# Technische Universität Berlin <br> Institut für Mathematik 

# ON TWO NUMERICAL METHODS FOR STATE-CONSTRAINED ELLIPTIC CONTROL PROBLEMS 

C. Meyer, u. Prüfert, and F. Tröltzsch

Preprint 2005/05

A linear-quadratic elliptic control problem with pointwise box constraints on the state is considered. The state-constraints are treated by a Lavrentiev type regularization. It is shown that the Lagrange multiplier associated with the regularized state-constraints are functions in $L^{2}$. Moreover, the convergence of the regularized controls is proven for regularization parameter tending to zero. To solve the problem numerically, an interior point method and a primal-dual active set strategy are implemented and treated in function space.

AMS subject classifications. 49J20, 49M20, 90C51, 65 K 10

# ON TWO NUMERICAL METHODS FOR STATE-CONSTRAINED ELLIPTIC CONTROL PROBLEMS* 

CHRISTIAN MEYER, UWE PRÜFERT, FREDI TRÖLTZSCH ${ }^{1}$


#### Abstract

A linear-quadratic elliptic control problem with pointwise box constraints on the state is considered. The state-constraints are treated by a Lavrentiev type regularization. It is shown that the Lagrange multiplier associated with the regularized state-constraints are functions in $L^{2}$. Moreover, the convergence of the regularized controls is proven for regularization parameter tending to zero. To solve the problem numerically, an interior point method and a primal-dual active set strategy are implemented and treated in function space.


Key words. Linear elliptic equations, quadratic optimal control problem, pointwise state constraints, interior point method, active set strategy

AMS subject classifications. 49J20, 49M20, 90C51, 65K10

1. Introduction. In this paper, we consider the numerical solution of the elliptic optimal control problem

$$
(\mathrm{P})\left\{\begin{array}{rl}
\text { minimize } & J(y, u)
\end{array}=\frac{1}{2} \int_{\Omega}\left(y-y_{d}\right)^{2} d x+\frac{\kappa}{2} \int_{\Omega}\left(u-u_{d}\right)^{2} d x ~ 子 \begin{array}{rl}
\text { subject to } \quad A y(x)=u(x) \quad \text { in } \Omega \\
\partial_{n} y(x)=0 \quad \text { on } \Gamma \\
\text { and } y_{a}(x) \leq y(x) \leq y_{b}(x) \quad \text { a.e. in } \Omega,
\end{array}\right.
$$

where $\Omega$ is a bounded domain and $\Gamma$ is the boundary of $\Omega$. Moreover, $\partial_{n}=\partial_{\vec{n}}$ denotes directional derivative with respect to the outward unit normal $\vec{n}$ and A is a uniformly elliptic differential operator. The functions $y_{d}, u_{d}, y_{a}$, and $y_{b}$ are given and $\kappa>0$ is a regularization parameter.
The main difficulty of the problem is the presence of pointwise state constraints. It is known from the Karush-Kuhn-Tucker theory in function spaces that the Lagrange multipliers associated with the state constraints are regular Borel measures. This fact is crucial both for the theory and for the the numerical solution.
There are different ideas to deal with the state-constraints numerically. For instance, the problem can be discretized and then solved by a primal-dual active set strategy applied in the finite dimensional space. The efficiency of this technique has been demonstrated by Bergounioux and Kunisch in [3]. On the other hand, interior point methods can be applied to the discretized problem as well, see Haddou et al. [1]. In the case of supremum-norm functional, also Grund and Rösch applied an interior point method to the discrete problem, see [5].

The situation is different when the problem is considered in function spaces. Primaldual active set strategies need the solution of equations such as $(y(u))(x)=d(x)$ on subsets of $\Omega$, where $d(x)=y_{a}(x)$ or $d(x)=y_{b}(x)$. The mapping $u \rightarrow y(u)$ is

[^0]compact, hence these equations for $u$ may cause effects of ill-posedness. It is well known from the theory of inverse problems that a Lavrentiev type regularization of the type $\lambda u+y=d$ is helpful to overcome this difficulty.

This is one reason to approximate the pointwise state constraints in (P) by

$$
\begin{equation*}
y_{a}(x) \leq \lambda u(x)+y(x) \leq y_{b}(x) \tag{1.1}
\end{equation*}
$$

A regularization of this type has several advantages. First, the associated Lagrange multipliers can assumed to be functions of $L^{2}(\Omega)$. This result has been shown for convex elliptic problems for a more general setting including also certain pointwise control constraints by Tröltzsch [10]. For (P), the proof of regularity of Lagrange multipliers is almost trivial, since box constraints on the control are missing, see Section 2 below.
Second, primal-dual active set strategies are be well defined in functions space for this type of regularized constraints. In this way, we are able to directly compare the performance of a primal-dual active set strategy and an interior point method.
Our paper complements the discussion of a semilinear version of $(\mathrm{P})$ in [8], where the existence of regular Lagrange multipliers, second-order sufficient optimality conditions and the application of an SQP method with primal-dual active set strategy for the quadratic subproblems have been discussed for fixed $\lambda>0$. Here, we concentrate on the convergence for $\lambda \downarrow 0$. Moreover, we briefly sketch the implementation of the active set method and an interior point method with classical continuation technique. With that part we continue the work in Prüfert et al. [9] on the applicaton of a classical interior point method in function spaces. The existence of a central path was shown there for a single state constraint. In the case of upper and lower bounds that is given here, the situation is so simple that we present the proof for convenience of the reader.

Throughout this paper, the domain $\Omega$ is a subset of $\mathbb{R}^{n}, n=2,3$, with a $C^{0,1}$-boundary $\Gamma$. As mentioned above, $A$ is an elliptic differential operator. More precisely, it has the form

$$
A y(x)=-\sum_{i, j=1}^{n} D_{i}\left(a_{i j}(x) D_{j} y(x)\right)+c(x) y(x)
$$

where $D_{i}$ denotes the partial derivative with respect to $x_{i}$. Here $c$ is a given function in $L^{\infty}(\Omega)$ with $c(x) \geq 0$ a.e., and $a_{i j} \in L^{\infty}(\Omega), i, j=1, \ldots, n$ satisfy the ellipticity condition

$$
\sum_{i, j=1}^{n} a_{i j}(x) \xi_{i} \xi_{j} \geq \theta|\xi|^{2} \quad \forall(x, \xi) \in \Omega \times \mathbb{R}^{n}
$$

with some positive constant $\theta$. Furthermore, the bounds $y_{a}$ and $y_{b}$ in (P) are fixed functions in $L^{\infty}(\Omega)$ with $y_{b}(x)-y_{a}(x) \geq c_{a d}>0$ a.e. in $\Omega$. The desired state $y_{d}$ and the function $u_{d}$ are defined in $L^{\infty}(\Omega)$.
With (1.1) at hand, we transform (P) into the following optimal control problem

$$
\left(\mathrm{P}_{\lambda}\right)\left\{\begin{align*}
& \text { minimize } J(y, u):=\frac{1}{2} \int_{\Omega}\left(y-y_{d}\right)^{2} d x+\frac{\kappa}{2} \int_{\Omega}\left(u-u_{d}\right)^{2} d x  \tag{1.2}\\
& \text { subject to } \quad A y(x)=u(x) \quad \text { in } \Omega \\
& \partial_{n} y(x)=0 \quad \text { on } \Gamma \\
& \text { and } y_{a}(x) \leq \lambda u(x)+y(x) \leq y_{b}(x) \quad \text { a.e. in } \Omega
\end{align*}\right.
$$

where $\lambda>0$ is a fixed regularization parameter. In the following, we show that $\left(\mathrm{P}_{\lambda}\right)$ admits Lagrange multipliers in $L^{2}(\Omega)$ and that the corresponding solution ( $\bar{u}_{\lambda}, \bar{y}_{\lambda}$ ) converges strongly to the solution of $(\mathrm{P})$ if $\lambda$ converges to zero.
For $n \leq 3,(1.2)$ admits for every $u \in L^{2}(\Omega)$ a unique solution $y \in H^{1}(\Omega) \cap L^{\infty}(\Omega)$ (see for instance [4]). Hence, we may introduce the control-to-state operator $G: L^{2}(\Omega) \rightarrow$ $H^{1}(\Omega) \cap L^{\infty}(\Omega)$ that assigns $y$ to $u$.

Notation. By $\|\cdot\|=\|\cdot\|_{L^{2}(\Omega)}$ and $(.,)=.(., .)_{L^{2}(\Omega)}$ we denote the natural norm and the associated inner product of $L^{2}(\Omega)$, respectively. For the $L^{\infty}(\Omega)$-norm, we abbreviatory write $\|\cdot\|_{\infty}=\|\cdot\|_{L^{\infty}(\Omega)}$. Furthermore, $I: L^{2}(\Omega) \rightarrow L^{2}(\Omega)$ is the identity. Given two normed spaces $U$ and $Y$ and a linear operator $S: U \rightarrow Y$, the associated adjoint operator is denoted by $S^{*}: Y^{*} \rightarrow U^{*}$. Throughout the paper, we say that $u \in$ $L^{2}(\Omega)$ is feasible for (P) if $y_{a}(x) \leq(G u)(x) \leq y_{b}(x)$ holds true a.e. in $\Omega$. Analogously, $u \in L^{2}(\Omega)$ is said to be feasible for $\left(\mathrm{P}_{\lambda}\right)$ if $y_{a}(x) \leq \lambda u(x)+(G u)(x) \leq y_{b}(x)$ is fulfilled a.e. in $\Omega$. By $E_{2}: H^{1}(\Omega) \cap L^{\infty}(\Omega) \rightarrow L^{2}(\Omega)$ we denote the embedding operator of $H^{1}(\Omega) \cap L^{\infty}(\Omega)$ in $L^{2}(\Omega)$, whereas $E_{\infty}$ denotes the analogous embedding operator with range in $L^{\infty}(\Omega)$.
2. First-order optimality conditions. If we consider the state $y$ as a function in $L^{2}(\Omega)$, then the associated solution operator of (1.2) is given by $S:=E_{2} G$. Since $E_{2}$ is compact, the same holds for $S: L^{2}(\Omega) \rightarrow L^{2}(\Omega)$.
The objective functional $f$ is strictly convex and lower semicontinuous. Therfore, the existence of solutions of $(\mathrm{P})$ and $\left(\mathrm{P}_{\lambda}\right)$, respectively, is obtained by standard methods. Moreover, the solutions are unique in both cases. However, considering first-order necessary optimality conditions, both optimal control problems behave different. As mentioned above, the Lagrange multipliers associated to the pure state-constraints in $(\mathrm{P})$ are in general regular Borel measures. Their singular part is concentrated on the boundary of the active set, see Bergounioux and Kunisch [2]. In contrast to that, we are able to prove the existence of regular Lagrange multipliers in $L^{2}(\Omega)$ in the case of $\left(\mathrm{P}_{\lambda}\right)$. To that end, we convert this problem into one with box-constraints on the control by substituting $v=\lambda u+y$. Thanks to the compactness of $S,(\lambda I+S)$ represents a Fredholm operator that has only countably many eigenvalues accumulating at 0 . Moreover, since $S$ is positive definite, the eigenvalues of $-S$ are negative. Thus, for every $\lambda>0$, the theory of Fredholm operators ensures that $(\lambda I+S)$ has a continuous inverse operator $B: L^{2}(\Omega) \rightarrow L^{2}(\Omega)$, i.e.

$$
\begin{equation*}
B v=(\lambda I+S)^{-1} v=u \tag{2.1}
\end{equation*}
$$

Therefore, $\left(\mathrm{P}_{\lambda}\right)$ can be transformed into the following optimization problem with
simple box constraints on the new control $v$

$$
(\mathrm{PV})\left\{\begin{aligned}
\text { minimize } & F(v)=\frac{1}{2}\left\|S B v-y_{d}\right\|^{2}+\frac{\kappa}{2}\left\|B v-u_{d}\right\|^{2} \\
\text { subject to } & y_{a}(x) \leq v(x) \leq y_{b}(x) \quad \text { a.e. in } \Omega
\end{aligned}\right.
$$

Since $F$ is continuously Fréchet-differentiable from $L^{2}(\Omega)$ to $\mathbb{R}$, the Riesz representation theorem implies that its derivative can be identified with a function in $L^{2}(\Omega)$. We denote this function by $g(x)$. Then, by standard arguments, one can show the existence of Lagrange multipliers $\nu_{\lambda}, \mu_{\lambda} \in L^{2}(\Omega)$ that are given by

$$
\begin{aligned}
& \nu_{\lambda}(x)=g(\bar{v})(x)_{+}=\frac{1}{2}(g(\bar{v})(x)+|g(\bar{v})(x)|) \\
& \mu_{\lambda}(x)=g(\bar{v})(x)_{-}=\frac{1}{2}(-g(\bar{v})(x)+|g(\bar{v})(x)|)
\end{aligned}
$$

where $\bar{v}$ denotes the unique optimal solution of (PV). Together with $\nu_{\lambda}$ and $\mu_{\lambda}, \bar{v}$ fulfills the following optimality system:

$$
\left.\begin{array}{c}
S^{*}\left(S B \bar{v}-y_{d}\right)+\kappa\left(B \bar{v}-u_{d}\right)+\left(B^{-1}\right)^{*} \mu_{\lambda}-\left(B^{-1}\right)^{*} \nu_{\lambda}=0  \tag{2.2}\\
\left(\nu_{\lambda}, y_{a}-\bar{v}\right)=\left(\mu_{\lambda}, \bar{v}-y_{b}\right)=0 \\
\nu_{\lambda}(x) \geq 0, \mu_{\lambda}(x) \geq 0, y_{a}(x) \leq \bar{v}(x) \leq y_{b}(x) \quad \text { a.e. in } \Omega .
\end{array}\right\}
$$

