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## Anran Hu

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Learning in Mean-Field Games and Continuous-Time Stochastic Control Problems

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Anran Hu


#### Abstract

Learning in Mean-Field Games and Continuous-Time Stochastic Control Problems


by

Anran Hu

Doctor of Philosophy in Engineering - Industrial Engineering and Operations Research
University of California, Berkeley
Professor Xin Guo, Chair

In recent years, there has been an ever-increasing demand for building reliable and versatile agents in applications arising from numerous fields including autonomous driving, supply chain, manufacturing, e-commerce and finance. To meet these challenging demands, researches in decision making systems have drawn upon a wide range of tools from applied probability, reinforcement learning (RL), stochastic control and game theory. This dissertation focuses on developing new methodologies and efficient algorithms with provable performance guarantees to deal with complex environments such as large population competitions and continuous-time systems.

The first part of this dissertation focuses on designing and analyzing RL algorithms for large population games. Large population games have appeared in many real-world problems. Examples include massive multiplayer online role-playing games, high frequency trading, and the sharing economy. However, in general, it becomes increasingly difficult to solve such problems as the number of players in the game grows. Mean field game (MFG) provides an ingenious and tractable aggregation approach to approximate the otherwise challenging N-player stochastic games. In Chapter 1, we present a general mean-field game (GMFG) framework for simultaneous learning and decision-making in stochastic games with a large population. It first establishes the existence of a unique Nash Equilibrium to this GMFG, and demonstrates that naively combining reinforcement learning with the fixed-point approach in classical MFGs yields unstable algorithms. It then proposes value-based and policy-based reinforcement learning algorithms (GMF-V and GMF-P, respectively) with smoothed policies, with analysis of their convergence properties and computational complexities. Experiments on an equilibrium product pricing problem demonstrate that GMF-V-Q and GMF-P-TRPO, two specific instantiations of GMF-V and GMF-P, respectively, with Q-learning and TRPO, are both efficient and robust in the GMFG setting. Moreover, their performance is superior in convergence speed, accuracy, and stability when compared with existing algorithms for multi-agent reinforcement learning in the $N$-player setting.

The second part of this dissertation focuses on designing and analyzing RL algorithms for continuous-time stochastic dynamical systems. As most physical systems in science and engineering evolve continuously in time, many real-world control tasks, such as those in aerospace, automotive industry and robotics, are naturally formulated in terms of continuoustime dynamical systems. Nevertheless, the mainstream RL algorithms have been designed for discrete-time systems, despite that they are widely applied to physical tasks in continuoustime systems. Continuous-time RL algorithms have also been developed in the past decades. But the theoretical guarantees of these works are limited to the asymptotic convergence and the non-asymptotic guarantees remain unknown. In Chapter 2, we take the first step towards designing algorithms with non-asymptotic guarantees for solving finite-time-horizon continuous-time linear quadratic (LQ) RL problems in an episodic setting, where both the state and control coefficients are unknown to the controller. We first propose a least-squares algorithm based on continuous-time observations and controls, and establish a logarithmic regret bound of magnitude $\mathcal{O}((\ln M)(\ln \ln M))$, with $M$ being the number of learning episodes. The analysis consists of two components: perturbation analysis, which exploits the regularity and robustness of the associated Riccati differential equation; and parameter estimation error, which relies on sub-exponential properties of continuous-time least-squares estimators. We further propose a practically implementable least-squares algorithm based on discrete-time observations and piecewise constant controls, which achieves similar logarithmic regret with an additional term depending explicitly on the time stepsizes used in the algorithm. In Chapter 3, we extend the results beyond linear-quadratic problems, where the unknown linear jump-diffusion process is controlled subject to non-smooth convex costs. We show that the associated linear-convex (LC) control problems admit Lipchitz continuous optimal feedback controls and further prove the Lipschitz stability of the feedback controls. The analysis relies on a stability analysis of the associated forward-backward stochastic differential equation. We then propose a least-squares algorithm which achieves a regret of the order $\mathcal{O}(\sqrt{N \ln N})$ on linear-convex learning problems with jumps, where $N$ is the number of learning episodes; the analysis leverages the Lipschitz stability of feedback controls and concentration properties of sub-Weibull random variables. Numerical experiment confirms the convergence and the robustness of the proposed algorithm.

To Bolin and Xiuzhen

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## Chapter 1

## A General Framework for Learning Mean-Field Games

### 1.1 Introduction

Motivating example. This paper is motivated by the following Ad auction problem for an advertiser. An Ad auction is a stochastic game on an Ad exchange platform among a large number of players, the advertisers. In between the time a web user requests a page and the time the page is displayed, usually within a millisecond, a Vickrey-type of second-best-price auction is run to incentivize interested advertisers to bid for an Ad slot to display advertisement. Each advertiser has limited information before each bid: first, her own valuation for a slot depends on some random conversion of clicks for the item; secondly, she, should she win the bid, only knows the reward after the users activities on the website are finished. In addition, she has a budget constraint in this repeated auction.

The question is, how should she bid in this online sequential repeated game when there is a large population of bidders competing on the Ad platform, with random conversions of clicks and rewards?

Besides Ad auctions, there are many other real-world problems involving a large number of players and uncertain systems. Examples include massive multi-player online role-playing games [91], high frequency tradings [106], and the sharing economy [79].

Our work. Motivated by these problems, we consider a general framework of simultaneous learning and decision-making in stochastic games with a large population. We formulate a general mean-field-game (GMFG) with incorporation of action distributions and (randomized) relaxed policies. This general framework can also be viewed as a generalized version of MFGs of extended McKean-Vlasov type [3], which is a different paradigm from the classical MFG. It is also beyond the scope of the existing reinforcement learning (RL) framework for Markov decision processes (MDP), as MDP is technically equivalent to a single player stochastic game.

On the theory front, this general framework differs from the existing MFGs. We establish under appropriate technical conditions the existence and uniqueness of the Nash equilibrium (NE) to this GMFG. On the computational front, we show that naively combining reinforcement learning with the three-step fixed-point approach in classical MFGs yields unstable algorithms. We then propose both value based and policy based reinforcement learning algorithms with smoothed policies (GMF-V and GMF-P, respectively), establish the convergence property and analyze the computational complexity (see Section 1.7 for all proof details). Finally, we apply GMF-V-Q and GMF-P-TRPO, which are two specific instantiations of GMF-V and GMF-P, respectively, with Q-learning and TRPO, to an equilibrium product pricing problem ${ }^{1}$. Both algorithms have demonstrated to be efficient and robust in the GMFG setting. Their performance is superior in terms of convergence speed, accuracy and stability, when compared with existing algorithms for multi-agent reinforcement learning in the $N$-player setting. Note that an earlier and preliminary version [76] has been published in NeurIPS. Nevertheless, the conference version focuses only on GMF-VQ, whereas this paper provides a new meta framework for learning mean-field-game which combines (1) the three-step fixed point approach, (2) the smoothing techniques, and (3) the single-agent algorithms with sample complexity guarantees in the sub-routine. This general framework incorporates both value-based algorithms and policy-based algorithms. In addition, the policy-based RL algorithm (GMF-P-TRPO) in this paper is the first globally convergent policy-based algorithm for solving mean-field-games. Numerical results show that it achieves similar performance as the Q-learning based algorithm (GMF-V-Q) in [76].

Related works. On learning large population games with mean-field approximations, [166] focuses on inverse reinforcement learning for MFGs without decision making, with its extension in [36] for agent-level inference; [167] studies an MARL problem with a first-order meanfield approximation term modeling the interaction between one player and all the other finite players, which has been generalized to the setting with partially observable states in [146]; and [95] and [168] consider model-based adaptive learning for MFGs in specific models (e.g., linear-quadratic and oscillator games). More recently, [115] studies the local convergence of actor-critic algorithms on finite time horizon MFGs, and [144] proposes a policy-gradient based algorithm and analyzes the so-called local NE for reinforcement learning in infinite time horizon MFGs. For learning large population games without mean-field approximation, see $[82,93]$ and the references therein. In the specific topic of learning auctions with a large number of advertisers, [28] and [92] explore reinforcement learning techniques to search for social optimal solutions with real-word data, and [90] uses MFGs to model the auction system with unknown conversion of clicks within a Bayesian framework.

However, none of these works consider the problem of simultaneous learning and decisionmaking in a general MFG framework. Neither do they establish the existence and uniqueness of the (global) NE, nor do they present model-free learning algorithms with complexity

[^0]analysis and convergence to the NE. Note that in principle, global results are harder to obtain compared to local results.

Following the conference version [76] of the current paper, various efforts have been made to extend our reinforcement learning work in [76] to more general MFG settings. These include linear-quadratic MFGs in both discrete-time setting [60, 153, 154] and in continuous-time setting [77, 163, 47], MFGs with general continuous state and/or action spaces [7], entropy regularized MFGs in discrete time [7, 164, 165, 41] and in continuous time [77], and non-stationary MFGs [117]. In particular, [41] interprets the softmax smoothing technique proposed in [76] from a smoothed equilibrium perspective. In addition, different frameworks based on monotonicity assumptions (instead of the contractivity assumption in [76]) have also been proposed, and fictitious play algorithms with policy and mean-field averaging [52, 131] and online mirror descent algorithms [129] have been proposed to solve MFGs under such assumptions. There are also some recent extensions to reinforcement learning of MFGs with strategic complementarity [105] and multiple agent types [66, 145]. These algorithms for reinforcement learning of MFGs have also been applied in economics [10], in finance [44], in animal behavior simulation [130], and in concave utility reinforcement learning [64]. In the meantime, the idea of simultaneous learning and decision making with mean-field interaction has been used for analyzing collaborative games with social optimal solution $[31,32,73,110,162,127,61,42,9]$.

Notations. Let $\left(\mathcal{X}, d_{\mathcal{X}}\right)$ be a metric space and $\mathcal{X}$ is equipped with the Borel $\sigma$-field $\mathcal{B}(\mathcal{X})$, meaning the $\sigma$-field generated by the open sets of $\mathcal{X}$. Denote $\mathcal{P}(\mathcal{X})$ for the set of (Borel) probability measures on $\mathcal{X} . \mathcal{W}_{p}$ denotes the Wasserstein distance of order $p$ such that
$W_{p}\left(\mu, \mu^{\prime}\right)=\inf \left\{\left(\int_{\mathcal{X} \times \mathcal{X}} d_{\mathcal{X}}^{p}\left(x, x^{\prime}\right) \nu\left(d x, d x^{\prime}\right)\right): \nu \in \mathcal{P}(\mathcal{X} \times \mathcal{X})\right.$ with marginals $\left.\mu, \mu^{\prime} \in \mathcal{P}(\mathcal{X})\right\}$.
$\mathcal{P}(\mathcal{X})$ is always equipped with $W_{1}\left(\mu, \mu^{\prime}\right)$. The Borel $\sigma$-field of $\mathcal{P}(\mathcal{X})$ is the $\sigma$-field induced by the evaluation $\mathcal{P}(\mathcal{X}) \ni \mu \mapsto \mu(C)$ for any Borel set $C \subset \mathcal{X}$. Note that the Borel $\sigma$-field of $\mathcal{P}(\mathcal{X})$ is generated by $W_{1}$. (See e.g. [157] and [100]).

Given two measurable spaces $(\mathcal{Y}, \mathcal{B}(\mathcal{Y}))$ and $(\mathcal{X}, \mathcal{B}(\mathcal{X}))$, we say a measure-valued function $f: \mathcal{Y} \rightarrow \mathcal{P}(\mathcal{X})$ is measurable if $\Lambda_{C} \circ f: \mathcal{Y} \rightarrow[0,1]$ is measurable for any $C \in \mathcal{B}(\mathcal{X})$, where $\Lambda_{C}: \mathcal{P}(\mathcal{X}) \ni \mu \mapsto \mu(C) \in[0,1]$.

### 1.2 Framework of General MFG (GMFG)

### 1.2.1 Background: classical $N$-player Markovian game and MFG

Let us first recall the classical $N$-player game. There are $N$ players in a game. At each step $t$, the state of player $i \quad(=1,2, \cdots, N)$ is $s_{t}^{i} \in \mathcal{S} \subseteq \mathbb{R}^{d}$ and she takes an action $a_{t}^{i} \in \mathcal{A} \subseteq \mathbb{R}^{p}$. Here $d, p$ are positive integers. The state space $\left(\mathcal{S}, d_{\mathcal{S}}\right)$ and the action space $\left(\mathcal{A}, d_{\mathcal{A}}\right)$ are two compact metric spaces, including the case of $\mathcal{S}$ and $\mathcal{A}$ being finite. Given
the current state profile of $N$-players $\boldsymbol{s}_{t}=\left(s_{t}^{1}, \ldots, s_{t}^{N}\right) \in \mathcal{S}^{N}$ and the action $a_{t}^{i}$, player $i$ will receive a reward $r^{i}\left(\boldsymbol{s}_{t}, a_{t}^{i}\right)$ sampled from a distribution $R^{i}\left(\boldsymbol{s}_{t}, a_{t}^{i}\right)$ and her state will change to $s_{t+1}^{i}$ according to a transition probability function $P^{i}\left(s_{t}, a_{t}^{i}\right)$. In particular, the probability transition $P^{i}: \mathcal{S}^{N} \times \mathcal{A} \rightarrow \mathcal{P}(\mathcal{S})$ and the distribution of the reward function $R^{i}: \mathcal{S}^{N} \times \mathcal{A} \rightarrow$ $\mathcal{P}\left(\left[0, \mathrm{R}_{\max }\right]\right)$ are both measurable functions with some constant $\mathrm{R}_{\max }>0$.

A Markovian game further restricts the admissible policy/control for player $i$ to be of the form $a_{t}^{i} \sim \pi_{t}^{i}\left(\boldsymbol{s}_{t}\right)$ with $\pi_{t}^{i}$ measurable. That is, $\pi_{t}^{i}: \mathcal{S}^{N} \rightarrow \mathcal{P}(\mathcal{A})$ maps each state profile $s \in \mathcal{S}^{N}$ to a randomized action. The accumulated reward (a.k.a. the value function) for player $i$, given the initial state profile $\boldsymbol{s}$ and the policy profile sequence $\boldsymbol{\pi}:=\left\{\boldsymbol{\pi}_{t}\right\}_{t=0}^{\infty}$ with $\boldsymbol{\pi}_{t}=\left(\pi_{t}^{1}, \ldots, \pi_{t}^{N}\right)$, is then defined as

$$
\begin{equation*}
V^{i}(s, \boldsymbol{\pi}):=\mathbb{E}\left[\sum_{t=0}^{\infty} \gamma^{t} r^{i}\left(s_{t}, a_{t}^{i}\right) \mid s_{0}=s\right], \tag{1.2.1}
\end{equation*}
$$

where $\gamma \in(0,1)$ is the discount factor, $a_{t}^{i} \sim \pi_{t}^{i}\left(s^{t}\right)$, and $s_{t+1}^{i} \sim P^{i}\left(\boldsymbol{s}_{t}, a_{t}^{i}\right)$. The goal of each player is to maximize her value function over all admissible policy sequences such that (1.2.1) is finite.

In general, this type of stochastic $N$-player game is notoriously hard to analyze, especially when $N$ is large [126]. Mean field game (MFG), pioneered by [87] and [104] in the continuous settings and later developed in [20, 68, 86, 109, 138] for discrete settings, provides an ingenious and tractable aggregation approach to approximate the otherwise challenging $N$-player stochastic games. The basic idea for an MFG goes as follows. Assume all players are identical, indistinguishable and interchangeable, when $N \rightarrow \infty$, one can view the limit of other players' states $\boldsymbol{s}_{t}^{-i}=\left(s_{t}^{1}, \ldots, s_{t}^{i-1}, s_{t}^{i+1}, \ldots, s_{t}^{N}\right)$ as a population state distribution $\mu_{t}$ with $\mu_{t}(s):=\lim _{N \rightarrow \infty} \frac{\sum_{j=1, j \neq i}^{N} \mathbf{I}_{s_{t}^{j}=s}}{N} .{ }^{2}$ Due to the homogeneity of the players, one can then focus on a single (representative) player. At time $t$, after the representative player chooses her action $a_{t}$ according to some policy $\pi_{t}$, she will receive reward $r\left(s_{t}, a_{t}, \mu_{t}\right)$ and her state will evolve under a controlled stochastic dynamics of a mean-field type $P\left(\cdot \mid s_{t}, a_{t}, \mu_{t}\right)$. Here the policy $\pi_{t}$ depends on both the current state $s_{t}$ and the current population state distribution $\mu_{t}$ such that $\pi_{t}: \mathcal{S} \times \mathcal{P}(\mathcal{S}) \rightarrow \mathcal{P}(\mathcal{A})$. Then, in mean-field limit, one may consider instead the following optimization problem,

$$
\begin{array}{ll}
\text { maximize }_{\boldsymbol{\pi}} & V(s, \boldsymbol{\pi}, \boldsymbol{\mu}):=\mathbb{E}\left[\sum_{t=0}^{\infty} \gamma^{t} r\left(s_{t}, a_{t}, \mu_{t}\right) \mid s_{0}=s\right] \\
\text { subject to } & s_{t+1} \sim P\left(s_{t}, a_{t}, \mu_{t}\right), \quad a_{t} \sim \pi_{t}\left(s_{t}, \mu_{t}\right),
\end{array}
$$

where $\boldsymbol{\pi}:=\left\{\pi_{t}\right\}_{t=0}^{\infty}$ denotes the policy sequence and $\boldsymbol{\mu}:=\left\{\mu_{t}\right\}_{t=0}^{\infty}$ the distribution flow.

### 1.2.2 General MFG (GMFG)

In the classical MFG setting, the reward and the dynamic for each player are known. They depend only on the state of the player $s_{t}$, the action of this particular player $a_{t}$, and

[^1]the population state distribution $\mu_{t}$. In contrast, in the motivating auction example, the reward and the dynamic are unknown; they rely on the actions of all players, as well as on $s_{t}$ and $\mu_{t}$.

We therefore define the following general MFG (GMFG) framework. At time $t$, after the representative player chooses her action $a_{t}$ according to some measurable policy $\pi: \mathcal{S} \times \mathcal{P}(\mathcal{S}) \rightarrow \mathcal{P}(\mathcal{A})$, she will receive a (possibly random) reward $r\left(s_{t}, a_{t}, \mathcal{L}_{t}\right)$ sampled from distribution $R\left(s_{t}, a_{t}, \mathcal{L}_{t}\right)$ and her state will evolve according to $P\left(\cdot \mid s_{t}, a_{t}, \mathcal{L}_{t}\right)$, with $\mathcal{L}_{t}=\mathbb{P}_{s_{t}, a_{t}} \in \mathcal{P}(\mathcal{S} \times \mathcal{A})$ the joint distribution of the state and the action, i.e., the population state-action pair. This joint distribution $\mathcal{L}_{t}$ has marginal distributions $\alpha_{t}$ for the population action and $\mu_{t}$ for the population state. Note the inclusion of $\alpha_{t}$ allows the reward and the dynamic to depend on all players' actions. Here $P: \mathcal{S} \times \mathcal{A} \times \mathcal{P}(\mathcal{S} \times \mathcal{A}) \rightarrow \mathcal{P}(\mathcal{S})$ and $R: \mathcal{S} \times \mathcal{A} \times \mathcal{P}(\mathcal{S} \times \mathcal{A}) \rightarrow \mathcal{P}\left(\left[0, \mathrm{R}_{\max }\right]\right)$ are measurable functions with some constant $\mathrm{R}_{\max }>0$. The objective of the player is to solve the following control problem:

$$
\begin{array}{ll}
\operatorname{maximize}_{\boldsymbol{\pi}} & V(s, \boldsymbol{\pi}, \mathcal{L}):=\mathbb{E}\left[\sum_{t=0}^{\infty} \gamma^{t} r\left(s_{t}, a_{t}, \mathcal{L}_{t}\right) \mid s_{0}=s\right]  \tag{GMFG}\\
\text { subject to } & s_{t+1} \sim P\left(s_{t}, a_{t}, \mathcal{L}_{t}\right), \quad a_{t} \sim \pi_{t}\left(s_{t}, \mu_{t}\right) .
\end{array}
$$

Here the expectation in the objective function is always taken for all randomness in the system. In addition, $\mathcal{L}:=\left\{\mathcal{L}_{t}\right\}_{t=0}^{\infty}$ which may be time dependent. That is, an infinite-time horizon MFG may have time-dependent NE solutions due to the mean information process in the MFG. This is fundamentally different from the theory of MDP where the optimal control, if exists uniquely, would be time independent in an infinite time horizon setting.

In this paper, we will analyze the existence of NE to GMFG. For ease of exposition, we will first focus on stationry NEs. Accordingly, for notational brevity, we abbreviate $\pi=\{\pi\}_{t=0}^{\infty}$ and $\mathcal{L}=\{\mathcal{L}\}_{t=0}^{\infty}$ as $\pi$ and $\mathcal{L}$, respectively. We will show in the end how this stationary constraint can be relaxed (cf. Section 1.9).

Definition 1.2.1 (Stationary NE for GMFGs). In (GMFG), a player-population profile ( $\pi^{\star}$, $\mathcal{L}^{\star}$ ) is called a stationary NE if

1. (Single player side) For any policy $\pi$ and any initial state $s \in \mathcal{S}$,

$$
\begin{equation*}
V\left(s, \pi^{\star}, \mathcal{L}^{\star}\right) \geq V\left(s, \pi, \mathcal{L}^{\star}\right) \tag{1.2.2}
\end{equation*}
$$

2. (Population side) $\mathbb{P}_{s_{t}, a_{t}}=\mathcal{L}^{\star}$ for all $t \geq 0$, where $\left\{s_{t}, a_{t}\right\}_{t=0}^{\infty}$ is the dynamics under the policy $\pi^{\star}$ starting from $s_{0} \sim \mu^{\star}$, with $a_{t} \sim \pi^{\star}\left(s_{t}, \mu^{\star}\right)$, $s_{t+1} \sim P\left(\cdot \mid s_{t}, a_{t}, \mathcal{L}^{\star}\right)$, and $\mu^{\star}$ being the population state marginal of $\mathcal{L}^{\star}$.

The single player side condition captures the optimality of $\pi^{\star}$, when the population side is fixed. The population side condition ensures the "consistency" of the solution: it guarantees that the state and action distribution flow of the single player does match the population state and action sequence $\mathcal{L}^{\star}:=\left\{\mathcal{L}^{\star}\right\}_{t=0}^{\infty}$.

### 1.2.3 Examples of GMFG

Here we provide three examples under the framework of GMFG.

A toy example. Take a two-state dynamic system with two choices of controls. The state space $\mathcal{S}=\{0,1\}$, the action space $\mathcal{A}=\{L, R\}$. Here the action $L$ means to move left and $R$ means to move right. The dynamic of the representative agent in the mean-field system $\left\{s_{t}\right\}_{t \geq 1}$ goes as follows: if the agent is in state $s_{t}$ and she takes action $a_{t}=L$ at time $t$, then $s_{t+1}=0$; if she takes action $a_{t}=R$, then $s_{t+1}=1$. At the end of each round, the agent will receive a reward $-W_{2}\left(\mu_{t}, B\right)-W_{2}\left(\beta_{t}\left(s_{t}, \cdot\right), B\right)$, which depends on all agents, where $W_{2}$ is the $\ell_{2}$-Wasserstein distance. Here $\mu_{t}(\cdot)$ denotes the state distribution of the mean-field population at time $t, \beta_{t}(s, \cdot):=\mathcal{L}_{t}(s, \cdot) / \mu_{t}(s)$ denotes the action distribution of the population in state $s(s=0,1)$ at time $t\left(\operatorname{set} \beta_{t}(s, \cdot):=(0.5,0.5)\right.$ when $\left.\mu_{t}(s)=0\right)$, and $B$ is a given Bernoulli distribution with parameter $p(0<p<1)$.

As a demonstrating example, here we provide the calculation for one stationary NE solution. Note that $-W_{2}(\mu, B) \leq 0$ for any distribution $\mu$ over $\mathcal{S}$. Similarly, $-W_{2}(\alpha, B) \leq$ 0 for any distribution $\alpha$ over $\mathcal{A}$. Hence for each policy $\boldsymbol{\pi}$, given population distribution flow $\mathcal{L}=\left\{\mathcal{L}_{t}\right\}_{t=1}^{\infty}$,

$$
\begin{equation*}
V(0, \boldsymbol{\pi}, \boldsymbol{L})=-\sum_{t=1}^{\infty} \gamma^{t} \mathbb{E}\left[W_{2}\left(\mu_{t}, B\right)+W_{2}\left(\beta_{t}\left(s_{t}, \cdot\right), B\right) \mid s_{0}=0\right] \leq 0 \tag{1.2.3}
\end{equation*}
$$

and

$$
\begin{equation*}
V(1, \boldsymbol{\pi}, \boldsymbol{L})=-\sum_{t=1}^{\infty} \gamma^{t} \mathbb{E}\left[W_{2}\left(\mu_{t}, B\right)+W_{2}\left(\beta_{t}\left(s_{t}, \cdot\right), B\right) \mid s_{0}=1\right] \leq 0 \tag{1.2.4}
\end{equation*}
$$

It is easy to check that $\mu^{\star}=(p, 1-p)$ and $\pi^{\star}\left(s, \mu^{\star}\right)=(p, 1-p)(s=0,1)$. is a pair of stationary mean-field solution. And $\mathcal{L}^{\star}$ is defined with $\mathcal{L}^{\star}(s, a)=\mu^{\star}(s) \pi^{\star}\left(a \mid s, \mu^{\star}\right)$ for any $s \in \mathcal{S}, a \in \mathcal{A}$, accordingly, where $\pi(a \mid s, \mu)$ is defined as the probability of taking action $a$ following the action distribution $\pi(s, \mu)$. In this case, the corresponding optimal value function is defined as

$$
V\left(0, \pi^{\star}, \mathcal{L}^{\star}\right)=V\left(1, \pi^{\star}, \mathcal{L}^{\star}\right)=0
$$

which reaches the upper bound in (1.2.3) and (1.2.4).

Repeated auction. Take a representative advertiser in the auction aforementioned in the motivating example in Section 1.1. Denote $s_{t} \in\left\{0,1,2, \cdots, s_{\max }\right\}$ as the budget of this player at time $t$, where $s_{\max } \in \mathbb{N}^{+}$is the maximum budget allowed on the Ad exchange with a unit bidding price. Denote $a_{t} \in\left\{0,1,2, \cdots, a_{\max }\right\}$ as the bid price submitted by this player, where $a_{\max }$ is the maximum bid set by the bidder, and $\alpha_{t}$ as the bidding/(action) distribution of the population. At time $t$, all advertisers are randomly divided into different groups and each group of advertisers competes for one slot to display their ads. Assuming
that there are $M$ advertisers in each group, then the representative advertiser competes with $M-1$ other representative players whose bidding prices are independently sampled from $\alpha_{t}$. Let $w_{t}^{M}$ denote whether the representative player wins the bid. Then if she takes action $a_{t}$, the probability she will win the bid is $\mathbb{P}\left(w_{t}^{M}=1\right)=F_{\alpha_{t}}\left(a_{t}\right)^{M-1}$, where $F_{\alpha_{t}}$ is the cumulative distribution function of a random variable $X \sim \alpha_{t}$.

If this advertiser does not win the bid, her reward $r_{t}=0$. If she wins, there are several components in her reward: $a_{t}^{M}$, the second best bid in a Vickrey auction, paid by the winning advertiser; $v_{t}$, the conversion of clicks of the slot; and $\rho$, the rate of penalty for overshooting if the payment $a_{t}^{M}$ exceeds her budget $s_{t}$. Therefore, at each time $t$, her reward with bid $a_{t}$ and budget $s_{t}$ is

$$
\begin{equation*}
r_{t}=\mathbf{I}_{\left\{w_{t}^{M}=1\right\}}\left[\left(v_{t}-a_{t}^{M}\right)-(1+\rho) \mathbf{I}_{\left\{s_{t}<a_{t}^{M}\right\}}\left(a_{t}^{M}-s_{t}\right)\right] \tag{1.2.5}
\end{equation*}
$$

where the first term is the profit of wining the auction and the second term is the penalty of overshooting. And the budget dynamics $s_{t}$ follows,

$$
s_{t+1}= \begin{cases}s_{t}, & w_{t}^{M} \neq 1  \tag{1.2.6}\\ s_{t}-a_{t}^{M}, & w_{t}^{M}=1 \text { and } a_{t}^{M} \leq s_{t} \\ 0, & w_{t}^{M}=1 \text { and } a_{t}^{M}>s_{t}\end{cases}
$$

That is, if this player does not win the bid, the budget remains the same; if she wins and has sufficient money to pay, her budget will decrease from $s_{t}$ to $s_{t}-a_{t}^{M}$; however, if she wins but does not have enough money to pay, her budget will be 0 after the payment and there will be a penalty in the reward function.

Notice that both distributions of $w_{t}^{M}$ and $a_{t}^{M}$ depend on the population distribution $\mathcal{L}_{t}$ (or more specifically $\alpha_{t}$ ). In fact, the reward function $r\left(s_{t}, a_{t}\right)=r_{t}$ and the transition probability $s_{t+1} \sim P\left(\cdot \mid s_{t}, a_{t}, \mathcal{L}_{t}\right)$ specified by (1.2.5) and (1.2.6) are fully characterized by the probabilities $\mathbb{P}\left(w_{t}^{M}=1, a_{t}^{M} \leq \cdot \mid s_{t}, a_{t}, \mathcal{L}_{t}\right)$ and $\mathbb{P}\left(w_{t}^{M}=0\right)$ (since $r_{t}=0$ and $s_{t+1}=s_{t}$ whenever $w_{t}^{M}=0$ ), with

$$
\mathbb{P}\left(w_{t}^{M}=1, a_{t}^{M} \leq x \mid s_{t}, a_{t}, \mathcal{L}_{t}\right)=F_{\alpha_{t}}\left(\min \left\{x, a_{t}\right\}\right)^{M-1}, \quad \mathbb{P}\left(w_{t}^{M}=0\right)=1-F_{\alpha_{t}}\left(a_{t}\right)^{M-1}
$$

Clearly the above model fits into the framework of (GMFG), with the following transition probability.

$$
\mathbb{P}\left(s^{\prime} \mid s, a, \mathcal{L}\right)= \begin{cases}F_{\alpha}(a)^{M-1}-F_{\alpha}(\min \{s, a\})^{M-1}, & s^{\prime}=0  \tag{1.2.7}\\ 1-F_{\alpha}(a)^{M-1}, & s^{\prime}=s \\ F_{\alpha}\left(\min \left\{s-s^{\prime}, a\right\}\right)^{M-1}-F_{\alpha}\left(\min \left\{s-s^{\prime}-1, a\right\}\right)^{M-1}, & 0<s^{\prime}<s\end{cases}
$$

where $\alpha$ is the action marginal of $\mathcal{L}$. The reward model can be explicitly written similarly.
In practice, one may modify the dynamics of $s_{t+1}$ with a non-negative random budget fulfillment $\Delta\left(s_{t+1}\right)$ after the auction clearing such that $\hat{s}_{t+1}=s_{t+1}+\Delta\left(s_{t+1}\right)[8,75]$.

Experiments of this repeated auction problem can be found in the conference version [76] of this paper, and will not be repeated here.

Equilibrium price. Another example, adapted from [74, Section 3] is to consider a large number (continuum) of homogeneous firms producing the same product under perfect competition, and the price of the product is determined endogenously by the supply-demand equilibrium [22]. Each firm, meanwhile, maintains a certain inventory level of the raw materials for production.

Given the homogeneity of the firms, it is sufficient to focus on a representative firm paired with the population distribution. In each period $t$, the representative firm decides a quantity $q_{t}$ to consume the raw materials for production and a quantity $h_{t}$ to replenish the inventory of raw materials. For simplicity, we assume each unit of the raw material is used to produce one unit of the product. Both the new products and ordered raw materials will be available at the end of this given period $t$. The representative agent makes decision based on her current inventory level of the raw material, denoted as $s_{t}$, which evolves according to

$$
\begin{equation*}
s_{t+1}=s_{t}-\min \left\{q_{t}, s_{t}\right\}+h_{t} . \tag{1.2.8}
\end{equation*}
$$

Note that if the firm overproduces and exceeds her current inventory capacity (i.e., $q_{t}>s_{t}$ ), then the firm will pay a cost for an emergency order of the raw material. Finally, the reward during this period $t$ is given by

$$
\begin{equation*}
r_{t}=\left(p_{t}-c_{0}\right) q_{t}-c_{1} q_{t}^{2}-c_{2} h_{t}-\left(c_{2}+c_{3}\right) \max \left\{q_{t}-s_{t}, 0\right\}-c_{4} s_{t} . \tag{1.2.9}
\end{equation*}
$$

Here $p_{t}$ is the selling price of the product of all firms; $c_{0}>0$ is the manufacturing cost and labor cost for making one unit of the product; $c_{1}>0$ is the quadratic cost which can be viewed as the transient price impact associated with the production level $q_{t} ; c_{2}>0$ is the cost of regular orders of the raw materials; $c_{3}>0$ is the additional cost for the emergency order of the raw materials; and finally, $c_{4}>0$ is the inventory cost.

