# OPTIMAL CONTROL OF PDEs IN A COMPLEX SPACE SETTING: APPLICATION TO THE SCHRÖDINGER EQUATION* 

M. SOLEDAD ARONNA ${ }^{\dagger}$, JOSEPH FRÉDÉRIC BONNANS ${ }^{\ddagger}$, AND AXEL KRÖNER $\S$


#### Abstract

In this paper we discuss optimality conditions for abstract optimization problems over complex spaces. We then apply these results to optimal control problems with a semigroup structure. As an application we detail the case when the state equation is the Schrödinger one, with pointwise constraints on the "bilinear" control. We derive first and second order optimality conditions and address, in particular, the case that the control enters affine in the cost function.


Key words. optimal control, partial differential equations, optimization in complex Banach spaces, second order optimality conditions, Goh transform, semigroup theory, Schrödinger equation, bilinear control systems

AMS subject classifications. 49J20, 49K20, 35J10, 93C20
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1. Introduction. In this paper we derive no-gap second order optimality conditions for optimal control problems in a complex Banach space setting with pointwise constraints on the control. This general framework includes, in particular, optimal control problems for the bilinear Schrödinger equation.

Let us consider $T>0, \Omega \subset \mathbb{R}^{n}$ an open bounded set, $n \in \mathbb{N}, Q:=(0, T) \times \Omega$. The Schrödinger equation is given by

$$
\begin{equation*}
i \dot{\Psi}(t, x)+\Delta \Psi(t, x)-u(t) B(x) \Psi(t, x)=0, \quad \Psi(x, 0)=\Psi_{0}(x) \tag{1.1}
\end{equation*}
$$

where $t \in(0, T), x \in \Omega$, and with $u:[0, T] \rightarrow \mathbb{R}$ the amplitude of the time-dependent electric field, $\Psi:[0, T] \times \Omega \rightarrow \mathbb{C}$ the wave function, and $B: \Omega \rightarrow \mathbb{R}$ the spatial profile. The system describes the position probability distribution of a quantum particle subject to the electric field, that will be considered as the control throughout this paper. The wave function $\Psi$ belongs to the unitary sphere in $L^{2}(\Omega ; \mathbb{C})$.

For $\alpha_{1} \in \mathbb{R}$ and $\alpha_{2} \geq 0$, the optimal control problem is given as

$$
\left\{\begin{align*}
\min J(u, \Psi):= & \frac{1}{2} \int_{\Omega}\left|\Psi(T)-\Psi_{d T}\right|^{2} \mathrm{~d} x+\frac{1}{2} \int_{Q}\left|\Psi-\Psi_{d}\right|^{2} \mathrm{~d} x \mathrm{~d} t  \tag{1.2}\\
& +\int_{0}^{T}\left(\alpha_{1} u(t)+\frac{1}{2} \alpha_{2} u(t)^{2}\right) \mathrm{d} t, \text { subject to (1.1) and } u \in \mathcal{U}_{\mathrm{ad}}
\end{align*}\right.
$$

[^0]with $\mathcal{U}_{\mathrm{ad}}:=\left\{u \in L^{\infty}(0, T): u_{m} \leq u(t) \leq u_{M}\right.$ a.e. in $\left.(0, T)\right\}, u_{m}, u_{M} \in \mathbb{R}, u_{m}<u_{M}$, and $|z|:=\sqrt{z \bar{z}}$ for $z \in \mathbb{C}$, and desired running and final states $\Psi_{d}:(0, T) \times \Omega \rightarrow \mathbb{C}$ and $\Psi_{d T}: \Omega \rightarrow \mathbb{C}$, respectively. The control of the Schrödinger equation is an important question in quantum physics. For the optimal control of semigroups, the reader is referred to Li amd Yao [38] and Li and Yong [39] and Fattorini and Frankowska [30] and Fattorini [29], and Goldberg and Tröltzsch [34]. In the context of optimal control of partial differential equations for systems in which the control enters affine in the cost function (we speak of control-affine problems), in a companion paper [5], we have extended the results of Bonnans [17] (about necessary and sufficient second order optimality conditions for a bilinear heat equation) to problems governed by general bilinear systems in a real Banach space setting, and presented applications for the heat and wave equation.

The contributions of this paper are as follow: (i) We extend to a complex Banach space setting the theory of optimality conditions (Bonnans and Shapiro [20, Chap. 2]) for an abstract optimization problem. (ii) We then turn to optimal control problems for a semigroup formulation of a dynamical system, and thanks to the complex structure, we express in a compact way the first order optimality conditions, especially the costate equation. (iii) We derive second order necessary and sufficient conditions, using the technique of Bonnans and Osmolovskiĭ [19]. (iv) In the case of problems with the Hamiltonian affine w.r.t. the control, we extend the second order necessary and sufficient conditions obtained in [5]. (v) The results are applied to the Schrödinger equation.

While the literature on optimal control of the heat equation is quite rich (see, e.g., the monograph by Tröltzsch [45]), much less is available for the optimal control of the Schrödinger equation. We list some references on optimal control of the Schrödinger equation and related topics. In Ito and Kunisch [36] necessary optimality conditions are derived and an algorithm is presented to solve the unconstrained problem; in Baudouin, Kavian, and Puel [9] regularity results for the Schrödinger equation with a singular potential are presented; further regularity results can be found in Baudouin and Salomon [10] and Boscain, Caponigro, and Sigalotti [22] and, in particular, in Ball, Marsden, and Slemrod [7]. For a minimum time problem and controllability problems for the Schrödinger equation see Beauchard et al. [14, 15, 13]. For second order analysis for control problems of control-affine ordinary differential systems see [1, 33]. About the case of optimal control of nonlinear Schrödinger equations of Gross-Pitaevskii type arising in the description of Bose-Einstein condensates, see Hintermüller et al. [35]; for sparse controls in quantum systems see Friesecke, Henneke, and Kunisch [32].

The paper is organized as follows. In section 2 necessary optimality conditions for general minimization problems in complex Banach spaces are formulated. In section 3 the abstract control problem is introduced in a semigroup setting and some basic calculus rules are established. In section 4 first order optimality conditions and in section 5 sufficient second order optimality conditions are presented; sufficient second order optimality conditions for singular problems are presented in section 6 , again in a general semigroup setting. Finally section 7 presents the application of the previous results to the control of the Schrödinger equation.

## 2. Optimality conditions in complex spaces.

2.1. Real and complex spaces. We consider complex Banach spaces which can be identified with the product of two identical real Banach spaces. That is, with a real Banach space $X$ we associate the complex Banach space $\bar{X}$ with elements $x_{c}$ represented in a unique way as $x_{c}=x_{1}+i x_{2}$, with $x_{1}, x_{2}$ in $X$ and $i=\sqrt{-1}$, and
the usual computing rules for complex variables, in particular, for $\gamma=\gamma_{1}+i \gamma_{2} \in \mathbb{C}$ with $\gamma_{1}$, $\gamma_{2}$ real, we define $\gamma x_{c}=\gamma_{1} x_{1}-\gamma_{2} x_{2}+i\left(\gamma_{2} x_{1}+\gamma_{1} x_{2}\right)$. We define the real and imaginary parts of $x_{c} \in \bar{X}$ by $\Re x_{c}:=x_{1}$ and $\Im x_{c}:=x_{2}$, respectively.

Let $X$ be a real Banach space and $\bar{X}$ the corresponding complex one. The dual (resp., antidual) of $X$ (resp., $\bar{X}$ ), i.e., the set of linear (resp., antilinear) forms, is denoted by $X^{*}$ (resp., $\bar{X}^{*}$ ). We denote by $\left\langle x^{*}, x\right\rangle_{X}$ the duality product between $x^{*} \in X^{*}$ and $x \in X$, and by $\left\langle x_{c}^{*}, x_{c}\right\rangle_{\bar{X}}$ the antiduality product (linear w.r.t. the first argument, and antilinear w.r.t. the second) between $x_{c}^{*} \in \bar{X}^{*}$ and $x_{c} \in \bar{X}$. Let $x:=\left(x_{1}, x_{2}\right) \in X \times X, x^{*}:=\left(x_{1}^{*}, x_{2}^{*}\right) \in X^{*} \times X^{*}$. Setting $x_{c}:=x_{1}+i x_{2}$ and $x_{c}^{*}:=x_{1}^{*}+i x_{2}^{*}$ observe that, due to linearity/antilinearity of $\langle\cdot, \cdot\rangle_{\bar{X}}$,

$$
\begin{equation*}
\left\langle x_{c}^{*}, x_{c}\right\rangle_{\bar{X}}=\left\langle x_{1}^{*}, x_{1}\right\rangle_{X}+\left\langle x_{2}^{*}, x_{2}\right\rangle_{X}+i\left(\left\langle x_{2}^{*}, x_{1}\right\rangle_{X}-\left\langle x_{1}^{*}, x_{2}\right\rangle_{X}\right), \tag{2.1}
\end{equation*}
$$

and therefore the "real" duality product in $X \times X$ given by

$$
\begin{equation*}
\left\langle x^{*}, x\right\rangle_{X \times X}:=\left\langle x_{1}^{*}, x_{1}\right\rangle_{X}+\left\langle x_{2}^{*}, x_{2}\right\rangle_{X} \tag{2.2}
\end{equation*}
$$

satisfies

$$
\begin{equation*}
\left\langle x^{*}, x\right\rangle_{X \times X}=\Re\left\langle x_{c}^{*}, x_{c}\right\rangle_{\bar{X}} . \tag{2.3}
\end{equation*}
$$

In what follows we drop the index $c$ for complex valued elements of Banach spaces.
2.2. First order optimality conditions in abstract optimization. We next address the questions of optimality conditions analogous to the ones obtained in the case of real Banach spaces [20]. Consider the problem

$$
\begin{equation*}
\operatorname{Min}_{u, x} f(u, x) ; \quad g(u, x) \in K_{g} ; \quad h(u, x) \in K_{h} . \tag{2.4}
\end{equation*}
$$

Here $U$ and $W$ are real Banach spaces, $\bar{X}$ and $\bar{Y}$ are complex Banach spaces, and $K_{g}$, $K_{h}$ are nonempty, closed convex subsets of $\bar{Y}$ and $W$, respectively. The mappings $f$, $g, h$ from $U \times \bar{X}$ to, respectively, $\mathbb{R}, \bar{Y}$, and $W$ are of class $C^{1}$. As said above, the complex space $\bar{X}$ is identified with the product $X \times X$ of real Banach spaces with dual $X^{*} \times X^{*}$.

We recall that the normal cone to the convex set $K_{h}$ at the point $\hat{w} \in K_{h}$ is defined by

$$
\begin{equation*}
N_{K_{h}}(\hat{w}):=\left\{w^{*} \in W^{*} ;\left\langle w^{*}, w-\hat{w}\right\rangle_{W} \leq 0 \text { for all } w \in K_{h}\right\} . \tag{2.5}
\end{equation*}
$$

The corresponding expression of normal cones in a complex setting is, say for $y_{c} \in K_{g}$,

$$
\begin{equation*}
N_{K_{g}}\left(y_{c}\right):=\left\{y_{c}^{*} \in \bar{Y}^{*} ; \Re\left\langle y_{c}^{*}, z_{c}-y_{c}\right\rangle_{\bar{Y}} \leq 0 \text { for all } z_{c} \in K_{g}\right\} . \tag{2.6}
\end{equation*}
$$

Let $\bar{X}, \bar{Y}$ be two complex spaces associated with the real Banach spaces $X$ and $Y$. The conjugate transpose of $A_{c} \in \mathcal{L}(\bar{X}, \bar{Y})$ is the operator $A_{c}^{*} \in \mathcal{L}\left(\bar{Y}^{*}, \bar{X}^{*}\right)$ defined by

$$
\begin{equation*}
\left\langle y_{c}^{*}, A_{c} x_{c}\right\rangle_{\bar{Y}}=\left\langle A_{c}^{*} y_{c}^{*}, x_{c}\right\rangle_{\bar{X}} \text { for all }\left(x_{c}, y_{c}^{*}\right) \text { in } \bar{X} \times \bar{Y}^{*} \tag{2.7}
\end{equation*}
$$

If $A_{c}=A_{1}+i A_{2}$ with $A_{1}$ and $A_{2}$ in $L(X, Y)$, then $A_{c}^{*}=A_{1}^{\top}-i A_{2}^{\top}$, where $\top$ denotes the transpose operator. The extension of $A \in L(U, \bar{Y})$ is $A_{c} \in L(\bar{U}, \bar{Y})$ defined by

$$
\begin{equation*}
A_{c}\left(u_{1}+i u_{2}\right):=A u_{1}+i A u_{2} . \tag{2.8}
\end{equation*}
$$

Then for $u \in U$ and $y_{c}^{*} \in \bar{Y}^{*}$ and using (2.7) we get that

$$
\begin{equation*}
\Re\left\langle y_{c}^{*}, A_{c} u\right\rangle_{\bar{Y}}=\Re\left\langle A_{c}^{*} y_{c}^{*}, u\right\rangle_{\bar{U}}=\left\langle\Re A_{c}^{*} y_{c}^{*}, u\right\rangle_{U} \tag{2.9}
\end{equation*}
$$

Coming back to problem (2.4), for $\lambda \in \bar{Y}^{*}$ and $\mu \in W^{*}$, the Lagrangian of the problem
is defined as

$$
\begin{equation*}
L(u, x, \lambda, \mu):=f(u, x)+\Re\langle\lambda, g(u, x)\rangle_{\bar{Y}}+\langle\mu, h(u, x)\rangle_{W} \tag{2.10}
\end{equation*}
$$

Lemma 2.1. The partial derivatives of the Lagrangian are as follows:

$$
\left\{\begin{align*}
\frac{\partial L}{\partial u} & =\frac{\partial f}{\partial u}+\Re\left(\frac{\partial g}{\partial u}^{*} \lambda\right)+{\frac{\partial h^{\prime}}{\partial u}}^{\top} \mu  \tag{2.11}\\
\frac{\partial L}{\partial x_{r}} & =\frac{\partial f}{\partial x_{r}}+\Re\left({\frac{\partial g^{*}}{\partial x}}^{*}\right)+{\frac{\partial h^{\prime}}{\partial x_{r}}}^{\top} \mu \\
\frac{\partial L}{\partial x_{i}} & =\frac{\partial f}{\partial x_{i}}+\Im\left(\frac{\partial g}{\partial x}^{*} \lambda\right)+{\frac{\partial h^{\top}}{\partial x_{i}}}^{\top} \mu
\end{align*}\right.
$$

