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# GREEDY SEARCH OF OPTIMAL APPROXIMATE SOLUTIONS 

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Dedicated to Professor Ronald DeVore on the occasion of his 80th birthday.


#### Abstract

In this paper we develop a procedure to deal with a family of parameter-dependent ill-posed problems, for which the exact solution in general does not exist. The original problems are relaxed by considering corresponding approximate ones, whose optimal solutions are well defined, where the optimality is determined by the minimal norm requirement. The procedure is based upon greedy algorithms that preserve, at least asymptotically, Kolmogorov approximation rates. In order to provide a-priori estimates for the algorithm, a Tychonoff-type regularization is applied, which adds an additional parameter to the model. The theory is developed in an abstract theoretical framework that allows its application to different kinds of problems. We present a specific example that considers a family of ill-posed elliptic problems. The required general assumptions in this case translate to rather natural uniform lower and upper bounds on coefficients of the considered operators.


## 1. Introduction

Greedy theory, inspired on the notion of Kolmogorov complexity and nonlinear approximation theory, has been extensively developed to provide optimal approximation rates for parameter depending problems, in particular Partial Differential Equations (PDE) $[5,18]$.

This is a relevant subject since, often in applications, the model mimicking the dynamics is not fully known and is subject to uncertainty, in particular, on some of the relevant parameters entering in the system such as diffusivity, Lamé coefficients etc.

A control theoretical counterpart was developed in [12]. The question addressed in that article referees to the classical problem of controllability, that of driving a dynamical system to a desired final configuration by the action of a suitable control. They did it in the context of parameter-depending problems, building a greedy algorithm allowing to guarantee an optimal approximation, in the sense of the Kolmogorov thickness, of the set of parameter-depending controls.

But, as observed in [12], an important case was left open. How can the greedy strategy be adapted when the system under consideration is not controllable, like it occurs for instance in parabolic PDEs in which, due to the strong time irreversibility of the model, only very smooth targets are reachable?

[^0]In this paper we consider these kind of problems in an abstract frame, motivated by that example, and characterized by the application of the greedy algorithm to linear systems in which the operators governing the system are not onto.

Inspired by the theory of approximate controllability for time-irreversible PDEs, developed by J. L. Lions, and by R. Glowinski and J. L. Lions in the numerical setting (e.g. $[3,14]$ ), and the classical technique of Tychonoff regularization for illposed (inverse) problems, we adopt a two-folded perspective that allows us to end up developing a greedy strategy ensuring that Kolmogorov complexity is reached, for quasi-solutions (those that assure the fulfillment of the system up to an $\varepsilon$ error) in those cases where the exact solution (corresponding to $\varepsilon$ equal to zero) does not exist in the given functional setting.

With that purpose we proceed in several steps:

1. We introduce a concept of optimal $\varepsilon$-solution, reminiscent of the theory of approximate controllability, that allows characterizing the solution of minimal norm assuring that the system is solved up to an $\varepsilon$ tolerance or error.
2. We link this $\varepsilon$-solution with a suitable Tychonoff regularization, in which the $\varepsilon$-error or tolerance, can be linked to the $\eta$-Tychonoff regularization parameter through a suitable nonlinear implicit equation.
3. We then adapt and apply the existing greedy methods for well posed parameter dependent problems, but in the context in which the number of free parameters is increased by one, to incorporate the $\eta$-Tychonoff parameter, that adds to the physical parameters on which the model depends.
4. We then use the nonlinear link of the $\varepsilon$-error and $\eta$-Tychonoff parameter to derive a greedy approximation result for the $\varepsilon$-problem.

Our abstract results apply to a wide variety of problems, such as deconvolution in image processing, time inversion of highly irreversible systems, like heat equations, or the approximate controllability of PDE.

Although we develop the theory in the context of bounded linear operators, it can also be extended to unbounded ones. This, in turn, allows us to study illposed elliptic problems in which the (exact) solution does in general not exist for an arbitrary right-hand side. In this way, we generalize the setting considered by R. DeVore and his collaborators, in which they consider a family of well-posed elliptic problems and develop the greedy algorithms to recover the corresponding solutions $[5,6]$.

The paper is organised as follows. The next section provides definition and analysis of optimal approximative solutions of a single $\varepsilon$-problem. Section 3 is devoted to construction of a greedy algorithm for solving a family of parameter dependent problems. Generalisation and application of the theory to unbounded operators is presented in Section 4, followed by a particular example related to the elliptic equation. The paper is closed by some concluding remarks and perspectives for further research.

## 2. Preliminaries on optimal approximate solutions of linear systems

Let $L$ be a bounded operator from $L(X, Y)$ with a dense range, where $X$ and $Y$ are infinite-dimensional, real Hilbert spaces. We consider the problem of solving
the system

$$
\begin{equation*}
L u=f \tag{2.1}
\end{equation*}
$$

for a given vector $f \in Y$. Of course, due to the fact that image of $L$ is only dense in $Y$, the above problem has no solution for every $f$. Thus we relax the problem and consider the approximate problem of finding $u$ such that

$$
\begin{equation*}
\|L u-f\|_{Y} \leq \varepsilon \tag{2.2}
\end{equation*}
$$

for some a-priori given $\varepsilon>0$.
Having assumed that the operator $L$ has a dense range, the set of approximate solutions satisfying (2.2) for any fixed $f$ is neither empty nor a singleton. Thus it is reasonable to choose the one which is optimal in some sense. We set the optimization criteria in this paper as the one of minimal norm. This leads us to the following problem

$$
\begin{equation*}
\min _{u \in X}\left\{\frac{1}{2}\|u\|_{X}^{2}: \quad\|L u-f\|_{Y} \leq \varepsilon\right\} \tag{2.3}
\end{equation*}
$$

whose solution we denote by $\tilde{u}$.
By exploring Fenchel-Rockafellar duality techniques (e.g. [15, §3.6]) the solution $\tilde{u}$ of the problem (2.3) can be obtained by solving the corresponding dual problem. This allows one to replace the original constrained problem by a non-constrained one. More precisely, the following theorem holds.

Proposition 2.1 ([13, Propositions $2 \& 4]$ ). The solution of problem (2.3) is given by

$$
\tilde{u}=L^{*} \tilde{v}
$$

where $\tilde{v} \in Y$ minimizes the dual problem

$$
\begin{equation*}
J(v)=\frac{1}{2}\left\|L^{*} v\right\|_{X}^{2}+\varepsilon\|v\|_{Y}-\langle f \mid v\rangle_{Y}, \quad v \in Y \tag{2.4}
\end{equation*}
$$

Here $L^{*} \in L(Y, X)$ stands for the adjoint operator of $L$, while $\langle\cdot \mid \cdot\rangle_{Y}$ denotes the scalar product in $Y$.

Note that the last proposition does not guarantee uniqueness of the solution to the dual problem (2.4). The latter includes the functional $J$ which is a non-standard one since it involves a non-smooth term of homogeneity one. Being strictly convex and continuous, it attains its minimal value at the unique point if it is coercive. According to the density assumption on $L$ and Hahn-Banach theorem, the adjoint operator $L^{*}$ is injective. But in general it is not coercive (this would correspond to the case of $L$ having a full rank). However, the $\varepsilon$ term entering the functional $J$ ensures its coercivity as shown by the following result.

Proposition 2.2. Functional $J$ defined by (2.4) is coercive, i.e. it satisfies

$$
\liminf _{\|v\|_{Y} \rightarrow \infty} \frac{J(v)}{\|v\|_{Y}} \geq \varepsilon
$$

Proof. The proof essentially follows the lines of the [9, Proposition 2.1] which treats the special case of $L$ being the heat operator.

