcal{G}_{j}\left(v, w_{v}\right) \\
& \geq c_{2,0}\|v\|_{\mathcal{X}}^{2}-\sum_{j \geq 1}\left|\mathcal{G}_{j}\left(v, w_{v}\right)\right| \\
& =c_{2,0}\|v\|_{\mathcal{X}}^{2}-c_{1,0} \sum_{j \geq 1}\left\|G_{j}\right\|_{\mathcal{X} \rightarrow \mathcal{Y}^{\prime}}\|v\|_{\mathcal{X}}^{2} \\
& =\left(c_{2,0}-c_{1,0} \sum_{j \geq 1}\left\|G_{j}\right\|_{\mathcal{X} \rightarrow \mathcal{Y}^{\prime}}\right)\|v\|_{\mathcal{X}}^{2} \\
& \geq c_{2,0}(1-\kappa)\|v\|_{\mathcal{X}}^{2} \\
& \geq \frac{c_{2,0}}{c_{1,0}}(1-\kappa)\|v\|_{\mathcal{X}}\left\|w_{v}\right\|_{\mathcal{Y}} \\
& =\gamma_{0}(1-\kappa)\|v\|_{\mathcal{X}}\left\|w_{v}\right\|_{\mathcal{Y}} .
\end{aligned}
$$

This implies (2.14). The stability condition (2.15) is verified analogously. The a-priori bound (2.17) follows then from (2.5) with the constant $\gamma=(1-\kappa) \gamma_{0}$.

From the preceding considerations, the following is readily verified.
Corollary 3. The affine parametric operator family (2.9) satisfies Assumption 1 with

$$
C_{0}=\frac{1}{(1-\kappa) \gamma_{0}} \quad \text { and } \quad b_{j}:=\frac{\left\|G_{j}\right\|_{\mathcal{X} \rightarrow \mathcal{Y}^{\prime}}}{(1-\kappa) \gamma_{0}} \quad \text { for all } j \geq 1
$$

2.3. Analytic dependence of solutions. We now establish that the dependence of the solution $q(\sigma)$ on $\sigma$ is analytic, with precise bounds on the growth of the partial derivatives. There holds

Theorem 4. Under Assumption 1, for every $f \in \mathcal{Y}^{\prime}$ and every $\sigma \in \mathcal{S}$ there exists a unique solution $q(\sigma) \in \boldsymbol{\mathcal { X }}$ of the parametric operator equation

$$
\begin{equation*}
G(\sigma) q(\sigma)=f \quad \text { in } \mathcal{Y}^{\prime} \tag{2.18}
\end{equation*}
$$

The parametric solution family $q(\sigma)$ depends analytically on the parameters, and the partial derivatives of the parametric solution family $q(\sigma)$ satisfy the bounds

$$
\begin{equation*}
\sup _{\sigma \in \mathcal{S}}\left\|\left(\partial_{\sigma}^{\nu} q\right)(\sigma)\right\|_{\mathcal{X}} \leq C_{0}\|f\|_{\mathcal{Y}^{\prime}}|\nu|!\tilde{b}^{\nu} \quad \text { for all } \nu \in \mathfrak{F} \tag{2.19}
\end{equation*}
$$

where $0!:=1$ and the sequence $\tilde{b}=\left(\tilde{b}_{j}\right)_{j \geq 1} \in \ell^{\mathfrak{p}}(\mathbb{N})$ is defined by

$$
\tilde{b}_{j}=b_{j} / \ln 2 \quad \text { for all } j \in \mathbb{N}
$$

Proof. Rather than proving (2.19), we prove the (slightly) stronger bound

$$
\begin{equation*}
\sup _{\sigma \in \mathcal{S}}\left\|\left(\partial_{\sigma}^{\nu} q\right)(\sigma)\right\|_{\mathcal{X}} \leq C_{0}\|f\|_{\mathcal{Y}^{\prime}} d_{|\nu|} b^{\nu} \quad \text { for all } \nu \in \mathfrak{F} \tag{2.20}
\end{equation*}
$$

where the sequence $d=\left(d_{n}\right)_{n \geq 0}$ is defined recursively by

$$
\begin{equation*}
d_{0}:=1, \quad d_{n}:=\sum_{i=0}^{n-1}\binom{n}{i} d_{i}, \quad n=1,2, \ldots \tag{2.21}
\end{equation*}
$$

The proof of (2.20) proceeds by induction with respect to $|\nu|$ : if $|\nu|=0, \nu=0$ and the assertion (2.20) follows from (2.7) and the a-priori bound (2.5). For $0 \neq \nu \in \mathfrak{F}$, we take the derivative $\partial_{\sigma}^{\nu}$ of the equation (2.18). Recalling for the (finitely supported) multiindices $\nu, \mu \in \mathfrak{F}$ their associated (finite) index sets $\mathfrak{n}, \mathfrak{m} \subset \mathbb{N}$ and abbreviate $n:=|\mathfrak{n}|=|\nu|, m:=|\mathfrak{m}|=|\mu|$, respectively, we find with the generalized product rule due to the $\sigma$-independence of $f$ the identity

$$
\sum_{\mathfrak{m} \in \mathfrak{P}(\mathfrak{n})} \partial_{\sigma}^{\mathfrak{n} \backslash \mathfrak{m}}(G(\sigma)) \partial_{\sigma}^{\mathfrak{m}}(q(\sigma))=0 \quad \text { for all } \sigma \in \mathcal{S}
$$

Here, $\mathfrak{P}(\mathfrak{n})$ denotes the power set of $\mathfrak{n} \subset \mathbb{N}$. Solving this identity for $\partial_{\sigma}^{\mathfrak{n}}(q(\sigma))$, we find

$$
G(\sigma)\left(\partial_{\sigma}^{\mathfrak{n}} q\right)(\sigma)=-\sum_{\mathfrak{m} \in \mathfrak{P}(\mathfrak{n}) \backslash\{\mathfrak{n}\}} \partial_{\sigma}^{\mathfrak{n} \backslash \mathfrak{m}}(G(\sigma)) \partial_{\sigma}^{\mathfrak{m}}(q(\sigma)) \quad \text { in } \quad \mathcal{Y}^{\prime}
$$

From the bounded invertibility of $G(\sigma)$, we get the recursion

$$
\begin{equation*}
\left(\partial_{\sigma}^{\nu} q\right)(\sigma)=-\sum_{\mathfrak{m} \in \mathfrak{P}(\mathfrak{n}) \backslash\{\mathfrak{n}\}}(G(\sigma))^{-1} \partial_{\sigma}^{\mathfrak{n} \backslash \mathfrak{m}}(G(\sigma)) \partial_{\sigma}^{\mathfrak{m}}(q(\sigma)) \quad \text { in } \quad \mathcal{Y}^{\prime} \tag{2.22}
\end{equation*}
$$

Taking the $\|\cdot\|_{\mathcal{X}}$ norm on both sides and using the triangle inequality, we find

$$
\begin{align*}
\left\|\left(\partial_{\sigma}^{\nu} q\right)(\sigma)\right\|_{\mathcal{X}} & \leq \sum_{\substack{\mathfrak{m} \in \mathfrak{P}(\mathfrak{n}) \backslash\{\mathfrak{n}\}}}\left\|(G(\sigma))^{-1} \partial_{\sigma}^{\mathfrak{n} \backslash \mathfrak{m}}(G(\sigma))\right\|_{\mathcal{L}(\mathcal{X}, \mathcal{X})}\left\|\partial_{\sigma}^{\mathfrak{m}}(q(\sigma))\right\|_{\mathcal{X}} \\
& \leq \sum_{m=0}^{n-1} \sum_{\substack{\mathfrak{m} \in \mathfrak{P}(\mathfrak{n}) \\
|\mathfrak{m}|=m}}\left\|(G(\sigma))^{-1} \partial_{\sigma}^{\mathfrak{n} \backslash \mathfrak{m}}(G(\sigma))\right\|_{\mathcal{L}(\mathcal{X}, \mathcal{X})}\left\|\partial_{\sigma}^{\mathfrak{m}}(q(\sigma))\right\|_{\mathcal{X}} \tag{2.23}
\end{align*}
$$

Now (2.20) for $n=|\nu|=1$ follows directly, upon using (2.8) for the singleton sets $\mathfrak{n}=\{j\}$.

We now proceed by induction with respect to $|\nu|$. We consider $\nu \in \mathfrak{F}$ such that $n=|\nu| \geq 2$ and assume that the assertion (2.20) has already been proved for all $\tilde{\nu} \in \mathfrak{F}$ such that $1 \leq|\tilde{\nu}|<n$. We then obtain from (2.23)

$$
\begin{aligned}
\left\|\left(\partial_{\sigma}^{\nu} q\right)(\sigma)\right\|_{\mathcal{X}} & \leq \sum_{\mathfrak{m} \in \mathfrak{P}(\mathfrak{n}) \backslash\{\mathfrak{n}\}}\left\|(G(\sigma))^{-1} \partial_{\sigma}^{\mathfrak{n} \backslash \mathfrak{m}}(G(\sigma))\right\|_{\mathcal{L}(\mathcal{X}, \mathcal{X})}\left\|\partial_{\sigma}^{\mathfrak{m}}(q(\sigma))\right\|_{\mathcal{X}} \\
& \leq \sum_{m=0}^{n-1} \sum_{\substack{\mathfrak{m} \in \mathfrak{P}(\mathfrak{n}) \\
|\mathfrak{m}|=m}} C_{0}\|f\|_{\mathcal{Y}^{\prime}} b^{\nu-\mu} d_{m} b^{\mu} \\
& =C_{0}\|f\|_{\mathcal{Y}^{\prime}} b^{\nu} \sum_{m=0}^{n-1}\binom{n}{m} d_{m} \\
& =C_{0}\|f\|_{\mathcal{Y}^{\prime}} b^{\nu} d_{n}
\end{aligned}
$$

which is $(2.20)$ for $\nu \in \mathfrak{F}$ such that $|\nu|=n$.
The assertion (2.19) now follows from (2.20) and the elementary inequality

$$
d_{n} \leq\left(\frac{1}{\ln 2}\right)^{n} n!\quad \text { for all } n \in \mathbb{N}
$$