Proof. Set $L^{\prime}(u, x, \lambda):=\Re\langle\lambda, g(u, x)\rangle_{\bar{Y}}$. It is enough to express the partial derivatives of this expression. We have that, skipping arguments,

$$
\begin{equation*}
\frac{\partial L^{\prime}}{\partial u} v=\Re\left\langle\lambda, \frac{\partial g}{\partial u} v\right\rangle_{\bar{Y}}=\Re\left\langle{\frac{\partial g^{*}}{\partial u}}^{*}, v\right\rangle_{U}=\left\langle\Re\left(\frac{\partial g}{\partial u}^{*} \lambda\right), v\right\rangle_{U} \tag{2.12}
\end{equation*}
$$

for all $v \in U$. We have used that setting $\frac{\partial g}{\partial u}=a+i b$ and $\lambda=\lambda_{r}+i \lambda_{i}$, then

$$
\begin{align*}
\left\langle\Re \left(\frac{\partial g}{}_{\partial u}\right.\right. & \lambda), v\rangle_{U} \tag{2.13}
\end{align*}=\left\langle\Re\left(a^{\top}-i b^{\top}\right)\left(\lambda_{r}+i \lambda_{i}\right), v\right\rangle_{U} .
$$

Now, for $z_{r} \in X$,

$$
\begin{equation*}
\frac{\partial L^{\prime}}{\partial x_{r}} z_{r}=\Re\left\langle\lambda, \frac{\partial g}{\partial x} z_{r}\right\rangle_{\bar{Y}}=\Re\left\langle\frac{\partial g}{\partial x}^{*} \lambda, z_{r}\right\rangle_{\bar{X}}=\left\langle\Re\left(\frac{\partial g}{\partial x}^{*} \lambda\right), z_{r}\right\rangle_{\bar{X}} \tag{2.14}
\end{equation*}
$$

and for all $z_{i} \in X$,
(2.15)
$\frac{\partial L^{\prime}}{\partial x_{i}} z_{i}=\Re\left\langle\lambda, \frac{\partial g}{\partial x_{i}} i z_{i}\right\rangle_{\bar{Y}}=-\Re\left\langle i \frac{\partial g^{*}}{\partial x_{i}} \lambda, z_{i}\right\rangle_{\bar{X}}=\Im\left\langle\frac{\partial g^{*}}{\partial x} \lambda, z_{i}\right\rangle_{\bar{X}}=\left\langle\Im\left(\frac{\partial g^{*}}{\partial x} \lambda\right), z_{i}\right\rangle_{\bar{X}}$.
The result follows.
Definition 2.2. (i) Let $(\hat{u}, \hat{x}) \in U \times \bar{X}$ satisfy the constraints of problem (2.4). Then we say then that $(\hat{u}, \hat{x})$ is a feasible point for (2.4).
(ii) An element $(\lambda, \mu)$ of $\bar{Y}^{*} \times W^{*}$ is called Lagrange multiplier associated with $(\hat{u}, \hat{x})$, if the following conditions are verified:

$$
\left\{\begin{array}{l}
D_{u} L(\hat{u}, \hat{x}, \lambda, \mu)=0, \quad D_{x} L(\hat{u}, \hat{x}, \lambda, \mu)=0  \tag{2.16}\\
\lambda \in N_{K_{g}}(g(\hat{u}, \hat{x})), \quad \mu \in N_{K_{h}}(h(\hat{u}, \hat{x}))
\end{array}\right.
$$

We call (2.16) the first order optimality system of problem (2.4).
(iii) Let $\mathbb{B}$ denote the unit ball of $\bar{Y} \times W$. A feasible point $(\hat{u}, \hat{x})$ of (2.4) is said to be qualified if, for some $\varepsilon>0$,
$\varepsilon \mathbb{B} \subset K_{g} \times K_{h}-(g(\hat{u}, \hat{x}), h(\hat{u}, \hat{x}))-\{D(g(\hat{u}, \hat{x}), h(\hat{u}, \hat{x}))(u-\hat{u}, x) ;(u, x) \in U \times \bar{X}\}$.

Lemma 2.3. Let $(\hat{u}, \hat{x})$ be a qualified local solution of problem (2.4), that is, ( $\hat{u}, \hat{x}$ ) is qualified and

$$
\begin{equation*}
f(\hat{u}, \hat{x}) \leq f(u, x) \text { for all feasible }(u, x), \text { close enough to }(\hat{u}, \hat{x}) \tag{2.18}
\end{equation*}
$$

Then with $(\hat{u}, \hat{x})$ is associated a nonempty and bounded set of Lagrange multipliers.
Proof. This is just an adaptation of the classical result in real spaces, due to Robinson [43]; see also Bonnans and Shapiro [20, Chap. 2].
3. The abstract control problem in a semigroup setting. Given a complex and reflexive Banach space $\overline{\mathcal{H}}$, we consider optimal control problems for equations of type

$$
\begin{equation*}
\dot{\Psi}+\mathcal{A} \Psi=f+u\left(\mathcal{B}_{1}+\mathcal{B}_{2} \Psi\right), \quad t \in(0, T), \quad \Psi(0)=\Psi_{0} \tag{3.1}
\end{equation*}
$$

where

$$
\begin{equation*}
\Psi_{0} \in \overline{\mathcal{H}}, \quad f \in L^{1}(0, T ; \overline{\mathcal{H}}), \quad \mathcal{B}_{1} \in \overline{\mathcal{H}}, \quad u \in L^{1}(0, T), \quad \mathcal{B}_{2} \in \mathcal{L}(\overline{\mathcal{H}}) \tag{3.2}
\end{equation*}
$$

and $\mathcal{A}$ is the generator of a strongly continuous semigroup on $\overline{\mathcal{H}}$, in the sense that, denoting by $e^{-t \mathcal{A}}$ the semigroup generated by $\mathcal{A}$, we have that

$$
\begin{equation*}
\operatorname{dom}(\mathcal{A}):=\left\{y \in \overline{\mathcal{H}} ; \quad \lim _{t \downarrow 0} \frac{y-e^{-t \mathcal{A}} y}{t} \text { exists }\right\} \tag{3.3}
\end{equation*}
$$

is dense and for $y \in \operatorname{dom}(\mathcal{A}), \mathcal{A} y$ is equal to the above limit. Then $\mathcal{A}$ is closed. Note that we choose to define $\mathcal{A}$ and not its opposite as the generator of the semigroup. We have then

$$
\begin{equation*}
\left\|e^{-t \mathcal{A}}\right\|_{\mathcal{L}(\overline{\mathcal{H}})} \leq c_{\mathcal{A}} e^{\lambda_{\mathcal{A}} t}, \quad t>0 \tag{3.4}
\end{equation*}
$$

for some positive $c_{\mathcal{A}}$ and $\lambda_{\mathcal{A}}$. For the semigroup theory in a complex space setting we refer to Dunford and Schwartz [28, Chap. VIII]. The mild solution of (3.1) is the function $\Psi \in C(0, T ; \overline{\mathcal{H}})$ such that, for all $t \in[0, T]$,

$$
\begin{equation*}
\Psi(t)=e^{-t \mathcal{A}} \Psi_{0}+\int_{0}^{t} e^{-(t-s) \mathcal{A}}\left(f(s)+u(s)\left(\mathcal{B}_{1}+\mathcal{B}_{2} \Psi(s)\right)\right) \mathrm{d} s \tag{3.5}
\end{equation*}
$$

This fixed-point equation (3.5) is well-posed in the sense that it has a unique solution in $C(0, T ; \overline{\mathcal{H}})$; see [5]. We recall that the conjugate transpose of $\mathcal{A}$ has domain
$\operatorname{dom}\left(\mathcal{A}^{*}\right):=\left\{\varphi \in \overline{\mathcal{H}}^{*} ;\right.$ for some $c>0:|\langle\varphi, \mathcal{A} y\rangle| \leq c\|y\|$ for all $\left.y \in \operatorname{dom}(\mathcal{A})\right\}$,
with antiduality product $\langle\cdot, \cdot\rangle:=\langle\cdot, \cdot\rangle_{\overline{\mathcal{H}}}$. Thus, $y \mapsto\langle\varphi, \mathcal{A} y\rangle$ has a unique extension to a linear continuous form over $\overline{\mathcal{H}}$, which by the definition is $\mathcal{A}^{*} \varphi$. This allows us to define weak solutions, extending to the complex setting the definition in [8].

Definition 3.1. We say that $\Psi \in C(0, T ; \overline{\mathcal{H}})$ is a weak solution of (3.1) if $\Psi(0)=$ $\Psi_{0}$ and, for any $\phi \in \operatorname{dom}\left(\mathcal{A}^{*}\right)$, the function $t \mapsto\langle\phi, \Psi(t)\rangle$ is absolutely continuous over $[0, T]$ and satisfies

$$
\begin{equation*}
\frac{\mathrm{d}}{\mathrm{~d} t}\langle\phi, \Psi(t)\rangle+\left\langle\mathcal{A}^{*} \phi, \Psi(t)\right\rangle=\left\langle\phi, f+u(t)\left(\mathcal{B}_{1}+\mathcal{B}_{2} \Psi(t)\right)\right\rangle \text { for a.a. } t \in[0, T] \tag{3.7}
\end{equation*}
$$

We recall the following result, an obvious extension to the complex setting of the corresponding result in [8].

Theorem 3.2. Let $\mathcal{A}$ be the generator of a strongly continuous semigroup. Then there is a unique weak solution of (3.7) that coincides with the mild solution.

So in the following we can use any of the two equivalent formulations (3.5) or (3.7). The control and state spaces are, respectively,

$$
\begin{equation*}
\mathcal{U}:=L^{1}(0, T), \quad \mathcal{Y}:=C(0, T ; \overline{\mathcal{H}}) \tag{3.8}
\end{equation*}
$$

Let $\hat{u} \in \mathcal{U}$ be given and $\hat{\Psi}$ a solution of (3.1). The linearized state equation at $(\hat{\Psi}, \hat{u})$, to be understood in the sense of mild solutions, is

$$
\begin{equation*}
\dot{z}(t)+\mathcal{A} z(t)=\hat{u}(t) \mathcal{B}_{2} z(t)+v(t)\left(\mathcal{B}_{1}+\mathcal{B}_{2} \hat{\Psi}(t)\right), \quad z(0)=0 \tag{3.9}
\end{equation*}
$$

where $v \in \mathcal{U}$. It is easily checked that given $v \in \mathcal{U}$, (3.9) has a unique solution denoted by $z[v]$, and that the mapping $u \mapsto \Psi[u]$ from $\mathcal{U}$ to $\mathcal{Y}$ is of class $C^{\infty}$ with $D \Psi[u] v=z[v]$.

The results above may allow us to prove higher regularity.
Definition 3.3 (restriction property). Let $E$ be a Banach space with norm denoted by $\|\cdot\|_{E}$ with continuous inclusion in $\overline{\mathcal{H}}$. Assume that the restriction of $e^{-t \mathcal{A}}$ to $E$ has an image in $E$, and that it is a continuous semigroup over this space. We let $\mathcal{A}^{\prime}$ denote its associated generator, and $e^{-t \mathcal{A}^{\prime}}$ the associated semigroup. By (3.3) we have that

$$
\begin{equation*}
\operatorname{dom}\left(\mathcal{A}^{\prime}\right):=\left\{y \in E ; \lim _{t \downarrow 0} \frac{e^{-t \mathcal{A}} y-y}{t} \text { exists }\right\} \tag{3.10}
\end{equation*}
$$

so that $\operatorname{dom}\left(\mathcal{A}^{\prime}\right) \subset \operatorname{dom}(\mathcal{A})$, and $\mathcal{A}^{\prime}$ is the restriction of $\mathcal{A}$ to $\operatorname{dom}\left(\mathcal{A}^{\prime}\right)$. We have that

$$
\begin{equation*}
\left\|e^{-t \mathcal{A}^{\prime}}\right\|_{\mathcal{L}(E)} \leq c_{\mathcal{A}^{\prime}} e^{\lambda_{\mathcal{A}^{\prime}} t} \tag{3.11}
\end{equation*}
$$

for some constants $c_{\mathcal{A}^{\prime}}$ and $\lambda_{\mathcal{A}^{\prime}}$. Assume that $\mathcal{B}_{1} \in E$, and denote by $\mathcal{B}_{2}^{\prime}$ the restriction of $\mathcal{B}_{2}$ to $E$, which is supposed to have an image in $E$ and to be continuous in the topology of $E$, that is,

$$
\begin{equation*}
\mathcal{B}_{1} \in E, \quad \mathcal{B}_{2}^{\prime} \in \mathcal{L}(E) \tag{3.12}
\end{equation*}
$$

In this case we say that $E$ has the restriction property.
3.1. Dual semigroup. Since $\overline{\mathcal{H}}$ is a reflexive Banach space it is known, e.g., [41, Chap. 1, Cor. 10.6], that $\mathcal{A}^{*}$ generates another strongly continuous semigroup called the dual semigroup on $\overline{\mathcal{H}}^{*}$, denoted by $e^{-t \mathcal{A}^{*}}$, which satisfies

$$
\begin{equation*}
\left(e^{-t \mathcal{A}}\right)^{*}=e^{-t \mathcal{A}^{*}} \tag{3.13}
\end{equation*}
$$

The reference [41] above assumes a real setting, but the arguments have an immediate extension to the complex one. Let $(z, p)$ be solution of the forward-backward system

$$
\left\{\begin{array}{rcl}
\text { (i) } & \dot{z}+\mathcal{A} z & =a z+b  \tag{3.14}\\
\text { (ii) } & -\dot{p}+\mathcal{A}^{*} p & =a^{*} p+g
\end{array}\right.
$$

where

$$
\left\{\begin{array}{l}
b \in L^{1}(0, T ; \overline{\mathcal{H}})  \tag{3.15}\\
g \in L^{1}\left(0, T ; \overline{\mathcal{H}}^{*}\right), \\
a \in L^{\infty}(0, T ; \mathcal{L}(\overline{\mathcal{H}}))
\end{array}\right.
$$

and for a.a. $t \in(0, T), a^{*}(t)$ is the conjugate transpose operator of $a(t)$, an element of $L^{\infty}\left(0, T ; \mathcal{L}\left(\overline{\mathcal{H}}^{*}\right)\right)$.