We argue by contradiction. We suppose there exist a sequence $\left(v_{n}\right)$ such that $\left\|v_{n}\right\|_{Y} \rightarrow \infty$ and

$$
\begin{equation*}
\liminf _{n} \frac{J\left(v_{n}\right)}{\left\|v_{n}\right\|_{Y}}<\varepsilon \tag{2.5}
\end{equation*}
$$

Denoting by $v_{n}^{0}=v_{n} /\left\|v_{n}\right\|_{Y}$ the corresponding normalized vectors, it follows that, up to a subsequence, $\left(v_{n}^{0}\right)$ converges weakly to some $v^{0} \in Y$. From (2.5) it follows

$$
\varepsilon>\frac{1}{2}\left\|v_{n}\right\|_{Y}\left\|L^{*} v_{n}^{0}\right\|_{X}^{2}+\varepsilon-\left\langle f \mid v_{n}^{0}\right\rangle_{Y} .
$$

As all the terms in the last relation, except $\left\|v_{n}\right\|_{Y}$, are bounded, it implies that $\left\|L^{*} v_{n}^{0}\right\|_{X} \rightarrow 0$. Consequently, $\left\|L^{*} v^{0}\right\|_{X}=0$, and the injectivity of $L^{*}$ implies $v^{0}=0$.

Thus we obtain

$$
\liminf _{n} \frac{J\left(v_{n}\right)}{\left\|v_{n}\right\|_{Y}} \geq \varepsilon-\lim _{n}\left\langle f \mid v_{n}^{0}\right\rangle_{Y}=\varepsilon
$$

which contradicts the initial assumption (2.5).
A detailed characterization of the unique minimizer of $J$ is provided by the next result.

Proposition 2.3. The minimizer $\tilde{v}$ equals zero if and only if $\|f\|_{Y} \leq \varepsilon$. Otherwise, it satisfies the Euler-Lagrange equation which has the form

$$
\begin{equation*}
L L^{*} \tilde{v}+\varepsilon \frac{\tilde{v}}{\|\tilde{v}\|_{Y}}-f=0 \tag{2.6}
\end{equation*}
$$

Proof. From the very definition (2.4) of functional $J$ it follows

$$
J(v) \geq \frac{1}{2}\left\|L^{*} v\right\|_{X}^{2}+\left(\varepsilon-\|f\|_{Y}\right)\|v\|_{Y} .
$$

Consequently, if $\|f\|_{Y} \leq \varepsilon$ the functional $J$ is nonnegative, obtaining its minimum for $\tilde{v}=0$.

Otherwise, the functional $J$ attains negative values as well. Indeed, take a sequence $v_{n}=f / n$ and calculate

$$
J\left(v_{n}\right)=\frac{1}{2 n^{2}}\left\|L^{*} f\right\|_{X}^{2}+\frac{1}{n}\left(\varepsilon\|f\|_{Y}-\|f\|_{Y}^{2}\right)
$$

For $n$ large enough, the right hand side of the last expression is dominated by the term $\|f\|_{Y}\left(\varepsilon-\|f\|_{Y}\right) / n<0$, and results in negative values of $J\left(v_{n}\right)$. Consequently the minimum $\tilde{v}$ differs from 0 .

As the functional $J$ is differentiable apart from the origin, its differential at $\tilde{v}$ equals zero, which results in the Euler-Lagrange equation (2.6).

The case $\|f\|_{Y} \leq \varepsilon$ is a trivial one, which we exclude from further analysis. In the non-trivial case $\left(\|f\|_{Y} \leq \varepsilon\right)$ the Euler-Lagrange equation implies that $\tilde{u}=L^{*} \tilde{v}$ is an eligible solution to the problem (2.2) which brings the system to the boundary of the target ball around $f$. Moreover, as already stated above, it is also the solution of the minimal norm among all $u \in X$ satisfying (2.2).

In this way, finding optimal approximative solution to (2.2) is equivalent to finding the minimizer of the functional $J$ given by (2.4), which, in turn, is equivalent to
solving the corresponding Euler-Lagrange equation (2.6). For these reasons, from now on, we shall concentrate on finding efficient methods for solving the latter equation.

Remark 2.4. It is interesting to note that problem (2.2) also allows for the finite dimensional solvability. More precisely, one can show that for any finite-dimensional subspace $E \subseteq Y$ and any target $f \in E$ there exists an approximative solution $u_{E}$ to (2.2) such that

$$
\begin{equation*}
P_{E}\left(L u_{E}\right)=P_{E}(f) \tag{2.7}
\end{equation*}
$$

where $P_{E}$ denotes the orthogonal projection to the space $E$. similarly, by $P_{E^{\perp}}=$ $I_{Y}-P_{E}$ we denote the projection to the orhogonal compliment of $E$.

Such result requires analysis of the functional

$$
J_{E}(v)=\frac{1}{2}\left\|L^{*} v\right\|_{X}^{2}+\varepsilon\left\|P_{E^{\perp}}(v)\right\|_{Y}-\langle f \mid v\rangle_{Y}, \quad v \in Y
$$

Its coercivity is proven in the manner similar to the proof of Lemma 2.2, which ensures existence of its unique minimizer $\tilde{v}_{E}$. Following the ideas of [20], one can show that $\tilde{u}_{E}=L^{*} \tilde{v}_{E}$ is the solution that satisfies both (2.2) and (2.7).

Indeed, take an arbitrary $z \in Y$. Then for any scalar $\lambda>0$ we have

$$
J_{E}\left(\tilde{v}_{E}\right)<J_{E}\left(\tilde{v}_{E}+\lambda z\right)
$$

By expanding the terms in the last relation it follows

$$
\begin{aligned}
\varepsilon\left\|P_{E^{\perp}}\left(\tilde{v}_{E}\right)\right\|_{Y}< & \frac{\lambda^{2}}{2}\left\|L^{*} z\right\|_{X}^{2}+\lambda\left\langle L^{*} \tilde{v}_{E} \mid L^{*} z\right\rangle_{X} \\
& +\varepsilon\left\|P_{E^{\perp}}\left(\tilde{v}_{E}+\lambda z\right)\right\|_{Y}-\lambda\langle f \mid z\rangle_{Y} .
\end{aligned}
$$

Consequently, by taking $\lambda>0$ we get

$$
\begin{aligned}
&\langle f \mid z\rangle_{Y}-\left\langle L^{*} \tilde{v}_{E} \mid L^{*} z\right\rangle_{X} \leq \varepsilon \liminf _{\lambda \rightarrow 0^{+}} \frac{\left\|P_{E^{\perp}}\left(\tilde{v}_{E}+\lambda z\right)\right\|_{Y}-\left\|P_{E^{\perp}}\left(\tilde{v}_{E}\right)\right\|_{Y}}{\lambda} \\
&=\varepsilon\left\|P_{E^{\perp}}(z)\right\|_{Y} .
\end{aligned}
$$

Similarly, by taking $\lambda<0$, the same procedure implies

$$
\left|\langle f \mid z\rangle_{Y}-\left\langle L^{*} \tilde{v}_{E} \mid L^{*} z\right\rangle_{X}\right| \leq \varepsilon\left\|P_{E^{\perp}}(z)\right\|_{Y}
$$

From here, by taking an arbitrary $z \in E$ the relation (2.7) follows. The approximation constraint (2.2) is then obtained by varying $z$ in the orthogonal complement of $E$.