2.4. Rates of $N$-term gpc approximation. The estimates (2.19) of the partial derivatives of $q(\sigma)$ with respect to $\sigma$ will be the basis for quantifying approximability of $q(\sigma)$ in the space $L^{2}(\mathcal{S}, \rho ; \mathcal{X})$. To this end, let $L_{n}(t)$ denote the Legendre polynomial of degree $n \geq 0$ in $(-1,1)$ which is normalized such that

$$
\int_{-1}^{1}\left|L_{n}(t)\right|^{2} \frac{d t}{2}=1
$$

Then $L_{0}=1$ and $\left\{L_{n}\right\}_{n \geq 0}$ is an orthonormal basis of $L^{2}(-1,1)$. For $\nu \in \mathfrak{F}$, denote $\nu!=\nu_{1}!\nu_{2}!\ldots$ and introduce the tensorized Legendre polynomials

$$
L_{\nu}(\sigma)=\prod_{j \geq 1} L_{\nu_{j}}\left(\sigma_{j}\right)
$$

Note that for each $\nu \in \mathfrak{F}$, there are only finitely many nontrivial factors in this product, and each $L_{\nu}(\sigma)$ depends only on finitely many of the $\sigma_{j}$. By construction, the countable collection $\left\{L_{\nu}(\sigma): \nu \in \mathfrak{F}\right\}$ is a Riesz basis, i.e. a dense, orthonormal family in $L^{2}(\mathcal{S}, \rho)$ : in particular, each $v \in L^{2}(\mathcal{S}, \rho ; \mathcal{X})$ admits an orthogonal expansion

$$
\begin{equation*}
v(\sigma)=\sum_{\nu \in \mathfrak{F}} v_{\nu} L_{\nu}(\sigma), \quad \text { where } \quad v_{\nu}:=\int_{\mathcal{S}} v(\sigma) L_{\nu}(\sigma) d \rho(\sigma) \in \mathcal{X} \tag{2.24}
\end{equation*}
$$

and there holds Parseval's equality

$$
\begin{equation*}
\|v\|_{L^{2}(\mathcal{S}, \rho ; \boldsymbol{\mathcal { X }})}^{2}=\sum_{\mu \in \mathfrak{F}}\left\|v_{\nu}\right\|_{\mathcal{X}}^{2} \tag{2.25}
\end{equation*}
$$

The Legendre representation (2.24) is the basis for the analysis of best $N$-term approximation rates. To this end, denote by $\Lambda \subset \mathfrak{F}$ a subset of cardinality $N=\# \Lambda<\infty$. Then, with $q_{\nu}$ denoting the Legendre coefficients of the solution $q(\sigma)$ of the parametric
operator equation (2.18), Parseval's identity (2.25) implies

$$
\begin{aligned}
\| q(\sigma) & -\sum_{\nu \in \Lambda} q_{\nu} L_{\nu}(\sigma) \|_{L^{2}(\mathcal{S}, \rho ; \boldsymbol{\mathcal { X }})}^{2} \\
& =\inf \left\{\left\|q(\sigma)-v_{\Lambda}\right\|_{L^{2}(\mathcal{S}, \rho ; \mathcal{X})}^{2}: v_{\Lambda} \in \operatorname{span}\left\{\sum_{\nu \in \Lambda} v_{\nu} L_{\nu}(\sigma)\right\}\right\} \\
& =\sum_{\nu \notin \Lambda}\left\|q_{\nu}\right\|_{\mathcal{X}}^{2} .
\end{aligned}
$$

Best $N$-term approximation rates in $\|\cdot\|_{L^{2}(\mathcal{S}, \rho ; \boldsymbol{\mathcal { X }})}$ will therefore follow from summability of the norms $\alpha_{\nu}=\left\|q_{\nu}\right\|_{\mathcal{X}}$ of the Legendre coefficients by Stechkin's Lemma whose's proof is elementary, see, e.g., $[\mathrm{DeV}]$.

Lemma 5. Let $0<\mathfrak{p} \leq \mathfrak{q}$ and $\alpha=\left(\alpha_{\nu}\right)_{\nu \in \mathfrak{F}}$ be a sequence in $\ell^{\mathfrak{p}}(\mathfrak{F})$. If $\mathfrak{F}_{N}$ is the set of indices corresponding to the $N$ largest values of $\left|\alpha_{\nu}\right|$, we have

$$
\left(\sum_{\nu \notin \mathfrak{F}_{N}}\left|\alpha_{\nu}\right|^{\mathfrak{q}}\right)^{1 / \mathfrak{q}} \leq\|\alpha\|_{\ell^{\mathfrak{p}}(\mathfrak{F})} N^{-r}
$$

where $r:=\frac{1}{\mathfrak{p}}-\frac{1}{\mathfrak{q}} \geq 0$.
We therefore need to address the $\mathfrak{p}$-summability of the $\|\cdot\|_{\mathcal{X}}$ norms of the Legendre coefficients $q_{\nu}$ of $q(\sigma)$. We first prove estimates for these coefficients.

Proposition 6. Let $0<\mathfrak{p} \leq 1$ and $b=\left(b_{j}\right)_{j \geq 1}$ be as in Assumption 1 above. Moreoever, let the sequence $d=\left(d_{j}\right)_{j \geq 1}$ be defined by $d_{j}:=\beta b_{j}$ where $\beta=1 /(\sqrt{3} \ln 2)$, and $\tilde{b}=\left(\tilde{b}_{j}\right)_{j \geq 1}$ be defined by $\tilde{b}_{j}:=b_{j} / \ln 2$. Under Assumption 1, we then have for all $\nu \in \mathfrak{F}$

$$
\begin{equation*}
\left\|q_{\nu}\right\|_{\mathcal{X}} \leq C_{0}\|f\|_{\mathcal{Y}^{\prime}} \frac{|\nu|!}{\nu!} d^{\nu} \tag{2.26}
\end{equation*}
$$

and

$$
\begin{equation*}
\left\|q_{\nu}\right\|_{\mathcal{X}}\left\|L_{\nu}\right\|_{L^{\infty}(\mathcal{S})} \leq C_{0}\|f\|_{\mathcal{Y}^{\prime}} \frac{|\nu|!}{\nu!} \tilde{b}^{\nu} \tag{2.27}
\end{equation*}
$$

Proof. In view of the representation (2.24) in terms of Legendre polynomials, the expansion coefficients $q_{\nu}$ of the solution $q(\sigma)$ of $(2.18)$ read for any $\nu \in \mathfrak{F}$

$$
\begin{equation*}
q_{\nu}=\int_{\mathcal{S}} q(\sigma) L_{\nu}(\sigma) d \rho(\sigma) \in \mathcal{X} \tag{2.28}
\end{equation*}
$$

Since $q(\sigma)$ depends analytically on $\sigma$, we can use repeated integration by parts to each of the one-dimensional integrals in (2.28), see the proof of Corollary 6.1 in [CDS1], to arrive at the a-priori estimate

$$
\left\|q_{\nu}\right\|_{\mathcal{X}} \leq \frac{\beta^{|\nu|}}{\nu!} \sup _{\sigma \in \mathcal{S}}\left\|\left(\partial_{\sigma}^{\nu} q\right)(\sigma)\right\|_{\mathcal{X}}
$$

Among others, such estimates allow to steer anisotropic sparse interpolation algorithms of Smolyak type.

Applying (2.19) to further estimate the right hand side immediately yields (2.26). Similarly, also the estimate (2.27) follows.

Based on the estimates in Proposition 6, we obtain the following result on convergence rates of best $N$-term polynomial approximations of the parametric solution $q(\sigma)$ of the parametric operator equation (2.18).

Theorem 7. Under Assumption 1 with some $0<\mathfrak{p} \leq 1$, there exists a sequence $\left(\Lambda_{N}\right)_{N \in \mathbb{N}} \subset \mathfrak{F}$ of index sets whose cardinality does not exceed $N$ and a constant $C>0$ independent of $N$ such that

$$
\begin{equation*}
\left\|q-q_{N}\right\|_{L^{2}(\mathcal{S}, \rho ; \mathcal{X})} \leq C N^{-r}, \quad r=\frac{1}{\mathfrak{p}}-\frac{1}{2} \tag{2.29}
\end{equation*}
$$

Here, $q_{N}:=q_{\Lambda_{N}}$ where $q_{\Lambda_{N}}$ denotes the sequence in $L^{2}(\mathcal{S}, \rho ; \mathcal{X})$ whose entries $q_{\nu}$ equal those of the sequence $q$ if $\nu \in \Lambda_{N} \subset \mathfrak{F}$ and which equal zero otherwise.

The proof of this theorem proceeds along the lines of the argument in [CDS1] for the parametric diffusion problem: we use the bounds (2.26) and (2.27), and Theorem 7.2 of [CDS1], i.e.,
for $0<\mathfrak{p} \leq 1: \quad\left(\frac{|\nu|!}{\nu!} \alpha^{\nu}\right)_{\nu \in \mathfrak{F}} \in \ell^{\mathfrak{p}}(\mathfrak{F}) \quad$ if and only if $\quad\|\alpha\|_{\ell^{1}(\mathbb{N})}<1$ and $\alpha \in \ell^{\mathfrak{p}}(\mathbb{N})$.
Applying this result to the sequences $\alpha=d$ and to $\alpha=\tilde{b}=(\ln 2)^{-1} b$, we obtain the $\mathfrak{p}$-summability, and, by referring to the Stechkin Lemma 5 with $\mathfrak{q}=2$ and the Parseval identity (2.25), the assertion (2.29) follows.

We next illustrate the scope of the foregoing abstract results with several concrete instances of PDE-constrained control problems: we consider problems constrained by parametric elliptic or parabolic PDE operators and different types of controls. In either case, we develop gpc approximation results by identifying the parametric control problem as particular case of the abstract parametric saddle point problem (2.18). Importantly, due to our formulation as a saddle point problem, the best $N$ term approximation rates obtained from Theorem 7 pertain to concurrent $N$-term approximation of state and control with the same set of active gpc coefficients.
3. Parametric Linear-Quadratic Elliptic Control Problems. We describe the setup of the control problem constrained by a linear parametric elliptic PDE by first addressing conditions on the PDE constraint as an operator equation with a parametric linear elliptic operator $A=A(\sigma)$ on a reflexive Banach space $Y$. Our standard example with be a scalar diffusion problem.

ASSUMPTION 3. For each fixed $\sigma \in \mathcal{S}$, the operator $A(\sigma) \in \mathcal{L}\left(Y, Y^{\prime}\right)$ is symmetric and boundedly invertible, i.e., $A(\sigma): Y \rightarrow Y^{\prime}$ is linear, self-adjoint, invertible and satisfies the continuity and coercivity estimates

$$
\begin{array}{rlr}
\left|\langle v, A(\sigma) w\rangle_{Y \times Y^{\prime}}\right| & \leq C_{A}\|v\|_{Y}\|w\|_{Y}, & v, w \in Y \\
\langle v, A(\sigma) v\rangle_{Y \times Y^{\prime}} \geq c_{A}\|v\|_{Y}^{2}, & v \in Y \tag{3.2}
\end{array}
$$

with some constants $0<c_{A} \leq C_{A}<\infty$ independent of $\sigma$.
These imply the estimates

$$
\begin{equation*}
c_{A}\|w\|_{Y} \leq\|A(\sigma) w\|_{Y^{\prime}} \leq C_{A}\|w\|_{Y} \quad \text { for any } w \in Y \tag{3.3}
\end{equation*}
$$

which, in terms of operator norms, may be expressed as

$$
\begin{equation*}
\|A(\sigma)\|_{Y \rightarrow Y^{\prime}}:=\sup _{w \in Y, w \not \equiv 0} \frac{\|A(\sigma) w\|_{Y^{\prime}}}{\|w\|_{Y}} \leq C_{A}, \quad\left\|A(\sigma)^{-1}\right\|_{Y^{\prime} \rightarrow Y} \leq c_{A}^{-1} \tag{3.4}
\end{equation*}
$$

If the precise format of the constants in (3.3) does not matter, we will abbreviate this as

$$
\begin{equation*}
\|A(\sigma) w\|_{Y^{\prime}} \sim\|w\|_{Y} \quad \text { for any } w \in Y \tag{3.5}
\end{equation*}
$$

and use $a \lesssim b$ or $a \gtrsim b$ for the corresponding one-sided estimates.
Some examples of operators $A$ and space $Y$ are provided next which satisfy Assumption 3 provided that (3.6) stated below holds. In all the following, $\Omega \subset \mathbb{R}^{\mathbb{d}}$ denotes a bounded domain with Lipschitz boundary $\partial \Omega$.