The mild solutions of (3.14), parameterized by $z(0)$ and $p(T)$, are $z \in C(0, T ; \overline{\mathcal{H}})$, $p \in C\left(0, T ; \overline{\mathcal{H}}^{*}\right)$, satisfying for a.a. $t \in(0, T)$,

$$
\begin{cases}\text { (i) } & z(t)=e^{-t \mathcal{A}} z(0)+\int_{0}^{t} e^{-(t-s) \mathcal{A}}(a(s) z(s)+b(s)) \mathrm{d} s  \tag{3.16}\\ \text { (ii) } & p(t)=e^{-(T-t) \mathcal{A}^{*}} p(T)+\int_{t}^{T} e^{-(s-t) \mathcal{A}^{*}}\left(a^{*}(s) p(s)+g(s)\right) \mathrm{d} s\end{cases}
$$

The following integration by parts lemma follows.
Lemma 3.4. Let $(z, p) \in C(0, T ; \overline{\mathcal{H}}) \times C\left(0, T ; \overline{\mathcal{H}}^{*}\right)$ satisfy (3.14)-(3.15). Then,

$$
\begin{equation*}
\langle p(T), z(T)\rangle+\int_{0}^{T}\langle g(t), z(t)\rangle \mathrm{d} t=\langle p(0), z(0)\rangle+\int_{0}^{T}\langle p(t), b(t)\rangle \mathrm{d} t \tag{3.17}
\end{equation*}
$$

Proof. This is an obvious extension of [5, Lemma 2] to the complex setting.
4. First order optimality conditions of optimal control problem. Let $q$ and $q_{T}$ be continuous quadratic forms over $\overline{\mathcal{H}}$, with associated symmetric and continuous operators $Q$ and $Q_{T}$ in $\mathcal{L}\left(\overline{\mathcal{H}}, \overline{\mathcal{H}}^{*}\right)$, such that $q(y)=\Re\langle Q y, y\rangle$ and $q_{T}(y)=$ $\Re\left\langle Q_{T} y, y\right\rangle$, where the operators $Q$ and $Q_{T}$ are self-adjoint, i.e.,

$$
\begin{equation*}
\langle Q x, y\rangle=\overline{\langle Q y, x\rangle} \quad \text { for all } x, y \text { in } \overline{\mathcal{H}} \tag{4.1}
\end{equation*}
$$

Observe that the derivative of $q$ at $y$ in direction $x$ is

$$
\begin{equation*}
D q(y) x=2 \Re\langle Q y, x\rangle \tag{4.2}
\end{equation*}
$$

Similar relations for $q_{T}$ hold. Given

$$
\begin{equation*}
\Psi_{d} \in L^{\infty}(0, T ; \overline{\mathcal{H}}) ; \quad \Psi_{d T} \in \overline{\mathcal{H}} \tag{4.3}
\end{equation*}
$$

we introduce the cost function, where $\alpha_{1} \in \mathbb{R}$ and $\alpha_{2} \geq 0$, assuming that $u \in L^{2}(0, T)$ if $\alpha_{2} \neq 0$,
$J(u, \Psi):=\int_{0}^{T}\left(\alpha_{1} u(t)+\frac{1}{2} \alpha_{2} u(t)^{2}\right) \mathrm{d} t+\frac{1}{2} \int_{0}^{T} q\left(\Psi(t)-\Psi_{d}(t)\right) \mathrm{d} t+\frac{1}{2} q_{T}\left(\Psi(T)-\Psi_{d T}\right)$.
The costate equation is

$$
\begin{equation*}
-\dot{p}+\mathcal{A}^{*} p=Q\left(\Psi-\Psi_{d}\right)+u \mathcal{B}_{2}^{*} p, \quad p(T)=Q_{T}\left(\Psi(T)-\Psi_{d T}\right) \tag{4.5}
\end{equation*}
$$

We denote by $p[u]$ the mild (backward) solution
$p(t)=e^{(t-T) \mathcal{A}^{*}} Q_{T}\left(\Psi(T)-\Psi_{d}(T)\right)+\int_{t}^{T} e^{(t-s) \mathcal{A}^{*}}\left(Q\left(\Psi(s)-\Psi_{d}(s)\right)+u(s) \mathcal{B}_{2}^{*} p(s)\right) \mathrm{d} s$.
The reduced cost is

$$
\begin{equation*}
F(u):=J(u, \Psi[u]) . \tag{4.7}
\end{equation*}
$$

The set of feasible controls is

$$
\begin{equation*}
\mathcal{U}_{a d}:=\left\{u \in \mathcal{U} ; u_{m} \leq u(t) \leq u_{M} \text { a.e. on }[0, T]\right\} \tag{4.8}
\end{equation*}
$$

with $u_{m}<u_{M}$ given real constants. The optimal control problem is

$$
\begin{equation*}
\operatorname{Min}_{u} F(u), \quad u \in \mathcal{U}_{a d} \tag{P}
\end{equation*}
$$

Definition 4.1. We say that $\hat{u} \in \mathcal{U}_{\text {ad }}$ is a minimum for $(\mathrm{P})$ if $F(\hat{u}) \leq F(u)$ for all $u \in \mathcal{U}_{a d}$. And $\hat{u} \in \mathcal{U}_{\text {ad }}$ is a weak minimum for (P) if there exists $\varepsilon>0$, such that $F(\hat{u})$ is a minimum of the set

$$
\left\{F(u) ; u \in \mathcal{U}_{a d},\|u-\hat{u}\|_{\infty}<\varepsilon\right\} .
$$

Given $\left(f, y_{0}\right) \in L^{1}(0, T ; \overline{\mathcal{H}}) \times \overline{\mathcal{H}}$, let $y\left[y_{0}, f\right]$ denote the mild solution of

$$
\begin{equation*}
\dot{y}(t)+\mathcal{A} y(t)=f(t), \quad t \in(0, T), \quad y(0)=y_{0} \tag{4.9}
\end{equation*}
$$

The compactness hypothesis is

$$
\left\{\begin{array}{l}
\text { for given } y_{0} \in \overline{\mathcal{V}} \subset \overline{\mathcal{H}}, \overline{\mathcal{V}} \subset \overline{\mathcal{H}} \text { subspace }  \tag{4.10}\\
\text { the mapping } f \mapsto \mathcal{B}_{2} y\left[y_{0}, f\right] \\
\text { is compact from } L^{2}(0, T ; \overline{\mathcal{V}}) \text { to } L^{2}(0, T ; \overline{\mathcal{H}})
\end{array}\right.
$$

Theorem 4.2. Let (4.10) hold. Then problem (P) has a nonempty set of minima.
Proof. The proof is similar to [5, Thm. 4].
We set

$$
\begin{equation*}
\Lambda(t):=\alpha_{1}+\alpha_{2} \hat{u}(t)+\Re\left\langle p(t), \mathcal{B}_{1}+\mathcal{B}_{2} \hat{\Psi}(t)\right\rangle \tag{4.11}
\end{equation*}
$$

ThEOREM 4.3. The mapping $u \mapsto F(u)$ is of class $C^{\infty}$ from $\mathcal{U}$ to $\mathbb{R}$ and we have that

$$
\begin{equation*}
D F(u) v=\int_{0}^{T} \Lambda(t) v(t) \mathrm{d} t \quad \text { for all } v \in \mathcal{U} \tag{4.12}
\end{equation*}
$$

Proof. That $F(u)$ and $J$ are of class $C^{\infty}$ follows from classical arguments based on the implicit function theorem, as in [5]. This also implies that, setting $\Psi:=\Psi[u]$ and $z:=z[u]$,

$$
\begin{align*}
D F(u) v= & \int_{0}^{T}\left(\alpha_{1}+\alpha_{2} u(t)\right) v(t) \mathrm{d} t+\int_{0}^{T} \Re\left\langle Q\left(\Psi(t)-\Psi_{d}(t)\right), z(t)\right\rangle \mathrm{d} t  \tag{4.13}\\
& +\Re\left\langle Q_{T}\left(\Psi(T)-\Psi_{d T}\right), z(T)\right\rangle
\end{align*}
$$

We deduce then (4.12) from Lemma 3.4.
Let, for $u \in \mathcal{U}_{a d}, I_{m}(u)$ and $I_{M}(u)$ be the associated contact sets defined, up to a zero-measure set, as

$$
\left\{\begin{array}{l}
I_{m}(u):=\left\{t \in(0, T): u(t)=u_{m}\right\}  \tag{4.14}\\
I_{M}(u):=\left\{t \in(0, T): u(t)=u_{M}\right\}
\end{array}\right.
$$

The first order optimality necessary condition is given as follows.
Proposition 4.4. Let $\hat{u}$ be a weak minimum of problem (P). Then, up to a set of measure zero, there holds

$$
\begin{equation*}
\{t ; \Lambda(t)>0\} \subset I_{m}(\hat{u}), \quad\{t ; \Lambda(t)<0\} \subset I_{M}(\hat{u}) . \tag{4.15}
\end{equation*}
$$

Proof. Proof is the same as in [5, Proposition 2].

## 5. Second order optimality conditions.

5.1. Technical results. Let $u$ belong to $\mathcal{U}$. Set $v:=u-\hat{u}, \hat{\Psi}:=\Psi[\hat{u}], \Psi:=\Psi[u]$, and $\delta \Psi:=\Psi-\hat{\Psi}$. Since $u \Psi-\hat{u} \hat{\Psi}=u \delta \Psi+v \hat{\Psi}$, we have that $\delta \Psi$ is the mild solution of

$$
\begin{equation*}
\frac{\mathrm{d}}{\mathrm{~d} t} \delta \Psi(t)+\mathcal{A} \delta \Psi(t)=\hat{u}(s) \mathcal{B}_{2} \delta \Psi(s)+v(t)\left(\mathcal{B}_{1}+\mathcal{B}_{2} \hat{\Psi}(t)+\mathcal{B}_{2} \delta \Psi(s)\right) \tag{5.1}
\end{equation*}
$$

Thus, $\eta:=\delta \Psi-z$ is a solution of

$$
\begin{equation*}
\dot{\eta}(t)+\mathcal{A} \eta(t)=\hat{u} \mathcal{B}_{2} \eta(t)+v(s) \mathcal{B}_{2} \delta \Psi(s) \tag{5.2}
\end{equation*}
$$

We get the following estimates.
Lemma 5.1. The linearized state $z$ solution of (3.9), the solution $\delta \Psi$ of (5.1), and $\eta=\delta \Psi-z$ the solution of (5.2) satisfy, whenever $v$ remains in a bounded set of $L^{1}(0, T)$,

$$
\begin{align*}
\|z\|_{L^{\infty}(0, T ; \overline{\mathcal{H}})} & =O\left(\|v\|_{1}\right)  \tag{5.3}\\
\|\delta \Psi\|_{L^{\infty}(0, T ; \overline{\mathcal{H}})} & =O\left(\|v\|_{1}\right)  \tag{5.4}\\
\|\eta\|_{L^{\infty}(0, T ; \overline{\mathcal{H}})} & =O\left(\|\delta \Psi v\|_{L^{1}(0, T ; \overline{\mathcal{H}})}\right)=O\left(\|v\|_{1}^{2}\right) \tag{5.5}
\end{align*}
$$

Proof. Proof is similar to the proof of Lemma 4 in [5].
For $(\hat{\Psi}, \hat{u})$ a solution of $(3.1), \hat{p}$ the corresponding solution of $(4.6), v \in L^{1}(0, T)$, and $z \in C(0, T ; \overline{\mathcal{H}})$, let us set

$$
\begin{equation*}
\mathcal{Q}(z, v):=\int_{0}^{T}\left(q(z(t))+\alpha_{2} v(t)^{2}+2 v(t) \Re\left\langle\hat{p}(t), \mathcal{B}_{2} z(t)\right\rangle\right) \mathrm{d} t+q_{T}(z(T)) \tag{5.6}
\end{equation*}
$$

We can refer to $\mathcal{Q}$ as the second variation of the Lagrangian.
Proposition 5.2. Let $u$ belong to $\mathcal{U}$. Set $v:=u-\hat{u}, \hat{\Psi}:=\Psi[\hat{u}], \Psi:=\Psi[u]$. Then

$$
\begin{equation*}
F(u)=F(\hat{u})+D F(\hat{u}) v+\frac{1}{2} \mathcal{Q}(\delta \Psi, v) . \tag{5.7}
\end{equation*}
$$

Proof. We can expand the cost function as follows:

$$
\begin{aligned}
F(u)= & F(\hat{u})+\frac{1}{2} \int_{0}^{T}\left(\alpha_{2} v(t)^{2}+q(\delta \Psi(t))\right) \mathrm{d} t+\frac{1}{2} q_{T}(\delta \Psi(T)) \\
& +\int_{0}^{T}\left(\alpha_{1}+\alpha_{2} \hat{u}(t)\right) v(t) \mathrm{d} t \\
& \left.+\Re\left(\int_{0}^{T}\left\langle Q\left(\hat{\Psi}(t)-\Psi_{d}(t)\right), \delta \Psi\right)\right\rangle \mathrm{d} t+\left\langle Q_{T}\left(\hat{\Psi}(T)-\Psi_{d}(T)\right), \delta \Psi(T)\right\rangle\right)
\end{aligned}
$$

Applying Lemma 3.4 to the pair $(\delta \Psi, \hat{p})$, where $z$ is a solution of the linearized equation (3.9), and using the expression of $\Lambda$ in (4.11), we obtain the result.

Corollary 5.3. We have that

$$
\begin{equation*}
F(u)=F(\hat{u})+D F(\hat{u}) v+\frac{1}{2} \mathcal{Q}(z, v)+O\left(\|v\|_{1}^{3}\right) \tag{5.9}
\end{equation*}
$$

where $z:=z[v]$.