Because of the equivalence of (PV) to $\left(\mathrm{P}_{\lambda}\right), \bar{u}_{\lambda}=B \bar{v}$ represents the optimal solution of $\left(\mathrm{P}_{\lambda}\right)$. With $B^{-1}=\lambda I+S$, the first equation in (2.2) is transformed into

$$
S^{*}\left(S \bar{u}_{\lambda}-y_{d}-\nu_{\lambda}+\mu_{\lambda}\right)+\kappa\left(\bar{u}_{\lambda}-u_{d}\right)+\lambda\left(\mu_{\lambda}-\nu_{\lambda}\right)=0
$$

Then, by substituting $\bar{y}_{\lambda}=S \bar{u}_{\lambda}$ and $p_{\lambda}:=S^{*}\left(\bar{y}_{\lambda}-y_{d}-\nu_{\lambda}+\mu_{\lambda}\right)$ in (2.2), we obtain the following optimality system for $\left(\mathrm{P}_{\lambda}\right)$ :

$$
\begin{array}{cc}
A \bar{y}_{\lambda}=\bar{u}_{\lambda} \quad \text { in } \Omega \quad A^{*} p_{\lambda}=\bar{y}_{\lambda}-y_{d}+\mu_{\lambda}-\nu_{\lambda} & \text { in } \Omega \\
\partial_{n} \bar{y}_{\lambda}=0 \quad \text { on } \Gamma \quad \partial_{n} p_{\lambda}=0 & \text { on } \Gamma \\
p_{\lambda}(x)+\kappa\left(\bar{u}_{\lambda}(x)-u_{d}(x)\right)+\lambda\left(\mu_{\lambda}(x)-\nu_{\lambda}(x)\right)=0 & \text { a.e. in } \Omega  \tag{2.3}\\
\left(\nu_{\lambda}, y_{a}-\lambda \bar{u}_{\lambda}-\bar{y}_{\lambda}\right)=\left(\mu_{\lambda}, \lambda \bar{u}_{\lambda}+\bar{y}_{\lambda}-y_{b}\right)=0 & \\
\nu_{\lambda}(x) \geq 0, \mu_{\lambda}(x) \geq 0, y_{a}(x) \leq \lambda \bar{u}_{\lambda}(x)+\bar{y}_{\lambda}(x) \leq y_{b}(x) & \text { a.e. in } \Omega .
\end{array}
$$

Here $A^{*}$ denotes the formal adjoint operator of $A$. The partial differential equation

$$
\begin{align*}
A^{*} p_{\lambda} & =\bar{y}_{\lambda}-y_{d}+\mu_{\lambda}-\nu_{\lambda} & & \text { in } \Omega \\
\partial_{n} p_{\lambda} & =0 & & \text { on } \Gamma \tag{2.4}
\end{align*}
$$

is called adjoint equation. Similar to (1.2), it admits a unique solution in $H^{1}(\Omega) \cap$ $L^{\infty}(\Omega)$ for every right hand side in $L^{2}(\Omega)$. Therefore, due to $\mu_{\lambda}, \nu_{\lambda} \in L^{2}(\Omega)$, we have $p_{\lambda} \in H^{1}(\Omega) \cap L^{\infty}(\Omega)$. As above, we introduce the solution operator to (2.4) by $G^{\circledast}: L^{2}(\Omega) \rightarrow H^{1}(\Omega) \cap L^{\infty}(\Omega)$. Notice that $G^{*}$ would transform $\left(H^{1}(\Omega) \cap L^{\infty}(\Omega)\right)^{*}$ into $L^{2}(\Omega)^{*}$. For the adjoint operator of $S$, we find $S^{*}=E_{2} G^{\circledast}$.
In this way, we have derived the following theorem:

ThEOREM 2.1. Let $\bar{u}_{\lambda}$ be the optimal solution of $\left(\mathrm{P}_{\lambda}\right)$ with associated state $\bar{y}_{\lambda}$, then there exist non-negative Lagrange multipliers $\nu_{\lambda} \in L^{2}(\Omega)$ and $\mu_{\lambda} \in L^{2}(\Omega)$ and an associated adjoint state $p_{\lambda} \in H^{1}(\Omega) \cap L^{\infty}(\Omega)$ such that the optimality system (2.3) is satisfied.
REMARK 2.2. Due to the convexity of the objective functional $J$, the optimality conditions in (2.3) are also sufficient.
3. Pass to the limit. In this section, we prove the convergence of the solutions of the regularized problem $\left(\mathrm{P}_{\lambda}\right)$ to the solution of the original problem (P). The theory is similar to the technique presented in [7]. However, here the situation is a little bit more difficult, since the state is bounded from above and below in our case, whereas in [7] only lower constraints are imposed on the state. This especially complicates the proof of Lemma 3.1 below.
In the following, the unique solution of $(\mathrm{P})$ is denoted by $\bar{u} \in L^{2}(\Omega)$ with associated state $\bar{y}$ and associated adjoint state $p$. Furthermore, we introduce the reduced objective functional $f$ by

$$
f(u):=\frac{1}{2}\left\|S u-y_{d}\right\|^{2}+\frac{\kappa}{2}\left\|u-u_{d}\right\|^{2}
$$

and a function $u_{\lambda}$ that is defined by

$$
\begin{equation*}
u_{\lambda}:=(\lambda I+S)^{-1} \bar{y} \tag{3.1}
\end{equation*}
$$

Notice that $u_{\lambda} \in L^{2}(\Omega)$ is well defined for all $\lambda>0$ because of the compactness of $S$ as described above. The feasibility of $\bar{u}$ for (P) yields $y_{a}(x) \leq \lambda u_{\lambda}(x)+\left(S u_{\lambda}\right)(x) \leq y_{b}(x)$ a.e. and thus, $u_{\lambda}$ is feasible for $\left(\mathrm{P}_{\lambda}\right)$. In the following, we will show that $u_{\lambda}$ converges to $\bar{u}$ as $\lambda \downarrow 0$. To that end, we introduce by $\left\{\lambda_{n}\right\}$ a sequence of positive numbers tending to zero and the sequence $\left\{u_{n}\right\}$ whose elements are defined by $u_{n}=\left(\lambda_{n} I+\right.$ $S)^{-1} \bar{y}$ according to (3.1).
Lemma 3.1. The sequence $\left\{u_{n}\right\}$ converges strongly in $L^{2}(\Omega)$ to $\bar{u}$, as $n \rightarrow \infty$.
Proof: By inserting $\bar{y}=S \bar{u}$ in (3.1), we obtain for a fixed, but arbitrary $\lambda$

$$
\begin{align*}
u_{\lambda}-\bar{u} & =(\lambda I+S)^{-1} S \bar{u}-(\lambda I+S)^{-1}(\lambda I+S) \bar{u} \\
& =(\lambda I+S)^{-1}(S+\lambda I-S) \bar{u}  \tag{3.2}\\
& =\lambda(\lambda I+S)^{-1} \bar{u} .
\end{align*}
$$

The set of eigenvectors of $S$, denoted by $v_{i}, i=1, \ldots, \infty$, represents an orthonormal basis of $L^{2}(\Omega)$. The associated eigenvalues of $S$ are denoted by $\mu_{i}, i \in \mathbb{N}$. We obtain for all $i \in \mathbb{N}$

$$
\begin{aligned}
(\lambda I+S)^{-1} v_{i} & =\frac{1}{\lambda+\mu_{i}}(\lambda I+S)^{-1}\left(\lambda+\mu_{i}\right) v_{i}=\frac{1}{\lambda+\mu_{i}}(\lambda I+S)^{-1}(\lambda I+S) v_{i} \\
& =\frac{1}{\lambda+\mu_{i}} v_{i}
\end{aligned}
$$

Since $\left\{v_{i}\right\}$ is an orthonormal basis of $L^{2}(\Omega)$, we have that $\bar{u}=\sum_{i=1}^{\infty}\left(\bar{u}, v_{i}\right) v_{i}$. There-
fore, (3.2) implies

$$
\begin{aligned}
u_{\lambda}-\bar{u}=\lambda(\lambda I+S)^{-1} \bar{u} & =\lambda \sum_{i=1}^{\infty}\left(\bar{u}, v_{i}\right)(\lambda I+S)^{-1} v_{i} \\
& =\sum_{i=1}^{\infty} \frac{\lambda}{\lambda+\mu_{i}}\left(\bar{u}, v_{i}\right) v_{i},
\end{aligned}
$$

and we obtain for the $L^{2}$-norm of $u_{\lambda}-\bar{u}$

$$
\left\|u_{\lambda}-\bar{u}\right\|^{2}=\left\|\sum_{i=1}^{\infty} \frac{\lambda}{\lambda+\mu_{i}}\left(\bar{u}, v_{i}\right) v_{i}\right\|^{2}=\sum_{i=1}^{\infty}\left(\frac{\lambda}{\lambda+\mu_{i}}\right)^{2}\left(\bar{u}, v_{i}\right)^{2} .
$$

Since $S$ is positive definite, all $\mu_{i}$ are positive. Therefore

$$
\begin{equation*}
\sum_{i=1}^{\infty}\left(\frac{\lambda}{\lambda+\mu_{i}}\right)^{2}\left(\bar{u}, v_{i}\right)^{2} \leq \sum_{i=1}^{\infty}\left(\bar{u}, v_{i}\right)^{2}=\|\bar{u}\|^{2} \tag{3.3}
\end{equation*}
$$

follows from the Bessel inequality. Consider the real valued functions

$$
\varphi_{i}(\lambda):=\left(\frac{\lambda}{\lambda+\mu_{i}}\right)^{2}\left(\bar{u}, v_{i}\right)^{2}
$$

which are continuous and zero at $\lambda=0$. The series $\sum_{i=1}^{\infty}\left(\bar{u}, v_{i}\right)^{2}$ majorizes the one of the left hand side in (3.3), which represents a function series with continuous functions. Thus, we are allowed to interchange summation and pass to the limit and obtain

$$
\lim _{n \rightarrow \infty}\left\|u_{n}-\bar{u}\right\|^{2}=\lim _{n \rightarrow \infty} \sum_{i=1}^{\infty} \varphi_{i}\left(\lambda_{n}\right)=\sum_{i=1}^{\infty} \varphi_{i}(0)=0
$$

Now, let $\left\{\bar{u}_{n}\right\}$ be the sequence of associated optimal solutions of $\left(\mathrm{P}_{\lambda_{n}}\right)$ with associated optimal states $\bar{y}_{n}=S \bar{u}_{n}$. Lemma 3.1 implies $f\left(u_{n}\right) \rightarrow f(\bar{u})$. Hence, the optimality of $\bar{u}_{n}$ and the feasibility of $u_{n}$ for $\left(\mathrm{P}_{\lambda_{n}}\right)$ yields $f\left(\bar{u}_{n}\right) \leq f\left(u_{n}\right) \leq f(\bar{u})+1$ for all sufficiently large $n$. Therefore, we have

$$
\left\|u_{n}\right\|^{2} \leq \frac{2}{\kappa}(f(\bar{u})+1)
$$

giving the uniform boundedness of $\left\{u_{n}\right\}$ in $L^{2}(\Omega)$. Thus we can select a weakly converging subsequence, $\bar{u}_{n_{k}} \rightharpoonup \tilde{u}$. Everything what follows is also valid for any other weakly converging subsequence. Thus, a known argument yields that w.l.o.g. $\bar{u}_{n} \rightharpoonup \tilde{u}$. Lemma 3.2. Let $\tilde{u}$ be the weak limit of $\left\{\bar{u}_{n}\right\}$. Then $\tilde{u}$ is feasible for $(\mathrm{P})$.
Proof: For every $\lambda_{n}>0$, the associated $\bar{u}_{n}$ is feasible for $\left(\mathrm{P}_{\lambda}\right)$ and hence fulfills the constraints

$$
y_{a}(x) \leq \lambda_{n} \bar{u}_{n}(x)+\bar{y}_{n}(x) \leq y_{b}(x) \quad \text { a.e. on } \Omega .
$$

The boundedness of $\left\|\bar{u}_{n}\right\|$ implies $\lambda_{n} \bar{u}_{n} \rightarrow 0$ in $L^{2}(\Omega)$. Furthermore, we have $\bar{y}_{n}=$ $S \bar{u}_{n} \rightarrow S \tilde{u}$ in $L^{2}(\Omega)$ due to the compactness of $S$ and the weak convergence of $\left\{\bar{u}_{n}\right\}$. Therefore, passing to the limit $n \rightarrow \infty, \tilde{u}$ is feasible for (P), i.e.