The structure of the optimal solution can be easily seen and analysed in terms of Fourier coefficients. To this effect, let us suppose that the operator $\Lambda=L L^{*}$ is diagonalisable, and denote by $\lambda_{k}$ the corresponding eigenvalues. Then the next formula follows directly from (2.6)

$$
\tilde{v}_{k}=\frac{f_{k}}{\lambda_{k}+\varepsilon /\|\tilde{v}\|_{Y}}=\frac{f_{k}-\varepsilon_{k}}{\lambda_{k}}
$$

where $\varepsilon_{k}=\varepsilon \tilde{v}_{k} /\|\tilde{v}\|_{Y}$, while $\tilde{v}_{k}$ and $f_{k}$ denote the $k$-th Fourier coefficient of vectors $\tilde{v}$ and $f$, respectively. Although the formula is not spectrally decomposed (expression
for $\tilde{v}_{k}$ requires knowledge of $\|\tilde{v}\|_{Y}$, i.e. of all the Fourier coefficients), lack of particular frequency in the target $f$ implies that the corresponding Fourier coefficient of the solution $\tilde{v}$ vanishes as well. In particular, if $f$ belongs to some finite-dimensional subspace spanned by finite number of eigenvectors of $\Lambda$, so does the solution $\tilde{v}$.

Furthermore, the expression for the solution is almost explicit, up to a scalar $\|\tilde{v}\|_{Y}$. Although it might look surprisingly at first, this is in accordance with known results for the optimal control problems for parabolic equations (e.g. [10,13]). Note that such problems can be written in the form (2.1) with $L$ being the control to (the final) state operator. If the control acts through initial data, then this can be considered as an inverse problem (of initial source identification). It is an important, but also numerically challenging issue due to the dissipative nature of such equations.

If the eigendecomposition of the operator $\Lambda$ is available, then the numerical procedure of calculating the optimal solution can be reduced to solving the equation for the unknown scalar by some suitable method. For general operators, with variable coefficients and/or acting on irregular domains such decomposition is not available or hard to construct. Another numerical approach should be used in that case, which usually employs some iterative method (e.g., conjugate gradient).

## 3. Parameter dependence

3.1. The problems setting and characterisation of the solutions. In the sequel we want to analyse a family of problems of the type

$$
\begin{equation*}
L_{\nu} u_{\nu}=f_{\nu} \tag{3.1}
\end{equation*}
$$

where $\nu$ is the parameter ranging over a compact, connected set $\mathcal{N} \subset \mathbf{R}^{d}, d \geq 1$.
The goal is to propose an efficient method for finding an optimal approximative solutions to (3.1) for a large number of parameters. To this effect, we make the following assumptions.
(A1) $L_{\nu}$ belongs to $L(X, Y)$ for every parameter, where $X$ and $Y$ are $\nu$-independent Hilbert spaces.
(A2) The associated adjoint operator $L_{\nu}^{*} \in L(Y, X)$ is injective for every $\nu \in \mathcal{N}$ and the family of operators $\Lambda_{\nu}=L_{\nu} L_{\nu}^{*}$ are uniformly bounded in $L(Y)$ from below by a positive self-adjoint operator $\Lambda$, i.e.

$$
\Lambda \leq \Lambda_{\nu}, \quad \nu \in \mathcal{N} .
$$

(The last inequality means that $\left\langle\left(\Lambda_{\nu}-\Lambda\right) v \mid v\right\rangle_{Y} \geq 0$ for every $v \in Y$.)
(A3) The right hand side vectors $f_{\nu}$ belong to a precompact subset of $Y$. In addition, we assume uniform boundedness from below, i.e. there exist a positive constant $f_{-}>\varepsilon$ such that $\left\|f_{\nu}\right\|_{Y} \geq f_{-}$. The aim of the last assumption is to exclude trivial solutions and singularities in the Euler-Lagrange equation (cf. Proposition 2.3).
(A4) The mappings $\nu \rightarrow\left(L_{\nu}, f_{\nu}\right) \in L(X, Y) \times Y$ are analytic.
According to the Hahn-Banach theorem, the injectivity assumption on $L_{\nu}^{*}$ implies that the image of $L_{\nu}$ is dense in $Y$. Thus, for a general target $f_{\nu}$, the problem (3.1) is not well posed. Therefore, as it was discussed in the previous section, we replace
it with the following constrained optimization problem

$$
\begin{equation*}
\min _{u \in X}\left\{\|u\|_{X} \mid\left\|L_{\nu} u-f_{\nu}\right\|_{Y} \leq \varepsilon\right\} \tag{3.2}
\end{equation*}
$$

Supposing that $\left\|f_{\nu}\right\|_{Y}$ is strictly larger than $\varepsilon$ for every parameter value, we know that the optimal solution $\tilde{u}_{\nu}$ equals $L_{\nu}^{*} \tilde{v}_{\nu}$, where $\tilde{v}_{\nu}$ is the the solution to the corresponding Euler-Lagrange equation

$$
\begin{equation*}
L_{\nu} L_{\nu}^{*} \tilde{v}_{\nu}+\varepsilon \frac{\tilde{v}_{\nu}}{\left\|\tilde{v}_{\nu}\right\|_{Y}}-f_{\nu}=0 \tag{3.3}
\end{equation*}
$$

The set of solutions we denote by $\tilde{v}(\mathcal{N})=\left\{\tilde{v}_{\nu}, \nu \in \mathcal{N}\right\}$. Equivalently, the solution $\tilde{v}_{\nu}$ can be characterised as the minimiser of the parameter-dependent functional $J_{\nu}$ defined by (2.4), with $L$ and $f$ being replaced by their parameter-dependent counterparts $L_{\nu}$ and $f_{\nu}$, respectively. The next result provides the boundedness of the set $\tilde{v}(\mathcal{N})$ in $Y$.

Lemma 3.1. The solutions $\tilde{v}_{\nu}$ are uniformly bounded in $Y$. More precisely, there exist positive constants $v_{+}, v_{-}$such that for every $\nu \in \mathcal{N}$ we have $v_{-} \leq\left\|\tilde{v}_{\nu}\right\|_{Y} \leq v_{+}$.
Proof. From the Euler-Lagrange equation (3.3) we get

$$
\left\|L_{\nu} L_{\nu}^{*} \tilde{\nu}_{\nu}\right\|_{Y} \geq\left\|f_{\nu}\right\|_{Y}-\varepsilon
$$

implying the required lower bound with $v_{-}=\left(f_{-}-\varepsilon\right) / L_{+}$, where $f_{-}>\varepsilon$ is the bound from the assumption (A3), and $L_{+}=\max _{\nu}\left\|L_{\nu} L_{\nu}^{*}\right\|_{L(Y)}$. The latter number exists as $\nu \rightarrow L_{\nu}$ is a smooth mapping on a compact set.

In order to obtain the upper bound, let us assume the contrary. We suppose there exists a sequence $\left(\nu_{n}\right)$ such that $\left\|\tilde{v}_{\nu_{n}}\right\|_{Y} \rightarrow \infty$. By multiplying the corresponding Euler-Lagrange equations with $v_{\nu_{n}}^{0}=\tilde{v}_{\nu_{n}} /\left\|\tilde{v}_{\nu_{n}}\right\|_{Y}$ one gets

$$
\begin{equation*}
\left\|v_{\nu_{n}}\right\|_{Y}\left\|L_{\nu_{n}}^{*} v_{\nu_{n}}^{0}\right\|_{X}^{2}+\varepsilon-\left\langle f_{\nu_{n}} \mid v_{\nu_{n}}^{0}\right\rangle_{Y}=0 \tag{3.4}
\end{equation*}
$$

Divergence of the sequence $\left(\left\|v_{\nu_{n}}\right\|_{Y}\right)$ implies that $\left\|L_{\nu_{n}}^{*} v_{\nu_{n}}^{0}\right\|_{X} \rightarrow 0$. By exploring the assumption (A2) and the sandwich theorem we obtain that $\left\langle\Lambda v_{\nu_{n}}^{0} \mid v_{\nu_{n}}^{0}\right\rangle_{Y} \rightarrow 0$.