The price $p_{t}$ is determined according to the supply-demand equilibrium on the market at each moment. On one hand, the normalized demand (per producer) on the market $D\left(p_{t}\right)$ follows ([74])

$$
\begin{equation*}
D\left(p_{t}\right):=d p_{t}^{-\sigma}, \tag{1.2.10}
\end{equation*}
$$

where $d$ denotes some benchmark demand level and $\sigma$ is the elasticity of demand that can be interpreted as the elasticity of substitution between the given product and any other good. On the other hand, the (average) supply in this market is given by the average production of all firms which follows $\mathbb{E}_{q_{t} \sim \pi_{t}}\left[q_{t}\right]$ under some policy $\pi_{t}$. If all firms are restricted to stationary policies (denoted as $\pi$ ), then this leads to a stationary equilibrium price $q$ which satisfies the supply-demand equilibrium:

$$
\begin{equation*}
\mathbb{E}_{q \sim \pi}[q]=d p^{-\sigma} \tag{1.2.11}
\end{equation*}
$$

To fit into the theoretical framework proposed in Section 1.2, we set $\mathcal{S}=\{0,1, \cdots, S\}$ and $\mathcal{A}=\{(q, h) \mid q \in\{0,1, \cdots, Q\}$ and $h \in\{0,1, \cdots, H\}\}$ for some positive integers $S, Q$ and $H$.

### 1.3 Solution for GMFGs

We now establish the existence and uniqueness of the stationary NE to (GMFG), by generalizing the classical fixed-point approach for MFGs to this GMFG setting. (See [87] and [104] for the classical case.) It consists of three steps.

Step A. Fix $\mathcal{L}$, (GMFG) becomes the classical single-player optimization problem. Indeed, with $\mathcal{L}$ fixed, the population state distribution $\mu$ is also fixed, and hence the space of admissible policies is reduced to the single-player case. Solving (GMFG) is now reduced to finding a policy $\pi_{\mathcal{L}}^{\star} \in \Pi:=\{\pi \mid \pi: \mathcal{S} \rightarrow \mathcal{P}(\mathcal{A})\}$ to maximize

$$
\begin{aligned}
V\left(s, \pi_{\mathcal{L}}, \mathcal{L}\right):= & \mathbb{E}\left[\sum_{t=0}^{\infty} \gamma^{t} r\left(s_{t}, a_{t}, \mathcal{L}\right) \mid s_{0}=s\right] \\
\text { subject to } & s_{t+1} \sim P\left(s_{t}, a_{t}, \mathcal{L}\right), \quad a_{t} \sim \pi_{\mathcal{L}}\left(s_{t}\right) .
\end{aligned}
$$

Notice that with $\mathcal{L}$ fixed, one can safely suppress the dependency on $\mu$ in the admissible policies.

Now given this fixed $\mathcal{L}$ and the solution $\pi_{\mathcal{L}}^{\star}$ to the above optimization problem, one can define a mapping from the fixed population distribution $\mathcal{L}$ to a chosen optimal randomized policy sequence. That is,

$$
\Gamma_{1}: \mathcal{P}(\mathcal{S} \times \mathcal{A}) \rightarrow \Pi
$$

such that $\pi_{\mathcal{L}}^{\star}=\Gamma_{1}(\mathcal{L})$. Note that the optimal policy of an MDP in general may not be unique. To ensure that $\Gamma_{1}$ is a single-valued instead of set-valued mapping, here $\Gamma_{1}$ includes a policy selection component to select a single optimal policy from the set of optimal policies for a given $\mathcal{L}$, which is guaranteed to exist by Zermelo's Axiom of Choice. For example, when the action space is finite, one can utilize the argmax-e operator and set the "maximizing" actions with equal probabilities (see Section 1.4.1 for the detailed definition). In addition, for non-degenerate linear-quadratic MFGs [60] and general MFGs where the Bellman mappings are strongly concave in actions [7] and the action space is convex in the Euclidean space, the optimal policy $\pi_{\mathcal{L}}^{\star}$ for a given $\mathcal{L}$ is unique under appropriate assumptions. Hence no policy selection is needed in such cases.

Note that this $\pi_{\mathcal{L}}^{\star}$ satisfies the single player side condition in Definition 1.2.1 for the population state-action pair $\mathcal{L}$,

$$
\begin{equation*}
V\left(s, \pi_{\mathcal{L}}^{\star}, \mathcal{L}\right) \geq V(s, \pi, \mathcal{L}) \tag{1.3.1}
\end{equation*}
$$

for any policy $\pi$ and any initial state $s \in \mathcal{S}$.
As in the MFG literature [87], a feedback regularity condition is needed for analyzing Step A.

Assumption 1. There exists a constant $d_{1} \geq 0$, such that for any $\mathcal{L}, \mathcal{L}^{\prime} \in \mathcal{P}(\mathcal{S} \times \mathcal{A})$,

$$
\begin{equation*}
D\left(\Gamma_{1}(\mathcal{L}), \Gamma_{1}\left(\mathcal{L}^{\prime}\right)\right) \leq d_{1} W_{1}\left(\mathcal{L}, \mathcal{L}^{\prime}\right) \tag{1.3.2}
\end{equation*}
$$

where

$$
\begin{equation*}
D\left(\pi, \pi^{\prime}\right):=\sup _{s \in \mathcal{S}} W_{1}\left(\pi(s), \pi^{\prime}(s)\right) \tag{1.3.3}
\end{equation*}
$$

and $W_{1}$ is the $\ell_{1}$-Wasserstein distance (a.k.a. earth mover distance) between probability measures [67, 132, 157].

Step B. Given $\pi_{\mathcal{L}}^{\star}$ obtained from Step A, update the initial $\mathcal{L}$ to $\mathcal{L}^{\prime}$ following the controlled dynamics $P\left(\cdot \mid s_{t}, a_{t}, \mathcal{L}\right)$.

Accordingly, for any admissible policy $\pi \in \Pi$ and a joint population state-action pair $\mathcal{L} \in \mathcal{P}(\mathcal{S} \times \mathcal{A})$, define a mapping $\Gamma_{2}: \Pi \times \mathcal{P}(\mathcal{S} \times \mathcal{A}) \rightarrow \mathcal{P}(\mathcal{S} \times \mathcal{A})$ as follows:

$$
\begin{equation*}
\Gamma_{2}(\pi, \mathcal{L}):=\hat{\mathcal{L}}=\mathbb{P}_{s_{1}, a_{1}}, \tag{1.3.4}
\end{equation*}
$$

where $a_{1} \sim \pi\left(s_{1}\right), s_{1} \sim \mu P\left(\cdot \mid \cdot, a_{0}, \mathcal{L}\right), a_{0} \sim \pi\left(s_{0}\right), s_{0} \sim \mu$, and $\mu$ is the population state marginal of $\mathcal{L}$.

One needs a standard assumption in this step.
Assumption 2. There exist constants $d_{2}, d_{3} \geq 0$, such that for any admissible policies $\pi, \pi_{1}, \pi_{2}$ and joint distributions $\mathcal{L}, \mathcal{L}_{1}, \mathcal{L}_{2}$,

$$
\begin{gather*}
W_{1}\left(\Gamma_{2}\left(\pi_{1}, \mathcal{L}\right), \Gamma_{2}\left(\pi_{2}, \mathcal{L}\right)\right) \leq d_{2} D\left(\pi_{1}, \pi_{2}\right)  \tag{1.3.5}\\
W_{1}\left(\Gamma_{2}\left(\pi, \mathcal{L}_{1}\right), \Gamma_{2}\left(\pi, \mathcal{L}_{2}\right)\right) \leq d_{3} W_{1}\left(\mathcal{L}_{1}, \mathcal{L}_{2}\right) . \tag{1.3.6}
\end{gather*}
$$

Step C. Repeat Step A and Step B until $\mathcal{L}^{\prime}$ matches $\mathcal{L}$.
This step is to ensure the population side condition. To ensure the convergence of the combined step one and step two, it suffices if $\Gamma: \mathcal{P}(\mathcal{S} \times \mathcal{A}) \rightarrow \mathcal{P}(\mathcal{S} \times \mathcal{A})$ with $\Gamma(\mathcal{L}):=$ $\Gamma_{2}\left(\Gamma_{1}(\mathcal{L}), \mathcal{L}\right)$ is a contractive mapping under the $W_{1}$ distance. Then by the Banach fixed point theorem and the completeness of the related metric spaces (cf. Section 1.10.1), there exists a unique stationary NE of the GMFG. That is,

Theorem 1.3.1 (Existence and Uniqueness of stationary GMFG solution). Given Assumptions 1 and 2, and assume $d_{1} d_{2}+d_{3}<1$. Then there exists a unique stationary $N E$ to (GMFG).

Proof. [Proof of Theorem 1.3.1] First by Definition 1.2 .1 and the definitions of $\Gamma_{i}(i=1,2)$, $(\pi, \mathcal{L})$ is a stationary NE iff $\mathcal{L}=\Gamma(\mathcal{L})=\Gamma_{2}\left(\Gamma_{1}(\mathcal{L}), \mathcal{L}\right)$ and $\pi=\Gamma_{1}(\mathcal{L})$, where $\Gamma(\mathcal{L})=$ $\Gamma_{2}\left(\Gamma_{1}(\mathcal{L}), \mathcal{L}\right)$. This indicates that for any $\mathcal{L}_{1}, \mathcal{L}_{2} \in \mathcal{P}(\mathcal{S} \times \mathcal{A})$,

$$
\begin{align*}
& W_{1}\left(\Gamma\left(\mathcal{L}_{1}\right), \Gamma\left(\mathcal{L}_{2}\right)\right)=W_{1}\left(\Gamma_{2}\left(\Gamma_{1}\left(\mathcal{L}_{1}\right), \mathcal{L}_{1}\right), \Gamma_{2}\left(\Gamma_{1}\left(\mathcal{L}_{2}\right), \mathcal{L}_{2}\right)\right) \\
& \leq W_{1}\left(\Gamma_{2}\left(\Gamma_{1}\left(\mathcal{L}_{1}\right), \mathcal{L}_{1}\right), \Gamma_{2}\left(\Gamma_{1}\left(\mathcal{L}_{2}\right), \mathcal{L}_{1}\right)\right)+W_{1}\left(\Gamma_{2}\left(\Gamma_{1}\left(\mathcal{L}_{2}\right), \mathcal{L}_{1}\right), \Gamma_{2}\left(\Gamma_{1}\left(\mathcal{L}_{2}\right), \mathcal{L}_{2}\right)\right)  \tag{1.3.7}\\
& \leq\left(d_{1} d_{2}+d_{3}\right) W_{1}\left(\mathcal{L}_{1}, \mathcal{L}_{2}\right)
\end{align*}
$$

And since $d_{1} d_{2}+d_{3} \in[0,1)$, by the Banach fixed-point theorem, we conclude that there exists a unique fixed-point of $\Gamma$, or equivalently, a unique stationary MFG solution to (GMFG).

Remark 1.3.1 (Existence and Uniqueness of the GMFG solution). (1) In general, there may multiple optimal policies in Step $A$ under a fixed mean-field information $\mathcal{L}$. In this case, the candidate fixed point(s) are the fixed point(s) of a set-valued map as described in [100]. To simplify the analysis, we specify a rule in Step A to select one optimal policy to ensure that $\Gamma$ is an injection.
(2) In the MFG literature, the uniqueness of the $M F G$ solution can be verified under the small parameter condition [29] or the monotonicity condition [104]. Our condition of $d_{1} d_{2}+d_{3}<1$ extends the small parameter condition in [29] for strict controls to relaxed controls.
(3) Finally, Theorem 1.3 .1 can be extended to a non-stationary setting, as will be shown in Section 1.9.

Remark 1.3.2. Assumptions 1 and 2 can be more explicit in specific problem settings.
For instance, when the action space is the Euclidean space or its convex subset, explicit conditions on $P$ and $r$ have been described for the linear-quadratic $M F G$ ( $L Q-M F G$ ) [60] and later generalized in [7].

When the action space is finite, the following lemma explicitly characterizes Assumption 2.

Lemma 1.3.2. Suppose that $\max _{s, a, \mathcal{L}, s^{\prime}} P\left(s^{\prime} \mid s, a, \mathcal{L}\right) \leq c_{1}$, and that $P\left(s^{\prime} \mid s, a, \cdot\right)$ is $c_{2}$-Lipschitz in $W_{1}$, i.e.,

$$
\begin{equation*}
\left|P\left(s^{\prime} \mid s, a, \mathcal{L}_{1}\right)-P\left(s^{\prime} \mid s, a, \mathcal{L}_{2}\right)\right| \leq c_{2} W_{1}\left(\mathcal{L}_{1}, \mathcal{L}_{2}\right) \tag{1.3.8}
\end{equation*}
$$

Then in Assumption 2, $d_{2}$ and $d_{3}$ can be chosen as

$$
\begin{equation*}
d_{2}=\frac{2 \operatorname{diam}(\mathcal{S}) \operatorname{diam}(\mathcal{A})|\mathcal{S}| c_{1}}{d_{\min }(\mathcal{A})} \tag{1.3.9}
\end{equation*}
$$

and $d_{3}=\frac{\operatorname{diam}(\mathcal{S}) \operatorname{diam}(\mathcal{A}) c_{2}}{2}$, respectively. Here $d_{\min }(\mathcal{A})=\min _{a \neq a^{\prime} \in \mathcal{A}}\left\|a-a^{\prime}\right\|_{2}$, which is guaranteed to be positive when $\mathcal{A}$ is finite.

When entropy regularization is introduced into the system (see e.g., [7, 164]), Assumption 1 can be reduced to boundedness and Lipschitz continuity conditions on P and r as in Lemma 1.3.2. Moreover, Theorem 1.3 .1 and all subsequent theoretical results hold whenever the composed mapping $\Gamma$ is contractive (in $W_{1}$ ), independent of Assumptions 1 or 2. In Section 1.8.2, we numerically verify that the $\Gamma$ mapping is contractive for various choices of the model parameters in our tested problems.

### 1.4 Naive algorithm and stabilization techniques

In this section, we design algorithms for the GMFG. Since the reward and transition distributions are unknown, this is simultaneously learning the system and finding the NE of
the game. We will focus on the case with finite state and action spaces, i.e., $|\mathcal{S}|,|\mathcal{A}|<\infty$. We will look for stationary (time independent) NEs. This stationarity property enables developing appropriate stationary reinforcement learning algorithms, suitable for an infinite time horizon game. Instead of knowing the transition probability $P$ and the reward $r$ explicitly, the algorithms we propose only assume access to a simulator oracle, which is described below. This is not restrictive in practice. For instance, in the ad auction example, one may adopt the bid recommendation perspective of the publisher, say Google, Facebook or Amazon, who acts as the auctioneer and owns the Ad slot inventory on its own Ad exchange platform. In this case, a high quality auction simulator is typically built and maintained by a team of the publisher. See also [144] for more examples.

Simulator oracle. For any policy $\pi \in \Pi$, given the current state $s \in \mathcal{S}$, for any population distribution $\mathcal{L}$, one can obtain a sample of the next state $s^{\prime} \sim P(\cdot \mid s, \pi(s), \mathcal{L})$, a reward $r=r(s, \pi(s), \mathcal{L})$, and the next population distribution $\mathcal{L}^{\prime}=\mathbb{P}_{s^{\prime}, \pi\left(s^{\prime}\right)}$. For brevity, we denote the simulator as $\left(s^{\prime}, r, \mathcal{L}^{\prime}\right)=\mathcal{G}(s, \pi, \mathcal{L})$. This simulator oracle can be weakened to fit the $N$-player setting, see Section 1.6.

In the following, we begin with a naive algorithm that simply combines the three-step fixed point approach with general RL algorithms, and demonstrate that this algorithm can be unstable (Section 1.4.1). We then propose some smoothing and projection techniques to resolve the issue (Section 1.4.2). In Section 1.5.1 and Section 1.5.2, we design general value-based and policy-based RL algorithms, and establish the corresponding convergence and complexity results. These two algorithms include most of the RL algorithms in the literature. We then illustrate by two concrete examples based on Q-learning and trust-region policy optimization algorithms.

### 1.4.1 Naive algorithm and its issue

We follow the three-step fixed-point approach described in Section 1.3. Notice the fact that with $\mathcal{L}$ fixed, Step A in Section 1.3 becomes a standard learning problem for an infinite horizon discounted MDP. More specifically, the MDP to be solved is $\mathcal{M}_{\mathcal{L}}=$ $\left(\mathcal{S}, \mathcal{A}, P_{\mathcal{L}}, r_{\mathcal{L}}, \gamma\right)$, where $P_{\mathcal{L}}\left(s^{\prime} \mid s, a\right)=P\left(s^{\prime} \mid s, a, \mathcal{L}\right)$ and $r_{\mathcal{L}}(s, a)=r(s, a, \mathcal{L})$. In general, for an $\operatorname{MDP} \mathcal{M}=(\mathcal{S}, \mathcal{A}, P, r, \gamma)$, for any policy $\pi$ one can define its value functions $V_{\mathcal{M}}^{\pi}(s)=$ $\mathbb{E}\left[\sum_{t=0}^{\infty} \gamma^{t} r\left(s_{t}, a_{t}\right) \mid s_{0}=s\right]$ and its Q-functions $Q_{\mathcal{M}}^{\pi}(s, a)=\mathbb{E}\left[\sum_{t=0}^{\infty} \gamma^{t} r\left(s_{t}, a_{t}\right) \mid s_{0}=s, a_{0}=a\right]$, where $s_{t}, a_{t}$ is the trajectory under policy $\pi$. One can also define the optimal Q-function as the unique solution of the Bellman equation:

$$
Q_{\mathcal{M}}^{\star}(s, a)=\mathbb{E}[r(s, a)]+\gamma \max _{a^{\prime}} \sum_{s^{\prime} \in \mathcal{S}} P\left(s^{\prime} \mid s, a\right) Q_{\mathcal{M}}^{\star}\left(s^{\prime}, a^{\prime}\right)
$$

for all $s, a$ and its optimal value function $V_{\mathcal{M}}^{\star}(s)=\max _{a} Q_{\mathcal{M}}^{\star}(s, a)$ for all $s$. We also use the shorthand $V_{\mathcal{L}}^{\star}=V_{\mathcal{M}_{\mathcal{L}}}^{\star}$ and $Q_{\mathcal{L}}^{\star}=Q_{\mathcal{M}_{\mathcal{L}}}^{\star}$ for notational brevity. Whenever the context is clear, we may omit $\mathcal{M}, \mathcal{L}$ and $\mathcal{M}_{\mathcal{L}}$ for notational convenience.

Given the optimal Q-function $Q_{\mathcal{L}}^{\star}$, one can obtain an optimal policy $\pi_{\mathcal{L}}^{\star}$ with $\pi_{\mathcal{L}}^{\star}(s)=$ $\operatorname{argmax}-\mathbf{e}\left(Q_{\mathcal{L}}^{\star}(s, \cdot)\right)$. Here the argmax-e operator is defined so that actions with equal maximum Q-values would have equal probabilities to be selected. Hereafter, we specify $\Gamma_{1}$ as a mapping to the aforementioned choice of the optimal policy, i.e., the s-component $\Gamma_{1}(\mathcal{L})_{s}=\operatorname{argmax}-\mathbf{e}\left(Q_{\mathcal{L}}^{\star}(s, \cdot)\right)$ for any $s \in \mathcal{S}$.

The population update in Step B can then be directly obtained from the simulator $\mathcal{G}$ following policy $\pi_{\mathcal{L}}^{\star}$. Combining these two steps leads to the following naive algorithm (Algorithm 1).

```
Algorithm 1 Naive Reinforcement Learning for GMFGs
    Input: Initial population state-action pair \(\mathcal{L}_{0}\)
    for \(k=0,1, \cdots\) do
        Obtain the optimal Q-function \(Q_{k}(s, a)=Q_{\mathcal{L}_{k}}^{\star}(s, a)\) of an MDP with dynamics
        \(P_{\mathcal{L}_{k}}\left(s^{\prime} \mid s, a\right)\) and reward distributions \(R_{\mathcal{L}_{k}}(s, a)\).
        Compute \(\pi_{k} \in \Pi\) with \(\pi_{k}(s)=\operatorname{argmax}-\mathbf{e}\left(Q_{k}(s, \cdot)\right)\).
        Sample \(s \sim \mu_{k}\), where \(\mu_{k}\) is the population state marginal of \(\mathcal{L}_{k}\), and obtain \(\mathcal{L}_{k+1}\) from
        \(\mathcal{G}\left(s, \pi_{k}, \mathcal{L}_{k}\right)\).
    end for
```

Unfortunately, in practice, one cannot obtain the exact optimal Q-function $Q_{k}$. In fact, invoking any commonly used RL algorithm with the simulator $\mathcal{G}$ leads to an approximation $\hat{Q}_{k}$ of the actual $Q_{k}$. This approximation error is then magnified by the discontinuous and sensitive argmax-e, which eventually leads to an unstable algorithm (see Figure 1.4 for an example of divergence). To see why argmax-e is not continuous, consider the following simple example. Let $x=(1,1)$, then $\operatorname{argmax}-\mathbf{e}(x)=(1 / 2,1 / 2)$. For any $\epsilon>0$, let $y_{\epsilon}=(1,1-\epsilon)$, then argmax-e $\left(y_{\epsilon}\right)=(1,0)$. Hence $\lim _{\epsilon \rightarrow 0} y_{\epsilon}=x$ but

$$
\lim _{\epsilon \rightarrow 0} \operatorname{argmax}-\mathbf{e}\left(y_{\epsilon}\right) \neq \operatorname{argmax}-\mathbf{e}(x) .
$$

This instability issue will be addressed by introducing smoothing and projection techniques.

### 1.4.2 Restoring stability

Smoothing techniques. To address the instability caused, we replace argmax-e with a smooth function that is a good approximation to argmax-e while being Lipschitz continuous. One such candidate is the softmax operator $\operatorname{softmax}_{c}: \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$, with

$$
\operatorname{softmax}_{c}(x)_{i}=\frac{\exp \left(c x_{i}\right)}{\sum_{j=1}^{n} \exp \left(c x_{j}\right)}, \quad i=1, \ldots, n
$$

for some positive constant $c$. The resulting policies are sometimes called Boltzmann policies, and are widely used in the literature of reinforcement learning [11, 78].

The softmax operator can be generalized to a wide class of operators. In fact, for positive constants $c, c^{\prime}>0$, one can consider a parametrized family $\mathcal{F}_{c, c^{\prime}} \subseteq\left\{f_{c, c^{\prime}}: \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}\right\}$ of all "smoothed" argmax-e's, i.e., all $f_{c, c^{\prime}}: \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ that satisfies the following two conditions:

- Condition 1: $f_{c, c^{\prime}}$ is $c$-Lipschitz, i.e., $\left\|f_{c, c^{\prime}}(x)-f_{c, c^{\prime}}(y)\right\|_{2} \leq c\|x-y\|_{2}$.
- Condition 2: $f_{c, c^{\prime}}$ is a good approximation of argmax-e, i.e.,

$$
\left\|f_{c, c^{\prime}}(x)-\operatorname{argmax}-\mathbf{e}(x)\right\|_{2} \leq 2 n \exp \left(-c^{\prime} \delta\right)
$$

where $\delta=x_{\text {max }}-\max _{x_{j}<x_{\max }} x_{j}, x_{\max }=\max _{i=1, \ldots, n} x_{i}$, and $\delta:=\infty$ when all $x_{j}$ are equal.

Notice that $\mathcal{F}_{c, c^{\prime}}$ is closed under convex combinations, i.e., if $f_{c, c^{\prime}}, g_{c, c^{\prime}} \in \mathcal{F}_{c, c^{\prime}}$, then for any $\lambda \in[0,1], \lambda f_{c, c^{\prime}}+(1-\lambda) g_{c, c}$ also satisfies the two conditions. Hence $\mathcal{F}_{c, c^{\prime}}$ is convex.

To have a better idea of what $\mathcal{F}_{c, c^{\prime}}$ looks like, we describe a subset $\mathcal{B}_{c, c^{\prime}}$ of $\mathcal{F}_{c, c^{\prime}}$ consisting of the generalized softmax operator softmax $\max _{h}: \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$, defined as

$$
\begin{equation*}
\operatorname{softmax}_{h}(x)_{i}=\frac{\exp \left(h\left(x_{i}\right)\right)}{\sum_{j=1}^{n} \exp \left(h\left(x_{j}\right)\right)}, \quad i=1, \ldots, n \tag{1.4.1}
\end{equation*}
$$

where $h: \mathbb{R} \rightarrow \mathbb{R}$ satisfies $c^{\prime}(x-y) \leq h(x)-h(y) \leq c(x-y)$ for any $x \geq y$. When $h$ is continuously differentiable, a sufficient condition is that $c^{\prime} \leq h^{\prime}(x) \leq c$. In particular, if $h(x) \equiv c x$ for some constant $c>0$, the operator reduces to the classical softmax operator, in which case we overload the notation to write $\boldsymbol{s o f t m a x}_{h}$ as $\boldsymbol{\operatorname { s o f t m a x }}_{c}$.

This operator is Lipschitz continuous and close to the argmax-e (see Lemmas 1.7.2 and 1.7.3 in Section 3.6), and in particular one can show that $\mathcal{B}_{c, c^{\prime}} \subseteq \mathcal{F}_{c, c^{\prime}}$. As a result, even though smoothed (e.g., Boltzmann) policies are not optimal, the difference between the smoothed and the optimal one can always be controlled by choosing a function $h$ with appropriate parameters $c, c^{\prime}$. Note that other smoothing operators (e.g., Mellowmax [11], which is a softmax operator with time-varying and problem dependent temperatures) may also be considered.

Error control in updating $\mathcal{L}$. Given the sub-optimality of the smoothed policy, one needs to characterize the difference between the optimal policy and the non-optimal ones. In particular, one can define the action gap between the best and the second best actions in terms of the Q -value as

$$
\delta^{s}(\mathcal{L}):=\max _{a^{\prime} \in \mathcal{A}} Q_{\mathcal{L}}^{\star}\left(s, a^{\prime}\right)-\max _{a \notin \operatorname{argmax}_{a \in \mathcal{A}} Q_{\mathcal{L}}^{\star}(s, a)} Q_{\mathcal{L}}^{\star}(s, a)>0 .
$$

Action gap is important for approximation algorithms [19], and is closely related to the problem-dependent bounds for regret analysis in reinforcement learning and multi-armed bandits, and advantage learning algorithms including A3C [116].

The problem is: in order for the learning algorithm to converge in terms of $\mathcal{L}$ (Theorems 1.5.1 and 1.5.5), one needs to ensure a definite differentiation between the optimal policy and the sub-optimal ones. This is problematic as the infimum of $\delta^{s}(\mathcal{L})$ over an infinite number of $\mathcal{L}$ can be 0 . To address this, the population distribution at step $k$, say $\mathcal{L}_{k}$, needs to be projected to a finite grid, called $\epsilon$-net. The relation between the $\epsilon$-net and action gaps is as follows:

For any $\epsilon>0$, there exist a positive function $\phi(\epsilon)$ and an $\epsilon$-net $S_{\epsilon}:=\left\{\mathcal{L}^{(1)}, \ldots, \mathcal{L}^{\left(N_{\epsilon}\right)}\right\}$ $\subseteq \mathcal{P}(\mathcal{S} \times \mathcal{A})$, with the properties that $\min _{i=1, \ldots, N_{\epsilon}} d_{T V}\left(\mathcal{L}, \mathcal{L}^{(i)}\right) \leq \epsilon$ for any $\mathcal{L} \in \mathcal{P}(\mathcal{S} \times \mathcal{A})$, and that $\max _{a^{\prime} \in \mathcal{A}} Q_{\mathcal{L}^{(i)}}^{\star}\left(s, a^{\prime}\right)-Q_{\mathcal{L}^{(i)}}^{\star}(s, a) \geq \phi(\epsilon)$ for any $i=1, \ldots, N_{\epsilon}, s \in \mathcal{S}$, and any $a \notin \operatorname{argmax}_{a \in \mathcal{A}} Q_{\mathcal{L}^{(i)}}^{\star}(s, a)$.

Here the existence of $\epsilon$-nets is trivial due to the compactness of the probability simplex $\mathcal{P}(\mathcal{S} \times \mathcal{A})$, and the existence of $\phi(\epsilon)$ comes from the finiteness of the action set $\mathcal{A}$.

In practice, $\phi(\epsilon)$ often takes the form of $D \epsilon^{\alpha}$ with $D>0$ and the exponent $\alpha>0$ characterizing the decay rate of the action gaps. In general, experiments are robust with respect to the choice of $\epsilon$-net.

In the next section, we propose value based and policy based algorithms for learning GMFG.

### 1.5 RL Algorithms for (stationary) GMFGs

### 1.5.1 Value-based algorithms

We start by introducing the following definition.
Definition 1.5.1 (Value-based Guarantee). For an arbitrary MDP $\mathcal{M}$, we say that an algorithm has a value-based guarantee with parameters $\left\{C_{\mathcal{M}}^{(i)}, \alpha_{1}^{(i)}, \alpha_{2}^{(i)}, \alpha_{3}^{(i)}, \alpha_{4}^{(i)}\right\}_{i=1}^{m}$, if for any $\epsilon, \delta>0$, after obtaining

$$
\begin{equation*}
T_{\mathcal{M}}(\epsilon, \delta)=\sum_{i=1}^{m} C_{\mathcal{M}}^{(i)}\left(\frac{1}{\epsilon}\right)^{\alpha_{1}^{(i)}}\left(\log \frac{1}{\epsilon}\right)^{\alpha_{2}^{(i)}}\left(\frac{1}{\delta}\right)^{\alpha_{3}^{(i)}}\left(\log \frac{1}{\delta}\right)^{\alpha_{4}^{(i)}} \tag{1.5.1}
\end{equation*}
$$

samples from the simulator oracle $\mathcal{G}$, with probability at least $1-2 \delta$, it outputs an approximate Q-function $\hat{Q}^{T_{\mathcal{M}}(\epsilon, \delta)}$ which satisfies $\left\|\hat{Q}^{T_{\mathcal{M}}(\epsilon, \delta)}-Q^{\star}\right\|_{\infty} \leq \epsilon$. Here the norm $\|\cdot\|_{\infty}$ is understood element-wisely.

### 1.5.1.1 GMF-V

We now state the first main algorithm (Algorithm 2). It applies to any algorithm Alg with a value-based guarantee.

Here $\operatorname{Proj}_{S_{\epsilon}}(\mathcal{L})=\operatorname{argmin}_{\mathcal{L}^{(1)}, \ldots, \mathcal{L}^{\left(N_{\epsilon}\right)}} d_{T V}\left(\mathcal{L}^{(i)}, \mathcal{L}\right)$. For computational tractability, it is sufficient to choose $S_{\epsilon}$ as a truncation grid so that projection of $\tilde{\mathcal{L}}_{k}$ onto the $\epsilon$-net reduces to

```
Algorithm 2 GMF-V (Alg, \(\left.f_{c, c^{\prime}}\right)\)
    Input: Initial \(\mathcal{L}_{0}, \epsilon\)-net \(S_{\epsilon}\), temperatures \(c, c^{\prime}>0\), tolerances \(\epsilon_{k}, \delta_{k}>0, k=0,1, \ldots\).
    for \(k=0,1, \cdots\) do
        Apply \(\operatorname{Alg}\) to find the approximate Q-function \(\hat{Q}_{k}^{\star}=\hat{Q}^{T_{k}}\) of the MDP \(\mathcal{M}_{\mathcal{L}_{k}}\), where
        \(T_{k}=T_{\mathcal{M}_{\mathcal{L}_{k}}}\left(\epsilon_{k}, \delta_{k}\right)\).
        Compute \(\pi_{k}(s)=f_{c, c^{\prime}}\left(\hat{Q}_{k}^{\star}(s, \cdot)\right)\).
        Sample \(s \sim \mu_{k}\left(\mu_{k}\right.\) is the population state marginal of \(\left.\mathcal{L}_{k}\right)\), obtain \(\tilde{\mathcal{L}}_{k+1}\) from
        \(\mathcal{G}\left(s, \pi_{k}, \mathcal{L}_{k}\right)\).
        Find \(\mathcal{L}_{k+1}=\operatorname{Proj}_{S_{\epsilon}}\left(\tilde{\mathcal{L}}_{k+1}\right)\)
    end for
```

truncating $\tilde{\mathcal{L}}_{k}$ to a certain number of digits. For instance, in experiments in Section 1.8, the number of digits is chosen to be 4. Appropriate choices of the hyper-parameters $c, c^{\prime}, \epsilon$ and tolerances $\epsilon_{k}, \delta_{k}(k \geq 0)$ are given in Theorems 1.5.1. Our experiment shows the algorithm is robust with respect to these hyper-parameters.