$$
y_{a}(x) \leq(S \tilde{u})(x) \leq y_{b}(x) \quad \text { a.e. in } \Omega
$$

since the set $\left\{y \in L^{2}(\Omega) \mid y_{a}(x) \leq y(x) \leq y_{b}(x)\right.$ a.e. in $\left.L^{2}(\Omega)\right\}$ is closed.
Now, we are able to prove our main result:
Theorem 3.3. The sequence of optimal solutions $\left\{\bar{u}_{n}\right\}$ of $\left(\mathrm{P}_{\lambda_{n}}\right)$ converges strongly in $L^{2}(\Omega)$ to the solution $\bar{u}$ of $(\mathrm{P})$, i.e.

$$
\bar{u}_{n} \rightarrow \bar{u} \quad, \quad n \rightarrow \infty
$$

Proof: Thanks to Lemma 3.1, i.e. the strong convergence of $u_{n}$ to $\bar{u}$ in $L^{2}(\Omega)$, the states $y_{n}=S u_{n}$ converge strongly in $L^{2}(\Omega)$ to $\bar{y}=S \bar{u}$. This implies

$$
\begin{equation*}
f\left(u_{n}\right) \rightarrow f(\bar{u}) \quad, \quad n \rightarrow \infty . \tag{3.4}
\end{equation*}
$$

Since $u_{n}=\left(\lambda_{n} I+S\right)^{-1} \bar{y}$ is feasible for $\left(\mathrm{P}_{\lambda_{n}}\right)$ and $\bar{u}_{n}$ is the optimal solution of $\left(\mathrm{P}_{\lambda_{n}}\right)$, $f\left(u_{n}\right) \geq f\left(\bar{u}_{n}\right)$ holds true for all $n \in \mathbb{N}$. On the other hand, the feasibilty of $\tilde{u}$ and the optimality of $\bar{u}$ for ( P ) imply $f(\tilde{u}) \geq f(\bar{u})$. Therefore, passing to the limit, (3.4) yields

$$
\begin{equation*}
f(\bar{u})=\lim _{n \rightarrow \infty} f\left(u_{n}\right) \geq \limsup _{n \rightarrow \infty} f\left(\bar{u}_{n}\right) \geq \liminf _{n \rightarrow \infty} f\left(\bar{u}_{n}\right) \geq f(\tilde{u}) \geq f(\bar{u}) \tag{3.5}
\end{equation*}
$$

since $f$ is weakly lower semicontinuous. Thus we get $f(\tilde{u})=f(\bar{u})$ and the strict convexity of $f$ implies

$$
\tilde{u}=\bar{u}
$$

and hence $\bar{u}_{n} \rightharpoonup \bar{u}$.
To show the strong convergence of $\left\{\bar{u}_{n}\right\}$, we will prove the norm convergence of $\left\|\bar{u}_{n}\right\|$ to $\|\bar{u}\|$. It follows from the convergence

$$
\lim _{n \rightarrow \infty} f\left(\bar{u}_{n}\right)=f(\bar{u})
$$

that is obtained from (3.5). Thus, by definition of $f$, we have

$$
\begin{aligned}
\lim _{n \rightarrow \infty}\left\|\bar{u}_{n}\right\|^{2} & =\lim _{n \rightarrow \infty} \frac{2}{\kappa}\left(f\left(\bar{u}_{n}\right)-\frac{1}{2}\left\|\bar{y}_{n}-y_{d}\right\|^{2}\right) \\
& =\frac{2}{\kappa}\left(f(\bar{u})-\frac{1}{2}\left\|\bar{y}-y_{d}\right\|^{2}\right)=\|\bar{u}\|^{2}
\end{aligned}
$$

where we again used $\bar{y}_{n} \rightarrow \bar{y}$ in $L^{2}(\Omega)$. It is well known that weak and norm convergence together yield strong convergence, i.e. $u_{n} \rightarrow \bar{u}$ for $n \rightarrow \infty$.
REmARK 3.4. Clearly, the states $y_{n}=S u_{n}$ converge strongly in $L^{2}(\Omega)$ to $\bar{y}=S \bar{u}$, too.
Next, we consider two different optimization methods for handling the regularized quadratic problem $\left(\mathrm{P}_{\lambda}\right)$ - an active set strategy and an interior point method.
4. Interior point method. This section is devoted to the depiction of an interior point algorithm for the solution of $\left(\mathrm{P}_{\lambda}\right)$. We follow the lines of [9] where the state is only bounded from below. However, here we have upper and lower bounds.

This simplifies the proof of existence of the central path. We think that it is worth to present this easier setting.
The idea of interior point methods is to transform problem $\left(\mathrm{P}_{\lambda}\right)$ into one without inequality contraints. To that end, we penalize the constraints by a logarithmic barrier term. For $\left(\mathrm{P}_{\lambda}\right)$, this amounts

$$
\left(\mathrm{P}_{\lambda}^{\varepsilon}\right)\left\{\begin{aligned}
\text { minimize } J_{\varepsilon}(y, u):= & \frac{1}{2}\left\|y-y_{d}\right\|^{2}+\frac{\kappa}{2}\left\|u-u_{d}\right\|^{2} \\
& -\varepsilon \sum_{i} \int_{\Omega}\left(\ln \left(\lambda u+y-y_{a}\right)+\ln \left(y_{b}-\lambda u-y\right)\right) d x \\
\text { subject to } A y(x)= & u(x) \quad \text { in } \Omega \\
\partial_{n} y(x)= & 0 \quad \text { on } \Gamma,
\end{aligned}\right.
$$

with $\varepsilon>0$. Introducing the solution operator $S=E_{2} G$ as defined in Section 2 and the operator $B$ defined by (2.1), we rewrite $\left(\mathrm{P}_{\lambda}^{\varepsilon}\right)$ as

$$
(\mathrm{Q})\left\{\begin{aligned}
\min F_{\varepsilon}(v):=\frac{1}{2}\left\|S B v-y_{d}\right\|^{2} & +\frac{\kappa}{2}\left\|B v-u_{d}\right\|^{2} \\
& -\varepsilon \int_{\Omega}\left(\ln \left(v-y_{a}\right)+\ln \left(y_{b}-v\right)\right) d x
\end{aligned}\right.
$$

Note that we supress the sub- and superscript and write ( Q ) instead of of $\left(\mathrm{Q}_{\lambda}^{\varepsilon}\right)$ to improve the readability.
The proof of existence of a solution of $(\mathrm{Q})$ is a little bit delicate, since the logarithmic barrier function in $F_{\varepsilon}(v)$ may tend to infinity as $v$ approaches the bounds $y_{a}$ or $y_{b}$. To compensate for this lack of continuity, we first restrict $v$ to a smaller set, where we can prove existence. To that end, we introduce for fixed $\tau>0$ and fixed $\varepsilon>0$ the auxilliary problem

$$
\left(\mathrm{Q}_{\tau}\right) \quad \min _{y_{a}+\tau \leq v \leq y_{b}-\tau} F_{\varepsilon}(v)
$$

where we again supress the indices $\varepsilon$ and $\lambda$. In the following, we will show that, for fixed $\varepsilon$ and $\lambda$ and sufficiently small $\tau$, the unique solution of the auxiliary problem $\left(\mathrm{Q}_{\tau}\right)$, denoted by $v_{\tau}$, represents the solution of $(\mathrm{Q})$, that will be called $v_{\lambda}^{\varepsilon}$ in all what follows.
Theorem 4.1. For all $0<\tau<c_{a d} / 2$ and for all $\varepsilon \geq 0$, problem $\left(\mathrm{Q}_{\tau}\right)$ has a unique solution $v_{\tau}$, and there is a constant $c$ such that $\left\|v_{\tau}\right\|_{\infty} \leq c$.
Proof: The admissible set associated to $\left(\mathrm{Q}_{\tau}\right)$ is defined by

$$
V_{a d}^{\tau}:=\left\{v \in L^{2}(\Omega) \mid y_{a}+\tau \leq v(x) \leq y_{b}-\tau \text { for a.a. } x \in \Omega\right\}
$$

where $\tau<c_{a d} / 2$ ensures that $V_{a d}^{\tau}$ is not empty. We notice that $F_{\varepsilon}$ is strictly convex and continuous on $V_{a d}^{\tau}$, therefore weakly lower semicontinuous. Moreover, $V_{a d}^{\tau}$ is convex, closed and bounded. Therefore, standard arguments show the existence of a unique solution $v_{\tau}$. Moreover, $\left\|v_{\tau}\right\|_{\infty}$ is uniformly bounded, since $y_{a} \leq v_{\tau} \leq y_{b}$ holds for all $\tau<c_{a d} / 2$.

For every $v \in V_{a d}^{\tau}$ and $t \in[0,1]$, the convexity of $V_{a d}^{\tau}$ yields $v_{\tau}+t\left(v-v_{\tau}\right) \in V_{a d}^{\tau}$. Obviously, $F_{\varepsilon}$ is not Gateaux-differentiable in $L^{2}$, since $F_{\varepsilon}(v+h t)$ may be undefined
for some $h \in L^{2}$ even for small $t>0$. However, it is directionally differentiable in the direction $v-v_{\tau}$, since $v_{\tau}+t\left(v-v_{\tau}\right) \in V_{a d}^{\tau}$. The optimality of $v_{\tau}$ gives

$$
\begin{equation*}
\frac{F_{\varepsilon}\left(v_{\tau}+t\left(v-v_{\tau}\right)\right)-F_{\varepsilon}\left(v_{\tau}\right)}{t} \geq 0 \tag{4.1}
\end{equation*}
$$

Passing to the limit $t \downarrow 0$, (4.1) implies for the directional derivative

$$
\begin{equation*}
F_{\varepsilon}^{\prime}\left(v_{\tau}\right)\left(v-v_{\tau}\right) \geq 0 \quad \forall v \in V_{a d}^{\tau} \tag{4.2}
\end{equation*}
$$

By the Riesz representation theorem, $F_{\varepsilon}^{\prime}\left(v_{\tau}\right)$ can be indentified with a function in $L^{2}(\Omega)$ that is denoted by $g_{\varepsilon}(x)$ in all what follows. With the definition of $F_{\varepsilon}$ in (Q) at hand, $g_{\varepsilon}$ is given by

$$
\begin{equation*}
g_{\varepsilon}=(S B)^{*}\left(S B v_{\tau}-y_{d}\right)+\kappa B^{*}\left(B v_{\tau}-u_{d}\right)-\frac{\varepsilon}{v_{\tau}-y_{a}}+\frac{\varepsilon}{y_{b}-v_{\tau}} \tag{4.3}
\end{equation*}
$$

Notice that $g_{\varepsilon}$ is well defined since $y_{a}(x)+\tau \leq v_{\tau}(x) \leq y_{b}(x)-\tau$. Then (4.2) is equivalent to

$$
\begin{equation*}
F_{\varepsilon}^{\prime}\left(v_{\tau}\right)\left(v-v_{\tau}\right)=\int_{\Omega} g_{\varepsilon}(x)\left(v(x)-v_{\tau}(x)\right) d x \geq 0 \quad \forall v \in V_{a d}^{\tau} \tag{4.4}
\end{equation*}
$$

Next, we substitute

$$
\begin{equation*}
p_{\tau}:=(S B)^{*}\left(S B v_{\tau}-y_{d}\right) \quad \text { and } \quad w_{\tau}:=\kappa B^{*}\left(B v_{\tau}-u_{d}\right) . \tag{4.5}
\end{equation*}
$$