On the other hand, denoting by $v^{0}$ a weak limit of $v_{\nu_{n}}^{0}$ (up to a subsequence), we have that

$$
\left\langle\Lambda v^{0} \mid v^{0}\right\rangle_{Y}=\left\|L^{*} v^{0}\right\|_{X}^{2} \leq \liminf _{n}\left\|L^{*} v_{\nu_{n}}^{0}\right\|_{X}^{2}=\lim _{n}\left\langle\Lambda v_{\nu_{n}}^{0} \mid v_{\nu_{n}}^{0}\right\rangle_{Y}=0
$$

The positivity of $\Lambda$ implies $v^{0}=0$.
Going back to (3.4) and exploring the relative precompactness of the set $\left\{f_{\nu} \mid \nu \in\right.$ $\mathcal{N}\}$, we obtain

$$
0>\varepsilon-\left\langle f_{\nu_{n}} \mid v_{\nu_{n}}^{0}\right\rangle_{Y} \rightarrow \varepsilon
$$

which provides the required contradiction.
In the next step we explore the smoothness properties of the solution mapping $\nu \rightarrow \tilde{v}_{\nu}$. The smooth dependence of the solutions on the parameter is a bit delicate due to the non-smooth term in the Euler-Lagrange equation. However, the solution vanishes only when the norm of the target is small, i.e. when $\left\|f_{\nu}\right\|_{Y}<\varepsilon$, which is the case we excluded from the analysis through the assumption (A3). Thus, in practice there is no singularity in the parameter-dependent Euler-Lagrange equation (3.3).

This observation plays a key role when analyzing the smooth dependence of the solutions on the parameter. Indeed, we claim that the solutions of (3.3) depend as smoothly on the parameter as the operators $L_{\nu}$ and vectors $f_{\nu}$ do. To see this, we apply the implicit function theorem (IFT). In order to do it successfully we have to analyse the nonlinear function

$$
F(v, \nu):=L_{\nu} L_{\nu}^{*} v+\varepsilon \frac{v}{\|v\|_{Y}}
$$

Writing the solution $\tilde{v}_{\nu}$ implicitly in terms of $\nu$ needs the computation of the derivative of $F$ with respect to $v$ and checking that it is invertible. This is indeed the case since

$$
\partial_{v} F(v, \nu):=L_{\nu} L_{\nu}^{*}+\varepsilon\left(\frac{I}{\|v\|_{Y}}-\frac{v \otimes v}{\|v\|_{Y}^{2}}\right)=L_{\nu} L_{\nu}^{*}+\frac{\varepsilon}{\|v\|_{Y}}\left(I_{Y}-P_{v}\right)
$$

where $I_{Y}$ is the identity, while $P_{v}$ is the orthogonal projection operator on the space spanned by $v$.

We emphasize, despite the apparently singular term arising in the denominator, that singularity does not actually occur since we are always working with nontrivial solutions $v \neq 0$. The IFT can then be applied without difficulty since the resulting operator $L_{\nu} L_{\nu}^{*}+\frac{\varepsilon}{\|\cdot\|_{Y}}\left(I_{Y}-P_{v}\right)$ is invertible. This can be easily seen by the Lax-Milgram lemma.

The direct application of the Lax-Milgram Lemma requires the coercivity of $\partial_{v} F(v, \nu)=L_{\nu} L_{\nu}^{*}+\frac{\varepsilon}{\|v\|_{Y}}\left(I_{Y}-P_{v}\right)$. However, in our setting in which only an approximate solution to the original system exists, $L_{\nu}$ can not be guaranteed to be coercive and neither the projection $I_{Y}-P_{v}$ is coercive. In fact, it degenerates along the direction $v$ by the very definition of the projection. Thus it is unclear whether the addition of the term $I_{Y}-P_{v}$ to the operator $\Lambda_{\nu}$ suffices to guarantee the invertibility of $\partial_{v} F(v, \nu)$. The positive answer is provided by the following lemma.

Lemma 3.2. Assume $Q=Q_{1}+Q_{2}$ is a bounded linear operator on a Hilbert space $H$, with $Q_{1}$ being a positive linear operator, while $Q_{2}$ is a self adjoint operator which vanishes on a finite dimensional subspace $V_{1}$ and is coercive on $V_{2}=V_{1}^{\perp}$. Then the operator $Q$ is coercive on $H$.

Proof. In order to prove the lemma we have to show that there exists a constant $c>0$ such that for every $v \in H$ it holds

$$
\langle Q v \mid v\rangle_{H} \geq c\|v\|_{H}^{2}
$$

We argue by contradiction. Assuming the contrary there exists a sequence $\left(v_{n}\right)$ such that

$$
\left\langle Q v_{n} \mid v_{n}\right\rangle_{H} \leq \frac{1}{n}\left\|v_{n}\right\|_{H}^{2}
$$

Let us introduce the decomposition $v_{n}=v_{n, 1}+v_{n, 2}$, where $v_{n, i} \in V_{i}, i=1,2$. By using the assumptions of the lemma, it follows

$$
\left\langle Q v_{n} \mid v_{n}\right\rangle_{H}=\left\langle Q_{1} v_{n} \mid v_{n}\right\rangle_{H}+\left\langle Q_{2} v_{n, 2} \mid v_{n, 2}\right\rangle_{H} \leq \frac{1}{n}\left\|v_{n}\right\|_{H}^{2}
$$

Dividing the last equation by $\left\|v_{n}\right\|_{H}^{2}$ and exploring the coercivity of $Q_{2}$ on the subspace $V_{2}$ we obtain

$$
\begin{equation*}
\left\langle Q_{1} v_{n}^{0} \mid v_{n}^{0}\right\rangle_{H}+c_{2} \frac{\left\|v_{n, 2}\right\|_{H}^{2}}{\left\|v_{n}\right\|_{H}^{2}} \leq \frac{1}{n}\left\|v_{n}\right\|_{H}^{2} \tag{3.5}
\end{equation*}
$$

where $v_{n}^{0}$ stands for the normalized vector $v_{n} /\left\|v_{n}\right\|_{H}$, while $c_{2}$ is the coercivity constant of $Q_{2}$ (on $V_{2}$ ).

The positivity of $Q_{1}$ implies $\frac{v_{n, 2}}{\left\|v_{n}\right\|_{H}} \rightarrow 0$ strongly in $H$. As $V_{1}$ is finite dimensional, we obtain the strong convergence of the whole normalized sequence $v_{n}^{0} \rightarrow v_{1}^{0} \in$ $V_{1} \backslash\{0\}$.

Finally, by passing to the limit in (3.5) we obtain $\left\langle Q_{1} v_{1}^{0} \mid v_{1}^{0}\right\rangle_{H}=0$ which contradicts the positivity of $Q_{1}$.

As discussed above, the last lemma allows application of IFT. In particular, we employ its analytic version (e.g. [16, §Appendix B]), by which we preserve the smoothness imposed by assumption (A4). We summarize the results of this discussion in the following proposition.

Proposition 3.3. The solutions to the Euler-Lagrange equation (3.3) preserve the regularity imposed on the mapping $\nu \rightarrow\left(L_{\nu}, f_{\nu}\right)$. In particular, if the latter is analytic, the same holds for the solution mapping.
3.2. Greedy approach and the Tychonoff regularization. The aim of this section is to develop an efficient numerical algorithm for reconstructing an arbitrary element of $\tilde{v}(\mathcal{N})$ corresponding to some given parameter value. In order to accomplish this task we rely on greedy algorithms which were introduced and analysed through the last two decades in the context of parametric PDEs. They serve as on of the most popular tools for construction of reduced basis (cf. [11] and the references therein).

The objective of the greedy approach is to approximate a compact set $\mathcal{K}$ (e.g. a family of solutions to parameter dependent problems) in a Banach space $Y$ by a linear subspace $V_{n}$ of (small) dimension $n$. The selection of its basis vectors is done gradually in the offline phase of the algorithm. Once the basis is known, the approximation for an arbitrary given element in $\mathcal{K}$ is computed in the online phase. Usually, the computational effort for the offline routine is much higher than for the online one, but it is performed only once.