Example 8.
(i) (Dirichlet problem with homogeneous Dirichlet boundary conditions)
$\langle v, A(\sigma) w\rangle_{Y \times Y^{\prime}}=\int_{\Omega}\left(a(\sigma) \nabla_{x} v \cdot \nabla_{x} w\right) d x, Y=H_{0}^{1}(\Omega)$.
In this and all the following examples, the coefficient $a(\sigma)$ is supposed to satisfy the uniform ellipticity assumption $\operatorname{UEA}\left(r_{a}, R_{a}\right)$ : there exist positive constants $r_{a}, R_{a}$ such that for all $x \in \Omega$ and all $\sigma \in \mathcal{S}$ it holds

$$
\begin{equation*}
0<r_{a} \leq a(x, \sigma) \leq R_{a}<\infty \tag{3.6}
\end{equation*}
$$

(ii) (Reaction-diffusion problem with possibly anisotropic diffusion with Neumann boundary conditions)

$$
\langle v, A(\sigma) w\rangle_{Y \times Y^{\prime}}=\int_{\Omega}\left(a(\sigma) \nabla_{x} v \cdot \nabla_{x} w+v w\right) d x, \quad Y=H^{1}(\Omega)
$$

Note that Assumption 3 is, due to the self-adjointness a special case of the conditions on the operator $G$ in Proposition 1 with $\mathcal{X}=\mathcal{Y}=Y$. Thus, this assumption implies that for any given deterministic $f \in Y^{\prime}$ and fixed $\sigma \in \mathcal{S}$, the operator equation

$$
\begin{equation*}
A(\sigma) y=f \tag{3.7}
\end{equation*}
$$

has a unique solution $y=y(\sigma) \in Y$.
3.1. Distributed or Neumann boundary control. Allowing an additional function $u=u(\sigma)$ on the right hand side of (3.7), we ask to steer the solution of such an equation towards a prescribed desired deterministic state $y_{*}$, under the condition that the effort on $u$ should be minimal. Consequently, we can define an optimal control problem with a functional of tracking type as follows: minimize for $\sigma \in \mathcal{S}$ over the state $y(\sigma)$ and the control $u(\sigma)$ the functional

$$
\begin{equation*}
\tilde{J}(y(\sigma), u(\sigma)):=\frac{1}{2}\left\|\mathfrak{T} y(\sigma)-y_{*}\right\|_{O}^{2}+\frac{\omega}{2}\|u(\sigma)\|_{U}^{2} \tag{3.8}
\end{equation*}
$$

subject to the linear operator equation

$$
\begin{equation*}
A(\sigma) y(\sigma)=f+\mathfrak{E} u(\sigma) . \tag{3.9}
\end{equation*}
$$

Here $\omega>0$ is a fixed constant which balances the least squares approximation of the states and the norm for the control and $\mathfrak{T}, \mathfrak{E}$ are some linear (trace and extension) operators described below.

We need to add some requirements on the norms used in (3.8). In view of Assumption 3 , in order for (3.9) to have a well-defined unique solution, we need to assure that either $y \in Y$ or $\mathfrak{E} u \in Y^{\prime}$. The latter is satisfied if the control space $U$ defining the penalty norm part of the functional is such that $U \subseteq Y^{\prime}$ with continuous embedding. Then the observation space $O$ defining the least squares part of the functional
(3.19) may be chosen as any $O \supseteq Y$. In this case, $\mathfrak{T}$ may be any continuous linear operator from $Y$ onto its range, i.e., $\|\mathfrak{T} v\|_{\text {range( } \mathfrak{T})} \lesssim\|v\|_{Y}$ for $v \in Y$ with range( $\mathfrak{T}$ ) continuously embedded in $O$. Alternatively, assuring $\mathfrak{T} y \in O$ and selecting $O \subseteq Y$ embedded continuously would allow for any choice of $U$.

There are two standard examples covered by this formulation which we have in mind (see [DK] for more general formulations). A distributed control problem is one where the control is exerted on all of the right hand side of (3.9), i.e., $\mathfrak{E}$ is just the identity. This case is perhaps rather of academic nature but serves as a good illustration for the essential mechanisms.

Example 9. (Dirichlet problem with distributed control)
Here the PDE constraints are given by the standard scalar second order Dirichlet problem with distributed control,

$$
\begin{align*}
-\partial_{x}\left(a(\sigma) \partial_{x}\right) y(\sigma) & =f+u(\sigma) & & \text { in } \Omega \\
y(\sigma) & =0 & & \text { on } \partial \Omega \tag{3.10}
\end{align*}
$$

which gives rise to the operator equation (3.9) with

$$
\begin{equation*}
\langle v, A(\sigma) w\rangle_{Y \times Y^{\prime}}=\int_{\Omega} a(\sigma) \nabla_{x} v \cdot \nabla_{x} w d x, \quad Y=H_{0}^{1}(\Omega), Y^{\prime}=H^{-1}(\Omega) \tag{3.11}
\end{equation*}
$$

and given $f \in Y^{\prime}$. Admissible choices for $O, U$ are the classical case $O=U=L_{2}(\Omega)$, see $[L]$, or the natural choice $O=Y$ and $U=Y^{\prime}$, in which case the operators $\mathfrak{T}, \mathfrak{E}$ are the canonical injections $\mathfrak{T}=I$, $\mathfrak{E}=I$. Many more possible choices covering, in particular, fractional Sobolev spaces, have been discussed in [DK], as well as including a class of Neumann problems with distributed control.

Example 10. (Reaction-diffusion problem with Neumann boundary control) Consider the second order Neumann problem in strong form

$$
\begin{align*}
-\partial_{x}\left(a(\sigma) \partial_{x}\right) y(\sigma)+y(\sigma) & =f & & \text { in } \Omega  \tag{3.12}\\
\left(a(\sigma) \nabla_{x} y(\sigma)\right) \cdot \mathbf{n} & =u(\sigma) & & \text { on } \partial \Omega
\end{align*}
$$

where $\mathbf{n}$ denotes the outward normal at $\partial \Omega$. Here the weak form is based on setting $Y=H^{1}(\Omega)$ and

$$
\begin{equation*}
\langle v, A(\sigma) w\rangle_{Y \times Y^{\prime}}=\int_{\Omega}\left(a(\sigma) \nabla_{x} v \cdot \nabla_{x} w+v w\right) d x \tag{3.13}
\end{equation*}
$$

and given $f \in Y^{\prime}$. Recall that for any $v \in H^{1}(\Omega)$, its normal trace $\mathbf{n} \cdot \nabla_{x} v$ to $\partial \Omega$ belongs to $H^{-1 / 2}(\partial \Omega)$. Thus, in order for the right hand side of (3.12) to be welldefined, the control $u$ must belong to $H^{-1 / 2}(\partial \Omega)$, i.e., the operator $\mathfrak{E}$ is the adjoint of the normal trace operator, or, $\mathfrak{E}: H^{-1 / 2}(\partial \Omega) \rightarrow Y^{\prime}$ is an extension operator to $\Omega$. The formulation of the constraint as an operator equation reads in this case

$$
\begin{equation*}
A(\sigma) y(\sigma)=f+\mathfrak{E} u(\sigma) \tag{3.14}
\end{equation*}
$$

As previously, one could choose $O$ to be a space defined on $\Omega$. However, a more frequent practical situation arises when one wants to achieve a prescribed state on some part of the boundary. Denote by $\Gamma_{\circ} \subseteq \partial \Omega$ an observation boundary with strictly positive $\mathfrak{d}$-1-dimensional measure and by $\mathfrak{T}: H^{1}(\Omega) \rightarrow H^{1 / 2}\left(\Gamma_{\circ}\right)$ the trace operator to this part of the boundary. Then an admissible choice is $O=H^{1 / 2}\left(\Gamma_{\circ}\right)$. As discussed
above, we need to require for the control that $u \in H^{-1 / 2}(\partial \Omega)$. For these choices, the functional (3.8) is of the form

$$
\begin{equation*}
\tilde{J}(y, u)=\frac{1}{2}\left\|\mathfrak{T} y-y_{*}\right\|_{H^{1 / 2}\left(\Gamma_{\circ}\right)}^{2}+\frac{\omega}{2}\|u\|_{H^{-1 / 2}(\partial \Omega)}^{2} \tag{3.15}
\end{equation*}
$$

The fractional trace norms appearing here in a natural form are often replaced, perhaps partly due to the difficulty of evaluating fractional order Sobolev norms numerically, by the classical choice $\Gamma_{\circ}=\partial \Omega$ and $O=U=L_{2}(\partial \Omega)[L]$; we hasten to add, however, that in the context of multiresolution discretizations in $\Omega$ and on $\partial \Omega$, fractional Sobolev norms can be realized numerically in optimal complexity (see, e.g., [DK, GK] and the references there).

One calls (3.8) with constraints (3.9) a linear-quadratic control problem: a quadratic functional is to be minimized subject to a linear equation coupling state and control. From an optimization point of view, the solution of this problem has a simple structure: on account of $\tilde{J}$ being convex, one only needs to consider the first order conditions for optimality. To derive these, for $\sigma \in \mathcal{S}$, in principle, the dual operator of $A(\sigma)$ comes into play which is defined by

$$
\begin{equation*}
\left\langle A(\sigma)^{*} v, w\right\rangle_{Y^{\prime} \times Y}:=\langle v, A(\sigma) w\rangle_{Y \times Y^{\prime}} \tag{3.16}
\end{equation*}
$$

that is, $A(\sigma)^{*} \in \mathcal{L}\left(Y, Y^{\prime}\right)$. Of course, since in Assumption $3 A(\sigma)$ was required to be self-adjoint for each fixed $\sigma \in \mathcal{S}$, we have $A(\sigma)^{*}=A(\sigma)$.

Note that in case of an unsymmetric $A(\sigma)$, the property to be boundedly invertible (3.5) immediately carries over to $A(\sigma)^{*}$, that is, for fixed $\sigma \in \mathcal{S}$ and any $v \in Y$, one has the mapping property

$$
\begin{equation*}
\left\|A(\sigma)^{*} v\right\|_{Y^{\prime}} \sim\|v\|_{Y} \tag{3.17}
\end{equation*}
$$

For ease of presentation in this paper, we select here the natural case $O=Y$ and $U=Y^{\prime}$ resulting in $\mathfrak{T}=I$ and $\mathfrak{E}=I$ for the trace and extension operators. The more general case which may involve Sobolev spaces with possibly fractional smoothness indices has been treated in [DK] for PDE-constrained control problems without parameters.

