# STATE CONSTRAINED CONTROL PROBLEMS IN BANACH LATTICES AND APPLICATIONS* 

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#### Abstract

This paper aims to study a family of deterministic optimal control problems in infinite-dimensional spaces. The peculiar feature of such problems is the presence of a positivity state constraint, which often arises in economic applications. To deal with such constraints, we set up the problem in a Banach lattice, not necessarily reflexive: a typical example is the space of continuous functions on a compact set. In this setting, which seems to be new in this context, we are able to find explicit solutions to the Hamilton-Jacobi-Bellman (HJB) equation associated to a suitable auxiliary problem and to write the corresponding optimal feedback control. Thanks to a type of infinite-dimensional Perron-Frobenius theorem, we use these results to gain information about the optimal paths of the original problem. This was not possible in the infinite-dimensional setting used in earlier works on this subject, where the state space was an $\mathrm{L}^{2}$ space.


Key words. optimal control in infinite dimension, dynamic programming, state constraints, Hamilton-Jacobi-Bellman equation, Banach lattice, AK model of economic growth

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1. Introduction. A typical feature of optimal control problems arising in economic applications-as well as in other fields-is the presence of positivity state and/or control constraints. This feature is very common and makes the problem nontrivial - even in the case where the state space has finite dimension. Indeed, on the one hand, if one approaches the problem by the maximum principle, then jumps of the co-state variables, which are associated to optimality conditions, may arise (see, e.g., Hartl, Sethi, and Vickson [32] for the finite-dimensional case and Fattorini [25, Chap. 11] for the infinite-dimensional case); on the other hand, if one employs dynamic programming techniques, then well-posedness and regularity results for the associated HJB equation may be hard to obtain (see, e.g., Capuzzo-Dolcetta and Lions [15] and Soner [38] for the finite-dimensional case; Cannarsa and Di Blasio [13], Cannarsa, Gozzi, and Soner [14], Faggian [24], and Kocan and Soravia [34] for the infinite-dimensional case; see also Katsoulakis [33], Calvia [11, 12], and Fabbri, Gozzi, and Swiech [23, Chap. 3] for the stochastic case).

In the last decades, various papers in the economic literature considered optimal control problems in which the state variable is infinite-dimensional. This typically happens when one has to take into account heterogeneity: for instance, the evolution of a key economic variable, such as capital, may depend not only on time but also on its spatial position or its vintage (see, e.g., Bambi et al. [4, 3], Boucekkine et al. $[6,7,8,9]$, Feichtinger et al. [29], and Fabbri and Gozzi [22]). The setting in which

[^0]such optimal control problems are studied in this area of literature is the separable Hilbert space $L^{2}$. Also in this case, positivity state constraints represent an essential feature of the problem. Handling such constraints in this infinite-dimensional setting turns out to be even more problematic than in the finite dimension, as the positive cone of $\mathrm{L}^{2}$ has an empty interior, and it is not obvious how to interpret the derivative in the associated HJB equation. Consequently, it is not obvious how to make sense of the optimal feedback map, which, as is well known, depends on the derivative of the solution to the HJB equation.

A possible way to deal with this difficulty (cf., e.g., Bambi et al. [4, 3], Boucekkine et al. $[6,7,9]$, and Fabbri and Gozzi [22]) is to study an auxiliary problem under a relaxed state constraint by allowing the state variable to evolve in a suitable halfspace, where the solution of the HJB equation can be found explicitly. If it is possible to find the optimal paths for the auxiliary problem, one can try to show, at least for some initial data (hopefully interesting for applications), that these paths satisfy the positivity state constraint and hence are optimal for the initial problem, too. This important property, i.e., the admissibility of the auxiliary optimal paths for the original problem, has been rigorously established in Boucekkine et al. $[6,7]$ only for the steady states of the problem that constitute a one-dimensional set of initial conditions. To the best of our knowledge, only in [3] this point is successfully addressed for a larger nontrivial set of initial conditions but at the price of a technical, involved, and tailor-made argument. Thus, the main motivation of our paper is to rigorously prove that admissibility of the auxiliary optimal paths for the original problem holds for a sufficiently rich set of initial conditions in the context of economic growth problems in time-space, such as the ones introduced by Boucekkine et al. [6, 7]: these results are collected in section 5 (cf. Theorem 5.12 and Corollary 5.14).

To achieve this goal, we consider the problems in a state space different from $\mathrm{L}^{2}$; since heterogeneity is usually modeled by a compact metric space $D$, the natural candidate to use is the space of continuous functions on $D$, endowed with the sup-norm and the standard pointwise partial order. This norm allows us to deal properly with pointwise constraints, such as the positivity state constraint that characterizes our optimization problems. However, we realized that, for the class of control problems with the features described above, many results that hold in the setting of continuous functions still hold in the more general setting of Banach lattices. Since we think that this more general setting can be useful for future applications, we build our theory in this abstract setup. More specifically, we consider the following:

- as state space $X$, a general separable Banach lattice of real-valued functions defined on a measure space $(D, \mu)$; the strictly positive orthant $X_{++}$is naturally defined and provides the required positivity state constraint;
- as state equation, the linear equation $x^{\prime}(t)=L x(t)-N c(t)$, where $x$ and $c$ are, respectively, the state and the control variable, and $L$ and $N$ are suitable operators;
- as objective functional, a type of discounted utility over an infinite horizon, i.e., a functional of the form $\int_{0}^{\infty} \mathrm{e}^{-\rho t} \mathcal{U}(c(t)) \mathrm{d} t$, where $\mathcal{U}$ is concave.

In this general setup, we define an appropriate auxiliary problem in which the state constraint is a suitable half-space of $X$ that strictly contains the positive cone. For this problem we prove a verification theorem (Theorem 3.1). Moreover, in the case when $X$ is either an $L^{p}$ space or the space of continuous functions, and $\mathcal{U}$ is homogeneous, we show that the HJB equation can be explicitly solved (Proposition 3.4) and,
consequently, that the optimal feedback control can be found (Theorem 3.9). ${ }^{1}$
Once the auxiliary problem in the half-space is explicitly solved, we are in a position to achieve our main goal, i.e., to prove a stability result (Theorem 3.11) that provides the basis to show admissibility of the auxiliary optimal paths for the original problem, at least when the initial conditions are in a neighborhood of the corresponding steady state (Theorem 5.12 and Corollary 5.14). To get these results, we need an important assumption that concerns the spectrum of the operator $L$ appearing in the state equation above (cf. point (i) of Assumption 3.2). This assumption can be verified thanks to Perron-Frobenius-type results in Banach spaces, which are provided in section 4.

The plan of the paper is the following:

- section 2 is devoted to presenting the general setup mentioned above;
- section 3 contains the main abstract results of the paper on the auxiliary problem: the verification theorem (subsection 3.1), the explicit solution to the auxiliary HJB equation (subsection 3.2), the stability result (subsection 3.3);
- section 4 introduces the Perron-Frobenius-type results guaranteeing the aforementioned key assumption needed in section 3;
- section 5 shows how the general theory exposed in the first four sections of the paper can be applied to a family of economic growth problems, including the one studied in Boucekkine et al. [6, 7].
- section 6 is devoted to discussing potential extensions of our approach to optimal control problems with positivity state constraints and nonlinear state equations.

2. The optimal control problem. Set $\mathbb{R}_{+}:=[0,+\infty)$, let $(D, \mu)$ be a measure space with a countably generated $\sigma$-algebra, ${ }^{2}$ and let $\left(X,|\cdot|_{X}, \leq_{X}\right)$ be a separable Banach lattice of real-valued functions defined on $D$. The symbols $|\cdot|_{X}$ and $\leq_{X}$ denote, respectively, the norm and the ordering on $X$.

Denote by $X^{\star}$ the topological dual of $X$, and let $\langle\cdot, \cdot\rangle$ be the dual pairing of $X, X^{\star}$. We have that $X^{\star}$ is an order-complete Banach lattice with its usual norm, denoted by $|\cdot|_{X^{\star}}$, and with the natural ordering $\leq_{X^{\star}}$ defined as follows (see, e.g., [2, p. 239]): given $\varphi^{\star}, \psi^{\star} \in X^{\star}$,

$$
\varphi^{\star} \leq_{X^{\star}} \psi^{\star} \quad \Longleftrightarrow \quad\left\langle f, \varphi^{\star}\right\rangle \leq\left\langle f, \psi^{\star}\right\rangle \quad \forall f \geq_{X} 0 .
$$

We consider the positive orthants (or positive cones) of $X, X^{\star}$, i.e., the sets

$$
X_{+}:=\left\{f \in X: f \geq_{x} 0\right\}, \quad X_{+}^{\star}:=\left\{\varphi^{\star} \in X^{\star}: \varphi^{\star} \geq X^{\star} 0\right\}
$$

and the strictly positive orthants (or strictly positive cones) of $X, X^{\star}$, i.e., the sets (see, e.g., [2, p. 119])

$$
\begin{aligned}
X_{++} & :=\left\{f \in X:\left\langle f, \varphi^{\star}\right\rangle>0 \quad \forall 0 \neq \varphi^{\star} \in X_{+}^{\star}\right\} \\
X_{++}^{\star} & :=\left\{\varphi^{\star} \in X^{\star}:\left\langle f, \varphi^{\star}\right\rangle>0 \quad \forall 0 \neq f \in X_{+}\right\}
\end{aligned}
$$

We write $f>_{X} 0$ if $f \in X_{++}$and $\varphi^{\star}>_{X^{\star}} 0$ if $\varphi^{\star} \in X_{++}^{\star}$.

[^1]Remark 2.1. The following spaces of functions, typical in applications, fall within the abstract setting above:
(i) $X=\mathrm{L}^{p}(D, \mu), p \in[1,+\infty)$, where $\mu$ is $\sigma$-finite. In this case $X$ is separable, $X^{\star}=\mathrm{L}^{q}(D, \mu)$, where $q$ satisfies $p^{-1}+q^{-1}=1$, and int $X_{+}=\emptyset$.
(ii) $X=\mathcal{C}(D)$, the space of real-valued continuous functions on a compact metric space $D$, equipped with the sup-norm; $\mu$ can be any Borel measure on $D$. In this case, $X$ is a Banach lattice with order unit, $X^{\star}$ is the space of regular Borel measures on $D$, and int $X_{+}=X_{++} \neq \emptyset$.
(iii) $X=\mathcal{C}_{0}(D)$, the space of real-valued continuous functions vanishing at infinity on a locally compact metric space $D$, equipped with the sup-norm; $\mu$ can be any Borel measure on $D$. In this case, $X^{\star}$ is the space of regular Borel measures on $D$, and int $X_{+}=\emptyset$.
(iv) Given a locally compact metric space $D$ and a continuous function $w: D \rightarrow$ $(0,+\infty)$, the space $X=\mathcal{C}_{w}(D)$ of real-valued continuous functions $f$ on $D$ such that $f w$ is bounded, equipped with the norm

$$
|f|_{X}:=\sup _{x \in D}|f(x) w(x)| .
$$

Also here $\mu$ can be any Borel measure on $D$. A typical case is $D=\mathbb{R}^{n}$ and $w(x)=\frac{1}{1+|x|^{\kappa}}$, which allows us to deal with continuous data with polynomial growth at infinity.
We introduce the optimal control problem in the space $X$ that we aim to study. Let $L: D(L) \subseteq X \rightarrow X$ be a (possibly) unbounded linear operator, and let $N: X \rightarrow X$ be a bounded linear operator. Given $x_{0} \in X$ and a control function $c \in \mathrm{~L}_{l o c}^{1}\left(\mathbb{R}_{+} ; X\right)$, we consider the following abstract state equation in $X$ :

$$
\left\{\begin{array}{l}
x^{\prime}(t)=L x(t)-N c(t), \quad t \geq 0,  \tag{2.1}\\
x(0)=x_{0} .
\end{array}\right.
$$

To stress the dependence of the solution to (2.1) on $x_{0} \in X$ and $c \in \mathrm{~L}_{l o c}^{1}\left(\mathbb{R}_{+} ; X\right)$, we will denote it by $x^{x_{0}, c}$. The following assumption will be in force throughout the paper.

Assumption 2.2.
(i) The linear operator $L: D(L) \subseteq X \rightarrow X$ is closed, densely defined, and generates a $C_{0}$-semigroup $\left\{\mathrm{e}^{t L}\right\}_{t \geq 0}$ in $X$. Moreover, the semigroup $\left\{\mathrm{e}^{t L}\right\}_{t \geq 0}$ preserves strict positivity, i.e., $\mathrm{e}^{\star \bar{L}}\left(X_{++}\right) \subseteq X_{++}$, for all $t \geq 0$.
(ii) The linear operator $N: X \rightarrow X$ is bounded and positive, i.e., $N\left(X_{+}\right) \subseteq X_{+}$.

Remark 2.3. It is worth noting that, by continuity, point (i) of Assumption 2.2 implies that $\left\{\mathrm{e}^{t L}\right\}_{t \geq 0}$ is also a positive semigroup; that is, it satisfies $\mathrm{e}^{t L}\left(X_{+}\right) \subseteq X_{+}$ for all $t \geq 0$. Sufficient conditions guaranteeing positivity of semigroups can be found, e.g., in [2, Chap. B-II, Thm. 1.6, Thm. 1.13, Chap. C-II, Thm. 1.2, Thm. 1.8]) and [16, Thm. 7.29 and Prop. 7.46].

According to [5, p. 129], for each $x_{0} \in X$ and $c \in \mathrm{~L}_{\text {loc }}^{1}\left(\mathbb{R}_{+} ; X\right)$, we call mild solution to (2.1) the function

$$
\begin{equation*}
x^{x_{0}, c}(t):=\mathrm{e}^{t L} x_{0}-\int_{0}^{t} \mathrm{e}^{(t-s) L} N c(s) \mathrm{d} s, \quad t \geq 0 . \tag{2.2}
\end{equation*}
$$

By [5, p. 131], the mild solution defined in (2.2) is also a weak solution to (2.1), i.e., it satisfies, for all $\varphi^{\star} \in D\left(L^{\star}\right)$,

$$
\begin{equation*}
\left\langle x^{x_{0}, c}(t), \varphi^{\star}\right\rangle=\left\langle x_{0}, \varphi^{\star}\right\rangle+\int_{0}^{t}\left\langle x^{x_{0}, c}(s), L^{\star} \varphi^{\star}\right\rangle \mathrm{d} s-\int_{0}^{t}\left\langle N c(s), \varphi^{\star}\right\rangle \mathrm{d} s, \quad t \geq 0 \tag{2.3}
\end{equation*}
$$

where $L^{\star}: D\left(L^{\star}\right) \subseteq X^{\star} \rightarrow X^{\star}$ denotes the adjoint of $L$.
We are interested in analyzing an optimal control problem in which the state variable, satisfying (2.1), remains in the strictly positive cone of $X$. This state constraint is expressed introducing the following class of admissible controls, which depend on $x_{0} \in X$ :

$$
\begin{equation*}
\mathcal{A}_{++}\left(x_{0}\right):=\left\{c \in \mathrm{~L}_{l o c}^{1}\left(\mathbb{R}_{+} ; X_{+}\right): x^{x_{0}, c}(t) \in X_{++} \forall t \geq 0\right\} \tag{2.4}
\end{equation*}
$$

For the reasons anticipated in the introduction and that will become clearer later, we are led to consider a larger class of admissible controls. Precisely, given $x_{0} \in X$ and $\varphi^{\star} \in X_{++}^{\star}$, we define the set

$$
\begin{align*}
\mathcal{A}_{++}^{\varphi^{\star}}\left(x_{0}\right) & :=\left\{c \in \mathrm{~L}_{l o c}^{1}\left(\mathbb{R}_{+} ; X_{+}\right):\left\langle x^{x_{0}, c}(t), \varphi^{\star}\right\rangle>0 \forall t \geq 0\right\}  \tag{2.5}\\
& =\left\{c \in \mathrm{~L}_{l o c}^{1}\left(\mathbb{R}_{+} ; X_{+}\right): x^{x_{0}, c}(t) \in X_{++}^{\varphi^{\star}} \forall t \geq 0\right\}
\end{align*}
$$

where $X_{++}^{\varphi^{\star}}$ is the open (infinite-dimensional) half-space of $X$ generated by $\varphi^{\star} \in X_{++}^{\star}$, i.e.,

$$
X_{++}^{\varphi^{\star}}:=\left\{f \in X: \quad\left\langle f, \varphi^{\star}\right\rangle>0\right\} .
$$

Note that for any $\varphi^{\star} \in X_{++}^{\star}$, we have the set inclusion $X_{++} \subseteq X_{++}^{\varphi^{\star}}$. In turn, this implies $\mathcal{A}_{++}\left(x_{0}\right) \subseteq \mathcal{A}_{++}^{\varphi^{\star}}\left(x_{0}\right)$ for any $x_{0} \in X$. Thus, we are relaxing the state constraint when passing from the former to the latter set of admissible controls. For future reference, we observe the following:

- both $\mathcal{A}_{++}\left(x_{0}\right)$ and $\mathcal{A}_{++}^{\varphi^{\star}}\left(x_{0}\right)$ are convex sets, due to linearity of the state equation;
- if $x_{0} \in X_{++}$, then $\mathcal{A}_{++}\left(x_{0}\right)$ and $\mathcal{A}_{++}^{\varphi^{\star}}\left(x_{0}\right)$ are nonempty, as the null control $c(\cdot) \equiv 0$ belongs to them;
- since $N$ is a positive operator, we have the monotonicity property:

$$
c_{1}(t) \leq_{X} c_{2}(t) \text { for almost all } t \geq 0 \Longrightarrow x^{x_{0}, c_{1}}(t) \geq_{X} x^{x_{0}, c_{2}}(t) \quad \forall t \geq 0
$$

and hence if $c_{1}(t) \leq_{X} c_{2}(t)$ for a.e. $t \geq 0$ and $c_{2} \in \mathcal{A}_{++}^{\varphi^{\star}}\left(x_{0}\right)$, then $c_{1} \in$ $\mathcal{A}_{++}^{\varphi^{\star}}\left(x_{0}\right)$.
We complete our setting by defining the functional to optimize. Let $u: D \times \mathbb{R}_{+} \rightarrow$ $\mathbb{R} \cup\{-\infty\}$ be a measurable function satisfying the following assumption, which will hold throughout the paper.

