# Approximately Optimal Control of MDPs over a Long Operation-Dependent Time Horizon and Application to Battery Energy Storage Systems 

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#### Abstract

This paper considers a class of optimal control problem for Markov Decision Processes (MDPs), with a time horizon which is both long and operation-dependent. The motivation for this study comes from the optimal management problem for grid-connected Battery Energy Storage Systems (BESSs). As investment represents a large portion of the total cost, and the BESSs can be used for a limited number of cycles only, it is essential to maximize the benefits of using a BESS over its lifetime, which, however, depends on its actual operation. First, we prove that the considered class of optimal control problems can be approximated by the minimization of the ratio of two infinite-horizon average-cost problems. We then characterize the optimal policy in terms of appropriate Bellman-type equations for infinite-horizon average-cost problems. Furthermore, we propose a relative-value-iteration-type algorithm and prove its convergence. Finally, we present some numerical results illustrating the efficiency of the proposed methods for the BESS application. Index Terms- Stochastic optimal control, Markov Decision Processes (MDPs), Energy systems, Optimization algorithms, Battery energy storage systems


## I. Introduction

This paper studies the optimal control of stochastic systems over their lifetime, in the case where the lifetime depends on their actual operation. This study's motivation comes from the study of the optimal operation of grid-connected Battery Energy Storage Systems (BESSs). Batteries are expected to play an essential role in the future power system. However, despite the continuing battery cost reduction, the initial investment costs for BESSs remain high and constitute one of the highest components of BESSs cost. Furthermore, batteries are degrading due to calendar aging and usage (cycling aging), and their lifetime depends on their actual operation. Thus, to optimize the profits from a BESS operation, it is essential to take into account the effects of BESS usage on its lifetime.

[^0]Battery aging is a complex nonlinear phenomenon, and various battery aging models have been proposed [1]. Still, most of the works on optimal management of grid-connected BESSs use simplified models. The articles [2]-[6] assume that there is a certain cost for using the battery, proportional to the initial investment cost. Other papers put certain constraints on battery usage [7]-[9]. Another relevant idea is to use separation of time scales [10], [11] (see also [12]). However, most of the works do not consider the impact of BESS usage on the problem's time horizon. Notable exceptions are [13], [14], which formulate the problem as a Stochastic Shortest Path (SSP) problem. To achieve this, however, the state space should be extended to include the total degradation up to the current time step. This leads to huge state space cardinalities and makes the analysis and optimization of detailed models computationally difficult.

In this work, we study the general MDP optimal control problem with long, operation-dependent time horizon, and approximately reformulate it into the problem of minimizing the ratio of two long-time average-cost criteria. This latter problem dates back to the 1960s, where problems with ratio objectives were studied in the context of maintenance/replacement scheduling [15], [16]. Derman in [15] proves the existence of optimal strategies within the class of memory-less, stationary policies, and reformulates the problem initially as a fractional linear program and then as a linear program. Ratio objectives also appear in the optimal control of Semi-Markov Decision Processes (SMDPs) [17], [18] (see also [19], [20]). There is also a large body of works applying Reinforcement Learning schemes to SMDPs (e.g. [21]-[24]). Recently, similar problems involving a ratio of two long-time average-cost functions were studied in the context of synthesizing systems that achieve an optimal trade-off between a given cost and reward model [25]-[27]. The primary technical difference of this work with the ones mentioned above, besides the fact that we start with a different original problem, is that in all these works, the cost function in the denominator can take only positive (or in some cases non-negative) values. However, there are interesting problems involving long, operation-dependent time horizon, for which the corresponding ratio objective has a denominator that could take both positive and negative values (see, for example, Section VII).

The technical contribution of the current work is the following. First, we show that the optimal control problems with long, operation dependent time horizon can be approximated by ratio cost problems. Second, we extend theory for the
optimal control of cost functions with ratio objectives to the case of indefinite denominators (namely, they could take both positive and negative values), and provide appropriate characterizations. Third, we propose an efficient, model-based algorithm and prove its (exponential) convergence. The last contribution is applying the proposed methods to a BESS example, which has a long horizon, affected by battery degradation.

The rest of this paper is organized as follows: Section II introduces some notation and presents some results of Stochastic Optimal Control theory. In Section III, the optimal control problem with long operation-dependent time horizon is stated and the approximate reduction to a ratio of two long-time average-cost problems is proved. In Section IV, the optimal policy is characterized in terms of a pair of coupled Bellmantype equations. Section V derives a simpler characterization of the optimal solution in terms of a single Bellman equation, presents a relative-value-iteration-type algorithm and proves its convergence. Section VI presents some numerical results, for the BESS application. Finally, Section VIII summarizes the main contributions of the current work.

## II. Background, Notation and Preliminary Results

We consider stochastic dynamics of the form:

$$
\begin{equation*}
x_{k+1}=f\left(x_{k}, u_{k}, w_{k}\right) \tag{1}
\end{equation*}
$$

where $x_{k}$ represents the state variable, $u_{k}$ the control action and $w_{k}$ a stochastic disturbance with known distribution depending only on $x_{k}$. We assume that $x_{k}$ has a finite state space $X=\{1, \ldots, n\}$ and the admissible set of control actions $u_{k}$ is a finite set $U=\left\{1, \ldots, n_{U}\right\}$. The controlled Markov chain notation, $p_{i j}(u)=P\left[x_{k+1}=j \mid x_{k}=i, u_{k}=u\right]$, will be also used. We assume that the distribution of the initial condition $x_{0}$ is known.

A general policy is a function of the past states and actions, and randomizes among the several possible actions. Particularly, we define the history at time $k$ as $\eta_{k}=$ $\left(x_{0}, u_{0}, \ldots, x_{k-1}, u_{k-1}, x_{k}\right)$. A policy is a sequence $\mu=$ $\left(\mu_{0}, \mu_{1}, \mu_{2}, \ldots\right)$, where $\mu_{k}$ is a function that maps the history $\eta_{k}$ to a probability distribution over the feasible actions. This is the most general class of policies considered in this work, and it will be called 'behavioral policies'. A policy $\mu$ is called stationary if each $\mu_{k}$ depends only on the current sate $x_{k}$ and $\mu_{k}=\mu_{0}$. A subset of stationary policies is the set of deterministic stationary policies i.e., policies where $\mu_{k}(x)$ assigns the entire probability to a single element of $U$, for all states $x \in X$. We denote the class of behavioral, stationary, and deterministic stationary policies by $\mathcal{U}_{b}, \mathcal{U}_{s}, \mathcal{U}_{d}$, respectively.

Let us state a couple of assumptions which are very common in the literature (see for example Bertsekas's textbook [19]).
Assumption 1: For the controlled Markov chain (1) it holds:
(a) (Weak Accessibility (WA)): The state space $X$ can be partitioned into two subsets $X_{c}$ and $X_{t}$ such that:

- The states in $X_{t}$ are transient for all stationary policies.
- For every pair of states $i, j \in X_{c}$, there is a stationary policy $\mu$ and an integer $k$ such that, under $\mu$, it holds:

$$
P\left(x_{k}=j \mid x_{0}=i\right)>0 .
$$

(b) (unichain assumption): For any stationary deterministic policy, the closed loop Markov chain has a single irreducible aperiodic ergodic class.
Let us note that the unichain assumption implies the weak accessibility assumption.

There are several possible formulations for the cost functions. This section, deals with the Long-Time Average (LTA) cost (known also as average-cost per stage or ergodic cost). We consider two different costs per stage $g_{1}\left(x_{k}, u_{k}\right)$ and $g_{2}\left(x_{k}, u_{k}\right)$, where the first represents an actual cost per stage and the second a degradation rate. The LTA cost is given by:
$\lambda_{s}(\mu, i)=\lim _{N \rightarrow \infty} \frac{1}{N} E\left[\sum_{k=1}^{N} g_{s}\left(x_{k}, u_{k}\right) \mid x_{0}=i, u_{k}=\mu_{k}\left(x_{k}\right)\right]$,
for $s=1,2$. Note that under Assumption 1(b), for every $\mu \in \mathcal{U}_{s}$, the cost $\lambda_{s}(\mu, i)$ does not depend on $i$. In this case we write $\lambda_{s}(\mu)$ in the place of $\lambda_{s}(\mu, i)$. Furthermore, under Assumption 1.a, the optimal cost for the LTA minimization problem, denoted by $\lambda_{s}^{\star}$, is the same for all initial states.

The optimal policy $\mu^{\star}$ is stationary, and under Assumption 1.a, is characterized by the corresponding Bellman equation:

$$
\begin{equation*}
\lambda_{s}+h_{s}(i)=\min _{u}\left\{g_{s}(i, u)+\sum_{j=1}^{n} p_{i j}(u) h_{s}(j)\right\} \tag{3}
\end{equation*}
$$

for all $i \in X$, where $h_{s}$ is a vector in $\mathbb{R}^{n}$ called bias (e.g. [19], [20]) and $\lambda_{s}$ is the minimum LTA cost.

Following Bertsekas ([19]), let us introduce the 'shorthand' notation for this problem. For a stationary policy $\mu$, denote by $\mathcal{T}_{s, \mu}$ the Bellman operator associated with the LTA cost $\lambda_{s}$ (with $s=1,2$ ), mapping a vector $h_{s} \in \mathbb{R}^{n}$ to a vector $\mathcal{T}_{s, \mu} h_{s} \in \mathbb{R}^{n}$, given by:

$$
\begin{equation*}
\left(\mathcal{T}_{s, \mu} h_{s}\right)(i)=g_{s}(i, \mu(i))+\sum_{j=1}^{n} p_{i j}(\mu(i)) h_{s}(j) \tag{4}
\end{equation*}
$$

If the LTA cost under a stationary policy $\mu$ (not necessarily optimal) is the same for all the initial conditions i.e., $\lambda_{s}(\mu, i)=\lambda_{s, \mu}$, for some real number $\lambda_{s, \mu}$, then it holds (Bellman equation for policy evaluation):

$$
\begin{equation*}
h_{s}+\lambda_{s, \mu} \mathbf{1}=\mathcal{T}_{s, \mu} h_{s} \tag{5}
\end{equation*}
$$

where 1 is the $n$-vector of ones. Conversely if a vector $h_{s} \in$ $\mathbb{R}^{n}$ and a scalar $\lambda_{s, \mu}$ satisfy (5), then the LTA cost under $\mu$ is equal to $\lambda_{s, \mu}$, for every initial state.

## A. Ratio of LTA Cost Problems

We are also interested in optimization problems, considering the ratio of two LTA costs, in the form:

$$
\begin{array}{ll}
\underset{\mu}{\operatorname{minimize}} & \frac{\lambda_{1}(\mu, i)}{\lambda_{2}(\mu, i)}  \tag{6}\\
\text { subject to } & \lambda_{2}(\mu, i)>0
\end{array}
$$

The following proposition gives a geometric characterization of the set of possible values of $\left(\lambda_{1}(\mu, i), \lambda_{2}(\mu, i)\right)$, for the different policies $\mu$. This result will be used to prove that there is a deterministic optimal policy solving (6).

Proposition 1: Under the weak accessibility assumption (Assumption 1.a), there is a convex compact polygon $D$ such that:
(i) Assume that for a policy $\mu \in \mathcal{U}_{b}$ the limits in (2) exist for $s=1,2$. Then, $\left(\lambda_{1}(\mu, i), \lambda_{2}(\mu, i)\right) \in D$, for all $i \in X$.
(ii) For every vertex of $D$, with coordinates $\left(\bar{\lambda}_{1}, \bar{\lambda}_{2}\right)$, there is a deterministic stationary policy $\mu \in \mathcal{U}_{d}$ such that, for all $i \in X, \lambda_{1}(\mu, i)=\bar{\lambda}_{1}$ and $\lambda_{2}(\mu, i)=\bar{\lambda}_{2}$.
(iii) For every point $\left(\bar{\lambda}_{1}, \bar{\lambda}_{2}\right) \in D$ there is a behavioral policy $\mu \in \mathcal{U}_{b}$ such that, for all $i \in X, \lambda_{1}(\mu, i)=\bar{\lambda}_{1}$ and $\lambda_{2}(\mu, i)=\bar{\lambda}_{2}$.
Proof: See Appendix A.
We call $D$ the 'feasible region'. If for a policy $\mu \in \mathcal{U}_{d}$ there is a vertex of $D$ with coordinates $\left(\bar{\lambda}_{1}, \bar{\lambda}_{2}\right)$ such that $\lambda_{1}(\mu, i)=\bar{\lambda}_{1}$ and $\lambda_{2}(\mu, i)=\bar{\lambda}_{2}$, for all $i \in X$, we call $\mu$ a 'corner policy'.

Remark 1: Results similar to Proposition 1 appear in the literature of constrained MDPs, under the unichain assumption (e.g. [29]). The proof here is needed, because we use the more general WA assumption.

## III. Problem Formulation and Approximate Infinite-Horizon Reformulation

In this section we present the general form of the optimal control problem with operation-dependent time horizon, as well as an infinite-horizon approximate reformulation. Consider a dynamics in the form (1). Assume that the system has initially a remaining life denoted by $R$ i.e., $R$ denotes the initial total degradation capacity of the system. During time period $k$, the remaining life is reduced by $g_{2}\left(x_{k}, u_{k}\right)$ (note that $g_{2}$ could be positive or negative). Thus, the time horizon of the problem is random and it is given by the stopping time:

$$
\begin{equation*}
T=\inf \left\{t: \sum_{k=1}^{t} g_{2}\left(x_{k}, u_{k}\right) \geq R\right\} \tag{7}
\end{equation*}
$$

We are interested in solving the problem:

$$
\begin{equation*}
\underset{\mu}{\operatorname{minimize}} J_{a}(\mu)=\lim _{N \rightarrow \infty} E\left[\sum_{k=0}^{T \wedge N-1} g_{1}\left(x_{k}, u_{k}\right)\right] / R \tag{8}
\end{equation*}
$$

where $T \wedge N=\min (T, N), \mu=\left(\mu_{1}, \mu_{2}, \ldots\right)$ and $u_{k}$ is computed according to $\mu_{k}\left(x_{k}\right)$. In the rest of the paper we refer to (8) as the 'original problem'. The cost function $J_{a}$ represents the total cost over the operation-dependent time horizon, scaled by the initial remaining life. Indeed, we expect the total cost $\lim _{N \rightarrow \infty} E\left[\sum_{k=0}^{T \wedge N-1} g_{1}\left(x_{k}, u_{k}\right)\right]$ to be proportional to $R$, for large $R$, and thus we use the scaling by $R$ to make the cost approximately independent of it.

Remark 2: The original Problem (8) can be written as a Stochastic Shortest Path problem. To this end, the system state space should be extended to include the cumulative degradation $y_{k}=\sum_{k^{\prime}=1}^{k} g_{2}\left(x_{k^{\prime}}, u_{k^{\prime}}\right)$. Let $Y$ denote the possible values of $y_{k}$. Then, a natural state space for (8) is
$\bar{X}=X \times Y$. For large $R$ the cardinality of the state space of the problem becomes huge. In the case where $g_{2}(x, u)$ can take negative values, $Y$ becomes countably infinite. In the next subsection, we propose an approximate reformulation of Problem (8) as a ratio of two LTA costs, which allows an efficient solution.

## A. Problem Reformulation

We then perform a series of approximate transformations to the original Problem (8), which eventually leads to a problem with ratio objectives. Assuming that $R$ is very large, we approximate the original problem (8) by its 'deterministic horizon' counterpart:

$$
\underset{T_{d} \in \mathbb{N}}{\operatorname{minimize}}\left\{\begin{array}{ll}
\min _{\mu} & E\left[\sum_{k=1}^{T_{d}} g_{1}\left(x_{k}, u_{k}\right)\right] / R  \tag{9}\\
\text { subject to } & \left|\sum_{k=1}^{T_{d}} E\left[g_{2}\left(x_{k}, u_{k}\right)\right]-R\right| \leq \bar{G}
\end{array}\right\}
$$

where $T_{d}$ is a deterministic decision variable and $\bar{G} \ll R$. The constraint represents the fact that the expected discrete degradation is approximately equal to $R$. The internal minimization of (9) i.e., the corresponding problem with a large fixed $T_{d}$, can be further approximated by its infinite-horizon counterpart. To this end, we use $\lambda_{1}(\mu, i)$ and $\lambda_{2}(\mu, i)$ defined in (2).

The random variable $T$ could take finite or infinite values. However, we are interested in cases where, under the optimal policy, the expected time horizon is finite and appropriate assumptions will be made in the following (Assumption 2) to ensure that this happens. There are two qualitatively distinct cases. In the first case, there is a policy $\mu$ under which the system is producing value i.e., $\lambda_{1}(\mu, i)<0$. In this case we are interested in delaying the stopping time $T$ (keeping the system alive), while enjoying negative costs. In the second case, every policy produces losses i.e, for every policy $\mu$ it holds $\lambda_{1}(\mu, i)>0$. In this case, we are interested in making the time horizon $T$ as short as possible.