To represent the Hilbert space norms in the optimization functional, we shall employ Riesz operators $R_{Y}: Y \rightarrow Y^{\prime}$ defined by

$$
\begin{equation*}
\left\langle v, R_{Y} w\right\rangle_{Y \times Y^{\prime}}:=(v, w)_{Y}, \quad v, w \in Y \tag{3.18}
\end{equation*}
$$

Defining $R_{Y^{\prime}}: Y^{\prime} \rightarrow Y$ correspondingly by $\left\langle v, R_{Y^{\prime}} w\right\rangle_{Y^{\prime} \times Y}:=(v, w)_{Y^{\prime}}$ for $v, w \in Y^{\prime}$, this implies $R_{Y^{\prime}}=R_{Y}^{-1}$ so that we can write both norms in the target functional in terms of one Riesz operator $R=R_{Y}$. Since the inner product $(\cdot, \cdot)_{Y}$ is symmetric, the Riesz operator $R$ is also symmetric.

Proposition 11. Necessary and sufficient for the linear-quadratic control problem to minimize for $\sigma \in \mathcal{S}$

$$
\begin{equation*}
J(y(\sigma), u(\sigma)):=\frac{1}{2}\left\|y(\sigma)-y_{*}\right\|_{Y}^{2}+\frac{\omega}{2}\|u(\sigma)\|_{Y^{\prime}}^{2} \tag{3.19}
\end{equation*}
$$

over all $(y(\sigma), u(\sigma)) \in Y \times Y^{\prime}$ subject to (3.9) are the Euler equations for the solution triple $(y(\sigma), p(\sigma), u(\sigma)) \in Y \times Y \times Y^{\prime}$

$$
\begin{align*}
A(\sigma) y(\sigma) & =f+u(\sigma) \\
A(\sigma)^{*} p(\sigma) & =R\left(y_{*}-y(\sigma)\right)  \tag{EE}\\
\omega R^{-1} u(\sigma) & =p(\sigma) \tag{3.20}
\end{align*}
$$

Proof. We present a proof of this well-known result only to bring out the roles of the Riesz operators; we skip in this proof the dependence of $\sigma$ for better readability. The derivation of (EE) is based on computing the zeroes of the first order variations of the Lagrangian functional

$$
\begin{equation*}
\operatorname{Lagr}(y, p, u):=J(y, u)+\langle p, A y-f-u\rangle_{Y \times Y^{\prime}} \tag{3.22}
\end{equation*}
$$

introducing a new variable $p \in Y$ called the Lagrangian or adjoint variable by which the constraints (3.9) are appended to the functional $J$, see, e.g., [L]. By inserting definition (3.19) and (3.18), the Lagrangian functional attains the form
$\operatorname{Lagr}(y, p, u)=\frac{1}{2}\left\langle y-y_{*}, R\left(y-y_{*}\right)\right\rangle_{Y \times Y^{\prime}}+\frac{\omega}{2}\left\langle u, R^{-1} u\right\rangle_{Y^{\prime} \times Y}+\langle p, A y-f-u\rangle_{Y \times Y^{\prime}}$.
The constraint (3.9) is recovered as the zero of the first variation of $\operatorname{Lagr}(y, p, u)$ in direction of $p$. Moreover, $\frac{\partial}{\partial u} \operatorname{Lagr}(y, p, u)=0$ yields $\omega R^{-1} u-p=0$. Finally,

$$
\begin{aligned}
\frac{\partial}{\partial y} \operatorname{Lagr}(y, p, u) & :=\lim _{\delta \rightarrow 0} \frac{\operatorname{Lagr}(y+\delta, p, u)-\operatorname{Lagr}(y, p, u)}{\delta} \\
& =\lim _{\delta \rightarrow 0} \frac{\frac{1}{2}\left\langle\delta, R\left(y-y_{*}\right)\right\rangle_{Y \times Y^{\prime}}+\frac{1}{2}\left\langle y-y_{*}, R \delta\right\rangle_{Y \times Y^{\prime}}+\langle p, A \delta\rangle_{Y \times Y^{\prime}}}{\delta} \\
& =\lim _{\delta \rightarrow 0} \frac{\left\langle\delta, R\left(y-y_{*}\right)\right\rangle_{Y \times Y^{\prime}}+\langle p, A \delta\rangle_{Y \times Y^{\prime}}}{\delta}
\end{aligned}
$$

by symmetry of $R$. Bringing $A$ on the other side of the dual form yields

$$
\frac{\partial}{\partial y} \operatorname{Lagr}(y, p, u)=R\left(y-y_{*}\right)+A^{*} p
$$

and therefore $\frac{\partial}{\partial y} \operatorname{Lagr}(y, p, u)=0$ if and only if (3.20) holds.
In our formulation, the design equation (3.21) expresses the control just as a weighted Riesz transformed adjoint state. For later analysis, it will help us to eliminate the control using (3.21) and write (EE) as the condensed Euler equations for the solution pair $(y(\sigma), p(\sigma)) \in Y \times Y$

$$
\begin{align*}
A(\sigma) y(\sigma) & =f+\frac{1}{\omega} R p(\sigma)  \tag{3.24}\\
A(\sigma)^{*} p(\sigma) & =R\left(y_{*}-y(\sigma)\right)
\end{align*}
$$

With the abbreviation $\hat{y}_{*}:=R y_{*} \in Y^{\prime}$, we write this as a coupled system to find for given $\left(f, \hat{y}_{*}\right) \in Y^{\prime} \times Y^{\prime}$ a solution pair $(y(\sigma), p(\sigma)) \in Y \times Y$ which solves

$$
\left(\begin{array}{cc}
A(\sigma) & -\frac{1}{\omega} R  \tag{3.25}\\
R & A(\sigma)^{*}
\end{array}\right)\binom{y(\sigma)}{p(\sigma)}=\binom{f}{\hat{y}_{*}} .
$$

Identifying the matrix operator appearing in this system with $G(\sigma)$ in the abstract problem in Section 2, we define the corresponding bilinear form $\mathcal{G}(\sigma ; \cdot, \cdot): \mathcal{X} \times \mathcal{X} \rightarrow \mathbb{R}$ where $\mathcal{X}:=Y \times Y, \mathcal{Y}:=\mathcal{X}$, for $q=(y, p), \tilde{q}=(\tilde{y}, \tilde{p}) \in \mathcal{X}$ by

$$
\begin{align*}
\mathcal{G}(\sigma ; q, \tilde{q}) & :=\left\langle q,\left(\begin{array}{cc}
A(\sigma) & -\frac{1}{\omega} R \\
R & A(\sigma)^{*}
\end{array}\right) \tilde{q}\right\rangle_{\mathcal{X} \times \mathcal{X}^{\prime}}  \tag{3.26}\\
& =\langle y, A(\sigma) \tilde{y}\rangle_{Y \times Y^{\prime}}-\frac{1}{\omega}\langle y, R \tilde{p}\rangle_{Y \times Y^{\prime}}+\langle p, R \tilde{y}\rangle_{Y \times Y^{\prime}}+\left\langle p, A(\sigma)^{*} \tilde{p}\right\rangle_{Y \times Y^{\prime}}
\end{align*}
$$

We equip the space $\mathcal{X}$ with the norm

$$
\begin{equation*}
\|q\|_{\mathcal{X}}=\left\|\binom{y}{p}\right\|_{Y \times Y}=\|y\|_{Y}+\|p\|_{Y} . \tag{3.27}
\end{equation*}
$$

Proposition 12. The parametric bilinear form $\mathcal{G}(\sigma ; \cdot, \cdot): \mathcal{X} \times \mathcal{X} \rightarrow \mathbb{R}$ is continuous on $\mathcal{X} \times \mathcal{X}$ for any constant weight $\omega>0$, and uniformly with respect to the parameter vector $\sigma$. It is coercive on $\mathcal{X}$ for $\omega=1$ with coercivity constant $c_{\mathcal{G}}:=\frac{c_{A}}{2}$ and $c_{A}$ from (3.2). For $0<\omega \leq 1$, it is coercive on $\mathcal{X}$ with a constant $c_{\mathcal{G}}$ defined below in (3.29). Moreover, it is symmetric for $\omega=1$.

Proof. The symmetry of $\mathcal{G}(\sigma ; \cdot, \cdot)$ for $\omega=1$ follows immediately from the representation (3.26) and by recalling that $A(\sigma)$ is self-adjoint. The continuity of $\mathcal{G}(\sigma ; \cdot, \cdot)$ results from the definition of $R(3.18)$ and applying Cauchy-Schwarz inequality, together with the continuity (3.1) of $A(\sigma)$, i.e., for any $q=(y, p), \tilde{q}=(\tilde{y}, \tilde{p}) \in \mathcal{X}$ we have from (3.26)

$$
\begin{align*}
|\mathcal{G}(\sigma ; q, \tilde{q})| & \leq\left|\langle y, A(\sigma) \tilde{y}\rangle_{Y \times Y^{\prime}}\right|+\left|\left\langle p, A(\sigma)^{*} \tilde{p}\right\rangle_{Y \times Y^{\prime}}\right|+\frac{1}{\omega}\left|(y, \tilde{p})_{Y}\right|+\left|(p, \tilde{y})_{Y}\right| \\
& \leq C_{A}\left(\|y\|_{Y}\|\tilde{y}\|_{Y}+\|p\|_{Y}\|\tilde{p}\|_{Y}\right)+\frac{1}{\omega}\|y\|_{Y}\|\tilde{p}\|_{Y}+\|p\|_{Y}\|\tilde{y}\|_{Y} \\
& \leq C_{A}\left(1+\frac{1}{\omega}\right)\|q\|_{\mathcal{X}}\|\tilde{q}\|_{\mathcal{X}} \quad=: C_{\mathcal{G}}\|q\| \boldsymbol{\mathcal { X }}\|\tilde{q}\|_{\mathcal{X}} . \tag{3.28}
\end{align*}
$$

As for the coercivity, for $q=(y, p) \in \mathcal{X}$, using the symmetry of $R$, its definition, the coercivity (3.2) and Cauchy-Schwarz' inequality, we infer for $0<\omega \leq 1$ (meaning that ( $1-\frac{1}{\omega} \leq 0$ )

$$
\begin{aligned}
\mathcal{G}(\sigma ; q, q) & =\langle y, A(\sigma) y\rangle_{Y \times Y^{\prime}}+\left\langle p, A(\sigma)^{*} p\right\rangle_{Y \times Y^{\prime}}+\left(1-\frac{1}{\omega}\right)\langle y, R p\rangle_{Y \times Y^{\prime}} \\
& \geq c_{A}\left(\|y\|_{Y}^{2}+\|p\|_{Y}^{2}\right)+\left(\frac{1}{\omega}-1\right)\left|(y, p)_{Y}\right| \\
& \geq c_{A}\left(\|y\|_{Y}^{2}+\|p\|_{Y}^{2}\right)+\left(\frac{1}{\omega}-1\right)\|y\|_{Y}\|p\|_{Y} .
\end{aligned}
$$

In case $\omega=1$, this immediately yields $\mathcal{G}(\sigma ; q, q) \geq \frac{c_{A}}{2}\|q\|_{\mathcal{X}}^{2}=c_{\mathcal{G}}\|q\|_{\mathcal{X}}^{2}$ for every $\sigma \in \mathcal{S}$. For $\omega<1$, we obtain

$$
\begin{align*}
\mathcal{G}(\sigma ; q, q) & \geq c_{A}\left(\|y\|_{Y}^{2}+\|p\|_{Y}^{2}\right)+\left(\frac{1}{\omega}-1\right)\|y\|_{Y}\|p\|_{Y} \\
& \geq \min \left\{\frac{c_{A}}{2}, \frac{1}{\omega}-1\right\}\|q\|_{\mathcal{X}}^{2}=: c_{\mathcal{G}}\|q\|_{\mathcal{X}}^{2} . \tag{3.29}
\end{align*}
$$

By the Theorem of Lax-Milgram, we therefore have, based on Proposition 12, the following result.