Assumption 2.4. The function $u: D \times \mathbb{R}_{+} \rightarrow \mathbb{R} \cup\{-\infty\}$ is such that $u(\theta, \cdot)$ is increasing and concave ${ }^{3}$ for all $\theta \in D$. Moreover, $u$ is bounded from either above or below. Without loss of generality we assume that either $u: D \times \mathbb{R}_{+} \rightarrow[-\infty, 0]$ or $u: D \times \mathbb{R}_{+} \rightarrow \mathbb{R}_{+}$.

[^2]Next, we consider the functionals

$$
\mathcal{U}(z):=\int_{D} u(\theta, z(\theta)) \mu(\mathrm{d} \theta), \quad z \in X_{+}
$$

and

$$
\begin{equation*}
\mathcal{J}(c):=\int_{0}^{\infty} \mathrm{e}^{-\rho t} \mathcal{U}(c(t)) \mathrm{d} t, \quad c \in \mathrm{~L}_{l o c}^{1}\left(\mathbb{R}_{+} ; X_{+}\right) \tag{2.6}
\end{equation*}
$$

where $\rho>0$ is a given discount factor. Notice that $\mathcal{U}$ and, consequently, $\mathcal{J}$ inherit concavity from $u$.

The optimal control problem we are interested in is

$$
\begin{equation*}
\text { maximize } \mathcal{J}(c) \text { over the set } \mathcal{A}_{++}\left(x_{0}\right), \quad x_{0} \in X_{++}, \tag{P}
\end{equation*}
$$

whose value function is

$$
V\left(x_{0}\right):=\sup _{c \in \mathcal{A}_{++}\left(x_{0}\right)} \mathcal{J}(c), \quad x_{0} \in X_{++}
$$

As anticipated, we relax the state constraint imposed by (2.4) introducing the family, parameterized by $\varphi^{\star} \in X_{++}^{\star}$, of auxiliary optimal control problems
$\left(P^{\varphi^{\star}}\right) \quad$ maximize $\mathcal{J}(c)$ over the set $\mathcal{A}_{++}^{\varphi^{\star}}\left(x_{0}\right), \quad x_{0} \in X_{++}$,
whose value function is

$$
V^{\varphi^{\star}}\left(x_{0}\right):=\sup _{c \in \mathcal{A}_{+++}^{\varphi^{\star}}\left(x_{0}\right)} \mathcal{J}(c), \quad x_{0} \in X_{++} .
$$

Remark 2.5. Since the state constraint imposed by (2.4) is stricter than the one imposed by $(2.5)$, we see that:
(i) $V^{\varphi^{\star}}\left(x_{0}\right) \geq V\left(x_{0}\right)$ for every $x_{0} \in X_{++}$;
(ii) if $\hat{c} \in \mathcal{A}_{++}^{\varphi^{\star}}\left(x_{0}\right)$ is optimal for $\left(P^{\varphi^{\star}}\right)$ and belongs to $\mathcal{A}_{++}\left(x_{0}\right)$, then it is optimal for $(P)$.
The reason for relaxing the state constraint and considering the new family of optimal control problems is that, as we will show in section 3, problem ( $P^{\varphi^{\star}}$ ) admits an explicit solution and an optimal feedback control. If it is possible to prove that this feedback control belongs to $\mathcal{A}_{++}\left(x_{0}\right)$, then the sufficient condition of Remark 2.5(ii) holds, and, consequently, this feedback control will be optimal for problem ( $P$ ). In the economic growth model studied in section 5 we will show that this is the case (provided that suitable assumptions hold), thanks to the stability result provided by Theorem 3.11.

It is worth noting that we cannot say ex ante that $V^{\varphi^{\star}}$ is finite. Even in simple one-dimensional cases (see, e.g., [30]), it may be always $+\infty$ or always $-\infty$. Sufficient conditions for finiteness will be provided later (see point (ii) of Assumption 3.2 and point (iii) of Remark 3.3).

We conclude this section by introducing the Hamilton-Jacobi-Bellman equation (HJB for short) associated to the optimal control problem ( $P^{\varphi^{\star}}$ ) indexed by $\varphi^{\star} \in$ $X_{++}^{\star}$,

$$
\begin{equation*}
\rho v(x)=\langle L x, \nabla v(x)\rangle+\mathcal{H}(\nabla v(x)), \quad x \in X_{++}^{\varphi^{\star}}, \tag{2.7}
\end{equation*}
$$

where the Hamiltonian function $\mathcal{H}$ is defined by

$$
\mathcal{H}\left(q^{\star}\right):=\sup _{z \in X_{+}} \mathcal{H}_{C V}\left(q^{\star} ; z\right), \quad q^{\star} \in X^{\star}
$$

and $\mathcal{H}_{C V}$ is the Current Value Hamiltonian function

$$
\mathcal{H}_{C V}\left(q^{\star} ; z\right):=\mathcal{U}(z)-\left\langle N z, q^{\star}\right\rangle, \quad z \in X_{+}, q^{\star} \in X^{\star}
$$

Notice that, without further assumptions, it may happen that $\mathcal{H}\left(q^{\star}\right)=+\infty$, for some $q^{\star} \in X^{\star}$. Finiteness of $\mathcal{H}\left(q^{\star}\right)$ is guaranteed at least in the following cases:

- If $\mathcal{U}$ is bounded above and $q^{\star} \in X_{+}^{\star}$;
- if $\mathcal{U}$ is bounded below and there exist constants $C>0, \alpha \in(0,1)$, and $k>0$, such that $\mathcal{U}(z) \leq C|z|_{X}^{\alpha}$ and $\left\langle N z, q^{\star}\right\rangle \geq k|z|_{X}, q^{\star} \in X^{\star}$.
In this paper we will consider classical solutions to HJB equation (2.7), according to the following definition.

Definition 2.6. A function $v \in \mathcal{C}^{1}\left(X_{++}^{\varphi^{\star}} ; \mathbb{R}\right)$ such that $\nabla v \in \mathcal{C}\left(X_{++}^{\varphi^{\star}} ; D\left(L^{\star}\right)\right)$ is called a classical solution to (2.7) on $X_{++}^{\varphi^{\star}}$ if

$$
\rho v(x)=\left\langle x, L^{\star} \nabla v(x)\right\rangle+\mathcal{H}(\nabla v(x)) \quad \forall x \in X_{++}^{\varphi^{\star}} .
$$

We point out that this notion of solution for an HJB equation in infinite dimension is very demanding, because of the required regularity of the solution itself. Nonetheless, we will see in section 3.2 that, under suitable assumptions, it is possible to find explicit solutions that verify the definition above. In more general cases, one can resort to other (weaker) notions, such as viscosity solutions (see section 6 for comments and references on this point).

Remark 2.7. Note that, if $v$ is a classical solution to HJB equation (2.7), then $\mathcal{H}(\nabla v(x))$ must be finite for every $x \in X_{++}^{\varphi^{\star}}$.
3. Verification theorem, explicit solutions, and stability. In this section we focus on the family of optimal control problems $\left(P^{\varphi^{\star}}\right)$. First, for each fixed $\varphi^{\star} \in X_{++}^{\star}$, we will provide a verification theorem in the general setting presented in section 2. Then, specializing our setting and suitably choosing $\varphi^{\star} \in X_{++}^{\star}$, we will provide an explicit solution to problem $\left(P^{\varphi^{\star}}\right)$ and give a stability result for this solution. These results will be crucial to studying our motivating economic application in section 5 .
3.1. Verification theorem. Typically, to prove a verification theorem for infinite horizon problems, a condition on the solution $v$ computed on the admissible trajectories when $t \rightarrow+\infty$ is needed. Given $x_{0} \in X$ and $\varphi^{\star} \in X_{++}^{\star}$, the (relaxed on integer numbers) condition that we shall use is

$$
\begin{equation*}
\lim _{k \in \mathbb{N}, k \rightarrow+\infty} \mathrm{e}^{-\rho k} v\left(x^{x_{0}, c}(k)\right)=0 \quad \forall c \in \mathcal{A}_{++}^{\varphi^{\star}}\left(x_{0}\right) \text { such that } \mathcal{J}(c)>-\infty \tag{3.1}
\end{equation*}
$$

ThEOREM 3.1 (verification). Let $\varphi^{\star} \in X_{++}^{\star}$ and $x_{0} \in X_{++}^{\varphi^{\star}}$. Let $v$ be a classical solution to (2.7), and assume that (3.1) holds. Then the following hold:
(i) $v\left(x_{0}\right) \geq V^{\varphi^{\star}}\left(x_{0}\right)$;
(ii) if, moreover, there exists $\hat{c} \in \mathcal{A}_{++}^{\varphi^{\star}}\left(x_{0}\right)$ such that, for a.e. $s \geq 0$, (3.2)

$$
\mathcal{H}\left(\nabla v\left(x^{x_{0}, \hat{c}}(s)\right)\right)=\mathcal{H}_{C V}\left(\nabla v\left(x^{x_{0}, \hat{c}}(s)\right) ; \hat{c}(s)\right) \Longleftrightarrow N^{\star} \nabla v\left(x^{x_{0}, \hat{c}}(s)\right) \in D^{+} \mathcal{U}(\hat{c}(s)),
$$

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where $D^{+} \mathcal{U}$ denotes the superdifferential of $\mathcal{U}$, then $v\left(x_{0}\right)=V^{\varphi^{\star}}\left(x_{0}\right)$ and $\hat{c}$ is optimal for $\left(P^{\varphi^{\star}}\right)$ starting at $x_{0}$, i.e., $\mathcal{J}(\hat{c})=V^{\varphi^{\star}}\left(x_{0}\right)$.
Proof. (i) Let $c \in \mathcal{A}_{++}^{\varphi^{\star}}\left(x_{0}\right)$ be such that $\mathcal{J}(c)>-\infty$. By the chain rule in infinite dimension (see [35]), we have, for every $t \geq 0$,

$$
\begin{aligned}
& \frac{\mathrm{d}}{\mathrm{~d} t}\left[\mathrm{e}^{-\rho t} v\left(x^{x_{0}, c}(t)\right)\right] \\
& =\mathrm{e}^{-\rho t}\left(-\rho v\left(x^{x_{0}, c}(t)\right)+\left\langle x^{x_{0}, c}(t), L^{\star} \nabla v\left(x^{x_{0}, c}(t)\right)\right\rangle-\left\langle N c(t), \nabla v\left(x^{x_{0}, c}(t)\right)\right\rangle\right) .
\end{aligned}
$$

Now we add and subtract $\mathrm{e}^{-\rho t} \mathcal{U}(c(t))$ to the right-hand side, use the fact that $v$ solves HJB, and integrate over $[0, t]$. We get, for every $t \geq 0$,

$$
\begin{aligned}
& \mathrm{e}^{-\rho t} v\left(x^{x_{0}, c}(t)\right)+\int_{0}^{t} \mathrm{e}^{-\rho s} \mathcal{U}(c(s)) \mathrm{d} s \\
& =v\left(x_{0}\right)+\int_{0}^{t} \mathrm{e}^{-\rho s}\left(-\mathcal{H}\left(\nabla v\left(x^{x_{0}, c}(s)\right)\right)+\mathcal{H}_{C V}\left(\nabla v\left(x^{x_{0}, c}(s) ; c(s)\right)\right)\right) \mathrm{d} s .
\end{aligned}
$$

Observe that, since $\mathcal{U}$ is concave (hence, sublinear from above) and $c(\cdot) \in \mathrm{L}_{l o c}^{1}\left(\mathbb{R}_{+} ; X_{+}\right)$, both sides of the above inequality are finite for every $t \geq 0$. Rearranging the terms and taking into account (3.2) and the definition of $\mathcal{H}$, we get, for every $t \geq 0$,

$$
\begin{equation*}
v\left(x_{0}\right) \geq \mathrm{e}^{-\rho t} v\left(x^{x_{0}, c}(t)\right)+\int_{0}^{t} \mathrm{e}^{-\rho s} \mathcal{U}(c(s)) \mathrm{d} s \tag{3.3}
\end{equation*}
$$

Since the sign of $\mathcal{U}$ is constant (cf. Assumption 2.4), we have that

$$
\begin{equation*}
\lim _{k \in \mathbb{N},}{ }_{k \rightarrow+\infty} \int_{0}^{k} \mathrm{e}^{-\rho s} \mathcal{U}(c(s)) \mathrm{d} s=\int_{0}^{\infty} \mathrm{e}^{-\rho s} \mathcal{U}(c(s)) \mathrm{d} s=\mathcal{J}(c) . \tag{3.4}
\end{equation*}
$$

Hence, passing to the limit in (3.3), as $k \rightarrow \infty, k \in \mathbb{N}$, and using (3.1), we deduce that

$$
v\left(x_{0}\right) \geq \mathcal{J}(c)
$$

Then, given the arbitrariness of $c \in \mathcal{A}_{++}^{\varphi^{\star}}\left(x_{0}\right)$ such that $\mathcal{J}(c)>-\infty$ and by definition of $V^{\varphi^{\star}}$, we immediately get the claim.
(ii) Notice that, by concavity of $\mathcal{U},(3.2)$ is equivalent to

$$
\begin{equation*}
\hat{c}(s) \in \underset{z \in X_{++}^{\varphi^{\star}}}{\operatorname{argmax}}\left\{\mathcal{U}(z)-\left\langle N z, \nabla v\left(x^{x_{0}, \hat{c}}(s)\right)\right\rangle\right\} \quad \text { for a.e. } s \geq 0, \tag{3.5}
\end{equation*}
$$

the usual closed loop condition for optimality. Hence, for $c=\hat{c}$ we have equality in (3.3), and therefore, passing to the limit, as $k \rightarrow \infty, k \in \mathbb{N}$, and using (3.1) and (3.4), we get the equality

$$
v\left(x_{0}\right)=\mathcal{J}(\hat{c}) .
$$

Since $\mathcal{J}(\hat{c}) \leq V^{\varphi^{\star}}\left(x_{0}\right)$, and since the reverse inequality holds by part (i) of the theorem, the claim follows.
3.2. Explicit solutions to HJB equation and optimal feedback control. In this section we are going to provide an explicit solution to the HJB equation (2.7) for problem $\left(P^{\varphi^{\star}}\right)$, for a specific choice of $\varphi^{\star} \in X_{++}^{\star}$. To this end, we need to specialize the setting of section 2 by introducing the following assumption that will be in force throughout this section.

Assumption 3.2. The set $X$ is either the space $\mathrm{L}^{p}(D, \mu), p \in[1,+\infty)$, where $\mu$ is a $\sigma$-finite measure, or the space $\mathcal{C}(D)$ of real-valued continuous functions on a compact metric space $D$, equipped with the sup-norm and with a Borel measure $\mu$.

Moreover, the linear operator $N$ in (2.1) and the function $u$ of Assumption 2.4 are explicitly given by

$$
\begin{align*}
{[N z](\theta) } & =\eta(\theta) z(\theta), \quad z \in X, \theta \in D  \tag{3.6}\\
u(\theta, \xi) & =\frac{\xi^{1-\gamma}}{1-\gamma} f(\theta), \quad \theta \in D, \xi \in \mathbb{R}_{+} \tag{3.7}
\end{align*}
$$

where $\eta, f: D \rightarrow(0,+\infty)$ are measurable functions and $\gamma \in(0,1) \cup(1,+\infty)$ is a fixed parameter.