Proof. We have that

$$
\begin{aligned}
\mathcal{Q}(\delta \Psi, v)-\mathcal{Q}(z, v)= & \Re\left(\int_{0}^{T}\langle Q(\delta \Psi(t)+z(t)), \eta(t)\rangle+2 v(t)\left\langle p(t), B_{2} \eta(t)\right\rangle \mathrm{d} t\right) \\
& +\Re\left\langle Q_{T}(\delta \Psi(T)+z(T)), \eta(T)\right\rangle .
\end{aligned}
$$

By (5.3)-(5.5) the difference above is of the order of $\|v\|_{1}^{3}$. The conclusion follows.
5.2. Second order necessary optimality conditions. Given a feasible control $u$, the critical cone is defined as

$$
C(u):=\left\{\begin{array}{l}
v \in L^{1}(0, T) \mid \Lambda(t) v(t)=0 \text { a.e. on }[0, T]  \tag{5.10}\\
v(t) \geq 0 \text { a.e. on } I_{m}(u), v(t) \leq 0 \text { a.e. on } I_{M}(u)
\end{array}\right\} .
$$

THEOREM 5.4. Let $\hat{u} \in \mathcal{U}_{\text {ad }}$ be a weak minimum of ( P ) and $\hat{p}$ be the corresponding costate. Then the second variation $\mathcal{Q}$ is positive semidefinite over the critical cone $C(\hat{u})$, i.e., there holds,

$$
\begin{equation*}
\mathcal{Q}(z[v], v) \geq 0 \quad \text { for all } v \in C(\hat{u}) \tag{5.11}
\end{equation*}
$$

Proof. The proof is similar to the one of Theorem 6 in [5].
5.3. Second order sufficient optimality conditions. In this subsection we assume that $\alpha_{2}>0$, and obtain second order sufficient optimality conditions. Consider the following condition of positive definiteness of $\mathcal{Q}$ : there exists $\alpha_{0}>0$ such that

$$
\begin{equation*}
\mathcal{Q}(z[v], v) \geq \alpha_{0} \int_{0}^{T} v(t)^{2} \mathrm{~d} t \quad \text { for all } v \in C(\hat{u}) \tag{5.12}
\end{equation*}
$$

Definition 5.5. We say that a weak minimum $\hat{u} \in \mathcal{U}_{\text {ad }}$ satisfies a quadratic growth condition if there exist $\varepsilon>0$ and $\varepsilon^{\prime}>0$ such that

$$
\begin{equation*}
F(u) \geq F(\hat{u})+\varepsilon\|u-\hat{u}\|_{2}^{2} \quad \text { for every } u \in \mathcal{U}_{\text {ad }} \text { with }\|u-\hat{u}\|_{\infty}<\varepsilon^{\prime} \tag{5.13}
\end{equation*}
$$

ThEOREM 5.6. Let $\hat{u} \in \mathcal{U}_{\text {ad }}$ satisfy the first order optimality conditions (4.15) of $(\mathrm{P}), \hat{p}$ being the corresponding costate, as well as the positive definiteness condition (5.12). Then $\hat{u}$ is a weak minimum of problem $(P)$ that satisfies the quadratic growth condition (5.13).

Proof. It suffices to adapt the arguments in, say, [18, Thm. 4.3] or Casas and Tröltzsch [25].

Using the technique of Bonnans and Osmolovskiĭ [19] we can actually deduce from Theorem 5.4 that $\hat{u}$ is a strong solution in the following sense (natural extension of the notion of strong solution in the sense of the calculus of variations).

DEFINITION 5.7. We say that a control $\hat{u} \in \mathcal{U}_{a d}$ is a strong minimum if there exists $\varepsilon>0$ such that, if $u \in \mathcal{U}_{a d}$ and $\|y[u]-y[\hat{u}]\|_{C(0, T ; \overline{\mathcal{H}})}<\varepsilon$, then $F(\hat{u}) \leq F(u)$.

In the context of optimal control of PDEs, sufficient conditions for strong optimality were recently obtained for elliptic state equations in Bayen, Bonnans, and Silva [11], and for parabolic equations by Bayen and Silva [12], and by Casas and Tröltzsch [25].

We consider the part of the Hamiltonian depending on the control:

$$
\begin{equation*}
H(t, u):=\alpha_{1} u+\frac{1}{2} \alpha_{2} u^{2}+u \Re\langle\hat{p}(t), \mathcal{B}(t)\rangle \tag{5.14}
\end{equation*}
$$

where $\mathcal{B}(t):=\mathcal{B}(t)_{1}+\mathcal{B}(t)_{2} \hat{\Psi}(t)$. The Hamiltonian inequality reads

$$
\begin{equation*}
H(t, \hat{u}(t)) \leq H(t, u) \quad \text { for all } u \in\left[u_{m}, u_{M}\right] \text { for a.a. } t \in[0, T] \tag{5.15}
\end{equation*}
$$

Since $\alpha_{2}>0, H(t, \cdot)$ is a strongly convex function, and therefore the Hamiltonian inequality follows from the first order optimality conditions and in addition we have the quadratic growth property

$$
\begin{equation*}
H(t, \hat{u}(t))+\frac{1}{2} \alpha_{2}(u-\hat{u}(t))^{2} \leq H(t, u) \quad \text { for all } u \in\left[u_{m}, u_{M}\right], \text { for a.a. } t \in[0, T] \tag{5.16}
\end{equation*}
$$

Lemma 5.8. Assume that $\alpha_{2}>0$. Let $\hat{u}$ be feasible and satisfy the first order optimality conditions (4.15). Let $\left(u_{k}\right)$ be a sequence of feasible controls such that the associated states $\hat{\Psi}_{k}:=\Psi\left[u_{k}\right]$ converge to $\hat{\Psi}$ in $C(0, T ; \overline{\mathcal{H}})$, and $\limsup _{k} F\left(u_{k}\right) \leq$ $F(\hat{u})$. Then $u_{k} \rightarrow \hat{u}$ in $L^{2}(0, T)$.

Proof. Since $\left(u_{k}\right)$ is bounded in $L^{\infty}(0, T)$, from the expression of the cost function of the optimal control problem in view of Theorem 4.3 and Corollary 5.3, it follows that

$$
\begin{equation*}
0 \geq \limsup _{k}\left(F\left(u_{k}\right)-F(\hat{u})\right)=\underset{k}{\lim \sup _{k}} \int_{0}^{T}\left(H\left(t, u_{k}(t)\right)-H(t, \hat{u}(t))\right) \mathrm{d} t \tag{5.17}
\end{equation*}
$$

Then the conclusion follows from the quadratic growth property (5.16).
For $u_{k}$ as in Lemma 5.8 we have

$$
\begin{equation*}
B_{k}:=\left\{t \in(0, T) ;\left|u_{k}(t)-\hat{u}(t)\right|>\sqrt{\left\|u_{k}-\hat{u}\right\|_{1}}\right\} ; \quad A_{k}:=(0, T) \backslash B_{k} \tag{5.18}
\end{equation*}
$$

Note that

$$
\begin{equation*}
\left|B_{k}\right| \leq \int_{0}^{T} \frac{\left|u_{k}(t)-\hat{u}(t)\right|}{\sqrt{\left\|u_{k}-\hat{u}\right\|_{1}}} \mathrm{~d} t=\sqrt{\left\|u_{k}-\hat{u}\right\|_{1}} \tag{5.19}
\end{equation*}
$$

Set, for a.a. $t$,

$$
\begin{equation*}
v_{k}^{A}(t):=\left(u_{k}(t)-\hat{u}(t)\right) \mathbf{1}_{A_{k}}(t), \quad v_{k}^{B}(t):=\left(u_{k}(t)-\hat{u}(t)\right) \mathbf{1}_{B_{k}}(t) \tag{5.20}
\end{equation*}
$$

We now extend to the semigroup setting the decomposition principle from [19], which has been extended to the elliptic setting by [11], and to the parabolic setting by [12].

Theorem 5.9 (decomposition principle). For $u_{k}$ as in Lemma 5.8 we have that $\left|B_{k}\right| \rightarrow 0$ and

$$
\begin{equation*}
F\left(u_{k}\right)=F\left(\hat{u}+v_{k}^{A}\right)+F\left(\hat{u}+v_{k}^{B}\right)-F(\hat{u})+o\left(\left\|u_{k}-\bar{u}\right\|_{2}^{2}\right), \tag{5.21}
\end{equation*}
$$

and also

$$
\begin{equation*}
F\left(\hat{u}+v_{k}^{B}\right)-F(\hat{u})=\int_{B^{k}}\left(H\left(t, u_{k}(t)\right)-H(t, \hat{u}(t))\right) \mathrm{d} t+o\left(\left\|u_{k}-\bar{u}\right\|_{2}^{2}\right) \tag{5.22}
\end{equation*}
$$

Proof. Remember the linearized state equation (3.9) whose solution is denoted by $z[v]$. Set

$$
\begin{equation*}
v_{k}:=u_{k}-\hat{u} ; \quad z_{k}:=z\left[v_{k}\right], \quad z_{k}^{A}:=z\left[v_{k}^{A}\right] ; \quad z_{k}^{B}:=z\left[v_{k}^{B}\right] . \tag{5.23}
\end{equation*}
$$

Since $A_{k} \cap B_{k}$ has null measure, we have that $z_{k}=z_{k}^{A}+z_{k}^{B}$. Also,

$$
\begin{equation*}
\left\|v_{k}^{B}\right\|_{1} \leq\left|B_{k}\right|^{1 / 2}\left\|v_{k}^{B}\right\|_{2}=o\left(\left\|v_{k}^{B}\right\|_{2}\right) \tag{5.24}
\end{equation*}
$$

since $\left|B_{k}\right| \rightarrow 0$ by Lemma 5.8. Then, in view of Lemma 5.1,

$$
\begin{equation*}
\left\|z_{k}^{B}\right\|_{C(0, T ; \overline{\mathcal{H}})}=O\left(\left\|v_{k}^{B}\right\|_{1}\right)=o\left(\left\|v_{k}^{B}\right\|_{2}\right) \tag{5.25}
\end{equation*}
$$

Combining with Corollary 5.3 and using the fact that $v_{k}^{A}(t) v_{k}^{B}(t)=0$ a.e., we deduce that

$$
\begin{align*}
F\left(u_{k}\right)-F(\hat{u})= & D F(\hat{u}) v_{k}+\frac{1}{2} \mathcal{Q}\left(v_{k}, z_{k}\right)+o\left(\left\|v_{k}\right\|_{2}^{2}\right)  \tag{5.26}\\
= & D F(\hat{u}) v_{k}+\frac{1}{2} \mathcal{Q}\left(v_{k}, z_{k}^{A}\right)+o\left(\left\|v_{k}\right\|_{2}^{2}\right) \\
= & D F(\hat{u}) v_{k}^{A}+\frac{1}{2} \mathcal{Q}\left(v_{k}^{A}, z_{k}^{A}\right)+D F(\hat{u}) v_{k}^{B}+\frac{1}{2} \alpha\left\|v_{k}^{B}\right\|_{2}^{2} \\
& +2 \int_{0}^{T} v_{k}^{B}(t) \Re\left\langle\hat{p}(t), \mathcal{B}_{2} z_{k}^{A}(t)\right\rangle \mathrm{d} t+o\left(\left\|v_{k}\right\|_{2}^{2}\right) \\
= & D F(\hat{u}) v_{k}^{A}+\frac{1}{2} \mathcal{Q}\left(v_{k}^{A}, z_{k}^{A}\right)+D F(\hat{u}) v_{k}^{B}+\frac{1}{2} \alpha\left\|v_{k}^{B}\right\|_{2}^{2}+o\left(\left\|v_{k}\right\|_{2}^{2}\right),
\end{align*}
$$

where we have used the fact that, by (5.24),

$$
\begin{equation*}
\left|\int_{0}^{T} v_{k}^{B}(t) \Re\left\langle\hat{p}(t), \mathcal{B}_{2} z_{k}^{A}(t)\right\rangle \mathrm{d} t\right|=O\left(\left\|v_{k}^{B}\right\|_{1}\left\|z_{k}^{A}\right\|_{C(0, T ; \overline{\mathcal{H}})}\right)=o\left(\left\|v_{k}\right\|_{2}^{2}\right) \tag{5.27}
\end{equation*}
$$

Now

$$
\begin{equation*}
F\left(\hat{u}+v_{k}^{A}\right)-F(\hat{u})=D F(\hat{u}) v_{k}^{A}+\frac{1}{2} \mathcal{Q}\left(v_{k}^{A}, z_{k}^{A}\right)+o\left(\left\|v_{k}^{A}\right\|_{2}^{2}\right) \tag{5.28}
\end{equation*}
$$

and by (5.25)

$$
\begin{equation*}
F\left(\hat{u}+v_{k}^{B}\right)-F(\hat{u})=D F(\hat{u}) v_{k}^{B}+\frac{1}{2} \alpha_{2}\left\|v_{k}^{B}\right\|_{2}^{2}+o\left(\left\|v_{k}^{B}\right\|_{2}^{2}\right) \tag{5.29}
\end{equation*}
$$

Combining the above relations we get the desired result.
Definition 5.10. We say that $\hat{u} \in \mathcal{U}_{\text {ad }}$ satisfies the quadratic growth condition for strong solutions if there exist $\varepsilon>0$ and $\varepsilon^{\prime}>0$ such that for any feasible control $u$,

$$
\begin{equation*}
F(\hat{u})+\varepsilon\|u-\hat{u}\|_{2}^{2} \leq F(u) \quad \text { whenever }\|\Psi[u]-\Psi[\hat{u}]\|_{C(0, T ; \overline{\mathcal{H}})}<\varepsilon^{\prime} \tag{5.30}
\end{equation*}
$$

Theorem 5.11. Let $\hat{u} \in \mathcal{U}_{\text {ad }}$ satisfy the first order necessary optimality condition (4.15), and the condition of positive definiteness of the second variation (5.12). Then $\hat{u}$ is a strong minimum that satisfies the quadratic growth for strong solutions.