Theorem 13. Under Assumption 3, for every $0<\omega \leq 1$ and for every $\sigma \in \mathcal{S}$, the control problem (3.25) admits a unique solution $q(\sigma)=(y(\sigma), u(\sigma)) \in \mathcal{X}$ for any given deterministic right hand side $g:=\left(f, \hat{y}_{*}\right) \in \mathcal{X}^{\prime}$.

If, morever, the parametric family $\{A(\sigma): \sigma \in \mathcal{S}\}$ depends on $\sigma$ in a affine fashion, i.e., if

$$
\begin{equation*}
A(\sigma)=A_{0}+\sum_{j \geq 1} \sigma_{j} A_{j}, \tag{3.30}
\end{equation*}
$$

the parametric matrix operator $G(\sigma) \in \mathcal{L}\left(\mathcal{X}, \mathcal{X}^{\prime}\right)$ satisfies Assumption 2 with $\mathcal{X}=$ $\mathcal{Y}=Y \times Y$.

Corollary 14. On account of Theorem 4, the preceding result establishes the simultaneous analyticity of state as well as of the costate, with respect to all parameters
and therefore, by (3.21), also of the control. Moreover, by Theorem 7, this implies sparsity of the tensorized Legendre expansion of the solution triple ( $y, p, u$ ) and therefore, in particular, best $N$-term gpc approximation rates of all these quantities in the $L^{2}(\mathcal{S}, \rho ; \mathcal{X})$ norm.

REMARK 15. Note that the affine dependence of the operator $G(\sigma)$ in (3.26), see Corollary 3, was crucial in being able to use the abstract results of Section 2. Analogous analytic dependence results also hold for control problems with certain more general parameter dependences.

Occasionally, it is useful to derive from (3.25) an equation for the control alone.
Proposition 16. Under Assumption 3, system (EE) reduces to the condensed equation for the control

$$
\begin{equation*}
\left(A(\sigma)^{-*} R A(\sigma)^{-1}+\omega R^{-1}\right) u(\sigma)=A(\sigma)^{-*} R\left(y_{*}-A(\sigma)^{-1} f\right) \tag{3.31}
\end{equation*}
$$

(using $A^{-*}:=\left(A^{*}\right)^{-1}$ ) which we abbreviate as

$$
\begin{equation*}
M(\sigma) u(\sigma)=m(\sigma) \tag{3.32}
\end{equation*}
$$

Proof. On account of Assumption 3, $A(\sigma) \in \mathcal{L}\left(Y, Y^{\prime}\right)$ is invertible uniformly with respect to $\sigma \in \mathcal{S}$ so that (3.9) can be expressed as

$$
\begin{equation*}
y(\sigma)=A(\sigma)^{-1} f+A(\sigma)^{-1} u(\sigma) \tag{3.33}
\end{equation*}
$$

Inserted into (3.20) this yields

$$
\begin{equation*}
A(\sigma)^{*} p(\sigma)=-R A(\sigma)^{-1} u(\sigma)+R\left(y_{*}-A(\sigma)^{-1} f\right) \tag{3.34}
\end{equation*}
$$

and, again by Assumption 3,

$$
p(\sigma)=-A(\sigma)^{-*} R A(\sigma)^{-1} u(\sigma)+A(\sigma)^{-*} R\left(y_{*}-A(\sigma)^{-1} f\right)
$$

Using the identity (3.21), we can eliminate $p(\sigma)$ and arrive at

$$
\omega R^{-1} u(\sigma)=-A(\sigma)^{-*} R A(\sigma)^{-1} u(\sigma)+A(\sigma)^{-*} R\left(y_{*}-A(\sigma)^{-1} f\right)
$$

which is just (3.31).
Remark 17. We observe that the condensed equation (3.31) contains the boundedly invertible, parametric Schur complement operator $M(\sigma)$; this operator, while being boundedly invertible, is not affine in the parameter vector $\sigma$ anymore. Therefore, the theory developed in Section 2 does not apply. Nevertheless, analytic parameter dependence can be inferred for $M(\sigma)$ from the structure of its definition, and analytic continuation as in [CDS2] will allow to infer directly analytic dependence and best $N$-term gpc approximation rates for the control $u(\sigma)$ without approximation of the state. As this requires introduction of complex extensions of all operators, forms and spaces involved, we do not address this in detail here.

REMARK 18. The setup of the class of control problems to minimize (3.19) subject to (3.9) is different from the stochastic control problems considered in [GLL] and papers cited therein. There the Neumann boundary control is assumed to be deterministic independent of the parameters, and the expectation of the objective functional is minimized. Moreover, the number of stochastic parameters is assumed to be finite.
3.2. Dirichlet boundary control. The last and perhaps practically most relevant example of control problems with a tracking-type functional and a stationary PDE as constraint concerns problems with Dirichlet boundary control: minimize for some given data $y_{*}$ the quadratic functional

$$
\begin{equation*}
J(y, u)=\frac{1}{2}\left\|y-y_{*}\right\|_{O}^{2}+\frac{\omega}{2}\|u\|_{U}^{2} \tag{3.35}
\end{equation*}
$$

where state $y$ and control $u$ are coupled through the linear elliptic boundary value problem

$$
\begin{align*}
-\nabla_{x} \cdot\left(a(\sigma) \nabla_{x} y\right)+y & =f & & \text { in } \Omega, \\
y & =u & & \text { on } \Gamma  \tag{3.36}\\
\left(a(\sigma) \nabla_{x} y\right) \cdot \mathbf{n} & =0 & & \text { on } \partial \Omega \backslash \Gamma .
\end{align*}
$$

Here $\Gamma \subset \partial \Omega$ denotes the control boundary assumed to be a set of positive Lebesgue measure on which the control is exerted. Of course, we could allow again for an observation boundary and trace to this boundary in (3.35) as in Example 10, see $[\mathrm{K}]$. We dispense with this generalization here and choose for the following simply $O=H^{1}(\Omega)$ and given observation $y_{*} \in H^{1}(\Omega)$. To formulate (3.36) in weak form, we define $A(\sigma)$ as in (3.13), and set $Y=H^{1}(\Omega)$. It is because of the appearance of the control $u$ as a Dirichlet boundary condition in (3.36) that this is referred to as a Dirichlet boundary control problem. As it will be required to allow for repeated updates of the control, this suggests to formulate the constraints (3.36) weakly as a saddle point problem itself which results from appending the Dirichlet boundary conditions by Lagrange multipliers as follows. The trace operator to $\Gamma, \mathfrak{T}: H^{1}(\Omega) \rightarrow$ $H^{1 / 2}(\Gamma)$ is surjective and defines a bilinear form

$$
\begin{equation*}
\langle\mathfrak{T} v, w\rangle_{H^{1 / 2}(\Gamma) \times\left(H^{1 / 2}(\Gamma)\right)^{\prime}}=\langle\mathfrak{T} v, w\rangle_{H^{1 / 2}(\Gamma) \times\left(H^{1 / 2}(\Gamma)\right)^{\prime}} \tag{3.37}
\end{equation*}
$$

on $H^{1}(\Omega) \times\left(H^{1 / 2}(\Gamma)\right)^{\prime}$. Setting $Q:=\left(H^{1 / 2}(\Gamma)\right)^{\prime}$, the PDE constraint (3.36) can be formulated weakly as a linear saddle point problem: find $\left(y_{1}, y_{2}\right) \in Y \times Q$ such that

$$
\left(\begin{array}{cc}
A(\sigma) & \mathfrak{T}^{*}  \tag{3.38}\\
\mathfrak{T} & 0
\end{array}\right)\binom{y_{1}(\sigma)}{y_{2}(\sigma)}=\binom{f}{u(\sigma)}
$$

holds. The trace operator $\mathfrak{T}: Y \rightarrow Q$ is continuous and surjective on the kernel of $A(\sigma)$ yielding that the linear saddle point operator

$$
B(\sigma):=\left(\begin{array}{cc}
A(\sigma) & \mathfrak{T}^{*}  \tag{3.39}\\
\mathfrak{T} & 0
\end{array}\right): Y \times Q \rightarrow Y^{\prime} \times Q^{\prime}
$$

is an isomorphism and one has the norm equivalence

$$
\begin{equation*}
\left\|B(\sigma)\binom{v_{1}}{v_{2}}\right\|_{Y^{\prime} \times Q^{\prime}} \sim\left\|\binom{v_{1}}{v_{2}}\right\|_{Y \times Q} \tag{3.40}
\end{equation*}
$$

see, e.g., $[\mathrm{K}]$. Thus, if again $A(\sigma)$ satisfies Assumption 3, we have assured that the saddle point operator $B(\sigma)$ for the PDE constraint defined in (3.38) also satisfies Assumption 3. Finally, we choose for the control in (3.35) the natural space $U=$ $H^{1 / 2}(\Gamma)$. For the control problem to minimize (3.35) subject to (3.38), the optimality conditions, derived analogously as in Proposition 11 are now to find for given $f \in$
$Y^{\prime}, y_{*} \in Y$ the quintuple $\left(y_{1}(\sigma), y_{2}(\sigma), p_{1}(\sigma), p_{2}(\sigma), u(\sigma)\right) \in \boldsymbol{\mathcal { X }} \times \boldsymbol{\mathcal { X }} \times Q$ for $\boldsymbol{\mathcal { X }}:=Y \times Q$ such that
(DEE)

$$
\begin{aligned}
B(\sigma)\binom{y_{1}(\sigma)}{y_{2}(\sigma)} & =\binom{f}{u(\sigma)} \\
B(\sigma)^{*}\binom{p_{1}(\sigma)}{p_{2}(\sigma)} & =\binom{-R_{Y}\left(y(\sigma)-y_{*}\right)}{0} \\
\omega R_{U} u(\sigma) & =p_{2}(\sigma)
\end{aligned}
$$

where $R_{Y}$ is defined as in (3.18) and $R_{U}$ accordingly for $(\cdot, \cdot)_{U}$. Setting $\hat{y}_{*}:=R_{Y} y_{*} \in$ $Y^{\prime}$ and using the design equation in (DEE) to eliminate $p_{2}(\sigma)$, we arrive at the saddle point system of saddle point problems similar to (3.25), to solve for $y(\sigma), p(\sigma)):=$ $\left(y_{1}(\sigma), y_{2}(\sigma), p_{1}(\sigma), p_{2}(\sigma)\right) \in \mathcal{X} \times \mathcal{X}$ the system

$$
G(\sigma):=\left(\begin{array}{cc}
B(\sigma) & \left(\begin{array}{cc}
0 & 0 \\
0 & -\frac{1}{\omega} R_{U}^{-1}
\end{array}\right)  \tag{3.41}\\
\left(\begin{array}{cc}
R_{Y} & 0 \\
0 & 0
\end{array}\right) & B(\sigma)^{*}
\end{array}\right)\binom{y(\sigma)}{p(\sigma)}=\left(\begin{array}{c}
f \\
0 \\
\hat{y}_{*} \\
0
\end{array}\right)=: g
$$