Furthermore, the following hold:
(i) There exists an eigenvector $b_{0}^{\star} \in X_{++}^{\star}$ for $L^{\star}: D\left(L^{\star}\right) \subseteq X^{\star} \rightarrow X^{\star}$ with eigenvalue $\lambda_{0}^{\star} \in \mathbb{R}$.
(ii) $\rho>\lambda_{0}^{\star}(1-\gamma)$, where $\rho$ is the discount factor appearing in (2.6).
(iii) If $\gamma>1$, then $\frac{\left(b_{0}^{\star}\right)^{1-\gamma}}{f} \in L^{\infty}\left(D, \mu ; \mathbb{R}_{+}\right)$.
(iv) If $X=\mathcal{C}(D)$, then the strictly positive measure $b_{0}^{\star}$ is absolutely continuous with respect to $\mu$, with density still denoted by $b_{0}^{\star}$. In addition, $b_{0}^{\star}, \eta, f \in$ $\mathcal{C}(D ;(0,+\infty))$.
(v) If $X=\mathrm{L}^{p}(D, \mu)$, then $\eta, f \in \mathrm{~L}^{\infty}(D, \mu ;(0,+\infty))$ and, moreover, (3.8)

$$
\int_{D} f(\theta)^{\frac{1}{\gamma}}\left(\eta(\theta) b_{0}^{\star}(\theta)\right)^{\frac{\gamma-1}{\gamma}} \mu(\mathrm{~d} \theta)<\infty, \quad \int_{D}\left(\frac{f(\theta)}{\eta(\theta) b_{0}^{\star}(\theta)}\right)^{p / \gamma} \mu(\mathrm{d} \theta)<\infty
$$

Remark 3.3. Let us comment as follows on the specific setting described above:
(i) It is crucial to assume the existence of the strictly positive eigenvector $b_{0}^{\star} \in$ $X_{++}^{\star}$, as required by point (i) of Assumption 3.2. Indeed, this enables us to find an explicit solution to HJB (2.7) associated to the auxiliary problem ( $\left.P^{b_{0}^{\star}}\right)$ in the half-space $X_{++}^{b_{0}^{\star}}$.
(ii) In the case $X=\mathcal{C}(D)$ with $D$ compact space, it is important that $b_{0}^{\star}$ is represented as a continuous function to guarantee well-posedness of the feedback operator $\Phi$ defined in (3.14), whence point (iv) of Assumption 3.2.
(iii) Point (ii) of Assumption 3.2 is needed to ensure finiteness of the solution of the HJB equation and hence of the value function.
(iv) Point (iii) of Assumption 3.2 is required to verify (3.1) in Lemma 3.8.
(v) Given the expression of $u$ and since $\eta$ and $f$, appearing in (3.6) and (3.7), respectively, are positive, point (ii) of Assumptions 2.2 and 2.4 are verified.
(vi) In the case $X=\mathrm{L}^{p}(D, \mu)$, the functions $\eta$ and $f$ do not need to be continuous but only essentially bounded. However, additional integrability conditions must hold: the first requirement of (3.8) allows us to make sense of the explicit solution of HJB; the second condition of (3.8) is used to make sense of the optimal feedback map. Both conditions are automatically verified when $X=\mathcal{C}(D)$ with $D$ compact, due to the compactness of $D$.
We are now ready to provide the explicit solution to (2.7) for the auxiliary problem $\left(P^{b_{0}^{\star}}\right)$.

Proposition 3.4. The function

$$
\begin{equation*}
v(x):=\alpha \frac{\left\langle x, b_{0}^{\star}\right\rangle^{1-\gamma}}{1-\gamma}, \quad x \in X_{++}^{b_{0}^{\star}} \tag{3.9}
\end{equation*}
$$

where

$$
\begin{equation*}
\alpha:=\gamma^{\gamma}\left(\frac{\int_{D} f(\theta)^{\frac{1}{\gamma}}\left(\eta(\theta) b_{0}^{\star}(\theta)\right)^{\frac{\gamma-1}{\gamma}} \mu(\mathrm{~d} \theta)}{\rho-\lambda_{0}^{\star}(1-\gamma)}\right)^{\gamma} \tag{3.10}
\end{equation*}
$$

is a classical solution to (2.7).
Proof. We point out, first, that the arguments provided here are the same regardless of the choice of $X$ in Assumption 3.2.

We start by noting that $v$ is differentiable, as a composition of a linear form with a power function, and that

$$
\nabla v(x)=\alpha\left\langle x, b_{0}^{\star}\right\rangle^{-\gamma} b_{0}^{\star}, \quad x \in X_{++}^{b_{0}^{\star}} .
$$

Due to Assumption 3.2(i), and since $\alpha>0$ by point (ii) and either point (iv) or point (v), $v$ enjoys the regularity required by Definition 2.6. Next, observe that by (3.7) and Assumption 3.2-either point (iv) or (v)—, the functional $\mathcal{U}$ is Fréchet differentiable, and

$$
\begin{equation*}
[\nabla \mathcal{U}(z)](\theta)=f(\theta) z(\theta)^{-\gamma}, \quad z \in X_{+}, \theta \in D \tag{3.11}
\end{equation*}
$$

The adjoint of the operator $N$, given in (3.6), is

$$
\begin{array}{ll}
N^{\star} q^{\star}(\theta)=\eta(\theta) q^{\star}(\theta), \quad \theta \in D, & \text { when } X=\mathrm{L}^{p}(D, \mu) \\
N^{\star} q^{\star}(\mathrm{d} \theta)=\eta(\theta) q^{\star}(\mathrm{d} \theta), \quad \theta \in D, & \text { when } X=\mathcal{C}(D)
\end{array}
$$

Therefore, taking into account point (iv) of Assumption 3.2, we get (regardless of the choice of $X$ )

$$
\begin{equation*}
N^{\star} \nabla v(x)=\alpha\left\langle x, b_{0}^{\star}\right\rangle^{-\gamma} \eta(\cdot) b_{0}^{\star}, \quad x \in X_{++}^{b_{0}^{\star}} . \tag{3.12}
\end{equation*}
$$

Putting together (3.11) and (3.12), we have

$$
\begin{array}{r}
\frac{\partial \mathcal{H}_{C V}}{\partial z}(\nabla v(x) ; z)=\nabla \mathcal{U}(z)-N^{\star} \nabla v(x)=f(\cdot) z^{-\gamma}-\alpha\left\langle x, b_{0}^{\star}\right\rangle^{-\gamma} \eta(\cdot) b_{0}^{\star}(\cdot) \\
x \in X_{++}^{b_{0}^{\star}}, \quad z \in X_{+}
\end{array}
$$

Since $z \mapsto \mathcal{H}_{C V}(\nabla v(x) ; z)$ is strictly concave, for each fixed $x \in X_{++}^{b_{0}^{\star}}$, the unique maximum point $\hat{z}(x)$ of this function is provided by

$$
\begin{equation*}
\hat{z}(x)=\underset{z \in X_{+}}{\operatorname{argmax}} \mathcal{H}_{C V}(\nabla v(x) ; z)=\left(\frac{f(\cdot)}{\alpha \eta(\cdot) b_{0}^{\star}(\cdot)}\right)^{\frac{1}{\gamma}}\left\langle x, b_{0}^{\star}\right\rangle, \quad x \in X_{++}^{b_{0}^{\star}}, \tag{3.13}
\end{equation*}
$$

whence

$$
\begin{aligned}
\mathcal{H}(\nabla v(x))=\mathcal{H}_{C V}(\nabla v(x) ; \hat{z}(x)) \\
\quad=\frac{\gamma\left\langle x, b_{0}^{\star}\right\rangle^{1-\gamma}}{1-\gamma} \int_{D} f(\theta)^{\frac{1}{\gamma}}\left(\alpha \eta(\theta) b_{0}^{\star}(\theta)\right)^{\frac{\gamma-1}{\gamma}} \mu(\mathrm{~d} \theta), \quad x \in X_{++}^{b_{0}^{\star}} .
\end{aligned}
$$

Plugging the latter expression into (2.7), we get the algebraic equation in $\alpha$,

$$
\frac{\rho}{1-\gamma} \alpha=\lambda_{0}^{\star} \alpha+\frac{\gamma}{(1-\gamma)} \alpha^{\frac{\gamma-1}{\gamma}} \int_{D} f(\theta)^{\frac{1}{\gamma}}\left(\eta(\theta) b_{0}^{\star}(\theta)\right)^{\frac{\gamma-1}{\gamma}} \mu(\mathrm{~d} \theta),
$$

which has a unique positive solution provided by (3.10).

Remark 3.5. It is easy to verify that the function $v$ does not change if we multiply $b_{0}^{\star}$ by any strictly positive constant. This is consistent with the fact that, if $k>0$ and $\varphi^{\star}=k b_{0}^{\star}$, then $X_{++}^{\varphi^{\star}}=X_{++}^{b_{0}^{\star}}$; hence, $\mathcal{A}_{++}^{\varphi^{\star}}\left(x_{0}\right)=\mathcal{A}_{++}^{b_{0}^{\star}}\left(x_{0}\right)$.

The following map provides the unique maximum point given in (3.13) as a function of the state:

$$
\begin{equation*}
[\Phi x](\theta)=\hat{z}(x)(\theta)=\left(\frac{f(\theta)}{\alpha \eta(\theta) b_{0}^{\star}(\theta)}\right)^{\frac{1}{\gamma}}\left\langle x, b_{0}^{\star}\right\rangle, \quad x \in X, \theta \in D . \tag{3.14}
\end{equation*}
$$

It is important to note that $\Phi: X \rightarrow X$ is a bounded and positive linear operator. This is the natural candidate for the optimal feedback map for problem $\left(P^{b_{0}^{\star}}\right)$, i.e., the map that will provide (as we will show) the optimal feedback control for this problem, given any initial condition $x_{0} \in X_{++}$. Therefore, we introduce the closed loop equation for problem $\left(P^{b_{0}^{\star}}\right)$ :

$$
\left\{\begin{array}{l}
x^{\prime}(t)=L x(t)-N \Phi x(t), \quad t \geq 0  \tag{3.15}\\
x(0)=x_{0} \in X
\end{array}\right.
$$

This linear equation admits a unique mild solution, denoted by $\hat{x}^{x_{0}}(\cdot)$, which is also a weak solution (see [5, pp. 129-131]). Notice that the operator $N \Phi: X \rightarrow X$, appearing in (3.15), is explicitly given by

$$
\begin{equation*}
[N \Phi x](\theta)=\eta(\theta)\left(\frac{f(\theta)}{\alpha \eta(\theta) b_{0}^{\star}(\theta)}\right)^{\frac{1}{\gamma}}\left\langle x, b_{0}^{\star}\right\rangle, \quad x \in X, \theta \in D . \tag{3.16}
\end{equation*}
$$

To study (3.15), let us define the operator

$$
\begin{equation*}
B:=L-N \Phi \tag{3.17}
\end{equation*}
$$

Since $B$ is a bounded perturbation of $L$, by [21, Chap. III, Thm. 1.3] we have that it is a closed operator with domain $D(B)=D(L)$ and that it generates a $C_{0}$-semigroup $\left\{\mathrm{e}^{t B}\right\}_{t \geq 0}$ on $X$. The next lemma will be useful in what follows.

Lemma 3.6. We have the following:
(i) The adjoint operator $(N \Phi)^{\star}: X^{\star} \rightarrow X^{\star}$ is explicitly given by
$(N \Phi)^{\star} q=\alpha^{-\frac{1}{\gamma}}\left(\int_{D} f(\theta)^{\frac{1}{\gamma}} \eta(\theta)^{1-\frac{1}{\gamma}} b_{0}^{\star}(\theta)^{-\frac{1}{\gamma}} q(\mathrm{~d} \theta)\right) b_{0}^{\star} \quad$ when $X=\mathcal{C}(D)$,
$(N \Phi)^{\star} q=\alpha^{-\frac{1}{\gamma}}\left(\int_{D} f(\theta)^{\frac{1}{\gamma}} \eta(\theta)^{1-\frac{1}{\gamma}} b_{0}^{\star}(\theta)^{-\frac{1}{\gamma}} q(\theta) \mu(\mathrm{d} \theta)\right) b_{0}^{\star} \quad$ when $X=\mathrm{L}^{p}(D, \mu)$.
(ii) Define

$$
\begin{equation*}
g:=\frac{\lambda_{0}^{\star}-\rho}{\gamma} \tag{3.20}
\end{equation*}
$$

Then $g$ is an eigenvalue of $B^{\star}$, and $b_{0}^{\star}$ is an eigenvector of $B^{\star}$ associated to $g$.

Proof. The proof of (i) follows from straightforward computations. For the proof of (ii) we observe that, using (i), (3.10), and (3.20), we get

$$
B^{\star} b_{0}^{\star}=L^{\star} b_{0}^{\star}-(N \Phi)^{\star} b_{0}^{\star}=\lambda_{0}^{\star} b_{0}^{\star}-\frac{\rho-\lambda_{0}^{\star}(1-\gamma)}{\gamma} b_{0}^{\star}=\left[\lambda_{0}^{\star}-\frac{\rho-\lambda_{0}^{\star}(1-\gamma)}{\gamma}\right] b_{0}^{\star}=g b_{0}^{\star} .
$$

Lemma 3.7. Let $x_{0} \in X_{++}^{b_{0}^{\star}}$, and let $\hat{x}^{x_{0}}$ be the unique mild (weak) solution to (3.15). Then,

$$
\begin{equation*}
\left\langle\hat{x}^{x_{0}}(t), b_{0}^{\star}\right\rangle=\left\langle x_{0}, b_{0}^{\star}\right\rangle \mathrm{e}^{g t} \quad \forall t \geq 0, \tag{3.21}
\end{equation*}
$$

where $g$ is as given by (3.20). Hence, in particular,

$$
\begin{equation*}
\hat{x}^{x_{0}}(t) \in X_{++}^{b_{0}^{\star}} \quad \forall t \geq 0 \tag{3.22}
\end{equation*}
$$

Proof. Testing the closed loop equation (3.15) against $b_{0}^{\star} \in D\left(L^{\star}\right)$, by Lemma 3.6 we get
$\frac{\mathrm{d}}{\mathrm{d} t}\left\langle\hat{x}^{x_{0}}(t), b_{0}^{\star}\right\rangle=\left\langle\hat{x}^{x_{0}}(t), L^{\star} b_{0}^{\star}\right\rangle-\left\langle N \Phi \hat{x}^{x_{0}}(t), b_{0}^{\star}\right\rangle=\left\langle\hat{x}^{x_{0}}(t),\left(L^{\star}-(N \Phi)^{*}\right) b_{0}^{\star}\right\rangle=g\left\langle\hat{x}^{x_{0}}(t), b_{0}^{\star}\right\rangle$.
The result follows by integrating this equation over $[0, t]$ and noting that $\left\langle x_{0}, b_{0}^{\star}\right\rangle>0$. $\square$
Lemma 3.8. Let $x_{0} \in X_{++}^{b_{0}^{\star}}$, and let $v$, defined in (3.9), be the solution to (2.7) associated to $\left(P^{b_{0}^{\star}}\right)$. Then,

$$
\begin{equation*}
\lim _{k \in \mathbb{N}, k \rightarrow+\infty} \mathrm{e}^{-\rho k} v\left(x^{x_{0}, c}(k)\right)=0 \quad \forall c \in \mathcal{A}_{++}^{b_{0}^{\star}}\left(x_{0}\right) \text { such that } \mathcal{J}(c)>-\infty \tag{3.23}
\end{equation*}
$$

Proof. Case 1: $\gamma \in(0,1)$. In this case $v$ is nonnegative, and this allows us to prove a stronger result, namely,

$$
\begin{equation*}
\lim _{t \rightarrow+\infty} \mathrm{e}^{-\rho t} v\left(x^{x_{0}, c}(t)\right)=0 \quad \forall c \in \mathcal{A}_{++}^{b_{0}^{\star}}\left(x_{0}\right) . \tag{3.24}
\end{equation*}
$$

Recalling that $x^{x_{0}, c}$ is a weak solution to the state equation (2.1), we can rewrite (2.3) with $\varphi^{\star}=b_{0}^{\star}$ and obtain

$$
\frac{\mathrm{d}}{\mathrm{~d} t}\left\langle x^{x_{0}, c}(t), b_{0}^{\star}\right\rangle=\lambda_{0}^{\star}\left\langle x^{x_{0}, c}(t), b_{0}^{\star}\right\rangle-\left\langle N c(t), b_{0}^{\star}\right\rangle \quad \forall c \in \mathcal{A}_{++}^{b_{0}^{\star}}\left(x_{0}\right), \quad \forall t \geq 0
$$

i.e.,

$$
\left\langle x^{x_{0}, c}(t), b_{0}^{\star}\right\rangle=\left\langle x_{0}, b_{0}^{\star}\right\rangle \mathrm{e}^{\lambda_{0}^{\star} t}-\int_{0}^{t} \mathrm{e}^{\lambda_{0}^{\star}(t-s)}\left\langle N c(s), b_{0}^{\star}\right\rangle \mathrm{d} s \quad \forall c \in \mathcal{A}_{++}^{b_{0}^{\star}}\left(x_{0}\right), \quad \forall t \geq 0
$$

Since $N$ is a positive operator and $b_{0}^{\star} \in X_{++}^{\star}$ (cf. point (ii) of Assumption 2.2 and point (i) of Assumption 3.2), we have

$$
0 \leq\left\langle x^{x_{0}, c}(t), b_{0}^{\star}\right\rangle \leq\left\langle x_{0}, b_{0}^{\star}\right\rangle \mathrm{e}^{\lambda_{0}^{\star} t} \quad \forall c \in \mathcal{A}_{++}^{b_{0}^{\star}}\left(x_{0}\right), \quad \forall t \geq 0 .
$$