A greedy approximation is optimal, where the optimality is to be understood in the sense of the Kolmogorov widths. The Kolmogorov $n$-width defines (theoretically) the best possible approximation error one could obtain by a subspace in $Y$ of a fixed dimension $n$. The results of $[2,7]$ show that the greedy approximation errors decay asymptotically with the same rate (exponential or polynomial) as the Kolmogorov widths.

Furthermore, the Kolmogorov widths are preserved under smooth (analytic) mappings ([4]). In such a way one can a-priori estimate approximation performance of a subspace constructed by a greedy method, by estimating Kolmogorov widths of a set of admissible parameters that generate a set of interest. However, it is important to emphasize that the very implementation of a greedy procedure does not
require analytic smoothness. Only, in that case, we lack the a-priori estimates on the approximation errors. In particular, this might happen if the parameter set $\mathcal{N}$ consists of a finite or a countable number of elements, which is not the case studied in this article.

In the development of greedy algorithms one of the main steps is the construction of some surrogate function which enables us to calculate the distance between unknown terms, or terms that are in general hard to calculate. Usually such a surrogate is provided by using some kind of residual. More, precisely, let us assume we have calculated $\tilde{v}_{\nu_{1}}$ for some parameter value $\nu_{1}$. We would like to check if we can use it to approximate some other optimal vector $\tilde{v}_{\nu}$. To this effect we plug $\tilde{v}_{\nu_{1}}$ into the equation satisfied by the latter term and define the residual

$$
R_{\nu} \tilde{v}_{\nu_{1}}:=L_{\nu} L_{\nu}^{*} \tilde{v}_{\nu_{1}}+\varepsilon N \tilde{v}_{\nu_{1}}-f_{\nu}
$$

where by $N$ we denote the normalization operator $N v=v /\|v\|_{Y}$. If the residual turns to be zero, due to the uniqueness of the solution it follows that $\tilde{v}_{\nu}=\tilde{v}_{\nu_{1}}$. In general we would like the residual to measure the distance between two optimal vectors. More precisely, we require estimates of the form

$$
\begin{equation*}
c_{-}\left\|\tilde{v}_{\nu}-\tilde{v}_{\nu_{1}}\right\|_{Y} \leq\left\|R_{\nu} \tilde{v}_{\nu_{1}}\right\|_{Y} \leq c_{+}\left\|\tilde{v}_{\nu}-\tilde{v}_{\nu_{1}}\right\|_{Y} \tag{3.6}
\end{equation*}
$$

where $c_{-}$and $c_{+}$are $\nu$-independent positive constants.
In order to obtain such kind of estimates let us rewrite the residual operator as

$$
\begin{equation*}
R_{\nu} \tilde{v}_{\nu_{1}}:=L_{\nu} L_{\nu}^{*}\left(\tilde{v}_{\nu_{1}}-\tilde{v}_{\nu}\right)+\varepsilon\left(N \tilde{v}_{\nu_{1}}-N \tilde{v}_{\nu}\right), \tag{3.7}
\end{equation*}
$$

where we explored the Euler-Lagrange equation (3.3). The upper bound in (3.6) now follows easily. For the first summand in (3.7) it is a direct consequence of the boundedness assumption, while for the last term in (3.7) it is obtained by using geometrical interpretation and the triangular inequality (cf. [1, §1.2]).

However, due to the structure of the normalization operator (which is constant along each half-ray emerging from the origin), one can easily check that the lower bound in (3.6) is equivalent to the coercivity of the operator $L_{\nu} L_{\nu}^{*}$. Of course, our assumptions on $L_{\nu}$ do not provide the required coercivity and we have to propose an alternative method. Note that here we can not use the approach applied in Lemma 3.2 as the normalization operator $N$ does not vanish on any non-trivial subspace.

In order to overcome the lack of coercivity, let us introduce a two-parameter family of linear problems

$$
\begin{equation*}
L_{\nu} L_{\nu}^{*} \tilde{w}(\nu, \eta)+\eta \tilde{w}(\nu, \eta)=f_{\nu} \tag{3.8}
\end{equation*}
$$

Here we have substituted the nonlinear normalization operator $N$ appearing in the Euler-Lagrange equation (3.3) by a linear term multiplied by the new parameter $\eta$. Let us note that the last equation coincides with the Euler-Lagrange equation obtained by minimization of the functional

$$
\begin{equation*}
J(v)=\frac{1}{2}\left\|L_{\nu}^{*} v\right\|_{X}^{2}+\eta\|v\|_{Y}^{2}-\left\langle f_{\nu} \mid v\right\rangle_{Y}, \quad v \in Y \tag{3.9}
\end{equation*}
$$

The non-smooth term appearing in the original functional (2.4) is here replaced by a quadratic one. This improves the coercivity properties of the problem, which is a standard benefit of the Tychonoff regularization.

The functional appearing in (3.9) corresponds to the penalization approach for approximation problems, where one forces the solution to approach the given target by letting the penalization constant blow up. It provides a smooth functional that is easier to handle, unlike the one in (2.4). However, such an approach does not allow an a-priori estimate of the deviation from the target, in particular, it does not ensure that it is smaller than the given approximation error $\varepsilon$. For this reason, the value of the parameter $\eta$ in (3.9) is not fixed, but is allowed to vary within a specific interval.

More precisely, we suppose the introduced parameter ranges within the interval $\left[\eta_{-}, \eta_{+}\right]=\left[\varepsilon / v_{+}, \varepsilon / v_{-}\right]$, where $v_{ \pm}$are bounds from Lemma 3.1. In such a way, for each parameter $\nu$ there exists an $\eta$ from the given range such that $\eta=\varepsilon /\left\|\tilde{v}_{\nu}\right\|_{Y}$, where $\tilde{v}_{\nu}$ is the corresponding solution of the Euler-Lagrange equation (3.3).

This implies that solving a family of two-parameter problems (3.8) for $(\nu, \eta)=$ $\mathcal{N} \times\left[\eta_{-}, \eta_{+}\right]$will also provide the solution of the original $\varepsilon$ problem.

In order to efficiently treat the auxiliary problem (3.8) we explore greedy algorithms introduced above. The problem is now linear, and it involves selfadjoint operators of the form $L_{\nu} L_{\nu}^{*}+\eta I$ which are uniformly bounded from below by $\eta_{-} I$. This allows one to consider the residual of the form

$$
R(\nu, \eta) w:=L_{\nu} L_{\nu}^{*} w+\eta w-f_{\nu}
$$

with $w$ being an arbitrary test function. Based on the above obtained bounds it directly follows

$$
c_{-}\|\tilde{w}(\nu, \eta)-w\|_{Y} \leq\|R(\nu, \eta) w\|_{Y} \leq c_{+}\|\tilde{w}(\nu, \eta)-w\|_{Y}
$$

Furthermore, employing the same kind of arguments based on the IFT and presented in the previous subsection, the solutions $\tilde{w}(\nu, \eta)$ to (3.8) preserve the smoothness of the mapping $\nu \rightarrow\left(L_{\nu}, f_{\nu}\right)$ at all levels. However, note that in this case the corresponding analysis is much simpler as there is no singularity in the two parameter equation (3.8) and we deal with uniformly coercive operators.

In particular if the operators $L_{\nu}$ and non-homogenous terms $f_{\nu}$ depend analytically on $\nu$, then the mapping $(\nu, \eta) \rightarrow \tilde{w}(\nu, \eta)$ is analytic as well. In this way, as discussed at the beginning of this section, one preserves the Kolmogorov widths of the two-parameter set $\mathcal{N} \times\left[\eta_{-}, \eta_{+}\right]$which are transferred to the manifold of solutions $\tilde{W}=\left\{\tilde{w}(\nu, \eta) \mid(\nu, \eta)=\mathcal{N} \times\left[\eta_{-}, \eta_{+}\right]\right\}(c f .[4])$.