Corollary 19. Together with Theorem 13, we have therefore established again the simultaneous analyticity of all the solution functions $y(\sigma), p(\sigma), u(\sigma)$ for the case that $A(\sigma)$ depends affinely on $\sigma$ according to (3.30). Moreover, applying again Theorem 7, we have established best $N$-term gpc approximation rates for the state, costate and control in the $L_{2}(\mathcal{S}, \rho ; \boldsymbol{\mathcal { X }})$ norm with the same rate $r$.
4. Parametric Linear-Quadratic Parabolic Control Problems. The preceding control problems were stationary, i.e., the equation of state was elliptic. We now show how control problems with parabolic equations of state fit into the the abstract results in Section 2. Accordingly, we introduce in the present section first an appropriate functional frame work for parabolic evolution problems, following [SS]. In view of Theorem 2, we verify in particular the stability conditions (2.10), (2.11) for the nominal parabolic operator $G_{0}$, in the corresponding spaces $X$ and $Y$ and establish its mapping properties and bounded invertibility. We then present examples of optimal control problems, following [GK].

The functional setting of the nominal problem is next used to formulate results for its parametric version and, in particular, for precise statements of smallness of perturbations. Sufficient conditions are once more given to cast the parametric parabolic control problem into the abstract theory of Section 2, implying in particular analytic dependence of state and controls on the parameter vector $\sigma$. Sufficient conditions on the perturbations to ensure best $N$-term convergence rates will be identified.
4.1. Space-Time Variational Formulations of Parabolic State Equations. Denote by $\Omega_{T}:=I \times \Omega$ with time interval $I:=(0, T)$ the time-space cylinder for functions $f=f(t, x)$ depending on time $t$ and space $x$. The parameter $T<\infty$ will always denote a finite time horizon. Let $Y$ be a dense subspace of $H:=L_{2}(\Omega)$ which is continuously embedded in $L_{2}(\Omega)$ and denote by $Y^{\prime}$ its topological dual. The associated dual form is denoted by $\langle\cdot, \cdot\rangle_{Y^{\prime} \times Y}$ or, shortly $\langle\cdot, \cdot\rangle$. Later we will use $\langle\cdot, \cdot\rangle$ also for duality pairings between function spaces on the time-space cylinder $\Omega_{T}$ with the precise meaning clear from the context. We consider abstract parabolic problems as developed, e.g., in [L, Chapter III, pp. 100]. Specifically, we assume given for a.e.
$t \in I$ and for $\sigma \in \mathcal{S}$ bilinear forms $a(\sigma, t ; \cdot, \cdot): Y \times Y \rightarrow \mathbb{R}$ so that $t \mapsto a_{0}(\sigma, t ; \cdot, \cdot)$ is measurable on $I$ and such that $a(\sigma, t ; \cdot, \cdot)$ is continuous and elliptic on $Y$, uniformly in $t \in I$ and in $\sigma \in \mathcal{S}$ : there exist constants $0<\alpha_{1} \leq \alpha_{2}<\infty$ independent of $t$ such that for a.e. $t \in I$ and for every $\sigma \in \mathcal{S}$

$$
\begin{align*}
|a(\sigma, t ; v, w)| & \leq \alpha_{2}\|v\|_{Y}\|w\|_{Y}, & & v, w \in Y  \tag{4.1}\\
a(\sigma, t ; v, v) & \geq \alpha_{1}\|v\|_{Y}^{2}, & & v \in Y
\end{align*}
$$

By the Riesz representation theorem, there exists a one-parameter family of bounded, linear operators $A(\sigma, t) \in \mathcal{L}\left(Y, Y^{\prime}\right)$ such that

$$
\begin{equation*}
\forall \sigma \in \mathcal{S}: \quad\langle A(\sigma, t) v, w\rangle:=a(\sigma, t ; v, w), \quad v, w \in Y \tag{4.2}
\end{equation*}
$$

Typically, $A(\sigma, t)$ will be a linear elliptic differential operator of second order on $\Omega$ and $Y$ will denote a function space on $\Omega$, such as, e.g., $Y=H_{0}^{1}(\Omega)$. We denote by $L_{2}(I ; Z)$ the space of all functions $v=v(t, x)$ for which for a.e. $t \in I$ one has $v(t, \cdot) \in Z$. Instead of $L_{2}(I ; Z)$, we will write this space as the (topological) tensor product of the two separable Hilbert spaces, $L_{2}(I) \otimes Z$, which, by [A, Theorem 12.6.1], can be identified.

For analytical purposes, linear parabolic evolution equations are often viewed as ordinary differential equations in $Y$ (see, e.g., [E]): given an initial condition $y_{0} \in H$ and right-hand side $f \in L_{2}\left(I ; Y^{\prime}\right)$, find $y(\sigma ; \cdot)$ in some function space on $\Omega_{T}$ such that

$$
\begin{align*}
\left\langle\frac{\partial y(\sigma ; t, \cdot)}{\partial t}, v\right\rangle+\langle A(\sigma, t) y(\sigma ; t, \cdot), v\rangle & =\langle f(t, \cdot), v\rangle & & \text { for all } v \in Y \text { and a.e. } t \in(0, T) \\
\langle y(0, \cdot), v\rangle & =\left\langle y_{0}, v\right\rangle & & \text { for all } v \in H \tag{4.3}
\end{align*}
$$

In order to cast such parabolic equations of state into the abstract setting of Section 2 and as basis for the recently developed space-time adaptive, compressive discretizations of such equations of state, however, space-time variational formulation for (4.3) are required. One such formulation is based on the Bochner type solution space

$$
\begin{align*}
\mathcal{X} & :=\left\{w \in L_{2}(I ; Y): \frac{\partial w(t, \cdot)}{\partial t} \in L_{2}\left(I ; Y^{\prime}\right)\right\}=L_{2}(I ; Y) \cap H^{1}\left(I ; Y^{\prime}\right) \\
& =\left(L_{2}(I) \otimes Y\right) \cap\left(H^{1}(I) \otimes Y^{\prime}\right) \tag{4.4}
\end{align*}
$$

equipped with the graph norm

$$
\begin{equation*}
\|w\|_{\mathcal{X}}^{2}:=\|w\|_{L_{2}(I ; Y)}^{2}+\left\|\frac{\partial w(t, \cdot)}{\partial t}\right\|_{L_{2}\left(I ; Y^{\prime}\right)}^{2} \tag{4.5}
\end{equation*}
$$

and the Bochner space of test functions

$$
\begin{equation*}
\mathcal{Y}:=L_{2}(I ; Y) \times H=\left(L_{2}(I) \otimes Y\right) \times H \tag{4.6}
\end{equation*}
$$

equipped, for $v=\left(v_{1}, v_{2}\right) \in \mathcal{Y}$, with the norm

$$
\begin{equation*}
\|v\|_{\mathcal{Y}}^{2}:=\left\|v_{1}\right\|_{L_{2}(I ; Y)}^{2}+\left\|v_{2}\right\|_{H}^{2} \tag{4.7}
\end{equation*}
$$

Note that $v_{1}=v_{1}(t, x)$ and $v_{2}=v_{2}(x)$. (We remark in passing that the choices (4.4) of spaces incorporates the initial condition as essential condition in the space; other possible formulations allow for the initial condition as natural condition, see [ChSt11] for details on such formulations which, in the present context of tracking type, high-dimensional parametric control problems, allow for completely analogous results).

Integration of (4.3) over $t \in I$ leads to the variational problem: given $f \in \mathcal{Y}^{\prime}$, for every $\sigma \in \mathcal{S}$ find a function $y(\sigma) \in \mathcal{Y}$

$$
\begin{equation*}
b(\sigma ; y(\sigma), v)=\langle f, v\rangle \quad \text { for all } v=\left(v_{1}, v_{2}\right) \in \mathcal{Y} \tag{4.8}
\end{equation*}
$$

where the bilinear form $b(\sigma ; \cdot, \cdot): \mathcal{X} \times \mathcal{Y} \rightarrow \mathbb{R}$ is defined by

$$
\begin{equation*}
b\left(\sigma ; w,\left(v_{1}, v_{2}\right)\right):=\int_{I}\left(\left\langle\frac{\partial w(t, \cdot)}{\partial t}, v_{1}(t, \cdot)\right\rangle+\left\langle A(\sigma, t) w(t, \cdot), v_{1}(t, \cdot)\right\rangle\right) d t+\left\langle w(0, \cdot), v_{2}\right\rangle \tag{4.9}
\end{equation*}
$$

and the right-hand side $\left\langle f+y_{0}, \cdot\right\rangle: \mathcal{Y} \rightarrow \mathbb{R}$ by

$$
\begin{equation*}
\left\langle f, v_{1}\right\rangle+\left\langle y_{0}, v_{2}\right\rangle:=\int_{I}\left\langle f(t, \cdot), v_{1}(t, \cdot)\right\rangle d t+\left\langle y_{0}, v_{2}\right\rangle \tag{4.10}
\end{equation*}
$$

for $v=\left(v_{1}, v_{2}\right) \in \mathcal{Y}$. It is well-known (see, e.g. [DL, Chapter XVIII, Sect. 3]) that the parametric operator family $\{B(\sigma): \sigma \in \mathcal{S}\}$ defined by the bilinear form $b(\sigma ; \cdot, \cdot)$ in (4.9) is a family of isomorphisms from $\mathcal{X}$ to $\mathcal{Y}^{\prime}$. In [SS], detailed bounds on the norms of the operator and its inverse were established. To prepare the ensuing formulation and regularity results on the parametric parabolic optimal control problem, we next formulate the corresponding result for the state equation (4.8). This result is again a special case of the abstract results, Theorem 4 and Theorem 7. Alternatively, it could be inferred from the abstract theory of parabolic evolution equations in [PS], subject to a requirement of continuity of $A(\sigma, t)$ with respect to $t \in[0, T]$, uniformly with respect to $\sigma \in \mathcal{S}$.