Hence, the following inequality holds for all $t \geq 0$ :

$$
\begin{aligned}
0 & \leq \mathrm{e}^{-\rho t} v\left(x^{x_{0}, c}(t)\right)=\alpha \mathrm{e}^{-\rho t} \frac{\left\langle x^{x_{0}, c}(t), b_{0}^{\star}\right\rangle^{1-\gamma}}{1-\gamma} \\
& \leq \frac{\alpha}{1-\gamma}\left\langle x_{0}, b_{0}^{\star}\right\rangle \mathrm{e}^{-\left(\rho-\lambda_{0}^{\star}(1-\gamma)\right) t} \quad \forall c \in \mathcal{A}_{++}^{b_{0}^{\star}}\left(x_{0}\right) .
\end{aligned}
$$

Noting that $\left\langle x_{0}, b_{0}^{\star}\right\rangle>0$ and using point (ii) of Assumption 3.2, we get (3.24).
Case 2: $\gamma>1$. In this case $v$ is nonpositive. Hence, for every $c \in \mathcal{A}_{++}^{b_{0}^{*}}\left(x_{0}\right)$, we have that

$$
\liminf _{k \in \mathbb{N}, k \rightarrow+\infty} \mathrm{e}^{-\rho k} v\left(x^{x_{0}, c}(k)\right) \leq \limsup _{k \in \mathbb{N}, k \rightarrow+\infty} \mathrm{e}^{-\rho k} v\left(x^{x_{0}, c}(k)\right) \leq 0 .
$$

Let $k \in \mathbb{N}$. Then, by (2.3) with $\varphi^{\star}=b_{0}^{\star}$ and since $c \in \mathcal{A}_{++}^{b_{0}^{\star}}\left(x_{0}\right)$,

$$
\begin{aligned}
0 & \leq\left\langle x^{x_{0}, c}(k+1), b_{0}^{\star}\right\rangle=\left\langle x^{x_{0}, c}(k), b_{0}^{\star}\right\rangle+\int_{k}^{k+1} \lambda_{0}^{\star}\left\langle x^{x_{0}, c}(s), b_{0}^{\star}\right\rangle \mathrm{d} s-\int_{k}^{k+1}\left\langle c(s), b_{0}^{\star}\right\rangle \mathrm{d} s \\
& =\left\langle x^{x_{0}, c}(k), b_{0}^{\star}\right\rangle \mathrm{e}^{\lambda_{0}^{\star}}-\int_{k}^{k+1} \mathrm{e}^{\lambda_{0}^{\star}(k+1-s)}\left\langle c(s), b_{0}^{\star}\right\rangle \mathrm{d} s .
\end{aligned}
$$

It follows that

$$
\mathrm{e}^{-\left|\lambda_{0}^{\star}\right|} \int_{k}^{k+1}\left\langle c(s), b_{0}^{\star}\right\rangle \mathrm{d} s \leq\left\langle x^{x_{0}, c}(k), b_{0}^{\star}\right\rangle \mathrm{e}^{\lambda_{0}^{\star}} ;
$$

therefore,

$$
\int_{k}^{k+1}\left\langle c(s), b_{0}^{\star}\right\rangle \mathrm{d} s \leq\left\langle x^{x_{0}, c}(k), b_{0}^{\star}\right\rangle \mathrm{e}^{2\left|\lambda_{0}^{\star}\right|} .
$$

Hence, by Jensen's inequality, monotonicity, concavity, and nonpositivity of the map $\xi \mapsto \frac{\xi^{1-\gamma}}{1-\gamma}$,

$$
\begin{aligned}
& \int_{k}^{k+1} \frac{\left\langle c(s), b_{0}^{\star}\right\rangle^{1-\gamma}}{1-\gamma} \mathrm{d} s \\
& \quad \leq \frac{1}{1-\gamma}\left(\int_{k}^{k+1}\left\langle c(s), b_{0}^{\star}\right\rangle \mathrm{d} s\right)^{1-\gamma} \leq \frac{\left\langle x^{x_{0}, c}(k), b_{0}^{\star}\right\rangle^{1-\gamma}}{1-\gamma} \mathrm{e}^{2\left|\lambda_{0}^{\star}\right|(1-\gamma)} \leq 0 .
\end{aligned}
$$

Therefore, multiplying by $\mathrm{e}^{-\rho(k+1)}$ and recalling that $\rho>0$,

$$
\begin{equation*}
c_{k}:=\int_{k}^{k+1} \mathrm{e}^{-\rho(k+1)} \frac{\left\langle c(s), b_{0}^{\star}\right\rangle^{1-\gamma}}{1-\gamma} \mathrm{d} s \leq \mathrm{e}^{-\rho k} \frac{\left\langle x^{x_{0}, c}(k), b_{0}^{\star}\right\rangle^{1-\gamma}}{1-\gamma} \mathrm{e}^{2\left|\lambda_{0}^{\star}\right|(1-\gamma)-\rho} \leq 0 . \tag{3.25}
\end{equation*}
$$

By point (iii) of Assumption 3.2, it follows that

$$
\begin{aligned}
-\infty & <\left|\frac{\left(b_{0}^{\star}\right)^{1-\gamma}}{f}\right|_{\infty} \mathcal{J}(c)=\left|\frac{\left(b_{0}^{\star}\right)^{1-\gamma}}{f}\right|_{\infty} \int_{0}^{\infty} \mathrm{e}^{-\rho s}\left(\int_{D} \frac{c(s, \theta)^{1-\gamma}}{1-\gamma} f(\theta) \mathrm{d} \theta\right) \mathrm{d} s \\
& \leq \int_{0}^{\infty} \mathrm{e}^{-\rho s}\left(\int_{D} \frac{c(s, \theta)^{1-\gamma}}{1-\gamma} \frac{b_{0}^{\star}(\theta)^{1-\gamma}}{f(\theta)} f(\theta) \mathrm{d} \theta\right) \mathrm{d} s \\
& =\int_{0}^{\infty} \mathrm{e}^{-\rho s}\left(\int_{D} \frac{c(s, \theta)^{1-\gamma} b_{0}^{\star}(\theta)^{1-\gamma}}{1-\gamma} \mathrm{d} \theta\right) \mathrm{d} s \\
& =\int_{0}^{\infty} \mathrm{e}^{-\rho s} \frac{\left\langle c(s), b_{0}^{\star}\right\rangle^{1-\gamma}}{1-\gamma} \mathrm{d} s=\sum_{k=0}^{\infty} \int_{k}^{k+1} \mathrm{e}^{-\rho s} \frac{\left\langle c(s), b_{0}^{\star}\right\rangle^{1-\gamma}}{1-\gamma} \mathrm{d} s \leq \sum_{k=0}^{\infty} c_{k} \leq 0 .
\end{aligned}
$$

Hence, $c_{k} \rightarrow 0$. Combining this fact with (3.25), we get

$$
0 \geq \liminf _{k \in \mathbb{N}, k \rightarrow \infty} \mathrm{e}^{-\rho k} v\left(x^{x_{0}, c}(k)\right)=\liminf _{k \in \mathbb{N}, k \rightarrow \infty} \mathrm{e}^{-\rho k} \frac{\left\langle x^{x_{0}, c}(k), b_{0}^{\star}\right\rangle^{1-\gamma}}{1-\gamma}=0 .
$$

Theorem 3.9. Let $x_{0} \in X_{++}$, and define the control

$$
\hat{c}(t, \theta):=\left(\Phi \hat{x}^{x_{0}}(t)\right)(\theta)=\left(\frac{f(\theta)}{\alpha \eta(\theta) b_{0}^{\star}(\theta)}\right)^{\frac{1}{\gamma}}\left\langle x_{0}, b_{0}^{\star}\right\rangle \mathrm{e}^{g t}, \quad t \geq 0, \theta \in D
$$

where $\hat{x}^{x_{0}}$ is the unique mild solution to the closed loop equation (3.15). Then, $\hat{c} \in$ $\mathcal{A}_{++}^{b_{0}^{\star}}\left(x_{0}\right)$, and it is optimal for problem $\left(P^{b_{0}^{\star}}\right)$ starting at $x_{0}$. Moreover, the value function $V^{b_{0}^{\star}}$ of problem $\left(P^{b_{0}^{\star}}\right)$ satisfies $V^{b_{0}^{\star}}\left(x_{0}\right)=v\left(x_{0}\right)$, where $v$ is the function defined in (3.9).

Proof. Recalling (3.21), $x^{x_{0}, \hat{c}}=\hat{x}^{x_{0}}$ by construction. Hence, by (3.22), $\hat{c} \in$ $\mathcal{A}_{++}^{b_{0}^{\star}}\left(x_{0}\right)$. Moreover, again by construction, $\hat{c}$ verifies the optimality condition (3.5). So, thanks to Lemma 3.8, the assumptions of Theorem 3.1 are verified, whence we obtain $V^{b_{0}^{\star}}\left(x_{0}\right)=v\left(x_{0}\right)$.
3.3. Steady states and stability of solutions. Let us begin by stating an important result concerning the operator $B$, defined in (3.17), that will motivate what follows.

Lemma 3.10 (see [21, Chap. V, Cor. 3.2]). If B generates an eventually compact $C_{0}$-semigroup on $X$, then the following properties hold:

1. The spectrum $\sigma(B)$ is either empty, finite, or countable and consists of poles of the resolvent of finite algebraic multiplicity only.
2. The set $\{\nu \in \sigma(B): \mathfrak{R e} \nu \geq r\}$ is finite, for any $r \in \mathbb{R}$.

As a consequence of Lemma 3.10, if $\left\{\mathrm{e}^{t B}\right\}_{t \geq 0}$ is eventually compact, and if $\sigma(B) \neq$ $\emptyset$, we have that $\sigma(B)=\left\{\nu_{0}, \nu_{1}, \ldots\right\}$, where $\mathfrak{R e} \nu_{k} \geq \mathfrak{R e} \nu_{k+1}$, for any $k \in \mathbb{N}$ and $\lim _{k \rightarrow \infty} \mathfrak{R e} \nu_{k}=-\infty$, provided that $\sigma(B)$ is infinite. Moreover, by [21, Chap. IV, Prop. 2.18(i)] we have that $\sigma(B)=\sigma\left(B^{\star}\right)$, and hence Lemma 3.6 entails that $g \in \sigma(B)$.

Note that $g$ is the exponential growth rate of the map $t \mapsto\left\langle\hat{x}^{x_{0}}(t), b_{0}^{\star}\right\rangle$ (cf. (3.21)), where $\hat{x}^{x_{0}}$ is the mild solution to the closed loop equation (3.15). This suggests studying convergence and stability of the detrended optimal paths $\left\{\mathrm{e}^{-g t} \hat{x}^{x_{0}}(t)\right\}_{t \geq 0}$ exploiting the spectral properties of operator $B$. In particular, we aim at using a Perron-Frobenius-type argument. To do so, we need that $g$ is the highest eigenvalue of $B$, i.e., that $g=\nu_{0}$. However, we cannot say, ex ante, for which $k \in \mathbb{N}$ we have $g=\nu_{k}$. In section 5 , we will see that, under appropriate assumptions, $g=\nu_{0}$.

Motivated by this argument, we are going to provide a stability property of the detrended optimal paths $\hat{x}_{\nu_{0}}^{x_{0}}$ for problem $\left(P^{b_{0}^{\star}}\right)$, where

$$
\hat{x}_{\nu_{0}}^{x_{0}}(t):=\mathrm{e}^{-\nu_{0} t} \hat{x}^{x_{0}}(t)=\mathrm{e}^{t\left(B-\nu_{0}\right)} x_{0}, \quad t \geq 0 .
$$

It is worth noting that $\hat{x}_{\nu_{0}}^{x_{0}}$ is the unique (mild) solution to the detrended closed loop equation

$$
\left\{\begin{array}{l}
x^{\prime}(t)=\left(L-\nu_{0}\right) x(t)-N \Phi x(t), \quad t \geq 0  \tag{3.26}\\
x(0)=x_{0} \in X
\end{array}\right.
$$

An element $\bar{x} \in X$ is called a steady state for the dynamical system (3.26) if

$$
\begin{equation*}
\hat{x}_{\nu_{0}}^{\bar{x}}(t)=\mathrm{e}^{t\left(B-\nu_{0}\right)} \bar{x}=\bar{x} \quad \forall t \geq 0, \tag{3.27}
\end{equation*}
$$

From [21, Chap. IV, Cor. 3.8] we deduce that

$$
\begin{equation*}
\bar{x} \in X \text { is a steady state for }(3.26) \Longleftrightarrow \bar{x} \in \operatorname{ker}\left(B-\nu_{0}\right) \tag{3.28}
\end{equation*}
$$

The next theorem provides a stability result under appropriate assumptions.

Theorem 3.11. Let $B$ generate an eventually compact $C_{0}$-semigroup on $X$, and assume that $\nu_{0}$ is a first order pole of the resolvent, that is, its algebraic multiplicity is 1 . Then, for any $x_{0} \in X$, there exist $M \geq 1, \epsilon>0$ such that

$$
\left|\hat{x}_{\nu_{0}}^{x_{0}}(t)-P x_{0}\right|_{X} \leq M \mathrm{e}^{-\epsilon t}\left|(1-P) x_{0}\right|_{X} \quad \forall t \geq 0,
$$

where $P$ is the spectral projection corresponding to the spectral set $\left\{\nu_{0}\right\} \subset \sigma(B)$. Moreover, Range $(P)=\operatorname{ker}\left(B-\nu_{0}\right)$, i.e., $P$ is a projection on $\operatorname{ker}\left(B-\nu_{0}\right)$, and $P x_{0}$ is a steady state for (3.26).

Proof. By [21, Chap. V, Cor. 3.3], there exist constants $\epsilon>0$ and $M \geq 1$ such that

$$
\begin{equation*}
\left|\mathrm{e}^{t\left(B-\nu_{0}\right)}-P\right|_{\mathcal{L}(X)} \leq M \mathrm{e}^{-\epsilon t} \quad \forall t \geq 0, \tag{3.29}
\end{equation*}
$$

where $|\cdot|_{\mathcal{L}(X)}$ denotes the operator norm on the space of linear continuous operators $\mathcal{L}(X)$ and where $P \in \mathcal{L}(X)$ is the residue of the resolvent at $\nu_{0}$.

From standard theory (see, e.g., [2, Chap. A-III, sect. 3] or [21, Chap. IV, sect. 1]), $P$ is the spectral projection corresponding to the spectral set $\left\{\nu_{0}\right\} \subset \sigma(B)$. Moreover, being $\nu_{0}$ a simple pole, we have that Range $(P)=\operatorname{ker}\left(B-\nu_{0}\right)$ (see, e.g., [2, Chap. A-III, sect. 3, p. 73] or [21, Chap. IV, sect. 1, p. 247]). By (3.28), this implies that $P x_{0}$ is a steady state for (3.26). Finally, using (3.28)-(3.29) and taking into account that $P^{2} x_{0}=P x_{0}$, we have that, for all $t \geq 0$,

$$
\begin{aligned}
\left|\hat{\nu}_{\nu_{0}}^{x_{0}}(t)-P x_{0}\right|_{X} & =\left|\mathrm{e}^{t\left(B-\nu_{0}\right)} x_{0}-\mathrm{e}^{t\left(B-\nu_{0}\right)} P x_{0}\right|_{X} \\
& =\left|\mathrm{e}^{t\left(B-\nu_{0}\right)}\left(x_{0}-P x_{0}\right)-P\left(x_{0}-P x_{0}\right)\right| X \leq M \mathrm{e}^{-\epsilon t}\left|(1-P) x_{0}\right|_{X} .
\end{aligned}
$$

Remark 3.12. For the reader's convenience, we recall that under the assumptions of Theorem 3.11 the spectral projection $P$ corresponding to the spectral set $\left\{\nu_{0}\right\} \subset$ $\sigma(B)$ is (see [21, Chap. IV, eq. (1.12)])

$$
\begin{equation*}
P:=-\frac{1}{2 \pi \mathrm{i}} \int_{\gamma}(B-\mu)^{-1} \mathrm{~d} \mu, \tag{3.30}
\end{equation*}
$$

where $\gamma$ can be taken as the simple curve given by the positively oriented boundary of a disk centered at $\nu_{0}$, with radius small enough so that it does not enclose any other points of $\sigma(B)$.
4. Existence of a strictly positive eigenvector of $L^{\star}$. In this section we state some results guaranteeing that the particularly relevant point (i) of Assumption 3.2 is satisfied. This fact is tightly linked to some properties of operator $L$, appearing in the state equation (2.1), that will be discussed. In the following, recall that we are working under Assumption 2.2.