Before we use a pointwise evaluation of (4.4) to show that $v_{\tau}=v_{\lambda}^{\varepsilon}$, we need the following Lemma that covers the boundedness of $p_{\tau}$ and $w_{\tau}$.
Lemma 4.2. For all $0<\tau<c_{a d} / 2$, there exist positive constants $c_{1}$ and $c_{2}$ such that $\left\|p_{\tau}\right\|_{\infty} \leq c_{1}$ and $\left\|w_{\tau}\right\|_{\infty} \leq c_{2}$ hold true a.e. in $\Omega$.
Proof: Let $z_{2}=B z_{1}$, then the definition of $B$ in (2.1) implies $z_{1}=S z_{2}+\lambda z_{2}$ and hence $\lambda z_{2}=z_{1}-S z_{2}=(I-S B) z_{1}$. It follows that $z_{2}=\frac{1}{\lambda}(I-S B) z_{1}$ and with $z_{2}=B z_{1}$, we get $B z_{1}=\frac{1}{\lambda}(I-S B) z_{1}$, what yields $B=\frac{1}{\lambda}(I-S B)$. Now let us consider $S$ as an operator with range in $L^{\infty}(\Omega)$ by introducing $\boldsymbol{S}=E_{\infty} G$, i.e. $\boldsymbol{S}: L^{2}(\Omega) \rightarrow L^{\infty}(\Omega)$. Then, together with $B: L^{2} \rightarrow L^{2}$, we obtain

$$
\boldsymbol{B}=\frac{1}{\lambda}(I-\boldsymbol{S} B)
$$

with $\boldsymbol{B}: L^{\infty} \rightarrow L^{\infty}$ continuously. Moreover, the commutativity of $S^{*}$ and $B^{*}$ implies $B^{*}=\frac{1}{\lambda}\left(I-S^{*} B^{*}\right)$. As above, we replace $S^{*}=E_{2} G^{\circledast}$ by $\boldsymbol{S}^{\circledast}=E_{\infty} G^{\circledast}$ with $S^{\circledast}: L^{2}(\Omega) \rightarrow L^{\infty}(\Omega)$. Thus, we may introduce the operator

$$
\boldsymbol{B}^{\circledast}=\frac{1}{\lambda}\left(I-\boldsymbol{S}^{\circledast} B^{*}\right)
$$

with $\boldsymbol{B}^{\circledast}: L^{\infty} \rightarrow L^{\infty}$ continuously. Hence, with the definitions in (4.5) at hand, we obtain

$$
\begin{aligned}
\left\|p_{\tau}\right\|_{\infty} & =\left\|\boldsymbol{B}^{\circledast} \boldsymbol{S}^{\circledast}\left(S B v_{\tau}-y_{d}\right)\right\|_{\infty} \\
& \leq\left\|\boldsymbol{B}^{\circledast}\right\|_{\mathcal{L}\left(L^{\infty}(\Omega)\right)}\left\|\boldsymbol{S}^{\circledast}\right\|_{\mathcal{L}\left(L^{2}(\Omega), L^{\infty}(\Omega)\right)}\left\|S B v_{\tau}-y_{d}\right\| \leq c_{1}
\end{aligned}
$$

and

$$
\begin{aligned}
\left\|w_{\tau}\right\|_{\infty} & =\left\|\kappa \boldsymbol{B}^{\circledast}\left(\boldsymbol{B} v_{\tau}-u_{d}\right)\right\|_{\infty} \\
& \leq \kappa\left\|\boldsymbol{B}^{\circledast}\right\|_{\mathcal{L}\left(L^{\infty}(\Omega)\right)}\left(\|\boldsymbol{B}\|_{\mathcal{L}\left(L^{\infty}(\Omega)\right)}\left\|v_{\tau}\right\|_{\infty}+\left\|u_{d}\right\|_{\infty}\right) \leq c_{2}
\end{aligned}
$$

where we used Theorem 4.1 for the boundedness of $v_{\tau}$.
In preparation of the proof of the next theorem, we define the following sets:

$$
\begin{aligned}
M_{+}(\tau) & :=\left\{x \in \Omega \mid g_{\varepsilon}(x)>0\right\} \\
M_{0}(\tau) & :=\left\{x \in \Omega \mid g_{\varepsilon}(x)=0\right\} \\
M_{-}(\tau) & :=\left\{x \in \Omega \mid g_{\varepsilon}(x)<0\right\}
\end{aligned}
$$

Theorem 4.3. For all sufficiently small $\tau>0$, the solution $v_{\tau}$ of $\left(\mathrm{Q}_{\tau}\right)$ is the unique solution $v_{\lambda}^{\varepsilon}$ of $(\mathrm{Q})$.
Proof: A pointwise evaluation of (4.4) yields

$$
g_{\varepsilon}(x) v_{\tau}(x)=\min _{y_{a}(x)+\tau \leq v \leq y_{b}(x)-\tau} g_{\varepsilon}(x) v
$$

with $v \in \mathbb{R}$. Hence, we have $v_{\tau}(x)=y_{a}(x)+\tau$ for almost all $x \in M_{+}(\tau)$ and $v_{\tau}(x)=y_{b}(x)-\tau$ for almost all $x \in M_{-}(\tau)$. Therefore, with the definition of $M_{+}(\tau)$ and $g_{\varepsilon}$, Lemma 4.2 implies

$$
\begin{equation*}
0<g_{\varepsilon}(x)=p_{\tau}(x)+w_{\tau}(x)-\frac{\varepsilon}{\tau}+\frac{\varepsilon}{y_{b}(x)-y_{a}(x)-\tau} \leq c_{1}+c_{2}-\frac{\varepsilon}{\tau}+\frac{2 \varepsilon}{c_{a d}} \tag{4.6}
\end{equation*}
$$

for almost every $x \in M_{+}(\tau)$. For $\tau \downarrow 0$, the right hand side in (4.6) tends to $-\infty$, a contradiction for sufficiently small $\tau>0$. Similarly, we have on $M_{-}(\tau)$

$$
0>g_{\varepsilon}(x)=p_{\tau}(x)+w_{\tau}(x)-\frac{\varepsilon}{y_{b}(x)-y_{a}(x)-\tau}+\frac{\varepsilon}{\tau} \geq-\left(c_{1}+c_{2}\right)-\frac{2 \varepsilon}{c_{a d}}+\frac{\varepsilon}{\tau} .
$$

Here, the right hand side tends to $\infty$ for $\tau \rightarrow 0$, leading to a contradiction too. Therefore, the sets $M_{+}(\tau)$ and $M_{-}(\tau)$ have measure zero for all sufficiently small $\tau>0$. Hence, if $\tau$ is sufficiently small, we have that $g_{\varepsilon}(x)=0$ holds a.e. on $\Omega$. This implies

$$
\int_{\Omega} g_{\varepsilon}(x) h(x) d x=F_{\varepsilon}^{\prime}\left(v_{\tau}\right) h=0 \quad \forall h \in L^{2}(\Omega)
$$

so that $v_{\tau}$ satisfies the necessary optimality conditions for the unconstrained problem (Q). By convexity, these necessary conditions are also sufficient for optimality. Uniqueness follows from strict convexity.
Remark 4.4. By Theorem 4.3, $\bar{u}_{\lambda}^{\varepsilon}:=B v_{\lambda}^{\varepsilon}$ and $\bar{y}_{\lambda}^{\varepsilon}:=S \bar{u}_{\lambda}^{\varepsilon}$ represent the optimal solution of $\left(\mathrm{P}_{\lambda}^{\varepsilon}\right)$.
In preparation of the numerical computations, we transform the nessesary conditions for (Q) given by

$$
\begin{equation*}
B^{*} S^{*}\left(S B v_{\lambda}^{\varepsilon}-y_{d}\right)+\kappa B^{*}\left(B v_{\lambda}^{\varepsilon}-u_{d}\right)-\frac{\varepsilon}{v_{\lambda}^{\varepsilon}-y_{a}}+\frac{\varepsilon}{y_{b}-v_{\lambda}^{\varepsilon}}=0 \tag{4.7}
\end{equation*}
$$

back to terms of the original problem $\left(\mathrm{P}_{\lambda}^{\varepsilon}\right)$. We apply the operator $\left(B^{*}\right)^{-1}=\left(\lambda I+S^{*}\right)$ to (4.7) and obtain

$$
S^{*}\left(S B v_{\lambda}^{\varepsilon}-y_{d}-\frac{\varepsilon}{v_{\lambda}^{\varepsilon}-y_{a}}+\frac{\varepsilon}{y_{b}-v_{\lambda}^{\varepsilon}}\right)+\kappa\left(B v_{\lambda}^{\varepsilon}-u_{d}\right)-\frac{\lambda \varepsilon}{v_{\lambda}^{\varepsilon}-y_{a}}+\frac{\lambda \varepsilon}{y_{b}-v_{\lambda}^{\varepsilon}}=0 .
$$

Now we substitute

$$
\nu_{\lambda}^{\varepsilon}=\frac{\varepsilon}{y_{b}-v_{\lambda}^{\varepsilon}} \quad \text { and } \quad \mu_{\lambda}^{\varepsilon}=\frac{\varepsilon}{y_{b}-v_{\lambda}^{\varepsilon}} .
$$

Notice that , $\nu_{\lambda}^{\varepsilon}(x)>0$ and $\mu_{\lambda}^{\varepsilon}>0$ hold true almost every where on $\Omega$, because of $y_{a}(x)<v_{\lambda}^{\varepsilon}(x)<y_{b}(x)$. Next, we set

$$
p_{\lambda}^{\varepsilon}=S^{*}\left(S B v_{\lambda}^{\varepsilon}-y_{d}-\frac{\varepsilon}{v_{\lambda}^{\varepsilon}-y_{a}}+\frac{\varepsilon}{y_{b}-v_{\lambda}^{\varepsilon}}\right) .
$$

Then, together with $\bar{y}_{\lambda}^{\varepsilon}=S B v_{\lambda}^{\varepsilon}$ and $\bar{u}_{\lambda}^{\varepsilon} B v_{\lambda}^{\varepsilon}$, we obtain the optimality system to $\left(\mathrm{P}_{\lambda}^{\varepsilon}\right)$ that is given by

$$
\begin{align*}
& A \bar{y}_{\lambda}^{\varepsilon}=\bar{u}_{\lambda}^{\varepsilon} \quad \text { in } \Omega \quad A^{*} p_{\lambda}^{\varepsilon}=\bar{y}_{\lambda}^{\varepsilon}-y_{d}+\mu_{\lambda}^{\varepsilon}-\nu_{\lambda}^{\varepsilon} \quad \text { in } \Omega \\
& \partial_{n} \bar{y}_{\lambda}^{\varepsilon}=0 \quad \text { on } \Gamma \quad \partial_{n} p_{\lambda}^{\varepsilon}=0 \quad \text { on } \Gamma \\
& p_{\lambda}^{\varepsilon}+\kappa\left(\bar{u}_{\lambda}^{\varepsilon}-u_{d}\right)+\lambda\left(\mu_{\lambda}^{\varepsilon}(x)-\nu_{\lambda}^{\varepsilon}(x)\right)=0 \quad \text { a.e. in } \Omega  \tag{4.8}\\
& \nu_{\lambda}^{\varepsilon}(x)\left(\bar{y}_{\lambda}^{\varepsilon}(x)+\lambda \bar{u}_{\lambda}^{\varepsilon}(x)-y_{a}(x)\right)=\varepsilon \quad \text { a.e. in } \Omega \\
& \mu_{\lambda}^{\varepsilon}(x)\left(y_{b}(x)-\bar{y}_{\lambda}^{\varepsilon}(x)-\lambda \bar{u}_{\lambda}^{\varepsilon}(x)\right)=\varepsilon \quad \text { a.e. in } \Omega \text {. }
\end{align*}
$$

4.1. Discretization. We start with the discretization of the state equation (1.2). Let $v \in V$ be an element of the space of test functions $V \subset H^{1}(\Omega)$. Multiplication of (1.2) with $v$ and integration by parts yield

$$
\begin{align*}
-\int_{\Omega} \sum_{i, j=1}^{n} a_{i j}(x) D_{j} y(x) D_{i} v(x)+ & c(x) y(x) v(x) d x  \tag{4.9}\\
& =\int_{\Omega} u(x) v(x) d x \text { for all } v \in V .
\end{align*}
$$

For a given triangulation $\tau_{h}(\Omega)$, we consider a finite dimensional subspace $V_{h}$ of $V$. Let $N$ denote the dimension of $V_{h}\left(\tau_{h}\right)$ and $\left\{\phi_{k}(x)\right\}, k=1, \ldots, N$, a basis of $V_{h}$. Then (4.9) implies that the variational equation is satisfied for all test functions $\phi_{k}(x) \in V_{h}\left(\tau_{h}\right)$, $k=1,2, \ldots, N$, i.e.

$$
\begin{align*}
-\int_{\Omega} \sum_{i, j=1}^{n} a_{i j}(x) D_{j} y(x) D_{i} \phi_{k}(x)+ & c(x) y(x) \phi_{k}(x) d x  \tag{4.10}\\
& =\int_{\Omega} u(x) \phi_{k}(x) d x, k=1, \ldots, N
\end{align*}
$$

For the discretization of (4.10), we discretize $y$ and $u$ by the same basis of $V_{h}$, i.e.

$$
\begin{equation*}
y(x)=\sum_{k=1}^{N} y_{k} \phi_{k}(x) \text { and } u(x)=\sum_{k=1}^{N} u_{k} \phi_{k}(x) . \tag{4.11}
\end{equation*}
$$