Suppose we have performed a greedy algorithm for the two-parameter problems (3.8). It selects a finite set of parameter pairs $\left(\nu_{i}, \eta_{i}\right), i=1 . . N$ and returns the corresponding solutions $w_{i}:=\tilde{w}\left(\nu_{i}, \eta_{i}\right)$ which constitute a reduced basis for the manifold $\tilde{W}$. In other words, for every value of $(\nu, \eta)$ there exist a set of linear coefficients $\alpha_{i}$ such that $\left\|\tilde{w}(\nu, \eta)-\sum \alpha_{i} w_{i}\right\|_{Y}<\delta$ where $\delta$ is a positive constant determined by the stopping criteria of the greedy algorithm.

In the next step we want to go back to the original $\varepsilon$ problem (3.3). More precisely, given an arbitrary value of the parameter $\nu$ we want to determine a set of coefficients $\alpha_{i}$ such that $v_{\nu}^{\star}=\sum \alpha_{i} w_{i}$ brings the system within $\varepsilon$ distance from the target $f_{\nu}$. The problem is feasible, as for $\eta=\varepsilon /\left\|\tilde{v}_{\nu}\right\|_{Y}$ the greedy algorithm provides a good approximation of the solution to (3.8), and consequently to the original Euler-Lagrange equation (3.3).

Therefore, let us propose the required approximation by projecting the target $f_{\nu}$ on the space spanned by $L_{\nu} L_{\nu}^{*} w_{i}$. More precisely, we determine the approximation coefficients $\alpha_{i}$ as solutions to the system

$$
\sum_{i=1}^{m} L_{\nu} L_{\nu}^{*} \alpha_{i} w_{i}=P_{\nu}^{m} f_{\nu}
$$

with $P_{\nu}^{m}$ denoting the orthogonal projection on the space spanned by vectors $L_{\nu} L_{\nu}^{*} w_{i}$, while $m$ stands for the number of parameters selected during the offline phase of the greedy procedure. The approximation of the solution $\tilde{u}_{\nu}$ to the constrained optimal control problem (3.2) is then given by

$$
\tilde{u}_{\nu} \approx \sum_{i=1}^{m} L_{\nu}^{*} \alpha_{i} w_{i}
$$

Such procedure results in an approximate solution that steers the system to the target $f_{\nu}$ as close as possible by means of the constructed reduced basis space $\left\{w_{i}, i=1 . . m\right\}$.

## 4. UnBounded operators

The theory developed in the previous section requires the operators of interest to be bounded. However, it can be generalised and applied to unbounded operators as well.

In this section we consider a family of linear unbounded operators $A_{\nu}$ on a Hilbert space $H$, where, as before, $\nu$ is the parameter ranging over a compact, connected set $\mathcal{N} \subset \mathbf{R}^{d}, d \geq 1$. We put the following hypothesis on the considered family.
(H1) $A_{\nu}$ are positive, self adjoint operators uniformly bounded from below, i.e. there exists $\alpha \in \mathbf{R}^{+}$such that $A_{\nu} \geq \alpha>0$ for every $\nu \in \mathcal{N}$;
(H2) The operators $A_{\nu}$ have common domains, i.e. there exists a subspace $D_{A} \subseteq$ $H$ such that $D\left(A_{\nu}\right)=D_{A}$ for every $\nu \in \mathcal{N}$;
(H3) the graph norms of $A_{\nu}$ are uniformly equivalent.
Remark 4.1. Instead of (H1) one can require the operators $A_{\nu}$ to be uniformly lower (or upper) bounded by an arbitrary constant. In order to simplify the presentation we restrict to the case of positive definite operators.

In addition we suppose that $H$ is densely and compactly embedded into some Hilbert space $Y$. We pose the problem of finding (an approximative) solution to the equation

$$
A_{\nu} x=y, \quad y \in Y
$$

Based on the assumption (H1) the image of operators $A_{\nu}$ equals $H$, but it is only dense in $Y$. This brings us to the problem of finding the optimal approximative solution discussed in previous sections. However, the theory we developed assumes bounded operators, which $A_{\nu}$ are not. In order to overcome this gap, we want to associate to each $A_{\nu}$ an operator $L_{\nu}: X \rightarrow Y$, where $X$ is some parameter independent Hilbert space which is still to be defined.

To this effect, let us denote by $X_{\nu}$ the space $D_{A}$ equipped with the norm

$$
\begin{equation*}
\|x\|_{X_{\nu}}=\left\|\left(\beta I_{H}-A_{\nu}\right) x\right\|_{H}, \tag{4.1}
\end{equation*}
$$

where $\beta<0$ is a scalar from the resolvent set of $A_{\nu}$, while $I_{H}$ stands for the identity on $H$. Note that the introduced space corresponds to the one denoted by $X_{1}$ in $[17, \S 2.10]$, and their norms are uniformly equivalent to the graph norms of $A_{\nu}$ ( [17, Proposition 2.10.1]). Based on the assumption (H3) it follows that the norms (4.1) are uniformly equivalent. For this reason in the sequel we use a common notation $X$ for all spaces $X_{\nu}$.

Let us denote by $\tilde{L}_{\nu}$ a family of operators from $X$ to $H$ defined by

$$
\tilde{L}_{\nu} x=A_{\nu} x
$$

By the definition of the space $X$, it is not difficult to check that the introduced operators are uniformly bounded and coercive. Indeed, we have

$$
\begin{equation*}
\left\|\tilde{L}_{\nu}\right\|_{L(X, H)}=\sup _{x \in D_{A}} \frac{\left\|\tilde{L}_{\nu} x\right\|_{H}}{\|x\|_{X}}=\sup _{x \in D_{A}} \frac{\left\|A_{\nu} x\right\|_{H}}{\left\|\left(\beta I_{H}-A_{\nu}\right) x\right\|_{H}} \leq 1 \tag{4.2}
\end{equation*}
$$

where the last inequality follows from the positivity of the operator $A_{\nu}$. Similarly, by exploring

$$
\begin{equation*}
\frac{\left\|\tilde{L}_{\nu} x\right\|_{H}}{\|x\|_{X}} \geq \frac{\left\|A_{\nu} x\right\|_{H}}{\|\beta x\|_{H}+\left\|A_{\nu} x\right\|_{H}} \geq \frac{\alpha}{\beta+\alpha} \tag{4.3}
\end{equation*}
$$

where $\alpha$ is the lower bound from the assumption (H1), one obtains the lower bound on $\tilde{L}_{\nu}$.

Thus we obtain that $\tilde{L}_{\nu}$ form a family of uniformly bounded operators in $L(X, H)$. Based on the assumption (H1) both the operators $A_{\nu}$ and $\tilde{L}_{\nu}$ are surjective into $H$. For the corresponding adjoints we have the following characterization.
Lemma 4.2. The adjoint operators $\tilde{L}_{\nu}^{*} \in L(H, X)$ are of the form $\tilde{L}_{\nu}^{*}=A_{\nu}\left(\beta I_{H}-\right.$ $\left.A_{\nu}\right)^{-2} \in L(H, X)$, satisfying the same lower and upper bounded estimates as the operators $\tilde{L}_{\nu} d o$.

Proof. By the definition of the adjoint, for $u, x \in X$ we have

$$
\begin{aligned}
\left\langle x \mid \tilde{L}_{\nu}^{*} u\right\rangle_{X} & =\left\langle\tilde{L}_{\nu} x \mid u\right\rangle_{H}=\left\langle A_{\nu} x \mid u\right\rangle_{H}=\left\langle x \mid A_{\nu} u\right\rangle_{H} \\
& =\left\langle\left(\beta I_{H}-A_{\nu}\right)^{-1} x \mid\left(\beta I_{H}-A_{\nu}\right)^{-1} A_{\nu} u\right\rangle_{X}
\end{aligned}
$$

From here we get

$$
\tilde{L}_{\nu}^{*} u=A_{\nu}\left(\beta I_{H}-A_{\nu}\right)^{-2} u, \quad u \in X
$$

As $X$ is dense in $H$, and the operators $A_{\nu}\left(\beta I_{H}-A_{\nu}\right)^{-1}$ are bounded on $H$, the last relation holds for an arbitrary $u \in H$. This provides the first part of the statement.