In the case where there exists a value producing policy, we assume that if for a policy $\mu$ we have $\lambda_{1}(\mu, i)<0$, then it also holds $\lambda_{2}(\mu, i)>0$. Otherwise, many of the system sample paths would produce $-\infty$ cost. Conversely, if all policies produce losses we assume that there is a policy $\mu$ such that $\lambda_{2}(\mu, i)>0$. Otherwise, all the policies will lead to infinite total cost. We will make thus the following assumption.
Assumption 2: For every initial condition $i \in X$ one of the following is true:
(a) (There exists a value producing policy) There is a policy $\mu$ such that $\lambda_{1}(\mu, i)<0$, and for every policy $\mu^{\prime}$ it holds: either $\lambda_{1}\left(\mu^{\prime}, i\right)>0$ or $\lambda_{2}\left(\mu^{\prime}, i\right)>0$ (see Figure 1.a).
(b) (All policies produce losses) For every policy $\mu$, it holds $\lambda_{1}(\mu, i)>0$. Furthermore, there is a policy $\mu^{\prime}$ such that $\lambda_{2}\left(\mu^{\prime}, i\right)>0$ (see Figure 1.b).


Fig. 1: The dark gray areas in parts (a) and (b) of this figure correspond to the feasible regions $D$, under Case (a) and Case (b) of Assumption 2 respectively. The light gray areas represent areas forbidden by Assumption 2. In both cases we denote by $F$ the point which is optimal for (10) (or equivalently for (11)). If the manager of the system ignores lifetime effects (degradation/regeneration) he/she will choose the point $F^{\prime}$, which is the one minimizing $\lambda_{1}$. In Case (a), it turns out that $F^{\prime}$ gives larger benefits per unit time, but much less time horizon compared with $F$, and thus it turns out to be inefficient. In Case (b), $F^{\prime}$ produces the smallest possible loss per unit time, but it leads to a negative $\lambda_{2}$ and thus on average these losses last forever and the total expected cost is infinite.

The 'deterministic horizon' Problem (9) can be approximately restated as:

$$
\begin{array}{ll}
\underset{T_{d} \in \mathbb{R}_{+}, \mu}{\operatorname{minimize}} & J_{b}(\mu, i)=T_{d} \lambda_{1}(\mu, i) / R \\
\text { subject to } & T_{d} \lambda_{2}(\mu, i)=R  \tag{10}\\
& \lambda_{2}(\mu)>0
\end{array}
$$

for large $R$. Let us call this problem the 'infinite-horizon approximate problem'. Substituting the constraint back to the objective, (10) can be equivalently stated as a 'ratio cost problem':

$$
\begin{array}{ll}
\underset{\mu}{\operatorname{minimize}} & J_{c}(\mu, i)=\lambda_{1}(\mu, i) / \lambda_{2}(\mu, i)  \tag{11}\\
\text { subject to } & \lambda_{2}(\mu)>0
\end{array}
$$

We then show that if an optimal policy for the 'ratio cost problem' is applied to the original Problem (8) it will be $\varepsilon$ optimal.

Theorem 1 (Problem Approximation): Suppose that Assumption 1.a and Assumption 2 hold. Then:
(i) Problem (11) is feasible and there is a corner optimal policy $\mu^{\star} \in \mathcal{U}_{d}$ minimizing $J_{c}$ for all $i$. Furthermore, the optimal value is the same for all $i \in X$.
(ii) For every $\varepsilon>0$, there is an $R_{0}$ such that any corner optimal policy for 'ratio cost problem' (11), is $\varepsilon$-optimal for the original Problem (8), for any $R>R_{0}$. Furthermore, if for a policy $\mu$, the original cost $J_{a}$ is finite, then the expected time horizon $E[T]$ is finite and the original cost may be written equivalently as:

$$
\begin{equation*}
J_{a}=E\left[\sum_{k=0}^{T} g_{1}\left(x_{k}, u_{k}\right)\right] / R \tag{12}
\end{equation*}
$$

Proof: See Appendix B.

Remark 3: This is our last reformulation and (11) is an important topic of the rest of this paper. While (11) has a much smaller state space than (8), dynamic programming is not directly applicable to it. In Sections IV and V we propose techniques to solve (11) efficiently.

Remark 4: Problems involving the ratio of two LTA costs have been already studied in the literature (e.g. [15] - [27]). In these works the authors assume that the cost function in the denominator $g_{2}$ can take only positive (or non-negative) values. Here we generalize the theory to the case where $g_{2}$ can also take negative values.

Remark 5: Case (b) of Assumption 2 can be reduced to case (a), using $\lambda_{1}^{\prime}=-\lambda_{2}$ and $\lambda_{2}^{\prime}=\lambda_{1}$. This transformation corresponds to a $90^{\circ}$ counterclockwise rotation of Figure 1.(b).

## IV. Fixed Point Characterization and the FPRVI Algorithm

In this section, we characterize the corner optimal policies of (11) in terms of a pair of coupled Bellman-type equations (fixed point characterization). This characterization is then used to propose a value iteration type algorithm.

## A. Fixed Point Characterization of the Optimal Policy

Proposition 2 (Fixed Point Characterization): Suppose that Assumption 1.a and Assumption 2.a hold.
(i) Assume that for a policy $\mu^{\star}$ there exist vectors $h_{1}^{\star}, h_{2}^{\star} \in$ $\mathbb{R}^{n}$ and scalars $\lambda_{1}^{\star}, \lambda_{2}^{\star} \in \mathbb{R}$ satisfying the fixed point equations:

$$
\begin{align*}
& \mathcal{T}_{1, \mu^{\star}} h_{1}^{\star}=\lambda_{1}^{\star} \mathbf{1}+h_{1}^{\star},  \tag{13}\\
& \mathcal{T}_{2, \mu^{\star}} h_{2}^{\star}=\lambda_{2}^{\star} \mathbf{1}+h_{2}^{\star},  \tag{14}\\
& \mu^{\star} \in \underset{\mu}{\arg \min }\left[\lambda_{2}^{\star} \mathcal{T}_{1, \mu} h_{1}^{\star}-\lambda_{1}^{\star} \mathcal{T}_{2, \mu} h_{2}^{\star}\right], \tag{15}
\end{align*}
$$

where $1 \in \mathbb{R}^{n}$ is the vector consisting of ones, in (15) the $\arg \mathrm{min}$ is considered for each component separately, and $T_{s, \mu}$ is given in (4). Then, $\mu^{\star}$ is optimal for (11), and the optimal value satisfies $\lambda^{\star}=\lambda_{1}^{\star} / \lambda_{2}^{\star}$.
(ii) Conversely, if $\mu^{\star}$ is a corner optimal policy, then it satisfies (13)-(15).
Proof: Let $\mu^{\star}$ be the optimal policy for (11) i.e., $\lambda_{2}\left(\mu^{\star}\right)>$ 0 and $\lambda_{1}(\mu, i) / \lambda_{2}(\mu, i) \geq \lambda_{1}\left(\mu^{\star}\right) / \lambda_{2}\left(\mu^{\star}\right)$, for every policy $\mu$ with $\lambda_{2}(\mu, i)>0$. By Assumption 2.a, the point $\left(\lambda_{1}\left(\mu^{\star}\right), \lambda_{2}\left(\mu^{\star}\right)\right)$ lies in the interior of the second quadrant (see Figure 2.a) and thus $\lambda_{2}\left(\mu^{\star}\right)>0$. Hence, if $\mu^{\star}$ is optimal then:

$$
\begin{equation*}
\lambda_{2}\left(\mu^{\star}\right) \lambda_{1}(\mu, i)-\lambda_{1}\left(\mu^{\star}\right) \lambda_{2}(\mu, i) \geq 0, \tag{16}
\end{equation*}
$$

for every $\mu$ and $i \in X$ with $\lambda_{2}(\mu, i)>0$. Furthermore, if there is an $i \in X$ and a policy $\mu$ such that $\lambda_{2}(\mu, i) \leq 0$ and $\lambda_{2}\left(\mu^{\star}\right) \lambda_{1}(\mu, i)-\lambda_{1}\left(\mu^{\star}\right) \lambda_{2}(\mu, i)<0$, then the feasible set $D$ intersects the third quadrant which contradicts Assumption 2 (see also Figure 2.a). Thus, if $\mu^{\star}$ is optimal then (16) holds true for all the policies $\mu$.

Conversely, assume that a policy $\mu^{\star}$ satisfies (16). Then from Assumption 2.a we conclude that $\left(\lambda_{1}\left(\mu^{\star}\right), \lambda_{2}\left(\mu^{\star}\right)\right)$ cannot be in the first quadrant, because there are points of $D$ in the second quadrant, and thus (16) cannot hold. We
then show that, $\left(\lambda_{1}\left(\mu^{\star}\right), \lambda_{2}\left(\mu^{\star}\right)\right)$ cannot be in the fourth quadrant. Assume the contrary i.e., $\lambda_{1}\left(\mu^{\star}\right)>0, \lambda_{2}\left(\mu^{\star}\right)<0$. Using Assumption 2.a we consider a point $\left(\lambda_{1}(\mu, i), \lambda_{2}(\mu, i)\right)$ in the second quadrant. Then, take the convex combination $\left[A\left(\lambda_{1}(\mu, i), \lambda_{2}(\mu, i)\right)+(1-A)\left(\lambda_{1}\left(\mu^{\star}\right), \lambda_{2}\left(\mu^{\star}\right)\right)\right] \in D$, where $A=\lambda_{1}\left(\mu^{\star}\right) /\left(\lambda_{1}\left(\mu^{\star}\right)-\lambda_{1}(\mu, i)\right)$. Observe that $A \lambda_{1}(\mu, i)+$ $(1-A) \lambda_{1}\left(\mu^{\star}\right)=0$. Furthermore:
$A \lambda_{2}(\mu, i)+(1-A) \lambda_{2}\left(\mu^{\star}\right)=\frac{\lambda_{1}\left(\mu^{\star}\right) \lambda_{2}(\mu, i)-\lambda_{2}\left(\mu^{\star}\right) \lambda_{1}(\mu, i)}{\lambda_{1}\left(\mu^{\star}\right)-\lambda_{1}(\mu, i)}$.

Thus, since $\mu^{\star}$ satisfies (16) for all $\mu, i$, it holds $A \lambda_{2}(\mu, i)+$ $(1-A) \lambda_{2}\left(\mu^{\star}\right) \leq 0$, which violates assumption 2.a. Finally, Assumption 2.a implies that $\left(\lambda_{1}\left(\mu^{\star}\right), \lambda_{2}\left(\mu^{\star}\right)\right)$ cannot belong to the third quadrant. Hence, $\left(\lambda_{1}\left(\mu^{\star}\right), \lambda_{2}\left(\mu^{\star}\right)\right)$ should belong to the second quadrant.

Dividing (16) by $\lambda_{2}(\mu, i) \lambda_{2}\left(\mu^{\star}\right)$ we conclude that $\lambda_{1}(\mu, i) / \lambda_{2}(\mu, i) \geq \lambda_{1}\left(\mu^{\star}\right) / \lambda_{2}\left(\mu^{\star}\right)$, for every policy $\mu$ with $\lambda_{2}(\mu, i)>0$.

Now (16) holds true for every $\mu$ if and only if:

$$
\begin{equation*}
\mu^{\star} \in \underset{\mu}{\arg \min }\left\{\lambda_{2}\left(\mu^{\star}\right) \lambda_{1}(\mu)-\lambda_{1}\left(\mu^{\star}\right) \lambda_{2}(\mu)\right\} \tag{17}
\end{equation*}
$$

(i) Let $\mu^{\star}$ a policy satisfying (13)-(15). Multiplying (13) by $\lambda_{2}^{\star}$ and subtracting (14) times $\lambda_{1}^{\star}$, we get:

$$
\begin{array}{r}
\lambda_{2}^{\star} h_{1}^{\star}(i)-\lambda_{1}^{\star} h_{2}^{\star}(i)=\lambda_{2}^{\star} g_{1}(i, u)-\lambda_{1}^{\star} g_{2}(i, u)+ \\
\quad+\sum_{j=1}^{n} p_{i j}(u)\left[\lambda_{2}^{\star} h_{1}^{\star}(j)-\lambda_{1}^{\star} h_{2}^{\star}(j)\right], \tag{18}
\end{array}
$$

for $u=\mu^{\star}(i)$. Thus, (18) and (15) imply (17). Hence, $\mu^{\star}$ is optimal.
(ii) Using Proposition 1.(ii) and Theorem 1 we conclude that there is an optimal policy $\mu^{\star}$ for (11) such that $\lambda_{1}\left(\mu^{\star}, i\right)$ and $\lambda_{2}\left(\mu^{\star}, i\right)$ do not depend on $i$. Then, using (5) and (17) we get (15).

Remark 6: Proposition 2 characterizes corner optimal policies. It is possible that there exist deterministic stationary optimal policies which are not corner. The following proposition shows that this does not happen in the generic case.

Proposition 3 (Non-Corner Optimal Policies are NonGeneric): Suppose that Assumption 1.a holds. Then, the set of vectors of the form $\left[g_{s}(x, u): x \in X, u \in U, s=1,2\right] \in$ $\mathbb{R}^{2 n n_{U}}$ such that there is a non-corner stationary deterministic optimal policy for (11) has Lebesgue measure zero.
Proof: See Appendix C.

## B. The FPRVI Algorithm

The fixed point characterization (13)-(15) can be used to derive a relative-value-iteration-type algorithm for Problem (11).


Fig. 2: The possible positions for $\left(\lambda_{1}\left(\mu^{\star}\right), \lambda_{2}\left(\mu^{\star}\right)\right)$. The gray area in these figures correspond to the feasible region $D$. The three different pictures illustrate three possible positions of $\left(\lambda_{1}\left(\mu^{\star}\right), \lambda_{2}\left(\mu^{\star}\right)\right)$. The arrows correspond to the minus gradient of the function $f\left(\lambda_{1}(\mu), \lambda_{2}(\mu)\right)=\lambda_{2}\left(\mu^{\star}\right) \lambda_{1}(\mu)-$ $\lambda_{1}\left(\mu^{\star}\right) \lambda_{2}(\mu)$. Part (a) of the figure corresponds to the correct position of the optimum. Parts (b) and (c) correspond to other candidate solutions in the first and fourth quadrant, which however cannot satisfy (16).

Starting from some initial guesses $\lambda_{1}^{0}, \lambda_{2}^{0} \in \mathbb{R}, h_{1}^{0}, h_{2}^{0} \in \mathbb{R}^{n}$ and an initial policy $\mu^{0}$, the iteration is given by:

$$
\begin{align*}
\lambda_{1}^{k+1}= & \left(\mathcal{T}_{1, \mu^{k}} h_{1}^{k}\right)(n)  \tag{19}\\
\lambda_{2}^{k+1}= & \left(\mathcal{T}_{2, \mu^{k}} h_{2}^{k}\right)(n)  \tag{20}\\
h_{1}^{k+1}(i)= & \left(\mathcal{T}_{1, \mu^{k}} h_{1}^{k}\right)(i)-\lambda_{1}^{k+1}  \tag{21}\\
h_{2}^{k+1}(i)= & \left(\mathcal{T}_{2, \mu^{k}} h_{2}^{k}\right)(i)-\lambda_{2}^{k+1}  \tag{22}\\
\mu^{k+1}(i)= & \arg \min _{u}\left[\lambda_{2}^{k} g_{1}(i, u)-\lambda_{1}^{k} g_{2}(i, u)+\right. \\
& \left.+\sum_{j=1}^{n} p_{i j}(u)\left(\lambda_{2}^{k} h_{1}^{k}(j)-\lambda_{1}^{k} h_{2}^{k}(j)\right)\right] \tag{23}
\end{align*}
$$

where $i=1, \ldots, n$ and $n$ is the last element of the state space $X$. We will refer to this algorithm as the Fixed Point Relative Value Iteration (FPRVI) algorithm. It is not difficult to see that a fixed point of (19)-(23) satisfies (13)-(15).

The algorithm is intuitively easy to derive, does not have any tuning parameters and works well for the numerical example analyzed (Section VI). However, the analysis of this algorithm is difficult. In the following section, we propose an alternative characterization in terms of a single Bellman equation, which leads to an algorithm with simpler convergence analysis. A comparison of the two algorithms is presented in Remark 8.

## V. Bellman Equation Characterization and the RVI Algorithm

## A. Bellman Equation Characterization of the Optimal Policy

In this section, we characterize the optimal policy and optimal value in terms of a single Bellman-type equation. This characterization leads to a simple value iteration type algorithm, described in the next subsection.