Moreover, we define the matrices

$$
\left.\begin{array}{rl}
K_{l k} & =\int_{\Omega} \sum_{i, j=1}^{n} a_{i j}\left(D_{j} \phi_{k}(x)\right) D_{i} \phi_{l}(x) d x \\
M_{l k}^{c} & =\int_{\Omega} c(x) \phi_{k}(x) \phi_{l}(x) d x  \tag{4.12}\\
M_{l k} & =\int_{\Omega} \phi_{k}(x) \phi_{l}(x) d x
\end{array}\right\}
$$

where $K$ is known as the stiffness matrix and $M$ as the Mass matrix. Then, inserting (4.11) in (4.10) yields together with (4.12)

$$
\begin{equation*}
\left(K+M^{c}\right) y_{h}=M u_{h} \tag{4.13}
\end{equation*}
$$

where $y_{h}$ resp. $u_{h}$ are the vectors of the coefficients of $y$ and $u$ with respect to the basis $\phi_{k}, k=1, \ldots, N$. , e.g. $u_{h}=\left(u_{1}, u_{2}, \ldots, u_{N}\right)^{\top}$. Note that for symmetric coefficients $a_{i j}(x)=a_{j i}(x), 1 \leq i j \leq n$ the matrix $K$ is symmetric, too. The adjoint equation in (4.8) is discretized analogously by

$$
\begin{equation*}
\left(K+M^{c}\right)^{\top} p_{h}=M\left(y_{h}-y_{d, h}+\mu_{h}-\nu_{h}\right) \tag{4.14}
\end{equation*}
$$

where $\nu_{h}$ and $\mu_{h}$ reprsent the coefficient vectors of $\mu$ and $\nu$, and $y_{d, h}$ denotes the vector of $y_{d}$ at the nodes of $\tau_{h}$, i.e. $y_{d, h}=\left(y_{d}\left(x_{1}\right), \ldots, y_{d}\left(x_{N}\right)\right)^{\top}$. For a pointwise evaluation of the last two equations in (4.8), we define

$$
\begin{equation*}
\Phi_{a}\left(\nu_{h}\right):=\operatorname{diag}\left(\nu_{h}\right) \quad \text { and } \quad \Psi_{a}\left(u_{h}, y_{h}\right):=\operatorname{diag}\left(y_{h}+\lambda u_{h}-y_{a, h}\right) \tag{4.15}
\end{equation*}
$$

with $\left(\operatorname{diag}\left(v_{h}\right)\right)_{i j}=v_{i} \delta_{i j}$ for an arbitrary $v_{h} \in \mathbb{R}^{N}$. Analogously, $\Phi_{b}$ and $\Psi_{b}$ are defined by

$$
\begin{equation*}
\Phi_{b}\left(\mu_{h}\right):=\operatorname{diag}\left(\mu_{h}\right) \quad \text { and } \quad \Psi_{b}\left(u_{h}, y_{h}\right):=\operatorname{diag}\left(y_{b, h}-y_{h}-\lambda u_{h}\right) \tag{4.16}
\end{equation*}
$$

Here, $y_{a, h}$ and $y_{b, h}$, denote the vector associated to $y_{a}$ and $y_{b}$, respectively at the nodes of $\tau_{h}$. Now we are able to define the finite dimensional approximation of the optimality system (4.8) to $\left(\mathrm{P}_{\lambda}^{\varepsilon}\right)$ : Let $\bar{z}_{h}:=\left(\bar{y}_{h}, \bar{u}_{h}, \bar{p}_{h}, \bar{\nu}_{h}, \bar{\mu}_{h}\right)^{\top} \in \mathbb{R}^{5 N}$ denote the approximation of $\left(\bar{y}_{\lambda}^{\varepsilon}, \bar{u}_{\lambda}^{\varepsilon}, p_{\lambda}^{\varepsilon}, \nu_{\lambda}^{\varepsilon}, \mu_{\lambda}^{\varepsilon}\right)$. Then $\bar{z}_{h}$ satisfies the following nonlinear system of equations

$$
F_{h}\left(\bar{z}_{h} ; \varepsilon\right)=\left(\begin{array}{c}
-\left(K+M^{c}\right) \bar{y}_{h}+M \bar{u}_{h}  \tag{4.17}\\
-\left(K+M^{c}\right)^{\top} \bar{p}_{h}+M\left(\bar{y}_{h}-y_{d, h}+\bar{\mu}_{h}-\bar{\nu}_{h}\right) \\
\bar{p}_{h}+\kappa\left(\bar{u}_{h}-u_{d, h}\right)+\lambda\left(\bar{\mu}_{h}-\bar{\nu}_{h}\right) \\
\Phi_{a}\left(\bar{\nu}_{h}\right)^{\top} \Psi_{a}\left(\bar{u}_{h}, \bar{y}_{h}\right)-\varepsilon \mathbf{1} \\
\Phi_{b}\left(\bar{\nu}_{h}\right)^{\top} \Psi_{b}\left(\bar{u}_{h}, \bar{y}_{h}\right)-\varepsilon \mathbf{1}
\end{array}\right)=0
$$

where $u_{d, h}$ denotes the vector associated to $u_{d}$ at the nodes of $\tau_{h}$ and $\mathbf{1}$ is defined by $\mathbf{1}:=(1)_{i=1}^{N}$. The function $F_{h}$ is continuously differentiable from $\mathbb{R}^{5 N} \times \mathbb{R}_{+}$to $\mathbb{R}^{5 N}$.
4.2. Interior point algorithms. With (4.17) at hand, we are in the position to formulate an interior point algorithm. By $\Delta_{h} z$ we denote the solution of the finite dimensional Newton equation associated to (4.17)

$$
\partial_{z} F_{h}\left(z_{h} ; \varepsilon\right) \Delta_{h} z=-F_{h}\left(z_{h} ; \varepsilon\right)
$$

where $\partial_{z} F_{h}$ denotes the JAcobian of $F_{h}$. With the definitions in (4.15) and (4.16), the Jacobian of $F_{h}$ is given by

$$
\partial_{z} F_{h}\left(z_{h} ; \varepsilon\right)=\left(\begin{array}{ccccc}
M & -\left(K+M^{c}\right) & 0 & 0 & 0  \tag{4.18}\\
0 & M & -\left(K+M^{c}\right)^{\top} & -M & M \\
\kappa I & 0 & I & -\lambda I & \lambda I \\
\lambda \Phi_{a} & \Phi_{a} & 0 & \Psi_{a} & 0 \\
-\lambda \Phi_{b} & -\Phi_{b} & 0 & 0 & \Psi_{b}
\end{array}\right)
$$

where $I$ denotes the $N \times N$-identity matrix. Notice that this Jacobian has a size of $5 N \times 5 N$, is sparse and not symmetric. Moreover, for $\lambda$ and $\kappa$ tending to zero, it tends to be ill conditioned. With (4.18) at hand, the interior point algorithm reads as follows:

Algorithm 1. [Classical continuation]

1. Initalization: choose $0<\sigma<1, \delta>0, \varepsilon^{0}>0$ and choose $z_{h}^{0}$ feasible.
2. For $k=1,2, \ldots$
$\varepsilon^{(k+1)}=\sigma \varepsilon^{(k)}$
solve

$$
\partial_{z} F_{h}\left(z_{h}^{(k)} ; \varepsilon^{(k+1)}\right) \Delta_{h} z^{(k)}=-F_{h}\left(z_{h}^{(k)} ; \varepsilon^{(k+1)}\right)
$$

up to a relative accuracy of

$$
\left\|\Delta_{h} z^{(k)}\right\| \leq \delta
$$

3. $z_{h}^{(k+1)}=z_{h}^{(k)}+\Delta_{h} z^{(k)}$

Algorithm 1 represents the simplest form of an interior point method. In case of a lower state constraint, the convergence of infinite dimensional counterpart of Algorithm 1 was discussed in [9]. There exist several other interior point algorithms for infinite dimensional problems like, for instance, short-step path following algorithms or affine scaling interior point algorithms. For further details, we refer to [13], [12] and [11].
5. Primal-dual active set strategy. This section is concerned with the description of an active set algorithm to solve the optimality system (2.3).
For the derivation of the active set strategy, we need the pointwise form of the complementary slackness condition in (2.3) that is given by

$$
\int_{\Omega} \nu_{\lambda}(x)\left(y_{a}(x)-\lambda \bar{u}_{\lambda}(x)-\bar{y}_{\lambda}(x)\right) d x=\int_{\Omega} \mu_{\lambda}(x)\left(\lambda \bar{u}_{\lambda}(x)+\bar{y}_{\lambda}(x)-y_{b}(x)\right) d x=0
$$

Because of $\nu_{\lambda}(x) \geq 0, \mu_{\lambda}(x) \geq 0$ and $y_{a}(x) \leq \lambda \bar{u}_{\lambda}(x)+\bar{y}_{\lambda}(x) \leq y_{b}(x)$, this implies

$$
\begin{align*}
\nu_{\lambda}(x)\left(y_{a}(x)-\lambda\right. & \left.\bar{u}_{\lambda}(x)-\bar{y}_{\lambda}(x)\right) \\
& =\mu_{\lambda}(x)\left(\lambda \bar{u}_{\lambda}(x)+\bar{y}_{\lambda}(x)-y_{b}(x)\right)=0 \quad \text { a.e. in } \Omega . \tag{5.1}
\end{align*}
$$

Given the optimal solution $\left(\bar{y}_{\lambda}, \bar{u}_{\lambda}\right)$ of $\left(\mathrm{P}_{\lambda}\right)$, we define the active and inactive sets up to sets of measure zero by

$$
\begin{align*}
\mathcal{A}_{a} & :=\left\{x \in \Omega \mid \lambda \bar{u}_{\lambda}(x)+\bar{y}_{\lambda}(x)-\nu_{\lambda}(x)<y_{a}(x)\right\} \\
\mathcal{A}_{b} & :=\left\{x \in \Omega \mid \lambda \bar{u}_{\lambda}(x)+\bar{y}_{\lambda}(x)+\mu_{\lambda}(x)>y_{b}(x)\right\}  \tag{5.2}\\
\mathcal{I} & :=\Omega \backslash\left\{\mathcal{A}_{a} \cup \mathcal{A}_{b}\right\} .
\end{align*}
$$

We rely on the following ASSUMPTION OF STRICT COMPLEMENTARITY:
(S) meas $\left\{x \in \Omega \mid y_{a}(x)-\lambda \bar{u}_{\lambda}(x)-\bar{y}_{\lambda}(x)=\nu_{\lambda}(x)=0\right\}=0$ and

$$
\text { meas }\left\{x \in \Omega \mid \lambda \bar{u}_{\lambda}(x)+\bar{y}_{\lambda}(x)-y_{b}(x)=\mu_{\lambda}(x)=0\right\}=0 .
$$

Under (S), the inequalities in (2.3) can be replaced by equalities on $\mathcal{A}_{a}, \mathcal{A}_{b}$, and $\mathcal{I}$, which is stated by follwoing lemma.
Lemma 5.1. Assume that (S) is fulfilled. Then, it follows that

$$
\begin{aligned}
& \lambda \bar{u}_{\lambda}(x)+\bar{y}_{\lambda}(x)=y_{a}(x), \quad \mu_{\lambda}(x)=0 \quad \text { a.e. on } \mathcal{A}_{a} \\
& \lambda \bar{u}_{\lambda}(x)+\bar{y}_{\lambda}(x)=y_{b}(x), \quad \nu_{\lambda}(x)=0 \quad \text { a.e. on } \mathcal{A}_{b} \\
& \nu_{\lambda}(x)=0, \quad \mu_{\lambda}(x)=0 \quad \text { a.e. on } \mathcal{I} .
\end{aligned}
$$

Proof: For convenience of the reader, we sketch the proof of this well known lemma. We distinct between the following cases:
$x \in \mathcal{A}_{a}$ : On $A_{a}$, we have $\lambda \bar{u}_{\lambda}(x)+\bar{y}_{\lambda}(x)-\nu_{\lambda}(x)<y_{a}(x)$ and hence, the feasibility of $\bar{u}_{\lambda}$ for $\left(\mathrm{P}_{\lambda}\right)$ yields $\nu_{\lambda}(x)>0$. The complemantary slackness condition (5.1) then gives

$$
\lambda \bar{u}_{\lambda}(x)+\bar{y}_{\lambda}(x)=y_{a}(x) \quad \text { and } \quad \mu_{\lambda}(x)=0 \quad \text { a.e. in } \mathcal{A}_{a} .
$$

$x \in \mathcal{A}_{b}$ : In this case, we have $\lambda \bar{u}_{\lambda}(x)+\bar{y}_{\lambda}(x)+\mu_{\lambda}(x)>y_{b}(x)$. Now the feasibility of $\bar{u}_{\lambda}$ for ( $\mathrm{P}_{\lambda}$ ) implies $\mu_{\lambda}(x)>0$, and, due to the complemantary slackness condition (5.1), we obtain

$$
\lambda \bar{u}_{\lambda}(x)+\bar{y}_{\lambda}(x)=y_{b}(x) \quad \text { and } \quad \nu_{\lambda}(x)=0 \quad \text { a.e. in } \mathcal{A}_{b} .
$$

$x \in \mathcal{I}$ : By the definition of $I$ in (5.2), we have $x \notin \mathcal{A}_{a}$ and hence

$$
\begin{equation*}
\lambda \bar{u}_{\lambda}(x)+\bar{y}_{\lambda}(x)-y_{a}(x) \geq \nu_{\lambda}(x) \quad \text { a.e. in } \mathcal{I} . \tag{5.3}
\end{equation*}
$$