Theorem 20. Assume that the parametric family $\left\{A(\sigma, t) \in \mathcal{L}\left(Y, Y^{\prime}\right): \sigma \in\right.$ $\mathcal{S}, t \in I\}$ satisfies Assumption 1 with $\mathcal{X}=\mathcal{Y}=Y$, uniformly for $t \in I$, i.e., $A(\sigma, t)$ is boundedly invertible with uniform (w.r. to $t \in I$ and $\sigma \in \mathcal{S}$ ) bound $C_{0}$ and there exists a sequence $b \in \ell^{\mathfrak{p}}(\mathbb{N})$ for some $0<\mathfrak{p} \leq 1$ such that

$$
\forall \nu \in \mathfrak{F}: \quad \sup _{t \in I} \sup _{\sigma \in \mathcal{S}}\left\|(A(0, t))^{-1}\left(\partial_{\sigma}^{\nu} A(\sigma, t)\right)\right\|_{\mathcal{L}(Y, Y)} \leq C_{0} b^{\nu}
$$

Then, for every $\sigma \in \mathcal{S}$, the parabolic evolution operator $B(\sigma) \in \mathcal{L}\left(\mathcal{X}, \mathcal{Y}^{\prime}\right)$ defined by $\langle B(\sigma) w, v\rangle:=b(\sigma ; w, v)$ for $w \in \mathcal{X}$ and for $v \in \mathcal{Y}$ with the parametric bilinear form $b(\sigma ; \cdot, \cdot)$ from (4.9) and with the choice of spaces $\mathcal{X}, \mathcal{Y}$ as in (4.4) and (4.6) is boundedly invertible: there exist constants $0<\beta_{1} \leq \beta_{2}<\infty$ such that

$$
\begin{equation*}
\sup _{\sigma \in \mathcal{S}}\|B(\sigma)\|_{\mathcal{X} \rightarrow \mathcal{Y}^{\prime}} \leq \beta_{2} \quad \text { and } \quad\left\|B(\sigma)^{-1}\right\|_{\mathcal{Y}^{\prime} \rightarrow \mathcal{X}} \leq \frac{1}{\beta_{1}} \tag{4.11}
\end{equation*}
$$

Moreoever, the parametric operator family $\{B(\sigma): \sigma \in \mathcal{S}\}$ satisfies Assumption 1. In particular, the parametric family $y(\sigma)$ in (4.8) of states satisfies the a-priori estimate

$$
\begin{equation*}
\forall \nu \in \mathfrak{F}: \quad \sup _{\sigma \in \mathbb{S}}\left\|\left(\partial_{\sigma}^{\nu} y\right)(\sigma)\right\|_{\mathcal{X}} \leq C_{0}\|f\|_{\mathcal{Y}^{\prime}}|\nu|!\tilde{b}^{\nu} \tag{4.12}
\end{equation*}
$$

and admits a Legendre expansion

$$
\begin{equation*}
y(\sigma)=\sum_{\nu \in \mathfrak{F}} y_{\nu}(\sigma) L_{\nu}(\sigma), \quad y_{\nu}=\int_{\sigma \in \mathcal{S}} y(\sigma) L_{\nu}(\sigma) \rho(d \sigma) \tag{4.13}
\end{equation*}
$$

which converges in $L^{2}(\mathcal{S}, \rho ; \mathcal{X})$. Morever, $\left(\left\|y_{\nu}\right\|_{\mathcal{X}}\right)_{\nu \in \mathfrak{F}} \in \ell^{\mathfrak{p}}(\mathfrak{F})$ and best $N$-term truncated Legendre expansions converge at rate $N^{-(1 / \mathfrak{p}-1 / 2)}$.

Proof. As proved in [SS], for every $\sigma \in \mathcal{S}$ the continuity constant $\beta_{2}$ and the inf-sup condition constant $\beta_{1}$ for $b(\sigma ; \cdot, \cdot)$ are independent of $\sigma \in \mathcal{S}$ and satisfy

$$
\begin{equation*}
\beta_{1} \geq \frac{\min \left(\alpha_{1} \alpha_{2}^{-2}, \alpha_{1}\right)}{\sqrt{2 \max \left(\alpha_{1}^{-2}, 1\right)+\varrho^{2}}}, \quad \beta_{2} \leq \sqrt{2 \max \left(1, \alpha_{2}^{2}\right)+\varrho^{2}} \tag{4.14}
\end{equation*}
$$

where $\alpha_{1}, \alpha_{2}$ are the constants from (4.1) bounding $A(\sigma, t)$ and $\varrho$ is defined as

$$
\begin{equation*}
\varrho:=\sup _{0 \neq w \in \mathcal{Y}} \frac{\|w(0, \cdot)\|_{H}}{\|w\|_{\mathcal{Y}}} . \tag{4.15}
\end{equation*}
$$

We like to recall from [DL, E] that $\mathcal{Y}$ is continuously embedded in $\mathcal{C}^{0}(I ; H)$ so that the pointwise in time initial condition in (4.3) is well-defined. From this it follows that the constant $\varrho$ is bounded uniformly for the choice of $\mathcal{Y} \hookrightarrow H$.

In the sequel, we will require the dual operator $B(\sigma)^{*}: \mathcal{Y} \rightarrow \mathcal{X}^{\prime}$ of $B(\sigma)$ which is defined formally by

$$
\begin{equation*}
\forall \sigma \in \mathcal{S}: \quad\langle B(\sigma) w, v\rangle=:\left\langle w, B(\sigma)^{*} v\right\rangle \tag{4.16}
\end{equation*}
$$

From the definition of the bilinear form (4.9) on $\mathcal{X} \times \mathcal{Y}$, it follows by integration by parts for the first term with respect to time and using the adjoint $A(\sigma, t)^{*}$ with respect to the duality pairing $Y^{\prime} \times Y$ that

$$
\begin{align*}
& b\left(\sigma ; w,\left(v_{1}, v_{2}\right)\right)= \int_{I}\left(\left\langle w(t, \cdot), \frac{\partial v_{1}(t, \cdot)}{\partial t}\right\rangle+\left\langle w(t, \cdot), A(\sigma, t)^{*} v_{1}(t, \cdot)\right\rangle\right) d t \\
& \quad+\left\langle w(0, \cdot), v_{2}\right\rangle+\left.\left\langle w(t, \cdot), v_{2}\right\rangle\right|_{0} ^{T} \\
&= \int_{I}\left(\left\langle w(t, \cdot), \frac{\partial v_{1}(t, \cdot)}{\partial t}\right\rangle+\left\langle w(t, \cdot), A(\sigma, t)^{*} v_{1}(t, \cdot)\right\rangle\right) d t \\
& \quad+\left\langle w(T, \cdot), v_{2}\right\rangle \\
&=:\left\langle w, B(\sigma)^{*} v\right\rangle . \tag{4.17}
\end{align*}
$$

Note that the first term of the right-hand side which involves $\frac{\partial}{\partial t} v_{1}(t, \cdot)$ is still welldefined with respect to $t$ as an element of $\mathcal{Y}^{\prime}$ on account of $w \in \mathcal{Y}$.

So far, we considered only the parabolic state equation and proved analyticity and polynomial approximation rates.

We now turn to perturbed, parametric state equations resulting from parametric uncertainty in the spatial operator $A(\sigma, t)$, and present in particular sufficient conditions on the perturbations of $A_{0}(t)$ in order for the perturbed state equation to fit into the general Assumption 2 and Theorem 2.
4.2. Tracking-type control problem constrained by a parametric, parabolic PDE. Recalling the situation from [GK], we wish to minimize, for some given target state $y_{*}$ and fixed end time $T>0$, the quadratic functional

$$
\begin{equation*}
J(y, u):=\frac{\omega_{1}}{2}\left\|y-y_{*}\right\|_{L_{2}(I ; O)}^{2}+\frac{\omega_{2}}{2}\left\|y(T, \cdot)-y_{*}(T, \cdot)\right\|_{O}^{2}+\frac{\omega_{3}}{2}\|u\|_{L_{2}(I ; U)}^{2} \tag{4.18}
\end{equation*}
$$

over the state $y(\sigma)=y(\sigma ; t, x)$ and over the control $u(\sigma)=u(\sigma ; t, x)$ subject to

$$
\begin{equation*}
B(\sigma) y(\sigma)=\mathfrak{E} u(\sigma)+\binom{f}{y_{0}} \quad \text { in } \mathcal{Y}^{\prime} \tag{4.19}
\end{equation*}
$$

where $B(\sigma)$ denotes the parametric, parabolic evolution operator defined by Theorem 20 and where $f \in \mathcal{Y}^{\prime}$ is given by (4.10). In (4.18), the real weight parameters $\omega_{1}, \omega_{2} \geq$ 0 are such that $\omega_{1}+\omega_{2}>0$ and $\omega_{3}>0$. The space $O$ by which the integral over $\Omega$ in the first two terms in (4.18) is indexed is to satisfy $O \supseteq Y$ with continuous embedding. Although there is in the wavelet framework great flexibility in choosing even fractional Sobolev spaces for $O$, for better readibility, we pick here $O=Y$. Moreover, in a general case we suppose that the operator $\mathfrak{E}$ is a linear operator $\mathfrak{E}: U \rightarrow \mathcal{Y}^{\prime}$ extending $\int_{I}\left\langle u(t, \cdot), v_{1}(t, \cdot)\right\rangle d t$ trivially, i.e., $\mathfrak{E} \equiv(I, 0)^{\top}$. For ease of presentation in the current setting, we choose again $U=Y^{\prime}$ similar to the stationary case in Section 3.1.

The tracking type control problem consists in minimizing the functional (4.18) subject to the parametric parabolic equation of state (4.19). We recall that the Riesz operator $R_{Y}: Y \rightarrow Y^{\prime}$ defined by

$$
\begin{equation*}
(v, z)_{Y}=:\left\langle v, R_{Y} z\right\rangle, \quad v, z \in Y \tag{4.20}
\end{equation*}
$$

maps $Y$ boundedly invertibly onto its dual $Y^{\prime}$. Since here $R_{U}=R_{Y}^{-1}$ as in Section 3.1, we write $R=R_{Y}$.