We next establish the convergence of the above GMF-V algorithm to an approximate Nash equilibrium of (GMFG), with complexity analysis.

Theorem 1.5.1 (Convergence and complexity of GMF-V). Assume the same assumptions as Theorem 1.3.1, and $f_{c, c^{\prime}} \subseteq \mathcal{F}_{c, c^{\prime}}$. Suppose that Alg has a value-based guarantee with parameters

$$
\left\{C_{\mathcal{M}}^{(i)}, \alpha_{1}^{(i)}, \alpha_{2}^{(i)}, \alpha_{3}^{(i)}, \alpha_{4}^{(i)}\right\}_{i=1}^{m}
$$

For any $\epsilon, \delta>0$, set $\delta_{k}=\delta / K_{\epsilon, \eta}, \epsilon_{k}=(k+1)^{-(1+\eta)}$ for some $\eta \in(0,1]\left(k=0, \ldots, K_{\epsilon, \eta}-1\right)$, and $c \geq c^{\prime} \geq \frac{\log (1 / \epsilon)}{\phi(\epsilon)}$. Then with probability at least $1-2 \delta$,

$$
W_{1}\left(\mathcal{L}_{K_{\epsilon, \eta}}, \mathcal{L}^{\star}\right) \leq C_{0} \epsilon
$$

Here $K_{\epsilon, \eta}:=\left\lceil 2 \max \left\{(\eta \epsilon / c)^{-1 / \eta}, \log _{d}(\epsilon / \max \{\operatorname{diam}(\mathcal{S}) \operatorname{diam}(\mathcal{A}), c\})+1\right\}\right\rceil$ is the number of outer iterations, and the constant $C_{0}$ is independent of $\delta, \epsilon$ and $\eta$.
Moreover, the total number of samples $T=\sum_{k=0}^{K_{\epsilon, \eta}-1} T_{\mathcal{M}_{\mathcal{L}_{k}}}\left(\delta_{k}, \epsilon_{k}\right)$ is bounded by

$$
\begin{equation*}
T \leq \sum_{i=1}^{m} \frac{2^{\alpha_{2}^{(i)}}}{2 \alpha_{1}^{(i)}+1} C_{\mathcal{M}}^{(i)} K_{\epsilon, \eta}^{2 \alpha_{1}^{(i)}+1}\left(K_{\epsilon, \eta} / \delta\right)^{\alpha_{3}^{(i)}}\left(\log \left(K_{\epsilon, \eta} / \delta\right)\right)^{\alpha_{2}^{(i)}+\alpha_{4}^{(i)}} \tag{1.5.2}
\end{equation*}
$$

The proof of Theorem 1.5.1 (in Section 1.7.4) depends on the Lipschitz continuity of the smoothing operator $f_{c, c^{\prime}}$, the closeness between $f_{c, c^{\prime}}$ and the argmax-e (Lemma 1.7.3 in Section 1.7.3), and the complexity of Alg provided by the value-based guarantee.

### 1.5.1.2 GMF-V-Q: GMF-V with Q-learning

As an example of the GMF-V algorithm, we describe algorithm GMF-V-Q, a Q-learning based GMF-V algorithm. For an $\operatorname{MDP} \mathcal{M}=(\mathcal{S}, \mathcal{A}, P, r, \gamma)$, the synchronous Q-learning
algorithm approximates the value iteration by stochastic approximation. At each step $l$, with state $s$ and action $a$, the system reaches state $s^{\prime}$ according to the controlled dynamics, and the Q-function approximation $Q_{l}$ is updated by

$$
\begin{equation*}
\hat{Q}^{l+1}(s, a)=\left(1-\beta_{l}\right) \hat{Q}^{l}(s, a)+\beta_{l}\left[r(s, a)+\gamma \max _{\bar{a}} \hat{Q}^{l}\left(s^{\prime}, \bar{a}\right)\right], \quad \forall s \in \mathcal{S}, a \in \mathcal{A} \tag{1.5.3}
\end{equation*}
$$

where $\hat{Q}^{0}(s, a)=C$ for some constant $C \in \mathbb{R}$ for any $s \in \mathcal{S}$ and $a \in \mathcal{A}$, and the step size $\beta_{l}$ can be chosen as ([53])

$$
\begin{equation*}
\beta_{l}=|l+1|^{-h} \tag{1.5.4}
\end{equation*}
$$

with $h \in(1 / 2,1)$.
The corresponding synchronous Q-learning based algorithm with the standard softmax operator is GMF-V-Q (Algorithm 3), and will be used in the experiment (Section 1.8).

```
Algorithm 3 Q-learning for GMFGs (GMF-V-Q)
    Input: Initial \(\mathcal{L}_{0}, \epsilon\)-net \(S_{\epsilon}\), tolerances \(\epsilon_{k}, \delta_{k}>0, k=0,1, \ldots\).
    for \(k=0,1, \cdots\) do
        Perform synchronous Q-learning with stepsizes (1.5.4) for \(T_{k}=T_{\mathcal{M}_{\mathcal{L}_{k}}}\left(\epsilon_{k}, \delta_{k}\right)\) itera-
        tions to find the approximate Q-function \(\hat{Q}_{k}^{\star}=\hat{Q}^{T_{k}}\) of the MDP \(\mathcal{M}_{\mathcal{L}_{k}}\).
        Compute \(\pi_{k} \in \Pi\) with \(\pi_{k}(s)=\operatorname{softmax}_{c}\left(\hat{Q}_{k}^{\star}(s, \cdot)\right)\).
        Sample \(s \sim \mu_{k}\left(\mu_{k}\right.\) is the population state marginal of \(\left.\mathcal{L}_{k}\right)\), obtain \(\tilde{\mathcal{L}}_{k+1}\) from
        \(\mathcal{G}\left(s, \pi_{k}, \mathcal{L}_{k}\right)\).
        Find \(\mathcal{L}_{k+1}=\operatorname{Proj}_{S_{\epsilon}}\left(\tilde{\mathcal{L}}_{k+1}\right)\)
    end for
```

Let us first recall the following sample complexity result for synchronous Q-learning method.

Lemma 1.5.2 ([53]: sample complexity of synchronous Q-learning). For an MDP, say $\mathcal{M}=(\mathcal{S}, \mathcal{A}, P, r, \gamma)$, suppose that the $Q$-learning algorithm takes step-sizes (1.5.4). Then $\left\|\hat{Q}^{T_{\mathcal{M}}(\delta, \epsilon)}-Q_{\mathcal{M}}^{\star}\right\|_{\infty} \leq \epsilon$ with probability at least $1-2 \delta$. Here $\hat{Q}^{T}$ is the $T$-th update in the $Q$-learning updates (1.5.3), $Q_{\mathcal{M}}^{\star}$ is the (optimal) $Q$-function, and

$$
T_{\mathcal{M}}(\epsilon, \delta)=\Omega\left(\left(\frac{V_{\max }^{2} \log \left(\frac{|\mathcal{S} \| \mathcal{A}| V_{\max }}{\delta \beta \epsilon}\right)}{\beta^{2} \epsilon^{2}}\right)^{\frac{1}{h}}+\left(\frac{1}{\beta} \log \frac{V_{\max }}{\epsilon}\right)^{\frac{1}{1-h}}\right)
$$

where $\beta=(1-\gamma) / 2, V_{\max }=R_{\max } /(1-\gamma)$, and $R_{\max }$ is such that a.s. $0 \leq r(s, a) \leq R_{\max }$.

This lemma implies immediately the value-based guarantee (as in Definition 1.5.1) and the convergence for GMF-V-Q. Similar results can be established for asynchronous Q-learning method, as shown in Section 1.10.2.

Corollary 1.5.3. The synchronous Q-learning algorithm with appropriate choices of stepsizes (cf. (1.5.4)) satisfies the value-based guarantee with parameters $\left\{\tilde{C}_{\mathcal{M}}^{(i)}, \alpha_{1}^{(i)}, \alpha_{2}^{(i)}, \alpha_{3}^{(i)}, \alpha_{4}^{(i)}\right\}_{i=1}^{3}$, where $C_{\mathcal{M}}^{(i)}(i=1,2,3)$ are constants depending on $|\mathcal{S}|,|\mathcal{A}|, V_{\max }, \beta$ and $h$, and

$$
\begin{aligned}
& \alpha_{1}^{(1)}=2 / h, \alpha_{2}^{(1)}=1 / h, \alpha_{3}^{(1)}=\alpha_{4}^{(1)}=0 \\
& \alpha_{1}^{(2)}=2 / h, \alpha_{2}^{(2)}=\alpha_{3}^{(2)}=0 \text { and } \alpha_{4}^{(2)}=1 / h ; \\
& \alpha_{1}^{(3)}=0, \alpha_{2}^{(3)}=1 /(1-h), \alpha_{3}^{(3)}=0 \text { and } \alpha_{4}^{(3)}=0 .
\end{aligned}
$$

In addition, assume the same assumptions as Theorem 1.3.1, then for Algorithm 3 with synchronous $Q$-learning method, with probability at least $1-2 \delta, W_{1}\left(\mathcal{L}_{K_{\epsilon, \eta}}, \mathcal{L}^{\star}\right) \leq C_{0} \epsilon$, where $K_{\epsilon, \eta}$ is defined as in Theorem 1.5.1. And the total number of samples $T=\sum_{k=0}^{K_{\epsilon, \eta}-1} T_{\mathcal{M}_{\mathcal{L}_{k}}}\left(\epsilon_{k}, \delta_{k}\right)$ is bounded by

$$
T \leq O\left(K_{\epsilon, \eta}^{\frac{4}{h}+1}\left(\log \frac{K_{\epsilon, \eta}}{\delta}\right)^{\frac{1}{h}}+\left(\log \frac{K_{\epsilon, \eta}}{\delta}\right)^{\frac{1}{1-h}}\right)
$$

### 1.5.2 Policy-based algorithms

In addition to algorithms with value-based guarantees ( $c f$. Definition 1.5.1), there are also numerous algorithms with policy-based guarantees.

Definition 1.5.2 (Policy-based Guarantee). For an arbitrary MDP $\mathcal{M}$, we say that an algorithm has a policy-based guarantee with parameters $\left\{C_{\mathcal{M}}^{(i)}, \alpha_{1}^{(i)}, \alpha_{2}^{(i)}, \alpha_{3}^{(i)}, \alpha_{4}^{(i)}\right\}_{i=1}^{m}$, if for any $\epsilon, \delta>0$, after obtaining

$$
\begin{equation*}
T_{\mathcal{M}}(\epsilon, \delta)=\sum_{i=1}^{m} C_{\mathcal{M}}^{(i)}\left(\frac{1}{\epsilon}\right)^{\alpha_{1}^{(i)}}\left(\log \frac{1}{\epsilon}\right)^{\alpha_{2}^{(i)}}\left(\frac{1}{\delta}\right)^{\alpha_{3}^{(i)}}\left(\log \frac{1}{\delta}\right)^{\alpha_{4}^{(i)}} \tag{1.5.5}
\end{equation*}
$$

samples from the simulator oracle $\mathcal{G}$, with probability at least $1-2 \delta$, it outputs an approximate policy $\pi_{T_{\mathcal{M}}(\epsilon, \delta)}$, which satisfies $V_{\mathcal{M}}^{\star}(s)-V_{\mathcal{M}}^{\pi_{T_{\mathcal{M}}(\epsilon, \delta)}}(s) \leq \epsilon, \forall s \in \mathcal{S}$.

### 1.5.2.1 GMF-P

Before we present policy-based RL algorithms, let us first establish a connection between policy-based and value-based guarantees.

To start, take any policy $\pi \in \Pi$, consider the following synchronous temporal difference (TD) iterations:

$$
\begin{equation*}
\tilde{Q}_{\pi}^{l+1}(s, a)=\left(1-\beta_{l}\right) \tilde{Q}_{\pi}^{l}(s, a)+\beta_{l}\left[r(s, a)+\gamma \tilde{Q}_{\pi}^{l}\left(s^{\prime}, a^{\prime}\right)\right], \quad \forall s \in \mathcal{S}, a \in \mathcal{A} \tag{1.5.6}
\end{equation*}
$$

where $a^{\prime} \sim \pi\left(s^{\prime}\right), \tilde{Q}_{\pi}^{0}(s, a)=C$ for some constant $C \in \mathbb{R}$ and any $s \in \mathcal{S}$ and $a \in \mathcal{A}$, and the step size $\beta_{l}=(l+1)^{-h}$ for some $h \in(1 / 2,1)$.

Then we have

Lemma 1.5.4. Suppose that the algorithm Alg satisfies a policy-based guarantee with parameters $\left\{C_{\mathcal{M}}^{(i)}, \alpha_{1}^{(i)}, \alpha_{2}^{(i)}, \alpha_{3}^{(i)}, \alpha_{4}^{(i)}\right\}_{i=1}^{m}$. Let $\tilde{Q}_{\pi}^{l}$ be defined by (1.5.6). Then for any $\delta \in(0,1)$ and $\epsilon>0$, with probability at least $1-2 \delta,\left\|\tilde{Q}_{\pi_{T_{\mathcal{M}}(\epsilon, \delta / 2)}^{l}}-Q_{\mathcal{M}}^{\star}\right\|_{\infty} \leq \epsilon$ if

$$
\begin{equation*}
l=\Omega\left(\left(\frac{V_{\max } \log \left(\frac{|\mathcal{S} \| \mathcal{A}| V_{\max }}{\delta \beta^{2} \epsilon}\right)}{\beta^{4} \epsilon^{2}}\right)^{1 / h}+\left(\frac{1}{\beta} \log \frac{V_{\max }}{\beta \epsilon}\right)^{1 /(1-h)}\right) \tag{1.5.7}
\end{equation*}
$$

where $V_{\max }=R_{\max } /(1-\gamma)$ and $\beta=(1-\gamma) / 2$.
Consequently, the algorithm Alg (combined with TD updates (1.5.6)) also has a valuebased guarantee with parameters $\left\{\tilde{C}_{\mathcal{M}}^{(i)}, \alpha_{1}^{(i)}, \alpha_{2}^{(i)}, \alpha_{3}^{(i)}, \alpha_{4}^{(i)}\right\}_{i=1}^{m+3}$, where $\tilde{C}_{\mathcal{M}}^{(i)}$ is some constant multiple of $C_{\mathcal{M}}^{(i)}(i=1, \ldots, m), \tilde{C}_{\mathcal{M}}^{(m+i)}(i=1,2,3)$ are constants depending on $V_{\max },|\mathcal{S}|$, $|\mathcal{A}|, \beta$ and $h$, and we have

$$
\begin{align*}
& \alpha_{1}^{(m+1)}=2 / h, \alpha_{2}^{(m+1)}=1 / h, \alpha_{3}^{(m+1)}=\alpha_{4}^{(m+1)}=0 \\
& \alpha_{1}^{(m+2)}=2 / h, \alpha_{2}^{(m+2)}=\alpha_{3}^{(m+2)}=0 \text { and } \alpha_{4}^{(m+2)}=1 / h  \tag{1.5.8}\\
& \alpha_{1}^{(m+3)}=0, \alpha_{2}^{(m+3)}=1 /(1-h), \alpha_{3}^{(m+3)}=0 \text { and } \alpha_{4}^{(m+3)}=0 .
\end{align*}
$$

The above lemma indicates that any algorithm with a policy-based guarantee also satisfies a value-based guarantee with similar parameters (when combined with the TD updates). The policy-based algorithm GMF-P (Algorithm 4) makes use of Lemma 1.5.4 to select the hyper-parameter $l$ so that the resulting $\tilde{Q}_{\pi_{T_{\mathcal{M}}(\epsilon, \delta / 2)}^{l}}$ forms a good value-based certificate.

```
Algorithm 4 GMF-P( \(\left.\operatorname{Alg}, \quad f_{c, c^{\prime}}\right)\)
    Input: Initial \(\mathcal{L}_{0}, \epsilon\)-net \(S_{\epsilon}\), temperatures \(c, c^{\prime}>0\), tolerances \(\epsilon_{k}, \delta_{k}>0, k=0,1, \ldots\).
    for \(k=0,1, \cdots\) do
        Apply Alg to find the approximate policy \(\hat{\pi}_{k}=\pi_{T_{k}}\) of the MDP \(\mathcal{M}_{k}:=\mathcal{M}_{\mathcal{L}_{k}}\), where
        \(T_{k}=T_{\mathcal{M}_{k}}\left(\epsilon_{k}, \delta_{k} / 2\right)\).
        Compute \(\tilde{Q}_{\tilde{\pi}_{k}}^{l_{k}}\) using TD updates (1.5.6) for MDP \(\mathcal{M}_{k}\), with \(l_{k}\) satisfying (1.5.7)
        (with \(\epsilon\) and \(\delta\) replaced by \(\epsilon_{k}\) and \(\delta_{k} / 2\), respectively).
        Compute \(\pi_{k}(s)=f_{c, c^{\prime}}\left(\tilde{Q}_{\tilde{\pi}_{k}}^{l_{k}}(s, \cdot)\right)\).
        Sample \(s \sim \mu_{k}\left(\mu_{k}\right.\) is the population state marginal of \(\left.\mathcal{L}_{k}\right)\), obtain \(\tilde{\mathcal{L}}_{k+1}\) from
        \(\mathcal{G}\left(s, \pi_{k}, \mathcal{L}_{k}\right)\).
        Find \(\mathcal{L}_{k+1}=\operatorname{Proj}_{S_{\epsilon}}\left(\tilde{\mathcal{L}}_{k+1}\right)\)
    end for
```

We next present the convergence property for the GMF-P algorithm by combining the proofs of Lemma 1.5.4 and Theorem 1.5.1.

Theorem 1.5.5 (Convergence and complexity of GMF-P). Assume the same assumptions as in Theorem 1.3.1, and in addition that $f_{c, c^{\prime}} \subseteq \mathcal{F}_{c, c^{\prime}}$. Suppose that Alg has a policy-based guarantee with parameters

$$
\left\{C_{\mathcal{M}}^{(i)}, \alpha_{1}^{(i)}, \alpha_{2}^{(i)}, \alpha_{3}^{(i)}, \alpha_{4}^{(i)}\right\}_{i=1}^{m}
$$

Then for any $\epsilon, \delta>0$, set $\delta_{k}=\delta / K_{\epsilon, \eta}, \epsilon_{k}=(k+1)^{-(1+\eta)}$ for some $\eta \in(0,1](k=$ $\left.0, \ldots, K_{\epsilon, \eta}-1\right)$, and $c \geq c^{\prime} \geq \frac{\log (1 / \epsilon)}{\phi(\epsilon)}$, with probability at least $1-2 \delta$,

$$
W_{1}\left(\mathcal{L}_{K_{\epsilon, \eta}}, \mathcal{L}^{\star}\right) \leq C_{0} \epsilon
$$

Here $K_{\epsilon, \eta}:=\left\lceil 2 \max \left\{(\eta \epsilon / c)^{-1 / \eta}, \log _{d}(\epsilon / \max \{\operatorname{diam}(\mathcal{S}) \operatorname{diam}(\mathcal{A}), c\})+1\right\}\right\rceil$ is the number of outer iterations, and the constant $C_{0}$ is independent of $\delta, \epsilon$ and $\eta$.
Moreover, the total number of samples $T=\sum_{k=0}^{K_{\epsilon, \eta}-1} T_{\mathcal{M}_{\mathcal{L}_{k}}}\left(\delta_{k}, \epsilon_{k}\right)$ is bounded by

$$
\begin{equation*}
T \leq \sum_{i=1}^{m+3} \frac{2^{\alpha_{2}^{(i)}}}{2 \alpha_{1}^{(i)}+1} \tilde{C}_{\mathcal{M}}^{(i)} K_{\epsilon, \eta}^{2 \alpha_{1}^{(i)}+1}\left(K_{\epsilon, \eta} / \delta\right)^{\alpha_{3}^{(i)}}\left(\log \left(K_{\epsilon, \eta} / \delta\right)\right)^{\alpha_{2}^{(i)}+\alpha_{4}^{(i)}} \tag{1.5.9}
\end{equation*}
$$

where the parameters $\left\{\tilde{C}_{\mathcal{M}}^{(i)}, \alpha_{1}^{(i)}, \alpha_{2}^{(i)}, \alpha_{3}^{(i)}, \alpha_{4}^{(i)}\right\}_{i=1}^{m+3}$ are defined in Lemma 1.5.4.

### 1.5.2.2 GMF-P-TRPO: GMF-P with TRPO

A special form of the GMF-P algorithm utilizes the trust region policy optimization (TRPO) algorithm [140, 141]. We call it GMF-P-TRPO.

Sample-based TRPO [141] assumes access to a $\nu$-restart model. That is, it can only access sampled trajectories and restarts according to the distribution $\nu$. Here we pick $\nu$ such that $C^{\pi^{\star}}:=\left\|\frac{d_{\mathrm{Unif}}^{\mathcal{S}}}{}\right\|_{\infty}=\max _{s \in \mathcal{S}}\left|\frac{d_{\mathrm{Unif}}^{\mathcal{S}}}{\pi_{\mathcal{S}}^{\star}(s)}\right|<\infty$, where $d_{\rho}^{\pi}=(1-\gamma) \rho\left(I-\gamma P^{\pi}\right)^{-1}$ and Unif $_{\mathcal{S}}$ is the uniform distribution on set $\mathcal{S}$. Sample-based TRPO samples $M_{0}$ trajectories per episode. The initial state $s_{0}$ at the beginning of each episode is sampled from $\nu$. In every trajectory $m\left(m=1,2, \cdots, M_{0}\right)$ of the $l$-th episode, it first samples $s_{m} \sim d_{\nu}^{\pi_{l}}$ and takes an action $a_{m} \sim \operatorname{Unif}_{\mathcal{A}}$ where $\operatorname{Unif}_{\mathcal{A}}$ is the uniform distribution on the set $\mathcal{A}$. Then, by following the current $\pi_{l}$, it estimates $Q^{\pi_{l}}\left(s_{m}, a_{m}\right)$ using a rollout. Denote this estimate as $\hat{Q}^{\pi_{l}}\left(s_{m}, a_{m}, m\right)$ and observe that it is (nearly) an unbiased estimator of $Q^{\pi_{l}}\left(s_{m}, a_{m}\right)$. We assume that each rollout runs sufficiently long so that the bias is sufficiently small. SampleBased TRPO updates the policy at the end of the $l$-th episode, by the following proximal problem

$$
\pi_{l+1} \in \arg \max _{\pi \in \Delta_{\mathcal{A}}^{|\mathcal{S}|}}\left\{\frac{1}{M_{0}} \sum_{m=1}^{M_{0}} \frac{1}{t_{l}(1-\gamma)} B_{w}\left(s_{m} ; \pi, \pi_{l}\right)+\left\langle\hat{\nabla} V^{\pi_{l}}[m], \pi\left(s_{m}\right)-\pi_{l}\left(s_{m}\right)\right\rangle\right\},
$$

where the estimation of the gradient is

$$
\hat{\nabla} V^{\pi_{l}}[m]:=\frac{1}{1-\gamma}|\mathcal{A}| \hat{Q}^{\pi_{l}}\left(s_{m}, \cdot, m\right) \circ \mathbf{I}_{\left\{\cdot=a_{m}\right\}}
$$

Given two policies $\pi_{1}$ and $\pi_{2}$, we denote their Bregman distance associated with a strongly convex function $w$ as $B_{w}\left(s ; \pi_{1}, \pi_{2}\right)=B_{w}\left(\pi_{1}(s), \pi_{2}(s)\right)$, where $B_{w}(x, y):=w(x)-w(y)-$ $\langle\nabla w(y), x-y\rangle$ and $\pi_{i}(s) \in P(\mathcal{A})(i=1,2)$. Denote $B_{w}\left(\pi_{1}, \pi_{2}\right) \in \mathbb{R}^{|\mathcal{S}|}$ as the corresponding state-wise vector. Here we consider two common cases for $w$ : when $w(x)=\frac{1}{2}\|x\|_{2}^{2}$ is the Euclidean distance, $B_{w}(x, y)=\frac{1}{2}\|x-y\|_{2}^{2}$; when $w(x)=H(x)$ is the negative entropy, $B_{w}(x, y)=d_{\mathrm{KL}}(x \| y)$. We refer to [141, Section 6.2] for more detailed discussion on Samplebased TRPO.

The above guarantee follows from the sample complexity result below by specifying $\mu:=$ Unif $_{\mathcal{S}}$. Notice that here for any $\mu \in \mathcal{P}(\mathcal{S})$, we define $V^{\star}(\mu):=\sum_{s \in \mathcal{S}} \mu(s) V^{\star}(s)$, and similarly $V^{\pi_{k}}(\mu):=\sum_{s \in \mathcal{S}} \mu(s) V^{\pi_{k}}(s)$.

The sample complexity of TRPO algorithm can be characterized as below.
Lemma 1.5.6 (Theorem 5 in [141]: sample complexity of TRPO). Let $\left\{\pi_{l}\right\}_{l \geq 0}$ be the sequence generated by Sample-Based TRPO, using

$$
M_{0} \geq \Omega\left(\frac{|\mathcal{A}|^{2} C^{2}(|\mathcal{S}| \log |\mathcal{A}|+\log 1 / \delta)}{(1-\gamma)^{2} \epsilon^{2}}\right)
$$

samples in each episode, with $t_{l}=\frac{(1-\gamma)}{C_{\omega, 1} C \sqrt{l+1}}$. Let $\left\{V_{\text {best }}^{N}\right\}_{N \geq 0}$ be the sequence of best achieved values, $V_{\text {best }}^{N}(\mu):=\max _{l=0,1, \cdots, N} V^{\pi_{l}}(\mu)$, where $\mu \in \mathcal{P}(\mathcal{S})$. Then with probability greater than $1-\delta$ for every $\epsilon>0$, the following holds for all $N \geq 1$ :

$$
V^{\star}(\mu)-V_{b e s t}^{N}(\mu) \leq O\left(\frac{C_{\omega, 1} C}{(1-\gamma)^{2} \sqrt{N}}+\frac{C^{\pi^{\star}} \epsilon}{(1-\gamma)^{2}}\right)
$$

Here $C>0$ is the upper bound on the reward function $r, C_{w, 1}=\sqrt{|\mathcal{A}|}$ in the euclidean case and $C_{w, 1}=1$ in the non-euclidean case, $C_{w, 2}=1$ for the euclidean case and $C_{w, 2}=|\mathcal{A}|^{2}$ for the non-euclidean case. Note that unlike the case of Q-learning, here we are only guaranteed to have certain iterate among iterations $0, \ldots, N$ that satisfy the desired suboptimality bound. Note that this is a common pattern of the theoretical results for policy optimization algorithms in the RL literature [4, 161], unless the (oracle) access to exact policy gradients is assumed [114]. For simplicity, hereafter we assume an oracle access to such an iterate after running TRPO. In practice, with additional (polynomial number of) samples, one can explicitly identify a single policy satisfying the desired bound with high probability; see e.g., the two-phase technique in [65].

Note that [141, Theorem 5] has both regularized version and unregularized version of TRPO. Here we only adopt the unregularized version which fits the framework of Algorithm 4. For more materials on regularized MDPs and reinforcement learning, we refer the readers to [122, 63, 48].

Based on the sample complexity in Lemma 1.5.6, the following policy-based guarantee for TRPO algorithm and the convergence result for GMF-P-TRPO can be obtained.

Corollary 1.5.7. Let $t_{l}=\frac{(1-\gamma)}{C_{\omega, 1} C \sqrt{l+1}}$, then TRPO algorithm satisfies the policy-based guarantee with parameters $\left\{\tilde{C}_{\mathcal{M}}^{(i)}, \alpha_{1}^{(i)}, \alpha_{2}^{(i)}, \alpha_{3}^{(i)}, \alpha_{4}^{(i)}\right\}_{i=1}^{2}$, where $C_{\mathcal{M}}^{(i)}(i=1,2)$ are constants depending on $|\mathcal{S}|,|\mathcal{A}|, V_{\max }, \beta$ and $h$, and we have:

$$
\begin{gathered}
\alpha_{1}^{(1)}=5 / 2, \alpha_{j}^{(1)}=0 \text { for } j=2,3,4, \\
\alpha_{1}^{(2)}=5 / 2, \alpha_{4}^{(2)}=1, \alpha_{2}^{(2)}=\alpha_{3}^{(2)}=0
\end{gathered}
$$

In addition, under same assumptions as Theorem 1.3.1, then for Algorithm 4 using TRPO method, with probability at least $1-2 \delta$, $W_{1}\left(\mathcal{L}_{K_{\epsilon, \eta}}, \mathcal{L}^{*}\right) \leq C_{0} \epsilon$, where $K_{\epsilon, \eta}$ is defined as in Theorem 1.5.1. And the total number of samples $T=\sum_{k=0}^{K_{\epsilon, \eta}-1} T_{\mathcal{M}_{\mathcal{L}_{k}}}\left(\delta_{k}, \epsilon_{k}\right)$ is bounded by

$$
T \leq O\left(K_{\epsilon, \eta}^{6}\left(\log \frac{K_{\epsilon, \eta}}{\delta}\right)+K_{\epsilon, \eta}^{\frac{4}{h}+1}\left(\log \frac{K_{\epsilon, \eta}}{\delta}\right)^{\frac{1}{h}}+\left(\log \frac{K_{\epsilon, \eta}}{\delta}\right)^{\frac{1}{1-h}}\right)
$$

### 1.6 Applications to $N$-player Games

In this section, we discuss a potential application of our modeling and approach to N player settings. To this end, we consider extensions of Algorithms 2 and 4 with weaker assumptions on the simulator access. In particular, we weaken the simulator oracle assumption in Section 1.4 as follows.

Weak simulator oracle. For each player $i$, given any policy $\pi \in \Pi$, the current state $s_{i} \in \mathcal{S}$, for any empirical population state-action distribution $\mathcal{L}_{N}$, one can obtain a sample of the next state $s_{i}^{\prime} \sim P_{\mathcal{L}_{N}}\left(\cdot \mid s_{i}, \pi\left(s_{i}\right)\right)=P\left(\cdot \mid s_{i}, \pi\left(s_{i}\right), \mathcal{L}_{N}\right)$ and a reward $r=r_{\mathcal{L}_{N}}\left(s_{i}, \pi\left(s_{i}\right)\right)=$ $r\left(s_{i}, \pi\left(s_{i}\right), \mathcal{L}_{N}\right)$. For brevity, we denote the simulator as $\left(s_{i}^{\prime}, r\right)=\mathcal{G}_{W}\left(s_{i}, \pi, \mathcal{L}_{N}\right)$.

We say that $\mathcal{L}_{N}$ is an empirical population state-action distribution of $N$-players if for each $s \in \mathcal{S}, a \in \mathcal{A}, \mathcal{L}_{N}(s, a)=\frac{1}{N} \sum_{i=1}^{N} \mathbf{I}_{s_{i}=s, a_{i}=a}$ for some state-action profile of $\left\{s_{i}, a_{i}\right\}_{i=1}^{N}$. Equivalently, this holds if $N \mathcal{L}_{N}(s, a)$ is a non-negative integer for each $s \in \mathcal{S}, a \in \mathcal{A}$, and $\sum_{s, a} \mathcal{L}_{N}(s, a)=1$. We denote the set of empirical population state-action distributions as $\operatorname{Emp}_{N}$.

RL algorithms with access only to $\mathcal{G}_{W}$. Compared to the original simulator oracle $\mathcal{G}$, the weak simulator $\mathcal{G}_{W}$ only accepts empirical population state-action distributions as inputs, and does not directly output the next (empirical) population state-action distribution.

To make use of the simulator $\mathcal{G}_{W}$, we modify Algorithm 2 and Algorithm 4 to algorithms (Algorithms 5 and 6). In particular, see Step 6 in Algorithm 5 and Step 7 in Algorithm 6 for generating empirical distributions from simulator $\mathcal{G}_{W}$.