Proposition 4 (Bellman Equation Characterization): Suppose that Assumption 1.a and Assumption 2.a hold. Then:
(i) If $\lambda^{\star}$ is the optimal value for the ratio cost function (11) then:

$$
\begin{equation*}
\min _{\mu}\left[\lambda_{1}(\mu)-\lambda^{\star} \lambda_{2}(\mu)\right]=0 \tag{24}
\end{equation*}
$$

(ii) There are at most two values of $\lambda^{\star}$ satisfying (24).
(iii) A policy $\mu^{\star}$ is optimal for (11) if and only if:

$$
\begin{equation*}
\mu^{\star} \in \underset{\mu}{\arg \min }\left[\lambda_{1}(\mu)-\lambda^{\star} \lambda_{2}(\mu)\right] \tag{25}
\end{equation*}
$$

where $\lambda^{\star}$ is the optimal value.
(iv) If $\lambda^{\star}$ is the optimal value for (11), then there exists a vector $h^{\star} \in \mathbb{R}^{n}$ satisfying the Bellman equation:

$$
\begin{equation*}
h^{\star}(i)=\min _{u}\left[g_{1}(i, u)-\lambda^{\star} g_{2}(i, u)+\sum_{j=1}^{n} h^{\star}(j) p_{i j}(u)\right] . \tag{26}
\end{equation*}
$$

Furthermore, among all the vectors $h$ satisfying (26), there is a unique one with $h(n)=0$. Conversely, assume that a vector $h^{\star}$ and a scalar $\lambda^{\star}$ satisfy the Bellman equation (26), along with $h(n)=0$ and denote by $\mu^{\star}$ the policy attaining the minimum in (26). If additionally $\lambda_{2}\left(\mu^{\star}\right)>0$, then $\lambda^{\star}$ is the optimal value and $\mu^{\star}$ is the optimal policy.
Proof: (i) Assume that $\lambda^{\star}, \mu^{\star}$ is the optimal value-optimal policy pair for (11). Then $\lambda^{\star}=\lambda_{1}\left(\mu^{\star}\right) / \lambda_{2}\left(\mu^{\star}\right)$. Equation (16) and the fact that $\lambda_{2}\left(\mu^{\star}\right)>0$ imply that $\min _{\mu}\left[\lambda_{1}(\mu)-\right.$ $\left.\lambda^{\star} \lambda_{2}(\mu)\right] \geq 0$. Observe that substituting $\mu^{\star}$ in the place of $\mu$ in the left-hand side of (16) we get 0 . Thus, (24) holds true.
(ii) Consider the feasible region, illustrated in Figure 1.a. Due to the fact that the feasible region $D$ is a two-dimensional compact convex set not containing the origin, there are at most two lines passing through the origin which support the feasible region (particularly the lines in Figure 2, (a) and (c)).
(iii) Observe that (17) and (25) are equivalent.
(iv) The direct part follows from (i). To prove the converse, consider a scalar $\lambda^{\star}$ and a vector $h^{\star}$ satisfying (26). Consider the problem of minimizing the LTA cost criterion:

$$
\begin{aligned}
& \lambda_{1}(\mu)-\lambda^{\star} \cdot \lambda_{2}(\mu)= \\
& \quad=E\left[\lim _{N \rightarrow \infty} \frac{1}{N} \sum_{k=1}^{N}\left(g_{1}\left(x_{k}, \mu\left(x_{k}\right)\right)-\lambda^{\star} g_{2}\left(x_{k}, \mu\left(x_{k}\right)\right)\right)\right] .
\end{aligned}
$$

For this criterion there is a unique pair of a scalar $\theta$ and a vector $h(1), \ldots, h(n)$ with $h(n)=0$ such that:

$$
\theta+h(i)=\min _{u}\left[g_{1}(i, u)-\lambda^{\star} g_{2}(i, u)+\sum_{j=1}^{n} h(j) p_{i j}(u)\right]
$$

for $i=1, \ldots, n$. Uniqueness and (26) imply that $\theta=0$ and thus (24) holds true. Inequality $\lambda_{2}\left(\mu^{\star}\right)>0$ and (ii) imply that $\lambda^{\star}$ is the optimal value. The characterization of the optimal policy follows from (iii).

Corollary 1: Suppose that Assumption 1.a and Assumption 2.b hold. Then (i)-(iv) of Proposition 4 hold true.

## B. A Relative Value Iteration Algorithm

In this section, we propose a Relative Value Iteration (RVI) algorithm for the problem, based on the single Bellman
equation characterization (26) and prove its convergence. The iteration is given by:

$$
\begin{align*}
\theta^{k+1} & =\min _{u}\left[g_{1}(n, u)-\lambda^{k} g_{2}(n, u)+\tau \sum_{j=1}^{n} h^{k}(j) p_{n j}(u)\right],  \tag{27}\\
h^{k+1}(i) & =\min _{u}\left[g_{1}(i, u)-\lambda^{k} g_{2}(i, u)+\tau \sum_{j=1}^{n} h^{k}(j) p_{i j}(u)\right]- \\
\lambda^{k+1} & =\lambda^{k}+\gamma \theta^{k+1},
\end{align*}
$$

where the initial conditions are $\lambda^{0}=0$ and $h^{0}=0, \gamma$ is a positive parameter (step size), $0<\tau \leq 1$ is a damping factor, $i=1, \ldots, n$, and $n$ is the last element of the state space $X$. The idea of the algorithm is the following. Let us assume that $\lambda$ is 'frozen' (is held fixed). Then (27)-(28) corresponds to a modified version of the relative-value-iteration-type scheme for the cost function:

$$
\begin{equation*}
E\left[\lim _{N \rightarrow \infty} \frac{1}{N} \sum_{k=1}^{N}\left[g_{1}\left(x_{k}, u_{k}\right)-\lambda g_{2}\left(x_{k}, u_{k}\right)\right]\right] \tag{30}
\end{equation*}
$$

The parameter $\tau$ is used to avoid cycling (periodic behavior) in (27)-(28) and ensure convergence (for $\tau=1$ we retrieve the usual relative value iteration scheme). Under (27)-(28), the value of $\theta^{k}$ converges to the minimum value of (30). On the other hand, it will be shown that the optimal value of (30) is strictly decreasing in $\lambda$, for $\lambda$ in an appropriate interval. Thus, allowing $\lambda$ to vary according to (29), if $\theta^{k}$ is positive, then $\lambda$ will increase and will eventually lead to a lower value of $\theta$. The following proposition shows the convergence of the algorithm.

Theorem 2 (Convergence of the Algorithm): Suppose that the Assumptions 1.b and 2.a hold and that $\tau<1$. Then, there is a value $\bar{\gamma}$ such that if $\gamma \leq \bar{\gamma}$, then the $\left(h^{k}, \lambda^{k}\right)$ converges to $\left(h^{\star} / \tau, \lambda^{\star}\right)$ at the rate of a geometric progression.
Proof: See Appendix D. The proof is based on ideas of singular perturbations (e.g. [30])

Remark 7: There are two basic differences of the RVI algorithm described in this section and the Reinforcement Learning algorithms for Semi-Markov Decision Processes (e.g. [21]- [24]). First this algorithm can handle denominators taking both positive and negative values. Second, since it does not use stochastic approximation it is able to achieve a faster (geometric) convergence rate. Despite the fact that this change of the convergence rate is intuitively trivial, the proof of convergence is quite different. With vanishing step-size the ODE method applies. On the other hand, here we develop a form of discrete time singular perturbation stability analysis.

Remark 8: Each iteration of the RVI algorithm ((27)-(29)) takes slightly less time and about half of the memory compared to the FPRVI algorithm (Section IV-B). Furthermore, RVI has guaranteed convergence. On the other hand, FPRVI computes both $\lambda_{1}$ and $\lambda_{2}$ and does not need any tuning.

## VI. Battery Energy Storage System Application

Consider a BESS connected to the grid ${ }^{1}$. The discrete time dynamics for the State of Energy (SoE) of the battery is given by:

$$
\begin{equation*}
x_{k+1}^{1}=x_{k}^{1}+u_{k} \tag{31}
\end{equation*}
$$

where $x_{k}^{1} \in[0,1]$ is the SoE i.e., the fraction of the energy stored in the BESS over the maximum possible stored energy, $u_{k} \in\left[-u^{M}, u^{M}\right]$ is the charging or discharging power.

At the beginning of each time period, before the decision $u_{k}$ is made, the grid operator sends a vector signal $x_{k}^{2}=$ $\left[l_{k} p_{k} d_{k}\right]^{T}$ to the BESS. In this vector, $l_{k}$ represents the power which the system requires from the battery to absorb or release (hereafter requested power), $p_{k}$ the price at which the battery owner will be paid for the energy exchange, and $d_{k}$ the penalty for deviating. The cost (which corresponds to the minus profit) for time period $k$ is given by:

$$
g_{1}\left(x_{k}^{2}, u_{k}\right)=d_{k}\left|u_{k}-l_{k}\right|-p_{k} u_{k}
$$

We assume that $p_{k}$ has the same sign with $l_{k}$ and that $d_{k}>0$. Thus the battery owner always has the motivation to use a $u_{k}$ such that $u_{k} l_{k} \geq 0$ i.e., to follow the sign of the grid operator signal. The cost $g_{1}\left(x_{k}^{2}, u_{k}\right)$ corresponds to the deviation penalty minus the revenues the BESS owner receives from the energy exchange. The vectors $x_{k}^{2}$ are modeled as i.i.d random variables (that is $x_{k}^{2}$ and $x_{k^{\prime}}^{2}$ are independent for $k \neq k^{\prime}$ and have the same distribution). For the consistency of the notation we write the evolution of the state variable $x^{2}$ as $x_{k+1}^{2}=w_{k}$. The system at time step $k$ has state variable $x_{k}=\left[x_{k}^{1}\left(x_{k}^{2}\right)^{T}\right]^{T}$ and the dynamics is written in compact form as $x_{k+1}=f\left(x_{k}, u_{k}, w_{k}\right)$.

In this example, we use a very simple Ah throughput model (e.g. [34]) to predict the battery lifetime. Particularly, we assume that after time interval $k$, the remaining useful life of the BESS is reduced by an amount of $g_{2}\left(x_{k}, u_{k}\right)=c_{1}+c_{2}\left|u_{k}\right|$. The first term corresponds to calendar aging and is independent of the battery use, while the second term corresponds to cycling aging. The lifetime of the BESS is thus given by:

$$
T=\min _{t}\left\{t: \sum_{k=1}^{t} g_{2}\left(x_{k}, u_{k}\right) \geq R\right\}
$$

where $R$ is a large positive constant denoting the initial remaining useful life of the battery. The problem of maximizing the revenues of the BESS over its lifetime can be stated in a form similar to (8):

$$
\begin{equation*}
\underset{\mu}{\operatorname{minimize}} \quad J_{a}=\lim _{N \rightarrow \infty} E\left[\sum_{k=1}^{T \wedge N-1} g_{1}\left(x_{k}, u_{k}\right)\right] / R \tag{32}
\end{equation*}
$$

where $\mu=\left(\mu_{1}, \ldots\right)$ and $u_{k}=\mu_{k}\left(x_{k}\right)$.
Let us observe that, if the SoE of the battery permits, the battery manager can always choose an action $u_{k}$ making the instant cost $g_{1}$ non-positive. Thus, it is in the best interest of the battery manager to try to make the lifetime $T$ long, while continuing profiting from its use. Since (31) is controllable and

[^1]$w_{k}$ is i.i.d., this model satisfies Assumption 1.a. Furthermore, since $g_{2}(u) \geq 0$ for all $u$, the model satisfies also assumption 2.a.

Remark 9: More complex degradation models, depending non-linearly on the power, on the state of charge of the battery, or its temperature can be considered. Furthermore, the requested energy, and the price may have some more complicated stochastic dynamics, depending probably on the time of the day. In this section, however, we use the simplest possible models to illustrate the application of the methods developed in the previous sections.

We then apply both FPRVI and RVI algorithms to this problem, compare their performance, compute the optimal policy and then study its sensitivity to the aging parameters.

## A. Optimal Control Computation

Let us first specify some parameters used in the numerical computations. The SoE $x_{k}^{1}$ is bounded in $[0,1]$, the maximum absolute power $u^{M}$ is 0.1 , the calendar aging constant $c_{1}$ is 0.01 and the cycling aging constant $c_{2}$ is 1 (the cycling aging is dominant). For simplicity, we assume that the requested power $l_{k}$ is distributed uniformly in $[-0.1,0.1]$ and that the price per unit energy $p_{k}$ and the deviation penalty are given by $p_{k}=100 l_{k}$ and $d_{k}=1.2 p_{k}+0.01$ respectively. This choice reduces the dimensionality of the state variable $x^{2}$ to one. We choose to discretize the SoE part of the state variable into 101 points and the $x^{2}$ part into 21 points. Thus, the state space has 2121 points.

We implemented both algorithms in Julia 0.6.2. Figure 3 compares their speed of convergence. Particularly, the vertical axis of the figure corresponds to the logarithm of the relative difference of the consecutive iterates of the algorithms. For RVI this corresponds to the quantity:

$$
\log _{10} \frac{\left\|h^{k+1}-h^{k}\right\|}{\left\|h^{k}\right\|}
$$

and for FPRVI to:

$$
\log _{10} \frac{\left\|\left[h_{1}^{k+1} h_{2}^{k+1}\right]-\left[h_{1}^{k} h_{2}^{k}\right]\right\|}{\left\|\left[h_{1}^{k} h_{2}^{k}\right]\right\|}
$$

Both algorithms converge linearly (geometrically). The runs were performed in an Intel Core i5 CPU $750 @ 2.67 \mathrm{GHz}, 4 \mathrm{~GB}$ RAM desktop PC. The run time ${ }^{2}$ for 200 iterations of RVI algorithm is 1.8 s , while 200 iterations of the FPRVI algorithm (Section IV-B) take 1.9s. The parameters for the RVI algorithm are $\tau=1$ (there is no need to use $\tau<1$, because there is no periodicity) and $\gamma=2$ (experimentally tuned).

Figure 4 illustrates the optimal control law, which is piecewise linear with respect to the SoE. It is interesting that at some combinations of SoE with requested power the optimal control is to absorb exactly the power requested i.e., to use $u_{k}=l_{k}$. In other combinations it is optimal to use zero control, in anticipation of possible future losses and battery aging. For example, for $l=0.05$ and for a $\operatorname{SoE} 80 \%$ the

[^2]

Fig. 3: The convergence of the suggested methods. The vertical axis represents the logarithm of the relative difference of two consecutive iterations. The solid line corresponds to the FPRVI algorithm (Section IV) and the dashed line to RVI algorithm. Numerical results suggest that both algorithms have a linear (geometric) convergence rate.


Fig. 4: The optimal control law for the battery Problem (32). The horizontal axis is the State of Energy $x^{1}$ and the different lines correspond to different values of the requested power $l$. We may observe that there is a symmetry around the point $(0.5,0)$. Particularly, for $l^{\prime}=-l$ and $x^{1^{\prime}}=1-x^{1}$ we have $\mu^{\star}\left(x^{1^{\prime}}, l^{\prime}\right)=-\mu^{\star}\left(x^{1}, l\right)$.
optimal action is to absorb no power. If the battery absorbed some energy then it would increase the probability that, at a subsequent time step, when the deviation penalty would probably be higher, it would not have the ability to follow the signal $l$, while at the same time it would certainly (with probability 1) contribute to the battery aging at the current time step.

We then show experimentally that, for a reasonable value of $R$, the horizon is approximately deterministic. Assume that $R=6000$, which corresponds to a battery capable of 3000 full cycles (a reasonable number of cycles for a stationary lithium ion battery storage system). For the optimal policy $\mu^{\star}$, the expected time horizon is $R / \lambda_{2}\left(\mu^{\star}\right)=103294$ time steps. On the other hand, running 100 Monte-Carlo simulations using the optimal policy and uniform initial conditions, we get a mean time horizon of 103298 steps and maximum absolute deviation of 554 steps or $0.53 \%$. The root mean square deviation is 192 steps or $0.18 \%$. These results agree with the intuition of the approximately deterministic time horizon.

## B. Quantifying Approximate Optimality

We then compare the value of the cost $J_{a}$ under the approximately optimal policy, $J_{a}\left(\mu^{\star}, x^{1}, x^{2}\right)$, computed in the previous subsection, to the minimum value of $J_{a}$.


Fig. 5: The difference between the cost of the approximately optimal policy $\mu^{\star}$ and the optimal policy $\bar{\mu}$, computed using the SSP reformulation.

Using Remark 2, we first transform the Problem (32) into a Stochastic Shortest Path (SSP) problem. Assume that all the parameters are as in the previous subsection. Consider an auxiliary state variable $y_{k}=\sum_{t=0}^{k} g_{2}\left(u_{k}\right)$. Since all the possible values of $g_{2}\left(u_{k}\right)$ are integer multiples of 0.01 , the state space for $y_{k}$ is $Y=\{0,0.01,0.02 \ldots, R\}$.