Proof. If the conclusion is false, then there exists a sequence $\left(u_{k}\right)$ of feasible controls such that $\Psi_{k} \rightarrow \hat{\Psi}$ in $C(0, T ; \overline{\mathcal{H}})$, where $\Psi_{k}:=\Psi\left[u_{k}\right]$, and $F\left(u_{k}\right) \leq F(\hat{u})+$ $o\left(\left\|u_{k}-\hat{u}\right\|_{2}^{2}\right)$. By Lemma 5.8, $u_{k} \rightarrow \hat{u}$ in $L^{2}(0, T)$. By the decomposition Theorem 5.9 and since $D F(\hat{u}) v_{k}^{B} \geq 0$, it follows that

$$
\begin{equation*}
\alpha_{2}\left\|v_{k}^{B}\right\|_{2}^{2}+F\left(\hat{u}+v_{k}^{A}\right)-F(\hat{u}) \leq o\left(\left\|v_{k}\right\|_{2}^{2}\right) \tag{5.31}
\end{equation*}
$$

We next distinguish two cases.
(a) Assume that $\left\|v_{k}^{A}\right\|_{2} /\left\|v_{k}\right\|_{2} \rightarrow 0$. We know that

$$
\begin{equation*}
F\left(\hat{u}+v_{k}^{A}\right)-F(\hat{u})=D F(\hat{u}) v_{k}^{A}+\frac{1}{2} \mathcal{Q}\left(v_{k}^{A}, z_{k}^{A}\right)+o\left(\left\|v_{k}^{A}\right\|_{2}^{2}\right) \tag{5.32}
\end{equation*}
$$

Since (by the first order optimality conditions) $D F(\hat{u}) v_{k}^{A} \geq 0$ and $\mathcal{Q}\left(v_{k}^{A}, z_{k}^{A}\right)=$ $O\left(\left\|v_{k}^{A}\right\|_{2}^{2}\right)=o\left(\left\|v_{k}\right\|_{2}^{2}\right)$ by hypothesis, it follows with (5.31) that $\left\|v_{k}^{B}\right\|_{2}^{2}=o\left(\left\|v_{k}\right\|_{2}^{2}\right)=$ $o\left(\left\|v_{k}^{B}\right\|_{2}^{2}\right)$ which gives a contradiction.
(b) Otherwise, $\liminf _{k}\left\|v_{k}^{A}\right\|_{2} /\left\|v_{k}\right\|_{2}>0$ (extracting if necessary a subsequence). It follows from (5.31) that

$$
\begin{equation*}
F\left(\hat{u}+v_{k}^{A}\right)-F(\hat{u}) \leq o\left(\left\|v_{k}^{A}\right\|_{2}\right) \tag{5.33}
\end{equation*}
$$

Since $\left\|v_{k}^{A}\right\|_{\infty} \rightarrow 0$, we obtain a contradiction with Theorem 5.4.
Remark 5.12. A shorter proof for Theorem 5.9 is obtained by combining Lemma 5.8 and the Taylor expansion in Corollary 5.3, which implies

$$
\begin{equation*}
F(u)=F(\hat{u})+D F(\hat{u}) v+\frac{1}{2} \mathcal{Q}(z, v)+O\left(\|v\|_{2}^{3}\right) \tag{5.34}
\end{equation*}
$$

from which we can state a sufficient condition for optimality in $L^{2}(0, T)$. On the other hand the present proof opens the way for dealing with nonquadratic (w.r.t. the control) Hamiltonian functions, as in [11].
6. Second order optimality conditions for singular problems. In this section we assume that $\alpha_{2}=0$, so that the control enters linearly in both the state equation and cost function. For such optimal control problems there is an extensive theory in the finite dimensional setting; see Kelley [37], Goh [33], Dmitruk [26, 27], Poggiolini and Stefani [42], Aronna et al. [2], and Frankowska and Tonon [31]; the case of additional scalar state constraints was considered in Aronna, Bonnans, and Goh [1].

In the context of optimal control of PDEs, there exist very few papers on sufficient optimality conditions for control-affine control problems; see Bergounioux and Tiba [16], Tröltzsch [44], Bonnans and Tiba [21], Casas [23] (and the related literature involving $L^{1}$ norms; see, e.g., Casas, Clason, and Kunisch [24]). As mentioned in the introduction, here we will follow the ideas in $[5,17]$ by using in an essential way the Goh transform [33].

Let $E_{1} \subset \mathcal{H}$ with continuous inclusion, having the restriction property (Definition 3.3). We can denote the restriction of $\mathcal{B}_{2}$ to $E_{1}$ by $\mathcal{B}_{2}$ with no risk of confusion. In the rest of the paper we make the following hypothesis:

$$
\begin{cases}\text { (i) } & \mathcal{B}_{1} \in \operatorname{dom}(\mathcal{A})  \tag{6.1}\\ \text { (ii) } & \mathcal{B}_{2} \operatorname{dom}(\mathcal{A}) \subset \operatorname{dom}(\mathcal{A}), \\ \mathcal{B}_{2}^{*} \operatorname{dom}\left(\mathcal{A}^{*}\right) \subset \operatorname{dom}\left(\mathcal{A}^{*}\right)\end{cases}
$$

with $\mathcal{B}_{i}^{k}:=\left(\mathcal{B}_{i}\right)^{k}$. So, we may define the operators below, with $\operatorname{domains} \operatorname{dom}(\mathcal{A})$ and $\operatorname{dom}\left(\mathcal{A}^{*}\right)$, respectively, for $k=1,2$ :

$$
\left\{\begin{align*}
{\left[\mathcal{A}, \mathcal{B}_{2}^{k}\right]: } & : \mathcal{A} \mathcal{B}_{2}^{k}-\mathcal{B}_{2}^{k} \mathcal{A},  \tag{6.2}\\
{\left[\left(\mathcal{B}_{2}^{k}\right)^{*}, \mathcal{A}^{*}\right]: } & :\left(\mathcal{B}_{2}^{k}\right)^{*} \mathcal{A}^{*}-\mathcal{A}^{*}\left(\mathcal{B}_{2}^{k}\right)^{*}
\end{align*}\right.
$$

We also suppose in the following that

$$
\begin{cases}\text { (i) } & \text { for } k=1,2,\left[\mathcal{A}, \mathcal{B}_{2}^{k}\right] \text { has a continuous extension to } E_{1}  \tag{6.3}\\ & \text { denoted by } M_{k} ; \\ \text { (ii) } & f \in L^{\infty}(0, T ; \overline{\mathcal{H}}), \quad M_{k}^{*} \hat{p} \in L^{\infty}\left(0, T ; \overline{\mathcal{H}}^{*}\right), k=1,2 \\ \text { (iii) } \hat{\Psi} \in L^{2}\left(0, T ; E_{1}\right), \quad\left[M_{1}, \mathcal{B}_{2}\right] \hat{\Psi} \in L^{\infty}(0, T ; \overline{\mathcal{H}})\end{cases}
$$

Remark 6.1. Point (6.1)(ii) implies

$$
\begin{equation*}
\mathcal{B}_{2}^{k} \operatorname{dom}(\mathcal{A}) \subset \operatorname{dom}(\mathcal{A}), \quad\left(\mathcal{B}_{2}^{k}\right)^{*} \operatorname{dom}\left(\mathcal{A}^{*}\right) \subset \operatorname{dom}\left(\mathcal{A}^{*}\right) \quad \text { for } k=1,2 . \tag{6.4}
\end{equation*}
$$

So, $\left[\mathcal{A}, \mathcal{B}_{2}\right]$ is well-defined as an operator with domain $\operatorname{dom}(\mathcal{A})$, and then point (6.3)(iii) makes sense.

We also assume that

$$
\begin{cases}\text { (i) } & \mathcal{B}_{2}^{2} f \in C(0, T ; \overline{\mathcal{H}}), \quad \Psi_{d} \in C(0, T ; \overline{\mathcal{H}}),  \tag{6.5}\\ \text { (ii) } & M_{k}^{*} \hat{p} \in C\left(0, T ; \overline{\mathcal{H}}^{*}\right), \quad k=1,2 .\end{cases}
$$

Given $w \in L^{2}(0, T)$, let $\xi \in C(0, T ; \overline{\mathcal{H}})$ be a (mild) solution of

$$
\begin{equation*}
\dot{\xi}+\mathcal{A} \xi=\hat{u} \mathcal{B}_{2} \xi+w b_{z}^{1}, \quad \xi(0)=0, \tag{6.6}
\end{equation*}
$$

where

$$
\begin{equation*}
b_{z}^{1}:=-\mathcal{B}_{2} f-M_{1} \hat{\Psi}-\mathcal{A} \mathcal{B}_{1} . \tag{6.7}
\end{equation*}
$$

By (6.1) and (6.3), $b_{z}^{1} \in L^{2}(0, T ; \overline{\mathcal{H}})$, so that $w b_{z}^{1} \in L^{1}(0, T ; \overline{\mathcal{H}})$, and (6.6) has a unique solution. Consider the space

$$
\begin{equation*}
W:=\left(L^{2}\left(0, T ; E_{1}\right) \cap C([0, T] ; \mathcal{H})\right) \times L^{2}(0, T) \times \mathbb{R} . \tag{6.8}
\end{equation*}
$$

We define the continuous quadratic forms over $W$ by

$$
\begin{equation*}
\widehat{\mathcal{Q}}(\xi, w, h)=\widehat{\mathcal{Q}}_{T}(\xi, h)+\widehat{\mathcal{Q}}_{a}(\xi, w)+\widehat{\mathcal{Q}}_{b}(w), \tag{6.9}
\end{equation*}
$$

where $\widehat{\mathcal{Q}}_{b}(w):=\int_{0}^{T} w^{2}(t) R(t) \mathrm{d} t$ and

$$
\begin{equation*}
\widehat{\mathcal{Q}}_{T}(\xi, h):=q_{T}(\xi(T)+h \mathcal{B}(T))+h^{2} \Re\left\langle\hat{p}(T), \mathcal{B}_{2} \mathcal{B}_{1}+\mathcal{B}_{2}^{2} \hat{\Psi}(T)\right\rangle+h \Re\left\langle\hat{p}(T), \mathcal{B}_{2} \xi(T)\right\rangle, \tag{6.10}
\end{equation*}
$$

$$
\begin{equation*}
\widehat{\mathcal{Q}}_{a}(\xi, w):=\Re \int_{0}^{T}\left(q(\xi)+2 w\langle Q \xi, \mathcal{B}\rangle+2 w\left\langle Q\left(\hat{\Psi}-\Psi_{d}\right), \mathcal{B}_{2} \xi\right\rangle-2 w\left\langle M_{1}^{*} \hat{p}, \xi\right\rangle\right) \mathrm{d} t \tag{6.11}
\end{equation*}
$$

with $R \in L^{\infty}(0, T)$ given by

$$
\left\{\begin{array}{c}
R(t):=q(\mathcal{B})+\Re\left\langle Q\left(\hat{\Psi}-\Psi_{d}\right), \mathcal{B}_{2} \mathcal{B}\right\rangle+\Re\langle\hat{p}(t), r(t)\rangle,  \tag{6.12}\\
r(t):=\mathcal{B}_{2}^{2} f(t)-\mathcal{A B}_{2} \mathcal{B}_{1}+2 \mathcal{B}_{2} \mathcal{A B}_{1}-\left[M_{1}, \mathcal{B}_{2}\right] \hat{\Psi} .
\end{array}\right.
$$

We write $P C_{2}(\hat{u})$ for the closure in the $L^{2} \times \mathbb{R}$-topology of the set

$$
\begin{equation*}
P C(\hat{u}):=\left\{(w, h) \in W^{1, \infty}(0, T) \times \mathbb{R} ; \dot{w} \in C(\hat{u}) ; w(0)=0, w(T)=h\right\} . \tag{6.13}
\end{equation*}
$$

The final value of $w$ becomes an independent variable when we consider this closure.
Definition 6.2 (singular arc). The control $\hat{u} \in \mathcal{U}_{a d}$ is said to have $a$ singular arc in a nonempty interval $\left(t_{1}, t_{2}\right) \subset[0, T]$ if, for all $\theta>0$, there exists $\varepsilon>0$ such that

$$
\begin{equation*}
\hat{u}(t) \in\left[u_{m}+\varepsilon, u_{M}-\varepsilon\right], \quad \text { for a.a. } t \in\left(t_{1}+\theta, t_{2}-\theta\right) . \tag{6.14}
\end{equation*}
$$

We may also say that $\left(t_{1}, t_{2}\right)$ is a singular arc itself. We call $\left(t_{1}, t_{2}\right)$ a lower boundary arc if $\hat{u}(t)=u_{m}$ for a.a. $t \in\left(t_{1}, t_{2}\right)$, and an upper boundary arc if $\hat{u}(t)=u_{M}$ for a.a. $t \in\left(t_{1}, t_{2}\right)$. We sometimes simply call them boundary arcs. We say that a boundary $\operatorname{arc}(c, d)$ is initial if $c=0$, and final if $d=T$.

Lemma 6.3. For $v \in L^{1}(0, T)$ and $w \in A C(0, T), w(t)=\int_{0}^{t} v(s) \mathrm{d} s$, there holds

$$
\begin{equation*}
\mathcal{Q}(z[v], v)=\widehat{\mathcal{Q}}(\xi[w], w, w(T)) \tag{6.15}
\end{equation*}
$$

We refer to $\widehat{\mathcal{Q}}$ as a transformed second variation.
Theorem 6.4 (second order necessary condition). Let $\hat{u} \in \mathcal{U}_{\text {ad }}$ be a weak minimum. Then, the transformed second variation $\widehat{\mathcal{Q}}$ is positive semidefinite on $P C_{2}(\hat{u})$, that is,

$$
\begin{equation*}
\widehat{\mathcal{Q}}(\xi[w], w, h) \geq 0 \quad \text { for all }(w, h) \in P C_{2}(\hat{u}) \tag{6.16}
\end{equation*}
$$

In addition, provided the mapping

$$
\begin{equation*}
w \mapsto \xi[w], \quad L^{2}(0, T) \rightarrow L^{2}(0, T ; \overline{\mathcal{H}}) \tag{6.17}
\end{equation*}
$$

is compact, we have that

$$
\begin{equation*}
R(t) \geq 0 \text { over singular arcs } \tag{6.18}
\end{equation*}
$$

Proof. Proof is similar to [5, Lemma 6 and Corollary 5].
In the following we assume that the following hypotheses hold:

1. finite structure:

$$
\left\{\begin{array}{l}
\text { there are finitely many boundary and singular maximal arcs }  \tag{6.19}\\
\text { and the closure of their union is }[0, T],
\end{array}\right.
$$

2. strict complementarity: for the control constraint (note that $\Lambda$ is a continuous function of time) (6.20)
$\left\{\begin{array}{l}\Lambda \text { has nonzero values over the interior of each boundary arc, and } \\ \text { at time } 0 \text { (resp., } T \text { ) if an initial (resp., final) boundary arc exists, }\end{array}\right.$
set

$$
\widehat{P C_{2}}(\hat{u}):=\left\{\begin{array}{l}
(w, h) \in L^{2}(0, T) \times \mathbb{R}, w \text { is constant over boundary arcs, }  \tag{6.21}\\
w=0 \text { over an initial boundary arc } \\
\text { and } w=h \text { over a terminal boundary arc }
\end{array}\right\}
$$

Recall the definition of $P C_{2}(\hat{u})$ given just before (6.13).
Proposition 6.5. Let (6.19)-(6.20) hold. Then

$$
\begin{equation*}
P C_{2}(\hat{u})=\left\{(w, h) \in \widehat{P C_{2}}(\hat{u}) ; w \text { is continuous at bang-bang junctions }\right\} \tag{6.22}
\end{equation*}
$$

Proof. The proof is a simplified version of the one of Proposition 4 in [1]. That result dealt with problems with both upper and lower bounds on the control, as well as state constraints, the latter being absent in the present setting.