We start our discussion with the general case where $X$ is any Banach lattice. Let $\sigma(L)$ be the spectrum of the operator $L$, and define the spectral bound of $L$,

$$
s_{L}:=\sup \{\mathfrak{R e} \lambda: \lambda \in \sigma(L)\},
$$

with the convention $\sup \emptyset=-\infty$. Due to Assumption 2.2, $s_{L}=\sup \{\lambda \in \mathbb{R}: \lambda \in \sigma(L)\}$ (see [16, Thms. 7.4 and 8.7]) and, if $\sigma(L) \neq \emptyset$, then $s_{L} \in \sigma(L)$. The peripheral or boundary spectrum of $L$ is the subset of $\mathbb{C}$ given by

$$
\sigma_{b}(L):=\left\{\lambda \in \sigma(L): \mathfrak{R e} \lambda=s_{L}\right\} .
$$

We introduce the following definitions (the first one can be generalized, as in $[2$, Chap. C-III, Def. 3.1] or [16, Prop. 7.6]).

Definition 4.1. The positive $C_{0}$-semigroup $\left\{\mathrm{e}^{t L}\right\}_{t \geq 0}$ is said to be irreducible if for all $f \in X_{+} \backslash\{0\}$ and $\varphi^{\star} \in X_{+}^{\star} \backslash\{0\}$, there exists $t_{0} \geq 0$ such that $\left\langle\mathrm{e}^{t_{0} L} f, \varphi^{\star}\right\rangle>0$.

Definition 4.2. A point $f \in X_{+}$is said to be a quasi-interior point (of $X_{+}$) if $\left\langle f, \varphi^{\star}\right\rangle>0$ for each $\varphi^{\star} \in X_{+}^{\star} \backslash\{0\}$.

Remark 4.3. If int $X_{+} \neq \emptyset$, the concepts of interior and quasi-interior points coincide; this is the case, for instance, when $X$ is the space of real-valued continuous functions on a compact set. Instead, in the case $X=\mathrm{L}^{p}(D, \mu)$, with $\mu$ a $\sigma$-finite measure and $p \in[1,+\infty)$, we have that $\operatorname{int} X_{+}=\emptyset$, while the set of quasi-interior points of $X_{+}$coincides with the set of $\mu$-a.e. strictly positive functions (see [2, p. 238]).

ThEOREM 4.4. Let $\left\{\mathrm{e}^{t L}\right\}_{t \geq 0}$ be irreducible, and assume that $s_{L}>-\infty$ is a pole of the resolvent of $L$. Then we have the following:
(i) There exists $\varphi_{0}^{\star} \in D\left(L^{\star}\right) \cap X_{++}^{\star}$ such that

$$
\begin{equation*}
L^{\star} \varphi_{0}^{\star}=s_{L} \varphi_{0}^{\star} \tag{4.1}
\end{equation*}
$$

(ii) $s_{L} \in \sigma(L)$ has algebraic and geometric multiplicity 1 , and there exists a unique (up to a multiplicative constant) quasi-interior point $f_{0} \in D(L) \cap X_{+}$ such that

$$
\begin{equation*}
L f_{0}=s_{L} f_{0} \tag{4.2}
\end{equation*}
$$

(iii) $\sigma_{b}(L)=s_{L}+\mathrm{i} \nu \mathbb{Z}$, for some $\nu \geq 0$;
(iv) $s_{L}$ is the only eigenvalue of $L$ admitting a positive eigenvector;
(v) $s_{L}$ is the only eigenvalue of $L^{\star}$ admitting a strictly positive eigenvector.

Proof. Items (i)-(iii) are proved, for instance, in [16, Thm. 8.17] (see also [2, Chap. C-III, Prop. 3.5]).

To prove item (iv), suppose that $\lambda \in \mathbb{C}$ is such that $L f=\lambda f$, where $f \in X_{+} \cap$ $D(L) \backslash\{0\}$. Then,

$$
\lambda\left\langle f, \varphi_{0}^{\star}\right\rangle=\left\langle L f, \varphi_{0}^{\star}\right\rangle=\left\langle f, L^{\star} \varphi_{0}^{\star}\right\rangle=s_{L}\left\langle f, \varphi_{0}^{\star}\right\rangle
$$

By [2, Chap. C-III, Prop. 3.5(a)], every positive eigenvector of $L$ must be a quasiinterior point of $X$, and hence $\left\langle f, \varphi_{0}^{\star}\right\rangle>0$; therefore, from the chain of equalities above we get $\lambda=s_{L}$.

Finally, to prove item (v), suppose that $\lambda^{\star} \in \mathbb{C}$ is such that $L^{\star} \varphi^{\star}=\lambda^{\star} \varphi^{\star}$, where $\varphi^{\star} \in X_{++}^{\star} \cap D\left(L^{\star}\right) \backslash\{0\}$. Then,

$$
\lambda^{\star}\left\langle f_{0}, \varphi^{\star}\right\rangle=\left\langle f_{0}, \lambda^{\star} \varphi^{\star}\right\rangle=\left\langle f_{0}, L^{\star} \varphi^{\star}\right\rangle=\left\langle L f_{0}, \varphi^{\star}\right\rangle=s_{L}\left\langle f_{0}, \varphi^{\star}\right\rangle
$$

Since $f_{0} \neq 0$ and $\varphi^{\star}$ is strictly positive, we have $\left\langle f_{0}, \varphi^{\star}\right\rangle>0$; therefore the chain of equalities above provides $\lambda^{\star}=s_{L}$.

Remark 4.5. Item (iii) of Theorem 4.4 entails that either $\sigma_{b}(L)=\left\{s_{L}\right\}$ or $\sigma_{b}(L)$ is an infinite unbounded set. If one is able to exclude the second case or to prove, at the least, that the intersection of $\sigma_{b}(L)$ with the point spectrum of $L$ is the singleton $\left\{s_{L}\right\}$, then, on the one hand, (4.2) implies that $s_{L}$ is the dominant eigenvalue of $L$; on the other hand, (4.1) implies that point (i) of Assumption 3.2 is satisfied.

Remark 4.6. The assumption that $s_{L}>-\infty$, i.e., that $\sigma(L) \neq \emptyset$, in Theorem 4.4 is essential, since there are examples of positive irreducible $C_{0}$-semigroups on Banach lattices such that $\sigma(L)=\emptyset$ (see, e.g. [2, Chap. C-III, Example 3.6]). Some conditions that, together with irreducibility, imply that this is not the case are stated in [2, Chap. C-III, Thm. 3.7].

Remark 4.7. Theorem 4.4 does not guarantee uniqueness (up to normalization) of a strictly positive eigenvector of $L^{\star}$ associated to $s_{L}$, since we do not know whether $s_{L}$ is a geometrically simple eigenvalue of $L^{\star}$.

We conclude our discussion by looking at the specific case where $X=\mathcal{C}(D)$, the space of real-valued continuous functions on a set $D$, that we suppose is a compact separable topological space. In this case, $X^{\star}$ is the space of bounded regular Borel measures on $D$. The following result is particularly useful in checking that point (i) of Assumption 2.2 is satisfied.

Proposition 4.8. If $X=\mathcal{C}(D)$, with $D$ a compact separable topological space, and if $\left\{\mathrm{e}^{t L}\right\}_{t \geq 0}$ is a positive $C_{0}$-semigroup, then $\left\{\mathrm{e}^{t L}\right\}_{t \geq 0}$ preserves strict positivity, i.e., $\mathrm{e}^{t L}\left(X_{++}\right) \subseteq X_{++}$for all $t \geq 0$. Moreover, the spectrum of $L$ is not empty, and $-\infty<s_{L} \in \sigma(L)$.

Proof. For the first part of the statement, see [2, Chap. B-II, Cor. 1.17]. A proof of the last assertion is given in [2, Chap. B-III, Thm. 1.1].

The following theorem provides a sufficient condition for ensuring that point (i) of Assumption 3.2 is verified.

Theorem 4.9. If $X=\mathcal{C}(D)$, with $D$ a compact separable topological space, and if $\left\{\mathrm{e}^{t L}\right\}_{t \geq 0}$ is a positive $C_{0}$-semigroup, the spectral and growth bounds of $L$ coincide, and there exists a positive probability measure $0 \neq \varphi^{\star} \in D\left(L^{\star}\right) \cap X_{+}^{\star}$ such that

$$
\begin{equation*}
L^{\star} \varphi^{\star}=s_{L} \varphi^{\star} \tag{4.3}
\end{equation*}
$$

Moreover, if $\left\{\mathrm{e}^{t L}\right\}_{t \geq 0}$ is irreducible, and if $s_{L}$ is a pole of the resolvent of $L$, then $\varphi^{\star} \in D\left(L^{\star}\right) \cap X_{++}^{\star}$, i.e., $\varphi^{\star}$ is strictly positive.

Proof. The proof of the first statement of the theorem can be found in [2, Chap. B-III, Thm. 1.6] (see also the discussion that follows). The last statement follows from Theorem 4.4, item (i) (see, also, [2, Chap. B-III, Prop. 3.5]).
5. Application to an economic growth problem with space heterogeneity. We apply the results of the previous sections to a family of economic growth problems in which the state variable, representing capital, is space heterogeneous, in the sense that it depends not only on time but also on the space location. Problems of this kind have been recently studied in the economic literature (see [6, 7, 8] for more details), embedding them in an infinite-dimensional state space of $L^{2}$ type. As mentioned in the introduction, this did not allow us to treat satisfactorily the presence of positivity constraints. Here, instead, we embed the problem in the space of continuous functions.

To introduce this family of economic growth problems, we specialize the abstract setting of section 2 with the following assumption that will be in force throughout this section.

Assumption 5.1.
(a) $D=\mathbf{S}^{1}:=\left\{\xi \in \mathbb{R}^{2}:|\xi|_{\mathbb{R}^{2}}=1\right\} \cong 2 \pi \mathbb{R} / \mathbb{Z}$. The space $\mathbf{S}^{1}$ is topologically equivalent to $[0,2 \pi] \subset \mathbb{R}$, whenever the endpoints of the latter interval are identified. Similarly, functions on $\mathbf{S}^{1}$ are identified with $2 \pi$-periodic functions on $\mathbb{R}$.
(b) $\mu$ is the Hausdorff measure on $\mathbf{S}^{1}$. Through the identification $\mathbf{S}^{1} \cong 2 \pi \mathbb{R} / \mathbb{Z}$, we will understand that $\mu$ is the Lebesgue measure on $[0,2 \pi]$. Accordingly, we will write $\mathrm{d} \theta$ instead of $\mu(\mathrm{d} \theta)$, when integrating functions with respect to $\mu$.
(c) $X=\mathcal{C}\left(\mathbf{S}^{1}\right)$, endowed with the sup-norm-i.e., $|x|_{X}=|x|_{\infty}:=\sup _{\theta \in D}|x(\theta)|-$ and with the pointwise ordering-i.e., $x \leq_{X} y \Longleftrightarrow x(\theta) \leq y(\theta)$ for all $\theta \in \mathbf{S}^{1}$.
(d) The operators $L$ and $N$ that define the state equation (2.1) and the function $u$ of Assumption 2.4 are explicitly given by

$$
\begin{equation*}
[L z](\theta)=\sigma \frac{\mathrm{d}^{2}}{\mathrm{~d} \theta^{2}} z(\theta)+A(\theta) z(\theta), \quad z \in X, \theta \in D \tag{5.1}
\end{equation*}
$$

with domain $D(L)=\mathcal{C}^{2}\left(\mathbf{S}^{1}\right)$,

$$
\begin{align*}
{[N z](\theta) } & =\eta(\theta) z(\theta), & & z \in X, \theta \in D  \tag{5.2}\\
u(\theta, \xi) & =\frac{\xi^{1-\gamma}}{1-\gamma} \eta(\theta)^{q}, & & \theta \in D, \xi \in \mathbb{R}_{+} \tag{5.3}
\end{align*}
$$

where $A, \eta \in \mathcal{C}\left(\mathbf{S}^{1} ;(0,+\infty)\right)$ are given functions, and $\sigma>0, q \geq 0, \gamma \in$ $(0,1) \cup(1,+\infty)$ are fixed constants.
Remark 5.2. Between the two possible choices of $X$ in Assumption 3.2, we are considering the set of real-valued continuous functions on a compact set. Note that the operator $N$ has the same expression given in (3.6), while $u$ is of the same type as in (3.7), with $f=\eta^{q}$. As already observed in point (v) of Remark 3.3, this implies that Assumption 2.2(ii) and Assumption 2.4 are verified. In the following, we will prove that point (i) of Assumption 2.2 and points (i)-(iv) of Assumption 3.2 are satisfied.

This specific setting corresponds to economic growth models, where the output of the economy is described by an AK production function, a one-dimensional geography is explicitly taken into account through a variable $\theta \in \mathbf{S}^{1}$, and capital diffusion over the geographical space is considered. Let us briefly summarize these problems (we refer the reader to $[6,7,8]$ for further details).

The state variable is the capital level $K(t, \theta)$ at time $t \geq 0$ and at location $\theta \in \mathbf{S}^{1}$. Its evolution depends on the consumption policy $C(t, \theta)$ that has to be determined by a social planner at each time $t \geq 0$ and at each location $\theta \in \mathbf{S}^{1}$ with the aim of maximizing a reward functional. More precisely, for any initial capital endowment $K_{0} \in \mathcal{C}\left(\mathbf{S}^{1} ;(0,+\infty)\right)$, the time-space evolution of the capital level obeys the following PDE:

$$
\left\{\begin{array}{l}
\frac{\partial K}{\partial t}(t, \theta)=\sigma \frac{\partial^{2} K}{\partial \theta^{2}}(t, \theta)+A(\theta) K(t, \theta)-\eta(\theta) C(t, \theta), \quad(t, \theta) \in \mathbb{R}_{+} \times \mathbf{S}^{1}  \tag{5.4}\\
K(0, \theta)=K_{0}(\theta), \quad \theta \in \mathbf{S}^{1}
\end{array}\right.
$$

The data appearing in this PDE are the functions $\eta, A$, where $\eta(\theta)$ represents the population density at location $\theta \in \mathbf{S}^{1}$, while $A(\theta)$ represents the technological level of the economy at location $\theta \in \mathbf{S}^{1}$, and a diffusion coefficient $\sigma$ that measures how fast the capital spreads over the geographical space. The optimization problem that the social planner aims to solve is

$$
\begin{equation*}
V_{++}\left(K_{0}\right):=\sup _{\mathcal{A}_{++}\left(K_{0}\right)} \int_{0}^{\infty} \mathrm{e}^{-\rho t}\left(\int_{\mathbf{S}^{1}} \frac{C(t, \theta)^{1-\gamma}}{1-\gamma} \eta(\theta)^{q} \mathrm{~d} \theta\right) \mathrm{d} t \tag{5.5}
\end{equation*}
$$

where $\rho>0$ is a fixed discount factor, the parameters $q \geq 0$ and $\gamma \in(0,1) \cup(1,+\infty)$ describe the preferences of the social planner, and

$$
\begin{equation*}
\mathcal{A}_{++}\left(K_{0}\right):=\left\{C \in \mathrm{~L}_{l o c}^{1}\left([0,+\infty) ; \mathcal{C}\left(\mathbf{S}^{1} ; \mathbb{R}^{+}\right)\right): K(t, \theta)>0 \forall(t, \theta) \in \mathbb{R}_{+} \times \mathbf{S}^{1}\right\} \tag{5.6}
\end{equation*}
$$

It is clear that under Assumption 5.1 this problem can be reformulated in the abstract setting introduced in section 2. More specifically, with the identifications $x_{0}:=K_{0}$, $x(t):=K(t, \cdot), c(t):=C(t, \cdot), \operatorname{PDE}(5.4)$ corresponds to the state equation (2.1), the optimization problem (5.5) coincides with problem ( $P$ ), and the class of admissible consumption plans (5.6) is precisely the class of admissible controls (2.4).

As in section 3, we can relax the state constraint (5.6) and consider the optimization problem

$$
\begin{equation*}
V^{b_{0}^{\star}}\left(K_{0}\right):=\sup _{\substack{\mathcal{A}_{++}^{b_{+}^{\star}}\left(K_{0}\right)}} \int_{0}^{\infty} \mathrm{e}^{-\rho t}\left(\int_{\mathbf{S}^{1}} \frac{C(t, \theta)^{1-\gamma}}{1-\gamma} \eta(\theta)^{q} \mathrm{~d} \theta\right) \mathrm{d} t, \tag{5.7}
\end{equation*}
$$

where $b_{0}^{\star}$ is the eigenvector of point (i) of Assumption 3.2, and

$$
\mathcal{A}_{++}^{b_{0}^{\star}}\left(K_{0}\right):=\left\{C \in \mathrm{~L}_{l o c}^{1}\left(\mathbb{R}^{+} ; \mathcal{C}\left(\mathbf{S}^{1} ; \mathbb{R}^{+}\right)\right):\left\langle K(t, \cdot), b_{0}^{\star}\right\rangle>0 \forall t \in \mathbb{R}_{+} \times \mathbf{S}^{1}\right\} .
$$

Exploiting the results proved in the previous sections, we are able to find an explicit solution to problem (5.7) and a feedback optimal control. Finally, studying the corresponding closed loop equation, we will prove optimality of this control for problem (5.5), at least for a suitable set of initial conditions $K_{0}$.