By using the obtained explicit expression for the adjoint operators, we have

$$
\frac{\left\|\tilde{L}_{\nu}^{*} u\right\|_{X}}{\|u\|_{H}}=\frac{\left\|A_{\nu}\left(\beta I_{H}-A_{\nu}\right)^{-1} u\right\|_{H}}{\|u\|_{H}}=\frac{\|A x\|_{H}}{\left\|\left(\beta I_{H}-A\right) x\right\|_{H}}
$$

where $x=\left(\beta I_{H}-A_{\nu}\right)^{-1} v \in X$. The required bounds now follow from (4.2) and (4.3).

Finally, in order to put us in the context of the preceding section, we introduce $L_{\nu} \in L(X, Y)$ defined as

$$
\begin{equation*}
L_{\nu}:=I \tilde{L}_{\nu} \tag{4.4}
\end{equation*}
$$

where $I$ stands for the inclusion operator $I: H \rightarrow Y$. As $H$ is compactly embedded into $Y$, the operator $L_{\nu}$ is bounded, injective, but not coercive operator with the dense image in $Y$. Moreover, due to the uniform coercivity of operators $\tilde{L}_{\nu}^{*}$ we have that

$$
L_{\nu} L_{\nu}^{*} \geq c I I^{*}
$$

This implies the assumption (A2) is fulfilled and consequently we fit the setting of the previous section.

## 5. A SPECIFIC EXAMPLE

We consider a family of Dirichlet Laplacians on $L^{2}(\Omega)$, where $\Omega$ is assumed to be an open, bounded set with a smooth boundary. More precisely, we define a family of unbounded operators on $L^{2}(\Omega)$

$$
\begin{equation*}
-\Delta_{\nu}=-\operatorname{div}\left(\mathbf{A}_{\nu} \nabla\right) \tag{5.1}
\end{equation*}
$$

accompanied by Dirichlet boundary conditions. For the coefficients $\mathbf{A}_{\nu}$ we assume that they depend smoothly on $\nu$ and they satisfy uniform boundedness and coercivity properties. More precisely, we assume $\mathbf{A}_{\nu} \in C^{1}(\Omega)$ are self adjoint matrix functions such that

$$
\mathbf{A}_{-}(x) \leq \mathbf{A}_{\nu}(x) \leq \mathbf{A}_{+}(x), \quad \nu \in \mathcal{N}, x \in \Omega
$$

for some bounded and coercive matrix functions $\mathbf{A}_{ \pm} \geq \alpha>0$. By using the Poincare inequality this immediately applies the uniform coercivity of the considered Laplacians, i.e. $-\Delta_{\nu} \geq \alpha$.

The domain of the Laplacian $-\Delta_{\nu}$ is parameter independent and coincides with $X=H^{2}(\Omega) \cap H_{0}^{1}(\Omega)$ for every $\nu$. Furthermore, by standard elliptic regularity results (e.g. $[8, \S 6.3]$ ), their graph norms are (uniformly) equivalent to $H^{2}$ norm. Consequently, the hypothesis (H1)-(H3) from the previous section are satisfied.

For the target space we take $Y=H^{-1}(\Omega)$ for which we have the dense and compact embedding $L^{2}(\Omega) \hookrightarrow H^{-1}(\Omega)$ (of course, any other space $H^{s}(\Omega)$, with $s<0$, will be appropriate in this context). Consequently, we introduce a sequence of bounded operators $L_{\nu} \in L(X, Y)$ defined by

$$
L_{\nu} u=-\operatorname{div}\left(\mathbf{A}_{\nu} \nabla u\right)
$$

The properties of the introduced operators can be examined through their matrix representation. To this effect, we explore the spectral decomposition of the Laplacian operator (5.1). In particular, there exists an orthonormal basis in $L^{2}(\Omega)$ consisting of eigenfunctions $\psi_{i, \nu}$ of $-\Delta_{\nu}$ such that

$$
-\Delta_{\nu} \sim \operatorname{diag}\left(\lambda_{1, \nu}, \lambda_{2, \nu}, \ldots\right)
$$

where $\left(\lambda_{i, \nu}\right)$ is a sequence of (positive) eigenvalues diverging to infinity.
Then it is not difficult to check that the matrix representation of the associated operator $L_{\nu}$ in the pair of basis $\hat{\psi}_{i, \nu}:=\psi_{i, \nu} / \lambda_{i, \nu} \in H^{2}(\Omega) \cap H_{0}^{1}(\Omega)$ and $\tilde{\psi}_{i, \nu}:=$ $\sqrt{\lambda_{i, \nu}} \psi_{i, \nu} \in H^{-1}(\Omega)$ has the form

$$
L_{\nu} \sim \operatorname{diag}\left(\frac{1}{\sqrt{\lambda_{1, \nu}}}, \frac{1}{\sqrt{\lambda_{2, \nu}}}, \ldots\right)
$$

Similarly, the same matrix representation form, in the reverse pairs of basis $\left\{\tilde{\psi}_{i, \nu}\right\}$ and $\left\{\hat{\psi}_{i, \nu}\right\}$, also holds for the adjoint operator $L_{\nu}^{*}$. Due to the properties of eigenvalue sequence, the operator $L_{\nu}^{*}$ is injective, but not coercive. This is equivalent to the statement that the image of the Laplacian is (only) dense in $H^{-1}(\Omega)$.

Having associated to Laplacians (5.1) a sequence of bounded operators $L_{\nu}$ with dense images, we can explore the approach developed in Section 3 for finding the optimal approximative solution. To this effect, we assume $f_{\nu}: \mathcal{N} \rightarrow H^{-1}(\Omega)$ is a smooth function with a precompact image and we consider a sequence of problems

$$
\begin{equation*}
-\operatorname{div}\left(\mathbf{A}_{\nu} \nabla\right) u=f_{\nu}, \tag{5.2}
\end{equation*}
$$

with solutions $u$ searched within the domain of the Laplacian, i.e. in $X=H^{2}(\Omega) \cap$ $H_{0}^{1}(\Omega)$. As $L_{\nu}(X)$ is only dense in $H^{-1}(\Omega)$, the problem in general does not admit a solution, and we relax it by considering an approximative one of the form (3.2). The optimal solution is thus of the form $\tilde{u}_{\nu}=L_{\nu}^{*} \tilde{v}_{\nu}$, where $\tilde{v}_{\nu}$ is the solution of the corresponding Euler-Lagrange equation (3.3). In the next lemma we provide the explicit expression of the adjoint operator $L_{\nu}^{*}$.
Lemma 5.1. The adjoint operator equals to $L_{\nu}^{*}=(-\Delta)^{-2}\left(-\Delta_{\nu}\right)(-\Delta)^{-1}$, where $-\Delta$ is the Laplacian with constant coefficients $\mathbf{A}=\mathbf{I}$.
Proof. Let us write the operator $L_{\nu}$ in the form (4.4), i.e. $L_{\nu}=I \tilde{L}_{\nu}$. Here $I$ stands for the inclusion operator $I: L^{2}(\Omega) \hookrightarrow H^{-1}(\Omega)$, while $\tilde{L}_{\nu}$ is the operator from $H^{2}(\Omega) \cap H_{0}^{1}(\Omega)$ to $L^{2}(\Omega)$ defined by $\tilde{L}_{\nu} u=-\operatorname{div}\left(\mathbf{A}_{\nu} \nabla u\right)$.