To simplify the numerical computations, we downscale the problem to $R=600$. Thus, $Y$ has 60001 elements and the extended state space $\bar{X}=X \times Y$ has $127,262,121$ elements. The Bellman equation for this problem is:

$$
\begin{align*}
& J_{\mathrm{SSP}}^{\star}\left(y, x^{1}, x^{2}\right)=\min _{u}\left\{\frac{g_{1}(x, u)}{R} I_{y+g_{2}(u)<R^{+}}\right. \\
& \left.\quad+\frac{1}{21} \sum_{x^{\prime 2}=1}^{21} J_{\mathrm{SSP}}^{\star}\left(y+g_{2}(u), x^{1}+u_{k}, x^{\prime 2}\right)\right\} \tag{33}
\end{align*}
$$

where $I$. is the indicator (characteristic) function. Furthermore, $J_{\mathrm{SSP}}^{\star}\left(y, x^{1}, x^{2}\right)=0$, for $y \geq R$. The Bellman equation has a special structure. Particularly, since always $g_{2}(u)>$ 0 , the value of $J_{\text {SSP }}^{\star}\left(y, x^{1}, x^{2}\right)$ depends only on values of $J_{\text {SSP }}^{\star}\left(y^{\prime}, x^{11}, x^{2}\right)$ with $y^{\prime}>y$. This allows us to solve (33) from $y=R$ down to $y=0$. Let us note that a similar property was used in [13].

Denote by $\bar{\mu}$ the optimal policy computed for the SSP problem. The optimal value $J_{a}^{\star}\left(x^{1}, x^{2}\right)$ of $J_{a}$ is equal to $J_{\mathrm{SSP}}^{\star}\left(0, x^{1}, x^{2}\right)$. Figure 5 illustrates the relative difference:

$$
\Delta=\frac{\max _{x^{1}, x^{2}}\left\{J_{a}\left(\mu^{\star}, x^{1}, x^{2}\right)-J_{a}\left(\bar{\mu}, x^{1}, x^{2}\right)\right\}}{\min _{x^{1}, x^{2}}\left|J_{a}\left(\bar{\mu}, x^{1}, x^{2}\right)\right|}
$$

for several values of $R$. For $R=600$ the value of the relative difference $\Delta$ is equal to $3 \times 10^{-5}$. Therefore, if instead of the approximately optimal policy $\mu^{\star}$ we used the actual optimal policy $\bar{\mu}$ the increase in our profits would be $0.003 \%$.

Now regarding the computation cost, each iteration of (33), requires almost the same number of operations as the algorithm (27)-(29). For (27)-(29) we need in this example 200 iterations, whereas for (33) we need 60000 iterations (for $R=6000$, as in the previous subsection, we would need 600,000 iterations). Thus, algorithm (27)-(29) needs almost 300 times less operations than (27)-(29) (3000 times, for $R=6000$ ).


Fig. 6: The optimal control law for different values of $c$. The value of the optimal controls, for negative $l$ can be deduced from the symmetry observation in Figure 4.

## C. Sensitivity to the Aging Parameters

Another interesting question is how the optimal control law changes when the proportion of the calendar aging over the cycling aging varies. We run the same optimization problem with parameters $c_{1}=c$ and $c_{2}=1-c$, for various values of $c$. The results are illustrated in Figure 6. For $c=1$, the optimal control law is identical to the one we would obtain by ignoring the aging and optimizing the long-term average cost $\lambda_{1}$. Indeed, calendar aging is independent of the control actions. As $c$ becomes smaller the cycling aging starts to dominate. We may observe from the numerical results that for smaller $c$ the optimal control action has always smaller absolute value.

Figure 7 illustrates the optimal control law in the form of a heat map, for three different values of $c$. For a small (in absolute value) requested power and for an almost full or an almost empty battery, we observe an 'opportunistic' charging or discharging. This behavior becomes more important, as the weight of the calendar aging becomes larger (or equivalently as we care less about aging). Furthermore, for small $c$, we observe that there is a region where the BESS uses $u=0$ i.e., does not charge or discharge, despite the existence of a non-zero signal $l$. This turquoise ( $u=0$ ) region shrinks and eventually disappears as the calendar aging becomes dominant. Another region is the upper right and lower left part of the heat maps, where the BESS is not able to follow the signal $l$ and thus we have saturation. Finally, there is a linear region where the BESS follows exactly the requested power signal.

For this sensitivity analysis we used a finer discretization with 201201 points to produce smoothly varying heat maps. Particularly we discretized $x^{1}$ variable in 1001 points and $x^{2}$ in 201 points.

## D. Comparing of the Proposed Solution to MPC

We then compare the control law computed in the previous subsections to a Stochastic Model Predictive Control (SMPC).


Fig. 7: The optimal control law for different values of $c$. Particularly, the first heat map corresponds to the case where cycling aging is dominant, the third plot to the case where the calendar aging is dominant (equivalently we do not care about aging) and the second to a case in between.

The control action $u_{k}$ at time step $k$ solves the problem:

$$
\begin{equation*}
\underset{\mu}{\operatorname{minimize}}\left[\sum_{k^{\prime}=k}^{k+N_{0}} g_{1}\left(x_{k}, u_{k}\right) \mid u_{k^{\prime}}=\mu_{k^{\prime}}\left(x_{k^{\prime}}\right)\right] \tag{34}
\end{equation*}
$$

where the optimization horizon $N_{0}$ may take several possible values. The parameters are as in subsection VI-A.

It is interesting that, if $N_{0}=1$, the control law becomes:
that is, $u_{k}$ follows the requested signal $l_{k}$ as close as possible, without violating the battery constraints. We call this the 'myopic' control law.

Figure 8 illustrates the cost $J_{a}$ for the SMPC control laws for several values of $N_{0}$. Interestingly, the cost $J_{a}$ is not decreasing as the horizon increases. This phenomenon occurs because (34) does not take into account the effect of the control law to the problem horizon.

The use of the approximately optimal control law increases the profits by $6.21 \%$ when compared with the 'myopic' control law (35), by $4.02 \%$ when compared to the SMPC with $N_{0}=$ 3 , which happens to result the optimal cost among SMPC controllers, and by $11.78 \%$ when compared with the SMPC with long horizon $N_{0}$.

## VII. Examples with Indefinite Denominator

In this section, we present briefly a couple of examples where the denominator $g_{2}$ can take both positive and negative values.

1) Maximum Range of an Electric Vehicle: Consider an electric vehicle with regenerative breaking and let us study the problem of designing a controller to achieve the longest possible distance, before recharging the battery. In contrast


Fig. 8: The cost of SMPC law for several values of the optimization horizon $N_{0}$, compared to the cost of the approximately optimal control law obtained in Subsection VI-A.
with the BESS example, here the system manager (driver) can choose actions that extend the remaining life of the system, for example he/she can use regenerative breaking in a downhill road. The velocity of the vehicle $x_{k}^{1}$ evolves according to Newton's second law: $x_{k+1}^{1}=x_{k}^{1}+u_{k}-F_{\text {Slope }}\left(x_{k}^{2}\right)-F_{\text {Friction }}\left(x_{k}^{1}\right)$, where we assume that the vehicle has unit mass, $F_{\text {Slope }}\left(x_{k}^{2}\right)$ denotes the force to the vehicle due to the road inclination, $x_{k}^{2}$ is the slope of the road, $F_{\text {Friction }}\left(x_{k}^{1}\right)$ is the force due to friction and $u_{k}$ is the force (from the road to the vehicle) due to the use of the accelerator or the break. We model the slope of the road $x_{k}^{2}$ as a sequence of dependent random variables evolving according to: $x_{k+1}^{2}=\left(1-\alpha x_{k}^{1}\right) x_{k}^{2}+\alpha x_{k}^{1} w_{k}$, where $w_{k}$ are i.i.d. random variables (not necessarily zero-mean) and $\alpha$ is a small positive constant. Let us denote the energy consumption per time step by $g_{2}\left(x_{k}, u_{k}\right)$.

If the mean slope of the road is non-negative, then we expect that the vehicle will consume eventually all the energy stored in its battery. The time, at which this will happen, is given by: $T=\inf _{t}\left\{t: \sum_{k=1}^{t} g_{2}\left(x_{k}, u_{k}\right) \geq R\right\}$, where $R$ is the amount of energy initially stored in the battery of the vehicle. The maximum range problem corresponds to the minimization of the following quantity: $\underset{u_{k}=\mu_{k}\left(x_{k}\right)}{\operatorname{minimize}}-E\left[\sum_{k=1}^{T} x_{k}^{1}\right]$.
2) Selling the Stock of a Non-Viable Retail Firm: Consider a small retail business, which does not have any profits for a long time, and decides to exit the market. This example studies the problem of optimizing the revenue from selling the remaining stock $R$ of products. Assume that if the business sets a price $u_{k}^{1} \in[0,1]$ at time $k$, then the number of products it sells per time step is given by $d\left(u_{k}^{1}\right)$. Furthermore, assume that at time step $k$, the business is able to buy products at a price $x_{k}$ which is random. However, as far as the firm is working, it has to pay a constant amount of money $c$ per unit time (e.g. rent). Denoting by $u_{k}^{2} \in[0,1]$ the number of products the business buys at time step $k$, the cost per stage at time step $k$ is given by $g_{1}\left(x_{k}, u_{k}^{1}, u_{k}^{2}\right)=x_{k} u_{k}^{2}-d\left(u_{k}^{1}\right) u_{k}^{1}+c$ and the stock after step $k$ is reduced by $g_{2}\left(x_{k}, u_{k}^{1}, u_{k}^{2}\right)=d\left(u_{k}^{1}\right)-u_{k}^{2}$. Therefore, the problem of optimizing firm's revenues, until the stock is over has the same structure of the previous problems.

Remark 10: In this problem, we have assumed that the
business is not viable. However, it is probably profitable to buy new products, if the price is low. There are policies which make the remaining stock increase in the long run. However, these policies produce long-term losses.

## VIII. CONCLUSION

The paper considered the problem of the optimal control of MDPs over a long and operation-dependent time horizon and reduced it approximately to a problem involving the minimization of the ratio of two long-time average-costs. The optimal control law was characterized by appropriate sets of Bellman-type equations. Particularly, two characterizations were given. In the first one (fixed point characterization), a control law is optimal if and only if there exist two vectors satisfying a pair of coupled Bellman equations. In the second characterization, a necessary condition is given in terms of a single Bellman equation. This condition also becomes sufficient if an additional inequality holds true. Based on each characterization, we proposed a value iteration-based algorithm (i.e. FPRVI and RVI algorithms). For the RVI algorithm we also proved its convergence.
The proposed techniques were then applied to the problem of the optimal management of a Battery Energy Storage System. The results show that the optimal control law consists of different behaviors, whose boundaries move according to the aging parameters. In the future it would be interesting to analyze more complex models for battery degradation or more accurate stochastic models to describe the grid operator signal.

## APPENDIX

## A. Proof of Proposition 1

Consider the set of points in $\mathbb{R}^{2}$ :

$$
D_{0}=\left\{\left(\lambda_{1}(\mu, i), \lambda_{2}(\mu, i)\right): \mu \in \mathcal{U}_{d}, i \in X\right\}
$$

and denote by $D$ its convex hull. Since the set of deterministic policies $\mathcal{U}_{d}$ is finite, $D_{0}$ is finite, and $D$ is a compact and convex.
(i) We use contradiction. Assume that there is a behavioral policy $\mu$, such that for some $i$, it holds $\left(\lambda_{1}(\mu, i), \lambda_{2}(\mu, i)\right) \notin$ $D$. Then, there is a line strictly separating $D$ and $\left(\lambda_{1}(\mu, i), \lambda_{2}(\mu, i)\right)$. Thus, there are constants $c_{1}, c_{2}$ such that:

$$
\begin{equation*}
c_{1} \lambda_{1}(\mu, i)+c_{2} \lambda_{2}(\mu, i)<c_{1} \tilde{\lambda}_{1}+c_{2} \tilde{\lambda}_{2}, \tag{36}
\end{equation*}
$$

for all $\left(\tilde{\lambda}_{1}, \tilde{\lambda}_{2}\right) \in D$.
Consider the cost function:

$$
\lim _{N \rightarrow \infty} \frac{1}{N} E\left[\sum_{k=1}^{N}\left(c_{1} g_{1}\left(x_{k}, u_{k}\right)+c_{2} g_{2}\left(x_{k}, u_{k}\right)\right) \mid x_{0}=i\right]
$$

and denote by $\tilde{\mu} \in \mathcal{U}_{d}$, its optimal policy. Then:

$$
\begin{equation*}
c_{1} \lambda_{1}(\tilde{\mu}, i)+c_{2} \lambda_{2}(\tilde{\mu}, i) \leq c_{1} \lambda_{1}(\mu, i)+c_{2} \lambda_{2}(\mu, i) \tag{37}
\end{equation*}
$$

On the other hand, $\tilde{\mu}$ is deterministic and thus $\left(\lambda_{1}(\tilde{\mu}, i), \lambda_{2}(\tilde{\mu}, i)\right) \in D_{0} \subset D$. Hence, (37) contradicts (36).
(ii) Denote by $D_{v}=\left\{A_{1}, \ldots, A_{m}\right\}$ the extremal points (vertices) of $D$ (see figure 9). Of course $D_{v} \subset D_{0}$. Observe


Fig. 9: The set $D$.
that $D$ is the convex hull of $D_{v}$. Consider an $A_{j} \in D_{v}$. Let $A_{j} B$ be a segment bisecting the angle $A_{j-1}^{A_{j} A_{j+1}}, l$ a line perpendicular to $A_{j} B$ and $\theta$ the angle of $A_{j} B$ from the horizontal axis. Since $A_{j}$ is an extremal point, the angle $A_{j-1}{\widehat{A_{j} A}}_{j+1}$ is strictly less than $180^{\circ}$. Thus, $A_{j}=\left(\bar{\lambda}_{1}, \bar{\lambda}_{2}\right)$ is the unique common point of line $l$ and set $D$, and the unique minimizer of the problem:

$$
\underset{\left(\tilde{\lambda}_{1}, \tilde{\lambda}_{2}\right) \in D}{\operatorname{minimize}}\left\{\cos \theta \tilde{\lambda}_{1}+\sin \theta \tilde{\lambda}_{2}\right\}
$$

Consider the cost function:

$$
\lim _{N \rightarrow \infty} \frac{1}{N} E\left[\sum_{k=1}^{N}\left(\cos \theta g_{1}\left(x_{k}, u_{k}\right)+\sin \theta g_{2}\left(x_{k}, u_{k}\right)\right)\right]
$$

and denote by $\mu \in \mathcal{U}_{d}$, its optimal policy. Weak accessibility implies that the optimal value is the same for all the initial conditions and thus the point $\left(\lambda_{1}(\mu, i), \lambda_{2}(\mu, i)\right)$ lies on line $l$ for all $i \in X$. Thus, since $\left(\lambda_{1}(\mu, i), \lambda_{2}(\mu, i)\right) \in D$, and $A_{j}$ is the unique common point of $l$ and $D, \lambda_{1}(\mu, i)=\bar{\lambda}_{1}$ and $\lambda_{2}(\mu, i)=\bar{\lambda}_{2}$, for all $i \in X$. This completes the proof of (ii).
(iii) For every vertex $A_{j}$ there is deterministic stationary policy $\mu^{j}$ such that if $A_{j}=\left(\bar{\lambda}_{1}^{j}, \bar{\lambda}_{2}^{j}\right)$, then $\lambda_{1}\left(\mu^{j}, i\right)=\bar{\lambda}_{1}^{j}$, $\lambda_{2}\left(\mu^{j}, i\right)=\bar{\lambda}_{2}^{j}$, for all $i \in X$.

Consider a point $\left(\bar{\lambda}_{1}, \bar{\lambda}_{2}\right)=A \in D$. Then, since $D$ is the convex hull of $D_{v}$, there are scalars $\alpha_{1}, \ldots, \alpha_{m}$ such that $\bar{\lambda}_{1}=\alpha_{1} \bar{\lambda}_{1}^{1}+\cdots+\alpha_{m} \bar{\lambda}_{1}^{m}$ and $\bar{\lambda}_{2}=\alpha_{1} \bar{\lambda}_{2}^{1}+\cdots+\alpha_{m} \bar{\lambda}_{2}^{m}$. Assuming that $m \geq 2$, let us describe a policy $\bar{\mu}=\left(\bar{\mu}_{1}, \bar{\mu}_{2}, \ldots\right)$ such that the limits in definition of $\lambda_{1}(\bar{\mu}, i), \lambda_{2}(\bar{\mu}, i)$ exist and $\lambda_{1}(\bar{\mu}, i)=\bar{\lambda}_{1}, \lambda_{2}(\bar{\mu}, i)=\bar{\lambda}_{2}$. Policy $\bar{\mu}$ is separated in different epochs. The $\nu$-th epoch lasts for $\nu$ time steps. That is, for $\nu(\nu-1) / 2<k \leq \nu(\nu+1) / 2$ we are in the $\nu-$ th epoch. During the $\nu$-th epoch $\bar{\mu}_{k}=\mu^{1}$ for the first $\beta_{\nu}^{1}=\left[\alpha_{1} \nu\right]$ time steps, then $\bar{\mu}_{k}=\mu^{2}$ for the next $\beta_{\nu}^{2}=\left[\alpha_{2} \nu\right]$ time steps and so on. Here [:] denotes the integer part. Finally, $\bar{\mu}_{k}=\mu^{m}$, for the last $\beta_{m}=\nu-\beta_{1}-\cdots-\beta_{m-1}$ time steps. Denote by $I_{\nu, j}$ the set of time steps during the $\nu-$ th period, where $\bar{\mu}_{k}=\mu^{j}$. That is, $I_{\nu, j}=\left\{\nu(\nu-1) / 2+\sum_{j^{\prime}=1}^{j-1} \beta_{\nu}^{j^{\prime}}+1, \ldots, \nu(\nu-1) / 2+\right.$ $\left.\sum_{j^{\prime}=1}^{j} \beta_{\nu}^{j^{\prime}}\right\}$.