One can observe that $\mathbf{E m p}_{N}$ already serves as an $1 / N$-net. So one can directly use it without additional projections. The definition of $\mathcal{L}_{k}$ also makes sure that $\mathcal{L}_{k} \in \mathbf{E m p}_{N}$ as required for the input of the weaker simulator.

```
Algorithm 5 GMF-VW( \(\left.\boldsymbol{A l g}, f_{c, c^{\prime}}\right)\) : weak simulator
    Input: Initial \(\mathcal{L}_{0}\), temperatures \(c, c^{\prime}>0\), tolerances \(\epsilon_{k}, \delta_{k}>0, k=0,1, \ldots\)
    for \(k=0,1, \cdots\) do
        Apply \(A l g\) to find the approximate Q-function \(\hat{Q}_{k}^{\star}=\hat{Q}^{T_{k}}\) of the MDP \(\mathcal{M}_{\mathcal{L}_{k}}\), where
        \(T_{k}=T_{\mathcal{M}_{\mathcal{L}_{k}}}\left(\epsilon_{k}, \delta_{k}\right)\).
        Compute \(\pi_{k}(s)=f_{c, c^{\prime}}\left(\hat{Q}_{k}^{\star}(s, \cdot)\right)\).
        for \(i=1,2, \cdots, N\) do
    6: \(\quad\) Sample \(s_{i} \stackrel{\text { i.i.d. }}{\sim} \mu_{k}\), then obtain \(s_{i}^{\prime}\) i.i.d. from \(\mathcal{G}_{W}\left(s_{i}, \pi_{k}, \mathcal{L}_{k}\right)\) and \(a_{i}^{\prime} \stackrel{\text { i.i.d. }}{\sim} \pi_{k}\left(s_{i}^{\prime}\right)\).
        end for
        Compute \(\mathcal{L}_{k+1}\) with \(\mathcal{L}_{k+1}(s, a)=\frac{1}{N} \sum_{i=1}^{N} \mathbf{I}_{s_{i}^{\prime}=s, a_{i}^{\prime}=a}\).
    end for
```

```
Algorithm 6 GMF-PW ( \(\boldsymbol{A l g}, \quad f_{c, c^{\prime}}\) ): weak simulator
    Input: Initial \(\mathcal{L}_{0}\), temperatures \(c, c^{\prime}>0\), tolerances \(\epsilon_{k}, \delta_{k}>0, k=0,1, \ldots\)
    for \(k=0,1, \cdots\) do
        Apply \(\operatorname{Alg}\) to find the approximate policy \(\hat{\pi}_{k}=\pi_{T_{k}}\) of the MDP \(\mathcal{M}_{k}:=\mathcal{M}_{\mathcal{L}_{k}}\), where
        \(T_{k}=T_{\mathcal{M}_{\mathcal{L}_{k}}}\left(\epsilon_{k}, \delta_{k} / 2\right)\).
        Compute \(\tilde{Q}_{\tilde{\pi}_{k}}^{l_{k}}\) using TD updates (1.5.6) for \(\operatorname{MDP} \mathcal{M}_{k}\), with \(l_{k}\) satisfying (1.5.7)
        (with \(\epsilon\) and \(\delta\) replaced by \(\epsilon_{k}\) and \(\delta_{k} / 2\), respectively).
        Compute \(\pi_{k}(s)=f_{c, c^{\prime}}\left(\hat{Q}_{\mathcal{M}_{k}, M_{k}, l_{k}}^{\hat{\pi}_{k}}(s, \cdot)\right)\).
        for \(i=1,2, \cdots, N\) do
    7: \(\quad\) Sample \(s_{i} \stackrel{\text { i.i.d. }}{\sim} \mu_{k}\), then obtain \(s_{i}^{\prime}\) i.i.d. from \(\mathcal{G}_{W}\left(s_{i}, \pi_{k}, \mathcal{L}_{k}\right)\) and \(a_{i}^{\prime} \stackrel{\text { i.i.d. }}{\sim} \pi_{k}\left(s_{i}^{\prime}\right)\).
        end for
        Compute \(\mathcal{L}_{k+1}\) with \(\mathcal{L}_{k+1}(s, a)=\frac{1}{N} \sum_{i=1}^{N} \mathbf{I}_{s_{i}^{\prime}=s, a_{i}^{\prime}=a}\).
    end for
```

Convergence results similar to Theorems 1.5.1 and 1.5.5 can be obtained for Algorithms 5 and 6 , respectively. (See Section 1.10.3.) Here the major difference is an additional $O(1 / \sqrt{N})$ term in the finite step error bound. It is worth mentioning that $O(1 / \sqrt{N})$ is consistent with the literature on MFG approximation errors of finite $N$-player games [87].

### 1.7 Proof of the main results

### 1.7.1 Proof of Lemma 1.3.2

In this section, we provide the proof of Lemma 1.3.2.

Proof. [Proof of Lemma 1.3.2] We begin by noticing that $\mathcal{L}^{\prime}=\Gamma_{2}(\pi, \mathcal{L})$ can be expanded and computed as follows:

$$
\begin{equation*}
\mu^{\prime}\left(s^{\prime}\right)=\sum_{s \in \mathcal{S}, a \in \mathcal{A}} \mu(s) P\left(s^{\prime} \mid s, a, \mathcal{L}\right) \pi(a \mid s), \quad \mathcal{L}^{\prime}\left(s^{\prime}, a^{\prime}\right)=\mu^{\prime}\left(s^{\prime}\right) \pi\left(a^{\prime} \mid s^{\prime}\right) \tag{1.7.1}
\end{equation*}
$$

where $\mu$ is the state marginal distribution of $\mathcal{L}$.
Now by the inequalities (1.10.3), we have

$$
\begin{align*}
& W_{1}\left(\Gamma_{2}\left(\pi_{1}, \mathcal{L}\right), \Gamma_{2}\left(\pi_{2}, \mathcal{L}\right)\right) \leq \operatorname{diam}(\mathcal{S} \times \mathcal{A}) d_{T V}\left(\Gamma_{2}\left(\pi_{1}, \mathcal{L}\right), \Gamma_{2}\left(\pi_{2}, \mathcal{L}\right)\right) \\
& =\frac{\operatorname{diam}(\mathcal{S} \times \mathcal{A})}{2} \sum_{s^{\prime} \in \mathcal{S}, a^{\prime} \in \mathcal{A}}\left|\sum_{s \in \mathcal{S}, a \in \mathcal{A}} \mu(s) P\left(s^{\prime} \mid s, a, \mathcal{L}\right)\left(\pi_{1}(a \mid s) \pi_{1}\left(a^{\prime} \mid s^{\prime}\right)-\pi_{2}(a \mid s) \pi_{2}\left(a^{\prime} \mid s^{\prime}\right)\right)\right| \\
& \leq \frac{\operatorname{diam}(\mathcal{S} \times \mathcal{A})}{2} \max _{s, a, \mathcal{L}, s^{\prime}} P\left(s^{\prime} \mid s, a, \mathcal{L}\right) \sum_{s, a, s^{\prime}, a^{\prime}} \mu(s)\left(\pi_{1}(a \mid s)+\pi_{2}(a \mid s)\right)\left|\pi_{1}\left(a^{\prime} \mid s^{\prime}\right)-\pi_{2}\left(a^{\prime} \mid s^{\prime}\right)\right| \\
& \leq \frac{\operatorname{diam}(\mathcal{S} \times \mathcal{A})}{2} \max _{s, a, \mathcal{L}, s^{\prime}} P\left(s^{\prime} \mid s, a, \mathcal{L}\right) \sum_{s^{\prime}, a^{\prime}}\left|\pi_{1}\left(a^{\prime} \mid s^{\prime}\right)-\pi_{2}\left(a^{\prime} \mid s^{\prime}\right)\right| \cdot(1+1) \\
& =2 \operatorname{diam}(\mathcal{S} \times \mathcal{A}) \max _{s, a, \mathcal{L}, s^{\prime}} P\left(s^{\prime} \mid s, a, \mathcal{L}\right) \sum_{s^{\prime}} d_{T V}\left(\pi_{1}\left(s^{\prime}\right), \pi_{2}\left(s^{\prime}\right)\right) \\
& \leq \frac{2 \operatorname{diam}(\mathcal{S} \times \mathcal{A}) \max _{s, a, \mathcal{L}, s^{\prime}} P\left(s^{\prime} \mid s, a, \mathcal{L}\right)|\mathcal{S}|}{d_{\min }(\mathcal{A})} D\left(\pi_{1}, \pi_{2}\right)=\frac{2 \operatorname{diam}(\mathcal{S}) \operatorname{diam}(\mathcal{A})|\mathcal{S}| c_{1}}{d_{\min }(\mathcal{A})} D\left(\pi_{1}, \pi_{2}\right) \tag{1.7.2}
\end{align*}
$$

Similarly, we have

$$
\begin{align*}
& W_{1}\left(\Gamma_{2}\left(\pi, \mathcal{L}_{1}\right), \Gamma_{2}\left(\pi, \mathcal{L}_{2}\right)\right) \leq \operatorname{diam}(\mathcal{S} \times \mathcal{A}) d_{T V}\left(\Gamma_{2}\left(\pi, \mathcal{L}_{1}\right), \Gamma_{2}\left(\pi, \mathcal{L}_{2}\right)\right) \\
& =\frac{\operatorname{diam}(\mathcal{S} \times \mathcal{A})}{2} \sum_{s^{\prime} \in \mathcal{S}, a^{\prime} \in \mathcal{A}}\left|\sum_{s \in \mathcal{S}, a \in \mathcal{A}} \mu(s) \pi(a \mid s) \pi\left(a^{\prime} \mid s^{\prime}\right)\left(P\left(s^{\prime} \mid s, a, \mathcal{L}_{1}\right)-P\left(s^{\prime} \mid s, a, \mathcal{L}_{2}\right)\right)\right| \\
& \leq \frac{\operatorname{diam}(\mathcal{S} \times \mathcal{A})}{2} \sum_{s, a, s^{\prime}, a^{\prime}} \mu(s) \pi(a \mid s) \pi\left(a^{\prime} \mid s^{\prime}\right)\left|P\left(s^{\prime} \mid s, a, \mathcal{L}_{1}\right)-P\left(s^{\prime} \mid s, a, \mathcal{L}_{2}\right)\right|  \tag{1.7.3}\\
& \leq \frac{\operatorname{diam}(\mathcal{S}) \operatorname{diam}(\mathcal{A}) c_{2}}{2} .
\end{align*}
$$

This completes the proof.

### 1.7.2 Proof of Lemma 1.5.4

For notation simplicity, in the following analysis we fix the MDP and omit the notation $\mathcal{M}$.

We begin by establishing the convergence rate of the synchronous TD updates (1.5.6).

Lemma 1.7.1. Take $\tilde{Q}_{\pi}^{l}$ from (1.5.6). Then for any $\delta \in(0,1)$ and $\epsilon>0$, with probability at least $1-\delta,\left\|\tilde{Q}_{\pi_{T(\epsilon, \delta)}^{l}}^{l}-Q^{\pi_{T(\epsilon, \delta)}}\right\|_{\infty} \leq \epsilon$ if

$$
\begin{equation*}
l=\Omega\left(\left(\frac{V_{\max } \log \left(\frac{|\mathcal{S} \| \mathcal{A}| V_{\max }}{\delta \beta \epsilon}\right)}{\beta^{2} \epsilon^{2}}\right)^{1 / h}+\left(\frac{1}{\beta} \log \frac{V_{\max }}{\epsilon}\right)^{1 /(1-h)}\right) \tag{1.7.4}
\end{equation*}
$$

where $V_{\max }=R_{\max } /(1-\gamma)$ and $\beta=(1-\gamma) / 2$.
The proof is adapted from that of [53, Theorem 2], with the max term in the Bellman operator modified to actions sampled from the current policy $\pi$. The details are omitted.

Proof. [Proof of Lemma 1.5.4] First, if $V^{\star}\left(s^{\prime}\right)-V^{\pi_{T(\epsilon, \delta / 2)}}\left(s^{\prime}\right) \leq \epsilon$, then

$$
\begin{align*}
\left|Q^{\pi_{T(\epsilon, \delta / 2)}}(s, a)-Q^{\star}(s, a)\right| & =\gamma\left|\sum_{s^{\prime} \in \mathcal{S}} P\left(s^{\prime} \mid s, a\right) V^{\pi_{T(\epsilon, \delta / 2)}}\left(s^{\prime}\right)-\sum_{s^{\prime} \in \mathcal{S}} P\left(s^{\prime} \mid s, a\right) V^{\star}\left(s^{\prime}\right)\right|  \tag{1.7.5}\\
& \leq \gamma \sum_{s^{\prime} \in \mathcal{S}} P\left(s^{\prime} \mid s, a\right)\left|V^{\pi_{T(\epsilon, \delta / 2)}}\left(s^{\prime}\right)-V^{\star}\left(s^{\prime}\right)\right| \leq \gamma \epsilon<\epsilon
\end{align*}
$$

for any $s \in \mathcal{S}, a \in \mathcal{A}$. Since $A l g$ is assumed to satisfying the policy-based guarantee, (1.7.5) holds with probability at least $1-\delta$.

In addition, by Lemma 1.7.1, whenever $l$ satisfies (1.5.7), with probability at least $1-$ $\delta / 2 \geq 1-\delta$,

$$
\begin{equation*}
\left\|\tilde{Q}_{\pi_{T(\epsilon, \delta / 2)}}^{l}-Q^{\pi_{T(\epsilon, \delta / 2)}}\right\|_{\infty} \leq(1-\gamma) \epsilon \tag{1.7.6}
\end{equation*}
$$

Combining (1.7.5) and (1.7.6), then for any $l$ satisfying (1.5.7), with probability at least $1-2 \delta$, we have

$$
\left\|\tilde{Q}_{\pi_{T(\epsilon, \delta / 2)}^{l}}^{l}-Q^{\star}\right\|_{\infty} \leq \gamma \epsilon+(1-\gamma) \epsilon=\epsilon
$$

The above result shows that for any $\delta \in(0,1)$ and $\epsilon>0$, after obtaining $T(\epsilon, \delta / 2)+|\mathcal{S} \| \mathcal{A}| l$ samples (with $l$ satisfying the lower bound (1.5.7)) from the simulator, with probability at least $1-2 \delta$, it outputs an approximate $Q$-function $\tilde{Q}_{\pi_{T(\epsilon, \delta / 2)}}$ which satisfies $\| \tilde{Q}_{\pi_{T(\epsilon, \delta / 2)}^{l}}-$ $Q^{\star} \|_{\infty} \leq \epsilon$. Thus Alg also has a value-based guarantee with parameters

$$
\begin{equation*}
\left\{\tilde{C}_{\mathcal{M}}^{(i)}, \alpha_{1}^{(i)}, \alpha_{2}^{(i)}, \alpha_{3}^{(i)}, \alpha_{4}^{(i)}\right\}_{i=1}^{m+3} \tag{1.7.7}
\end{equation*}
$$

specified in (1.5.8). Here the first $m$ groups of parameters come from $T(\epsilon, \delta / 2)$ while the last three groups of parameters come from $|\mathcal{S} \| \mathcal{A}| l$ (with the lower bound (1.5.7) of $l$ plugged in here).

### 1.7.3 Proof of $\mathcal{B}_{c, c^{\prime}} \subseteq \mathcal{F}_{c, c^{\prime}}$

Lemma 1.7.2. Suppose that $h: \mathbb{R} \rightarrow \mathbb{R}$ satisfies $h(a)-h(b) \leq c(a-b)$ for any $a \geq b \in \mathbb{R}$. Then the softmax function $\operatorname{softmax}_{h}$ is $c$-Lipschitz, i.e., $\left\|\operatorname{softmax}_{h}(x)-\operatorname{softmax}_{h}(y)\right\|_{2} \leq$ $c\|x-y\|_{2}$ for any $x, y \in \mathbb{R}^{n}$.

Proof. [Proof of Lemma 1.7.2] Notice that $\operatorname{softmax}_{h}(x)=\operatorname{softmax}(\tilde{h}(x))$, where

$$
\operatorname{softmax}(x)_{i}=\frac{\exp \left(x_{i}\right)}{\sum_{j=1}^{n} \exp \left(x_{j}\right)}(i=1, \ldots, n)
$$

is the standard softmax function and $\tilde{h}(x)_{i}=h\left(x_{i}\right)$ for $i=1, \ldots, n$. Now since softmax is 1-Lipschitz continuous ( $c f$. [62, Proposition 4]), and $\tilde{h}$ is $c$-Lipschitz continuous, we conclude that the composition softmax $\circ \tilde{h}$ is $c$-Lipschitz continuous.

Notice that for a finite set $\mathcal{X} \subseteq \mathbb{R}^{k}$ and any two (discrete) distributions $\nu, \nu^{\prime}$ over $\mathcal{X}$, we have

$$
\begin{equation*}
W_{1}\left(\nu, \nu^{\prime}\right) \leq \operatorname{diam}(\mathcal{X}) d_{T V}\left(\nu, \nu^{\prime}\right)=\frac{\operatorname{diam}(\mathcal{X})}{2}\left\|\nu-\nu^{\prime}\right\|_{1} \leq \frac{\operatorname{diam}(\mathcal{X}) \sqrt{|\mathcal{X}|}}{2}\left\|\nu-\nu^{\prime}\right\|_{2} \tag{1.7.8}
\end{equation*}
$$

where in computing the $\ell_{1}$-norm, $\nu, \nu^{\prime}$ are viewed as vectors of length $|\mathcal{X}|$.
Lemma 1.7.2 implies that for any $x, y \in \mathbb{R}^{|\mathcal{X}|}$, when $\operatorname{softmax}_{c}(x)$ and $\operatorname{softmax}_{c}(y)$ are viewed as probability distributions over $\mathcal{X}$, we have

$$
W_{1}\left(\operatorname{softmax}_{c}(x), \operatorname{softmax}_{c}(y)\right) \leq \frac{\operatorname{diam}(\mathcal{X}) \sqrt{|\mathcal{X}|} c}{2}\|x-y\|_{2} \leq \frac{\operatorname{diam}(\mathcal{X})|\mathcal{X}| c}{2}\|x-y\|_{\infty}
$$

Lemma 1.7.3. Suppose that $h: \mathbb{R} \rightarrow \mathbb{R}$ satisfies $c^{\prime}(a-b) \leq h(a)-h(b)$ for any $a \leq b \in \mathbb{R}$. Then for any $x \in \mathbb{R}^{n}$, the distance between the $\operatorname{softmax}_{h}$ and the argmax-e mapping is bounded by

$$
\left\|\operatorname{softmax}_{h}(x)-\operatorname{argmax}-\mathbf{e}(x)\right\|_{2} \leq 2 n \exp \left(-c^{\prime} \delta\right),
$$

where $\delta=x_{\max }-\max _{x_{j}<x_{\max }} x_{j}, x_{\max }=\max _{i=1, \ldots, n} x_{i}$, and $\delta:=\infty$ when all $x_{j}$ are equal.
Similar to Lemma 1.7.2, Lemma 1.7.3 implies that for any $x \in \mathbb{R}^{|\mathcal{X}|}$, viewing $\operatorname{softmax}_{h}(x)$ as probability distributions over $\mathcal{X}$ leads to

$$
W_{1}\left(\operatorname{softmax}_{h}(x), \operatorname{argmax}-\mathbf{e}(x)\right) \leq \operatorname{diam}(\mathcal{X})|\mathcal{X}| \exp (-c \delta)
$$

Proof. [Proof of Lemma 1.7.3] Without loss of generality, assume that $x_{1}=x_{2}=\cdots=x_{m}=$ $\max _{i=1, \ldots, n} x_{i}=x^{\star}>x_{j}$ for all $m<j \leq n$. Then

$$
\begin{gathered}
\operatorname{argmax}-\mathbf{e}(x)_{i}= \begin{cases}\frac{1}{m}, & i \leq m, \\
0, & \text { otherwise } .\end{cases} \\
\operatorname{softmax}_{h}(x)_{i}= \begin{cases}\frac{e^{h\left(x^{\star}\right)}}{m e^{h\left(x^{\star}\right)}+\sum_{j=m+1}^{n} e^{h\left(x_{j}\right)}}, & i \leq m, \\
\frac{e^{h\left(x_{i}\right)}}{m e^{h\left(x^{\star}\right)}+\sum_{j=m+1}^{n} e^{h\left(x_{j}\right)}}, & \text { otherwise. }\end{cases}
\end{gathered}
$$

Therefore

$$
\begin{aligned}
& \left\|\operatorname{softmax}_{h}(x)-\operatorname{argmax}-\mathbf{e}(x)\right\|_{2} \leq\left\|\operatorname{softmax}_{h}(x)-\operatorname{argmax}-\mathbf{e}(x)\right\|_{1} \\
& \quad=m\left(\frac{1}{m}-\frac{e^{h\left(x^{\star}\right)}}{m e^{h\left(x^{\star}\right)}+\sum_{j=m+1}^{n} e^{h\left(x_{j}\right)}}\right)+\frac{\sum_{i=m+1}^{n} e^{h\left(x_{i}\right)}}{m e^{h\left(x^{\star}\right)}+\sum_{j=m+1}^{n} e^{h\left(x_{j}\right)}} \\
& \quad=\frac{2 \sum_{i=m+1}^{n} e^{h\left(x_{i}\right)}}{m e^{h\left(x^{\star}\right)}+\sum_{i=m+1}^{n} e^{h\left(x_{i}\right)}}=\frac{2 \sum_{i=m+1}^{n} e^{-c^{\prime} \delta_{i}}}{m+\sum_{i=m+1}^{n} e^{-c \delta_{i}}} \\
& \quad \leq \frac{2}{m} \sum_{i=m+1}^{n} e^{-c^{\prime} \delta_{i}} \leq \frac{2(n-m)}{m} e^{-c^{\prime} \delta} \leq 2 n e^{-c^{\prime} \delta},
\end{aligned}
$$

with $\delta_{i}=x^{\star}-x^{i}$.
We are now ready to present the proofs of Theorems 1.5.1 and 1.5.5.

### 1.7.4 Proof of Theorems 1.5.1 and 1.5.5

Proof. [Proof of Theorem 1.5.1] Here we prove the case when we are using GMF-V and Alg has a value-based guarantee. Define $\hat{\Gamma}_{1}^{k}\left(\mathcal{L}_{k}\right):=f_{c, c^{\prime}}\left(\hat{Q}_{k}^{\star}\right)$. In the following, $\pi=f_{c, c^{\prime}}\left(Q_{\mathcal{L}}\right)$ is understood as the policy $\pi$ with $\pi(s)=f_{c, c^{\prime}}\left(Q_{\mathcal{L}}(s, \cdot)\right)$. Let $\mathcal{L}^{\star}$ be the population state-action pair in a stationary NE of (GMFG). Then $\pi_{k}=\hat{\Gamma}_{1}^{k}\left(\mathcal{L}_{k}\right)$. Denoting $d:=d_{1} d_{2}+d_{3}$, we see

$$
\begin{aligned}
& W_{1}\left(\tilde{\mathcal{L}}_{k+1}, \mathcal{L}^{\star}\right)=W_{1}\left(\Gamma_{2}\left(\pi_{k}, \mathcal{L}_{k}\right), \Gamma_{2}\left(\Gamma_{1}\left(\mathcal{L}^{\star}\right), \mathcal{L}^{\star}\right)\right) \\
& \quad \leq W_{1}\left(\Gamma_{2}\left(\Gamma_{1}\left(\mathcal{L}_{k}\right), \mathcal{L}_{k}\right), \Gamma_{2}\left(\Gamma_{1}\left(\mathcal{L}^{\star}\right), \mathcal{L}^{\star}\right)\right)+W_{1}\left(\Gamma_{2}\left(\Gamma_{1}\left(\mathcal{L}_{k}\right), \mathcal{L}_{k}\right), \Gamma_{2}\left(\hat{\Gamma}_{1}^{k}\left(\mathcal{L}_{k}\right), \mathcal{L}_{k}\right)\right) \\
& \leq W_{1}\left(\Gamma\left(\mathcal{L}_{k}\right), \Gamma\left(\mathcal{L}^{\star}\right)\right)+d_{2} D\left(\Gamma_{1}\left(\mathcal{L}_{k}\right), \hat{\Gamma}_{1}^{k}\left(\mathcal{L}_{k}\right)\right) \\
& \leq\left(d_{1} d_{2}+d_{3}\right) W_{1}\left(\mathcal{L}_{k}, \mathcal{L}^{\star}\right)+d_{2} D\left(\operatorname{argmax}-\mathbf{e}\left(Q_{\mathcal{L}_{k}}^{\star}\right), f_{c, c^{\prime}}\left(\hat{Q}_{k}^{\star}\right)\right) \\
& \leq d W_{1}\left(\mathcal{L}_{k}, \mathcal{L}^{\star}\right)+d_{2} D\left(f_{c, c^{\prime}}\left(\hat{Q}_{k}^{\star}\right), f_{c, c^{\prime}}\left(Q_{\mathcal{L}_{k}}^{\star}\right)\right) \\
& \quad+d_{2} D\left(\operatorname{argmax}-\mathbf{e}\left(Q_{\mathcal{L}_{k}}^{\star}\right), f_{c, c^{\prime}}\left(Q_{\mathcal{L}_{k}}^{\star}\right)\right) \\
& \leq d W_{1}\left(\mathcal{L}_{k}, \mathcal{L}^{\star}\right)+\frac{c d_{2} \operatorname{diam}(\mathcal{A})|\mathcal{A}|}{2}\left\|\hat{Q}_{k}^{\star}-Q_{\mathcal{L}_{k}}^{\star}\right\|_{\infty}+d_{2} D\left(\operatorname{argmax}-\mathbf{e}\left(Q_{\mathcal{L}_{k}}^{\star}\right), f_{c, c^{\prime}}\left(Q_{\mathcal{L}_{k}}^{\star}\right)\right) .
\end{aligned}
$$

Since $\mathcal{L}_{k} \in S_{\epsilon}$ by the projection step, by Lemma 1.7.3 and the algorithm $A l g$ has a policybased guarantee, with the choice of $T_{k}=T_{\mathcal{M}_{\mathcal{L}_{k}}}\left(\delta_{k}, \epsilon_{k}\right)$ ), we have, with probability at least $1-2 \delta_{k}$,

$$
\begin{equation*}
W_{1}\left(\tilde{\mathcal{L}}_{k+1}, \mathcal{L}^{\star}\right) \leq d W_{1}\left(\mathcal{L}_{k}, \mathcal{L}^{\star}\right)+\frac{c d_{2} \operatorname{diam}(\mathcal{A})|\mathcal{A}|}{2} \epsilon_{k}+d_{2} \operatorname{diam}(\mathcal{A})|\mathcal{A}| e^{-c^{\prime} \phi(\epsilon)} \tag{1.7.9}
\end{equation*}
$$

Finally, with probability at least $1-2 \delta_{k}$,

$$
\begin{aligned}
W_{1}\left(\mathcal{L}_{k+1}, \mathcal{L}^{\star}\right) & \leq W_{1}\left(\tilde{\mathcal{L}}_{k+1}, \mathcal{L}^{\star}\right)+W_{1}\left(\tilde{\mathcal{L}}_{k+1}, \operatorname{Proj}_{S_{\epsilon}}\left(\tilde{\mathcal{L}}_{k+1}\right)\right) \\
& \leq d W_{1}\left(\mathcal{L}_{k}, \mathcal{L}^{\star}\right)+\frac{c d_{2} \operatorname{diam}(\mathcal{A})|\mathcal{A}|}{2} \epsilon_{k}+d_{2} \operatorname{diam}(\mathcal{A})|\mathcal{A}| e^{-c^{\prime} \phi(\epsilon)}+\epsilon
\end{aligned}
$$

This implies that with probability at least $1-2 \sum_{k=0}^{K-1} \delta_{k}$,

$$
\begin{align*}
W_{1}\left(\mathcal{L}_{K}, \mathcal{L}^{\star}\right) \leq & d^{K} W_{1}\left(\mathcal{L}_{0}, \mathcal{L}^{\star}\right)+\frac{c d_{2} \operatorname{diam}(\mathcal{A})|\mathcal{A}|}{2} \sum_{k=0}^{K-1} d^{K-k} \epsilon_{k}  \tag{1.7.10}\\
& +\frac{\left(d_{2} \operatorname{diam}(\mathcal{A})|\mathcal{A}| e^{-c^{\prime} \phi(\epsilon)}+\epsilon\right)\left(1-d^{K}\right)}{1-d}
\end{align*}
$$

Since $\epsilon_{k}$ is summable, we have $\sup _{k \geq 0} \epsilon_{k}<\infty$,

$$
\sum_{k=0}^{K-1} d^{K-k} \epsilon_{k} \leq \frac{\sup _{k \geq 0} \epsilon_{k}}{1-d} d^{\lfloor(K-1) / 2\rfloor}+\sum_{k=\lceil(K-1) / 2\rceil}^{\infty} \epsilon_{k} .
$$

Now plugging in $K=K_{\epsilon, \eta}$, with the choice of $\delta_{k}$ and $c=\frac{\log (1 / \epsilon)}{\phi(\epsilon)}$, and noticing that $d \in[0,1)$, we have with probability at least $1-2 \delta$,

$$
\begin{align*}
W_{1}\left(\mathcal{L}_{K_{\epsilon, \eta}}, \mathcal{L}^{\star}\right) \leq & d^{K_{\epsilon, \eta}} W_{1}\left(\mathcal{L}_{0}, \mathcal{L}^{\star}\right) \\
& +\frac{c d_{2} \operatorname{diam}(\mathcal{A})|\mathcal{A}|}{2}\left(\frac{\sup _{k \geq 0} \epsilon_{k}}{1-d} d^{\left\lfloor\left(K_{\epsilon, \eta}-1\right) / 2\right\rfloor}+\sum_{k=\left\lceil\left(K_{\epsilon, \eta}-1\right) / 2\right\rceil}^{\infty} \epsilon_{k}\right)  \tag{1.7.11}\\
& +\frac{\left(d_{2} \operatorname{diam}(\mathcal{A})|\mathcal{A}|+1\right) \epsilon}{1-d}
\end{align*}
$$

Setting $\epsilon_{k}=(k+1)^{-(1+\eta)}$, then when $K_{\epsilon, \eta} \geq 2\left(\log _{d}(\epsilon / c)+1\right)$,

$$
\frac{\sup _{k \geq 0} \epsilon_{k}}{1-d} d^{\left\lfloor\left(K_{\epsilon, \eta}-1\right) / 2\right\rfloor} \leq \frac{\epsilon / c}{1-d}
$$

Similarly, when $K_{\epsilon, \eta} \geq 2(\eta \epsilon / c)^{-1 / \eta}$,

$$
\sum_{k=\left\lceil\frac{K_{\epsilon, \eta}-1}{2}\right\rceil}^{\infty} \epsilon_{k} \leq \epsilon / c .
$$

Finally, when $K_{\epsilon, \eta} \geq \log _{d}(\epsilon /(\operatorname{diam}(\mathcal{S}) \operatorname{diam}(\mathcal{A}))), d^{K_{\epsilon, \eta}} W_{1}\left(\mathcal{L}_{0}, \mathcal{L}^{\star}\right) \leq \epsilon$, since $W_{1}\left(\mathcal{L}_{0}, \mathcal{L}^{\star}\right) \leq$ $\operatorname{diam}(\mathcal{S} \times \mathcal{A})=\operatorname{diam}(\mathcal{S}) \operatorname{diam}(\mathcal{A})$.