Analogously to the derivation of the system (EE) in Section 3.1, we can derive the first order necessary conditions consisting of the primal system together with the costate or adjoint equations and the design equation. For a unification of notation, it will be useful to define

$$
y_{1}(\sigma):=y(\sigma), \quad y_{2}(\sigma):=y(\sigma ; 0)
$$

and, since the adjoint state also requires the state to be evaluated at the finite end point (sometimes also denoted as finite horizon) $T, y_{3}(\sigma):=y(\sigma ; T)$. Then the necessary conditions for optimality read:
find the solution tuple $\left(y_{1}(\sigma), y_{2}(\sigma), y_{3}(\sigma), p_{1}(\sigma), p_{2}(\sigma), u(\sigma)\right) \in \mathcal{X} \times Y \times Y \times \mathcal{X} \times Y \times \mathcal{Y}^{\prime}$ as

$$
\begin{align*}
B_{1}(\sigma) y_{1}(\sigma) & =u(\sigma)+f \\
B_{2}(\sigma) y_{2}(\sigma) & =y_{0} \\
B_{1}(\sigma)^{*} p_{1}(\sigma)+\omega_{1} R_{Y} y_{1}(\sigma) & =\omega_{1} R_{Y} y_{*}  \tag{4.21}\\
B_{2}(\sigma)^{*} p_{2}(\sigma)+\omega_{2} R_{Y} y_{3}(\sigma) & =\omega_{2} R_{Y} y_{*}(T) \\
\omega_{3} u(\sigma) & =R_{Y} p_{1}(\sigma) .
\end{align*}
$$

Here $B_{1}(\sigma), B_{2}(\sigma)$ are the linear operators defined by the first and second dual forms in (4.9), respectively, with 'dual' $B_{1}(\sigma)^{*}, B_{2}(\sigma)^{*}$ defined according to (4.17). Note that the appearance of the Lagrange multipliers $p_{1}(\sigma), p_{2}(\sigma)$ is caused by appending the parabolic constraints (4.19) to the functional (4.18). Thus, the variable $p_{1}(\sigma)$ is the adjoint state $p_{1}(\sigma)=p(\sigma ; t, x)$, and $p_{2}(\sigma)$ may be interpreted as evaluating $p$ at the end point $T$, i.e., $p_{2}(\sigma)=p(\sigma ; T, x)$. For presentation purposes, we also define $p_{3}(\sigma)=p(\sigma ; 0, x)$. Eliminating $u(\sigma)=\omega_{3}^{-1} R_{Y} p_{1}(\sigma)$ from the design equation and abbreviating

$$
\hat{y}_{*}:=R_{Y} y_{*} \quad \text { and } \quad \hat{y}_{*}(T):=R_{Y} y_{*}(T),
$$

and

$$
\left.\widehat{y(\sigma)}=\left(y_{1}(\sigma), y_{2}(\sigma), y_{3}(\sigma)\right), \quad \widehat{p(\sigma)}\right)=\left(p_{1}(\sigma), p_{2}(\sigma), p_{3}(\sigma)\right)
$$

we arrive at the coupled system

$$
\begin{align*}
G(\sigma)\binom{\widehat{y(\sigma)}}{\widehat{p(\sigma)}} & :=\left(\begin{array}{cc}
\widehat{B(\sigma)} & \operatorname{diag}\left(-\frac{1}{\omega_{3}} R_{Y}, 0,0\right) \\
\left(\begin{array}{ccc}
\frac{1}{\omega_{1}} R_{Y} & 0 & 0 \\
0 & 0 & \frac{1}{\omega_{2}} R_{Y} \\
0 & 0 & 0
\end{array}\right) \quad \widehat{B(\sigma)}^{*}
\end{array}\right)\left(\begin{array}{c}
f\left(\frac{y(\sigma)}{p(\sigma)}\right)
\end{array}\right)  \tag{4.22}\\
& =\left(\begin{array}{c}
f \\
y_{0} \\
0 \\
\omega_{1} \hat{y}_{*} \\
\omega_{2} \hat{y}_{*}(T) \\
0
\end{array}\right)=: g
\end{align*}
$$

where $\widehat{B(\sigma)}:=\operatorname{diag}\left(B_{1}(\sigma), B_{2}(\sigma), 0\right)$.
THEOREM 21. Let for every $t \in[0, T]$ the parametric family of spatial operators $\left\{A(\sigma, t) \in \mathcal{L}\left(Y, Y^{\prime}\right): \sigma \in \mathcal{S}\right\}$ satisfies Assumption 1. Then for every $\sigma \in \mathcal{S}$, the tracking type control problem (4.21) can be written as a parametric saddle point operator equation $G(\sigma)(\widehat{y(\sigma)}, \widehat{p(\sigma)})=g$ for the solution tuple $(\widehat{y(\sigma)}, \widehat{p(\sigma)}) \in \boldsymbol{\mathcal { X }}$ with $G(\sigma) \in \mathcal{L}\left(\mathcal{X}, \mathcal{Y}^{\prime}\right)$ where the space $\mathcal{X}=\mathcal{Y}$ is given by

$$
\begin{equation*}
\mathcal{X}=\mathcal{X} \times Y \times Y \tag{4.23}
\end{equation*}
$$

Moreover, for $\omega_{1}+\omega_{2}>0$, $\omega_{3}>0$, the parametric saddle point operator $\mathbf{G}(\sigma) \in$ $\mathcal{L}(\mathcal{X}, \mathcal{Y})$ in (4.22) is boundedly invertible for all $\sigma \in \mathcal{S}$ and satisfies Assumption 1. The parametric family of state-control pairs $\mathcal{S} \ni \sigma \mapsto\left(\frac{\overline{y(\sigma)}}{p(\sigma)}\right) \in L^{2}(\mathcal{S}, \rho ; \boldsymbol{\mathcal { X }})$ depends analytically on $\sigma \in \mathcal{S}$ and admits a concurrent Legendre expansion

$$
\begin{equation*}
\left(\frac{\widehat{y(\sigma)}}{\overrightarrow{p(\sigma)}}\right)=\sum_{\nu \in \mathfrak{F}} L_{\nu}(\sigma)\binom{y_{\nu}}{p_{\nu}}, \quad\binom{y_{\nu}}{p_{\nu}} \in \mathcal{X} \tag{4.24}
\end{equation*}
$$

Furthermore, the parametric Legendre expansion is sparse in the sense that, if the elliptic operator family $A(\sigma, t)$ is a uniformly in $t \in[0, T] \mathfrak{p}$-analytic operator family in the sense that, for every fixed $t \in[0, T]$, Assumption 1 and, in particular, (2.8), holds, then the coefficient sequence in (4.24) is $\mathfrak{p}$-summable. This means that

$$
\left(\left\|\binom{y_{\nu}}{p_{\nu}}\right\|_{\mathcal{X}}\right)_{\nu \in \mathfrak{F}} \in \ell^{\mathfrak{p}}(\mathfrak{F})
$$

for the same value of $\mathfrak{p}$ and, for every $N \in \mathbb{N}$, there exists an index set $\Lambda \subset \mathfrak{F}$ of cardinality not exceeding $N$ such that the $N$-term truncated Legendre expansion

$$
\binom{y_{N}(\sigma)}{p_{N}(\sigma)}:=\sum_{\nu \in \Lambda} L_{\nu}(\sigma)\binom{y_{\nu}}{p_{\nu}}, \quad\binom{y_{\nu}}{p_{\nu}} \in \mathcal{X}
$$

approximates concurrently the state and the control on the entire parameter domain $\mathcal{S}$ at rate $N^{-(1 / \mathfrak{p}-1 / 2)}$ in $L^{2}(\mathcal{S}, \rho ; \mathcal{X})$.
5. Conclusion. We have proved, for control problems constrained by linear elliptic and parabolic PDEs which depend on possibly countably infinitely many paramters, analytic parameter dependence of the state, co-state and of the control. The parameter dependence was allowed to be more general than affine. The particular case of affine dependence arises, for example, in state equations with random coefficients which are parametrized in terms of Karhunen-Loève expansions as in [ST]. We have quantified the analytic dependence of (co)state and control. Specifically, we established that these quantities allow expansions in terms of tensorized "polynomial chaos" type bases which are sparse, their sparsity being quantified in terms of $\mathfrak{p}$-summability of the coefficient sequences. This sparsity result is the analytical foundation for the development of sparse tensor discretizations of these problems which allow adaptive Galerkin approximations of (co)state and control on the entire (possibly infinite-dimensional) parameter space, following [G], combined with appropriate discretizations in space and time, following [DK, GK, K]. Details of this will be reported in $[\mathrm{KS}]$.

## REFERENCES

[A] J.P. Aubin, Applied Functional Analysis, 2nd ed., Wiley, 2000.
[BF] F. Brezzi and M. Fortin, Mixed and Hybrid Finite Element Methods, Springer Verlag, Berlin, 1991.
[ChSt11] N. Chegini and R. Stevenson, Adaptive wavelet schemes for parabolic problems: sparse matrices and numerical results, SIAM J. Numer. Anal. 49 (2011), 182-212.
[CDS1] A. Cohen, R. DeVore and Ch. Schwab, Convergence rates of best $N$-term Galerkin approximation for a class of elliptic sPDEs, Found. Comput. Math. 10 (6), 2010, 615-646.
[CDS2] A. Cohen, R. DeVore and Ch. Schwab, Analytic regularity and polynomial approximation of parametric and stochastic elliptic PDEs, Analysis and Applications 9 (1), 2011, 1-37.
[DK] W. Dahmen and A. Kunoth, Adaptive wavelet methods for linear-quadratic elliptic control problems: Convergence rates, SIAM J. Contr. Optim. 43(5) (2005) 1640-1675.
[DL] R. Dautray and J.-L. Lions, Mathematical Analysis and Numerical Methods for Science and Technology, Vol.5: Evolution Problems, Springer, 2000.
[DeV] R. DeVore, Nonlinear approximation, Acta Numerica 7 (1998), 51-150.
[DeVK] R. DeVore, A. Kunoth (Eds.), Multiscale, Nonlinear and Adaptive Approximation, Springer, 2009.
[E] L.C. Evans, Partial Differential Equations, AMS, Providence, RI, 1998.
[G] C.J. Gittelson, Adaptive stochastic Galerkin methods: Beyond the elliptic case. Report 2011-12, Seminar for Applied Mathematics, ETH Zürich (in review).
[GK] M.D. Gunzburger and A. Kunoth, Space-time adaptive wavelet methods for control problems constrained by parabolic evolution equations, SIAM J. Contr. Optim. 49(3) (2011) 11501170.
[GLL] M.D. Gunzburger, H.-C. Lee and J. Lee, Error estimates of stochastic optimal Neumann boundary control problems, SIAM J. Numer. Anal. 49(4) (2011), 1532-1552.
[K] A. Kunoth, Adaptive wavelet schemes for an elliptic control problem with Dirichlet boundary control, Numer. Algor. 39 1-3 (2005), 199-220.
[KS] A. Kunoth, and Ch. Schwab, Sparse, adaptive tensor Galerkin approximations of stochastic control problems. (In preparation 2012).
[L] J.L. Lions, Optimal Control of Systems Governed by Partial Differential Equations, Springer, Berlin, 1971.
[PS] J. Prüss and R. Schnaubelt, Solvability and maximal regularity of parabolic evolution equations with coefficients continuous in time, J. Math. Anal. Appl., 256 (2001) 405-430.
[SG11] Ch. Schwab and C.J. Gittelson, Sparse tensor discretizations of high-dimensional parametric and stochastic PDEs, Acta Numerica 20 (2011), Cambridge University Press, Cambridge, UK.
[SS] Ch. Schwab and R. Stevenson, Space-time adaptive wavelet methods for parabolic evolution equations, Math. Comp. 78, (2009) 1293-1318.
[ST] Ch. Schwab and R.A. Todor, Karhúnen-Loève approximation of random fields by generalized fast multipole methods, Journal of Computational Physics 217 (2006), 100-122.

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    ${ }^{\dagger}$ Institut für Mathematik, Universität Paderborn, Warburger Str. 100, 33098 Paderborn, Germany

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    ${ }^{\dagger}$ Institut für Mathematik, Universität Paderborn, Warburger Str. 100, 33098 Paderborn, Germany, kunoth@math.uni-paderborn.de,
    www2.math.uni-paderborn.de/ags/kunoth/group/angelakunoth.html
    $\ddagger$ Seminar für Angewandte Mathematik, ETH Zürich, ETHZ HG G57.1, CH 8092 Zürich, Switzerland