Letting $\mathcal{T}_{B B}$ denote the set of bang-bang junctions, we assume in addition that

$$
\begin{equation*}
R(t)>0, \quad t \in \mathcal{T}_{B B} \tag{6.23}
\end{equation*}
$$

Consider the following uniform positive definiteness conditions on the transformed second variation: there exists $\alpha>0$ such that

$$
\begin{array}{ll}
\widehat{\mathcal{Q}}(\xi[w], w, h) \geq \alpha\left(\|w\|_{2}^{2}+h^{2}\right) & \text { for all }(w, h) \in P C_{2}(\hat{u}) \\
\widehat{\mathcal{Q}}(\xi[w], w, h) \geq \alpha\left(\|w\|_{2}^{2}+h^{2}\right) & \text { for all }(w, h) \in \widehat{P C_{2}}(\hat{u}) \tag{6.25}
\end{array}
$$

Since $P C_{2}(\hat{u}) \subset \widehat{P C_{2}}(\hat{u}),(6.25)$ implies (6.24).

Definition 6.6. We say that $\hat{u} \in \mathcal{U}_{\text {ad }}$ satisfies the weak quadratic growth condition if there exist $\varepsilon>0$ and $\varepsilon^{\prime}>0$ such that for any $u \in \mathcal{U}_{\text {ad }}$, setting $v:=u-\hat{u}$ and $w(t):=\int_{0}^{T} v(s) \mathrm{d} s$, we have

$$
\begin{equation*}
F(u) \geq F(\hat{u})+\varepsilon\left(\|w\|_{2}^{2}+w(T)^{2}\right) \quad \text { if }\|v\|_{1}<\varepsilon^{\prime} \tag{6.26}
\end{equation*}
$$

The word "weak" makes reference to the fact that the growth is obtained for the $L^{2}$ norm of $w$, and not the one of $v$.

Theorem 6.7. Assume that (6.1), (6.3), and (6.5) hold, as well as (6.19)-(6.20) and (6.23).
(i) Let $\hat{u} \in \mathcal{U}_{\text {ad }}$ satisfy the first order necessary optimality conditions (4.15). Then, the uniform positive definiteness on $\widehat{P C}_{2}(\hat{u})$ in (6.25) implies the weak quadratic growth (6.26).
(ii) Conversely, for a weak minimum $\hat{u} \in \mathcal{U}_{\text {ad }}$, the weak quadratic growth condition (6.26) implies the uniform positive definiteness on $P C_{2}(\hat{u})$ in (6.24).

Proof. Proof is similar to the one in [5, Thm. 8], taking into account the erratum [4].

Remark 6.8. Under the assumptions of the previous theorem, if no bang-bang switch occurs, $P C_{2}(\hat{u})=\widehat{P C_{2}}(\hat{u}),(6.24)$ is equivalent to the quadratic growth condition (6.26), and the necessary and sufficient conditions have no gap.

## 7. Application to the Schrödinger equation.

7.1. Statement of the problem. The equation is formulated first in an informal way. Let $\Omega \subset \mathbb{R}^{n}, n \in \mathbb{N}$, open and bounded, and $T>0$. The state equation, with $\Psi=\Psi(t, x)$, is

$$
\left\{\begin{align*}
\dot{\Psi}(t, x)-i \sum_{j, k=1}^{n} \frac{\partial}{\partial x_{k}}\left[a_{j k}(x) \frac{\partial \Psi(t, x)}{\partial x_{j}}\right] & =-i u b_{2} \Psi(t, x)+f & & \text { in }(0, T) \times \Omega  \tag{7.1}\\
\Psi(0, x) & =\Psi_{0} & & \text { in } \Omega, \\
\Psi(t, x) & =0 & & \text { on }(0, T) \times \partial \Omega
\end{align*}\right.
$$

with

$$
\begin{equation*}
\Psi_{0} \in \bar{V}, \quad b_{2}^{k} \in W^{2, \infty}(\Omega), k=1,2, \quad f \in L^{2}(0, T ; \bar{V}) \cap C(0, T ; \bar{H}) \tag{7.2}
\end{equation*}
$$

and the complex valued spaces $\bar{H}:=L^{2}(\Omega ; \mathbb{C})$ and $\bar{V}:=H_{0}^{1}(\Omega ; \mathbb{C})$. Note that although $f$ is usually equal to zero, it is useful to introduce it, since the sensitivity of the solution w.r.t. the right-hand side plays a role in the numerical analysis. Here the $a_{j k}$ are $C^{1}$ functions over $\bar{\Omega}$ that satisfy, for each $x \in \bar{\Omega}$, the symmetry hypothesis $a_{j k}=a_{k j}$ for all $j, k$ as well as the following coercivity hypothesis, that for some $\nu>0$,

$$
\begin{equation*}
\sum_{j, k=1}^{n} a_{j k}(x) \xi_{j} \bar{\xi}_{k} \geq \nu|\xi|^{2} \quad \text { for all } \xi \in \mathbb{C}^{n}, x \in \Omega \tag{7.3}
\end{equation*}
$$

We apply the abstract setting with $\overline{\mathcal{H}}=\bar{H}$. Consider the unbounded operator in $\bar{H}$ defined by

$$
\begin{equation*}
\left(\mathcal{A}_{0} \Psi\right)(t, x):=-\sum_{j, k=1}^{n} \frac{\partial}{\partial x_{k}}\left[a_{j k}(x) \frac{\partial \Psi(t, x)}{\partial x_{j}}\right], \quad(t, x) \in(0, T) \times \Omega \tag{7.4}
\end{equation*}
$$

with domain $\operatorname{dom}\left(\mathcal{A}_{0}\right):=\bar{H}^{2}(\Omega) \cap \bar{V}$, where $\bar{H}^{2}(\Omega)$ denotes the complex valued Sobolev space $H^{2}(\Omega, \mathbb{C})$. One easily checks that this operator is self-adjoint, i.e., equal to the conjugate transpose. The $\operatorname{PDE}$ (7.1) enters in the semigroup framework, with generator

$$
\begin{equation*}
\left(\mathcal{A}_{\bar{H}} \Psi\right):=i \mathcal{A}_{0} \Psi \quad \text { for all } \Psi \in \bar{H} \tag{7.5}
\end{equation*}
$$

Lemma 7.1. The operator $\mathcal{A}_{\overline{\mathcal{H}}}$, with domain $\operatorname{dom}\left(\mathcal{A}_{\overline{\mathcal{H}}}\right):=\bar{H}^{2}(\Omega) \cap \bar{V}$, is the generator of a unitary semigroup and (7.1) has a mild solution $\Psi \in C(0, T ; \bar{H})$.

Proof. By the Hille-Yosida theorem (Pazy [41]), $\mathcal{A}_{\overline{\mathcal{H}}}$ is the generator of a contracting semigroup iff, for all $\lambda>0,\left(\lambda I+\mathcal{A}_{\overline{\mathcal{H}}}\right)$ has a continuous inverse that satisfies

$$
\begin{equation*}
\left\|\left(\lambda I+\mathcal{A}_{\overline{\mathcal{H}}}\right)^{-1}\right\|_{\mathcal{L}(\overline{\mathcal{H}})} \leq 1 / \lambda \tag{7.6}
\end{equation*}
$$

This is easily checked. In addition, the operator $\mathcal{A}_{\overline{\mathcal{H}}}$ being the opposite of its conjugate transpose it follows that the semigroup is norm preserving.

We define then the following sesquilinear form over $\bar{V}$ :

$$
\begin{equation*}
a(y, z):=\sum_{j, k=1}^{n} \int_{\Omega} a_{j k}(x) \frac{\partial y}{\partial x_{j}} \frac{\partial \bar{z}}{\partial x_{k}} \mathrm{~d} x \quad \text { for all } y, z \text { in } \bar{V} \tag{7.7}
\end{equation*}
$$

which is self-adjoint in the sense that

$$
\begin{equation*}
\overline{a(y, z)}=a(z, y) \tag{7.8}
\end{equation*}
$$

Furthermore, for $y, z$ in $\operatorname{dom}\left(\mathcal{A}_{0}\right)$ we have that

$$
\begin{equation*}
\left\langle\mathcal{A}_{0} y, z\right\rangle_{\bar{H}}=a(y, z)=\overline{a(z, y)}=\left\langle y, \mathcal{A}_{0} z\right\rangle_{\bar{H}} \tag{7.9}
\end{equation*}
$$

so that $\mathcal{A}_{0}$ is also self-adjoint.

### 7.2. Link to variational setting and regularity for Schrödinger equation.

 We introduce the function space$$
\begin{equation*}
\mathcal{X}:=L^{\infty}(0, T ; \bar{V}) \cap H^{1}\left(0, T ; \bar{V}^{\prime}\right) \tag{7.10}
\end{equation*}
$$

endowed with the natural norm

$$
\begin{equation*}
\|\Psi\|_{\mathcal{X}}:=\|\Psi\|_{L^{\infty}(0, T ; \bar{V})}+\|\Psi\|_{H^{1}\left(0, T ; \bar{V}^{\prime}\right)} \tag{7.11}
\end{equation*}
$$

There holds the following weak convergence result.
Lemma 7.2. Let $\left(\Psi_{k}\right)$ be a bounded sequence in $\mathcal{X}$. Then there exists $\Psi \in \mathcal{X}$ such that a subsequence of $\Psi_{k}$ converges to $\Psi$ strongly in $L^{2}(0, T ; \bar{H})$, and weakly in $L^{2}(0, T ; \bar{V})$ and $H^{1}\left(0, T ; \bar{V}^{\prime}\right)$. Finally, if $u_{k}$ weakly* converges to $u$ in $L^{\infty}(0, T)$, then

$$
\begin{equation*}
u_{k} b_{2} \Psi_{k} \rightarrow u b_{2} \Psi \quad \text { weakly in } L^{2}(0, T ; \bar{H}) \tag{7.12}
\end{equation*}
$$

Proof. By the Aubin-Lions lemma [6], $\mathcal{X}$ is compactly embedded into $L^{2}(0, T ; \bar{H})$. Thus, extracting a subsequence if necessary, we may assume that $\Psi_{k}$ converges in $L^{2}(0, T ; \bar{H})$ to some $\Psi$. Since $\Psi_{k}$ is bounded in the Hilbert spaces $L^{2}(0, T ; \bar{V})$ and $H^{1}\left(0, T ; \bar{V}^{\prime}\right)$, reextracting a subsequence if necessary, we may assume that it also weakly converges in these spaces.

Let $C_{R}$ denote the closed ball of $L^{\infty}(0, T, \bar{V})$ of radius $R$. This is a closed subset of $L^{2}(0, T, \bar{V})$ that, for large enough $R$, contains the sequence $\Psi_{k}$. Since any closed convex set is weakly closed, $\Psi \in C_{R}$. Thus $\Psi \in \mathcal{X}$. That (7.12) holds follows from the joint convergence of $u_{k}$ in $L^{\infty}(0, T)$ (endowed with the weak* topology), and of $\Psi_{k}$ in $L^{2}(0, T ; \bar{H})$.

The variational solution of (7.1) is given as $\Psi \in \mathcal{X}$ satisfying, for a.a. $t \in(0, T)$,

$$
\begin{equation*}
\langle\dot{\Psi}(t), z\rangle_{\bar{V}}+i a(\Psi(t), z)+i u(t)\left\langle b_{2} \Psi, z\right\rangle_{\bar{H}}=\langle f(t), z\rangle_{\bar{V}} \text { for all } z \in \bar{V}, \tag{7.13}
\end{equation*}
$$

and $\Psi(0)=\Psi_{0} \in \bar{V}$.
For $\left(f, b_{2}, u, \Psi_{0}\right) \in L^{2}(0, T ; \bar{V}) \times W^{1, \infty}(\Omega) \times L^{\infty}(\Omega) \times \bar{V}$ we set

$$
\begin{align*}
\kappa\left[f, b_{2}, u, \Psi_{0}\right]= & \|f\|_{L^{1}(0, T ; \bar{V})}^{2}+\left\|\Psi_{0}\right\|_{V}^{2} \\
& +\|u\|_{L^{\infty}(0, T)}^{2}\left\|\nabla b_{2}\right\|_{L^{\infty}(\Omega)}^{2}\left(\|f\|_{L^{2}(0, T ; \bar{H})}^{2}+\left\|\Psi_{0}\right\|_{\bar{V}}^{2}\right) . \tag{7.14}
\end{align*}
$$

There holds the following existence and regularity result for the unique solution of (7.13) (cf. [40]).

Theorem 7.3. Let $\left(f, b_{2}, u, \Psi_{0}\right) \in L^{2}(0, T ; \bar{V}) \times W^{1, \infty}(\Omega) \times L^{\infty}(\Omega) \times \bar{V}$. Then there exists $c_{0}>0$ independent of $\left(f, b_{2}, u, \Psi_{0}\right)$ such that (7.13) has a unique variational solution $\Psi$ in $\mathcal{X}$, that satisfies the estimates

$$
\begin{gather*}
\|\Psi\|_{C(0, T ; \bar{H})} \leq c_{0}\left(\|f\|_{L^{1}(0, T ; \bar{H})}+\left\|\Psi_{0}\right\|_{\bar{H}}\right),  \tag{7.15}\\
\|\Psi\|_{C(0, T ; \bar{V})}+\|\dot{\Psi}(t)\|_{L^{2}\left(0, T ; \bar{V}^{\prime}\right)} \leq c_{0} \kappa\left[f, b_{2}, u, \Psi_{0}\right] . \tag{7.16}
\end{gather*}
$$

Proof. Since $\Omega$ is bounded, there exists a Hilbert basis of $H_{0}^{1}(\Omega)\left(w_{j}, \lambda_{j}\right), j \in \mathbb{N}$, of (real) eigenvalues and nonnegative eigenvectors of the operator $\mathcal{A}_{0}$ (with, by the definition, homogeneous Dirichlet conditions), i.e.,

$$
\begin{equation*}
-\sum_{j, k=1}^{n} \frac{\partial}{\partial x_{k}}\left[a_{j k}(x) \frac{\partial w_{j}(x)}{\partial x_{j}}\right]=\lambda_{j} w_{j}(x), \quad j=1, \ldots, w_{j} \in H_{0}^{1}(\Omega), \quad \lambda_{j} \in \mathbb{R}_{+} \tag{7.17}
\end{equation*}
$$