Let us consider the Hilbert space $\mathrm{L}^{2}\left(\mathbf{S}^{1}\right)$ of square integrable functions $x: \mathbf{S}^{1} \rightarrow \mathbb{R}$ endowed with inner product and induced norm

$$
\begin{aligned}
(x, y)_{\mathrm{L}^{2}\left(\mathbf{S}^{1}\right)} & :=\int_{\mathbf{S}^{1}} x(\theta) y(\theta) \mathrm{d} \theta, \quad x, y \in \mathrm{~L}^{2}\left(\mathbf{S}^{1}\right) ; \\
|x|_{\mathrm{L}^{2}\left(\mathbf{S}^{1}\right)} & :=\left(\int_{\mathbf{S}^{1}} x(\theta)^{2} \mathrm{~d} \theta\right)^{\frac{1}{2}}, \quad x \in \mathrm{~L}^{2}\left(\mathbf{S}^{1}\right) .
\end{aligned}
$$

Each element $f \in \mathrm{~L}^{2}\left(\mathbf{S}^{1}\right)$ can be seen as an element of $\mathcal{C}\left(\mathbf{S}^{1}\right)^{\star}$ through $\langle x, f\rangle:=$ $(x, f)_{\mathrm{L}^{2}\left(\mathbf{S}^{1}\right)}$, and we have the continuous embeddings $\mathcal{C}\left(\mathbf{S}^{1}\right) \hookrightarrow \mathrm{L}^{2}\left(\mathbf{S}^{1}\right) \hookrightarrow \mathcal{C}\left(\mathbf{S}^{1}\right)^{\star}$. The next proposition ensures that also point (i) of Assumption 2.2 is satisfied (for point (ii) of the same assumption; cf. Remark 5.2).

Proposition 5.3. By [2, Chap. B-III, Examples 2.14(a) and 3.4(e)] and [2, Chap. B-II, Cor. 1.17], we have the following results:
(i) The operator $L: D(L)=\mathcal{C}^{2}\left(\mathbf{S}^{1}\right) \subset \mathcal{C}\left(\mathbf{S}^{1}\right) \rightarrow \mathcal{C}\left(\mathbf{S}^{1}\right)$, defined in (5.1), is a closed operator generating a positive irreducible $C_{0}$-semigroup in the space $\mathcal{C}\left(\mathbf{S}^{1}\right)$.
(ii) Assumption 2.2 is satisfied.

Our first aim is to use Theorem 3.9 to provide an explicit solution to problem (5.7). To do so, we first need to ensure that Assumption 3.2 is verified.

By standard results from the Sturm-Liouville theory with periodic boundary conditions (see, e.g., [17, Chap. 7, Thms. 2.1 and 4.1; Chap. 8, Thm. 3.1]; cf. also [10]), we have the following results.

Proposition 5.4.
(i) The eigenvalues of $L$ are real and form a countable discrete set $\left\{\lambda_{k}\right\}_{k \in \mathbb{N}}$ such that $\lambda_{0}>\lambda_{1}>\cdots$, and $\lim _{k \rightarrow \infty} \lambda_{k}=-\infty$.
(ii) The eigenvalues $\lambda_{k}$, for $k \geq 1$, have geometric multiplicity 2 , while the highest eigenvalue, $\lambda_{0}$, is simple. The corresponding (normalized) eigenvectors $\left\{b_{k}\right\}_{k \in \mathbb{N}}$ form an orthonormal basis of $\mathrm{L}^{2}\left(\mathbf{S}^{1}\right)$.
(iii) Each $b_{k}$ has exactly $k$ zeros in $\mathbf{S}^{1}$; without loss of generality, we take $b_{0}>0$; this is, up to a positive multiplicative constant, the unique strictly positive eigenvector of $L$.
(iv) For all $x \in \mathcal{C}\left(\mathbf{S}^{1}\right)$,

$$
\begin{equation*}
x=\sum_{k=0}^{\infty}\left(x, b_{k}\right)_{\mathrm{L}^{2}\left(\mathbf{S}^{1}\right)} b_{k}:=\lim _{n \rightarrow \infty} \sum_{k=0}^{n}\left(x, b_{k}\right)_{\mathrm{L}^{2}\left(\mathbf{S}^{1}\right)} b_{k} \quad \text { in }|\cdot|_{\infty} \tag{5.8}
\end{equation*}
$$

(v) For all $\mu \in \mathbb{C} \backslash\left\{\lambda_{k}\right\}_{k \in \mathbb{N}}$, the operator $L-\mu: \mathcal{C}^{2}\left(\mathbf{S}^{1}\right) \rightarrow \mathcal{C}\left(\mathbf{S}^{1}\right)$ is invertible, and the inverse $(L-\mu)^{-1}: \mathcal{C}\left(\mathbf{S}^{1}\right) \rightarrow \mathcal{C}^{2}\left(\mathbf{S}^{1}\right)$ is bounded and admits the representation

$$
\begin{equation*}
(L-\mu)^{-1} x=\sum_{k=0}^{\infty} \frac{\left(b_{k}, x\right)_{\mathrm{L}^{2}\left(\mathbf{S}^{1}\right)}}{\lambda_{k}-\mu} b_{k}:=\lim _{n \rightarrow \infty} \sum_{k=0}^{n} \frac{\left(b_{k}, x\right)_{\mathrm{L}^{2}\left(\mathbf{S}^{1}\right)}}{\lambda_{k}-\mu} b_{k} \quad \text { in }|\cdot|_{\infty} . \tag{5.9}
\end{equation*}
$$

In particular the resolvent set of $L$ is $\varrho(L)=\mathbb{C} \backslash\left\{\lambda_{k}\right\}_{k \geq 0}$.
Remark 5.5. The last assertion of Proposition 5.4 entails that the spectrum of operator $L$ is $\sigma(L)=\left\{\lambda_{k}\right\}_{k \geq 0}$. Moreover, using (5.8) and the fact that $L$ is closed, we have that

$$
\begin{equation*}
L x=\sum_{k \in \mathbb{N}}\left(x, b_{k}\right)_{\mathrm{L}^{2}\left(\mathbf{S}^{1}\right)} \lambda_{k} b_{k} \quad \text { in }|\cdot|_{\infty}, \quad x \in \mathcal{C}^{2}\left(\mathbf{S}^{1}\right) \tag{5.10}
\end{equation*}
$$

From the facts above, we have that $s_{L}=\lambda_{0}$ is the dominant eigenvalue of the operator $L$, and Theorem 4.9 guarantees that there exists a strictly positive eigenvector $b_{0}^{\star} \in D\left(L^{\star}\right) \cap \mathcal{C}\left(\mathbf{S}^{1}\right)_{++}^{\star}$ of $L^{\star}$ associated to the eigenvalue $\lambda_{0}^{\star}=\lambda_{0}$. Moreover, thanks to point (v) of Theorem 4.4, $\lambda_{0}$ is the only eigenvalue of $L^{\star}$ admitting a strictly positive eigenvector. Hence, point (i) of Assumption 3.2 is satisfied. The next proposition characterizes $b_{0}^{\star}$.

Proposition 5.6. Every strictly positive eigenvector of $L^{\star}$ associated to $\lambda_{0}$ is in the form $b_{0}^{\star}=\kappa b_{0}$ for some $\kappa>0$.

Proof. Take $x \in \mathcal{C}\left(\mathbf{S}^{1}\right)$. Using (5.10) and by uniform convergence, we get

$$
\left\langle L x, b_{0}\right\rangle=\sum_{k=0}^{\infty}\left(x, b_{k}\right)_{\mathrm{L}^{2}\left(\mathbf{S}^{1}\right)} \lambda_{k}\left\langle b_{k}, b_{0}\right\rangle=\lambda_{0}\left(x, b_{0}\right)_{\mathrm{L}^{2}\left(\mathbf{S}^{1}\right)}\left|b_{0}\right|_{\mathrm{L}^{2}\left(\mathbf{S}^{1}\right)}^{2}=\left\langle x, \lambda_{0} b_{0}\right\rangle
$$

Since $x \in \mathcal{C}\left(\mathbf{S}^{1}\right)$ is arbitrary, we obtain that $b_{0} \in D\left(L^{\star}\right)$ and $L^{\star} b_{0}=\lambda_{0} b_{0}$.
In the following, we set $\kappa:=1$, i.e., we take $b_{0}^{\star}:=b_{0}$. The result above, together with continuity and positivity of the function $\eta$ (recall also that $f=\eta^{q}$ in this section), shows that point (iv) of Assumption 3.2 is satisfied. Moreover, thanks to the compactness of $D=\mathbf{S}^{1}$, we also get that point (iii) of Assumption 3.2 is verified. Finally, assuming that point (ii) of Assumption 3.2 holds, i.e., that

$$
\begin{equation*}
\rho>\lambda_{0}(1-\gamma) \tag{5.11}
\end{equation*}
$$

all the hypotheses of Theorem 3.9 are verified, and hence, by making explicit the latter for the economic problem introduced above, we get the following.

Theorem 5.7. We have the following facts.
(i) The value function of problem (5.7) is

$$
\begin{equation*}
V^{b_{0}^{\star}}\left(K_{0}\right)=\frac{\alpha}{1-\gamma}\left(\int_{\mathbf{S}^{1}} K_{0}(\theta) b_{0}(\theta) \mathrm{d} \theta\right)^{1-\gamma} \tag{5.12}
\end{equation*}
$$

where

$$
\begin{equation*}
\alpha:=\left[\frac{\gamma}{\rho-\lambda_{0}(1-\gamma)} \int_{\mathbf{S}^{1}} \eta(\theta)^{\frac{q+\gamma-1}{\gamma}} b_{0}(\theta)^{\frac{\gamma-1}{\gamma}} \mathrm{~d} \theta\right]^{\gamma} \tag{5.13}
\end{equation*}
$$

(ii) the control

$$
\begin{equation*}
\hat{C}(t, \theta):=\left(\int_{\mathbf{S}^{1}} K_{0}(\xi) b_{0}(\xi) \mathrm{d} \xi\right)\left(\alpha b_{0}(\theta)\right)^{-\frac{1}{\gamma}} \eta(\theta)^{\frac{q-1}{\gamma}} \mathrm{e}^{g t} \tag{5.14}
\end{equation*}
$$

where $g:=\frac{\lambda_{0}-\rho}{\gamma}$, is optimal for problem (5.7), and the corresponding optimal capital $\hat{K}$ solves, for $(t, \theta) \in \mathbb{R}_{+} \times \mathbf{S}^{1}$, the linear integro-PDE

$$
\left\{\begin{array}{l}
\frac{\partial K}{\partial t}(t, \theta)=\sigma \frac{\partial^{2} K}{\partial \theta^{2}}(t, \theta)+A(\theta) K(t, \theta)  \tag{5.15}\\
\quad-\left(\alpha b_{0}(\theta)\right)^{-\frac{1}{\gamma}}\left(\int_{\mathbf{S}^{1}} K(t, \xi) b_{0}(\xi) \mathrm{d} \xi\right) \eta(\theta)^{\frac{q+\gamma-1}{\gamma}} \\
K(0, \theta)=K_{0}(\theta)
\end{array}\right.
$$

Remark 5.8. The value $g:=\frac{\lambda_{0}-\rho}{\gamma}$ is the optimal growth rate of the economy.
Our second aim is to use the stability result provided by Theorem 3.11 to show that the optimal control $\hat{C}$ for problem (5.7), defined in (5.14), is optimal also for problem (5.5), at least for a suitable set of initial conditions $K_{0}$. To achieve this, we need to analyze the semigroup generated by the operator $B=L-N \Phi$. By (5.15), its explicit expression is, for all $x \in \mathcal{C}^{2}\left(\mathbf{S}^{1}\right)$,

$$
\begin{equation*}
(B x)(\theta)=\sigma \frac{\mathrm{d}^{2} x}{\mathrm{~d} \theta^{2}}(\theta)+A(\theta) x(\theta)-\left(\alpha b_{0}(\theta)\right)^{-\frac{1}{\gamma}}\left(\int_{\mathbf{S}^{1}} x(\xi) b_{0}(\xi) \mathrm{d} \xi\right) \eta(\theta)^{\frac{q+\gamma-1}{\gamma}}, \theta \in \mathbf{S}^{1} \tag{5.16}
\end{equation*}
$$

Lemma 5.9. $\left\{\mathrm{e}^{t B}\right\}_{t \geq 0}$ is (immediately) compact.
Proof. We start by showing that the semigroup generated by the operator $L$ on $X$ is (immediately) compact. In fact, $L$ generates an analytic semigroup (see, e.g., [21, Chap. VI, sect. 4]), and the resolvent of $L$ is nonempty; indeed, it is the complement set of the spectrum of $L$, which is a pure point spectrum $\sigma(L)=\left\{\lambda_{0}, \lambda_{1}, \ldots\right\}$ by point (i) of Proposition 5.4.

From [21, Chap. II, eq. (4.26)] we deduce that $\left\{\mathrm{e}^{t L}\right\}_{t \geq 0}$ is (immediately) norm continuous. Then, since the canonical injection $\iota:\left(\mathcal{C}^{2}\left(\mathbf{S}^{1}\right),|\cdot|_{\mathcal{C}}{ }^{2}\left(\mathbf{S}^{1}\right)\right) \rightarrow\left(\mathcal{C}\left(\mathbf{S}^{1}\right),|\cdot|_{\infty}\right)$ is compact by the Ascoli-Arzelà theorem, we can apply [21, Chap. II, Prop. 4.25] and deduce that $L$ has a compact resolvent. Finally, we learn from [21, Chap. II, Thm. 4.29] that $\left\{\mathrm{e}^{t L}\right\}_{t \geq 0}$ is (immediately) compact.

Next, notice that $B$ is an additive perturbation of $L$, obtained by subtracting from it the operator $N \Phi$, which is bounded. Therefore, from [21, Chap. III, Prop. 1.16(i)], we deduce that $\left\{\mathrm{e}^{t B}\right\}_{t \geq 0}$ is (immediately) compact.

Next, we prove that, under suitable assumptions, $g$ is a simple eigenvalue of $B$, and we characterize its associated eigenvector. To ease notation, let us define the constant

$$
\begin{equation*}
\alpha_{0}:=\alpha^{\frac{1}{1-\gamma}}=\left[\frac{\gamma}{\rho-\lambda_{0}(1-\gamma)} \int_{\mathbf{S}^{1}} \eta(\theta)^{\frac{q+\gamma-1}{\gamma}} b_{0}(\theta)^{\frac{\gamma-1}{\gamma}} \mathrm{~d} \theta\right]^{\frac{\gamma}{1-\gamma}} \tag{5.17}
\end{equation*}
$$

and the function

$$
\begin{equation*}
\beta(\theta):=\alpha_{0} b_{0}(\theta), \quad \theta \in \mathbf{S}^{1} . \tag{5.18}
\end{equation*}
$$

Using (5.17) and (5.18), we can rewrite the expression for operator $N \Phi$, appearing in (5.16) and defined in (3.16), as follows:

$$
\begin{equation*}
(N \Phi) x=\left(\alpha b_{0}\right)^{-\frac{1}{\gamma}}\left(\int_{\mathbf{S}^{1}} x(\xi) b_{0}(\xi) \mathrm{d} \xi\right) \eta^{\frac{q+\gamma-1}{\gamma}}=\beta^{-\frac{1}{\gamma}} \eta^{\frac{q+\gamma-1}{\gamma}}\langle x, \beta\rangle, \quad x \in \mathcal{C}\left(\mathbf{S}^{1}\right) \tag{5.19}
\end{equation*}
$$

Notice that, under our assumptions, $\beta^{-\frac{1}{\gamma}} \eta^{\frac{q+\gamma-1}{\gamma}} \in \mathcal{C}\left(\mathbf{S}^{1}\right)$, and hence it admits a Fourier series expansion with uniform convergence with respect to $\left\{b_{k}\right\}_{k \in \mathbb{N}}$, with coefficients:

$$
\begin{array}{ll}
\beta_{0}:=\left(b_{0}, \beta^{-\frac{1}{\gamma}} \eta^{\frac{q+\gamma-1}{\gamma}}\right)_{\mathrm{L}^{2}\left(\mathbf{S}^{1}\right)}=\frac{\mu_{0}}{\alpha_{0}}, & \text { where } \mu_{0}:=\left(\beta, \beta^{-\frac{1}{\gamma}} \eta^{\frac{q+\gamma-1}{\gamma}}\right)_{\mathrm{L}^{2}\left(\mathbf{S}^{1}\right)} \\
\beta_{k}:=\left(b_{k}, \beta^{-\frac{1}{\gamma}} \eta^{\frac{q+\gamma-1}{\gamma}}\right)_{\mathrm{L}^{2}\left(\mathbf{S}^{1}\right)}, & k \geq 1 . \tag{5.21}
\end{array}
$$

Using (5.17), it is immediate to see that $\mu_{0}=\lambda_{0}-g$. Hence,

$$
\beta^{-\frac{1}{\gamma}} \eta^{\frac{q+\gamma-1}{\gamma}}=\frac{\lambda_{0}-g}{\alpha_{0}} b_{0}+\sum_{k=1}^{\infty} \beta_{k} b_{k},
$$

with convergence in $\mathcal{C}\left(\mathbf{S}^{1}\right)$. Consider the formal series

$$
\begin{equation*}
\frac{b_{0}}{\alpha_{0}}+\sum_{k=1}^{\infty} \frac{\beta_{k}}{\lambda_{k}-g} b_{k} \tag{5.22}
\end{equation*}
$$

Proposition 5.10. Let (5.11) hold. The series (5.22) converges in $\mathcal{C}\left(\mathbf{S}^{1}\right)$ and defines a function of $\mathcal{C}^{2}\left(\mathbf{S}^{1}\right)$.