Due to the complementary slackness condition (5.1), equality can only occur in (5.3) if $\nu_{\lambda}(x)=\lambda \bar{u}_{\lambda}(x)+\bar{y}_{\lambda}(x)-y_{a}(x)=0$, which contradicts assumption (S). Therefore, the inequality in (5.3) is strict. Thanks to $\nu_{\lambda}(x) \geq 0$, this implies $\lambda \bar{u}_{\lambda}(x)+\bar{y}_{\lambda}(x)>y_{a}(x)$ and hence $\nu_{\lambda}(x)=0$, because of the complementary slackness condition. A similar discussion for $x \notin \mathcal{A}_{b}$ finally gives

$$
\nu_{\lambda}(x)=\mu_{\lambda}(x)=0 \quad \text { a.e. in } \mathcal{I} .
$$

With Lemma 5.1 at hand, the optimality system (2.3) can be transformed into

$$
\begin{array}{ccc}
A \bar{y}_{\lambda}=\bar{u}_{\lambda} & \text { in } \Omega \quad A^{*} p_{\lambda}=\bar{y}_{\lambda}-y_{d}+\mu_{\lambda}-\nu_{\lambda} & \text { in } \Omega \\
\partial_{n} \bar{y}_{\lambda}=0 & \text { on } \Gamma \quad \partial_{n} p_{\lambda}=0 & \text { on } \Gamma \\
p_{\lambda}(x)+\kappa\left(\bar{u}_{\lambda}(x)-u_{d}(x)\right)+\lambda\left(\mu_{\lambda}(x)-\nu_{\lambda}(x)\right)=0 & \text { a.e. in } \Omega  \tag{5.4}\\
\lambda \bar{u}_{\lambda}(x)+\bar{y}_{\lambda}(x)=y_{a}(x), \mu_{\lambda}(x)=0 & \text { a.e. on } \mathcal{A}_{a} \\
\lambda \bar{u}_{\lambda}(x)+\bar{y}_{\lambda}(x)=y_{b}(x), \nu_{\lambda}(x)=0 & \text { a.e. on } \mathcal{A}_{b} \\
\nu_{\lambda}(x)=\mu_{\lambda}(x)=0 & \text { a.e. on } \mathcal{I} .
\end{array}
$$

Discretization of (5.4). As before, the numerical solution is indicated by the subscript $h$. We use the same basis functions for the discretization of $u, y, p, \nu$ and $\mu$ in (4.11). Then the partial differential equations in (5.4) are discretized in the same way as in Section 4.1. Thus, we obtain (4.13) for the discrete version of the state equation and (4.14) for the discrete adjoint equation. A pointwise evaluation of the third equation in (5.4) at the nodes of $\tau_{h}$ yields

$$
\begin{equation*}
\kappa u_{i}+p_{i}+\lambda\left(\mu_{i}-\nu_{i}\right)=\kappa u_{d}\left(x_{i}\right) \quad i=1, \ldots, N . \tag{5.5}
\end{equation*}
$$

For the discretization of the remaining equations in (5.4), we introduce the following index sets that represent the discrete counter parts of the active sets defined in (5.2),

$$
\begin{align*}
\mathcal{A}_{a, h} & :=\left\{i \in\{1, \ldots, N\} \mid \lambda u_{i}+y_{i}-\nu_{i}<y_{a}\left(x_{i}\right)\right\} \\
\mathcal{A}_{b, h} & :=\left\{i \in\{1, \ldots, N\} \mid \lambda u_{i}+y_{i}+\mu_{i}>y_{b}\left(x_{i}\right)\right\}  \tag{5.6}\\
\mathcal{I}_{h} & :=\{1, \ldots, N\} \backslash\left\{\mathcal{A}_{a, h} \cup \mathcal{A}_{b, h}\right\} .
\end{align*}
$$

These definitions allow a pointwise evaluation of the equations on $\mathcal{A}_{a}, \mathcal{A}_{b}$ and $\mathcal{I}$ in (5.4). To that end, we define the matrix $E_{a} \in \mathbb{R}^{N \times N}$ by

$$
E_{a, i j}=\left\{\begin{array}{l}
1, \text { if } i=j \text { and } i \in \mathcal{A}_{a, h} \\
0, \text { otherwise }
\end{array}\right.
$$

and introduce $E_{b}$ analogously. Thus, the pointwise discrete version of the equation $\lambda \bar{u}_{\lambda}(x)+\bar{y}_{\lambda}(x)=y_{a}(x)$ a.e. on $\mathcal{A}_{a}$ is given by

$$
\begin{equation*}
E_{a}\left(\lambda u_{h}+y_{h}\right)=E_{a} y_{a, h} . \tag{5.7}
\end{equation*}
$$

Similarly, the equation $\nu_{\lambda}(x)=0$ a.e. on $\mathcal{A}_{b} \cup \mathcal{I}$ is discretized by

$$
\begin{equation*}
\left(I-E_{a}\right) \nu_{h}=0 . \tag{5.8}
\end{equation*}
$$

An addition of (5.7) and (5.8) yields

$$
\begin{equation*}
E_{a}\left(\lambda u_{h}+y_{h}\right)+\left(I-E_{a}\right) \nu_{h}=E_{a} y_{a, h} \tag{5.9}
\end{equation*}
$$

Together with an analogous equation for $\mathcal{A}_{b, h}$, the discrete versions of the PDEs, and (5.5), we obtain the following $5 N \times 5 N$-linear system of equations

$$
\left(\begin{array}{ccccc}
M & -\left(K+M^{c}\right) & 0 & 0 & 0  \tag{5.10}\\
0 & M & -\left(K+M^{c}\right)^{\top} & -M & M \\
\kappa I & 0 & I & -\lambda I & \lambda I \\
\lambda E_{a} & E_{a} & 0 & I-E_{a} & 0 \\
\lambda E_{b} & E_{b} & 0 & 0 & I-E_{b}
\end{array}\right)\left(\begin{array}{c}
u_{h} \\
y_{h} \\
p_{h} \\
\nu_{h} \\
\mu_{h}
\end{array}\right)=\left(\begin{array}{c}
0 \\
M y_{d, h} \\
\kappa u_{d, h} \\
E_{a} y_{a, h} \\
E_{b} y_{b, h}
\end{array}\right)
$$

that represents the discrete version of (5.4). Notice that the coefficient matrix in (5.10) has a same structure as the Jacobian in (4.18) arising from the interior point method. Similar to the matrix in (4.18), it tends to be ill-conditioned for $\lambda, \kappa \downarrow 0$.

Active set algorithm. The primal dual active set algorithm proceeds as follows. We denote by $w_{h}$ the solution vector of (5.10), i.e. $w_{h}=\left(u_{h}, y_{h}, p_{h}, \nu_{h}, \mu_{h}\right)$.

## Algorithm 2.

1. Define initial sets $\mathcal{A}_{a, h}^{(0)} \subset\{1, \ldots, N\}$ and $\mathcal{A}_{b, h}^{(0)} \subset\{1, \ldots, N\}$ with $\mathcal{A}_{a, h}^{(0)} \cap \mathcal{A}_{b, h}^{(0)}=$ $\emptyset$. Set $\mathcal{I}_{h}^{(0)}=\{1, \ldots, N\} \backslash\left\{\mathcal{A}_{b, h}^{(0)} \cup \mathcal{A}_{a, h}^{(0)}\right\}$ and $k=0$.
2. Find $w_{h}^{(k)}$ by solving (5.10).
3. Set

$$
\begin{aligned}
\mathcal{A}_{a, h}^{(k+1)} & =\left\{i \in\{1, \ldots, N\} \mid \lambda u_{i}^{(k)}+y_{i}^{(k)}-\nu_{i}^{(k)}<y_{a}\left(x_{i}\right)\right\} \\
\mathcal{A}_{b, h}^{(k+1)} & =\left\{i \in\{1, \ldots, N\} \mid \lambda u_{i}^{(k)}+y_{i}^{(k)}+\mu_{i}^{(k)}>y_{b}\left(x_{i}\right)\right\} \\
\mathcal{I}_{h}^{(k+1)} & :=\{1, \ldots, N\} \backslash\left\{\mathcal{A}_{a, h}^{(k+1)} \cup \mathcal{A}_{b, h}^{(k+1)}\right\} .
\end{aligned}
$$

4. If $\mathcal{A}_{a, h}^{(k+1)}=\mathcal{A}_{a, h}^{(k)}$ and $\mathcal{A}_{b, h}^{(k+1)}=\mathcal{A}_{b, h}^{(k)}$ then STOP, else:

Update $k=k+1$ and goto 2.
The termination condition in step 4 is justified by the following theorem. We introduce the discrete version of the optimality system (2.3) with the complementary slackness condition in the pointwise form (5.1) that is given by

$$
\begin{gather*}
\left(K+M^{c}\right) \bar{y}_{h}=M \bar{u}_{h} \quad\left(K+M^{c}\right)^{\top} \bar{p}_{h}=M\left(\bar{y}_{h}-y_{d, h}+\bar{\mu}_{h}-\bar{\nu}_{h}\right) \\
\kappa \bar{u}_{h}+\bar{p}_{h}+\lambda\left(\bar{\mu}_{h}-\bar{\nu}_{h}\right)=0  \tag{5.11}\\
\bar{\nu}_{i}\left(\bar{y}_{a}\left(x_{i}\right)-\lambda \bar{u}_{i}-\bar{y}_{i}\right)=\bar{\mu}_{i}\left(\lambda \bar{u}_{i}+\bar{y}_{i}-y_{b}\left(x_{i}\right)\right)=0 \quad, i=1, \ldots, N \\
y_{a}\left(x_{i}\right) \leq \lambda \bar{u}_{i}+\bar{y}_{i} \leq y_{b}\left(x_{i}\right), \bar{\nu}_{i} \geq 0, \bar{\mu}_{i} \geq 0 \quad, i=1, \ldots, N,
\end{gather*}
$$

where $\bar{u}_{h}, \bar{y}_{h}, \bar{p}_{h}, \bar{\nu}_{h}$ and $\bar{\mu}_{h}$ again denote the discret optimal solution.
ThEOREM 5.2. If $\mathcal{A}_{a, h}^{(k+1)}=\mathcal{A}_{a, h}^{(k)}$ and $\mathcal{A}_{b, h}^{(k+1)}=\mathcal{A}_{b, h}^{(k)}$ for some $k \in \mathbb{N}$ then the associated solution of (5.10), denoted by $w_{h}^{(k)}$, satisfies the discrete optimality system (5.11).

For the proof of this theorem, we refer to results of Kunisch and Rösch [6], that can easily be adapted to our case.
6. Numerical tests. We tested the two algorithms by two examples. Generally, we consider the following optimal control problem

$$
(\mathrm{PT})\left\{\begin{aligned}
\text { minimize } & J(y, u):=\frac{1}{2}\left\|y-y_{d}\right\|^{2}+\frac{\kappa}{2}\left\|u-u_{d}\right\|^{2} \\
\text { subject to }-\Delta y(x)+y(x)=u(x) & \text { in } \Omega \\
\partial_{n} y(x)=0 & \text { on } \Gamma \\
\text { and } h(y) \geq 0 & \text { a.e. in } \Omega,
\end{aligned}\right.
$$

with $h(y)=\left(y-y_{a}, y_{b}-y\right)^{\top}$ in the first example and $h(y)=y_{b}-y$ in the second one. This problem fits into our problem setting with $A y=-\Delta y+y$. In both examples, we take the unit circle $B(0,1) \subset \mathbb{R}^{2}$ for the domain $\Omega$. The associated exact solutions are given in polar coordinates and only depend on the radius that is given by $r=$ $\|x\|_{2}=\sqrt{x_{1}^{2}+x_{2}^{2}}$.
6.1. Example with regular Lagrange multipliers in $L^{2}(\Omega)$. In the first example, $h$ is given by $h(y)=\left(y-y_{a}, y_{b}-y\right)^{\top}$ and thus hence we consider (PT) with lower and upper state constraints, i.e.

$$
y_{a}(x) \leq y(x) \leq y_{b}(x) \quad \text { a.e. in } \Omega .
$$

The Lagrange multipliers associated to such constraints are in general Borel measures, with a singular part concentrated on the boundaries of the active sets, see Bergouniuox and Kunisch [2]. Since, in our examples, the boundaries of the active sets do not coincide with $\Gamma$, the optimality system is given by

$$
\begin{align*}
& -\Delta \bar{y}+\bar{y}=\bar{u} \quad \text { in } \Omega \quad-\Delta p+p=\bar{y}-y_{d}+\mu-\nu \quad \text { in } \Omega \\
& \partial_{n} \bar{y}=0 \quad \text { on } \Gamma \quad \partial_{n} p=0 \quad \text { on } \Gamma \\
& p(x)+\kappa\left(\bar{u}(x)-u_{d}(x)\right)=0 \quad \text { a.e. in } \Omega \\
& \int_{\Omega}\left(y_{a}-\bar{y}\right) d \nu=\int_{\Omega}\left(\bar{y}-y_{b}\right) d \mu=0  \tag{6.1}\\
& \nu \geq 0, \mu \geq 0, y_{a}(x) \leq \bar{y}(x) \leq y_{b}(x) \quad \text { a.e. in } \Omega .
\end{align*}
$$