Let us derive some inequalities for the duration of $I_{\nu, j}$. For $j=1, \ldots, m-1$ it holds:

$$
0 \leq \alpha_{j} \nu-\beta_{\nu}^{j}<1
$$

Adding these inequalities and substituting $\sum_{j=1}^{m-1} \alpha_{j}=1-\alpha_{m}$ and $\sum_{j=1}^{m-1} \beta_{\nu}^{j}=\nu-\beta_{\nu}^{m}$ we get:

$$
0 \leq \beta_{\nu}^{m}-\nu \alpha_{m}<m-1
$$

Therefore, recalling that $m \geq 2$, for all $j=1, \ldots, m$ it holds:

$$
\begin{equation*}
\left|\beta_{\nu}^{j}-\nu \alpha_{j}\right|<m-1 \tag{38}
\end{equation*}
$$

Consider the sequences:

$$
S_{\nu, s}=E\left[\left.\frac{1}{\nu} \sum_{k=\nu(\nu-1) / 2+1}^{\nu(\nu+1) / 2} g_{s}\left(x_{k}, u_{k}\right) \right\rvert\, u_{k}=\mu_{k}\left(x_{k}\right)\right]
$$

for $s=1,2$. Then, $S_{\nu, s}$ can be written as:

$$
\begin{equation*}
S_{\nu, s}=\frac{1}{\nu} \sum_{j=1}^{m} \beta_{\nu}^{j} \bar{S}_{\nu, s}^{j} \tag{39}
\end{equation*}
$$

where:

$$
\bar{S}_{\nu, s}^{j}=E\left[\left.\frac{1}{\beta_{\nu}^{j}} \sum_{k \in I_{\nu, j}} g_{s}\left(x_{k}, u_{k}\right) \right\rvert\, u_{k}=\mu^{j}\left(x_{k}\right)\right] .
$$

We then prove the following claim:
Claim: The sequence $S_{\nu, s}$ converges to $\bar{\lambda}_{s}$ for $s=1,2$.
Combining this with (39) and (38), and using the fact that $\sum_{j=1}^{m} \beta_{\nu}^{j}=\nu$, we get:

$$
\begin{align*}
\left|S_{\nu, s}-\bar{\lambda}_{s}\right| & \leq \sum_{j=1}^{m}\left|\frac{\beta_{\nu}^{j}}{\nu} \bar{S}_{\nu, s}^{j}-\alpha_{j} \bar{\lambda}_{s}^{j}\right|=\frac{1}{\nu} \sum_{j=1}^{m}\left|\beta_{\nu}^{j} \bar{S}_{\nu, s}^{j}-\nu \alpha_{j} \bar{\lambda}_{s}^{j}\right| \\
& \leq \frac{1}{\nu} \sum_{j=1}^{m}\left[\beta_{\nu}^{j}\left|\bar{S}_{\nu, s}^{j}-\bar{\lambda}_{s}^{j}\right|+\left|\bar{\lambda}_{s}^{j}\right|\left|\beta_{\nu}^{j}-\nu \alpha_{j}\right|\right] \\
& \leq \max _{j}\left|\bar{S}_{\nu, s}^{j}-\bar{\lambda}_{s}^{j}\right|+\frac{m(m-1)}{\nu} \max _{j}\left|\bar{\lambda}_{s}^{j}\right| \tag{40}
\end{align*}
$$

Inequality (38) implies that $\beta_{\nu}^{j} \geq \nu \alpha_{j}-m+1$. Thus, $\bar{S}_{\nu, s}^{j} \rightarrow$ $\bar{\lambda}_{s}^{j}$, for all $j=1, \ldots, m, s=1,2$ and all $x_{0}$. Hence, $S_{\nu, s} \rightarrow$ $\bar{\lambda}_{s}$, which completes the proof of the claim.

Now, consider the sequence:

$$
\begin{equation*}
\Xi_{s, i, N}=\frac{1}{N} E\left[\sum_{k=1}^{N} g_{s}\left(x_{k}, u_{k}\right) \mid x_{0}=i\right] \tag{41}
\end{equation*}
$$

To prove (iii) it is sufficient to show that $\Xi_{s, i, N} \rightarrow \bar{\lambda}_{s}$, as $N \rightarrow \infty$, for all $i=1, \ldots, n$ and for $s=1,2$. First, observe that the subsequence $\Xi_{s, i, \bar{\nu}(\bar{\nu}+1) / 2}$ converges to the desired limit, as $\bar{\nu} \rightarrow \infty$. Indeed since:

$$
\begin{equation*}
\Xi_{s, i, \bar{\nu}(\bar{\nu}+1) / 2}=\frac{2}{\bar{\nu}(\bar{\nu}+1)} \sum_{\nu=1}^{\bar{\nu}} \nu S_{j, \nu} \tag{42}
\end{equation*}
$$

The claim implies that $\Xi_{s, i, \bar{\nu}(\bar{\nu}+1) / 2} \rightarrow \bar{\lambda}_{s}$.
Let us then examine $\Xi_{s, i, N}$ for $N$ between two successive evaluations of the subsequence. If $\bar{\nu}(\bar{\nu}+1) / 2 \leq N<(\bar{\nu}+$ 1) $(\bar{\nu}+2) / 2$ then:

$$
\begin{aligned}
& \Xi_{s, i, N}-\Xi_{s, i, \bar{\nu}(\bar{\nu}+1) / 2}=\frac{1}{N} \frac{\bar{\nu}(\bar{\nu}+1)}{2} \Xi_{s, i, \bar{\nu}(\bar{\nu}+1) / 2}+ \\
& \quad+\frac{1}{N} E\left[\left.\sum_{k=\frac{\bar{\nu}(\bar{\nu}+1)}{2}+1}^{N} g_{s}\left(x_{k}, u_{k}\right) \right\rvert\, x_{0}=i\right]-\Xi_{s, i, \bar{\nu}(\bar{\nu}+1) / 2}
\end{aligned}
$$

Thus:

$$
\begin{aligned}
& \left|\Xi_{s, i, N}-\Xi_{s, i, \bar{\nu}(\bar{\nu}+1) / 2}\right| \leq\left|\frac{\bar{\nu}(\bar{\nu}+1)}{2 N}-1\right| \Xi_{s, i, \bar{\nu}(\bar{\nu}+1) / 2}+ \\
& \quad+\frac{\bar{\nu} \max _{i^{\prime}, u} g_{s}\left(i^{\prime}, u\right)}{N}
\end{aligned}
$$

Therefore, the variation of $\Xi_{s, i, N}$ for $N$ between $\bar{\nu}(\bar{\nu}+1) / 2$ and $(\bar{\nu}+1)(\bar{\nu}+2) / 2$ tends to zero as $\bar{\nu} \rightarrow \infty$.

## B. Proof of Theorem 1

(i) Problem (11) is feasible, because for any $i \in X$ there is a policy $\bar{\mu}$ with $\lambda_{2}(\bar{\mu}, i)>0$. Due to Proposition 1 , for all $i$ Problem (11) has the same value with:

$$
\begin{equation*}
\underset{\left(\bar{\lambda}_{1}, \bar{\lambda}_{2}\right) \in D}{\operatorname{minimize}}\left\{\bar{\lambda}_{1} / \bar{\lambda}_{2}: \bar{\lambda}_{2}>0\right\} . \tag{43}
\end{equation*}
$$

Assumption 2 and the compactness of $D$ imply that this problem has a finite value (exists and it is greater than $-\infty$ ). Furthermore, since $D$ is a compact polygon and $\bar{\lambda}_{1} / \bar{\lambda}_{2}$ is quasi-linear function (e.g. [35]), $\bar{\lambda}_{1} / \bar{\lambda}_{2}$ is minimized on a vertex of $D$. Therefore, there is a deterministic optimal policy, according to Proposition 1.
(ii) Denote by $\lambda^{\star}$ the optimal value of (11). The proof proceeds in two discrete steps.

First, we show that $J_{a}$ cannot be much smaller than $\lambda^{\star}$, for a large value of $R$ (Lemma 3). Conversely, it is shown that a corner optimal policy for (10) has at most a slightly higher cost when applied to (8) (Lemma 4).

Throughout the proof we use repeatedly Doob's (sub-) martingale optional stopping theorem (e.g. [36]).

Theorem 3 (Doob's optional stopping theorem): Let $v_{t}$ be an $\mathcal{F}_{t}$-submartingale and $T$ a stopping time such that $E[T]<$ $\infty$. Assume further that there is a constant $C$ such that $\mid v_{t+1}-$ $v_{t} \mid<C$, for all $t$. Then:

$$
E\left[v_{T}\right] \geq E\left[v_{0}\right]
$$

The last relation holds as equality if $v_{t}$ is a martingale.
We use also the following lemma the proof of which is almost identical to that of the optional stopping theorem. For the shake of completeness we present its proof.

Lemma 1: Let $Z_{k}$ be a sequence of random variables and assume that there is a constant $C$ such that $\left|Z_{k}\right|<C$. Assume that $T$ is a random time with $E[T]<\infty$. Then:

$$
\begin{equation*}
\lim _{N \rightarrow \infty} E\left[\sum_{k=1}^{T \wedge N} Z_{k}\right]=E\left[\sum_{k=1}^{T} Z_{k}\right] . \tag{44}
\end{equation*}
$$

Equation (44) asserts also the existence of the limit.
Proof: Since $E[T]<\infty$, the random variable $Y_{N}=$ $\sum_{k=1}^{T \wedge N} Z_{k}$ converges almost surely to $Y=\sum_{k=1}^{T} Z_{k}$, as $N \rightarrow \infty$. Furthermore, for all $N$ the random variable $Y_{N}$ is bounded above by:

$$
\left|Y_{N}\right| \leq C \sum_{k=1}^{\infty} I_{T \geq k}=C T
$$

where $I$. is the indicator (characteristic) function. But, since $E[T]<\infty$, the dominated convergence theorem applies and it holds:

$$
\lim _{N \rightarrow \infty} E\left[Y_{N}\right]=E[Y]
$$

which is exactly (44). Additionally, dominated convergence theorem proves also the existence of the limit.

For an $n$-vector $h$, let us introduce the span semi-norm:

$$
\begin{equation*}
\operatorname{sp}(h)=\max _{i} h(i)-\min _{i} h(i) . \tag{45}
\end{equation*}
$$

Lemma 2: Suppose that assumptions 1 and 2 hold. Consider any behavioral policy $\mu=\left(\mu_{0}, \mu_{1}, \ldots\right)$. If under $\mu$ :

$$
\underline{\mathrm{L}}=\liminf _{N \rightarrow \infty} E\left[\sum_{k=0}^{T \wedge N-1} g_{1}\left(x_{k}, u_{k}\right)\right]<\infty
$$

then:
(i) The expected horizon is finite i.e., $E[T]<\infty$.
(ii) The limit in $J_{a}$ exists and satisfies:

$$
\begin{align*}
J_{a}\left(x_{0}, \mu\right) & =\lim _{N \rightarrow \infty} E\left[\sum_{k=0}^{T \wedge N-1} g_{1}\left(x_{k}, u_{k}\right) / R\right]  \tag{46}\\
& =E\left[\sum_{k=0}^{T-1} g_{1}\left(x_{k}, u_{k}\right) / R\right] \tag{47}
\end{align*}
$$

Therefore, the limit in the definition of $J_{a}$ exists always (either finite or infinite).
Proof: Consider closed third quadrant:

$$
Q_{3}=\left\{(x, y) \in \mathbb{R}^{2}: x \leq 0, y \leq 0\right\}
$$

and the feasible set $D$. The set $D$ is compact and convex and $Q_{3}$ is closed and convex. According to Assumption 2 $D \cap Q_{3}=\emptyset$. Thus, there is a strictly separating line between $D$ and $Q_{3}$. This implies that there are constants $c_{1}, c_{2}, c_{3}>0$ such that:

$$
\begin{equation*}
c_{1} \lambda_{1}(\bar{\mu}, i)+c_{2} \lambda_{2}(\bar{\mu}, i) \geq c_{3} \tag{48}
\end{equation*}
$$

for all the stationary policies $\bar{\mu}$.
Consider the cost per stage $\bar{g}(i, u)=c_{1} g_{1}(i, u)+$ $c_{2} g_{2}(i, u)-c_{3}$. Then, due to (48), the minimal LTA cost for $\bar{g}$ is non-negative. Thus, there is a pair $\bar{\lambda} \geq 0, \bar{h} \in \mathbb{R}^{n}$ such that:

$$
\begin{equation*}
\bar{\lambda}+\bar{h}(i)=\min _{u}\left[\bar{g}(i, u)+\sum_{j=1}^{n} p_{i j}(u) \bar{h}(j)\right] \tag{49}
\end{equation*}
$$

Denote by $\mathcal{F}_{t}$ the $\sigma$ algebra generated by $\left(x_{0}, u_{0}, \ldots, x_{t}, u_{t}\right)$, and consider the stochastic process:

$$
v_{t}=\bar{h}\left(x_{t}\right)+\sum_{k=0}^{t-1} \bar{g}\left(x_{k}, u_{k}\right) .
$$

Observe that $v_{t}$ is a $\mathcal{F}_{t}$-submartingale. Indeed:

$$
\begin{aligned}
E\left[v_{t+1}-v_{t} \mid \mathcal{F}_{t}\right] & =-\bar{h}\left(x_{t}\right)+\bar{g}\left(x_{t}, u_{t}\right)+\sum_{j=1}^{n} p_{i j}\left(u_{t}\right) \bar{h}(j) \\
& \geq \bar{\lambda} \geq 0
\end{aligned}
$$

where for the first inequality we used the Bellman equation (49). Let us consider the stopped process $v_{T_{t}}$, where $T_{t}=$ $T \wedge t$. Since, $E\left[T_{t}\right] \leq t$, the optional stopping theorem implies that $E\left[v_{T_{t}}\right] \geq E\left[v_{0}\right]$, or equivalently:

$$
\begin{array}{r}
E\left[c_{1} \sum_{k=0}^{T_{t}-1} g_{1}\left(x_{k}, u_{k}\right)+c_{2} \sum_{k=0}^{T_{t}-1} g_{2}\left(x_{k}, u_{k}\right)-c_{3}\left(T_{t}-1\right)\right] \geq \\
\geq E\left[h\left(x_{0}\right)-h\left(x_{T_{t}}\right)\right]
\end{array}
$$

Furthermore:

$$
E\left[h\left(x_{0}\right)-h\left(x_{T_{t}}\right)\right] \geq-\operatorname{sp}(\bar{h}),
$$

where $\operatorname{sp}(h)$ is defined in (45). Thus:

$$
\begin{array}{r}
E\left[T_{t}\right] \leq \frac{c_{1}}{c_{3}} E\left[\sum_{k=0}^{T_{t}-1} g_{1}\left(x_{k}, u_{k}\right)\right]+\frac{c_{2}}{c_{3}} E\left[\sum_{k=0}^{T_{t}-1} g_{2}\left(x_{k}, u_{k}\right)\right]+ \\
+1+\operatorname{sp}(\bar{h}) / c_{3} .
\end{array}
$$

Observe that $\sum_{k=0}^{T_{t}-1} g_{2}\left(x_{k}, u_{k}\right) \leq R$. Hence:

$$
E\left[T_{t}\right] \leq \frac{c_{1}}{c_{3}} S_{t}+\Delta
$$

where $\Delta=\frac{c_{2}}{c_{3}} R+1+\operatorname{sp}(\bar{h}) / c_{3}$ and $S_{t}=$ $E\left[\sum_{k=0}^{T_{t}-1} g_{1}\left(x_{k}, u_{k}\right)\right]$.

There is a subsequence $S_{t_{\nu}}$ which attains the limit inferior, i.e., $S_{t_{\nu}} \rightarrow \underline{\mathrm{L}}<\infty$. Thus, there is a $\nu_{0}$ such that $S_{t_{\nu}} \leq \underline{\mathrm{L}}+1$, for all $\nu \geq \nu_{0}$. Observe that, since $T_{t}$ is non-decreasing in $t$, it holds:

$$
E\left[T_{t}\right] \leq \frac{c_{1}}{c_{3}}(\underline{\mathrm{~L}}+1)+\Delta
$$

for all $t$. Thus, due to monotone convergence theorem we conclude that:

$$
E[T] \leq \frac{c_{1}}{c_{3}}(\underline{\mathrm{~L}}+1)+\Delta<\infty
$$

This completes the proof of (i) of the lemma. Then, (ii) is a consequence of Lemma 1.