In summary, if $K_{\epsilon, \eta}=\left\lceil 2 \max \left\{(\eta \epsilon / c)^{-1 / \eta}, \log _{d}(\epsilon / \max \{\operatorname{diam}(\mathcal{S}) \operatorname{diam}(\mathcal{A}), c\})+1\right\}\right\rceil$, then with probability at least $1-2 \delta$,

$$
\begin{equation*}
W_{1}\left(\mathcal{L}_{K_{\epsilon, \eta}}, \mathcal{L}^{\star}\right) \leq\left(1+\frac{d_{2} \operatorname{diam}(\mathcal{A})|\mathcal{A}|(2-d)}{2(1-d)}+\frac{\left(d_{2} \operatorname{diam}(\mathcal{A})|\mathcal{A}|+1\right)}{1-d}\right) \epsilon=O(\epsilon) \tag{1.7.12}
\end{equation*}
$$

Finally, if we are using GMF-V and have assumed that $A l g$ satisfies a value-based guarantee with parameters $\left\{C_{\mathcal{M}}^{(i)}, \alpha_{1}^{(i)}, \alpha_{2}^{(i)}, \alpha_{3}^{(i)}, \alpha_{4}^{(i)}\right\}_{i=1}^{m}$, plugging in $\epsilon_{k}$ and $\delta_{k}$ into $T_{\mathcal{M}_{\mathcal{L}}}\left(\delta_{k}, \epsilon_{k}\right)$, and noticing that $k \leq K_{\epsilon, \eta}$ and $\sum_{k=0}^{K_{\epsilon, \eta}-1}(k+1)^{\alpha} \leq \frac{K_{\epsilon, \eta}^{\alpha+1}}{\alpha+1}$, we have

$$
\begin{align*}
T & =\sum_{k=0}^{K_{\epsilon, \eta}} \sum_{i=1}^{m} C_{\mathcal{M}}^{(i)}\left(\frac{1}{\epsilon_{k}}\right)^{\alpha_{1}^{(i)}}\left(\log \frac{1}{\epsilon_{k}}\right)^{\alpha_{2}^{(i)}}\left(\frac{1}{\delta_{k}}\right)^{\alpha_{3}^{(i)}}\left(\log \frac{1}{\delta_{k}}\right)^{\alpha_{4}^{(i)}} \\
& =\sum_{k=0}^{K_{\epsilon, \eta}} \sum_{i=1}^{m}(1+\eta)^{\alpha_{2}^{(i)}} C_{\mathcal{M}}^{(i)}(k+1)^{\alpha_{1}^{(i)}(1+\eta)}(\log (k+1))^{\alpha_{2}^{(i)}}\left(K_{\epsilon, \eta} / \delta\right)^{\alpha_{3}^{(i)}}\left(\log \left(K_{\epsilon, \eta} / \delta\right)\right)^{\alpha_{4}^{(i)}} \\
& \leq \sum_{i=1}^{m} \frac{(1+\eta)^{\alpha_{2}^{(i)}}}{\alpha_{1}^{(i)}(1+\eta)+1} C_{\mathcal{M}}^{(i)} K_{\epsilon, \eta}^{\alpha_{1}^{(i)}(1+\eta)+1}\left(\log \left(K_{\epsilon, \eta}+1\right)\right)^{\alpha_{2}^{(i)}}\left(K_{\epsilon, \eta} / \delta\right)_{3}^{\alpha_{3}^{(i)}}\left(\log \left(K_{\epsilon, \eta} / \delta\right)\right)^{\alpha_{4}^{(i)}} \\
& \leq \sum_{i=1}^{m} \frac{2^{\alpha_{2}^{(i)}}}{2 \alpha_{1}^{(i)}+1} C_{\mathcal{M}}^{(i)} K_{\epsilon, \eta}^{2 \alpha_{\alpha}^{(i)}+1}\left(K_{\epsilon, \eta} / \delta\right)^{\alpha_{3}^{(i)}}\left(\log \left(K_{\epsilon, \eta} / \delta\right)\right)^{\alpha_{2}^{(i)}+\alpha_{4}^{(i)}}, \tag{1.7.13}
\end{align*}
$$

which completes the proof of the value-based case.
Proof. [Proof of Theorem 1.5.5] If we use GMF-P and assume that Alg has the policy-based guarantee, then by Lemma 1.5.4,

$$
\begin{equation*}
\mathbb{P}\left(\left\|\tilde{Q}_{\tilde{\pi}_{k}}^{l_{k}}-Q_{\mathcal{L}_{k}}^{\star}\right\|_{\infty}>\epsilon\right) \leq 2 \delta . \tag{1.7.14}
\end{equation*}
$$

Hence one can simply replace $\hat{Q}_{k}^{\star}$ by $\tilde{Q}_{\tilde{\pi}_{k}}^{l_{k}}$ in the proof of Theorem 1.5.1, and obtain the same bound on $W_{1}\left(\mathcal{L}_{K_{\epsilon, \eta}}, \mathcal{L}^{\star}\right)(c f$. (1.7.12)). The only difference is that in each iteration, the required number of samples $T_{\mathcal{M}_{\mathcal{L}}}$ now has parameters $\left\{\tilde{C}_{\mathcal{M}}^{(i)}, \alpha_{1}^{(i)}, \alpha_{2}^{(i)}, \alpha_{3}^{(i)}, \alpha_{4}^{(i)}\right\}_{i=1}^{m+3}$ as defined in Lemma 1.5.4. Hence repeating the proof of (1.7.13) leads to (1.5.9).

### 1.8 Experiments

In this section, we report the performance of the proposed GMF-V-Q Algorithm and GMF-P-TRPO Algorithm with an equilibrium pricing model (see Section 1.2.3). The objectives of the experiments include 1) testing the convergence and stability of both GMF-V-Q
and GMF-P-TRPO in the GMFG setting, 2) empirically verifying the contractive property of mapping $\Gamma$, and 3) comparing GMF-V-Q and GMF-P-TRPO with existing multi-agent reinforcement learning algorithms, including the Independent Learner (IL) algorithm [150, 85] and the MF-Q ${ }^{3}$ algorithm [167]. Another set of experiments for the repeated auction model (see Section 1.2.3) is demonstrated in the short version [76].

### 1.8.1 Set-up and parameter configuration

We introduce two testing environments in our numerical experiments, one is the GMFG environment with a continuum of agents (i.e., infinite number of agents) descried in Section 1.2.3 and the other one is an N-player environment with a weak simulator.

Equilibrium price as an $N$-player game. We also consider an $N$-player game version of the equilibrium price model, which is the GMFG version described above with an $N$-player weak simulator oracle as described in Section 1.6. In particular, Take $N$ companies. At each time $t$, company $i$ decides a quantity $q_{t}^{i}$ for production and a quantity $h_{t}^{i}$ to replenish the inventory. Let $s_{t}^{i}$ denote the current inventory level of company $i$ at time $t$. Then similar to Section 1.2.3, the inventory level evolves according to

$$
s_{t+1}^{i}=s_{t}^{i}-\min \left\{q_{t}^{i}, s_{t}^{i}\right\}+h_{t}
$$

and the reward of company $i$ at time $t$ is given by

$$
r_{t}^{i}=\left(p_{t}-c_{0}\right) q_{t}^{i}-c_{1}\left(q_{t}^{i}\right)^{2}-c_{2} h_{t}^{i}-\left(c_{2}+c_{3}\right) \max \left\{q_{t}^{i}-s_{t}^{i}, 0\right\}-c_{4} s_{t}^{i} .
$$

Here $p_{t}$, the price of the product at time $t$, is determined according to the supply-demand equilibrium on the market. The total supply is $\sum_{i=1}^{N} q_{t}^{i}$, while the total demand is assumed to be $d_{N} p_{t}^{-\sigma}$, where $d_{N}=d N$ is supposed to be linearly growing as $N$ grows, i.e., the number of customers grows proportionally to the number of producers in the market. Then by equating supply and demand, we obtain that

$$
\frac{1}{N} \sum_{i=1}^{N} q_{t}^{i}=d p_{t}^{-\sigma}
$$

and by taking the limit $N \rightarrow \infty$, we obtain the mean-field counterpart (1.2.11).
In this setting, accordingly, we test the performance of GMF-VW-Q, which is GMFVW (Algorithm 5) with synchronous Q-learning and the standard softmax operator ( $c f$. Algorithm 7) and GMF-PW-TRPO, which is GMF-PW (Algorithm 6) with TRPO and the standard softmax operator. ${ }^{4}$