Consider the associated Faedo-Galerkin discretization method; that is, let $\left\{\bar{V}_{k}\right\}$ be the finite dimensional subspaces of $\bar{V}$ generated by the (complex combinations of the) $w_{j}$ for $j \leq k$. The corresponding approximate solution $\Psi_{k}(t)=\sum_{j=1}^{k} \psi_{k}^{j}(t) w_{j}$ of (7.1) with $\psi_{k}^{j}(t) \in \mathbb{C}$, is defined as the solution of

$$
\begin{equation*}
\left\langle\dot{\Psi}_{k}(t), w_{j}\right\rangle_{\bar{H}}+i a\left(\Psi_{k}(t), w_{j}\right)+i u(t)\left\langle b_{2} \Psi_{k}(t), w_{j}\right\rangle_{\bar{H}}=\left\langle f(t), w_{j}\right\rangle_{\bar{H}} \tag{7.18}
\end{equation*}
$$

for $j=1, \ldots, k$ and $t \in[0, T]$, with initial condition

$$
\begin{equation*}
\psi_{k}^{j}(0)=\left(\Psi_{0}, w_{j}\right) \quad \text { for } j=1, \ldots, k . \tag{7.19}
\end{equation*}
$$

For each $k \in \mathbb{N}$, the above equations are a system of linear ordinary differential equations that has a unique solution $\psi_{k}=\left(\psi_{k}^{1}, \ldots, \psi_{k}^{k}\right) \in C\left(0, T ; \mathbb{C}^{k}\right)$. It follows that for any $\Phi(t)=\sum_{j=1}^{k} \phi^{j}(t) w_{j}$ (where $\phi^{j}(t) \in L^{1}(0, T)$ for $j=1, \ldots, k$ ) we have that

$$
\begin{equation*}
\left\langle\dot{\Psi}_{k}(t), \Phi(t)\right\rangle_{\bar{H}}+i a\left(\Psi_{k}(t), \Phi(t)\right)+i u(t)\left\langle b_{2} \Psi_{k}(t), \Phi(t)\right\rangle_{\bar{H}}=\langle f(t), \Phi(t)\rangle_{\bar{H}} . \tag{7.20}
\end{equation*}
$$

We derive a priori estimates by using different test functions $\Phi$ :

1. Testing with $\Phi(t)=\Psi_{k}(t)$ gives
(7.21)
$\left\langle\dot{\Psi}_{k}(t), \Psi_{k}(t)\right\rangle_{\bar{H}}+i a\left(\Psi_{k}(t), \Psi_{k}(t)\right)+i u(t)\left\langle b_{2} \Psi_{k}(t), \Psi_{k}(t)\right\rangle_{\bar{H}}=\left\langle f(t), \Psi_{k}(t)\right\rangle_{\bar{H}}$.
Taking the real part in both sides in (7.21) we obtain

$$
\begin{equation*}
\frac{1}{2} \frac{\mathrm{~d}}{\mathrm{~d} t}\left\|\Psi_{k}(t)\right\|_{\bar{H}}^{2} \leq C_{1}\|f(t)\|_{\bar{H}}\left\|\Psi_{k}(t)\right\|_{\bar{H}} \leq C_{2}\left(\|f(t)\|_{\bar{H}}^{2}+\left\|\Psi_{k}(t)\right\|_{\bar{H}}^{2}\right) . \tag{7.22}
\end{equation*}
$$

By Gronwall's inequality we get the following estimate:

$$
\begin{equation*}
\left\|\Psi_{k}\right\|_{L^{\infty}(0, T ; \bar{H})}^{2} \leq C_{3}\left(\|f\|_{L^{1}(0, T ; \bar{H})}^{2}+\left\|\Psi_{k}(0)\right\|_{\tilde{H}^{2}}^{2}\right) \tag{7.23}
\end{equation*}
$$

2. Testing with $\Phi(t)=\sum_{j=1}^{k} \lambda_{j} \psi_{k}^{j}(t) w_{j}=\mathcal{A}_{0} \Psi_{k}(t)$ gives
(7.24)
$\left\langle\dot{\Psi}_{k}(t), \mathcal{A}_{0} \Psi_{k}(t)\right\rangle_{\bar{H}}+i a\left(\Psi_{k}(t), \mathcal{A}_{0} \Psi_{k}(t)\right)+i u(t)\left(b_{2} \Psi_{k}(t)-f(t), \mathcal{A}_{0} \Psi_{k}(t)\right)_{\bar{H}}=0$.
Applying (7.9) (in both directions) we get

$$
\begin{align*}
i\left\langle\mathcal{A}_{0} \Psi_{k}(t), \mathcal{A}_{0} \Psi_{k}(t)\right\rangle_{\bar{H}} & +a\left(\dot{\Psi}_{k}(t), \Psi_{k}(t)\right)  \tag{7.25}\\
& +i u(t) a\left(b_{2} \Psi_{k}(t), \Psi_{k}(t)\right)-a\left(f(t), \Psi_{k}(t)\right)=0
\end{align*}
$$

Since $a(\cdot, \cdot)$ is self-adjoint we have that

$$
\begin{align*}
\frac{\mathrm{d}}{\mathrm{~d} t} a\left(\Psi_{k}(t), \Psi_{k}(t)\right) & =a\left(\Psi_{k}(t), \dot{\Psi}_{k}(t)\right)+a\left(\dot{\Psi}_{k}(t), \Psi_{k}(t)\right) \\
& =2 \Re\left(a\left(\Psi_{k}(t), \dot{\Psi}_{k}(t)\right)\right) \tag{7.26}
\end{align*}
$$

So, taking the real parts in (7.25) we get, using Young's inequality and the coercivity of $a(\cdot, \cdot)$ over $\bar{V}$,

$$
\begin{align*}
\frac{1}{2} \frac{\mathrm{~d}}{\mathrm{~d} t} a\left(\Psi_{k}(t), \Psi_{k}(t)\right) & =-\Re\left(a\left(\Psi_{k}(t), i u(t) b_{2} \Psi_{k}(t)-f(t)\right)\right) \\
& \leq c\left\|\Psi_{k}(t)\right\|_{\bar{V}}\left(\left\|\Psi_{k}(t)\right\|_{\bar{V}}+\|f(t)\|_{\bar{V}}\right)  \tag{7.27}\\
& \leq c^{\prime}\left(a\left(\Psi_{k}(t), \Psi_{k}(t)\right)+\|f(t)\|_{\bar{V}}\right)
\end{align*}
$$

So, by Gronwall's estimate and using (7.23),

$$
\begin{equation*}
\left\|\Psi_{k}\right\|_{L^{\infty}(0, T ; \bar{V})} \leq c_{0} \kappa\left[f, b_{2}, u, \Psi_{0}\right] . \tag{7.28}
\end{equation*}
$$

3. Any $\Phi \in \bar{V}$ can be written as $\Phi=\Phi^{1}+\Phi^{2}$ with $\Phi^{1} \in \bar{V}_{j}$ and $\Phi^{2}$ orthogonal to $\bar{V}_{j}$ in both spaces $\bar{H}$ and $\bar{V}$. Recall the notation for the dual and antidual pairing introduced in section 4 . Then

$$
\begin{equation*}
\left\langle\dot{\Psi}_{k}(t), \Phi\right\rangle_{\bar{V}}=\left\langle\dot{\Psi}_{k}(t), \Phi\right\rangle_{\bar{H}}=\left\langle\dot{\Psi}_{k}(t), \Phi^{1}\right\rangle_{\bar{H}}=\left\langle\dot{\Psi}_{k}(t), \Phi^{1}\right\rangle_{\bar{V}} \tag{7.29}
\end{equation*}
$$

It follows from (7.20) that there exists $c^{\prime \prime}>0$ such that, when $\|\Phi\|_{\bar{V}} \leq 1$,

$$
\begin{equation*}
\left\langle\dot{\Psi}_{k}(t), \Phi\right\rangle_{\bar{V}} \leq c^{\prime \prime}\left(\left\|\Psi_{k}(t)\right\|_{\bar{V}}+\|u\|_{L^{\infty}(0, T)}\left\|b_{2}\right\|_{L^{\infty}(\Omega)}\left\|\Psi_{k}(t)\right\|_{\bar{H}}+\|f(t)\|_{\bar{H}}\right) \tag{7.30}
\end{equation*}
$$

Combining with the above estimates we obtain

$$
\begin{equation*}
\left\|\dot{\Psi}_{k}\right\|_{L^{2}\left(0, T ; \bar{V}^{\prime}\right)} \leq c_{0} \kappa\left[f, b_{2}, u, \Psi_{0}\right] . \tag{7.31}
\end{equation*}
$$

By Lemma 7.2 a subsequence of $\left(\Psi_{k}\right)$ strongly converges in $L^{2}(0, T ; \bar{H})$ and weakly in $L^{2}(0, T ; \bar{V}) \cap H^{1}\left(0, T ; \bar{V}^{\prime}\right)$, while $u b_{2} \Psi_{k} \rightarrow u b_{2} \Psi$ weakly in $L^{2}(0, T ; \bar{H})$. Passing to the limit in (7.20) we obtain that $\Psi$ is the solution of the Schrödinger equation. That $\Psi$ is unique, belongs to $\mathcal{X}$, and satisfies (7.15), (7.16), and (7.31) follows from the same techniques as those used in the study of the Faedo-Galerkin approximation.

Lemma 7.4. For $\left(f, b_{2}, u, \Psi_{0}\right) \in L^{2}(0, T ; \bar{V}) \times W^{1, \infty}(\Omega) \times L^{\infty}(\Omega) \times \bar{V}$ the mild solution coincides with the variational solution.

Proof. That the variational and mild solutions coincide can be shown by an argument similar to [5, Lemma 10].

The corresponding data of the abstract theory are $\mathcal{B}_{1} \in \bar{H}$ equal to zero and $\mathcal{B}_{2} \in \mathcal{L}(\bar{H})$ defined by $\left(\mathcal{B}_{2} \Psi\right)(x):=-i b_{2}(x) \Psi(x)$ for $\Psi$ in $\bar{H}$ and $x \in \Omega$. The cost function is, given $\alpha_{1} \in \mathbb{R}$,

$$
\begin{array}{rl}
J(u, \Psi):=\alpha_{1} \int_{0}^{T} & u(t) \mathrm{d} t+\frac{1}{2} \int_{(0, T) \times \Omega}\left(\Psi(t, x)-\Psi_{d}(t, x)\right)^{2} \mathrm{~d} x \mathrm{~d} t  \tag{7.32}\\
& +\frac{1}{2} \int_{\Omega}\left(\Psi(T, x)-\Psi_{d T}(x)\right)^{2} \mathrm{~d} x
\end{array}
$$

We assume that

$$
\begin{equation*}
\Psi_{d} \in C(0, T ; \bar{V}), \quad \Psi_{d T} \in \bar{V} \tag{7.33}
\end{equation*}
$$

For $u \in L^{1}(0, T)$, write the reduced cost as $F(u):=J(u, \Psi[u])$. The optimal control problem is, $\mathcal{U}_{a d}$ being defined in (4.8),

$$
\begin{equation*}
\operatorname{Min} F(u), \quad u \in \mathcal{U}_{a d} \tag{7.34}
\end{equation*}
$$

7.3. Compactness for the Schrödinger equation. To prove the existence of an optimal control of $(\mathrm{P})$ we have to verify the compactness hypothesis (4.10).

Proposition 7.5. Problem (P) for (7.1) and cost function (7.32) has a nonempty set of minima.

Proof. This follows from Theorem 4.2, whose compactness hypothesis holds thanks to Lemma 7.2.
7.4. Commutators. Given $\Psi \in \operatorname{dom}\left(\mathcal{A}_{\overline{\mathcal{H}}}\right)$, we have by (7.5) that

$$
\begin{equation*}
M_{1} \Psi=-\sum_{j, k=1}^{n}\left(\frac{\partial b_{2}}{\partial x_{k}}\left[a_{j k} \frac{\partial \Psi}{\partial x_{j}}\right]+\frac{\partial}{\partial x_{k}}\left[a_{j k} \Psi \frac{\partial b_{2}}{\partial x_{j}}\right]\right) . \tag{7.35}
\end{equation*}
$$

As expected, this commutator is a first order differential operator that has a continuous extension to the space $\bar{V}$. In a similar way we can check that $\left[M_{1}, \mathcal{B}_{2}\right]$ is the "zero order" operator given by

$$
\begin{equation*}
\left[M_{1}, \mathcal{B}_{2}\right] \Psi=2 i \sum_{j, k=1}^{n} a_{j, k} \frac{\partial b_{2}}{\partial x_{j}} \frac{\partial b_{2}}{\partial x_{k}} \Psi \tag{7.36}
\end{equation*}
$$

Remark 7.6. In the case of the Laplace operator, i.e., when $a_{j k}=\delta_{j k}$, we find that for $\Psi \in \bar{V}$

$$
\begin{equation*}
M_{1} \Psi=-2 \nabla b_{2} \cdot \nabla \Psi-\Psi \Delta b_{2} ; \quad\left[M_{1}, \mathcal{B}_{2}\right] \Psi=2 i \Psi\left|\nabla b_{2}\right|^{2} \tag{7.37}
\end{equation*}
$$

and then for $p \in \bar{V}$ we have

$$
\begin{equation*}
M_{1}^{*} p=2 \nabla b_{2} \cdot \nabla \bar{p}+\bar{p} \Delta b_{2} \tag{7.38}
\end{equation*}
$$

Similarly, we have

$$
\left\{\begin{align*}
M_{2} \Psi & =2 i \nabla b_{2}^{2} \cdot \nabla \Psi+i \Psi \Delta b_{2}^{2}  \tag{7.39}\\
{\left[M_{2}, \mathcal{B}_{2}\right] \Psi } & =-2 i \Psi\left|\nabla b_{2}^{2}\right|^{2} \\
M_{2}^{*} p & =-i\left(2 \nabla b_{2}^{2} \cdot \nabla \bar{p}+\bar{p} \Delta b_{2}^{2}\right)
\end{align*}\right.
$$

7.5. Analysis of optimality conditions. For the sake of simplicity we only discuss the case of the Laplace operator. The costate equation is then

$$
\begin{equation*}
-\dot{p}+i \Delta p=\Psi-\Psi_{d}+i u b_{2} p \text { in }(0, T) \times \Omega, \quad p(T)=\Psi(T)-\Psi_{d T} \tag{7.40}
\end{equation*}
$$

Remembering the expression of $b_{z}^{1}$ in (6.7), we obtain that the equation for $\xi:=\xi_{z}$ introduced in (6.6) reduces to

$$
\begin{equation*}
\dot{\xi}-i \Delta \xi=-i \hat{u} b_{2} \xi+w\left(i b_{2} f+2 \nabla b_{2} \cdot \nabla \Psi+\Psi \Delta b_{2}\right) \quad \text { in }(0, T) \times \Omega, \quad \xi(0)=0 \tag{7.41}
\end{equation*}
$$