Lemma 3: For every $\varepsilon>0$ there is an $R_{0}$ such that for any $R \geq R_{0}$ and any behavioral policy $\mu=\left(\mu_{0}, \mu_{1}, \ldots\right)$ it holds:

$$
\lambda^{\star} \leq J_{a}\left(x_{0}, \mu\right)+\varepsilon,
$$

where $\lambda^{\star}$ is the optimal value of (11).
Proof: Throughout this proof we assume that $u_{k}$ is chosen according to $\mu_{k}\left(x_{k}\right)$. In the case where $J_{a}\left(x_{0}, \mu\right)=\infty$, the lemma is trivially true. Thus, we assume that $J_{a}\left(x_{0}, \mu\right)<\infty$. According to Proposition 4 and Corollary $1^{3}$, there is a vector $h^{\star}$ satisfying:

$$
\begin{equation*}
h^{\star}(i)=\min _{u}\left[g_{1}(i, u)-\lambda^{\star} g_{2}(i, u)+\sum_{j=1}^{n} p_{i j}(u) h^{\star}(j)\right] . \tag{50}
\end{equation*}
$$

Consider the stochastic process:

$$
v_{t}=h^{\star}\left(x_{t}\right)+\sum_{k=0}^{t-1}\left[g_{1}\left(x_{k}, u_{k}\right)-\lambda^{\star} g_{2}\left(x_{k}, u_{k}\right)\right]
$$

Observe that $v_{t}$ is an $\mathcal{F}_{t}$-submartingale. Indeed, due to (50) it holds:

$$
\begin{aligned}
& E\left[v_{t+1}-v_{t} \mid \mathcal{F}_{t}\right]=-h^{\star}\left(x_{t}\right)+ \\
& \quad+\left[g_{1}\left(x_{t}, u_{t}\right)-\lambda^{\star} g_{2}\left(x_{t}, u_{t}\right)+\sum_{j=1}^{n} p_{i j}\left(u_{t}\right) h^{\star}(j)\right] \geq 0 .
\end{aligned}
$$

Using optional stopping theorem we get $E\left[v_{T_{t}}\right] \geq E\left[v_{0}\right]=$ $E\left[h\left(x_{0}\right)\right]$. Therefore:
$E\left[h^{\star}\left(x_{t}\right)+\sum_{k=0}^{T_{t}-1}\left[g_{1}\left(x_{k}, u_{k}\right)-\lambda^{\star} g_{2}\left(x_{k}, u_{k}\right)\right]\right] \geq E\left[h\left(x_{0}\right)\right]$.
${ }^{3}$ Let us first note that the proof of Proposition 2 and Proposition 4 do not depend on the proof of (ii) of this proposition and thus this argument is not circular.

Thus:

$$
\begin{equation*}
\left.E\left[\sum_{k=0}^{T_{t}-1} g_{1}\left(x_{k}, u_{k}\right)\right] \geq \lambda^{\star} E\left[\sum_{k=0}^{T_{t}-1} g_{2}\left(x_{k}, u_{k}\right)\right]\right]-\operatorname{sp}(h) . \tag{51}
\end{equation*}
$$

We examine the two cases of Assumption 2 separately. Case 1: Assumption 2.a holds true.
In this case $\lambda^{\star}<0$ and thus $\sum_{k=0}^{T_{t}-1} g_{2}\left(x_{k}, u_{k}\right)<R$ implies:

$$
\lambda^{\star} E\left[\sum_{k=0}^{T_{t}-1} g_{2}\left(x_{k}, u_{k}\right)\right]>\lambda^{\star} R
$$

Combining this with (51) we get:

$$
\lambda^{\star} \leq E\left[\sum_{k=1}^{T_{t}-1} g_{1}\left(x_{k}, u_{k}\right) / R\right]+\operatorname{sp}(h) / R .
$$

Taking the limit of the right-hand side we complete the proof of the lemma for Case 1.
Case 2: Assumption 2.b holds true.
Taking the limits in (51) we get:

$$
\begin{equation*}
\left.J_{a}\left(\mu, x_{0}\right) \geq \lambda^{\star} \lim _{t \rightarrow \infty} E\left[\sum_{k=0}^{T_{t}-1} g_{2}\left(x_{k}, u_{k}\right)\right]\right]-\operatorname{sp}(h) \tag{52}
\end{equation*}
$$

Using Lemma 1 we get:

$$
\begin{aligned}
\left.\lim _{t \rightarrow \infty} E\left[\sum_{k=0}^{T_{t}-1} g_{2}\left(x_{k}, u_{k}\right)\right]\right]=E & {\left.\left[\sum_{k=0}^{T-1} g_{2}\left(x_{k}, u_{k}\right)\right]\right] \geq } \\
& \geq R-\max _{i, u} g_{2}(i, u)
\end{aligned}
$$

Combining this inequality with (52) and observing that $\lambda^{\star}>0$ we get:

$$
\lambda^{\star} \leq J_{a}\left(x_{0}, \mu\right)+\frac{\max _{i, u} g_{2}(i, u) \lambda^{\star}}{R}+\frac{\operatorname{sp}(h)}{R}
$$

This completes the proof.
We then prove the converse:
Lemma 4: Let $\mu^{\star}, \lambda^{\star}$ be a corner optimal policy and the optimal value for (11). Then, for any $\varepsilon>0$, there is an $R_{0}$ such that:

$$
J_{a}\left(\mu^{\star}, x_{0}\right) \leq \lambda^{\star}+\varepsilon
$$

for all $x_{0}$ and all $R \geq R_{0}$.
Proof: (i) Since $\mu^{\star}$ is a corner policy, there is a pair $h_{2}^{\star} \in \mathbb{R}^{n}$, $\lambda_{2}^{\star}>0$ satisfying:

$$
\lambda_{2}^{\star}+h_{2}^{\star}(i)=g_{2}\left(i, \mu^{\star}(i)\right)+\sum_{j=1}^{n} p_{i j}\left(\mu^{\star}(i)\right) h_{2}^{\star}(j)
$$

Then, stochastic process:

$$
v_{t}=h_{2}^{\star}\left(x_{t}\right)+\sum_{k=0}^{t-1}\left[g_{2}\left(x_{k}, u_{k}\right)-\lambda_{2}^{\star}\right]
$$

is a martingale. Similarly to the previous lemma we have $E\left[v_{T_{t}}\right]=E\left[v_{0}\right]=E\left[h_{2}^{\star}\left(x_{0}\right)\right]$. Hence:
$E\left[\sum_{k=0}^{T_{t}-1} g_{2}\left(x_{k}, u_{k}\right)-\lambda_{2}^{\star}\left(T_{t}-1\right)\right]=E\left[h_{2}^{\star}\left(x_{0}\right)\right]-E\left[h_{2}^{\star}\left(x_{t}\right)\right]$.

Thus,

$$
E\left[T_{t}\right] \leq \frac{1}{\lambda_{2}^{\star}}\left(R+\operatorname{sp}\left(h_{2}^{\star}\right)\right)+1
$$

and monotone convergence theorem shows that $E[T]<\infty$.
For the optimal policy optimal value pair $\mu^{\star}, \lambda^{\star}$, there is a vector $h^{\star} \in \mathbb{R}^{n}$ such that:

$$
h^{\star}(i)=g_{1}\left(i, \mu^{\star}(i)\right)-\lambda^{\star} g_{2}\left(i, \mu^{\star}(i)\right)+\sum_{j=1}^{n} p_{i j}\left(\mu^{\star}(i)\right) h^{\star}(j)
$$

Then, stochastic process:

$$
v_{t}=h^{\star}\left(x_{t}\right)+\sum_{k=0}^{t-1}\left[g_{1}\left(x_{k}, u_{k}\right)-\lambda^{\star} g_{2}\left(x_{k}, u_{k}\right)\right]
$$

is a martingale. Similarly to the previous lemma we have $E\left[v_{T_{t}}\right]=E\left[v_{0}\right]=E\left[h^{\star}\left(x_{0}\right)\right]$. Hence:

$$
\begin{array}{r}
E\left[\sum_{k=0}^{T_{t}-1}\left[g_{1}\left(x_{k}, u_{k}\right)-\lambda^{\star} g_{2}\left(x_{k}, u_{k}\right)\right]\right]= \\
=E\left[h^{\star}\left(x_{0}\right)\right]-E\left[h^{\star}\left(x_{t}\right)\right]
\end{array}
$$

Therefore:

$$
\begin{equation*}
E\left[\sum_{k=0}^{T_{t}-1} g_{1}\left(x_{k}, u_{k}\right)\right] \leq \lambda^{\star} E\left[\sum_{k=0}^{T_{t}-1} g_{2}\left(x_{k}, u_{k}\right)\right]+\operatorname{sp}\left(h^{\star}\right) \tag{53}
\end{equation*}
$$

Case 1: Assumption 2.b holds.
In this case, using that $\sum_{k=0}^{T_{t}-1} g_{2}\left(x_{k}, u_{k}\right) \leq R$, and that $\lambda^{\star}>$ 0 , (53) becomes:

$$
\begin{equation*}
E\left[\sum_{k=0}^{T_{t}-1} g_{1}\left(x_{k}, u_{k}\right)\right] \leq \lambda^{\star} R+\operatorname{sp}\left(h^{\star}\right) \tag{54}
\end{equation*}
$$

and dividing by $R$ and taking the limit we arrive to the desired conclusion.
Case 2: Assumption 2.a holds.
Similarly with Case 2 of Lemma 3, we have:

$$
\begin{aligned}
\left.\lim _{t \rightarrow \infty} E\left[\sum_{k=0}^{T_{t}-1} g_{2}\left(x_{k}, u_{k}\right)\right]\right]=E & {\left.\left[\sum_{k=0}^{T-1} g_{2}\left(x_{k}, u_{k}\right)\right]\right] \geq } \\
& \geq R-\max _{i, u} g_{2}(i, u)
\end{aligned}
$$

Therefore:

$$
\begin{aligned}
J_{a}\left(x_{0}, \mu^{\star}\right)= & \lim _{t \rightarrow \infty} E\left[\sum_{k=0}^{T_{t}-1} g_{1}\left(x_{k}, u_{k}\right) / R\right] \leq \\
& \leq \lambda^{\star}+\frac{\operatorname{sp}\left(h^{\star}\right)-\lambda^{\star} \max _{i, u} g_{2}(i, u)}{R}
\end{aligned}
$$

This completes the proof.
Proof of Theorem 1: Combining lemmas 3 and 4, for large $R$ and any policy $\mu$, we get:

$$
J_{a}\left(\mu^{\star}, x_{0}\right) \leq \lambda^{\star}+\varepsilon \leq J_{a}\left(\mu, x_{0}\right)+2 \varepsilon .
$$

Thus, $\mu^{\star}$ is $\varepsilon$-optimal for large $R$.


Fig. 10: The set $D$ and a line with $\bar{\lambda}_{1} / \bar{\lambda}_{2}$ constant.

## C. Proof of Proposition 3

We show that in the generic case any deterministic static optimal policy is a corner policy. We first show that for the most of the values of the costs $g_{1}, g_{2}$ there is a unique optimum in the Problem (43).

Lemma 5: Suppose that the controlled Markov chain is weakly accessible (Assumption 1.a holds). Then, the set of vectors of the form $\left[g_{s}(x, u): x \in X, u \in U, s=1,2\right] \in$ $\mathbb{R}^{2 n n_{U}}$ such that Problem (43) has multiple minima has Lebesgue measure zero.
Proof: Observe that the value of $\bar{\lambda}_{1} / \bar{\lambda}_{2}$ is constant, in the $\left(\bar{\lambda}_{1}, \bar{\lambda}_{2}\right)$ plane, on half-lines starting from the origin. Therefore, if (43) has multiple minima then they correspond to a line segment i.e., the intersection of the half-line with $D$. Furthermore, observe that $\bar{\lambda}_{1} / \bar{\lambda}_{2}=\cot \theta$, where $\theta$ is the angle of the half line with the horizontal axis (see Figure 11). Thus, since $\cot (\cdot)$ is strictly decreasing in $(0, \pi)$, if (43) has multiple minima they correspond to an edge of $D$. Therefore, there are two vertices of $D$ co-linear with the origin.

Using Proposition 1, if (43) has multiple minima then there is a pair of policies $\mu, \mu^{\prime} \in \mathcal{U}_{d}$ and a pair of states $i, i^{\prime} \in X$ such that:

$$
\begin{array}{cl}
\frac{\lambda_{1}(\mu, i)}{\lambda_{2}(\mu, i)}= & \frac{\lambda_{1}\left(\mu^{\prime}, i^{\prime}\right)}{\lambda_{2}\left(\mu^{\prime}, i^{\prime}\right)}  \tag{55}\\
\lambda_{1}(\mu, i) \neq \lambda_{1}\left(\mu^{\prime}, i^{\prime}\right), \quad & \lambda_{2}(\mu, i), \lambda_{2}\left(\mu^{\prime}, i^{\prime}\right)>0
\end{array}
$$

The set of states $X$ and the set of stationary deterministic strategies $\mathcal{U}_{d}$ are finite. Thus, it is sufficient to show that for any pair of states $i, i^{\prime}$ and any pair of strategies $\mu, \mu^{\prime} \in \mathcal{U}_{d}$, (55) does not hold in the generic case.

Denote by $\mathfrak{g}_{s}=\left[g_{s}(1,1), \ldots, g_{s}\left(1, n_{U}\right), \ldots\right.$, $\left.g_{s}(n, 1), \ldots, g_{s}\left(n, n_{U}\right)\right]$, where $s=1,2$, and $n, n_{U}$ the cardinalities of the state space and the set of actions. Then, using [20] (Theorem 8.1.1), for a state $\bar{i} \in X$ and a policy $\bar{\mu} \in \mathcal{U}_{d}$ we have:

$$
\lambda_{s}(\bar{\mu}, \bar{i})=e_{\bar{i}}^{T} P_{\bar{\mu}}^{\star} S_{\bar{\mu}} \mathfrak{g}_{s}
$$

where $e_{\bar{i}}$ an $n$-vector with $\bar{i}$-th entry equal to one and all the other entries equal to zero, $P_{\bar{\mu}}^{\star}$ is the limiting matrix of the Markov chain under the policy $\bar{\mu}$ and $S_{\bar{\mu}}$ is an $n \times n|U|$ matrix such that $S_{\bar{\mu}} \mathfrak{g}_{s}=\left[g_{s}(1, \bar{\mu}(1)), \ldots, g_{s}(n, \bar{\mu}(1))\right]^{T}$. Since the limiting matrix $P_{\bar{\mu}}^{\star}$ is a stochastic matrix and $S_{\bar{\mu}}$ has all its entries equal to zero except of a single entry on each row which is equal to one, the $1 \times n n_{U}$ row vector:

$$
d_{\bar{\mu}, \bar{i}}=e_{\bar{i}}^{T} P_{\bar{\mu}}^{\star} S_{\bar{\mu}}
$$

has non-negative entries summing to one.
Assume a pair of policies $\mu, \mu^{\prime}$ and a pair of states $i, i^{\prime}$. If (55) holds then:

$$
\begin{equation*}
\frac{d_{\mu, i} \mathfrak{g}_{1}}{d_{\mu^{\prime}, i^{\prime}} \mathfrak{g}_{1}}=\frac{d_{\mu, i} \mathfrak{g}_{2}}{d_{\mu^{\prime}, i^{\prime}} \mathfrak{g}_{2}} \neq 1 \tag{56}
\end{equation*}
$$

If $d_{\mu, i}=d_{\mu^{\prime}, i^{\prime}}$ then (55) does not holds. Thus, assume $d_{\mu, i} \neq d_{\mu^{\prime}, i^{\prime}}$.