Proof. By (5.11), $g$ is not an eigenvalue of $L$; moreover, $\beta^{-\frac{1}{\gamma}} \eta^{\frac{q+\gamma-1}{\gamma}} \in \mathcal{C}\left(\mathbf{S}^{1}\right)$. Using (5.9), (5.20) and (5.21) and recalling that $\mu_{0}=\lambda_{0}-g$, we get

$$
\begin{equation*}
(L-g)^{-1}\left(\beta^{-\frac{1}{\gamma}} \eta^{\frac{q+\gamma-1}{\gamma}}\right)=\sum_{k=0}^{\infty} \frac{\left(b_{k}, \beta^{-\frac{1}{\gamma}} \eta^{\frac{q+\gamma-1}{\gamma}}\right)_{\mathrm{L}^{2}\left(\mathbf{S}^{1}\right)}}{\lambda_{k}-g} b_{k}=\frac{b_{0}}{\alpha_{0}}+\sum_{k=1}^{\infty} \frac{\beta_{k}}{\lambda_{k}-g} b_{k} . \square \tag{5.23}
\end{equation*}
$$

Given the proposition above, we set

$$
\begin{equation*}
w:=\frac{b_{0}}{\alpha_{0}}+\sum_{k=1}^{\infty} \frac{\beta_{k}}{\lambda_{k}-g} b_{k} \in \mathcal{C}^{2}\left(\mathbf{S}^{1}\right) \tag{5.24}
\end{equation*}
$$

We are ready to state the following result on the spectrum of operator $B$.
Proposition 5.11. Let (5.11) hold, and assume that $g>\lambda_{1}$. Then, we have the following facts:
(i) The spectrum of the operator $B$ is given by $\sigma(B)=\{g\} \cup\left\{\lambda_{k}\right\}_{k \in \mathbb{N} \backslash\{0\}}$, and hence $g$ is the dominant eigenvalue of $B$;
(ii) $g$ is a simple eigenvalue of $B$;
(iii) $\operatorname{ker}(B-g)=\operatorname{Span}\{w\}$, where $w$ is given as in (5.24).
(iv) The spectral projection $P$ corresponding to the spectral set $\{g\} \subset \sigma(B)$ is the bounded and linear operator:

$$
\begin{equation*}
P x=\langle x, \beta\rangle w, \quad x \in \mathcal{C}\left(\mathbf{S}^{1}\right) \tag{5.25}
\end{equation*}
$$

Proof. (i) We start by showing that $g \in \sigma(B)$. By (5.11), $g$ is not an eigenvalue of $L$. Thus, [21, Chap. IV, Prop. 4.2] entails that $g$ is an eigenvalue of $B=L-N \Phi$ if and only if 1 is an eigenvalue of $(L-g)^{-1} N \Phi$, i.e., if and only if the equation

$$
\begin{equation*}
(L-g)^{-1} N \Phi x=x \tag{5.26}
\end{equation*}
$$

admits a nontrivial solution $x \in \mathcal{C}\left(\mathbf{S}^{1}\right)$. By (5.19) and (5.23), $(L-g)^{-1} N \Phi x=$ $\langle x, \beta\rangle w$, for all $x \in \mathcal{C}\left(\mathbf{S}^{1}\right)$, and hence (5.26) becomes $\langle x, \beta\rangle w=x$. By (5.24) we have $\langle w, \beta\rangle=(w, \beta)_{\mathrm{L}^{2}\left(\mathbf{S}^{1}\right)}=\left(b_{0}, b_{0}\right)_{\mathrm{L}^{2}\left(\mathbf{S}^{1}\right)}=1$. Therefore, we obtain that $w$ is a nontrivial solution to (5.26), and thus $g$ is an eigenvalue of $B$.

Next, we prove that $\lambda_{0} \notin \sigma(B)$. Suppose, by contradiction, that $\lambda_{0}$ is an eigenvalue of $B$ and that $x$ is an associated eigenfunction. Let $x_{k}:=\left(x, b_{k}\right)_{\mathrm{L}^{2}\left(\mathbf{S}^{1}\right)}$. By (5.10), recalling (5.19) and that $\beta^{-\frac{1}{\gamma}} \eta^{\frac{q+\gamma-1}{\gamma}}$ admits a Fourier series expansion whose coefficients are given in (5.20)-(5.21), we get

$$
0=B x-\lambda_{0} x=\sum_{k=0}^{\infty} x_{k} \lambda_{k} b_{k}-\left(\frac{\mu_{0}}{\alpha_{0}} b_{0}+\sum_{k=1}^{\infty} \beta_{k} b_{k}\right) x_{0} \alpha_{0}-\sum_{k=0}^{\infty} \lambda_{0} x_{k} b_{k}
$$

This equation is satisfied if and only if

$$
\left\{\begin{array}{l}
\lambda_{0} x_{0}-\mu_{0} x_{0}=\lambda_{0} x_{0}, \\
\lambda_{k} x_{k}-x_{0} \alpha_{0} \beta_{k}=\lambda_{0} x_{k}, \quad k \in \mathbb{N} \backslash\{0\}
\end{array}\right.
$$

Since $\mu_{0}=\lambda_{0}-g$, which entails $\mu_{0} \neq 0$ by (5.11), the first equation implies $x_{0}=0$. Substituting into the second equation, we get $x_{k}\left(\lambda_{k}-\lambda_{0}\right)=0$ for all $k \geq 1$. Since $x \neq 0$, there must be an index $k \geq 1$ such that $\lambda_{k}=\lambda_{0}$, a contradiction.

Finally, to show that $\lambda_{k}, k \geq 1$, are the only other eigenvalues of $B$, it is enough to prove that $\mu \in \mathbb{C} \backslash\left\{g, \lambda_{0}\right\}$ is an eigenvalue of $B$ if and only if $\mu$ is an eigenvalue of $L$. This is readily verified with computations similar to those above.
(ii) We start by recalling that, since $\left\{\mathrm{e}^{t B}\right\}_{t \geq 0}$ is a compact semigroup (Lemma 5.9), we know from Lemma 3.10 that $g$ is a pole of the resolvent of $B$ of finite algebraic multiplicity. Let $h \in \mathcal{C}\left(\mathbf{S}^{1}\right)$ and $\mu \in \mathbb{C} \backslash \sigma(B)$, and set $h_{k}:=\left(h, b_{k}\right)_{\mathrm{L}^{2}\left(\mathbf{S}^{1}\right)}$ for $k \in \mathbb{N}$. We claim that

$$
\begin{equation*}
(B-\mu)^{-1} h=(L-\mu)^{-1}\left(h+\frac{h_{0} \alpha_{0}}{g-\mu} \beta^{-\frac{1}{\gamma}} \eta^{\frac{q+\gamma-1}{\gamma}}\right) . \tag{5.27}
\end{equation*}
$$

First, note that, since $\mu \neq g$, we have $h+\frac{h_{0} \alpha_{0}}{g-\mu} \beta^{-\frac{1}{\gamma}} \eta^{\frac{q+\gamma-1}{\gamma}} \in \mathcal{C}\left(\mathbf{S}^{1}\right)$, so the right-hand side above is well defined and belongs to $\mathcal{C}^{2}\left(\mathbf{S}^{1}\right)$. Next, recalling (5.9) and given that
$N \Phi$ is a bounded linear operator, we get

$$
\begin{align*}
& N \Phi(L-\mu)^{-1}\left(h+\frac{h_{0} \alpha_{0}}{g-\mu} \beta^{-\frac{1}{\gamma}} \eta^{\frac{q+\gamma-1}{\gamma}}\right)  \tag{5.28}\\
= & \beta^{-\frac{1}{\gamma}} \eta^{\frac{q+\gamma-1}{\gamma}}\left\langle\sum_{k=0}^{\infty} \frac{h_{k}}{\lambda_{k}-\mu} b_{k}, \beta\right\rangle+\frac{h_{0} \alpha_{0}}{g-\mu} \beta^{-\frac{1}{\gamma}} \eta^{\frac{q+\gamma-1}{\gamma}}\left\langle\frac{\mu_{0}}{\alpha_{0}\left(\lambda_{0}-\mu\right)} b_{0}+\sum_{k=1}^{\infty} \frac{\beta_{k}}{\lambda_{k}-\mu} b_{k}, \beta\right\rangle \\
= & \frac{h_{0} \alpha_{0}\left(g-\mu+\mu_{0}\right)}{(g-\mu)\left(\lambda_{0}-\mu\right)} \beta^{-\frac{1}{\gamma}} \eta^{\frac{q+\gamma-1}{\gamma}} .
\end{align*}
$$

Therefore, recalling that $\mu_{0}=\lambda_{0}-g$, we obtain

$$
\begin{aligned}
& (B-\mu)(L-\mu)^{-1}\left(h+\frac{h_{0} \alpha_{0}}{g-\mu} \beta^{-\frac{1}{\gamma}} \eta^{\frac{q+\gamma-1}{\gamma}}\right) \\
= & (L-\mu)(L-\mu)^{-1}\left(h+\frac{h_{0} \alpha_{0}}{g-\mu} \beta^{-\frac{1}{\gamma}} \eta^{\frac{q+\gamma-1}{\gamma}}\right) \\
& -N \Phi(L-\mu)^{-1}\left(h+\frac{h_{0} \alpha_{0}}{g-\mu} \beta^{-\frac{1}{\gamma}} \eta^{\frac{q+\gamma-1}{\gamma}}\right)=h .
\end{aligned}
$$

This shows (5.27), i.e., it provides an explicit expression of the resolvent of $B$. Finally, recalling that the operator $(L-\mu)^{-1}$ is well defined for $\mu=g$, since $g$ is not an eigenvalue of $L$ under (5.11) and under the assumption $g>\lambda_{1}$, we deduce from (5.27) that $g$ is a simple pole of the resolvent of $B$.
(iii) Using (5.19) and (5.23)-(5.24), it can be immediately proved that $w$ is an eigenfunction of $B$ associated to $g$.
(iv) We use the explicit formula for $P$ given in Remark 3.12. By (3.30) and (5.27) we have that for any $x \in \mathcal{C}\left(\mathbf{S}^{1}\right)$, and setting $x_{k}:=\left(x, b_{k}\right)_{\mathrm{L}^{2}\left(\mathbf{S}^{1}\right)}$,

$$
\begin{equation*}
P x=-\frac{1}{2 \pi \mathrm{i}} \int_{\gamma}(B-\mu)^{-1} x \mathrm{~d} \mu=-\frac{1}{2 \pi \mathrm{i}} \int_{\gamma}(L-\mu)^{-1}\left(x+\frac{x_{0} \alpha_{0}}{g-\mu} \beta^{-\frac{1}{\gamma}} \eta^{\frac{q+\gamma-1}{\gamma}}\right) \mathrm{d} \mu \tag{5.29}
\end{equation*}
$$

where we take $\gamma$ to be the simple curve given by the positively oriented boundary of the disk centered at $g$, with radius small enough so that it does not encircle any of the eigenvalues $\left\{\lambda_{k}\right\}_{k \in \mathbb{N}}$ of $L$, which are poles of the resolvent $(L-\mu)^{-1}$. Substituting (5.9) into (5.29), and recalling (5.20) and (5.21), we get

$$
\begin{align*}
P x= & -\frac{1}{2 \pi \mathrm{i}} \int_{\gamma} \sum_{k \in \mathbb{N}} \frac{x_{k}}{\lambda_{k}-\mu} b_{k} \mathrm{~d} \mu-\frac{x_{0} \alpha_{0}}{2 \pi \mathrm{i}} \int_{\gamma} \sum_{k \in \mathbb{N}} \frac{\left(b_{k}, \beta^{-\frac{1}{\gamma}} \eta^{\frac{q+\gamma-1}{\gamma}}\right)_{\mathrm{L}^{2}\left(\mathbf{S}^{1}\right)}}{\left(\lambda_{k}-\mu\right)(g-\mu)} b_{k} \mathrm{~d} \mu \\
= & -\sum_{k \in \mathbb{N}} x_{k} b_{k} \frac{1}{2 \pi \mathrm{i}} \int_{\gamma} \frac{1}{\lambda_{k}-\mu} \mathrm{d} \mu-x_{0} \mu_{0} b_{0} \frac{1}{2 \pi \mathrm{i}} \int_{\gamma} \frac{1}{\left(\lambda_{0}-\mu\right)(g-\mu)} \mathrm{d} \mu  \tag{5.30}\\
& -\frac{x_{0} \alpha_{0}}{2 \pi \mathrm{i}} \sum_{k \in \mathbb{N} \backslash\{0\}} \beta_{k} b_{k} \int_{\gamma} \frac{1}{\left(\lambda_{k}-\mu\right)(g-\mu)} \mathrm{d} \mu .
\end{align*}
$$

Since any $\lambda_{k}, k \in \mathbb{N}$, lies outside of the curve $\gamma$, we have $\int_{\gamma} \frac{1}{\lambda_{k}-\mu} \mathrm{d} \mu=0$. Moreover, since the residues at $g$ of the functions $\frac{1}{\left(\lambda_{k}-\mu\right)(g-\mu)}$, for $k \in \mathbb{N}$, are equal to $-\frac{1}{\lambda_{k}-g}$, and recalling that $\gamma$ is a simple curve, we get

$$
\int_{\gamma} \frac{1}{\left(\lambda_{k}-\mu\right)(g-\mu)} \mathrm{d} \mu=-\frac{2 \pi \mathrm{i}}{\lambda_{k}-g}, \quad k \in \mathbb{N}
$$

Therefore, since $\lambda_{0}-g=\mu_{0},(5.30)$ becomes

$$
P x=\frac{x_{0} \mu_{0}}{\lambda_{0}-g} b_{0}+\sum_{k \in \mathbb{N} \backslash\{0\}} \frac{x_{0} \alpha_{0} \beta_{k}}{\lambda_{k}-g} b_{k}=x_{0} \alpha_{0}\left(\frac{b_{0}}{\alpha_{0}}+\sum_{k \in \mathbb{N} \backslash\{0\}} \frac{\beta_{k}}{\lambda_{k}-g} b_{k}\right)=\langle x, \beta\rangle w
$$

By the previous results, we are ready to apply Theorem 3.11 to the economic problem studied here and obtain the following.

THEOREM 5.12. Let (5.11) hold, and assume that $g>\lambda_{1}$. Let $K_{0} \in \mathcal{C}\left(\mathbf{S}^{1}\right)_{++}$, and denote by $\hat{K}^{K_{0}}$ the solution to the linear integro-PDE (5.15). Then,

$$
\begin{equation*}
\left|\hat{K}_{g}^{K_{0}}(t, \cdot)-\left\langle K_{0}, \beta\right\rangle w\right|_{\infty} \leq M \mathrm{e}^{-\left(g-\lambda_{1}\right) t}\left|K_{0}-\left\langle K_{0}, \beta\right\rangle w\right|_{\infty} \quad \forall t \geq 0 \tag{5.31}
\end{equation*}
$$

where $\beta$ is defined by (5.18), $\hat{K}_{g}^{K_{0}}(t, \cdot):=\mathrm{e}^{-g t} \hat{K}^{K_{0}}(t, \cdot)$, and

$$
M:=1+|w|_{\infty} \int_{0}^{2 \pi} \beta(\theta) \mathrm{d} \theta
$$

In particular,

$$
\hat{K}_{g}^{K_{0}}(t, \cdot) \xrightarrow{t \rightarrow+\infty}\left\langle K_{0}, \beta\right\rangle w \quad \text { in } \mathcal{C}\left(\mathbf{S}^{1}\right) .
$$

Proof. By Proposition 5.11 and Theorem 3.11, we get

$$
\begin{equation*}
\left|\hat{K}_{g}^{K_{0}}(t, \cdot)-\left\langle K_{0}, \beta\right\rangle w\right|_{\infty} \leq M \mathrm{e}^{-\varepsilon t}\left|K_{0}-\left\langle K_{0}, \beta\right\rangle w\right|_{\infty} \quad \forall t \geq 0 \tag{5.32}
\end{equation*}
$$

We now compute explicitly the constants $M$ and $\varepsilon$ appearing in (5.32). From the proof of Theorem 3.11 we know that these two constants come from the estimate provided in (3.29), which we can now refine.