[^2]```
Algorithm 7 Q-learning for GMFGs (GMF-VW-Q): weak simulator
    Input: Initial \(\mathcal{L}_{0}, \epsilon\)-net \(S_{\epsilon}\), tolerances \(\epsilon_{k}, \delta_{k}>0, k=0,1, \ldots\).
    for \(k=0,1, \cdots\) do
        Perform Q-learning with hyper-parameters in Lemma 1.5.2 for \(T_{k}=T_{\mathcal{M}_{\mathcal{L}_{k}}}\left(\epsilon_{k}, \delta_{k}\right)\)
        iterations to find the approximate Q-function \(\hat{Q}_{k}^{\star}=\hat{Q}^{T_{k}}\) of the MDP \(\mathcal{M}_{\mathcal{L}_{k}}\).
        Compute \(\pi_{k} \in \Pi\) with \(\pi_{k}(s)=\operatorname{softmax}_{c}\left(\hat{Q}_{k}^{\star}(s, \cdot)\right)\).
        for \(i=1,2, \cdots, N\) do
    6: \(\quad\) Sample \(s_{i} \stackrel{\text { i.i.d. }}{\sim} \mu_{k}\), then obtain \(s_{i}^{\prime}\) i.i.d. from \(\mathcal{G}_{W}\left(s_{i}, \pi_{k}, \mathcal{L}_{k}\right)\) and \(a_{i}^{\prime} \stackrel{\text { i.i.d. }}{\sim} \pi_{k}\left(s_{i}^{\prime}\right)\).
        end for
        Compute \(\mathcal{L}_{k+1}\) with \(\mathcal{L}_{k+1}(s, a)=\frac{1}{N} \sum_{i=1}^{N} \mathbf{I}_{s_{i}^{\prime}=s, a_{i}^{\prime}=a}\).
    end for
```

Parameters. The model parameters are (unless otherwise specified): $\gamma=0.2, d=50$ and $\sigma=2 . S=Q=H=10$ and hence $|\mathcal{S}|=10$ and $|\mathcal{A}|=100 . c_{0}=0.5, c_{1}=0.1, c_{2}=0.5$, $c_{3}=0.2$ and $c_{4}=0.2$.

The algorithm parameters are (unless otherwise specified): the temperature parameter is set as $c=4.0$ and the learning rate is set as $\eta=0.01{ }^{5}$. For simplicity, we set the inner iteration $T_{k}$ to be $100 \times|\mathcal{S}| \times|\mathcal{A}|$. The $90 \%$-confidence intervals are calculated with 20 sample paths.

### 1.8.2 Performance evaluation in the GMFG setting.

Our experiments show that GMF-V-Q and GMF-P-TRPO Algorithms are efficient and robust.

Performance metric. We adopt the following metric to measure the difference between a given policy $\pi$ and an NE (here $\epsilon_{0}>0$ is a safeguard, and is taken as 0.1 in the experiments):

$$
C_{M F}(\pi)=\frac{\max _{\pi^{\prime}} \mathbb{E}_{s \sim \mu}\left[V\left(s, \pi^{\prime}, \mathcal{L}\right)\right]-V(s, \pi, \mathcal{L})}{\left|\max _{\pi^{\prime}} \mathbb{E}_{s \sim \mu}\left[V\left(s, \pi^{\prime}, \mathcal{L}\right)\right]\right|+\epsilon_{0}}
$$

Here $\mu$ is the invariant distribution of the transition matrix $P^{\pi}$, where

$$
P^{\pi}\left(s, s^{\prime}\right)=\sum_{a \in \mathcal{A}} P\left(s^{\prime} \mid s, a\right) \pi(a \mid s)
$$

for $s, s^{\prime} \in \mathcal{S}$, and $\mathcal{L}(s, a)=\mu(s) \pi(a \mid s)$ for $s, a \in \mathcal{S} \times \mathcal{A}$. Note that in the equilibrium product pricing model we are considering here, the transition model $P$ is independent of the

[^3]mean-field term $\mathcal{L}$, and hence we write $P\left(s^{\prime} \mid s, a\right)=P\left(s^{\prime} \mid s, a, \mathcal{L}\right)$. In general, an additional mean-field matching error term needs to be added into the definition of $C_{M F}(\pi)$. Clearly $C_{M F}(\pi) \geq 0$, and $C_{M F}\left(\pi^{\star}\right)=0$ if and only if $\left(\pi^{\star}, \mathcal{L}^{\star}\right)$ is an NE where $\mathcal{L}^{\star}$ is the invariant distribution of $P^{\pi^{\star}}$. A similar metric without normalization has been adopted in [41].

Contractiveness of mapping $\Gamma$. As explained in Remark 1.3.1 from Section 1.3, the contractiveness property of $\Gamma$ is the key for establishing the uniqueness of MFG solution and hence the convergence of the GMFG algorithm. To empirically verify whether this property holds for the equilibrium price example, we plot the value of $\frac{\left\|\Gamma\left(\mathcal{L}_{1}\right)-\Gamma\left(\mathcal{L}_{2}\right)\right\|_{1}}{\left\|\mathcal{L}_{1}-\mathcal{L}_{1}\right\|_{1}}$ for randomly generated state-action distributions $\mathcal{L}_{1}$ and $\mathcal{L}_{2}$. Technically speaking, $\Gamma$ is contractive and there exists a unique MFG solution if the value of $\frac{\left\|\Gamma\left(\mathcal{L}_{1}\right)-\Gamma\left(\mathcal{L}_{2}\right)\right\|_{1}}{\left\|\mathcal{L}_{1}-\mathcal{L}_{2}\right\|_{1}}$ is smaller than one for all choices of $\mathcal{L}_{1}$ and $\mathcal{L}_{2}$.

W observe from Figure 1.1 that, with various of choices of different model parameters, the quantity $\frac{\left\|\Gamma\left(\mathcal{L}_{1}\right)-\Gamma\left(\mathcal{L}_{2}\right)\right\|_{1}}{\left\|\mathcal{L}_{1}-\mathcal{L}_{2}\right\|_{1}}$ is always smaller than 0.3 indicating that $\Gamma$ is contractive.


Figure 1.1: Histogram of $\frac{\left\|\Gamma\left(\mathcal{L}_{1}\right)-\Gamma\left(\mathcal{L}_{2}\right)\right\|_{1}}{\left\|\mathcal{L}_{1}-\mathcal{L}_{2}\right\|_{1}}$ under various settings $\left(\mathcal{L}_{1}\right.$ and $\mathcal{L}_{2}$ are randomly sampled according to the uniform distribution).

Convergence and stability. Both GMF-V-Q and GMF-P-TRPO are efficient and robust. First, both GMF-V-Q and GMF-P-TRPO converge within about 5 outer iterations; secondly, as the number of inner iterations increases, the error decreases (Figure 1.2); and finally, the convergence is robust with respect to both the change of number of states and
actions (Figure 1.3). The performance of GMF-V-Q is (slightly) more stable than GMF-P-TRPO with a smaller variance across 20 repeated experiments (see Figure 1.2a versus Figure 1.2b or Figure 1.3a versus Figure 1.3b). This is due to the fact that GMF-P-TRPO uses asynchronous updates, which leads to slightly less stable performance compared to GMF-V-Q, which uses synchronous updates.

In contrast, the Naive algorithms, i.e., GMF-V-Q without smoothing (denoted as GMF-V-Q-nonsmoothing) and GMF-P-TRPO without smoothing (denoted as GMF-P-TRPOnonsmoothing), do not converge even with 50 outer iterations and $200 \times|\mathcal{S}| \times|\mathcal{A}|$ inner iterations within each outer iteration. In particular, GMF-V-Q-nonsmoothing and GMF-P-TRPO-nonsmoothing present different unstable behaviors (see Figure 1.4). The joint distribution $\mathcal{L}_{t}$ from GMF-V-Q-nonsmoothing keeps fluctuating (Figure 1.4a) whereas the joint distribution $\mathcal{L}_{t}$ from GMF-P-TRPO (without smoothing) is trapped around the initialization which is far away from the true equilibrium distribution (Figure 1.4 b ).


Figure 1.2: Convergence with different number of inner iterations $(|\mathcal{A}|=100$ and $|\mathcal{S}|=10)$.


Figure 1.3: Convergence with different size of state space and action space.


Figure 1.4: Fluctuations of algorithms without smoothing (Dotted black line: theoretical value of the equilibrium price).

Model verification and interpretation of equilibrium scenario. In Figures 1.5 and 1.6, we run both algorithms for 20 outer iterations with the same number of inner iterations $(100,000=100 \times|\mathcal{A}| \times|\mathcal{S}|)$ within each outer iteration. The final equilibrium inventory distribution and production distribution from both algorithms are close to each other.

(a) Equilibrium inventory dis- (b) Equilibrium production (c) Equilibrium price at the
tribution.
distribution.

Figure 1.5: GMF-V-Q versus GMF-P-TRPO ( $\sigma=1.3$ and one trajectory).


Figure 1.6: GMF-V-Q versus GMF-P-TRPO ( $\sigma=2.0$ and one trajectory).

Note that the demand elasticity $\sigma$ captures how sensitive demand for a product is compared to the changes in other economic factors, such as price or income. When $\sigma$ is increased from 1.3 to 2.0 indicating that the demand is more sensitive to price rise, the equilibrium price decreases from 10.9 to 7.3 (see Figures 1.5 c and 1.6 c ) and the distribution of the equilibrium production level is centered towards smaller values (see Figures 1.5b and 1.6b). The equilibrium inventory level has a huge mass at 0 when $\sigma=1.3$. This implies that producers do not keep large inventories and pay the inventory cost in the equilibrium. On the other hand, the equilibrium inventory is more uniformly distributed when $\sigma=2$.

### 1.8.3 Performance evaluation in the N-player setting

Performance metric. Similar to the performance metric introduced in Section 1.8.2 for the GMFG setting, we adopt the following metric to measure the difference between a given policy $\pi$ and an NE under THE N-player setting (here $\epsilon_{0}>0$ is a safeguard, and is taken as 0.1 in the experiments):

$$
C(\boldsymbol{\pi})=\frac{1}{N|\mathcal{S}|^{N}} \sum_{i=1}^{N} \sum_{\boldsymbol{s} \in \mathcal{S}^{N}} \frac{\max _{\pi^{i}} V^{i}\left(\boldsymbol{s},\left(\boldsymbol{\pi}^{-i}, \pi^{i}\right)\right)-V^{i}(\boldsymbol{s}, \boldsymbol{\pi})}{\left|\max _{\pi^{i}} V^{i}\left(\boldsymbol{s},\left(\boldsymbol{\pi}^{-i}, \pi^{i}\right)\right)\right|+\epsilon_{0}} .
$$

Clearly $C(\boldsymbol{\pi}) \geq 0$, and $C\left(\boldsymbol{\pi}^{\star}\right)=0$ if and only if $\boldsymbol{\pi}^{\star}$ is an NE. Policy $\arg \max _{\pi_{i}} V_{i}\left(\boldsymbol{s},\left(\boldsymbol{\pi}^{-i}, \pi_{i}\right)\right)$ is called the best response to $\pi^{-i}$. A similar metric without normalization has been adopted in [128].

Existing algorithms for $N$-player games. To test the effectiveness of GMF-VW-Q for approximating $N$-player games, we next compare GMF-VW-Q with the IL algorithm and the MF-Q algorithm. The IL algorithm [150] considers $N$ independent players and each player solves a decentralized reinforcement learning problem ignoring other players in the system. The MF-Q algorithm [167] extends the NASH-Q Learning algorithm for the $N$-player game introduced in [84], adds the aggregate actions $\left(\overline{\boldsymbol{a}}_{-i}=\frac{\sum_{j \neq-} a_{j}}{N-1}\right)$ from the opponents, and works for the class of games where the interactions are only through the average actions of $N$ players.

Results and analysis. Our experiment (Figure 1.7) shows that GMF-VW-Q and GMF-PW-TRPO achieve similar performance, and both of them are superior in terms of convergence rate, accuracy, and stability for approximating an $N$-player game. In general, both algorithms converge faster than IL and MF-Q and achieve the smallest errors.

For instance, when $N=20$, IL Algorithm converges with the largest error 0.220 . The error from MF-Q is 0.101 , smaller than IL but still bigger than the error from GMF-VW-Q. The GMF-VW-Q and GMF-PW-TRPO converge with the lowest error 0.065. Moreover, as $N$ increases, the error of GMF-VW-Q and GMF-PW-TRPO decease while the errors of both MF-Q and IL increase significantly. As $|\mathcal{S}|$ and $|\mathcal{A}|$ increase, GMF-VW-Q and GMF-PW-TRPO are robust with respect to this increase of dimensionality, while both MF-Q and

IL clearly suffer from the increase of the dimensionality with decreased convergence rate and accuracy. Therefore, GMF-VW-Q and GMF-PW-TRPO are more scalable than IL and MF-Q, when the system is complex and the number of players $N$ is large.


Figure 1.7: Learning accuracy based on $C(\boldsymbol{\pi})$.

### 1.9 Extension: Existence and uniqueness for non-stationary NE of GMFGs

In this section, we describe the setting of non-stationary NE for GMFGs and establish the corresponding results of existence and uniqueness.

Definition 1.9.1 (NE for GMFGs). In (GMFG), a player-population profile $\left(\boldsymbol{\pi}^{\star}, \mathcal{L}^{\star}\right):=$ $\left(\left\{\pi_{t}^{\star}\right\}_{t=0}^{\infty},\left\{\mathcal{L}_{t}^{\star}\right\}_{t=0}^{\infty}\right)$ is called an NE if

1. (Single player side) Fix $\mathcal{L}^{\star}$, for any policy sequence $\boldsymbol{\pi}:=\left\{\pi_{t}\right\}_{t=0}^{\infty}$ and initial state $s \in \mathcal{S}$,

$$
\begin{equation*}
V\left(s, \boldsymbol{\pi}^{\star}, \mathcal{L}^{\star}\right) \geq V\left(s, \boldsymbol{\pi}, \mathcal{L}^{\star}\right) \tag{1.9.1}
\end{equation*}
$$

2. (Population side) $\mathbb{P}_{s_{t}, a_{t}}=\mathcal{L}_{t}^{\star}$ for all $t \geq 0$, where $\left\{s_{t}, a_{t}\right\}_{t=0}^{\infty}$ is the dynamics under the policy sequence $\boldsymbol{\pi}^{\star}$ starting from $s_{0} \sim \mu_{0}^{\star}$, with $a_{t} \sim \pi_{t}^{\star}\left(s_{t}, \mu_{t}^{\star}\right), s_{t+1} \sim P\left(\cdot \mid s_{t}, a_{t}, \mathcal{L}_{t}^{\star}\right)$, and $\mu_{t}^{\star}$ being the population state marginal of $\mathcal{L}_{t}^{\star}$.

Step A. Fix $\mathcal{L}:=\left\{\mathcal{L}_{t}\right\}_{t=0}^{\infty}$, (GMFG) becomes the classical optimization problem. Indeed, with $\mathcal{L}$ fixed, the population state distribution sequence $\boldsymbol{\mu}:=\left\{\mu_{t}\right\}_{t=0}^{\infty}$ is also fixed, hence the space of admissible policies is reduced to the single-player case. Solving (GMFG) is now reduced to finding a policy sequence $\pi_{t, \mathcal{L}}^{\star} \in \Pi:=\{\pi \mid \pi: \mathcal{S} \rightarrow \mathcal{P}(\mathcal{A})\}$ over all admissible $\boldsymbol{\pi}_{\mathcal{L}}=\left\{\pi_{t, \mathcal{L}}\right\}_{t=0}^{\infty}$, to maximize

$$
\begin{array}{ll}
V\left(s, \boldsymbol{\pi}_{\mathcal{L}}, \mathcal{L}\right):= & \mathbb{E}\left[\sum_{t=0}^{\infty} \gamma^{t} r\left(s_{t}, a_{t}, \mathcal{L}_{t}\right) \mid s_{0}=s\right] \\
\text { subject to } & s_{t+1} \sim P\left(s_{t}, a_{t}, \mathcal{L}_{t}\right), \quad a_{t} \sim \pi_{t, \mathcal{L}}\left(s_{t}\right)
\end{array}
$$

Notice that with $\mathcal{L}$ fixed, one can safely suppress the dependency on $\mu_{t}$ in the admissible policies. Moreover, given this fixed $\mathcal{L}$ sequence and the solution $\pi_{\mathcal{L}}^{\star}:=\left\{\pi_{t, \mathcal{L}}^{\star}\right\}_{t=0}^{\infty}$, one can define a mapping from the fixed population distribution sequence $\mathcal{L}$ to an optimal randomized policy sequence. That is,

$$
\Gamma_{1}:\{\mathcal{P}(\mathcal{S} \times \mathcal{A})\}_{t=0}^{\infty} \rightarrow\{\Pi\}_{t=0}^{\infty}
$$

such that $\boldsymbol{\pi}_{\mathcal{L}}^{\star}=\Gamma_{1}(\mathcal{L})$. Note that this $\boldsymbol{\pi}_{\mathcal{L}}^{\star}$ sequence satisfies the single player side condition in Definition 1.9.1 for the population state-action pair sequence $\mathcal{L}$. That is, $V\left(s, \boldsymbol{\pi}_{\mathcal{L}}^{\star}, \mathcal{L}\right) \geq$ $V(s, \boldsymbol{\pi}, \mathcal{L})$, for any policy sequence $\boldsymbol{\pi}=\left\{\pi_{t}\right\}_{t=0}^{\infty}$ and any initial state $s \in \mathcal{S}$.

Accordingly, a similar feedback regularity condition is needed in this step.
Assumption 3. There exists a constant $d_{1} \geq 0$, such that for any $\mathcal{L}, \mathcal{L}^{\prime} \in\{\mathcal{P}(\mathcal{S} \times \mathcal{A})\}_{t=0}^{\infty}$,

$$
\begin{equation*}
D\left(\Gamma_{1}(\mathcal{L}), \Gamma_{1}\left(\mathcal{L}^{\prime}\right)\right) \leq d_{1} \mathcal{W}_{1}\left(\mathcal{L}, \mathcal{L}^{\prime}\right) \tag{1.9.2}
\end{equation*}
$$

where

$$
\begin{align*}
D\left(\boldsymbol{\pi}, \boldsymbol{\pi}^{\prime}\right) & :=\sup _{s \in \mathcal{S}} \mathcal{W}_{1}\left(\boldsymbol{\pi}(s), \boldsymbol{\pi}^{\prime}(s)\right)=\sup _{s \in \mathcal{S}} \sup _{t \in \mathbb{N}} W_{1}\left(\pi_{t}(s), \pi_{t}^{\prime}(s)\right) \\
\mathcal{W}_{1}\left(\mathcal{L}, \mathcal{L}^{\prime}\right) & :=\sup _{t \in \mathbb{N}} W_{1}\left(\mathcal{L}_{t}, \mathcal{L}_{t}^{\prime}\right) \tag{1.9.3}
\end{align*}
$$

and $W_{1}$ is the $\ell_{1}$-Wasserstein distance between probability measures.

Step B. Based on the analysis in Step A and $\boldsymbol{\pi}_{\mathcal{L}}^{\star}=\left\{\pi_{t, \mathcal{L}}^{\star}\right\}_{t=0}^{\infty}$, update the initial sequence $\mathcal{L}$ to $\mathcal{L}^{\prime}$ following the controlled dynamics $P\left(\cdot \mid s_{t}, a_{t}, \mathcal{L}_{t}\right)$.

Accordingly, for any admissible policy sequence $\boldsymbol{\pi} \in\{\Pi\}_{t=0}^{\infty}$ and a joint population stateaction pair sequence $\mathcal{L} \in\{\mathcal{P}(\mathcal{S} \times \mathcal{A})\}_{t=0}^{\infty}$, define a mapping $\Gamma_{2}:\{\Pi\}_{t=0}^{\infty} \times\{\mathcal{P}(\mathcal{S} \times \mathcal{A})\}_{t=0}^{\infty} \rightarrow$ $\{\mathcal{P}(\mathcal{S} \times \mathcal{A})\}_{t=0}^{\infty}$ as follows:

$$
\begin{equation*}
\Gamma_{2}(\boldsymbol{\pi}, \mathcal{L}):=\hat{\mathcal{L}}=\left\{\mathbb{P}_{s_{t}, a_{t}}\right\}_{t=0}^{\infty} \tag{1.9.4}
\end{equation*}
$$

where $s_{t+1} \sim \mu_{t} P\left(\cdot \mid \cdot, a_{t}, \mathcal{L}_{t}\right), a_{t} \sim \pi_{t}\left(s_{t}\right), s_{0} \sim \mu_{0}$, and $\mu_{t}$ is the population state marginal of $\mathcal{L}_{t}$.

One also needs a similar assumption in this step.
Assumption 4. There exist constants $d_{2}, d_{3} \geq 0$, such that for any admissible policy sequences $\boldsymbol{\pi}, \boldsymbol{\pi}^{1}, \boldsymbol{\pi}^{2}$ and joint distribution sequences $\mathcal{L}, \mathcal{L}^{1}, \boldsymbol{L}^{2}$,

$$
\begin{gather*}
\mathcal{W}_{1}\left(\Gamma_{2}\left(\boldsymbol{\pi}^{1}, \mathcal{L}\right), \Gamma_{2}\left(\boldsymbol{\pi}^{2}, \mathcal{L}\right)\right) \leq d_{2} D\left(\boldsymbol{\pi}^{1}, \boldsymbol{\pi}^{2}\right)  \tag{1.9.5}\\
\mathcal{W}_{1}\left(\Gamma_{2}\left(\boldsymbol{\pi}, \mathcal{L}^{1}\right), \Gamma_{2}\left(\boldsymbol{\pi}, \mathcal{L}^{2}\right)\right) \leq d_{3} \mathcal{W}_{1}\left(\boldsymbol{L}^{1}, \boldsymbol{L}^{2}\right) \tag{1.9.6}
\end{gather*}
$$

Similarly, Assumption 4 can be reduced to Lipschitz continuity and boundedness of the transition dynamics $P$ under certain conditions.

Step C. Repeat Step A and Step B until $\mathcal{L}^{\prime}$ matches $\mathcal{L}$.
This step is to take care of the population side condition. To ensure the convergence of the combined step $A$ and step $B$, it suffices if $\Gamma:\{\mathcal{P}(\mathcal{S} \times \mathcal{A})\}_{t=0}^{\infty} \rightarrow\{\mathcal{P}(\mathcal{S} \times \mathcal{A})\}_{t=0}^{\infty}$ is a contractive mapping under the $\mathcal{W}_{1}$ distance, with $\Gamma(\mathcal{L}):=\Gamma_{2}\left(\Gamma_{1}(\mathcal{L}), \mathcal{L}\right)$. Then by the Banach fixed point theorem and the completeness of the related metric spaces, there exists a unique NE to the GMFG.

In summary, we have
Theorem 1.9.1 (Existence and Uniqueness of GMFG solution). Given Assumptions 3 and 4, and assuming that $d_{1} d_{2}+d_{3}<1$, there exists a unique NE to (GMFG).

The proof of Theorem 1.9 .1 can be established by modifying appropriately the fixed-point approach for the stationary GMFG in Theorem 1.3.1.

### 1.10 Appendix

### 1.10.1 Distance metrics and completeness

This section reviews some basic properties of the Wasserstein distance. It then proves that the metrics defined in the main text are indeed distance functions and define complete metric spaces.
$\ell_{1}$-Wasserstein distance and dual representation. The $\ell_{1}$ Wasserstein distance over $\mathcal{P}(\mathcal{X})$ for $\mathcal{X} \subseteq \mathbb{R}^{k}$ is defined as

$$
\begin{equation*}
W_{1}\left(\nu, \nu^{\prime}\right):=\inf _{M \in \mathcal{M}\left(\nu, \nu^{\prime}\right)} \int_{\mathcal{X} \times \mathcal{X}}\|x-y\|_{2} \mathrm{~d} M(x, y) . \tag{1.10.1}
\end{equation*}
$$

where $\mathcal{M}\left(\nu, \nu^{\prime}\right)$ is the set of all measures (couplings) on $\mathcal{X} \times \mathcal{X}$, with marginals $\nu$ and $\nu^{\prime}$ on the two components, respectively.

The Kantorovich duality theorem enables the following equivalent dual representation of $W_{1}$ :

$$
\begin{equation*}
W_{1}\left(\nu, \nu^{\prime}\right)=\sup _{\|f\|_{L} \leq 1}\left|\int_{\mathcal{X}} f d \nu-\int_{\mathcal{X}} f d \nu^{\prime}\right|, \tag{1.10.2}
\end{equation*}
$$

where the supremum is taken over all 1-Lipschitz functions $f$, i.e., $f$ satisfying $|f(x)-f(y)| \leq$ $\|x-y\|_{2}$ for all $x, y \in \mathcal{X}$.

The Wasserstein distance $W_{1}$ can also be related to the total variation distance via the following inequalities [67]:

$$
\begin{equation*}
d_{\min }(\mathcal{X}) d_{T V}\left(\nu, \nu^{\prime}\right) \leq W_{1}\left(\nu, \nu^{\prime}\right) \leq \operatorname{diam}(\mathcal{X}) d_{T V}\left(\nu, \nu^{\prime}\right) \tag{1.10.3}
\end{equation*}
$$

where $d_{\min }(\mathcal{X})=\min _{x \neq y \in \mathcal{X}}\|x-y\|_{2}$, which is guaranteed to be positive when $\mathcal{X}$ is finite.

When $\mathcal{S}$ and $\mathcal{A}$ are compact, for any compact subset $\mathcal{X} \subseteq \mathbb{R}^{k}$, and for any $\nu, \nu^{\prime} \in \mathcal{P}(\mathcal{X})$, $W_{1}\left(\nu, \nu^{\prime}\right) \leq \operatorname{diam}(\mathcal{X}) d_{T V}\left(\nu, \nu^{\prime}\right) \leq \operatorname{diam}(\mathcal{X})<\infty$, where $\operatorname{diam}(\mathcal{X})=\sup _{x, y \in \mathcal{X}}\|x-y\|_{2}$ and $d_{T V}$ is the total variation distance. Moreover, one can verify

Lemma 1.10.1. Both $D$ and $\mathcal{W}_{1}$ are distance functions, and they are finite for any input distribution pairs. In addition, both $\left(\{\Pi\}_{t=0}^{\infty}, D\right)$ and $\left(\{\mathcal{P}(\mathcal{S} \times \mathcal{A})\}_{t=0}^{\infty}, \mathcal{W}_{1}\right)$ are complete metric spaces.

These facts enable the usage of Banach fixed-point mapping theorem for the proof of existence and uniqueness (Theorems 1.9.1 and 1.3.1).

Proof. [Proof of Lemma 1.10.1] It is known that for any compact set $\mathcal{X} \subseteq \mathbb{R}^{k},\left(\mathcal{P}(\mathcal{X}), W_{1}\right)$ defines a complete metric space [25]. Since $W_{1}\left(\nu, \nu^{\prime}\right) \leq \operatorname{diam}(\mathcal{X})$ is uniformly bounded for any $\nu, \nu^{\prime} \in \mathcal{P}(\mathcal{X})$, we know that $\mathcal{W}_{1}\left(\mathcal{L}, \mathcal{L}^{\prime}\right) \leq \operatorname{diam}(\mathcal{X})$ and $D\left(\boldsymbol{\pi}, \boldsymbol{\pi}^{\prime}\right) \leq \operatorname{diam}(\mathcal{X})$ as well, so they are both finite for any input distribution pairs. It is clear that they are distance functions based on the fact that $W_{1}$ is a distance function.

Finally, we show the completeness of the two metric spaces $\left(\{\Pi\}_{t=0}^{\infty}, D\right)$ and $(\{\mathcal{P}(\mathcal{S} \times$ $\left.\mathcal{A})\}_{t=0}^{\infty}, \mathcal{W}_{1}\right)$. Take $\left(\{\Pi\}_{t=0}^{\infty}, D\right)$ for example. Suppose that $\pi^{k}$ is a Cauchy sequence in $\left(\{\Pi\}_{t=0}^{\infty}, D\right)$. Then for any $\epsilon>0$, there exists a positive integer $N$, such that for any $m, n \geq N$,

$$
\begin{equation*}
D\left(\boldsymbol{\pi}^{n}, \boldsymbol{\pi}^{m}\right) \leq \epsilon \Longrightarrow W_{1}\left(\pi_{t}^{n}(s), \pi_{t}^{m}(s)\right) \leq \epsilon \text { for any } s \in \mathcal{S}, t \in \mathbb{N} \tag{1.10.4}
\end{equation*}
$$

which implies that $\pi_{t}^{k}(s)$ forms a Cauchy sequence in $\left(\mathcal{P}(\mathcal{A}), W_{1}\right)$, and hence by the completeness of $\left(\mathcal{P}(\mathcal{A}), W_{1}\right), \pi_{t}^{k}(s)$ converges to some $\pi_{t}(s) \in \mathcal{P}(\mathcal{A})$. As a result, $\boldsymbol{\pi}^{n} \rightarrow \boldsymbol{\pi} \in\{\Pi\}_{t=0}^{\infty}$ under metric $D$, which shows that $\left(\{\Pi\}_{t=0}^{\infty}, D\right)$ is complete.

The completeness of $\left(\{\mathcal{P}(\mathcal{S} \times \mathcal{A})\}_{t=0}^{\infty}, \mathcal{W}_{1}\right)$ can be proved similarly.
The same argument for Lemma 1.10 .1 shows that both $D$ and $W_{1}$ are distance functions and are finite for any input distribution pairs, with both $(\Pi, D)$ and $\left(\mathcal{P}(\mathcal{S} \times \mathcal{A}), W_{1}\right)$ again complete metric spaces.

### 1.10.2 Bounds for GMF-V-Q using asynchronous Q-learning

In the main text, we have shown the results by using synchronous Q-learning algorithm. Here for the completeness, we also show the corresponding results for asynchronous Qlearning algorithm.

For asynchronous Q-learning algorithm, at each step $l$ with the state $s$ and an action $a$, the system reaches state $s^{\prime}$ according to the controlled dynamics and the Q-function approximation $Q_{l}$ is updated according to

$$
\begin{equation*}
\hat{Q}^{l+1}(s, a)=\left(1-\beta_{l}(s, a)\right) \hat{Q}^{l}(s, a)+\beta_{l}(s, a)\left[r(s, a)+\gamma \max _{\bar{a}} \hat{Q}^{l}\left(s^{\prime}, \bar{a}\right)\right] \tag{1.10.5}
\end{equation*}
$$

where $\hat{Q}^{0}(s, a)=C$ for some constant $C \in \mathbb{R}$ for any $s \in \mathcal{S}$ and $a \in \mathcal{A}$, and the step size $\beta_{l}(s, a)$ can be chosen as ([53])

$$
\beta_{l}(s, a)= \begin{cases}|\#(s, a, l)+1|^{-h}, & (s, a)=\left(s_{l}, a_{l}\right)  \tag{1.10.6}\\ 0, & \text { otherwise }\end{cases}
$$

with $h \in(1 / 2,1)$. Here $\#(s, a, l)$ is the number of times up to time $l$ that one visits the state-action pair $(s, a)$. The algorithm then proceeds to choose action $a^{\prime}$ based on $\hat{Q}^{l+1}$ with appropriate exploration strategies, including the $\epsilon$-greedy strategy.

Lemma 1.10.2 ([53]: sample complexity of asynchronous Q-learning). For an MDP, say $\mathcal{M}=(\mathcal{S}, \mathcal{A}, P, r, \gamma)$, suppose that the $Q$-learning algorithm takes step-sizes (1.5.4). Also suppose that the covering time of the state-action pairs is bounded by $L$ with probability at least $1-p$ for some $p \in(0,1)$. Then $\left\|\hat{Q}^{T_{\mathcal{M}}(\delta, \epsilon)}-Q_{\mathcal{M}}^{\star}\right\|_{\infty} \leq \epsilon$ with probability at least $1-2 \delta$. Here $\hat{Q}^{T}$ is the $T$-th update in the $Q$-learning updates (1.5.3), $Q_{\mathcal{M}}^{\star}$ is the (optimal) $Q$-function, and

$$
T_{\mathcal{M}}(\delta, \epsilon)=\Omega\left(\left(\frac{L \log _{p}(\delta)}{\beta} \log \frac{V_{\max }}{\epsilon}\right)^{\frac{1}{1-h}}+\left(\frac{\left(L \log _{p}(\delta)\right)^{1+3 h} V_{\max }^{2} \log \left(\frac{|\mathcal{S} \| \mathcal{A}| V_{\max }}{\delta \beta \epsilon}\right)}{\beta^{2} \epsilon^{2}}\right)^{\frac{1}{h}}\right)
$$

where $\beta=(1-\gamma) / 2, V_{\max }=R_{\max } /(1-\gamma)$, and $R_{\max }$ is such that a.s. $0 \leq r(s, a) \leq R_{\max }$.

Here the covering time $L$ of a state-action pair sequence is defined to be the number of steps needed to visit all state-action pairs starting from any arbitrary state-action pair. Also notice that the $l_{\infty}$ norm above is defined in an element-wise sense, i.e., for $M \in \mathbb{R}^{|\mathcal{S}| \times|\mathcal{A}|}$, we have $\|M\|_{\infty}=\max _{s \in \mathcal{S}, a \in \mathcal{A}}|M(s, a)|$.

Corollary 1.10.3 (Value-based guarantee of asynchronous Q-learning algorithm). The asynchronous Q-learning algorithm with appropriate choices of step-sizes (cf. (1.5.4)) satisfies the following value-based guarantee, where $C_{\mathcal{M}}^{(i)}(i=1,2,3)$ are constants depending on $|\mathcal{S}|,|\mathcal{A}|, V_{\max }, \beta$ and $h$, and we have:

$$
\begin{array}{r}
\alpha_{2}^{(1)}=\alpha_{4}^{(1)}=\frac{1}{1-h}, \alpha_{1}^{(1)}=\alpha_{3}^{(2)}=0, \\
\alpha_{1}^{(2)}=\frac{2}{h}, \alpha_{4}^{(2)}=\frac{2+3 h}{h}, \alpha_{j}^{(2)}=0 \text { for } j=2,3, \\
\alpha_{1}^{(3)}=\frac{2}{h}, \alpha_{2}^{(3)}=\frac{1}{h}, \alpha_{4}^{(3)}=\frac{1+3 h}{h}, \alpha_{3}^{(3)}=0
\end{array}
$$

In addition, assume the same assumptions as Theorem 1.3.1, then for Algorithm 3 with asynchronous $Q$-learning method, with probability at least $1-2 \delta, W_{1}\left(\mathcal{L}_{K_{\epsilon, \eta}}, \mathcal{L}^{\star}\right) \leq C_{0} \epsilon$, where $K_{\epsilon, \eta}$
is defined as in Theorem 1.5.1. And the total number of samples $T=\sum_{k=0}^{K_{\epsilon, \eta}-1} T_{\mathcal{M}_{\mathcal{L}_{k}}}\left(\delta_{k}, \epsilon_{k}\right)$ is bounded by

$$
T \leq O\left(K_{\epsilon, \eta}^{\frac{4}{h}+1}\left(\log \frac{K_{\epsilon, \eta}}{\delta}\right)^{3+\frac{2}{h}}+\left(\log \frac{K_{\epsilon, \eta}}{\delta}\right)^{\frac{2}{1-h}}\right)
$$

### 1.10.3 Weak simulator

In this section, we state the counterpart of Theorems 1.5.1 and 1.5.5 for Algorithms 5 and 6 , respectively. Notice that here the major difference is the additional $O(1 / \sqrt{N})$ term.

We first (re)state the relation between $\mathbf{E m p}_{N}$ (which serves as a $1 / N$-net) and action gaps:

For any positive integer $N$, there exist a positive constant $\phi_{N}>0$, with the property that $\max _{a^{\prime} \in \mathcal{A}} Q_{\mathcal{L}}^{\star}\left(s, a^{\prime}\right)-Q_{\mathcal{L}}^{\star}(s, a) \geq \phi_{N}$ for any $\mathcal{L} \in \boldsymbol{E m p}_{N}, s \in \mathcal{S}$, and any $\notin$ $\operatorname{argmax}_{a \in \mathcal{A}} Q_{\mathcal{L}}^{\star}(s, a)$.

Now we are ready to state the convergence results.
Theorem 1.10.4 (Convergence and complexity of GMF-VW). Assume the same assumptions as Theorem 1.3.1. Suppose that Alg has a value-based guarantee with parameters

$$
\left\{C_{\mathcal{M}}^{(i)}, \alpha_{1}^{(i)}, \alpha_{2}^{(i)}, \alpha_{3}^{(i)}, \alpha_{4}^{(i)}\right\}_{i=1}^{m}
$$

For any $\epsilon, \delta>0$, set $\delta_{k}=\delta / K_{\epsilon, \eta}, \epsilon_{k}=(k+1)^{-(1+\eta)}$ for some $\eta \in(0,1]\left(k=0, \ldots, K_{\epsilon, \eta}-\right.$ 1), and $c^{\prime}=c=\frac{\log (1 / \epsilon)}{\phi_{N}} .{ }^{6}$ Then with probability at least $1-4 \delta$,

$$
W_{1}\left(\mathcal{L}_{K_{\epsilon, \eta}}, \mathcal{L}^{\star}\right) \leq C \epsilon+\frac{\operatorname{diam}(\mathcal{S}) \operatorname{diam}(\mathcal{A})|\mathcal{S}||\mathcal{A}|}{2(1-d)} \sqrt{\frac{1}{2 N} \log \left(|\mathcal{S}||\mathcal{A}| K_{\epsilon, \eta} / \delta\right)}
$$

Here $K_{\epsilon, \eta}:=\left\lceil 2 \max \left\{(\eta \epsilon / c)^{-1 / \eta}, \log _{d}(\epsilon / \max \{\operatorname{diam}(\mathcal{S}) \operatorname{diam}(\mathcal{A}), c\})+1\right\}\right\rceil$ is the number of outer iterations, and the constant $C$ is independent of $\delta, \epsilon$ and $\eta$.

Moreover, the total number of samples $T=\sum_{k=0}^{K_{\epsilon}, \eta} T_{\mathcal{M}_{\mathcal{L}_{k}}}\left(\delta_{k}, \epsilon_{k}\right)$ is bounded by

$$
\begin{equation*}
T \leq \sum_{i=1}^{m} \frac{2^{\alpha_{2}^{(i)}}}{2 \alpha_{1}^{(i)}+1} C_{\mathcal{M}}^{(i)} K_{\epsilon, \eta}^{2 \alpha_{1}^{(i)}+1}\left(K_{\epsilon, \eta} / \delta\right)^{\alpha_{3}^{(i)}}\left(\log \left(K_{\epsilon, \eta} / \delta\right)\right)^{\alpha_{2}^{(i)}+\alpha_{4}^{(i)}} \tag{1.10.7}
\end{equation*}
$$

Theorem 1.10.5 (Convergence and complexity of GMF-PW). Assume the same assumptions as in Theorem 1.3.1. Suppose that Alg has a policy-based guarantee with parameters

$$
\left\{C_{\mathcal{M}}^{(i)}, \alpha_{1}^{(i)}, \alpha_{2}^{(i)}, \alpha_{3}^{(i)}, \alpha_{4}^{(i)}\right\}_{i=1}^{m}
$$

[^4]Then for any $\epsilon, \delta>0$, set $\delta_{k}=\delta / K_{\epsilon, \eta}, \epsilon_{k}=(k+1)^{-(1+\eta)}$ for some $\eta \in(0,1] \quad(k=$ $\left.0, \ldots, K_{\epsilon, \eta}-1\right)$, and $c^{\prime}=c=\frac{\log (1 / \epsilon)}{\phi_{N}},^{7}$ with probability at least $1-4 \delta$,

$$
W_{1}\left(\mathcal{L}_{K_{\epsilon, \eta}}, \mathcal{L}^{\star}\right) \leq C \epsilon+\frac{\operatorname{diam}(\mathcal{S}) \operatorname{diam}(\mathcal{A})|\mathcal{S}||\mathcal{A}|}{2(1-d)} \sqrt{\frac{1}{2 N} \log \left(|\mathcal{S} \| \mathcal{A}| K_{\epsilon, \eta} / \delta\right)}
$$

Here $K_{\epsilon, \eta}:=\left\lceil 2 \max \left\{(\eta \epsilon / c)^{-1 / \eta}, \log _{d}(\epsilon / \max \{\operatorname{diam}(\mathcal{S}) \operatorname{diam}(\mathcal{A}), c\})+1\right\}\right\rceil$ is the number of outer iterations, and the constant $C$ is independent of $\delta, \epsilon$ and $\eta$.

Moreover, the total number of samples $T=\sum_{k=0}^{K_{\epsilon, \eta}-1} T_{\mathcal{M}_{\mathcal{L}_{k}}}\left(\delta_{k}, \epsilon_{k}\right)$ is bounded by

$$
\begin{equation*}
T \leq \sum_{i=1}^{m+1} \frac{2^{\alpha_{2}^{(i)}}}{2 \alpha_{1}^{(i)}+1} \tilde{C}_{\mathcal{M}}^{(i)} K_{\epsilon, \eta}^{2 \alpha_{\eta}^{(i)}+1}\left(K_{\epsilon, \eta} / \delta\right)^{\alpha_{3}^{(i)}}\left(\log \left(K_{\epsilon, \eta} / \delta\right)\right)^{\alpha_{2}^{(i)}+\alpha_{4}^{(i)}} \tag{1.10.8}
\end{equation*}
$$

where the parameters $\left\{\tilde{C}_{\mathcal{M}}^{(i)}, \alpha_{1}^{(i)}, \alpha_{2}^{(i)}, \alpha_{3}^{(i)}, \alpha_{4}^{(i)}\right\}_{i=1}^{m+1}$ are defined in Lemma 1.5.4.
The key to the proof of Theorems 1.10.4 and 1.10.5 is the following lemma, which follows from the Hoeffding inequality.

Lemma 1.10.6. The expectation

$$
\mathbb{E}\left[\mathcal{L}_{k+1}\left(s^{\prime}, a^{\prime}\right)\right]=\pi_{k}\left(s^{\prime}, a^{\prime}\right) \sum_{s \in \mathcal{S}} \sum_{a \in \mathcal{A}} \mu_{k}(s) P\left(s^{\prime} \mid s, a, \mathcal{L}_{k}\right) \pi_{k}(s, a)=\Gamma_{2}\left(\pi_{k}, \mathcal{L}_{k}\right) .
$$

In addition, we have that for any $t>0, s^{\prime} \in \mathcal{S}$ and $a^{\prime} \in \mathcal{A}$,

$$
\begin{equation*}
\mathbb{P}\left(\left|\mathcal{L}_{k+1}\left(s^{\prime}, a^{\prime}\right)-\mathbb{E}\left[\mathcal{L}_{k+1}\left(s^{\prime}, a^{\prime}\right)\right]\right| \geq t\right) \leq 2 \exp \left(-2 N t^{2}\right) \tag{1.10.9}
\end{equation*}
$$

The above lemma essentially states that the iterates $\mathcal{L}_{k+1}$ of Algorithms 5 and 6 are very close to the " $\tilde{\mathcal{L}}_{k+1}$ " obtained from the (strong) simulator $\mathcal{G}\left(s, \pi_{k}, \mathcal{L}_{k}\right)$ with $s \sim \mu_{k}$ following line 5 in Algorithm 2 and line 6 in Algorithm 4. This bridges the gap between the weak and the strong simulators. In particular, by noticing that

$$
\begin{aligned}
W_{1}\left(\mathcal{L}_{k+1}, \mathcal{L}^{\star}\right) & \leq W_{1}\left(\mathcal{L}_{k+1}, \Gamma_{2}\left(\pi_{k}, \mathcal{L}_{k}\right)\right)+W_{1}\left(\Gamma_{2}\left(\pi_{k}, \mathcal{L}_{k}\right), \Gamma_{2}\left(\Gamma_{1}\left(\mathcal{L}^{\star}\right), \mathcal{L}^{\star}\right)\right) \\
& =W_{1}\left(\mathcal{L}_{k+1}, \mathbb{E}\left[\mathcal{L}_{k+1}\right]\right)+W_{1}\left(\Gamma_{2}\left(\pi_{k}, \mathcal{L}_{k}\right), \Gamma_{2}\left(\Gamma_{1}\left(\mathcal{L}^{\star}\right), \mathcal{L}^{\star}\right)\right) \\
& \leq \frac{\operatorname{diam}(\mathcal{S}) \operatorname{diam}(\mathcal{A})|\mathcal{S} \| \mathcal{A}|}{2}\left\|\mathcal{L}_{k+1}-\mathbb{E}\left[\mathcal{L}_{k+1}\right]\right\|_{\infty}+W_{1}\left(\Gamma_{2}\left(\pi_{k}, \mathcal{L}_{k}\right), \Gamma_{2}\left(\Gamma_{1}\left(\mathcal{L}^{\star}\right), \mathcal{L}^{\star}\right)\right)
\end{aligned}
$$

one can bound the first term with high probability via (1.10.9). The second term is then bounded in exactly the same way as the proofs for Theorems 1.5.1 and 1.5.5, and hence we omit the details.

[^5]
### 1.10.4 Adaption of MF-Q

Note that MF-Q is designed for global states and coupled local actions, while in our equilibrium price example we have coupled local (private) states and decoupled local actions. To suit this setting, we adapt MF-Q by replacing the mean-field action term with the meanfield state term. In addition, our comparison is within the tabular setting, which is the setting theoretically analyzed in [167]. In this case, the $Q$ functions in MF-Q are also functions of only states and actions but not the mean-field terms, which are iteratively changing together with the $Q$-functions and claimed to be converging to a unique point. This is not the case for the function approximation setting (where $Q$-functions in MF-Q are indeed dependent on the mean-field terms and function approximation enables learning and generalizing such dependencies via experience replay). We leave it a future work to extend and compare our algorithms with MF-Q in the deep approximation setting.

## Chapter 2

## Logarithmic Regret for Episodic Continuous-Time Linear-Quadratic Reinforcement Learning over a Finite-Time Horizon

### 2.1 Introduction

Reinforcement learning (RL) for linear quadratic (LQ) control problems has been one of the most active areas for both the control and the reinforcement learning communities. Over the last few decades, significant progresses have been made in the discrete-time setting.

### 2.1.1 Discrete-time RL

In the area of adaptive control with unknown dynamics parameters, the goal is to find optimal stationary policy that stabilizes the unknown dynamics and minimizes the long term average cost ([88, 103]). For an infinite-time horizon LQ system, it has been shown that persistent excitation conditions [72] are critical to the parameter identification. Meanwhile, algorithms with asymptotic convergence in both the parameter estimation and the optimal control have been developed in [69], [99] and [30]: the first one assumes that costs only depend on state variables and the other two consider both state and control costs and use a cost-biased least-squared estimation method. See [55, 57] and references therein for recent developments of (randomised) adaptive control algorithms for LQ systems.

Following the seminal works of [14, 13] and [124], non-asymptotic regret bound analysis for RL algorithms has been one of the main topics, and has been developed for tabular Markov decision problems.

The non-asymptotic analysis of adaptive LQ problem by [1] utilizes the Optimism in the Face of Uncertainty principle to construct a sequence of improving confidence regions for
the unknown model parameters, and solves a non-convex constrained optimization problem for each confidence region; their algorithm achieves an $\mathcal{O}(\sqrt{T})$ regret bound, with $T$ being the number of time steps. To reduce the computational complexity and to avoid the non-convexity issue, [2] and [125] propose Thompson-sampling-based algorithms and derive $\mathcal{O}(\sqrt{T})$ regret bounds in the Bayesian setting; [45] proposes a robust adaptive control algorithm to solve a convex sub-problem in each step and achieves an $\mathcal{O}\left(T^{2 / 3}\right)$ regret bound. The gap between these regret bounds is removed by [112] and [40] via two different approaches for the same $\mathcal{O}(\sqrt{T})$ frequentist regret bound. Later, [142] establishes a lower bound on the regret of order $\mathcal{O}\left(\sqrt{d_{u}^{2} d_{x} T}\right)$, where $d_{u}$ and $d_{x}$ are the dimensions of the actions and the states, and shows that a simple variant of certainty equivalent control matches the lower bound in both $T$ and the dimensions. Similar regret bounds have also been established under different settings and assumptions, such as [35] in the adversarial setting and [101] without a stabilizing controller at the early stages of agent-environment interaction.

All the analyses are in discrete-time with an infinite time horizon. In all these problems, adaptive control algorithms are shown to achieve logarithmic regret bounds when additional information regarding the parameters of the system (often referred to as identifiability conditions) is available. Indeed, $[56,58]$ prove that certainty equivalent adaptive regulator achieves logarithmic regret bounds if the system parameter satisfies certain sparsity or low-rankness conditions. [34] establishes logarithmic regret bounds when either the state transition matrix is unknown, or when the state-action transition matrix is unknown and the optimal policy is non-degenerate. In partially observable linear dynamical systems, which takes linear-quadratic Gaussian problem as a special case, [102] proposes an algorithm with a logarithmic regret bound, under the assumption that one has access to a set in which all controllers persistently excite the system to approximate the optimal control. Logarithmic regret bounds in the adversarial setting with known dynamics parameters have been established in [5, 59].

### 2.1.2 Continuous-time RL.

Most real-world control systems, such as those in aerospace, automotive industry and robotics, are naturally continuous-time dynamical systems.

So are their related physical tasks, such as inverted pendulum problems, cart-pole balancing problems, and legged robot problems. Continuous-time finite-time horizon LQ control problems can be found in portfolio optimization [160], algorithmic trading [33], production management of exhaustible resources [71], and biological movement systems [16].

Analysis for continuous-time LQ-RL and general RL problems, however, is fairly limited. The primary approach is to develop learning algorithms after discretizing both the time and the space spaces, and establish the convergence as discretization parameters tend to zero. For instance, [120] proposes a policy gradient algorithm and shows the convergence of the policy gradient estimate to the true gradient. [121, 119] design learning algorithms by discretizing Bellman equations of the underlying control problems and prove the asymptotic convergence of their algorithms. For the LQ system, attentions have been mostly on algorithms designs,
including the integral reinforcement learning algorithm in [118], and the policy iteration algorithm in [136]. Yet, very little is known regarding the convergence rate or the regret bound of all these algorithms. Indeed, despite the natural analytical connection between LQ control and RL, the best known theoretical work for continuous-time LQ-RL is still due to [50], where an asymptotically sublinear regret for an ergodic model has been derived via a weighted least-squares-based estimation approach. Nevertheless, the exact order of the regret bound has not been studied.

Issues and challenges from non-asymptotic analysis. It is insufficient and improper to rely solely on the analysis and algorithms for the discrete-time RL to solve the continuoustime problems. There is a mismatch between the algorithms timescale for the former and the underlying systems timescale for the latter. When designing algorithms that make observations and take actions at discrete time points, it is important to take the model mismatch into consideration. For instance, the empirical studies in [149] suggest that vanilla $Q$-learning methods exhibit degraded performance as the time stepsize decreases, while a proper scaling of learning rates with stepsize leads to more robust performance.

The questions are therefore: A) How to quantify the precise impacts of the observation stepsize and action stepsize on algorithm performance? B) How to derive non-asymptotic regret analysis for learning algorithms in continuous-time LQ-RL (or general RL) system, analogous to the discrete-time LQ-RL counterpart?

There are technical reasons behind the limited theoretical progress in the continuous-time domain for RL, including LQ-RL. In addition to the known difficulty for analyzing stochastic control problems, the learning component compounds the problem complexity and poses new challenges.

For instance, the counterpart in the continuous-time problem to the algebraic equations in [112] for the discrete-time version is the regularity and stability of the continuous-time Riccati equation and the regularity of feedback controls. While Riccati equation and its robustness and existence and uniqueness of optimal controls have been well studied in the control literature, regularity of feedback controls with respect to underlying models is completely new for control theory and crucial for algorithm design and its robustness analysis. Moreover, deriving the exact order of the regret bound requires developing new and different techniques than those used for the asymptotic regret analysis in [50].

Our work and contributions. This paper studies finite-time horizon continuous-time LQ-RL problems in an episodic setting.

- It first proposes a greedy least-squares algorithm based on continuous-time observations and controls. At each iteration, the algorithm estimates the unknown parameters by a regularized least-squares estimator based on observed trajectories, then designs linear feedback controls via the Riccati differential equation for the estimated model. It identifies conditions under which the unknown state transition matrix and state-action transition matrix are uniquely identifiable under the optimal policies. (Remark 2.2.1 and Proposition
2.2.1). By exploiting the identifiability of coefficients, this continuous-time least-squares algorithm is shown to have a logarithmic regret of the magnitude $\mathcal{O}((\ln M)(\ln \ln M))$, with $M$ being the number of learning episodes (Theorem 2.2.2). To the best of our knowledge, this is the first non-asymptotic logarithmic regret bound for continuous-time LQ-RL problems with unknown state and control coefficients.
- It then proposes a practically implementable least-squares algorithm based on discretetime observations and controls. At each iteration, the algorithm estimates the unknown parameters by observing continuous-time trajectories at discrete time points, then designs a piecewise constant linear feedback control via Riccati difference equations for an associated discrete-time LQ-RL problem. It shows that the regret of the discrete-time least-squares algorithm is of the magnitude $\mathcal{O}\left((\ln M)(\ln \ln M)+\sum_{\ell=0}^{\ln M} 2^{\ell} \tau_{\ell}^{2}\right)$, where $\tau_{\ell}$ is the time stepsize used in the $(\ell+1)$-th update of model parameters (Theorem 2.2.3).
Our analysis shows that scaling the regularization parameter of the discrete-time leastsquares estimator with respect to time stepsize is critical for a robust performance of the algorithm in different timescales (Remark 2.2.3). To the best of our knowledge, this is the first discrete-time algorithm with rigorous regret bound for continuous-time LQ-RL problems.

Different from the least-squares algorithms for the ergodic LQ problems (see e.g., [50, 112]), our continuous-time least-squares algorithm constructs feedback controls via Riccati differential equations instead of the algebraic equations in [112]. Here, the regularity and stability of the continuous-time Riccati equation is analyzed in order to establish the robustness of feedback controls.

Moreover, our analysis for the estimation error exploits extensively the sub-exponential tail behavior of the least-squares estimators. This probabilistic approach differs from the asymptotic sublinear regret analysis in [50]; it establishes the exact order of the logarithmic regret bound by the concentration inequality for the error bound.

In addition, our analysis also exploits an important self-exploration property of finite-time horizon continuous-time LQ-RL problems, for which the time-dependent optimal feedback matrices ensure that the optimal state and control processes span the entire parameter space. This property allows us to design exploration-free learning algorithms with logarithmic regret bounds. Furthermore, we provide explicit conditions on models that guarantees the successful identification of the unknown parameters with optimal feedback policies. This is in contrast to the identification conditions for logarithmic regret bounds in discrete-time infinite-time-horizon LQ problems. Our conditions apply to arbitrary finite time-horizon problems, without imposing sparsity or low-rankness conditions on system parameters as in $[56,58]$ or requiring these parameters to be partially known to the controller as in $[34,59]$.

Finally, our analysis provides the precise parameter estimation error in terms of the sample size and time stepsize, and quantifies the performance gap between applying a piecewiseconstant policy from an incorrect model and applying the optimal policy. The misspecification error scales linearly with respect to the stepsize, and the performance gap depends
quadratically with respect to the time stepsize and the magnitude of parameter perturbations. Our analysis is based on the first-order convergence of Riccati difference equations and a uniform sub-exponential tail bound of discrete-time least-squares estimators.

Notation. For each $n \in \mathbb{N}$, we denote by $I=I_{n}$ the $n \times n$ identity matrix, and by $\mathbb{S}_{0}^{n}$ (resp. $\mathbb{S}_{+}^{n}$ ) the space of symmetric positive semidefinite (resp. definite) matrices. We denote by $|\cdot|$ the Euclidean norm of a given Euclidean space, by $\|\cdot\|_{2}$ the matrix norm induced by Euclidean norms, and by $A^{\top}$ and $\operatorname{tr}(A)$ the transpose and trace of a matrix $A$, respectively. For each $T>0$, filtered probability space $\left(\Omega, \mathcal{F}, \mathbb{F}=\left\{\mathcal{F}_{t}\right\}_{t \in[0, T]}, \mathbb{P}\right)$ satisfying the usual condition and Euclidean space $(E,|\cdot|)$, we introduce the following spaces:

- $C([0, T] ; E)$ is the space of continuous functions $\phi:[0, T] \rightarrow E$ satisfying $\|\phi\|_{C([0, T] ; E)}=$ $\sup _{t \in[0, T]}\left|\phi_{t}\right|<\infty ;$
- $C^{1}([0, T] ; E)$ is the space of continuously differentiable functions $\phi:[0, T] \rightarrow E$ satisfying $\|\phi\|_{C^{1}([0, T] ; E)}=\sup _{t \in[0, T]}\left(\left|\phi_{t}\right|+\left|\phi_{t}^{\prime}\right|\right)<\infty ;$
- $\mathcal{S}^{2}(E)$ is the space of $E$-valued $\mathbb{F}$-progressively measurable càdlàg processes $X: \Omega \times$ $[0, T] \rightarrow E$ satisfying $\|X\|_{\mathcal{S}^{2}(E)}=\mathbb{E}\left[\sup _{t \in[0, T]}\left|X_{t}\right|^{2}\right]^{1 / 2}<\infty$;
- $\mathcal{H}^{2}(E)$ is the space of $E$-valued $\mathbb{F}$-progressively measurable processes $X: \Omega \times[0, T] \rightarrow E$ satisfying $\|X\|_{\mathcal{H}^{2}(E)}=\mathbb{E}\left[\int_{0}^{T}\left|X_{t}\right|^{2} \mathrm{~d} t\right]^{1 / 2}<\infty$.
For notation simplicity, we denote by $C \in[0, \infty)$ a generic constant, which depends only on the constants appearing in the assumptions and may take a different value at each occurrence.


### 2.2 Problem formulation and main results

### 2.2.1 Linear-quadratic reinforcement learning problem

In this section, we consider the linear-quadratic reinforcement learning (LQ-RL) problem, where the drift coefficient of the state dynamics is unknown to the controller.

More precisely, let $T \in(0, \infty)$ be a given terminal time, $W$ be an $n$-dimensional standard Brownian motion defined on a complete probability space $(\Omega, \mathcal{F}, \mathbb{P})$, and $\mathbb{F}=\left(\mathcal{F}_{t}\right)_{t \in[0, T]}$ be the filtration generated by $W$ augmented by the $\mathbb{P}$-null sets. Let $x_{0} \in \mathbb{R}^{n}$ be a given initial state and $\left(A^{\star}, B^{\star}\right) \in \mathbb{R}^{n \times n} \times \mathbb{R}^{n \times d}$ be fixed but unknown matrices, consider the following problem:

$$
\begin{equation*}
\inf _{U \in \mathcal{H}^{2}\left(\mathbb{R}^{d}\right)} J^{\theta^{\star}}(U), \quad \text { with } \quad J^{\theta^{\star}}(U)=\mathbb{E}\left[\int_{0}^{T}\left(\left(X_{t}^{\theta^{\star}, U}\right)^{\top} Q X_{t}^{\theta^{\star}, U}+\left(U_{t}\right)^{\top} R U_{t}\right) \mathrm{d} t\right], \tag{2.2.1}
\end{equation*}
$$

where for each $U \in \mathcal{H}^{2}\left(\mathbb{R}^{d}\right)$, the process $X^{\theta^{\star}, U} \in \mathcal{S}^{2}\left(\mathbb{R}^{n}\right)$ satisfies the following controlled dynamics associated with the parameter $\theta^{\star}=\left(A^{\star}, B^{\star}\right)^{\top}$ :

$$
\begin{equation*}
\mathrm{d} X_{t}=\left(A^{\star} X_{t}+B^{\star} U_{t}\right) \mathrm{d} t+\mathrm{d} W_{t}, \quad t \in[0, T] ; \quad X_{0}=x_{0}, \tag{2.2.2}
\end{equation*}
$$

with given matrices $Q \in \mathbb{S}_{0}^{n}$ and $R \in \mathbb{S}_{+}^{d}$. Note that we assume the loss functional (2.2.1) only involves a time homogeneous running cost to allow a direct comparison with infinitetime horizon RL problems (see e.g., [51]), but similar analysis can be performed if the cost
functions are time inhomogeneous, a terminal cost is included, or the Brownian motion $W$ in (2.2.2) is scaled by an known nonsingular diffusion matrix.

If the parameter $\theta^{\star}=\left(A^{\star}, B^{\star}\right)^{\top}$ are known to the controller, then (2.2.1)-(2.2.2) reduces to the classical LQ control problems. In this case, it is well known that (see e.g., [169] and the references therein), the optimal control $U^{\theta^{*}}$ of $(2.2 .1)-(2.2 .2)$ is given in a feedback form by

$$
\begin{equation*}
U_{t}^{\theta^{\star}}=\psi^{\theta^{\star}}\left(t, X_{t}^{\theta^{\star}}\right), \quad \text { with } \psi^{\theta^{\theta^{*}}}(t, x)=K_{t}^{\theta^{\star}} x, \forall(t, x) \in[0, T] \times \mathbb{R}^{n} \tag{2.2.3}
\end{equation*}
$$

where $K_{t}^{\theta^{\star}}=-R^{-1}\left(B^{\star}\right)^{\top} P_{t}^{\theta^{\star}}$ for all $t \in[0, T],\left(P_{t}^{\theta^{*}}\right)_{t \in[0, T]}$ solves the Riccati equation

$$
\begin{equation*}
\frac{\mathrm{d}}{\mathrm{~d} t} P_{t}+\left(A^{\star}\right)^{\top} P_{t}+P_{t} A^{\star}-P_{t}\left(B^{\star} R^{-1}\left(B^{\star}\right)^{\top}\right) P_{t}+Q=0, \quad t \in[0, T] ; \quad P_{T}=0 \tag{2.2.4}
\end{equation*}
$$

and $X^{\theta^{\star}}$ is the state process governed by the following dynamics:

$$
\begin{equation*}
\mathrm{d} X_{t}=\left(A^{\star} X_{t}+B^{\star} K_{t}^{\theta^{\star}} X_{t}\right) \mathrm{d} t+\mathrm{d} W_{t}, \quad t \in[0, T] ; \quad X_{0}=x_{0} \tag{2.2.5}
\end{equation*}
$$

To solve the LQ-RL problem (2.2.1)-(2.2.2) with unknown $\theta^{\star}$, the controller searches for the optimal control while simultaneously learning the system, i.e., the matrices $A^{\star}, B^{\star}$. In an episodic (also known as reset or restart) learning framework, the controller improves her knowledge of the underlying dynamics $X_{t}$ through successive learning episodes, in order to find a control that is close to the optimal one.

Mathematically, it goes as follows. Let $M \in \mathbb{N}$ be the total number of learning episodes. In the $i$-th learning episode, $i=1, \ldots, M$, a feedback control $\psi^{i}$ is exercised, and the state process $X^{\psi^{i}}$ evolves according to the dynamics (2.2.2) controlled by the policy $\psi^{i}$ :

$$
\begin{equation*}
\mathrm{d} X_{t}=\left(A^{\star} X_{t}+B^{\star} \psi^{i}\left(t, X_{t}\right)\right) \mathrm{d} t+\mathrm{d} W_{t}^{i}, \quad t \in[0, T] ; \quad X_{0}=x_{0} \tag{2.2.6}
\end{equation*}
$$

Here $W^{i}, i=1,2, \ldots, M$ are independent $n$-dimensional Brownian motions defined on the same probability space $(\Omega, \mathcal{F}, \mathbb{P})$. The (expected) cost of learning in the $i$-th episode is then given by

$$
\begin{equation*}
J^{\theta^{\star}}\left(U^{\psi^{i}}\right)=\mathbb{E}\left[\int_{0}^{T}\left(\left(X_{t}^{\psi^{i}}\right)^{\top} Q X_{t}^{\psi^{i}}+\left(U_{t}^{\psi^{i}}\right)^{\top} R U_{t}^{\psi^{i}}\right) \mathrm{d} t\right], \text { with } U_{t}^{\psi^{i}}:=\psi^{i}\left(t, X_{t}^{\psi^{i}}\right), t \in[0, T] \tag{2.2.7}
\end{equation*}
$$

and the (expected) regret of learning up to $M \in \mathbb{N}$ episodes (with the sequence of controls $\left.\left(U^{\psi^{i}}\right)_{i=1}^{M}\right)$ is defined as follows:

$$
\begin{equation*}
R(M)=\sum_{i=1}^{M}\left(J^{\theta^{\star}}\left(U^{\psi^{i}}\right)-J^{\theta^{\star}}\left(U^{\theta^{\star}}\right)\right) \tag{2.2.8}
\end{equation*}
$$

where $J^{\theta^{\star}}\left(U^{\theta^{\star}}\right)$ is the optimal cost of (2.2.1)-(2.2.2) when $A^{\star}, B^{\star}$ are known. Intuitively, the regret characterizes the cumulative loss from taking sub-optimal policies in all episodes.

In the following, we shall propose several least-squares-based learning algorithms to solve (2.2.1)-(2.2.2), and prove that they achieve logarithmic regrets if $\theta^{\star}$ is identifiable (see Remark 2.2.1 for details).

### 2.2.2 Continuous-time least-squares algorithm and its regret bound

In this section, we consider a continuous-time least-squares algorithm, which chooses the optimal feedback control based on the current estimation of the parameter, and updates the parameter estimation based on the whole trajectories of the state dynamics.

More precisely, let $\theta=(A, B)^{\top} \in \mathbb{R}^{(n+d) \times n}$ be the current estimate of the unknown parameter $\theta^{\star}$, then the controller would exercise the optimal feedback control $\psi^{\theta}$ for (2.2.1)(2.2.2) with $\theta^{\star}$ replaced by $\theta$, i.e.,

$$
\begin{equation*}
\psi^{\theta}(t, x)=K_{t}^{\theta} x, \quad K_{t}^{\theta}:=-R^{-1} B^{\top} P_{t}^{\theta}, \quad \forall(t, x) \in[0, T] \times \mathbb{R}^{n} \tag{2.2.9}
\end{equation*}
$$

where $P^{\theta}$ satisfies the Riccati equation (3.5.1) with $\theta^{\star}$ replaced by $\theta$ :

$$
\begin{equation*}
\frac{\mathrm{d}}{\mathrm{~d} t} P_{t}+A^{\top} P_{t}+P_{t} A-P_{t}\left(B R^{-1} B^{\top}\right) P_{t}+Q=0, \quad t \in[0, T] ; \quad P_{T}=0 \tag{2.2.10}
\end{equation*}
$$

This leads to the state process $X^{\psi^{\theta}}$ satisfying (cf. (2.2.6)):

$$
\begin{equation*}
\mathrm{d} X_{t}=\left(A^{\star} X_{t}+B^{\star} \psi^{\theta}\left(t, X_{t}\right)\right) \mathrm{d} t+\mathrm{d} W_{t}, \quad t \in[0, T] ; \quad X_{0}=x_{0} \tag{2.2.11}
\end{equation*}
$$

We proceed to derive an $\ell_{2}$-regularized least-squares estimation for $\theta^{\star}$ based on sampled trajectories of $X^{\psi^{\theta}}$. Observing from (2.2.11) that

$$
Z_{t}^{\psi^{\theta}}\left(\mathrm{d} X_{t}^{\psi^{\theta}}\right)^{\top}=Z_{t}^{\psi^{\theta}}\left(Z_{t}^{\psi^{\theta}}\right)^{\top} \theta^{\star} \mathrm{d} t+Z_{t}^{\psi^{\theta}}\left(\mathrm{d} W_{t}\right)^{\top}, \quad \text { with } Z_{t}^{\psi^{\theta}}=\binom{X^{\psi^{\theta}}}{\psi^{\theta}\left(t, X_{t}^{\psi^{\theta}}\right)} \text { for all } t \in[0, T] .
$$

Hence the martingale property of the Itô integral implies that

$$
\begin{equation*}
\theta^{\star}=\left(\mathbb{E}\left[\int_{0}^{T} Z_{t}^{\psi^{\theta}}\left(Z_{t}^{\psi^{\theta}}\right)^{\top} \mathrm{d} t\right]\right)^{-1} \mathbb{E}\left[\int_{0}^{T} Z_{t}^{\psi^{\theta}}\left(\mathrm{d} X_{t}^{\psi^{\theta}}\right)^{\top}\right] \tag{2.2.12}
\end{equation*}
$$

provided that $\mathbb{E}\left[\int_{0}^{T} Z_{t}^{\psi^{\theta}}\left(Z_{t}^{\psi^{\theta}}\right)^{\top} \mathrm{d} t\right]$ is invertible. This suggests a practical rule to improve one's estimate $\theta$ for the true parameter $\theta^{\star}$, by replacing the expectations in (2.2.12) with empirical averages over independent realizations. More precisely, let $m \in \mathbb{N}$ and

$$
\left(X_{t}^{\psi^{\theta}, i}, \psi^{\theta}\left(t, X_{t}^{\psi^{\theta}, i}\right)\right)_{t \in[0, T]}, \quad i=1, \ldots, m
$$

be trajectories of $m$ independent realizations of the state and control processes, we shall update the estimate $\theta$ by the following rule, inspired by (2.2.12):

$$
\begin{equation*}
\theta \longleftarrow\left(\frac{1}{m} \sum_{i=1}^{m} \int_{0}^{T} Z_{t}^{\psi^{\theta}, i}\left(Z_{t}^{\psi^{\theta}, i}\right)^{\top} \mathrm{d} t+\frac{1}{m} I\right)^{-1}\left(\frac{1}{m} \sum_{i=1}^{m} \int_{0}^{T} Z_{t}^{\psi^{\theta}, i}\left(\mathrm{~d} X_{t}^{\psi^{\theta}, i}\right)^{\top}\right), \tag{2.2.13}
\end{equation*}
$$

where $Z_{t}^{\psi^{\theta}, i}:=\binom{X_{t}^{\psi^{\theta}, i}}{\psi^{\theta}\left(t, X_{t}^{\psi^{\theta}, i}\right)}$ for all $t \in[0, T]$ and $i=1, \ldots, m$, and $I$ is the $(n+d) \times(n+d)$ identity matrix.

The regularization term $\frac{1}{m} I$ in (2.2.13) guarantees the required matrix inverse and vanishes as $m \rightarrow \infty$. The estimator (2.2.13) can be equivalently expressed as an $\ell_{2}$-regularized least-squares estimator, as pointed out in [51] for the ergodic LQ-RL problem.

We summarize the continuous-time least-squares algorithm as follows.