Then for calculating the adjoint we use the relation $L_{\nu}^{*}=\tilde{L}_{\nu}^{*} I^{*}$. Similarly as in Lemma 4.2 one obtains that $\tilde{L}_{\nu}^{*}=(-\Delta)^{-2}\left(-\Delta_{\nu}\right)$. Thus it remains to express the adjoint of the inclusion. To this effect let us note

$$
\left\langle u \mid I^{*} v\right\rangle_{L^{2}}=\langle I u \mid v\rangle_{H^{-1}}=\left\langle(-\Delta)^{-1 / 2} u \mid(-\Delta)^{-1 / 2} v\right\rangle_{L^{2}},
$$

from where we get $I^{*}=(-\Delta)^{-1}$, which completes the proof.
Based on the last lemma, the Euler-Lagrange equation (3.3) for the problem (5.2) can be rewritten as

$$
\begin{equation*}
\left(-\Delta_{\nu}\right)(-\Delta)^{-2}\left(-\Delta_{\nu}\right)(-\Delta)^{-1} \tilde{v}_{\nu}+\varepsilon N \tilde{v}_{\nu}-f_{\nu}=0 \tag{5.3}
\end{equation*}
$$

where, as before, $N$ stands for the normalization operator $N v=v /\|v\|_{H^{-1}}$.
Of course, in numerical calculations one should employ some standard discretisation procedure (based on finite differences or finite elements), and reduce the last equation to a finite dimensional, algebraic one. The numerical resolution using gradient methods for a specific realization of the model has been previously developed in [3]. Once the duality in Propositions 2.1-2.3 has been adopted in the computational method, its implementation does not differ significantly from the classical problems, that in our case would correspond to operators $L$ which, instead of being dense, are assumed to have a full range.

However, as we have explained, our goal is not to solve the equation (5.3) for every value of the parameter. Instead, we employ a greedy procedure which allows us to obtain a reduced basis by which we can approximate any $\tilde{v}_{\nu}$ with some apriori given precision. The reduced basis is constructed during the offline phase,
and consists of solutions to two-parameter problem (3.8) for some carefully selected parameters' values.

To the best of authors' knowledge, so far the greedy methods were addressed only to the exact solutions. On the other hand, the approximate solution problem is relevant and natural in several contexts such as image processing, control of timedependent irreversible processes, etc. Our contribution is the analog of the greedy theory for exact solutions in the context of the approximate one (e.g. [5] ).

The complete analysis of the computational cost is beyond the scope of this paper, but the well-known conclusions for the exact solution problem apply in this case as well (cf. $[12, \S 5]$ ). In particular, the implementation requires extensive offline work, which is the most expensive part of the algorithm. Once this is done, the method is optimal since it leads to sharp approximation rates.

The cost of the online part of the algorithm is of order $m C$, where $C$ is the cost of applying (a finite dimensional approximation of ) the operator $L_{\nu}$ to an arbitrary vector, while $m$ is the number of parameters selected during the offline phase. Consequently, the cost reduction obtained by choosing the greedy algorithm depends linearly on the ratio between the number of selected snapshots and the system dimension.

Practical application of the method is justified when the approximate equation needs to be solved for a wide class of parameter-dependent problems.

## 6. Conclusion

In this article we develop a procedure for handling a class of parameter dependent, ill-posed problems, that, in general, do not allow exact solution. The original problem is relaxed by considering corresponding approximate problems, whose optimal solutions are well defined, where the optimality is determined through the minimal norm requirement.

The procedure is based upon reduced basis methods, in particular upon greedy algorithms, by which one constructs a reduced basis during the offline phase. An approximation of the solution for a specific parameter is constructed in the online phase as a suitable linear combination of reduced basis vectors. In order to provide apriori estimates for the algorithm, a Tychonoff-kind regularization is applied, which adds an additional parameter to the model.

The theory is developed in a rather general theoretical framework, which allows its application to different kinds of problems. As a specific example we consider a family of ill-posed elliptic problems. The required general assumptions in this case translate to rather natural uniform lower and upper bounds on coefficients of the considered operators.

Other potential applications would include approximate controllability, inverse problems related to high dissipative systems, like heat equations, deconvolution in image processing etc. Beside identifying the class of suitable operators for each of these problems, it would also be interesting to perform corresponding numerical simulations and verify the efficiency of the method on particular examples.

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## References

[1] N. Antonić and M. Lazar, Parabolic H-measures, Journal of Functional Analysis 265 (2013), 1190-1239.
[2] P. Binev, A. Cohen, W.Dahmen, R. DeVore, G. Petrova and P. Wojtaszczyk, Convergence rates for greedy algorithms in reduced basis methods, SIAM J. Math. Anal. 43 (2011), 1457-1472.
[3] C. Carthel, R. Glowinski and J. Lions, On exact and approximate boundary controllabilities for the heat equation: a numerical approach, J. Optim. Theory Appl. 82 (1994), 429-484.
[4] A. Cohen and R. DeVore, Kolmogorov widths under holomorphic mappings, IMA Journal on Numerical Analysis 36 (2016), 1-12.
[5] A. Cohen and R.DeVore, Approximation of high-dimensional parametric PDEs, Acta Numerica 24 (2015) 1-159.
[6] R. DeVore, The theoretical foundation of reduced basis methods in: Model Reduction and Approximation, P. Benner, A. Cohen, M. Ohlberger, and K. Willcox (eds.), SIAM, Philadelphia, PA, 2017, pp. 137-168.
[7] R. DeVore, G. Petrova and P.Wojtaszczyk, Greedy algorithms for reduced bases in Banach spaces, Constr. Approx. 37 (2013), 455-466.
[8] L. Evans, Partial Differential Equations, American Mathematical Society, Providence, RI, 2010.
[9] C. Fabre, J. P. Puel and E. Zuazua, Approximate controllability for the semilinear heat equation Proc. Roy. Soc. Edinburgh 125A (1995), 31-61.
[10] L. Grubišić, M. Lazar, I. Nakić and M. Tautenhahn, Optimal control of parabolic equations a spectral calculus based approach, submitted, 2021.
[11] M. Lazar and J. Lohéac, Control of parameter dependent systems, in: Numerical Control and beyond, E. Trélat and E. Zuazua (eds.), Handbook of Numerical Analysis, vol. 22, Elsevier, 2022.
[12] M. Lazar and E. Zuazua, Greedy controllability of finite dimensional linear systems, Automatica 74 (2016), 327-340.
[13] M. Lazar and C. Molinari, Optimal distributed control of the heat-type equations by spectral decomposition, Optimal Control, Applications \& Methods 42 (2021), 891-926.
[14] J. L. Lions, Remarks on approximate controllability, J. Anal. Math. 59 (1992), 103-116.
[15] J. Peypouquet, Convex Optimization in Normed Spaces: Theory, Methods and Examples, Springer Briefs in Optimization, 2015.
[16] J. Pöschel and E. Trubowitz, Inverse Spectral Theory, Academic Press, San Diego, 1987.
[17] M. Tucsnak and G. Weiss, Observation and Control for Operator Semigroups, BirkhauserVerlag AG, Berlin, 2009.
[18] K. Veroy, C. Prud'Homme, D. Rovas, and A. Patera, A posteriori error bounds for reduced-basis approximation of parametrized noncoercive and nonlinear elliptic partial differential equations, in: 16th AIAA Computational Fluid Dynamics Conference. Orlando, United States, 2003.
[19] E. Zuazua Approximate controllability for linear parabolic equations with rapidly oscillating coefficients, Control and Cybernetics 23 (1994), 793-801.
[20] E. Zuazua, Finite dimensional null controllability for the semilinear heat equation, Journal de Mathématiques Pures et Appliquées 76 (1997), 237-264.

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