For the sake of clarity, define $S(t):=\mathrm{e}^{t B}, t \geq 0$. Since $w$ is an eigenvector of $B$ with eigenvalue $g$, it follows from [21, Thm. 3.6, Chap. IV, p. 276] that $\mathrm{e}^{g t}$ is an eigenvalue of $S(t)$ with eigenvector $w$. Using the fact that Range $P=\operatorname{Span}\{w\}$, we have $S(t) P=\mathrm{e}^{g t} P$ for all $t \geq 0$. Moreover, since Range $(1-P)=\operatorname{ker} P$, we have $S(t)(1-P)=S(t)_{\mid \operatorname{ker} P}(1-P)$, for all $t \geq 0$, where $S(t)_{\mid \text {ker } P}$ indicates the restriction of the operator $S(t)$ to ker $P$. Hence, denoting by $|\cdot|_{\mathcal{L}\left(\mathcal{C}\left(\mathbf{S}^{1}\right)\right)}$ the operator norm on the space of linear continuous operators $\mathcal{L}\left(\mathcal{C}\left(\mathbf{S}^{1}\right)\right)$, we get

$$
\begin{align*}
& \left|\mathrm{e}^{-g t} S(t)-P\right|_{\mathcal{L}\left(\mathcal{C}\left(\mathbf{S}^{1}\right)\right)}=\left|\mathrm{e}^{-g t}(S(t)-S(t) P)\right|_{\mathcal{L}\left(\mathcal{C}\left(\mathbf{S}^{1}\right)\right)}=\mathrm{e}^{-g t}|S(t)-S(t) P|_{\mathcal{L}\left(\mathcal{C}\left(\mathbf{S}^{1}\right)\right)}  \tag{5.33}\\
= & \mathrm{e}^{-g t}|S(t)(1-P)|_{\mathcal{L}\left(\mathcal{C}\left(\mathbf{S}^{1}\right)\right)} \leq \mathrm{e}^{-g t}\left|S(t)_{\mid \operatorname{ker} P}\right|_{\mathcal{L}\left(\mathcal{C}\left(\mathbf{S}^{1}\right)\right)}|(1-P)|_{\mathcal{L}\left(\mathcal{C}\left(\mathbf{S}^{1}\right)\right)}, \quad t \geq 0
\end{align*}
$$

We now estimate $\left|S(t)_{\mid \text {ker } P}\right|_{\mathcal{L}\left(\mathcal{C}\left(\mathbf{S}^{1}\right)\right)}$ and $|(1-P)|_{\mathcal{L}\left(\mathcal{C}\left(\mathbf{S}^{1}\right)\right)}$.
We begin with the first quantity. We aim to prove, with the aid of the Hille-Yosida theorem, that $\left|S(t)_{\mid \text {ker } P}\right|$ satisfies

$$
\begin{equation*}
\left|S(t)_{\mid \operatorname{ker} P}\right|_{\mathcal{L}\left(\mathcal{C}\left(\mathbf{S}^{1}\right)\right)} \leq \mathrm{e}^{\lambda_{1} t}, \quad t \geq 0 \tag{5.34}
\end{equation*}
$$

From the discussion on subspace semigroups given in [21, paragraph II.2.3], we deduce that the infinitesimal generator of the semigroup $\left\{S(t)_{\mid \operatorname{ker} P}\right\}_{t \geq 0}$ is the operator $B_{\mid \text {ker } P} \equiv L_{\mid \text {ker } P}$, whose domain is $D\left(B_{\mid \text {ker } P}\right)=D\left(L_{\mid \text {ker } P}\right)=D(L) \cap$ ker $P$. Spectra of $B_{\mid \text {ker } P}$ and $L_{\mid \text {ker } P}$ must coincide, and clearly $\sigma\left(L_{\mid \text {ker } P}\right) \subset \sigma(L)$. In particular, recalling that ker $P$ is the closed subspace of $\mathcal{C}\left(\mathbf{S}^{1}\right)$ of functions $x$ such that $\langle x, \beta\rangle=0$, we
have that $\sigma\left(L_{\mid k e r P}\right)=\left\{\lambda_{k}\right\}_{k \in \mathbb{N} \backslash\{0\}}$. Indeed, on the one hand, $\lambda_{0}$ cannot be an eigenvalue of $L_{\mid \text {ker } P}$, since the associated eigenvector $b_{0}$ does not satisfy $\left\langle b_{0}, \beta\right\rangle=0$; on the other hand, the remaining eigenvalues of $L$ are also eigenvalues of $L_{\mid \operatorname{ker} P}$, since the corresponding eigenvectors belong to ker $P$ and hence are in $D\left(L_{\mid \operatorname{ker} P}\right)=D\left(B_{\mid \operatorname{ker} P}\right)$. This shows that the resolvent set of $B_{\mid k e r} P$ satisfies $\rho\left(B_{\mid \mathrm{ker} P}\right) \supset\left(\lambda_{1},+\infty\right)$. Then, in order to apply the Hille-Yosida theorem, it suffices to verify that

$$
\begin{equation*}
\left|\left(\mu-\lambda_{1}\right)\left(B_{\mid \operatorname{ker} P}-\mu\right)^{-1}\right|_{\mathcal{L}\left(\mathcal{C}\left(\mathbf{S}^{1}\right)\right)} \leq 1 \quad \text { for } \mu>\lambda_{1} \tag{5.35}
\end{equation*}
$$

From [17, Chap. 7, sect. 3], we see that either $\left|\left(B_{\mid \text {ker } P}-\mu\right)^{-1}\right|_{\mathcal{L}\left(\mathcal{C}\left(\mathbf{S}^{1}\right)\right)}$ or $-\mid\left(B_{\mid \text {ker } P}-\right.$ $\mu)\left.^{-1}\right|_{\mathcal{L}\left(\mathcal{C}\left(\mathbf{S}^{1}\right)\right)}$ must be an eigenvalue of $\left(B_{\mid \text {ker } P}-\mu\right)^{-1}$ and that

$$
\sigma\left(\left(B_{\mid \operatorname{ker} P}-\mu\right)^{-1}\right)=\left\{\frac{1}{\lambda_{k}-\mu}\right\}_{k \in \mathbb{N} \backslash\{0\}}
$$

Therefore, recalling that $\left\{\lambda_{k}\right\}_{k \in \mathbb{N} \backslash\{0\}}$ is a decreasing sequence and that $\mu>\lambda_{1}$, we obtain $\left|\left(B_{\mid \text {ker } P}-\mu\right)^{-1}\right|_{\mathcal{L}\left(\mathcal{C}\left(\mathbf{S}^{1}\right)\right)} \leq \frac{1}{\mu-\lambda_{1}}$, and hence $\left|\left(\mu-\lambda_{1}\right)\left(B_{\mid \text {ker } P}-\mu\right)^{-1}\right|_{\mathcal{L}\left(\mathcal{C}\left(\mathbf{S}^{1}\right)\right)} \leq 1$.

Since we have $\rho\left(B_{\mid \operatorname{ker} P}\right) \supset\left(\lambda_{1},+\infty\right)$ and (5.35), the assumptions of the HilleYosida theorem are verified, and we get (5.34).

Finally, an estimate for the quantity $|(1-P)|_{\mathcal{L}\left(\mathcal{C}\left(\mathbf{S}^{1}\right)\right)}$ can be obtained as follows. Let $x \in \mathcal{C}\left(\mathbf{S}^{1}\right)$ with $|x|_{\mathcal{C}\left(\mathbf{S}^{1}\right)}=1$ and recall that, under the standing assumptions, $\beta \in \mathcal{C}\left(\mathbf{S}^{1}\right)_{++}$; then,

$$
|(1-P) x|_{\infty} \leq 1+|\langle x, \beta\rangle||w|_{\infty} \leq 1+|w|_{\infty} \int_{0}^{2 \pi} \beta(\theta) \mathrm{d} \theta
$$

Plugging this estimate into (5.33), together with (5.34), we obtain the claim.
Remark 5.13. Theorem 5.12 deserves some comments. The claim could be proved in the setting $X=\mathrm{L}^{2}\left(\mathbf{S}^{1}\right)$-the setting of $[7,8]$-but it would not have useful consequences from the point of view of the application. Indeed, the norm of $L^{2}\left(\mathbf{S}^{1}\right)$ cannot control pointwise constraints such as the one prescribed by (5.6). Therefore, this result cannot be used to prove that, at least for a suitable set of initial data $K_{0}$, the optimal control $\hat{C}$ for problem (5.7), defined in (5.14), is also optimal for problem (5.5). In the aforementioned references, it is just assumed that $\hat{K}$, the solution to (5.15), verifies the state constraint (5.6). In this way, the authors of [7, 8] automatically have that $\hat{C}$ is optimal for problem (5.5). In other words, this important issue is skipped from the theoretical point of view and verified only in numerical exercises. In a different context (delay equations), this issue is approached from the theoretical point of view in [4] (although the proof has a gap) and in [3], where the argument used to get the result is very involved.

Instead, Theorem 5.12 has an important consequence in our $X=\mathcal{C}\left(\mathbf{S}^{1}\right)$ setting: indeed, if

$$
\begin{equation*}
\underline{w}:=\inf _{\mathbf{S}^{1}} w>0 \quad \text { and } \quad\left(1+|w|_{\infty} \int_{0}^{2 \pi} \beta(\theta) \mathrm{d} \theta\right)\left|K_{0}-\left\langle K_{0}, \beta\right\rangle w\right|_{\infty} \leq\left|\left\langle K_{0}, \beta\right\rangle\right| \underline{w} \tag{5.36}
\end{equation*}
$$

then we deduce that $\hat{K}_{g}^{K_{0}}(t, \cdot)>0$ for all $t \in \mathbb{R}_{+}$, and the optimal control $\hat{C}$ given in (5.14) is such that $\hat{C} \in \mathcal{A}_{++}\left(K_{0}\right)$; hence, by Remark 2.5 , it is optimal for problem (5.5).

We conclude this work by specializing the result above to the case of homogeneous data, i.e., $A(\cdot) \equiv A>0$ and $\eta(\cdot) \equiv 1$. This corresponds to the setting presented in [6], where, as repeatedly said, the problem is embedded in $X=\mathrm{L}^{2}\left(\mathbf{S}^{1}\right)$. In this case the sequence of eigenvalues and eigenvectors of $L$ is fully explicit. In particular, we have that $\lambda_{0}=A, \lambda_{1}=A-\sigma, b_{0}=\frac{1}{\sqrt{2 \pi}} \mathbf{1}_{\mathbf{S}^{1}}$. Rewriting the statement of Theorem 5.12 using $\lambda_{0}, \lambda_{1}$, and $b_{0}$ from above, we get the following corollary.

Corollary 5.14. Let $K_{0} \in X_{++}$, and let $\hat{K}^{K_{0}}$ be the solution to the linear integro-PDE (5.15). Assume that $A(1-\gamma)<\rho<A(1-\gamma)+\sigma \gamma$. Then,
$\left|\hat{K}_{g}^{K_{0}}(t)-\frac{1}{2 \pi} \int_{\mathbf{S}^{1}} K_{0}(\theta) \mathrm{d} \theta\right|_{\mathcal{C}\left(\mathbf{S}^{1}\right)} \leq 2 \mathrm{e}^{-(g-A+\sigma) t}\left|K_{0}-\frac{1}{2 \pi} \int_{\mathbf{S}^{1}} K_{0}(\theta) \mathrm{d} \theta\right|_{\mathcal{C}\left(\mathbf{S}^{1}\right)} \quad \forall t \geq 0$, where $\hat{K}_{g}^{K_{0}}(t):=\mathrm{e}^{-g t} \hat{K}^{K_{0}}(t)$. In particular,

$$
\hat{K}_{g}^{K_{0}}(t) \xrightarrow{t \rightarrow+\infty} \frac{1}{2 \pi} \int_{\mathbf{S}^{1}} K_{0}(\theta) \mathrm{d} \theta \quad \text { in } \mathcal{C}\left(\mathbf{S}^{1}\right),
$$

and if

$$
\begin{equation*}
2\left|K_{0}-\frac{1}{2 \pi} \int_{\mathbf{S}^{1}} K_{0}(\theta) \mathrm{d} \theta\right|_{\mathcal{C}\left(\mathbf{S}^{1}\right)} \leq \frac{1}{2 \pi} \int_{\mathbf{S}^{1}} K_{0}(\theta) \mathrm{d} \theta, \tag{5.38}
\end{equation*}
$$

we have $\hat{K}_{g}(t) \in X_{++}$for all $t \in \mathbb{R}_{+}$. Hence, the optimal control for problem ( $\left.P^{\frac{1}{\sqrt{2 \pi}} \mathbf{1}^{1}}\right)$, i.e.,

$$
\begin{equation*}
\hat{C}(t, \theta)=\mathrm{e}^{g t} \frac{A-g}{2 \pi} \int_{\mathbf{S}^{1}} K_{0}(\xi) \mathrm{d} \xi, \tag{5.39}
\end{equation*}
$$

belongs to $\mathcal{A}_{++}\left(K_{0}\right)$, and therefore by Remark 2.5 it is optimal for problem $(P)$.
Proof. We need only rewrite (5.31) in the present setting to get (5.37). The remaining statements follow immediately (see also Remark 5.13). Note that, in the present setting, the coefficients $\beta_{k}$ defined in (5.21) vanish for all $k \geq 1$; hence,

$$
w=\frac{b_{0}}{\alpha_{0}} .
$$

Therefore, since $b_{0}=\frac{1}{\sqrt{2 \pi}} \mathbf{1}_{\mathbf{S}^{1}}$, we get

$$
\left\langle K_{0}, \beta\right\rangle w=\alpha_{0}\left\langle K_{0}, b_{0}\right\rangle \frac{b_{0}}{\alpha_{0}}=\frac{1}{2 \pi} \int_{\mathbf{S}^{1}} K_{0}(\theta) \mathrm{d} \theta .
$$

Moreover, the constant $M$ appearing in (5.31) satisfies

$$
M=1+|w|_{\infty} \int_{0}^{2 \pi} \beta(\theta) \mathrm{d} \theta=1+\left|\frac{b_{0}}{\alpha_{0}}\right|_{\infty} \int_{0}^{2 \pi} \alpha_{0} b_{0}(\theta) \mathrm{d} \theta=2 .
$$

Finally, recalling that $\lambda_{1}=A-\sigma$, we rewrite (5.31) to get (5.37).
6. Conclusions and potential extensions. In this paper we studied a class of optimal control problems in infinite dimension with a positivity state constraint, motivated by economic applications. With respect to the area of literature dealing with the $L^{2}$ setting, we provided a more general abstract framework that includes the
space of continuous functions. This allows us to treat rigorously the positivity state constraint by introducing a suitable auxiliary optimal control problem and verifying the admissibility of its optimal path for the original problem. As explained in the introduction, this was an issue that in previous works was left out at the theoretical level and checked only numerically (see Remark 5.13). In section 5 we showed, instead, that by formulating the problem in the space of continuous functions it is possible to deal successfully with the positivity state constraint, thanks to the stability results given in section 3.3.

Our approach is based on the particular structure of the problem that allows us to find explicit solutions to the HJB equation of the auxiliary problem, as shown in section 3.2. Indeed, the state equation (2.1) is linear, in (3.6) we choose the bounded and positive linear operator $N$ appearing in (2.1) to be of multiplicative type, and the function $u$ in (3.7) is a power utility function.

Nonetheless, it may be possible to adapt our approach to treat optimal control problems with a positivity state constraint, where nonlinear state equations and/or nonhomogeneous utility functions appear, as in various economic models. To this end, we suggest a possible scheme for future work in the case where a semilinear state equation is considered, i.e., where an additional term $F(x(t))$, depending only on the state variable $x$, appears in the state equation (2.1).
(i) A nonlinear semigroup can be associated to semilinear abstract equations, and it is necessary to impose the fact that it preserves strict positivity, similarly to the requirement of point (i) of Assumption 2.2.
(ii) The pointwise positivity state constraint still makes the HJB equation associated to the optimal control problem very difficult to study. Therefore, we can consider an auxiliary problem in a half-space containing the positive cone of the state space $X$.
(iii) It is now necessary to prove existence, uniqueness, and regularity of the solutions to the HJB equation of the auxiliary optimal control problem: this may be very hard, but there are results of this type in some cases; see, e.g., [28, 26, 27], where the concept of viscosity solutions is used.
(iv) After the previous point is settled, a verification theorem should be established, as we did in Theorem 3.1: its proof does not require finding explicit solutions to the HJB equation of the auxiliary optimal control problem; see again [28, 26, 27].
(v) Finally, one should prove that under suitable assumptions the solution to the auxiliary optimal control problem remains in the positive cone of the state space $X$, as happens in our setting if the a priori estimate (5.36) holds, by virtue of Theorem 5.12 . This could be achieved by extending the stability results given in section 3.3 to the case of nonlinear semigroups.

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[^1]:    ${ }^{1}$ It is worth noting that HJB equations in Banach spaces have been scarcely studied in the literature (see, e.g., Addona, Bandini, and Masiero [1], Fuhrman, Masiero, and Tessitore [31], and Masiero [36, 37] for mild solutions to second order semilinear HJB equations; see Crandall and Lions [18, 19, 20] and Soner [39] for viscosity solutions to first order HJB equations).
    ${ }^{2}$ In most applied examples, $D$ has a topological structure. However, at this abstract stage, this is not needed.

[^2]:    ${ }^{3}$ From an economic perspective, these two properties entail that $u(\theta, \cdot)$ is a utility function for any $\theta \in D$.