```
Algorithm 8 Continuous-time least-squares algorithm
    Input: Choose an initial estimation \(\theta_{0}\) of \(\theta^{\star}\) and numbers of learning episodes
    \(\left\{m_{\ell}\right\}_{\ell \in \mathbb{N} \cup\{0\}}\).
    for \(\ell=0,1, \cdots\) do
        Obtain the feedback control \(\psi^{\theta_{\ell}}\) as (2.2.9) with \(\theta=\theta_{\ell}\).
        Execute the feedback control \(\psi^{\theta_{\ell}}\) for \(m_{\ell}\) independent episodes, and collect the trajec-
        tory data \(\left(X_{t}^{\psi^{\theta_{\ell}, i}}, \psi^{\theta_{\ell}}\left(t, X_{t}^{\psi^{\theta_{\ell}, i}}\right)\right)_{t \in[0, T]}, i=1, \ldots, m_{\ell}\).
        Obtain an updated estimation \(\theta_{\ell+1}\) by using (2.2.13) and the \(m_{\ell}\) trajectories collected
        above.
    end for
```

Note that Algorithm 8 operates in cycles, with $m_{\ell}$ the number of episodes in the $\ell$-th cycle. Hence, the regret of learning up to $M$ episodes (cf. (2.2.8)) can be upper bounded by the accumulated regret at the end of the $L$-th cycle, where $L$ is the smallest integer such that $\sum_{\ell=0}^{L} m_{\ell} \geq M$.

In this section, we analyze the regret of Algorithm 8 based on the following assumptions of the learning problem (2.2.1)-(2.2.2).
H.1. (1) $T \in(0, \infty), n, d \in \mathbb{N}, x_{0} \in \mathbb{R}^{n}, A^{\star} \in \mathbb{R}^{n \times n}, B^{\star} \in \mathbb{R}^{n \times d}, Q \in \mathbb{S}_{0}^{n}$ and $R \in \mathbb{S}_{+}^{d}$.
(2) $\left\{v \in \mathbb{R}^{d} \mid\left(K_{t}^{\theta^{\star}}\right)^{\top} v=0, \forall t \in[0, T]\right\}=\{0\}$, with $K^{\theta^{\star}}$ defined in (2.2.3).

Before discussing the regret of Algorithm 8, we make the following remark of (H.1).
Remark 2.2.1 (Self-exploration of finite-time horizon RL problems). (H.1(1)) is the standard assumption for finite-time horizon LQ-RL problems (see e.g., [80]), except that H.1(1) allows $Q$ to be positive semidefinite, which is important for costs depending on partial states. (H.15) corresponds to the identifiability of the true parameter $\theta^{\star}$ by executing the optimal policy $K^{\theta^{\star}}$. In fact, as shown in Proposition 2.3.10, under (H.1(1)), (H.15) is equivalent to the following statement:
(2') if $u \in \mathbb{R}^{n}$ and $v \in \mathbb{R}^{d}$ satisfy $u^{\top} X^{\theta^{\star}}+v^{\top} U^{\theta^{\star}}=0$ for $\mathrm{d} \mathbb{P} \otimes \mathrm{d}$ t-almost everywhere in $\Omega \times[0, T]$, then $u=0$ and $v=0$, where $X^{\theta^{\star}}$ and $U^{\theta^{\star}}$ are the optimal state and control processes of (2.2.1)-(2.2.2) defined by (2.2.5) and (2.2.3), respectively,

Item (2') indicates an important self-exploration property of finite-time horizon continuoustime RL problems. In particular, the time-dependent optimal feedback matrix $K^{\theta^{\star}}$ and the non-degenerate noises guarantee the non-degeneracy of the space spanned by $X^{\theta^{\star}}$ and $U^{\theta^{\star}}$,
enabling learning the parameters sufficiently well. This self-exploration property is critical for our design of exploration-free learning algorithms for (2.2.1)-(2.2.2) with a logarithmic regret (see Theorems 2.2.2 and 2.2.3).

One can easily show that (H.15) holds if the optimal policy $\left(K^{\theta^{*}}\right)_{t \in[0, T]}$ is nondegenerate, i.e., $\sup _{t \in[0, T]} \lambda_{\min }\left(\left(K_{t}^{\theta^{*}}\right)\left(K_{t}^{\theta^{*}}\right)^{\top}\right)>0$. Similar nondegeneracy condition has been imposed in [34] for discrete-time ergodic LQ-RL problems. In particular, by assuming that the optimal stationary policy satisfies $\lambda_{\min }\left(K^{\star}\left(K^{\star}\right)^{\top}\right)>0$ (along with other controllablity conditions), they propose learning algorithms with a logarithmic regret, under the assumption that only the control coefficient $B^{\star}$ is unknown. In contrast, we allow both the state coefficient $A^{\star}$ and the control coefficient $B^{\star}$ to be unknown.

Moreover, the following proposition provides sufficient conditions of (H.15), which are special cases of Proposition 2.3.11.

Proposition 2.2.1. Let $n, d \in \mathbb{N}, Q \in \mathbb{S}_{0}^{n}$ and $R \in \mathbb{S}_{+}^{d}$.
(1) If $\left(B^{\star}\right)^{\top} Q B^{\star} \in \mathbb{S}_{+}^{d}$, then (H.15) holds for all $T>0$.
(2) Assume that the algebraic Riccati equation $\left(A^{\star}\right)^{\top} P+P A^{\star}-P\left(B^{\star} R^{-1}\left(B^{\star}\right)^{\top}\right) P+Q=0$ admits a unique maximal solution $P_{\infty}^{\star} \in \mathbb{S}_{+}^{n}$. Let $K_{\infty}^{\star}=-R^{-1}\left(B^{\star}\right)^{\top} P_{\infty}^{\star}$, and for each $T>0$, let $P^{\star,(T)} \in C\left([0, T] ; \mathbb{S}_{0}^{n}\right)$ be defined in (3.5.1). Assume that $\lim _{T \rightarrow \infty} P_{0}^{\star,(T)}=P_{\infty}^{\star}$ and $K_{\infty}^{\star}\left(K_{\infty}^{\star}\right)^{\top} \in \mathbb{S}_{+}^{d}$. Then there exists $T_{0}>0$, such that (H.15) holds for all $T \geq T_{0}$.

Proposition 2.2.1 provides two sets of conditions for (H.15) under two different scenarios: Item (1) applies to an arbitrary finite $T>0$, and Item (2) only applies to sufficiently large $T$. Item (2) assumes the asymptotic behavior of solutions to Riccati differential equations, which can be ensured by the stabilizability of the pair $\left(A^{\star}, B^{\star}\right)$ and detectability of the pair $\left(A^{\star}, Q^{1 / 2}\right)$ (see [24, Theorems 10.9 and 10.10]). Note that our subsequent analysis is based on (H.1), and does not require stabilizability assumptions.

Remark 2.2.2 (Stabilizability of $\left(A^{\star}, B^{\star}\right)$ and dependence on $\left.T\right)$. Since the $L Q-R L$ problem (2.2.1)-(2.2.2) is over the time horizon $[0, T]$ with a fixed $T<\infty$, in general one does not need additional conditions on $\left(A^{\star}, B^{\star}\right)$ for the well-definedness of (2.2.1)-(2.2.2). If $T=\infty$, then some stabilizability/controllability conditions of $\left(A^{\star}, B^{\star}\right)$ may be required for (2.2.1)-(2.2.2) to ensure a well-defined solution (see e.g., [46]). Under these conditions, different algorithms have been shown to achieve sub-linear regret with respect to the number of decision steps (see e.g., [112, 40]), and even logarithmic regrets provided that further identifiability assumptions are satisfied (see e.g., [56, 58, 34, 102]); see Section 2.1.1 for more details. For $T<\infty$, the regrets of learning algorithms for (2.2.1)-(2.2.2) in general depend exponentially on the time horizon $T$ (e.g., the constants $C_{0}, C^{\prime}$ in Theorem 2.2.2), as the moments of the optimal state process $X^{\theta^{\star}}$ and control process $U^{\theta^{\star}}$ may grow exponentially with respect to $T$. It would be interesting to quantify the precise dependence of the regret bounds on $T$. This would entail deriving precise a priori bounds of solutions to (2.2.10) and
estimating the norm $\left\|\left(\mathbb{E}\left[\int_{0}^{T} Z_{t}^{\theta^{\star}}\left(Z_{t}^{\theta^{\star}}\right)^{\top} \mathrm{d} t\right]\right)^{-1}\right\|_{2}$ in terms of $\left(A^{\star}, B^{\star}, Q, T\right)$, and is left for future research.

We are now ready to state the main result of this section, which shows that the regret of Algorithm 8 grows logarithmically with respect to the number of episodes.

Theorem 2.2.2. Suppose (H.1) holds and let $\theta_{0}=\left(A_{0}, B_{0}\right)^{\top} \in \mathbb{R}^{(n+d) \times d}$ such that (H.15) holds with $\theta_{0}$. Then there exists a constant $C_{0}>0$ such that for all $C \geq C_{0}$, and $\delta \in\left(0, \frac{3}{\pi^{2}}\right)$, if one sets $m_{0}=C(-\ln \delta)$ and $m_{\ell}=2^{\ell} m_{0}$ for all $\ell \in \mathbb{N}$, then with probability at least $1-\frac{\pi^{2} \delta}{3}$, the regret of Algorithm 8 given by (2.2.8) satisfies

$$
R(M) \leq C^{\prime}((\ln M)(\ln \ln M)+(-\ln \delta)(\ln M)), \quad \forall M \in \mathbb{N},
$$

where $C^{\prime}$ is a constant independent of $M$ and $\delta$.
To simplify the presentation, we analyze the performance of Algorithm 8 by assuming the number of learning episodes $\left\{m_{\ell}\right\}_{\ell}$ is doubled between two successive updates of the estimation of $\theta^{\star}$. Similar regret results can be established for Algorithm 8 with different choices of $\left\{m_{\ell}\right\}_{\ell}$. Under this specific choice of $\left\{m_{\ell}\right\}_{\ell}$, for any $M \in \mathbb{N}$, Algorithm 8 splits $M$ episodes into $L=\left\lceil\log _{2}\left(\frac{M}{m_{0}}+1\right)\right\rceil-1$ cycles, where the $\ell$-th cycle, $\ell=0,1, \ldots, L-1$, contains $m_{\ell}$ episodes, and the remaining $M-\sum_{\ell=0}^{L-1} m_{\ell}$ episodes are in the last cycle.

Sketched proof of Theorem 2.2.2. We outline the key steps of the proof of Theorem 2.2.2, and present the detailed arguments to Section 2.3.3. By exploiting the regularity and robustness of solutions to (2.2.10), we prove that the performance gap $J^{\theta^{\star}}\left(U^{\psi^{\theta}}\right)-J^{\theta^{\star}}\left(U^{\theta^{\star}}\right)$ is of the magnitude $\mathcal{O}\left(\left|\theta-\theta^{\star}\right|^{2}\right.$ ), for all $a$-priori bounded $\theta$ (Proposition 2.3.8). We then establish a uniform sub-exponential property for the (deterministic and stochastic) integrals in (2.2.12), which along with (H.15) and Bernstein's inequality leads to the following estimate of the parameter estimation error: for all $\delta \in(0,1 / 2)$, all sufficiently large $m \in \mathbb{N}$, and all $\theta$ sufficiently close to $\theta^{\star}$,

$$
\begin{equation*}
\left|\hat{\theta}-\theta^{\star}\right| \leq \mathcal{O}\left(\sqrt{\frac{-\ln \delta}{m}}+\frac{-\ln \delta}{m}+\frac{(-\ln \delta)^{2}}{m^{2}}\right), \quad \text { with probability } 1-2 \delta, \tag{2.2.14}
\end{equation*}
$$

where $\hat{\theta}$ is generated by (2.2.13) with $\psi^{\theta}$ (Proposition 3.3.12). Then for each $\delta>0$, applying (2.2.14) with $\delta_{\ell}=\delta /(\ell+1)^{2}$ for all $\ell \in \mathbb{N} \cup\{0\}$ shows that with probability $1-2 \sum_{\ell=0}^{\infty} \delta_{\ell}=$ $1-\frac{\pi^{2} \delta}{3}$,

$$
\begin{equation*}
\left|\hat{\theta}_{\ell+1}-\theta^{\star}\right|^{2} \lesssim \frac{-\ln \delta_{\ell}}{m_{\ell}}+\frac{\left(-\ln \delta_{\ell}\right)^{2}}{m_{\ell}^{2}}+\frac{\left(-\ln \delta_{\ell}\right)^{4}}{m_{\ell}^{4}}, \quad \forall \ell \in \mathbb{N}, \tag{2.2.15}
\end{equation*}
$$

where $\lesssim$ means the inequality holds with a multiplicative constant independent of $\delta$ and $\ell$. By the quadratic performance gap and the choice of $\left\{m_{\ell}\right\}_{\ell}$, the regret of Algorithm

8 up to the $M$-th episode can be bounded by the regret at the end of $L$-th cycle with $L=\left\lceil\log _{2}\left(\frac{M}{m_{0}}+1\right)\right\rceil-1$ :

$$
\begin{equation*}
R(M) \lesssim \sum_{\ell=0}^{L} m_{\ell}\left|\theta_{\ell}-\theta^{\star}\right|^{2} \lesssim \sum_{\ell=0}^{L}\left(-\ln \delta_{\ell}\right)\left(1+\frac{-\ln \delta_{\ell}}{m_{\ell}}+\frac{\left(-\ln \delta_{\ell}\right)^{3}}{m_{\ell}^{3}}\right) \tag{2.2.16}
\end{equation*}
$$

Observe that the choices of $\left\{\delta_{\ell}\right\}_{\ell}$ and $\left\{m_{\ell}\right\}_{\ell}$ ensure that $\sup _{\delta \in\left(0, \frac{3}{\pi^{2}}\right), \ell \in \mathbb{N}} \frac{-\ln \delta_{\ell}}{m_{\ell}}<\infty$. Hence, the right-hand side of (2.2.16) is of the magnitude $\mathcal{O}\left(\sum_{\ell=0}^{L}\left(-\ln \delta_{\ell}\right)\right)$, which along with the choices of $\delta_{\ell}$ and $L$ leads to the desired regret bound; see the end of Section 2.3.3 for more details.

### 2.2.3 Discrete-time least-squares algorithm and its regret bound

Note that Algorithm 8 in Section 2.2.2 requires executing feedback controls and observing corresponding state trajectories continuously. A common practice to solve continuous-time RL problems is by assuming that at each learning episode the dynamics only evolves in discrete time, and then estimate parameters according to discrete-time RL algorithms (see e.g., $[121,119,120,149])$. As the true dynamics evolves continuously, it is necessary to quantify the impact of reaction stepsize on the algorithm performance.

In this section, we analyze the performance of the above procedure for solving (2.2.1)(2.2.2). We adapt regularized least-squares algorithms for discrete-time LQ problems to the present setting, and establish their regret bounds in terms of the discretization stepsize. Our analysis shows that a proper scaling of the regularization term in the least-squares estimation in terms of stepsize is critical for a robust performance with respect to different timescales.

More precisely, for a given cycle (i.e., the index $\ell$ in Algorithm 8), let $\theta=(A, B)^{\top} \in$ $\mathbb{R}^{(n+d) \times n}$ be the current estimate of $\theta^{\star}$ in (2.2.2), and let $\left\{t_{i}\right\}_{i=0}^{N}, N \in \mathbb{N}$, be a uniform partition of $[0, T]$ with stepsize $\tau=T / N$. We then assume that (2.2.1)-(2.2.2) is piecewise constant between any two grid points $\left\{t_{i}\right\}_{i=0}^{N}$, choose actions and make observations every $\tau$, and update the estimated parameter based on these observations. To this end, we consider the following discrete-time LQ control problem with parameter $\theta$ :

$$
\begin{equation*}
\inf _{U \in \mathcal{H}_{N}^{2}\left(\mathbb{R}^{d}\right)} J_{N}(U), \quad \text { with } J_{N}(U)=\mathbb{E}\left[\sum_{i=0}^{N-1}\left(\left(X_{t_{i}}^{U, \tau}\right)^{\top} Q X_{t_{i}}^{U, \tau}+U_{t_{i}}^{\top} R U_{t_{i}}\right) \tau\right] \tag{2.2.17}
\end{equation*}
$$

where $\mathcal{H}_{N}^{2}\left(\mathbb{R}^{d}\right)=\left\{U \in \mathcal{H}^{2}\left(\mathbb{R}^{d}\right) \mid U_{t}=U_{t_{i}}, t \in\left[t_{i}, t_{i+1}\right), i=0, \ldots, N-1\right\}$, and $\left(X_{t_{i}}^{U, \tau}\right)_{i=0}^{N-1}$ are defined by

$$
\begin{equation*}
X_{t_{i+1}}^{U, \tau}-X_{t_{i}}^{U, \tau}=\left(A X_{t_{i}}^{U, \tau}+B U_{t_{i}}\right) \tau+W_{t_{i+1}}-W_{t_{i}}, \quad i=0, \ldots N-1 ; \quad X_{0}^{U, \tau}=x_{0} \tag{2.2.18}
\end{equation*}
$$

Note that for simplicity, our strategy is constructed by assuming a discrete-time dynamics arising from an Euler discretization of (2.2.2) (with the estimated parameter $\theta$ ); similar analysis can be performed with a high-order approximation of (2.2.1)-(2.2.2).

It is well-known that (see e.g., [24]), the optimal control of (2.2.17)-(2.2.18) is given by the following feedback form:

$$
\begin{equation*}
U_{t}=\psi^{\theta, \tau}\left(t, X_{t}^{U, \tau}\right), \quad \text { with } \psi^{\theta, \tau}(t, x)=K_{t}^{\theta, \tau} x, \quad \forall(t, x) \in[0, T) \times \mathbb{R}^{n} \tag{2.2.19}
\end{equation*}
$$

where $K^{\theta, \tau}:[0, T) \rightarrow \mathbb{R}^{d \times n}$ is the piecewise constant function (with stepsize $\tau=T / N$ ) defined by

$$
\begin{align*}
P_{t_{i-1}}^{\theta, \tau}= & \tau Q+(I+\tau A)^{\top} P_{t_{i}}^{\theta, \tau}(I+\tau A)-(I+\tau A)^{\top} P_{t_{i}}^{\theta, \tau} \tau B\left(R+\tau B^{\top} P_{t_{i}}^{\theta, \tau} B\right)^{-1} B^{\top} P_{t_{i}}^{\theta, \tau}(I+\tau A), \\
& \forall i=0, \ldots, N-1 ; \quad P_{T}^{\theta, \tau}=0 \\
K_{t}^{\theta, \tau}=- & \left(R+\tau B^{\top} P_{t_{i+1}}^{\theta, \tau} B\right)^{-1} B^{\top} P_{t_{i+1}}^{\theta, \tau}(I+\tau A), \quad t \in\left[t_{i}, t_{i+1}\right), i=0, \ldots, N-1 . \tag{2.2.20}
\end{align*}
$$

We then implement the piecewise constant strategy $\psi^{\theta, \tau}$ defined in (2.2.19) on the original system (2.2.2) for $m$ episodes, and update the estimated parameter $\theta$ by observing (2.2.2) with stepsize $\tau=T / N$. More precisely, let $X^{\psi^{\theta, \tau}} \in \mathcal{S}^{2}\left(\mathbb{R}^{n}\right)$ be the state process associated with $\psi^{\theta, \tau}$ :

$$
\begin{equation*}
\mathrm{d} X_{t}=\left(A^{\star} X_{t}+B^{\star} K_{t_{i}}^{\theta, \tau} X_{t}\right) \mathrm{d} t+\mathrm{d} W_{t}, \quad t \in\left[t_{i}, t_{i+1}\right], i=0, \ldots, N-1 ; \quad X_{0}=x_{0}, \tag{2.2.21}
\end{equation*}
$$

and $\left(X_{t}^{\psi^{\theta, \tau}, j}\right)_{t \in[0, T]}, j=1, \ldots, m, m \in \mathbb{N}$, be $m$ independent trajectories of $X^{\psi^{\theta, \tau}} \in \mathcal{S}^{2}\left(\mathbb{R}^{n}\right)$, we update the parameter $\theta$ according to the following discrete-time least-squares estimator:

$$
\begin{equation*}
\theta \leftarrow \underset{\theta \in \mathbb{R}^{(n+d) \times n}}{\arg \min } \sum_{j=1}^{m} \sum_{i=0}^{N-1}\left\|X_{t_{i+1}}^{\psi^{\theta, \tau}, j}-X_{t_{i}}^{\psi^{\theta, \tau}, j}-\tau \theta^{\top} Z_{t_{i}}^{\psi^{\theta, \tau}, j}\right\|_{2}^{2}+\tau \operatorname{tr}\left(\theta^{\top} \theta\right), \tag{2.2.22}
\end{equation*}
$$

with $Z_{t_{i}}^{\psi^{\theta, \tau}, j}:=\binom{X_{t_{i}}^{\psi^{\theta, \tau}, j}}{K_{t_{i}, \tau}^{\theta, \tau} X_{t_{i}}^{\psi^{\theta, \tau}, j}}$ for all $i, j$. The update (2.2.22) is consistent with the agent's assumption that the state evolves according to (2.2.18) between two grid points. Setting the derivative (with respect to $\theta$ ) of the right-hand side of (2.2.22) to zero leads to

$$
-\sum_{j=1}^{m} \sum_{i=0}^{N-1} \tau Z_{t_{i}}^{\psi^{\theta, \tau}, j}\left(\left(X_{t_{i+1}}^{\psi^{\theta, \tau}, j}-X_{t_{i}}^{\psi^{\theta, \tau}, j}\right)^{\top}-\tau\left(Z_{t_{i}}^{\psi^{\theta, \tau}, j}\right)^{\top} \theta\right)+\tau \theta=0 .
$$

Dividing both sides by $\tau / m$ and rearranging the terms give the following equivalent expression of the discrete-time least squares estimator (2.2.22):

$$
\begin{equation*}
\theta \longleftarrow\left(\frac{1}{m} \sum_{j=1}^{m} \sum_{i=0}^{N-1} Z_{t_{i}}^{\psi^{\theta, \tau}, j}\left(Z_{t_{i}}^{\psi^{\theta, \tau}, j}\right)^{\top} \tau+\frac{1}{m} I\right)^{-1}\left(\frac{1}{m} \sum_{j=1}^{m} \sum_{i=0}^{N-1} Z_{t_{i}}^{\psi^{\theta, \tau}, j}\left(X_{t_{i+1}}^{\psi^{\theta, \tau}, j}-X_{t_{i}}^{\psi^{\theta, \tau}, j}\right)^{\top}\right) . \tag{2.2.23}
\end{equation*}
$$

Remark 2.2.3 (Scaling hyper-parameters with timescales). In principle, when applying discrete-time RL algorithms in a continuous environment, it is critical to adopt a proper scaling of the hyper-parameters for a robust performance with respect to different timescales. Indeed, scaling the regularization term $\operatorname{tr}\left(\theta^{\top} \theta\right)$ in (2.2.22) with respect to the stepsize $\tau$ is essential for the robustness of (2.2.23) for all small stepsizes $\tau$. If one updates $\theta$ by minimizing the following $\ell_{2}$-regularized loss function with a given hyper-parameter $\alpha<1$ such that

$$
\begin{equation*}
\underset{\theta \in \mathbb{R}^{(n+d) \times n}}{\arg \min } \sum_{j=1}^{m} \sum_{i=0}^{N-1}\left\|X_{t_{i+1}}^{\psi^{\theta, \tau}, j}-X_{t_{i}}^{\psi^{\theta, \tau}, j}-\tau \theta^{\top} Z_{t_{i}}^{\psi^{\theta, \tau}, j}\right\|_{2}^{2}+\tau^{\alpha} \operatorname{tr}\left(\theta^{\top} \theta\right), \tag{2.2.24}
\end{equation*}
$$

then the corresponding discrete-time estimator is given by
$\theta^{\tau}:=\left(\frac{1}{m} \sum_{j=1}^{m} \sum_{i=0}^{N-1} Z_{t_{i}}^{\psi^{\theta, \tau}, j}\left(Z_{t_{i}}^{\psi^{\theta, \tau}, j}\right)^{\top} \tau+\frac{1}{\tau^{1-\alpha} m} I\right)^{-1}\left(\frac{1}{m} \sum_{j=1}^{m} \sum_{i=0}^{N-1} Z_{t_{i}}^{\psi^{\theta, \tau}, j}\left(X_{t_{i+1}}^{\psi^{\theta, \tau}, j}-X_{t_{i}}^{\psi^{\theta, \tau}, j}\right)^{\top}\right)$.
Observe that for any given $m \in \mathbb{N}$, the estimator $\theta^{\tau}$ degenerates to zero as the stepsize $\tau$ tends to zero. Hence, to ensure the viability of $\theta^{\tau}$ across different timescales, the number of episodes $m$ has to increase appropriately when $\tau$ tends to zero. In contrast, by choosing $\alpha=1$ in (2.2.24), (2.2.23) admits a continuous-time limit (2.2.13) as $\tau \rightarrow 0$, and leads to a learning algorithm in which the episode numbers and the time stepsize can be chosen independently (see Theorem 2.2.3).

We now summarize the discrete-time least-squares algorithm as follows.