The quadratic forms $\mathcal{Q}$ and $\widehat{\mathcal{Q}}$ defined in (5.6) and (6.9) are as follows. First

$$
\begin{equation*}
\mathcal{Q}(z, v)=\int_{0}^{T}\left(\|z(t)\|_{\bar{H}}^{2}+2 v(t) \Re\left\langle\hat{p}(t), b_{2} z(t)\right\rangle_{\bar{H}}\right) \mathrm{d} t+\|z(T)\|_{\bar{H}}^{2} \tag{7.42}
\end{equation*}
$$

and, second,

$$
\begin{equation*}
\widehat{\mathcal{Q}}(\xi, w, h)=\widehat{\mathcal{Q}}_{T}(\xi, h)+\widehat{\mathcal{Q}}_{a}(\xi, w)+\widehat{\mathcal{Q}}_{b}(w), \quad \widehat{\mathcal{Q}}_{b}(w):=\int_{0}^{T} w^{2}(t) R(t) \mathrm{d} t \tag{7.43}
\end{equation*}
$$

Here $R \in C(0, T)$ and

$$
\begin{equation*}
\widehat{\mathcal{Q}}_{T}(\xi, h):=\left\|\xi(T)-i h b_{2} \hat{\Psi}(T)\right\|_{\bar{H}}^{2}-h^{2} \Re\left\langle\hat{p}(T), b_{2}^{2} \hat{\Psi}(T)\right\rangle_{\bar{H}}+h \Re\left\langle i \hat{p}(T), b_{2} \xi(T)\right\rangle_{\bar{H}} \tag{7.44}
\end{equation*}
$$

$$
\begin{align*}
\widehat{\mathcal{Q}}_{a}(\xi, w) & :=\int_{0}^{T}\left(\|\xi\|_{\bar{H}}^{2}+2 w \Re\left(i\left\langle\xi, b_{2} \hat{\Psi}\right\rangle_{\bar{H}}+i\left\langle\hat{\Psi}-\Psi_{d}, b_{2} \xi\right\rangle_{\bar{H}}-\left\langle M_{1}^{*} \hat{p}, \xi\right\rangle_{\bar{H}}\right)\right) \mathrm{d} t  \tag{7.45}\\
46) \quad R(t) & \left.:=\left\|b_{2} \hat{\Psi}\right\|_{\bar{H}}^{2}-\Re\left\langle\hat{\Psi}-\Psi_{d}, b_{2}^{2} \hat{\Psi}\right\rangle_{\bar{H}}+\left.\Re\left\langle\hat{p}(t),-b_{2}^{2} f(t)-2 i\right| \nabla b_{2}\right|^{2} \hat{\Psi}\right\rangle_{\bar{H}} \tag{7.46}
\end{align*}
$$

THEOREM 7.7 (second order necessary and sufficient conditions). Let $\hat{u} \in \mathcal{U}_{\text {ad }}$.
(i) If $\hat{u}$ is a weak minimum then the second order necessary conditions (6.16) and (6.18) hold.
(ii) Let $\hat{u}$ satisfy the first order necessary optimality conditions (4.15) and assume the hypotheses of Theorem 6.7. Then, the uniform positive definiteness condition on $\widehat{P C}_{2}(\hat{u})$ in (6.25) implies the weak quadratic growth (6.26).
(iii) Conversely, if $\hat{u}$ is a weak minimum satisfying the weak quadratic growth condition, (6.26) implies the uniform positive definiteness on $P C_{2}(\hat{u})$ in (6.24).

Proof. (i) Conditions (6.1)(i) and (ii) are satisfied with (7.2). Since we have

$$
\begin{equation*}
\overline{\left[-i \Delta,\left(-i b_{2}\right)^{k}\right]} \hat{\Psi}=-(-i)^{k-1}\left(\Delta b_{2}^{k} \hat{\Psi}+2 \nabla b_{2}^{k} \nabla \hat{\Psi}\right), \quad k=1,2 \tag{7.47}
\end{equation*}
$$

i.e., the commutator is a first order differential operator and has an extension to the space $\bar{V}$, we obtain (6.3)(i) with $E_{1}=\bar{V}$. (6.3)(ii) and (iii) follow from the regularity assumptions in (7.2) and (7.33).

The compactness hypothesis (6.17) for

$$
\begin{equation*}
w \mapsto \xi[w], \quad L^{2}(0, T) \rightarrow L^{2}(0, T ; \bar{H}) \tag{7.48}
\end{equation*}
$$

follows from $(7.2)$, since $\xi[w] \in L^{2}(0, T ; \bar{V}) \cap H^{1}\left(0, T ; \bar{V}^{\prime}\right)$ which is compactly embedded in $L^{2}(0, T ; \bar{H})$ by Aubin's lemma [6].

Thus, item (i) of the current theorem follows from Theorem 6.4.
(ii) and (iii). We apply Theorem 6.7. We already checked hypotheses (6.1) and (6.3) in this proof. Condition (6.5) follows also from the assumptions in (7.2) and (7.33).

Remark 7.8. It is not difficult to extend such results for more general differential operators of the type, where the $a_{j k}$ are as before, $b \in L^{\infty}(\Omega)^{n}$, and $c \in L^{\infty}(\Omega)$ :

$$
\begin{equation*}
\left(\mathcal{A}_{\overline{\mathcal{H}}} \Psi\right)(t, x)=-i \sum_{j, k=1}^{n} \frac{\partial}{\partial x_{k}}\left[a_{j k}(x) \frac{\partial}{\partial x_{j}} \Psi(t, x)\right]+\sum_{j=1}^{n} \frac{\partial\left(b_{j}(x) \Psi(t, x)\right)}{\partial x_{j}}+c \Psi(t, x) \tag{7.49}
\end{equation*}
$$

## REFERENCES

[1] M. S. Aronna, J. F. Bonnans, and B. S. Goh, Second order analysis of control-affine problems with scalar state constraint, Math. Program., 160 (2016), pp. 115-147.
[2] M. S. Aronna, J. F. Bonnans, A. V. Dmitruk, and P. A. Lotito, Quadratic order conditions for bang-singular extremals, Numer. Algebra, Control Optim., 2 (2012), pp. 511-546.
[3] M. S. Aronna, J. F. Bonnans, and A. Kröner, Optimal control of bilinear systems in a complex space setting, IFAC-PapersOnLine, 50 (2017), pp. 2872-2877.
[4] M. S. Aronna, J. F. Bonnans, and A. Kröner, Correction to: Optimal control of infinite dimensional bilinear systems: Application to the heat and wave equations, Math. Program., 170 (2018), pp. 569-570.
[5] M. S. Aronna, J. F. Bonnans, and A. Kröner, Optimal control of infinite dimensional bilinear systems: Application to the heat and wave equations, Math. Program., 168 (2018), pp. 717-757.
[6] J.-P. Aubin, Un théorème de compacité, C. R. Acad. Sci. Paris, 256 (1963), pp. 5042-5044.
[7] J. M. Ball, J. E. Marsden, and M. Slemrod, Controllability for distributed bilinear systems, SIAM J. Control Optim., 20 (1982), pp. 575-597.
[8] J. M. Ball, Strongly continuous semigroups, weak solutions, and the variation of constants formula, Proc. Amer. Math. Soc., 63 (1977), pp. 370-373.
[9] L. Baudouin, O. Kavian, and J.-P. Puel, Regularity for a Schrödinger equation with singular potentials and application to bilinear optimal control, J. Differential Equations, 216 (2005), pp. 188-222.
[10] L. Baudouin and J. Salomon, Constructive solution of a bilinear optimal control problem for a Schrödinger equation, Systems Control Lett., 57 (2008), pp. 453-464.
[11] T. Bayen, J. F. Bonnans, and F. J. Silva, Characterization of local quadratic growth for strong minima in the optimal control of semi-linear elliptic equations, Trans. Amer. Math. Soc., 366 (2014), pp. 2063-2087.
[12] T. Bayen and F. J. Silva, Second Order Analysis for Strong Solutions in the Optimal Control of Parabolic Equations, SIAM J. Control Optim., 54 (2016), pp. 819-844.
[13] K. Beauchard, Contribution à létude de la contrôlabilité et de la stabilisation de l'équation de Schrödinger, Ph.D. thesis, Université Paris-Sud, Paris, 2005.
[14] K. Beauchard, J. M. Coron, and H. Teismann, Minimal time for the bilinear control of Schrödinger equations, Systems Control Lett., 71 (2014), pp. 1-6.
[15] K. Beauchard and M. Morancey, Local controllability of $1 D$ Schrödinger equations with bilinear control and minimal time, Math. Control Relat. Fields, 4 (2014), pp. 125-160.
[16] M. Bergounioux and D. Tiba, General optimality conditions for constrained convex control problems, SIAM J. Control Optim., 34 (1996), pp. 698-711.
[17] J. F. Bonnans, Optimal control of a semilinear parabolic equation with singular arcs, Optim. Methods Softw., 29 (2014), pp. 964-978.
[18] J. F. Bonnans and P. Jaisson, Optimal control of a parabolic equation with time-dependent state constraints, SIAM J. Control Optim., 48 (2010), pp. 4550-4571.
[19] J. F. Bonnans and N. P. Osmolovskĭ̆, Second-order analysis of optimal control problems with control and initial-final state constraints, J. Convex Anal., 17 (2010), pp. 885-913.
[20] J. F. Bonnans and A. Shapiro, Perturbation Analysis of Optimization Problems, Springer Ser. Oper. Res., Springer, New York, 2000.
[21] J. F. Bonnans and D. Tiba, Control problems with mixed constraints and application to an optimal investment problem, Math. Rep. (Bucur.), 11 (2009), pp. 293-306.
[22] U. Boscain, M. Caponigro, and M. Sigalotti, A weak spectral condition for the controllability of the bilinear Schrödinger equation with application to the control of a rotating planar molecule, Comm. Math. Phys., 311 (2012), pp. 423-455.
[23] E. CASAS, Second order analysis for bang-bang control problems of PDEs, SIAM J. Control Optim., 50 (2012), pp. 2355-2372.
[24] E. Casas, C. Clason, and K. Kunisch, Parabolic control problems in measure spaces with sparse solutions, SIAM J. Control Optim., 51 (2013), pp. 28-63.
[25] E. Casas and F. Tröltzsch, Second-order optimality conditions for weak and strong local solutions of parabolic optimal control problems, Vietnam J. Math., 44 (2016), pp. 181-202.
[26] A. V. Dmitruk, Quadratic conditions for a weak minimum for singular regimes in optimal control problems, Sov. Math. Dokl., 18 (1977), pp. 418-422.
[27] A. V. Dmitruk, Quadratic conditions for the Pontryagin minimum in an optimal control problem linear with respect to control. II. Theorems on the relaxing of constraints on the equality, Izv. Akad. Nauk SSSR Ser. Mat., 51 (1987), pp. 812-832.
[28] N. Dunford and J. Schwartz, Linear Operators, Vol. I, Interscience, New York, 1958.
[29] H. O. Fattorini, Infinite dimensional linear control systems, North-Holland Math. Stud. 201, Elsevier, Amsterdam, 2005.
[30] H. O. Fattorini and H. Frankowska, Necessary conditions for infinite-dimensional control problems, Math. Control Signals Systems, 4 (1991), pp. 41-67.
[31] H. Frankowska and D. Tonon, The Goh necessary optimality conditions for the Mayer problem with control constraints, 52nd IEEE Conference on Decision and Control, IEEE, Piscataway, NJ, 2013, pp. 538-543.
[32] G. Friesecke, F. Henneke, and K. Kunisch, Frequency-sparse optimal quantum control, Math. Control Relat. Fields, 8 (2018), pp. 155-176.
[33] B. S. Goh, The second variation for the singular Bolza problem, J. SIAM Control, 4 (1966), pp. 309-325.
[34] H. Goldberg and F. Tröltzsch, Second-order sufficient optimality conditions for a class of nonlinear parabolic boundary control problems, SIAM J. Control Optim., 31 (1993), pp. 1007-1025.
[35] M. Hintermüller, D. Marahrens, P. A. Markowich, and C. Sparber, Optimal bilinear control of Gross-Pitaevskii equations, SIAM J. Control Optim., 51 (2013), pp. 2509-2543.
[36] K. Ito and K. Kunisch, Optimal bilinear control of an abstract Schrödinger equation, SIAM J. Control Optim., 46 (2007), pp. 274-287.
[37] H. J. Kelley, A second variation test for singular extremals, AIAA Journal, 2 (1964), pp. 13801382.
[38] X. Li and Y. Yao, Maximum principle of distributed parameter systems with time lags, Lect. Notes Control Inf. Sci. 75, Springer, New York, 1985, pp. 410-427.
[39] X. Li and J. Yong, Necessary conditions for optimal control of distributed parameter systems, SIAM J. Control Optim., 29 (1991), pp. 895-908.
[40] J.-L. Lions and E. Magenes, Non-Homogeneous Boundary Value Problems and Applications. Vol. I, Grundlehren Math. Wiss. 181, Springer, New York, 1972.
[41] A. Pazy, Semigroups of Linear Operators and Applications to Partial Differential Equations, Appl. Math. Sci. 44, Springer, New York, 1983.
[42] L. Poggiolini and G. Stefani, Sufficient optimality conditions for a bang-singular extremal in the minimum time problem, Control Cybernet., 37 (2008), pp. 469-490.
[43] S. M. Robinson, First order conditions for general nonlinear optimization, SIAM J. Appl. Math., 30 (1976), pp. 597-607.
[44] F. Tröltzsch, Regular Lagrange multipliers for control problems with mixed pointwise controlstate constraints, SIAM J. Optim., 15 (2004), pp. 616-634.
[45] F. Tröltzsch, Optimal Control of Partial Differential Equations, Grad. Stud. Math. 112, American Mathematical Society, Providence, RI, 2010.


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    ${ }^{\dagger}$ EMAp/FGV, Rio de Janeiro 22250-900, Brazil (soledad.aronna@fgv.br).
    ${ }^{\ddagger}$ Inria Saclay and CMAP, Ecole Polytechnique, CNRS, Université Paris Saclay, 91128 Palaiseau, France (Frederic.Bonnans@inria.fr).
    ${ }^{\S}$ Institute for Mathematics, Humboldt-Universität zu Berlin, 10099 Berlin, Germany. CMAP, Ecole Polytechnique, CNRS, Université Paris Saclay, and Inria, France (Axel.Kroener@math.huberlin.de).