In this example, we construct $\mu$ and $\nu$ such that $d \nu=\nu(x) d x$ and $d \mu=\mu(x) d x$ with nonnegative functions $\mu, \nu \in L^{\infty}(\Omega)$.
Choosing $\bar{y}(r)=-r^{6}+3 r^{4}-3 r^{2}+1$ for the optimal state, the state equation in (6.1) implies

$$
\bar{u}(r)=-\Delta \bar{y}(r)+\bar{y}(r)=-r^{6}+39 r^{4}-51 r^{2}+13
$$

To fulfill the state constraints, we define

$$
y_{a}(r)=\left\{\begin{array}{l}
\bar{y}(r), \bar{y}(r) \leq c_{a} \\
c_{a},
\end{array} \quad \bar{y}(r)>c_{a} \quad \text { and } \quad y_{b}(r)=\left\{\begin{array}{ll}
\bar{y}(r), & \bar{y}(r) \geq c_{b} \\
c_{b} & , \bar{y}(r)<c_{b}
\end{array},\right.\right.
$$

with $c_{a}=0.3$ and $c_{b}=0.7$. Furthermore, with

$$
\nu(r)=\left\{\begin{array}{ll}
c_{a}-\bar{y}(r)+1, & \bar{y}(r) \leq c_{a} \\
0 & , \bar{y}(r)>c_{a}
\end{array} \quad \text { and } \quad \mu(r)=\left\{\begin{array}{ll}
\bar{y}(r)-c_{b}+1, & \bar{y}(r) \geq c_{b} \\
0 & , \bar{y}(r)<c_{b}
\end{array},\right.\right.
$$

the complementary slackness condition in (6.1) are satisfied. For these Lagrange multipliers, we have $\mu_{a}, \mu_{b} \in L^{\infty}(\Omega) \subset L^{2}(\Omega)$. Therefore, the complementary slackness conditions in (6.1) can be replaced by $\left(\nu, \bar{y}-y_{b}\right)=\left(\mu, y_{a}-y\right)=0$. Moreover, we define the adjoint state by

$$
p(r)=r^{4}-2 r+1
$$

Notice that both $p$ and $\bar{y}$ fulfill the homogeneous Neumann boundary conditions. To satisfy the adjoint equation, $y_{d}$ must be defined by

$$
\begin{aligned}
y_{d}(r) & =\bar{y}(r)+\Delta p(r)-p(r)+\mu(r)-\nu(r) \\
& = \begin{cases}-2 r^{6}+5 r^{4}+10 r^{2}+2 r-7-c_{a}, & \bar{y}(r) \leq c_{a} \\
-r^{6}+2 r^{4}+13 r^{2}+2 r-8 \\
-2 r^{6}+5 r^{4}+10 r^{2}+2 r-7-c_{b}, & , \bar{y}(r) \geq c_{a}<\end{cases}
\end{aligned}
$$

Finally, the optimality condition gives

$$
\begin{aligned}
u_{d}(r) & =\bar{u}(r)+\frac{1}{\kappa} p(r) \\
& =-r^{6}+\left(39+\frac{1}{\kappa}\right) r^{4}-51 r^{2}-\frac{2}{\kappa} r+13+\frac{1}{\kappa}
\end{aligned}
$$

The functions $y_{d}, u_{d}, y_{a}$, and $y_{b}$ are shown in Figures 6.1-6.3.

6.2. Example with Lagrange multiplier in $C^{*}(\bar{\Omega})$. In this example, only the upper state constraint is imposed on $y$, i.e.

$$
y(x) \leq y_{b}(x) \quad \text { a.e. in } \Omega
$$

In this case, the optimality system reads as follows

$$
\begin{align*}
& -\Delta \bar{y}+\bar{y}=\bar{u} \quad \text { in } \Omega \quad-\Delta p+p=\bar{y}-y_{d}+\mu \quad \text { in } \Omega \\
& \partial_{n} \bar{y}=0 \quad \text { on } \Gamma \quad \partial_{n} p=0 \quad \text { on } \Gamma \\
& p(x)+\kappa\left(\bar{u}(x)-u_{d}(x)\right)=0 \quad \text { a.e. in } \Omega \\
& \int_{\Omega}\left(\bar{y}-y_{b}\right) d \mu=0  \tag{6.2}\\
& \mu \geq 0, \bar{y}(x) \leq y_{b}(x) \quad \text { a.e. in } \Omega .
\end{align*}
$$

To construct an example with $\mu \in C^{*}(\bar{\Omega})$, we consider the fundamential solution $\Phi$ of Poisson's equation in $\mathbb{R}^{2}$,

$$
\Phi(r):=-\frac{1}{2 \pi} \log (r)
$$

for $r>0$. It is known that in $\mathbb{R}^{2}$

$$
-\Delta \Phi=\delta_{0},
$$

where $\delta_{0}$ denotes the Dirac measure on $\mathbb{R}^{2}$ concentrated in $r=0$. Notice that $\delta_{0} \in$ $C^{*}(\bar{\Omega})$ but $\delta_{0} \notin H^{1}(\Omega)^{*}$. With the fundamental solution, the optimal adjoint state is given by

$$
p(r)=\frac{1}{4 \pi} r^{2}+\Phi(r)=\frac{1}{4 \pi} r^{2}-\frac{1}{2 \pi} \log (r) .
$$

One can easily verify that $p$ satisfies the homogeneous Neumann boundary conditions on $\Gamma=\partial B(0,1)$. Moreover, we set

$$
\bar{y} \equiv 4 \quad \text { and } \quad \bar{u}=-\Delta \bar{y}+\bar{y} \equiv 4 .
$$

The upper bound in the state constraint is defined by

$$
y_{b}(r)=r+4 .
$$

Therefore, the optimal state touches the bound only in the point $r=0$, see also Figure 6.6. Hence, a possible Lagrange multiplier, satisfying the complementary slackness conditions, is given by

$$
\mu=\delta_{0}
$$

and thus $\mu$ represents a regular Borel measure. From the adjoint equation, we get

$$
y_{d}(r)=\bar{y}(r)+\Delta p(r)-p(r)+\mu=4+\frac{1}{\pi}-\frac{1}{4 \pi} r^{2}+\frac{1}{2 \pi} \log (r) .
$$

Finally, the optimality condition implies

$$
u_{d}(r)=\bar{u}(r)+\frac{1}{\kappa} p(r)=4+\frac{1}{4 \pi \kappa} r^{2}-\frac{1}{2 \pi \kappa} \log (r) .
$$

Figures 6.4 and 6.5 show the desired state $y_{d}$ and the control shift $u_{d}$ for this example.

6.3. Numerical results. Each algorithm was tested by both examples with nine different values of $\lambda$ each. For the numerical investigations, we used unstructured grids that were refined at the boundaries of the active sets. In the first example, the mesh was additionally refined at $\partial B$ and, in the second test case, at $r=0$. All computations were performed using Matlab on a PC with a 2.8 GHz processor.

Example 1. In the first example, described in Section 6.1, the Tikhonov regularization parameter was fixed at $\kappa=10^{-4}$. Figures $6.7-6.11$ show the numerical solution computed by the active set algorithm on a grid with $\mathrm{N}=29272$ nodes and $\lambda=10^{-4}$. Here and in the following, the superscript "as" marks results that were computed with Algorithm 2, whereas resuts of the Algorithms 1 are denoted by the superscript "ip".


Fig. 6.7. Control $u_{h}^{\text {as }}$


Fig. 6.9. Lagrange multiplier $\mu_{h}^{\text {as }}$


Fig. 6.11. Adjoint state $p_{h}^{\text {as }}$

The figures show that, the numerical errors in $u_{h}^{\text {as }}, \nu_{h}^{\text {as }}$, and $\mu_{h}^{\text {as }}$ are quite large compared with the errors in $y_{h}^{\mathrm{as}}$ and $p_{h}^{\mathrm{as}}$. This is also visible in the Tables 6.1 and 6.2. A possible explanation is that $y_{h}$ and $p_{h}$ are smooth as the discrete solutions of linear PDEs.

To describe the accuracy of the algorithms for $\lambda \downarrow 0$, the relative errors of $u, y, p$, and the Lagrange multipliers are displayed in the Tables 6.1-6.4. For the control, the
relative error used here is defined by

$$
e_{u}:=\frac{\left\|\bar{u}-\bar{u}_{h}\right\|}{\|\bar{u}\|} \approx \sqrt{\frac{\left(\overline{\boldsymbol{u}}-\overline{\boldsymbol{u}}_{\boldsymbol{h}}\right)^{\top} M\left(\overline{\boldsymbol{u}}-\overline{\boldsymbol{u}}_{\boldsymbol{h}}\right)}{\overline{\boldsymbol{u}}^{\top} M \overline{\boldsymbol{u}}}} .
$$

Here, $\bar{u}$ denotes the exact optimal control, $\bar{u}_{h}$ the discrete optimal control, and $\overline{\boldsymbol{u}}$ and $\overline{\boldsymbol{u}}_{\boldsymbol{h}}$, respectively, the vector of values at the nodes of $\tau_{h}$, i.e. for instance $\overline{\boldsymbol{u}}=$ $\left(\bar{u}\left(x_{1}\right), \ldots, \bar{u}\left(x_{N}\right)\right)^{\top}$. The errors $e_{y}, e_{p}, e_{\nu}$, and $e_{\mu}$ are defined analogously.
As an indicator for the performance of the algorithms, we used the parameter \#es that denotes the number of linear systems of equations that have to be solved during the respective iterations. Since the coefficient matrices that are defined in (4.18) and (5.10), respectively, possess the same size and a quite similar structure and the solution of the associated linear systems of equations represent the main effort of both algorithms, \#es represents a suitable value to compare the different algorithms.

Table 6.1
Example 1: Interior point algorithm with $N=29272$

| $\lambda$ | $\# \mathrm{es}$ | $e_{u}^{\mathrm{ip}}$ | $e_{y}^{\mathrm{ip}}$ | $e_{p}^{\mathrm{ip}}$ | $e_{\nu}^{\mathrm{ip}}$ | $e_{\mu}^{\mathrm{ip}}$ |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $1 \mathrm{e}-2$ | 9 | $7.0644 \mathrm{e}-01$ | $1.0898 \mathrm{e}-01$ | $2.2005 \mathrm{e}-01$ | $1.4511 \mathrm{e}-01$ | $2.9594 \mathrm{e}-01$ |
| $1 \mathrm{e}-3$ | 13 | $3.3134 \mathrm{e}-01$ | $1.4279 \mathrm{e}-02$ | $2.3352 \mathrm{e}-02$ | $4.8361 \mathrm{e}-02$ | $7.7850 \mathrm{e}-02$ |
| $1 \mathrm{e}-4$ | 18 | $4.5059 \mathrm{e}-02$ | $1.5656 \mathrm{e}-03$ | $2.3878 \mathrm{e}-04$ | $4.1617 \mathrm{e}-02$ | $5.2152 \mathrm{e}-02$ |
| $1 \mathrm{e}-5$ | 20 | $3.8185 \mathrm{e}-02$ | $5.1407 \mathrm{e}-04$ | $4.5432 \mathrm{e}-05$ | $9.1956 \mathrm{e}-02$ | $1.5128 \mathrm{e}-01$ |
| $1 \mathrm{e}-6$ | 20 | $3.8571 \mathrm{e}-02$ | $4.4392 \mathrm{e}-04$ | $4.0144 \mathrm{e}-05$ | $1.0915 \mathrm{e}-01$ | $1.8786 \mathrm{e}-01$ |
| $1 \mathrm{e}-7$ | 20 | $3.8609 \mathrm{e}-02$ | $4.3790 \mathrm{e}-04$ | $4.0231 \mathrm{e}-05$ | $1.1103 \mathrm{e}-01$ | $1.9190 \mathrm{e}-01$ |
| $1 \mathrm{e}-8$ | 20 | $3.8613 \mathrm{e}-02$ | $4.3731 \mathrm{e}-04$ | $4.0247 \mathrm{e}-05$ | $1.1122 \mathrm{e}-01$ | $1.9230 \mathrm{e}-01$ |
| $1 \mathrm{e}-9$ | 20 | $3.8613 \mathrm{e}-02$ | $4.3725 \mathrm{e}-04$ | $4.0248 \mathrm{e}-05$ | $1.1124 \mathrm{e}-01$ | $1.9234 \mathrm{e}-01$ |
| 0.0 | 20 | $3.8613 \mathrm{e}-02$ | $4.3724 \mathrm{e}-04$ | $4.0249 \mathrm{e}-05$ | $1.1124 \mathrm{e}-01$ | $1.9235 \mathrm{e}-01$ |