Denote by by $\mathfrak{g}_{1}^{1}$ the first entry of $\mathfrak{g}_{1}$ and by $\mathfrak{g}_{1}^{\text {Rest }}$ the rest of the entries. Similarly, denote by $d_{\mu, i}^{1}, d_{\mu^{\prime}, i^{\prime}}^{1}$ the first entries $d_{\mu, i}, d_{\mu^{\prime}, i^{\prime}}$ and $d_{\mu, i}^{\text {Rest }}, d_{\mu^{\prime}, i^{\prime}}^{\text {Rest }}$, the rest of the entries. Observe that for all the values of $\mathfrak{g}_{1}^{\text {Rest }}$ and $\mathfrak{g}_{2}$ there is at most one value of $\mathfrak{g}_{1}^{1}$ such that (56) holds which, whenever is defined, is given by:

$$
f\left(\mathfrak{g}_{1}^{\text {Rest }}, \mathfrak{g}_{2}\right)=\frac{\left(d_{\mu, i} \mathfrak{g}_{2}\right) d_{\mu^{\prime}, i^{\prime}}^{\text {Rest }} \mathfrak{g}_{1}^{\text {Rest }}-\left(d_{\mu^{\prime}, i^{\prime}} \mathfrak{g}_{2}\right) d_{\mu, i}^{\text {Rest }} \mathfrak{g}_{1}^{\text {Rest }}}{\left(d_{\mu^{\prime}, i^{\prime}} \mathfrak{g}_{2}\right) d_{\mu, i}^{1}-\left(d_{\mu, i} \mathfrak{g}_{2}\right) d_{\mu^{\prime}, i^{\prime}}^{1}}
$$

Thus, the set of points in $\mathbb{R}^{2 n n_{U}}$, such that (55) holds is a subset of the graph of $f$. Therefore, it has Lebesgue measure zero ([37], exercise 3.10.52).

To prove the proposition, assume that $\mu^{\star}$ is such a policy. Proposition A implies that $\lambda_{1}\left(\mu^{\star}, i\right) / \lambda_{2}\left(\mu^{\star}, i\right)$ has the same (optimal) value for all $i$. But, since, $\mu^{\star}$ is not a corner policy, there is a pair $i, i^{\prime} \in X$ such that $\lambda_{1}(\mu, i) \neq \lambda_{1}\left(\mu, i^{\prime}\right)$. This corresponds to a case where Problem (43) has multiple minima.

## D. Proof of Theorem 2

The proof is based on a discrete version of singular perturbation theory (e.g. [30]). The proof is follows 4 steps. First we show that for $\lambda$ close to $\lambda^{\star}$, the optimal policy $\mu^{\lambda}$ for (30) is such that $\lambda_{2}\left(\mu^{\lambda}\right)>0$. Then, for a fixed (frozen) value of $\lambda$, we construct a Lyapunov function for the dynamics (27),(28). Then, we prove some technical lemmas for the Lyapunov function. Finally, we use these lemmas to construct a composite Lyapunov function for (27)-(28).

1) Positivity of $\lambda_{2}\left(\mu^{\lambda}\right)$ : Assumption 1.a and the compactness of the feasible region (see Figure 1) imply that there is a positive constant $\varepsilon_{\lambda_{2}}$ such that for each policy $\mu$ either $\lambda_{1}(\mu) \geq 0$ or $\lambda_{2}(\mu) \geq \varepsilon_{\lambda_{2}}$. We state first the following lemma.

Lemma 6: There is a positive scalar $\delta_{\lambda}$ such that if $\lambda>\lambda^{\star}-$ $\delta_{\lambda}$ then the policy $\mu^{\lambda}$ minimizing (30) satisfies $\lambda_{2}\left(\mu^{\lambda}\right) \geq \varepsilon_{\lambda_{2}}$. Proof : Refers to Figure 11. Let $F$ be optimal point of the feasible set $D$ (i.e. having the lowest $\lambda_{1} / \lambda_{2}$ ) closest to the origin. Recall that $D$ is a convex polygon. Let $Q$ be the next vertex of $D$ counting counter-clockwise. From Proposition 4, $O F$ is perpendicular to the vector $\left(-1, \lambda^{\star}\right)$. Observe, there is a vector $\left(-1, \lambda^{\star}-\delta_{\lambda}\right)$, such that the vectors $\left(-1, \lambda^{\star}\right)$ and $\left(-1, \lambda^{\star}-\delta_{\lambda}\right)$ have angle less than or equal to $\widehat{O F Q}$. Thus, for $\lambda>\lambda^{\star}-\delta_{\lambda}$ the extreme points of $D$ in the direction $(-1, \lambda)$ will have $\lambda_{2}>0$. The degenerate cases, where $D$ is a point or an interval follow easily.
2) Construction of the Lyapunov Function: Let us introduce some notation. The norm symbol $\|\cdot\|$ denotes the infinity norm. For a scalar $\lambda$, denote by $\mathcal{T}_{\lambda}$ the Bellman operator:

$$
\left(\mathcal{T}_{\lambda} h\right)(i)=\min _{u}\left[g_{1}(i, u)-\lambda g_{2}(i, u)+\tau \sum_{j=1}^{n} h(j) p_{i j}(u)\right]
$$



Fig. 11: The situation in Lemma 6.
and by $F_{\lambda}$ the corresponding relative value iteration operator:

$$
F_{\lambda} h=\mathcal{T}_{\lambda} h-\left(\mathcal{T}_{\lambda} h\right)(n) \cdot \mathbf{1}+(1-\tau) h
$$

Let $h^{\lambda}$ be the unique fixed point of $F_{\lambda}$. It is not difficult to see that $\left(\lambda, h^{\lambda} \tau\right)$ satisfy the Bellman equation for (30). Denote also by $\bar{G}$ the maximum value of $\left|g_{2}(i, u)\right|$.

For a pair of scalars $\lambda$ and $\lambda^{\prime}$ and a pair of $n$-vectors $h$ and $h^{\prime}$, operator $T_{\lambda}$ satisfies the following pair of inequalities:

$$
\begin{equation*}
\left\|\mathcal{T}_{\lambda} h-\mathcal{T}_{\lambda^{\prime}} h\right\| \leq \bar{G}\left|\lambda-\lambda^{\prime}\right|, \quad\left\|\mathcal{T}_{\lambda} h-\mathcal{T}_{\lambda} h^{\prime}\right\| \leq \tau\left\|h-h^{\prime}\right\| \tag{57}
\end{equation*}
$$

Similarly:

$$
\begin{equation*}
\left\|F_{\lambda} h-F_{\lambda^{\prime}} h\right\| \leq 2 \bar{G}\left|\lambda-\lambda^{\prime}\right|, \quad\left\|F_{\lambda} h-F_{\lambda} h^{\prime}\right\| \leq 2\left\|h-h^{\prime}\right\| \tag{58}
\end{equation*}
$$

Inequalities (57), (58) will be repeatedly used throughout the proof.

Denote by $\phi_{\lambda, h}^{k}$ the vector defined by:

$$
\begin{align*}
& \phi_{\lambda, h}^{k+1}=F_{\lambda} \phi_{\lambda, h}^{k}  \tag{59}\\
& \phi_{\lambda, h}^{0}=h
\end{align*}
$$

Let $q_{\lambda, h}^{k}$ be the difference of two successive iterates of $\phi_{\lambda, h}$ :

$$
q_{\lambda, h}^{k}=\phi_{\lambda, h}^{k+1}-\phi_{\lambda, h}^{k}
$$

Finally, recall the the span semi-norm $\operatorname{sp}(\cdot)$ of an $n$-vector $h$ introduced in (45). As an intermediate step for proving the convergence of the usual value iteration method it has been shown that there is a positive integer $m$ and an $\varepsilon>0$ such that:

$$
\begin{equation*}
\operatorname{sp}\left(q_{\lambda, h}^{k}\right) \leq(1-\varepsilon) \operatorname{sp}\left(q_{\lambda, h}^{k-m}\right) \tag{60}
\end{equation*}
$$

for all $k \geq m$ (see for example [19], [20]). More importantly, $m$ and $\varepsilon$ depend only on the properties of the controlled Markov chain and not on the cost function. Thus, (60) is satisfied for all $\lambda$, for the same values of $m$ and $\varepsilon$. Using (60), we can bound the evolution of $\operatorname{sp}\left(q_{\lambda, h}^{k}\right)$, in the form:

$$
\operatorname{sp}\left(q_{\lambda, h}^{k}\right) \leq B(h, \lambda) \xi^{k}
$$

where $\xi=(1-\varepsilon)^{1 / m}$ and:

$$
\begin{equation*}
B(h, \lambda) \geq \max _{k=0, \ldots, m} \operatorname{sp}\left(q_{\lambda, h}^{k}\right) /(1-\varepsilon) \tag{61}
\end{equation*}
$$

Based on this bound, a convergence rate for (59) can be derived. It holds $q_{\lambda, h}^{k}(n)=0$. Thus, $\min _{i}\left\{q_{\lambda, h}^{k}(i)\right\} \leq 0$ and $\left\|q_{\lambda, h}^{k}\right\| \leq \operatorname{sp}\left(q_{\lambda, h}^{k}\right)$. Therefore,

$$
\begin{equation*}
\left\|\phi_{\lambda, h}^{k}-h^{\lambda}\right\| \leq \sum_{t=k}^{\infty}\left\|q_{\lambda, h}^{t}\right\| \leq \frac{B(h, \lambda) \xi^{k}}{1-\xi} \tag{62}
\end{equation*}
$$

We then derive a formula for $B(h, \lambda)$ which does not depend on $\lambda$. To this end, observe that for a vector $h^{\prime}$ it holds:

$$
\left\|F_{\lambda} h^{\prime}-h^{\prime}\right\| \leq\left\|F_{\lambda} h^{\prime}-F_{\lambda} h^{\lambda}\right\|+\left\|F_{\lambda} h^{\lambda}-h^{\prime}\right\|
$$

Applying the second inequality of (58) to the first term of the right-hand side and recalling that $F_{\lambda} h^{\lambda}=h^{\lambda}$ we get:

$$
\left\|F_{\lambda} h^{\prime}-h^{\prime}\right\| \leq 3\left\|h^{\prime}-h^{\lambda}\right\|
$$

Thus:
$\left\|\phi_{\lambda, h}^{k+1}-h^{\lambda}\right\| \leq\left\|\phi_{\lambda, h}^{k+1}-\phi_{\lambda, h}^{k}\right\|+\left\|\phi_{\lambda, h}^{k}-h^{\lambda}\right\| \leq 4\left\|\phi_{\lambda, h}^{k}-h^{\lambda}\right\|$.
Using iteratively this inequality, we get:
$\left\|\phi_{\lambda, h}^{k+1}-h^{\lambda}\right\| \leq 4^{k}\left\|h-h^{\lambda}\right\|, \quad\left\|\phi_{\lambda, h}^{k+1}-\phi_{\lambda, h}^{k}\right\| \leq 3 \cdot 4^{k-1}\left\|h-h^{\lambda}\right\|$.
Furthermore, $\operatorname{sp}\left(q_{\lambda, h}^{k}\right) \leq 2\left\|q_{\lambda, h}^{k}\right\|$. Thus, (61) is satisfied with:

$$
\begin{equation*}
B(h, \lambda)=6 \cdot 4^{m-1}\left\|h-h^{\lambda}\right\|=B_{0}\left\|h-h^{\lambda}\right\| . \tag{64}
\end{equation*}
$$

Define $K$ as the minimum positive integer such that $\rho=$ $B_{0} \xi^{K} /(1-\xi)<1$. Then the function:

$$
\begin{equation*}
V(h, \lambda)=\sum_{k=0}^{K-1}\left\|\phi_{\lambda, h}^{k}-h^{\lambda}\right\|^{2} \tag{65}
\end{equation*}
$$

is a Lyapunov function for the dynamics (59). Indeed:

$$
\begin{align*}
& V\left(F_{\lambda} h, \lambda\right)-V(h, \lambda)=\left\|\phi_{\lambda, h}^{k}-h^{\lambda}\right\|^{2}-\left\|h-h^{\lambda}\right\|^{2} \leq \\
& \quad \leq-(1-\rho)\left\|h-h^{\lambda}\right\|^{2} \tag{66}
\end{align*}
$$

where the last inequality holds true due to (62).
3) Technical Lemmas: The proof of Theorem 2 depends on the following four lemmas.

Lemma 7: There is a positive constant $L_{1}$ such that:

$$
\left\|h^{\lambda}-h^{\lambda^{\prime}}\right\| \leq L_{1}\left|\lambda-\lambda^{\prime}\right|
$$

for any $\lambda, \lambda^{\prime}$
Proof: The limit of the dynamics $\phi_{\lambda, h^{\prime}}^{k+1}=F_{\lambda} \phi_{\lambda, h^{\prime}}^{k}$ is $h^{\lambda}$ irrespectively of the initial condition $h^{\prime}$. Using $h^{\lambda^{\prime}}$ as the initial condition, we get $h^{\lambda}=\lim _{k \rightarrow \infty} \phi_{h^{\lambda^{\prime}}, \lambda}^{k}$.
Claim: If for some constant $\Gamma$ it holds $\left\|h-h^{\lambda^{\prime}}\right\| \leq \Gamma$ then:

$$
\begin{align*}
\left\|F_{\lambda} h-h\right\| & \leq 2 \bar{G}\left|\lambda-\lambda^{\prime}\right|+3 \Gamma,  \tag{67}\\
\left\|F_{\lambda} h-h^{\lambda^{\prime}}\right\| & \leq 2 \bar{G}\left|\lambda-\lambda^{\prime}\right|+4 \Gamma . \tag{68}
\end{align*}
$$

To prove the claim, observe that:

$$
\begin{aligned}
\left\|F_{\lambda} h-h\right\| \leq\left\|F_{\lambda} h-F_{\lambda^{\prime}} h\right\| & +\left\|F_{\lambda^{\prime}} h-F_{\lambda^{\prime}} h^{\lambda^{\prime}}\right\|+ \\
& +\left\|F_{\lambda^{\prime}} h^{\lambda^{\prime}}-h\right\|
\end{aligned}
$$

Using both inequalities in (58) we get (67). Inequality (68) is proved using (67) and triangle inequality, which completes the proof of the claim.

We then apply (67), (68) in $\phi_{h^{\lambda^{\prime}}, \lambda}^{k}$. For $k=0$, we have $\left\|\phi_{h^{\lambda^{\prime}}, \lambda}^{0}-h^{\lambda^{\prime}}\right\|=0$. Thus, applying (68) recursively we get:

$$
\left\|\phi_{h^{\lambda^{\prime}}, \lambda}^{k}-h^{\lambda^{\prime}}\right\| \leq 2 \bar{G}\left|\lambda-\lambda^{\prime}\right| \frac{4^{k}-1}{4-1}<\bar{G}\left|\lambda-\lambda^{\prime}\right| 4^{k}
$$

Hence, applying (67) it holds:

$$
\begin{aligned}
& \operatorname{sp}\left(q_{\lambda, h^{\lambda^{\prime}}}^{k}\right) \leq 2\left\|F_{\lambda} \phi_{h^{\lambda^{\prime}}, \lambda}^{k}-\phi_{h^{\lambda^{\prime}}, \lambda}^{k}\right\| \leq \\
& \quad \leq 4 \bar{G}\left|\lambda-\lambda^{\prime}\right|+6 \Gamma \bar{G}\left|\lambda-\lambda^{\prime}\right| 4^{k}<2 \Gamma \bar{G}\left|\lambda-\lambda^{\prime}\right| 4^{k+1}
\end{aligned}
$$

Then:

$$
\max _{k=0, \ldots, m} \operatorname{sp}\left(q_{\lambda, h}^{k}\right) /(1-\varepsilon) \leq \overline{2} G\left|\lambda-\lambda^{\prime}\right| 5^{m+1} /(1-\varepsilon)
$$

Therefore, using (62) with $k=0$ and the fact that $h^{\lambda}=$ $\lim _{k \rightarrow \infty} \phi_{h^{\lambda^{\prime}}, \lambda}^{k}$ we get:

$$
\left\|h^{\lambda}-h^{\lambda^{\prime}}\right\| \leq \frac{\bar{G} 5^{m+1}}{(1-\varepsilon)(1-\xi)}\left|\lambda-\lambda^{\prime}\right|=L_{1}\left|\lambda-\lambda^{\prime}\right|
$$

which completes the proof of the lemma.
Lemma 8: Denote by $v(h, \lambda)=\left(T_{\lambda} h\right)(n)$. Then, if $\lambda, \lambda^{\prime}>$ $\lambda^{\star}-\delta_{\lambda}$ and $h, h^{\prime} \in \mathbb{R}^{n}$, it holds:

$$
\begin{align*}
& \left|v(h, \lambda)-v\left(h^{\prime}, \lambda\right)\right| \leq\left\|h-h^{\prime}\right\|  \tag{69}\\
& \varepsilon_{\lambda_{2}}\left|\lambda-\lambda^{\prime}\right| \leq\left|v\left(h^{\lambda}, \lambda\right)-v\left(h^{\lambda^{\prime}}, \lambda^{\prime}\right)\right| \leq \bar{G}\left|\lambda-\lambda^{\prime}\right|,  \tag{70}\\
& \varepsilon_{\lambda_{2}}\left|\lambda-\lambda^{\star}\right| \leq\left|v\left(h^{\lambda}, \lambda\right)\right| \leq \bar{G}\left|\lambda-\lambda^{\star}\right|  \tag{71}\\
& \left(\lambda-\lambda^{\star}\right) v\left(h^{\lambda}, \lambda\right) \leq-\varepsilon_{\lambda_{2}}\left(\lambda-\lambda^{\star}\right)^{2} \tag{72}
\end{align*}
$$