Table 6.2
Example 1: Active set algorithm with $N=29272$

| $\lambda$ | $\# \mathrm{es}$ | $e_{u}^{\text {as }}$ | $e_{y}^{\text {as }}$ | $e_{p}^{\text {as }}$ | $e_{\nu}^{\text {as }}$ | $e_{\mu}^{\text {as }}$ |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $1 \mathrm{e}-2$ | 4 | $7.1296 \mathrm{e}-01$ | $1.0893 \mathrm{e}-01$ | $2.2004 \mathrm{e}-02$ | $1.4513 \mathrm{e}-01$ | $2.9596 \mathrm{e}-01$ |
| $1 \mathrm{e}-3$ | 7 | $3.3598 \mathrm{e}-01$ | $1.4291 \mathrm{e}-02$ | $2.3351 \mathrm{e}-03$ | $4.8289 \mathrm{e}-02$ | $7.7192 \mathrm{e}-02$ |
| $1 \mathrm{e}-4$ | 11 | $4.5674 \mathrm{e}-02$ | $1.5158 \mathrm{e}-03$ | $2.3883 \mathrm{e}-04$ | $4.7154 \mathrm{e}-02$ | $5.8081 \mathrm{e}-02$ |
| $1 \mathrm{e}-5$ | 23 | $4.6090 \mathrm{e}-02$ | $4.7228 \mathrm{e}-04$ | $4.9732 \mathrm{e}-05$ | $3.1060 \mathrm{e}-01$ | $4.2250 \mathrm{e}-01$ |
| $1 \mathrm{e}-6$ | 33 | $4.8178 \mathrm{e}-02$ | $4.0856 \mathrm{e}-04$ | $4.9575 \mathrm{e}-05$ | $4.9601 \mathrm{e}-01$ | $6.3486 \mathrm{e}-01$ |
| $1 \mathrm{e}-7$ | 33 | $4.8302 \mathrm{e}-02$ | $4.0350 \mathrm{e}-04$ | $5.0274 \mathrm{e}-05$ | $5.2879 \mathrm{e}-01$ | $6.6125 \mathrm{e}-01$ |
| $1 \mathrm{e}-8$ | 33 | $4.8314 \mathrm{e}-02$ | $4.0301 \mathrm{e}-04$ | $5.0353 \mathrm{e}-05$ | $5.3240 \mathrm{e}-01$ | $6.6408 \mathrm{e}-01$ |
| $1 \mathrm{e}-9$ | 33 | $4.8316 \mathrm{e}-02$ | $4.0296 \mathrm{e}-04$ | $5.0361 \mathrm{e}-05$ | $5.3277 \mathrm{e}-01$ | $6.6436 \mathrm{e}-01$ |
| 0.0 | 33 | $4.8316 \mathrm{e}-02$ | $4.0296 \mathrm{e}-04$ | $5.0361 \mathrm{e}-05$ | $5.3281 \mathrm{e}-01$ | $6.6439 \mathrm{e}-01$ |

The Tables 6.1 and 6.2 show that both algorithms achieve a similar accuracy even though the interior point method averages slighty smaller errors. In both methods, the errors of all quantities are significantly reduced from $\lambda=10^{-2}$ to $\lambda=10^{-4}$, but stagnate or are even increased for smaller values of $\lambda$. Especially $e_{\nu}^{\text {as }}$ and $e_{\mu}^{\text {as }}$ increase up to errors of $53 \%$ and $66 \%$, respectively. However, considering the results for $\lambda=10^{-4}$, both algorithms provide errors lower than $6 \%$ also for $\nu$ and $\mu$. In this
sense, a choice of $\lambda=10^{-4}$ seems to be optimal for both algorithms if a sufficiently accurate approximation of all quantities including the Lagrange multipliers is desired. A further decrease of $\lambda$ only improves $e_{y}$ and $e_{p}$ significantly, but $s$ the errors of the discrete Lagrange mutlipliers.

We observe that the iteration numbers and thus number \#es of solved linear systems of equations increase with a reduction of $\lambda$. However, similar to the development of the errors, \#es remains static for $\lambda \leq 10^{-5}$ in case of the interior point method and for $\lambda \leq 10^{-6}$ in case of the active set algorithm. The range between the minimal and maximal number of solved linear systems varies between 9 and 20 for Algorithm 1 and between 4 and 33 for Algorithm 2. Thus the interior point method seems to be less sensitive with respect to the regularization parameter $\lambda$ than the active set algorithm. On the other hand, for the optimal value $\lambda=10^{-4}$, the active set algorithm is slightly more efficient since \#es amounts 11 in this case, whereas 18 linear systems of equations have to be solved in the interior point interation.
Examle 2. As described in Section 6.2, the Lagrange multiplier in the second example is the Dirac measure. For the computations, we fixed $\kappa=1.0$ and used a mesh with 21993 nodes that was refined at $r=0$ to deal with the singularity of $p$ and $\mu$ at this point. Figures 6.12-6.15 show the numerical solution for $\lambda=10^{-4}$.


FIG. 6.12. Control $u_{h}^{\text {ip }}$


Fig. 6.14. Adjoint state $p_{h}^{\mathrm{ip}}$

FIG. 6.13. State $y_{h}^{\text {ip }}$


Fig. 6.15. Lagrange multiplier $\mu_{h}^{\text {ip }}$

We observe that the Lagrange multiplier approximates the Dirac measure. As in the first example, the two algorithms were compared by the relative errors and the number
\#es of solved linear systems of equations. Since the exact Lagrange multiplier does not belong to $L^{2}(\Omega)$, Tables 6.3 and 6.4 only contain $e_{u}, e_{y}$, and $e_{p}$.

Table 6.3
Example 1: Interior point algorithm with $N=21993$

| $\lambda$ | $\# \mathrm{es}$ | $e_{u}^{\mathrm{ip}}$ | $e_{y}^{\mathrm{ip}}$ | $e_{p}^{\mathrm{ip}}$ |
| :---: | :---: | :---: | :---: | :---: |
| $1 \mathrm{e}-02$ | 30 | $1.9825 \mathrm{e}-02$ | $1.8589 \mathrm{e}-03$ | $1.0539 \mathrm{e}-01$ |
| $1 \mathrm{e}-03$ | 28 | $8.9583 \mathrm{e}-03$ | $1.4051 \mathrm{e}-03$ | $7.5260 \mathrm{e}-02$ |
| $1 \mathrm{e}-04$ | 26 | $7.2597 \mathrm{e}-04$ | $4.4813 \mathrm{e}-05$ | $5.6160 \mathrm{e}-02$ |
| $1 \mathrm{e}-05$ | 24 | $1.6252 \mathrm{e}-03$ | $3.3719 \mathrm{e}-05$ | $5.6161 \mathrm{e}-02$ |
| $1 \mathrm{e}-06$ | 25 | $1.8601 \mathrm{e}-03$ | $1.8320 \mathrm{e}-05$ | $5.6162 \mathrm{e}-02$ |
| $1 \mathrm{e}-07$ | 25 | $1.8836 \mathrm{e}-03$ | $1.6640 \mathrm{e}-05$ | $5.6163 \mathrm{e}-02$ |
| $1 \mathrm{e}-08$ | 25 | $1.8859 \mathrm{e}-03$ | $1.6471 \mathrm{e}-05$ | $5.6163 \mathrm{e}-02$ |
| $1 \mathrm{e}-09$ | 25 | $1.8862 \mathrm{e}-03$ | $1.6454 \mathrm{e}-05$ | $5.6163 \mathrm{e}-02$ |
| 0.0 | 25 | $1.8862 \mathrm{e}-03$ | $1.6452 \mathrm{e}-05$ | $5.6163 \mathrm{e}-02$ |

Table 6.4
Example 2: Active set algorithm with $N=21993$

| $\lambda$ | $\#$ es | $e_{u}^{\text {as }}$ | $e_{y}^{\text {as }}$ | $e_{p}^{\text {as }}$ |
| :---: | :---: | :---: | :---: | :---: |
| $1 \mathrm{e}-2$ | 8 | $1.9826 \mathrm{e}-02$ | $1.8593 \mathrm{e}-03$ | $1.0539 \mathrm{e}-01$ |
| $1 \mathrm{e}-3$ | 12 | $8.9583 \mathrm{e}-03$ | $1.4053 \mathrm{e}-03$ | $7.5259 \mathrm{e}-02$ |
| $1 \mathrm{e}-4$ | 21 | $7.2613 \mathrm{e}-04$ | $4.4649 \mathrm{e}-05$ | $5.6158 \mathrm{e}-02$ |
| $1 \mathrm{e}-5$ | 17 | $1.6252 \mathrm{e}-03$ | $3.3565 \mathrm{e}-05$ | $5.6162 \mathrm{e}-02$ |
| $1 \mathrm{e}-6$ | 28 | $1.8600 \mathrm{e}-03$ | $1.8167 \mathrm{e}-05$ | $5.6162 \mathrm{e}-02$ |
| $1 \mathrm{e}-7$ | 75 | $1.8834 \mathrm{e}-03$ | $1.6487 \mathrm{e}-05$ | $5.6162 \mathrm{e}-02$ |
| $1 \mathrm{e}-8$ | 83 | $1.8858 \mathrm{e}-03$ | $1.6318 \mathrm{e}-05$ | $5.6162 \mathrm{e}-02$ |
| $1 \mathrm{e}-9$ | 75 | $1.8861 \mathrm{e}-03$ | $1.6301 \mathrm{e}-05$ | $5.6162 \mathrm{e}-02$ |
| 0.0 | 78 | $1.8861 \mathrm{e}-03$ | $1.6300 \mathrm{e}-05$ | $5.6162 \mathrm{e}-02$ |

As a solution of a PDE, $p$ is smooth. Nevertheless, the error $e_{p}$ is significantly larger than $e_{u}$ and $e_{y}$ in both algorithms. A possible explanation for this fact could be that the exact solutions $\bar{y}=\bar{u} \equiv 4$ are identically constant. Hence, the state equation is exactly satisfied by $\bar{y}, \bar{u}$. also in the finite dimensional setting.
In this example, the two algorithms behave similarly to the first test case. The difference in the accuracy of both algorithms is marginal, since the relative errors are nearly identical. As above, we observe that the errors stagnate or even increase if $\lambda \leq 10^{-5}$ in case of $u_{h}$ and $p_{h}$ and $\lambda \leq 10^{-7}$ in case of $y_{h}$. Concerning the control $u_{h}$, the best approximation is achieved for $\lambda=10^{-4}$ in both algorithms.

The performance of the algorithms is similar to the first example. Again the active set algorithm is more sensitive with respect to $\lambda$ than the interior point method. For $\lambda \downarrow 0$, \#es increase significantly in the active set algorithm, while the effort of the interior point algorithm remains nearly constant. In contrast to this, the active set algorithm requires less iterations than the interior point method for larger values of $\lambda$. This is also true for $\lambda=10^{-4}$, where the best approximation of $u_{h}$ is achieved with both methods.

Comparing the accuracy of the two methods, the difference between both methods is negligible. However, they slightly differ in the peformance: the interior point method is less sensitive to $\lambda$, whereas the number of iterations of the active set algorithm increases as $\lambda \downarrow 0$. The active set algorithm is less expensive than the interior point algorithm for larger values of $\lambda$, i.e. $\lambda \geq 10^{-4}$ in the first and $\lambda \geq 10^{-5}$ in the second example. Larger values of $\lambda$ lead to a better approximation of the Lagrange multipliers in the first example and to the control in the second example. This shows the benefit of the regularization of pointwise state constraints.
Acknowledgement. The authors are grateful to Prof. B. Hofmann (TU Chemnitz) for pointing out the idea of proving Lemma 3.1.

## REFERENCES

[1] M. Bergounioux, M. Haddou, M. Hintermüller, and K. Kunisch. A comparison of a moreauyosida based active strategy and interior point methods for constrained optimal control problems. SIAM J. on Optimization, 11:495-521, 2000.
[2] M. Bergounioux and K. Kunisch. On the structure of the lagrange multiplier for stateconstrained optimal control problems. Systems and Control Letters, 48:169-176, 2002.
[3] M. Bergounioux and K. Kunisch. Primal-dual active set strategy for state constrained optimal control problems. Computational Optimization and Applications, 22:193-224, 2002.
[4] E. Casas. Control of an elliptic problem with pointwise state constraints. SIAM J. Control and Optimization, 4:1309-1322, 1986.
[5] T. Grund and A. Rösch. Optimal control of a linear elliptic equation with a supremum-norm functional. Opt. Meth. Software, 15:299-329, 2001.
[6] K. Kunisch and A. Rösch. Primal-dual strategy for constrained parabolic optimal control problems. SIAM Journal on Optimization, 13:321-334, 2002.
[7] C. Meyer, A. Rösch, and F. Tröltzsch. Optimal control problems of PDEs with regularized pointwise state constraints. Preprint 14, Inst. of Math., TU Berlin, 2004. To appear in Computational Optimization and Applications.
[8] C. Meyer and F. Tröltzsch. On an elliptic optimal control problem with pointwise mixed control-state constraints. In A. Seeger, editor, Recent Advances in Optimization., Lectures Notes in Economics and Mathematical Systems, 2004. To appear.
[9] U. Prüfert, F. Tröltzsch, and M. Weiser. The convergence of an interior point method for an elliptic control problem with mixed control-state contraints. Preprint 36, Inst. of Math., TU Berlin, 2004.
[10] F. Tröltzsch. Regular lagrange multipliers for control problems with mixed pointwise controlstate constraints. SIAM Journal on Optimization, 2003. To appear.
[11] M. Ulbrich and S. Ulbrich. Superlinear convergence of affine-scaling interior-point newton methods for infinite-dimensional nonlinear problems with pointwise bounds. SIAM Journal on Control and Optimization, 38(6):1938-1984, 2000.
[12] M. Ulbrich, S. Ulbrich, and M. Heinkenschloss. Global convergence of trust-region interior-point algorithms for infinite-dimensional nonconvex minimization subject to pointwise bounds. SIAM J. Control Optim, 37:731-764, 1999.
[13] M. Weiser. Interior point methods in function space. Technical Report 03-35, Zuse Inst. Berlin, 2003.


[^0]:    *Supported by the DFG Research Center Matheon "Mathematics for key technologies" in Berlin.
    ${ }^{1}$ Institut für Mathematik, Technische Universität Berlin, D-10623 Berlin, Str. des 17. Juni 136, Germany.