Proof: Inequality (69) is a direct consequence of (57). To prove (70) consider a $\lambda^{\prime}>\lambda$. Recall that $v\left(h^{\lambda}, \lambda\right)$ and $v\left(h^{\lambda^{\prime}}, \lambda^{\prime}\right)$ are the minimum values for problems in the form (30). Denote by $\mathcal{M}$ the class of policies $\mu$ that satisfy $\lambda_{2}(\mu) \geq \varepsilon_{\lambda_{2}}$. Under any policy $\mu \in \mathcal{M}$ it holds:

$$
\begin{aligned}
& \lambda_{1}(\mu)-\lambda^{\prime} \cdot \lambda_{2}(\mu)+\left(\lambda^{\prime}-\lambda\right) \min _{\mu^{\prime} \in \mathcal{M}}\left[\lambda_{2}\left(\mu^{\prime}\right)\right] \leq \lambda_{1}(\mu)- \\
& -\lambda \cdot \lambda_{2}(\mu) \leq \lambda_{1}(\mu)-\lambda^{\prime} \cdot \lambda_{2}(\mu)+\left(\lambda^{\prime}-\lambda\right) \max _{\mu^{\prime} \in \mathcal{M}}\left[\lambda_{2}\left(\mu^{\prime}\right)\right]
\end{aligned}
$$

Taking the minimum with respect to $\mu$, we get:

$$
\begin{aligned}
v\left(h^{\lambda^{\prime}}, \lambda^{\prime}\right)+ & \left(\lambda^{\prime}-\lambda\right) \min _{\mu^{\prime} \in \mathcal{M}}\left[\lambda_{2}\left(\mu^{\prime}\right)\right] \leq v\left(h^{\lambda}, \lambda\right) \leq \\
& \leq v\left(h^{\lambda^{\prime}}, \lambda^{\prime}\right)+\left(\lambda^{\prime}-\lambda\right) \max _{\mu^{\prime}}\left[\lambda_{2}\left(\mu^{\prime}\right)\right]
\end{aligned}
$$

Furthermore, $\varepsilon_{\lambda_{2}} \leq \min _{\mu^{\prime}}\left[\lambda_{2}\left(\mu^{\prime}\right)\right] \leq \min _{\mu^{\prime}}\left[\lambda_{2}\left(\mu^{\prime}\right)\right] \leq \bar{G}$. Thus,

$$
\begin{equation*}
\varepsilon_{\lambda_{2}}\left(\lambda-\lambda^{\prime}\right) \leq v\left(h^{\lambda}, \lambda\right)-v\left(h^{\lambda^{\prime}}, \lambda^{\prime}\right) \leq \bar{G}\left(\lambda-\lambda^{\prime}\right) \tag{73}
\end{equation*}
$$

if $\lambda^{\prime}>\lambda$. Similarly:

$$
\begin{equation*}
\bar{G}\left(\lambda-\lambda^{\prime}\right) \leq v\left(h^{\lambda}, \lambda\right)-v\left(h^{\lambda^{\prime}}, \lambda^{\prime}\right) \leq \varepsilon_{\lambda_{2}}\left(\lambda-\lambda^{\prime}\right) \tag{74}
\end{equation*}
$$

if $\lambda^{\prime} \leq \lambda$. The last two equations prove (70).
To prove (71) recall that $v\left(h^{\lambda^{\star}}, \lambda^{\star}\right)=0$. Inequality (72) follows similarly by multiplying (73) or (74) by $\lambda-\lambda^{\star}$.

Lemma 9: There are positive constants $L_{2}, L_{3}$ such that:

$$
\begin{equation*}
\left|V(h, \lambda)-V\left(h, \lambda^{\prime}\right)\right| \leq L_{2}\left|\lambda-\lambda^{\prime}\right|^{2}+L_{3}\left|\lambda-\lambda^{\prime}\right|\left\|h-h^{\lambda}\right\| \tag{75}
\end{equation*}
$$

Proof: It holds:
$\left|V(h, \lambda)-V\left(h, \lambda^{\prime}\right)\right| \leq \sum_{k=0}^{K-1}\left|\left\|\phi_{\lambda, h}^{k}-h^{\lambda}\right\|^{2}-\left\|\phi_{\lambda^{\prime}, h}^{k}-h^{\lambda^{\prime}}\right\|^{2}\right|$.
Let us write each term of the summation as:
$\left|\left\|\phi_{\lambda, h}^{k}-h^{\lambda}\right\|-\left\|\phi_{\lambda^{\prime}, h}^{k}-h^{\lambda^{\prime}}\right\|\right|\left(\left\|\phi_{\lambda, h}^{k}-h^{\lambda}\right\|+\left\|\phi_{\lambda^{\prime}, h}^{k}-h^{\lambda^{\prime}}\right\|\right)$.
Applying twice the triangle inequality, the first factor of (76) can be bounded by:

$$
\left|\left\|\phi_{\lambda, h}^{k}-h^{\lambda}\right\|-\left\|\phi_{\lambda^{\prime}, h}^{k}-h^{\lambda^{\prime}}\right\|\right| \leq\left\|\phi_{\lambda, h}^{k}-\phi_{\lambda^{\prime}, h}^{k}\right\|+\left\|h^{\lambda}-h^{\lambda^{\prime}}\right\| .
$$

Claim: There is a positive constant $L_{4}$ such that:

$$
\left\|\phi_{\lambda, h}^{k}-\phi_{\lambda^{\prime}, h}^{k}\right\| \leq L_{4}\left|\lambda-\lambda^{\prime}\right|
$$

for $k=0, \ldots, K$.
To prove the claim, observe that $\left\|\phi_{\lambda, h}^{0}-\phi_{\lambda^{\prime}, h}^{0}\right\|=0$ and that for any vectors $h^{\prime}, h^{\prime \prime}$, it holds:

$$
\begin{aligned}
\left\|F_{\lambda} h^{\prime}-F_{\lambda^{\prime}} h^{\prime \prime}\right\| & \leq\left\|F_{\lambda} h^{\prime}-F_{\lambda^{\prime}} h^{\prime}\right\|+\left\|F_{\lambda^{\prime}} h^{\prime}-F_{\lambda^{\prime}} h^{\prime \prime}\right\| \\
& \leq 2 \bar{G}\left|\lambda-\lambda^{\prime}\right|+2\left\|h^{\prime}-h^{\prime \prime}\right\|
\end{aligned}
$$

To prove the claim it is sufficient to apply recursively the last inequality for $k=0, \ldots, K$.

Using the claim, we conclude that:

$$
\begin{equation*}
\left|\left\|\phi_{\lambda, h}^{k}-h^{\lambda}\right\|-\left\|\phi_{\lambda^{\prime}, h}^{k}-h^{\lambda^{\prime}}\right\|\right| \leq\left(L_{1}+L_{4}\right)\left|\lambda-\lambda^{\prime}\right|, \tag{77}
\end{equation*}
$$

for $k=0, \ldots, K$.
Using (62), (64) to the second factor of (76) we get:

$$
\begin{aligned}
\left\|\phi_{\lambda, h}^{k}-h^{\lambda}\right\|+\| \phi_{\lambda^{\prime}, h}^{k} & -h^{\lambda^{\prime}} \| \leq \frac{B_{0}}{1-\xi}\left(\left\|h-h^{\lambda}\right\|+\left\|h-h^{\lambda^{\prime}}\right\|\right) \\
& \leq \frac{B_{0}}{1-\xi}\left(2\left\|h-h^{\lambda}\right\|+L_{1}\left|\lambda-\lambda^{\prime}\right|\right)
\end{aligned}
$$

Combining the last inequality with (76) and (??) we conclude to the desired result.

Lemma 10: If $\gamma^{k}$ is small enough, then for a $\delta<\delta_{\lambda} / 2$, there is a $k_{0}$ such that the state $\left(\lambda^{k}, h^{k}\right)$ satisfies $\left\|h^{k}-h^{\lambda_{k}}\right\| \leq \delta$ and $\lambda^{k}>\lambda^{\star}-\delta$, for all $k>k_{0}$.

Proof: Refer to Figure 12. For some $k$, denote by $h=h^{k}$, $h^{+}=h^{k+1}, \lambda=\lambda^{k}, \lambda^{+}=\lambda^{k+1}$ and assume that $0>\lambda>$ $\lambda^{\star}-\delta_{\lambda}$. The Lyapunov function (65) satisfies:

$$
\begin{align*}
& V\left(h^{+}, \lambda^{+}\right)-V(h, \lambda)= \\
& =V\left(h^{+}, \lambda^{+}\right)-V\left(h^{+}, \lambda\right)+V\left(h^{+}, \lambda\right)-V(h, \lambda) \\
& \leq L_{2} \gamma^{2} v^{2}(h, \lambda)+L_{3} \gamma v(h, \lambda)\left\|h-h^{\lambda}\right\|-(1-\rho)\left\|h-h^{\lambda}\right\|^{2} \tag{78}
\end{align*}
$$

where (29), (75), and the fact that $v(h, \lambda)=\left(T_{\lambda} h\right)(n)$ were used. Furthermore, using (69), (71) we get:

$$
v(h, \lambda) \leq\left\|h-h^{\lambda}\right\|+v\left(h^{\lambda}, \lambda\right) \leq\left\|h-h^{\lambda}\right\|+\bar{G}\left|\lambda-\lambda^{\star}\right| .
$$

Substituting back to (77), we get:
$V\left(h^{+}, \lambda^{+}\right)-V(h, \lambda) \leq-\left[(1-\rho)-L_{2} \gamma^{2}+L_{3} \gamma\right]\left\|h-h^{\lambda}\right\|^{2}+$ $+\left(L_{2} \gamma+L_{3}\right) \gamma \bar{G}\left|\lambda^{\star}-\delta_{\lambda}\right|| | h-h^{\lambda} \|+L_{2} \gamma^{2} \bar{G}\left|\lambda^{\star}-\delta_{\lambda}\right|^{2}$.

The last inequality implies that if $\gamma$ is small enough, the dynamics (27)-(29) will enter the gray area in Figure 12 and remain there as soon as $\lambda>\lambda_{\star}$.

It will be show in the next subsection that the circle in Figure 12 is positively invariant (using Laypunov function(78)). Furthermore, for $\gamma$ small there is no possibility of moving from the right to the left of the cycle within the gray area in a single step. Since, the circle is positively invariant, $\lambda$ will remain greater than or equal to $\lambda^{\star}-\delta_{\lambda}$.
4) Construction of the composite Lyapunov function: Now consider the dynamics (27)-(29) and the Lyapunov function candidate:

$$
\begin{equation*}
V_{c}(h, \lambda)=V(h, \lambda)+\left(\lambda-\lambda^{\star}\right)^{2} . \tag{79}
\end{equation*}
$$

For $h=h^{k}, \lambda=\lambda^{k}$ and $h^{+}=h^{k+1}, \lambda^{+}=\lambda^{k+1}$ and assuming that $\lambda>\lambda^{\star}-\delta$ we have:

$$
\begin{align*}
& V_{c}\left(h^{+}, \lambda^{+}\right)-V_{c}(h, \lambda)=\left(V\left(h^{+}, \lambda^{+}\right)-V\left(h^{+}, \lambda\right)\right)+ \\
& +\left(V\left(h^{+}, \lambda\right)-V(h, \lambda)\right)+\left(\left(\lambda^{+}-\lambda^{\star}\right)^{2}-\left(\lambda-\lambda^{\star}\right)^{2}\right) \tag{80}
\end{align*}
$$

The first term of (79) is bounded above using Lemma 9 and (63):

$$
\begin{aligned}
& V\left(h^{+}, \lambda^{+}\right)-V\left(h^{+}, \lambda\right) \leq \\
& \quad \leq L_{2}(\gamma v(h, \lambda))^{2}+L_{3} \gamma|v(h, \lambda)|\left\|h^{+}-h^{\lambda}\right\| \\
& \quad \leq L_{2}(\gamma v(h, \lambda))^{2}+4 L_{3} \gamma|v(h, \lambda)|\left\|h-h^{\lambda}\right\|
\end{aligned}
$$

The second term of the right-hand side of (79) is bounded above by $-(1-\rho)\left\|h-h^{\lambda}\right\|^{2}$, due to (66). The last term may be rewritten as $\left[(\gamma v(h, \lambda))^{2}+2 \gamma v(h, \lambda)\left(\lambda-\lambda^{\star}\right)\right]$. Thus:

$$
\begin{align*}
& V_{c}\left(h^{+}, \lambda^{+}\right)-V_{c}(h, \lambda) \leq \\
& \quad \leq\left(L_{2}+1\right)(\gamma v(h, \lambda))^{2}+4 L_{3} \gamma|v(h, \lambda)|\left\|h-h^{\lambda}\right\|- \\
& \quad-(1-\rho)\left\|h-h^{\lambda}\right\|^{2}+2 \gamma v(h, \lambda)\left(\lambda-\lambda^{\star}\right) \tag{81}
\end{align*}
$$

Using $v(h, \lambda)=v\left(h^{\lambda}, \lambda\right)+\left(v(h, \lambda)-v\left(h^{\lambda}, \lambda\right)\right)$, along with (69), (71) and (72) we get:

$$
\begin{align*}
& (v(h, \lambda))^{2} \leq \bar{G}^{2}\left|\lambda-\lambda^{\star}\right|^{2}+ \\
& \quad+\left\|h-h^{\lambda}\right\|^{2}+2 \bar{G}\left|\lambda-\lambda^{\star}\right|\left\|h-h^{\lambda}\right\|  \tag{82}\\
& v(h, \lambda)\left(\lambda-\lambda^{\star}\right) \leq \\
& \leq-\bar{G}^{\prime}\left(\lambda-\lambda^{\star}\right)^{2}+\left\|h-h^{\lambda}\right\|\left|\lambda-\lambda^{\star}\right| \tag{83}
\end{align*}
$$

Substituting (81), (82) and the right-hand side inequality of (71) into (80) we get:

$$
\begin{aligned}
& V_{c}\left(h^{+}, \lambda^{+}\right)-V_{c}(h, \lambda) \leq \\
& \leq-\left\|h-h^{\lambda}\right\|^{2}\left[(1-\rho)+\bar{G}^{2} \gamma^{2}\left(L_{2}+1\right)\right]+ \\
& +\left\|h-h^{\lambda}\right\|\left|\lambda-\lambda^{\star}\right|\left[2 \gamma+2 \bar{G} \gamma^{2}\left(L_{2}+1\right)+4 L_{3} \gamma \bar{G}\right]- \\
& -\left|\lambda-\lambda^{\star}\right|^{2}\left[2 \bar{G}^{\prime} \gamma-\bar{G}^{2} \gamma^{2}\left(L_{2}+1\right)^{2}\right]
\end{aligned}
$$

which for appropriate positive constants $C_{1}, \ldots, C_{5}$, can be written as:

$$
\begin{aligned}
& V_{c}\left(h^{+}, \lambda^{+}\right)-V_{c}(h, \lambda) \leq-\left[\left\|h-h^{\lambda}\right\|\left|\lambda-\lambda^{\star}\right|\right] . \\
& \quad \cdot\left[\begin{array}{cc}
(1-\rho)+\gamma^{2} C_{1} & -C_{2} \gamma-\gamma^{2} C_{3} \\
-C_{2} \gamma-\gamma^{2} C_{3} & C_{4} \gamma-\gamma^{2} C_{5}
\end{array}\right]\left[\begin{array}{c}
\left\|h-h^{\lambda}\right\| \\
\left|\lambda-\lambda^{\star}\right|
\end{array}\right] .
\end{aligned}
$$

Let us denote by $A(\gamma)$ the matrix in the right-hand side of the last equation. It is not difficult to see that there is a positive


Fig. 12: The Lyapunov argument for Theorem 2
constant $\bar{\gamma}$ such that for $\gamma<\bar{\gamma}$ the matrix $A(\gamma)$ is positive definite. Furthermore, it is not difficult to see that fixing a $\underline{\gamma}$, such that $0 \leq \underline{\chi} \leq \bar{\gamma}$, there is a positive constant $\varepsilon_{0}$ such that $A(\gamma) \succeq \varepsilon_{0} I$, for all $\mathcal{1} \leq \gamma \leq \bar{\gamma}$. Therefore, $\lambda^{k} \rightarrow \lambda^{\star}$ and $h^{k} \rightarrow \bar{h}^{\lambda \star}=h^{\star} / \tau$ at a rate of a geometric progression.

Now due to Lemma 10, the state $\left(\lambda^{k}, h^{k}\right)$ will enter the gray region in Figure 12. In this region, $\lambda>\lambda^{\star}-\delta_{\lambda}$ and the Lyapunov analysis above holds true.

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[^1]:    ${ }^{1}$ Usual applications are frequency regulation, peak shaving, energy arbitrage, microgrid operation in island mode, etc. (e.g. [31], [32], [33]) .

[^2]:    ${ }^{2}$ In the implementation of both algorithms, we used a monotonicity property of the optimal control law, which is specific to this problem. This reduces the running time, but all the other results are not affected.