```
Algorithm 9 Discrete-time least-squares algorithm
    Input: Choose an initial estimation \(\theta_{0}\) of \(\theta^{\star}\), numbers of learning episodes \(\left\{m_{\ell}\right\}_{\ell \in \mathbb{N} \cup\{0\}}\)
    and numbers of intervention points \(\left\{N_{\ell}\right\}_{\ell \in \mathbb{N} \cup\{0\}}\).
    for \(\ell=0,1, \cdots\) do
        Obtain the piecewise constant control \(\psi^{\theta_{\ell}, \tau_{\ell}}\) as (2.2.19) with \(\tau=T / N_{\ell}\) and \(\theta=\theta_{\ell}\).
        Execute the control \(\psi^{\theta_{\ell}, \tau_{\ell}}\) for \(m_{\ell}\) independent episodes, and collect the data \(X_{t_{i}}^{\psi^{\theta_{\ell}, \tau_{\ell}, j}}\),
        \(i=0, \ldots, N_{\ell}, j=1, \ldots, m_{\ell}\).
        Obtain an updated estimation \(\theta_{\ell+1}\) by using (2.2.23) and the data
        \(\left(X_{t_{i}}^{\psi_{\ell}, \tau_{\ell}, j}\right)_{i=0, \ldots, N_{\ell}, j=1, \ldots, m_{\ell}}\).
    end for
```

Again, as the $\ell$-th cycle of Algorithm 9 contains $m_{\ell}$ episodes, for each $M \in \mathbb{N}$, the regret of learning up to $M$ episodes (cf. (2.2.8)) can be upper bounded by the accumulated regret at the end of the $L$-th cycle, where $L$ is the smallest integer such that $\sum_{\ell=0}^{L} m_{\ell} \geq M$. The following theorem is an analogue of Theorem 2.2.2 for Algorithm 9.

Theorem 2.2.3. Suppose (H.1) holds and let $\theta_{0}=\left(A_{0}, B_{0}\right)^{\top} \in \mathbb{R}^{(n+d) \times d}$ such that (H.15) holds with $\theta_{0}$. Then there exists $C_{0}>0$ and $n_{0} \in \mathbb{N}$ such that for all $C \geq C_{0}$, and $\delta \in\left(0, \frac{3}{\pi^{2}}\right)$,
if one sets $m_{0}=C(-\ln \delta)$, $m_{\ell}=2^{\ell} m_{0}$ and $N_{\ell} \geq n_{0}$ for all $\ell \in \mathbb{N} \cup\{0\}$, then with probability at least $1-\frac{\pi^{2} \delta}{3}$, the regret of Algorithm 9 given by (2.2.8) satisfies

$$
\begin{equation*}
R(M) \leq C^{\prime}\left((\ln M)(\ln \ln M)+(-\ln \delta)(\ln M)+(-\ln \delta) \sum_{\ell=0}^{\ln M} 2^{\ell} N_{\ell}^{-2}\right), \quad \forall M \in \mathbb{N} \tag{2.2.25}
\end{equation*}
$$

where $C^{\prime}$ is a constant independent of $M, \delta$ and $\left(N_{\ell}\right)_{\ell \in \mathbb{N} \cup\{0\}}$.

Remark 2.2.4. Theorem 2.2.3 provides a general regret bound of Algorithm 9 with any time discretization steps $\left\{N_{\ell}\right\}_{\ell \geq 0}$, where $N_{\ell}$ is the number of intervention points in the $\ell$ th cycle. Compared with Algorithm 8, the regret of Algorithm 9 has an additional term $(-\ln \delta) \sum_{\ell=0}^{\ln M} 2^{\ell} N_{\ell}^{-2}$ : for each learning episode, one achieves a sub-optimal loss by adjusting her policy in the discrete time and also suffers from model misspecification error in parameter estimation from discrete-time observations. Specifically,

- if the time discretization step is fixed for all cycles, i.e., $N_{\ell}=T / \tau$ for all $\ell$, then the last term of (2.2.25) is of the magnitude:

$$
\mathcal{O}\left((-\ln \delta) \sum_{\ell=0}^{\ln M} 2^{\ell} N_{\ell}^{-2}\right)=\mathcal{O}\left((-\ln \delta) \tau^{2} \sum_{\ell=0}^{\ln M} 2^{\ell}\right)=\mathcal{O}\left((-\ln \delta) \tau^{2} M\right)
$$

and consequently Algorithm 9 achieves a sub-optimal linear regret;

- if the time discretization step of the $\ell$-th cycle increases exponentially in terms of $\ell$, e.g., $N_{\ell}=\sqrt{2}^{\ell} N_{0}$ for $\ell=1, \ldots, \ln M$, then the last term of (2.2.25) is of the magnitude:

$$
\mathcal{O}\left((-\ln \delta) \sum_{\ell=0}^{\ln M} 2^{\ell} N_{\ell}^{-2}\right)=\mathcal{O}\left((-\ln \delta) \sum_{\ell=0}^{\ln M} N_{0}^{-2}\right)=\mathcal{O}((-\ln \delta) \ln M)
$$

which guarantees that the regret of Algorithm 9 is still logarithmic in $M$.

Sketched proof of Theorem 2.2.3. We point out the main differences between the proofs of Theorems 2.2.2-2.2.3, and give the detailed proof of Theorem 2.2.3 in Section 2.3.4. Compared with Theorem 2.2.2, the essential challenges in proving Theorem 2.2.3 are to quantify the precise dependence of the performance gap and the parameter estimation error on the stepsize. To this end, we first prove a first-order convergence of (2.2.20) to (2.2.9) as the stepsize tends to zero. Then by exploiting the affine structure of (2.2.21), we establish the following quadratic performance gap for a piecewise constant policy $\psi^{\theta, \tau}$ (Proposition 2.3.12):

$$
\begin{equation*}
J^{\theta^{\star}}\left(U^{\psi^{\theta, \tau}}\right)-J^{\theta^{\star}}\left(U^{\theta^{\star}}\right) \leq C\left(\left|\theta-\theta^{\star}\right|^{2}+\tau^{2}\right) \tag{2.2.26}
\end{equation*}
$$

The analysis of the parameter estimation error is somewhat involved, as the state trajectories are merely $\alpha$-Hölder continuous in time with $\alpha<1 / 2$. By leveraging the analytic expression of $X^{\psi^{\theta, \tau}}$, we first show the first-order convergence of $\hat{\theta}^{\tau}$ to $\theta^{\star}$ with

$$
\begin{equation*}
\hat{\theta}^{\tau}:=\left(\mathbb{E}\left[\sum_{i=0}^{N-1} Z_{t_{i}}^{\psi^{\theta, \tau}}\left(Z_{t_{i}}^{\psi^{\theta, \tau}}\right)^{\top} \tau\right]\right)^{-1}\left(\mathbb{E}\left[\sum_{i=0}^{N-1} Z_{t_{i}}^{\psi^{\theta, \tau}}\left(X_{t_{i+1}}^{\psi^{\theta, \tau}}-X_{t_{i}}^{\psi^{\theta, \tau}}\right)^{\top}\right]\right) . \tag{2.2.27}
\end{equation*}
$$

We then prove that (2.2.23) enjoys a uniform sub-exponential tail bound for all $\theta$ close to $\theta^{\star}$ and small $\tau$. Comparing (2.2.23) with (2.2.27) and applying the above results allow for bounding the estimation error of (2.2.23) by (2.2.14) with an additional $\mathcal{O}(\tau)$ term (Proposition 2.3.14).

### 2.3 Proofs of Theorems 2.2.2 and 2.2.3

To simplify the notation, for any given $N, m \in \mathbb{N}$ and control $\psi:[0, T] \times \mathbb{R}^{n} \rightarrow \mathbb{R}^{d}$ that is affine in the spatial variable, we introduce the following random variables associated with continuous-time observations:

$$
\begin{align*}
V^{\psi} & =\int_{0}^{T} Z_{t}^{\psi}\left(Z_{t}^{\psi}\right)^{\top} \mathrm{d} t, & Y^{\psi} & =\int_{0}^{T} Z_{t}^{\psi}\left(\mathrm{d} X_{t}^{\psi}\right)^{\top}, \\
V^{\psi, m} & =\frac{1}{m} \sum_{j=1}^{m} \int_{0}^{T} Z_{t}^{\psi, j}\left(Z_{t}^{\psi, j}\right)^{\top} \mathrm{d} t, & Y^{\psi, m} & =\frac{1}{m} \sum_{j=1}^{m} \int_{0}^{T} Z_{t}^{\psi, j}\left(\mathrm{~d} X_{t}^{\psi, j}\right)^{\top}, \tag{2.3.1}
\end{align*}
$$

and the random variables associated with discrete-time observations with stepsize $\tau=T / N$ :

$$
\begin{align*}
V^{\psi, \tau} & =\sum_{i=0}^{N-1} Z_{t_{i}}^{\psi}\left(Z_{t_{i}}^{\psi}\right)^{\top} \tau, & Y^{\psi, \tau} & =\sum_{i=0}^{N-1} Z_{t_{i}}^{\psi}\left(X_{t_{i+1}}^{\psi}-X_{t_{i}}^{\psi}\right)^{\top} \\
V^{\psi, \tau, m} & =\frac{1}{m} \sum_{j=1}^{m} \sum_{i=0}^{N-1} Z_{t_{i}}^{\psi, j}\left(Z_{t_{i}}^{\psi, j}\right)^{\top} \tau, & Y^{\psi, \tau, m} & =\frac{1}{m} \sum_{j=1}^{m} \sum_{i=0}^{N-1} Z_{t_{i}}^{\psi, j}\left(X_{t_{i+1}}^{\psi, j}-X_{t_{i}}^{\psi, j}\right)^{\top}, \tag{2.3.2}
\end{align*}
$$

where $X^{\psi}$ is the state process associated with the parameter $\theta^{\star}$ and the control $\psi$ (cf. (2.2.6)), $Z_{t}^{\psi}=\binom{X^{\psi}}{\psi\left(t, X_{t}^{\psi}\right)}$ for all $t \in[0, T]$, and $\left(X^{\psi, j}, Z^{\psi, j}\right)_{j=1}^{m}$ are independent copies of $\left(X^{\psi}, Z^{\psi}\right)$.

### 2.3.1 Convergence and stability of Riccati equations and feedback controls

Lemma 2.3.1. Suppose (H.1(1)) holds. Then for all $\theta=(A, B)^{\top} \in \mathbb{R}^{(n+d) \times n}$, the Riccati equation

$$
\begin{equation*}
\frac{\mathrm{d}}{\mathrm{~d} t} P_{t}+A^{\top} P_{t}+P_{t} A-P_{t} B R^{-1} B^{\top} P_{t}+Q=0, \quad t \in[0, T] ; \quad P_{T}=0 . \tag{2.3.3}
\end{equation*}
$$

admits a unique solution $P^{\theta} \in C\left([0, T] ; \mathbb{R}^{n \times n}\right)$. Moreover, the map $\mathbb{R}^{(n+d) \times n} \ni \theta \mapsto P^{\theta} \in$ $C^{1}\left([0, T] ; \mathbb{R}^{n \times n}\right)$ is continuously differentiable.

Proof. It has been shown in [169, Corollary 2.10 on p. 297] that under (H.1(1)), for all $\theta=(A, B)^{\top} \in \mathbb{R}^{(n+d) \times n}$, (2.3.3) admits a unique solution $P^{\theta} \in C\left([0, T] ; \mathbb{R}^{n \times n}\right)$ such that $P_{t}^{\theta} \in \mathbb{S}_{0}^{n}$ for all $t \in[0, T]$. It remains to to prove the continuous differentiability of $\theta \mapsto P^{\theta}$.

To this end, consider the Banach spaces $\mathbb{X}=\mathbb{R}^{(n+d) \times n} \times C^{1}\left([0, T] ; \mathbb{R}^{n \times n}\right)$ and $\mathbb{Y}=$ $C\left([0, T] ; \mathbb{R}^{n \times n}\right) \times \mathbb{R}^{n \times n}$, and the operator $\Phi: \mathbb{X} \rightarrow \mathbb{Y}$ defined by

$$
\mathbb{X} \ni(\theta, P) \mapsto \Phi(\theta, P):=\left(F(\theta, P), P_{T}\right) \in \mathbb{Y}
$$

where $F(\theta, P)_{t}=\frac{\mathrm{d}}{\mathrm{d} t} P_{t}+A^{\top} P_{t}+P_{t} A-P_{t} B R^{-1} B^{\top} P_{t}+Q$ for all $t \in[0, T]$. Observe that for all $\theta \in \mathbb{R}^{(n+d) \times n}, \Phi\left(P^{\theta}, \theta\right)=0$. Moreover, one can easily show that for any $(P, \theta) \in \mathbb{X}$, $\Phi$ is continuously Fréchet differentiable at $(P, \theta)$, and the partial derivative $\frac{\partial}{\partial P} \Phi(\theta, P)$ : $C^{1}\left([0, T] ; \mathbb{R}^{n \times n}\right) \rightarrow \mathbb{Y}$ is a bounded linear operator such that for all $\tilde{P} \in C^{1}\left([0, T] ; \mathbb{R}^{n \times n}\right)$,

$$
\frac{\partial}{\partial P} \Phi(\theta, P)(\tilde{P})=\binom{\left(\frac{\mathrm{d}}{\mathrm{~d} t} \tilde{P}_{t}+A^{\top} \tilde{P}_{t}+\tilde{P}_{t} A-\tilde{P}_{t} B R^{-1} B^{\top} P_{t}-P_{t} B R^{-1} B^{\top} \tilde{P}_{t}\right)_{t \in[0, T]}}{\tilde{P}_{T}} \in \mathbb{Y}
$$

Classical well-posedness results of linear differential equations and the boundedness of $P$ imply that $\frac{\partial}{\partial P} \Phi(P, \theta): C^{1}\left([0, T] ; \mathbb{R}^{n \times n}\right) \rightarrow \mathbb{Y}$ has a bounded inverse (and hence a bijection). Thus, applying the implicit function theorem (see [38, Theorem 7.13-1]) to $\Phi$ proves that $\mathbb{R}^{(n+d) \times n} \ni \theta \mapsto P^{\theta} \in C^{1}\left([0, T] ; \mathbb{R}^{n \times n}\right)$ is continuously differentiable.

The following lemma establishes the stability of the Riccati difference operator, which is crucial for the subsequent convergence analysis.

Lemma 2.3.2. Suppose (H.1(1)) holds. For each $\theta=(A, B)^{\top} \in \mathbb{R}^{(n+d) \times n}$ and $N \in \mathbb{N}$, let $\tau=T / N$ and the function $\Gamma_{\tau}^{\theta}: \mathbb{S}_{0}^{n} \rightarrow \mathbb{S}_{0}^{n}$ such that for all $P \in \mathbb{S}_{0}^{n}$,

$$
\begin{equation*}
\Gamma_{\tau}^{\theta}(P):=\tau Q+(I+\tau A)^{\top} P(I+\tau A)-(I+\tau A)^{\top} P \tau B\left(R+\tau B^{\top} P B\right)^{-1} B^{\top} P(I+\tau A) \tag{2.3.4}
\end{equation*}
$$

Then for all $P, P^{\prime} \in \mathbb{S}_{0}^{n}$,
(1) $\left\|\Gamma_{\tau}^{\theta}(P)\right\|_{2} \leq \tau\|Q\|_{2}+\left(1+\tau\|A\|_{2}\right)^{2}\|P\|_{2}$,
(2) $\left\|\Gamma_{\tau}^{\theta}(P)-\Gamma_{\tau}^{\theta}\left(P^{\prime}\right)\right\|_{2} \leq\left(1+\tau\left\|R^{-1}\right\|_{2}\|B\|_{2}^{2} \max \left\{\|P\|_{2},\left\|P^{\prime}\right\|_{2}\right\}\right)^{2}\left(1+\tau\|A\|_{2}\right)^{2}\left\|P-P^{\prime}\right\|_{2}$.

Proof. Item (1) follows directly from the definition of $\Gamma_{\tau}^{\theta}$ and the identity that $\left\|\Gamma_{\tau}^{\theta}(P)\right\|_{2}=$ $\sup \left\{x^{\top} \Gamma_{\tau}^{\theta}(P) x\left|x \in \mathbb{R}^{n},|x|=1\right\}\right.$. We now prove Item (2). Let $\delta P=P-P^{\prime}$ and $\delta \Gamma(P)=$ $\Gamma_{\tau}^{\theta}(P)-\Gamma_{\tau}^{\theta}\left(P^{\prime}\right)$, by [24, Lemma 10.1],

$$
\delta \Gamma(P)=F^{\top} \delta P F-F^{\top} \delta P \tau B\left(\tau B^{\top} P \tau B+\tau R\right)^{-1} \tau B^{\top} \delta P F
$$

with $F=\left(I-\tau B\left(\tau B^{\top} P^{\prime} \tau B+\tau R\right)^{-1} \tau B^{\top} P^{\prime}\right)(I+\tau A)$. Thus for all $x \in \mathbb{R}^{n}, x^{\top} \delta \Gamma(P) x \leq$ $\|\delta P\|_{2}\|F\|_{2}^{2}|x|^{2}$, which along with $\left\|\left(\tau B^{\top} P^{\prime} B+R\right)^{-1}\right\|_{2} \leq\left\|R^{-1}\right\|_{2}$ implies

$$
x^{\top} \delta \Gamma(P) x \leq\|\delta P\|_{2}\left(1+\tau\left\|R^{-1}\right\|_{2}\|B\|_{2}^{2}\left\|P^{\prime}\right\|_{2}\right)^{2}\left(1+\tau\|A\|_{2}\right)^{2}|x|^{2}, \quad x \in \mathbb{R}^{n} .
$$

Hence, interchanging the roles of $P$ and $P^{\prime}$ in the above inequality and taking the supremum over $x \in \mathbb{R}^{n}$ lead to the desired estimate.

The following proposition establishes the first-order convergence of the Riccati difference equation and the associated feedback controls, as the stepsize tends to zero.

Proposition 2.3.3. Suppose (H.1(1)) holds and let $\Theta$ be a bounded subset of $\mathbb{R}^{(n+d) \times n}$. For each $\theta=(A, B)^{\top} \in \Theta$ and $N \in \mathbb{N}$, let $\left(P_{i}^{\theta, \tau}\right)_{i=0}^{N}$ such that $P_{N}^{\theta, \tau}=0$ and $P_{i}^{\theta, \tau}=\Gamma_{\tau}^{\theta}\left(P_{i+1}^{\theta, \tau}\right)$ for all $i=0, \ldots, N-1$, with $\Gamma_{\tau}^{\theta}$ defined in (2.3.4) with $\tau=T / N$. Then there exists a constant $C \geq 0$ such that for all $\theta \in \Theta, N \in \mathbb{N}$,

$$
\sup _{i=0, \ldots, N-1} \sup _{t \in[i \tau,(i+1) \tau)}\left(\left\|P_{t}^{\theta}-P_{i}^{\theta, \tau}\right\|_{2}+\left\|K_{t}^{\theta}-K_{i}^{\theta, \tau}\right\|_{2}\right) \leq C \tau
$$

where $P^{\theta} \in C^{1}\left([0, T] ; \mathbb{R}^{n \times n}\right)$ satisfies (2.3.3), $K_{t}^{\theta}=-R^{-1} B^{\top} P_{t}^{\theta}$ for all $t \in[0, T]$ and $K_{i}^{\theta, \tau}=$ $-\left(R+\tau B^{\top} P_{i+1}^{\theta, \tau} B\right)^{-1} B^{\top} P_{i+1}^{\theta, \tau}(I+\tau A)$ for all $i=0, \ldots, N-1$.

Proof. Throughout this proof, we shall fix $\theta \in \Theta, N \in \mathbb{N}$, let $t_{i}=i \tau$ for all $i=0, \ldots, N$, and denote by $C$ a generic constant independent of $N$ and $\theta$. By the continuity of the map $\theta \mapsto P^{\theta}$ (Lemma 2.3.1) and the boundedness of $\Theta$, there exists a constant $C$ such that $\left\|P^{\theta}\right\|_{C^{1}\left([0, T] ; \mathbb{R}^{n \times n}\right)} \leq C$ for all $\theta \in \Theta$, which implies $\left\|P_{t}^{\theta}-P_{s}^{\theta}\right\|_{2} \leq C|t-s|$ for all $t, s \in[0, T]$. Consequently, it suffices to prove $\left\|P_{t_{i}}^{\theta}-P_{i}^{\theta, \tau}\right\|_{2}+\left\|K_{t_{i+1}}^{\theta}-K_{i}^{\theta, \tau}\right\|_{2} \leq C \tau$ for all $i=0, \ldots, N-1$.

We start by making two important observations. By Lemma 2.3.2 Item (1), $\left\|P_{i}^{\theta, \tau}\right\|_{2} \leq$ $\tau C+(1+C \tau)\left\|P_{i+1}^{\theta, \tau}\right\|_{2}$ for all $i=0, \ldots, N-1$, which along with Gronwall's inequality gives $\left\|P_{i}^{\theta, \tau}\right\|_{2} \leq C$ for all $i=0, \ldots, N$. Moreover, by (2.3.4), for all $P \in \mathbb{S}_{0}^{n}$,

$$
\begin{aligned}
\Gamma_{\tau}^{\theta}(P)= & \tau Q+P+\tau\left(A^{\top} P+P A\right)+\tau^{2} A^{\top} P A \\
& -\tau\left(P B\left(R+\tau B^{\top} P B\right)^{-1} B^{\top} P+\tau\left(A^{\top} H+H A^{\top}\right)+\tau^{2} A^{\top} H A\right),
\end{aligned}
$$

with $H:=P B\left(R+\tau B^{\top} P B\right)^{-1} B^{\top} P$. Hence for any given $i=0, \ldots, N-1$, we see from (2.3.3) that

$$
\begin{aligned}
& P_{t_{i}}^{\theta}-\Gamma_{\tau}^{\theta}\left(P_{t_{i+1}}^{\theta}\right) \\
&= \int_{t_{i}}^{t_{i+1}}\left(A^{\top}\left(P_{t}^{\theta}-P_{t_{i+1}}^{\theta}\right)+\left(P_{t}^{\theta}-P_{t_{i+1}}^{\theta}\right) A\right) \mathrm{d} t-\int_{t_{i}}^{t_{i+1}}\left(P_{t}^{\theta} B R^{-1} B^{\top} P_{t}^{\theta}-P_{t_{i+1}}^{\theta} B R^{-1} B^{\top} P_{t_{i+1}}^{\theta}\right) \mathrm{d} t \\
&-\int_{t_{i}}^{t_{i+1}}\left(P_{t_{i+1}}^{\theta} B R^{-1} B^{\top} P_{t_{i+1}}^{\theta}-P_{t_{i+1}}^{\theta} B\left(R+\tau B^{\top} P_{t_{i+1}}^{\theta} B\right)^{-1} B^{\top} P_{t_{i+1}}^{\theta}\right) \mathrm{d} t \\
&+\tau^{2}\left(-A^{\top} P_{t_{i+1}}^{\theta} A+A^{\top} H_{i+1}^{\theta}+H_{i+1}^{\theta} A^{\top}+\tau A^{\top} H_{i+1}^{\theta} A\right),
\end{aligned}
$$

with $H_{i+1}^{\theta}=P_{t_{i+1}}^{\theta} B\left(R+\tau B^{\top} P_{t_{i+1}}^{\theta} B\right)^{-1} B^{\top} P_{t_{i+1}}^{\theta}$. Since $\left\|P^{\theta}\right\|_{C^{1}\left([0, T] ; \mathbb{R}^{n \times n}\right)} \leq C$ and $R \in \mathbb{S}_{+}^{d}$, we have $\left\|P_{t_{i}}^{\theta}-\Gamma_{\tau}^{\theta}\left(P_{t_{i+1}}^{\theta}\right)\right\|_{2} \leq C \tau^{2}$ for all $i=0, \ldots, N-1$.

We are ready to show $\max _{i=0, \ldots, N-1}\left(\left\|P_{t_{i}}^{\theta}-P_{i}^{\theta, \tau}\right\|_{2}+\left\|K_{t_{i+1}}^{\theta}-K_{i}^{\theta, \tau}\right\|_{2}\right) \leq C \tau$. For any given $i=0, \ldots, N-1$, by Lemma 2.3.2 Item (2) and the uniform boundedness of $\left(P_{t_{i}}^{\theta}\right)_{i=0}^{N}$ and

$$
\begin{aligned}
& \left(P_{i}^{\theta, \tau}\right)_{i=0}^{N}, \\
& \left\|P_{t_{i}}^{\theta}-P_{i}^{\theta, \tau}\right\|_{2} \leq\left\|P_{t_{i}}^{\theta}-\Gamma_{\tau}^{\theta}\left(P_{t_{i+1}}^{\theta}\right)\right\|_{2}+\left\|\Gamma_{\tau}^{\theta}\left(P_{t_{i+1}}^{\theta}\right)-\Gamma_{\tau}^{\theta}\left(P_{i+1}^{\theta, \tau}\right)\right\|_{2} \\
& \leq C \tau^{2}+\left(1+\tau C \max \left\{\left\|P_{t_{i+1}}^{\theta}\right\|_{2},\left\|P_{i+1}^{\theta, \tau}\right\|_{2}\right\}\right)^{2}(1+\tau)\left\|P_{t_{i+1}}^{\theta}-P_{i+1}^{\theta, \tau}\right\|_{2} \\
& \leq C \tau^{2}+(1+\tau C)\left\|P_{t_{i+1}}^{\theta}-P_{i+1}^{\theta, \tau}\right\|_{2},
\end{aligned}
$$

which along with Gronwall's inequality and $P_{T}^{\theta}=P_{N}^{\theta, \tau}=0$ shows the desired convergence rate of $\left(P_{i}^{\theta, \tau}\right)_{i=1}^{N}$. Furthermore, for all $i=0, \ldots, N-1$,

$$
\begin{aligned}
\left\|K_{t_{i+1}}^{\theta}-K_{i}^{\theta, \tau}\right\|_{2} \leq & \left\|\left(R^{-1}-\left(R+\tau B^{\top} P_{i+1}^{\theta, \tau} B\right)^{-1}\right) B^{\top} P_{t_{i+1}}^{\theta}\right\|_{2} \\
& +\left\|\left(R+\tau B^{\top} P_{i+1}^{\theta, \tau} B\right)^{-1} B^{\top}\left(P_{t_{i+1}}^{\theta}-P_{i+1}^{\theta, \tau}(I+\tau A)\right)\right\|_{2} \leq C \tau
\end{aligned}
$$

from the facts that $\left\|P_{t_{i}}^{\theta}\right\|_{2} \leq C,\left\|P_{i}^{\theta, \tau}\right\|_{2} \leq C$ and $\left\|P_{t_{i}}^{\theta}-P_{i}^{\theta, \tau}\right\|_{2} \leq C \tau$ for all $i$.

### 2.3.2 Concentration inequalities for least-squares estimators

In this section, we analyze the concentration behavior of the least-squares estimators (2.2.13) and (2.2.23). We first recall the definition of sub-exponential random variables (see e.g., [158]).

Definition 2.3.1. A random variable $X$ with mean $\mu=\mathbb{E}[X]$ is $(\nu, b)$-sub-exponential for $\nu, b \in[0, \infty)$ if $\mathbb{E}\left[e^{\lambda(X-\mu)}\right] \leq e^{\nu^{2} \lambda^{2} / 2}$ for all $|\lambda|<1 / b$.

Note that a $(\nu, 0)$-sub-exponential random variable is usually called a sub-Gaussian random variable. It is well-known that products of sub-Gaussian random variables are subexponential, and the class of sub-exponential random variables forms a vector space. Moreover, sub-exponential random variables enjoy the following concentration inequality (also known as Bernstein's inequality; see e.g., [158, Equation 2.18 p. 29]).

Lemma 2.3.4. Let $m \in \mathbb{N}, \nu, b \in[0, \infty)$ and $\left(X_{i}\right)_{i=1}^{m}$ be independent $(\nu, b)$-sub-exponential random variables with $\mu=\mathbb{E}\left[X_{i}\right]$ for all $i=1, \ldots, m$. Then for all $\epsilon \geq 0$,

$$
\mathbb{P}\left(\left|\frac{1}{m} \sum_{i=1}^{m} X_{i}-\mu\right| \geq \epsilon\right) \leq 2 \exp \left(-\min \left\{\frac{m \epsilon^{2}}{2 \nu^{2}}, \frac{m \epsilon}{2 b}\right\}\right) .
$$

The following lemma shows double iterated Itô integrals are sub-exponential random variables.

Lemma 2.3.5. Let $L \geq 0$ and $g, h:[0, T] \times[0, T] \rightarrow \mathbb{R}^{n \times n}$ be measurable functions such that $|g(t, s)| \leq L$ and $|h(t, s)| \leq L$ for all $t, s \in[0, T]$. Then there exist $\nu, b \in[0, \infty)$, depending polynomially on $L, n, T$, such that
(1) $\int_{0}^{T}\left(\int_{0}^{t} g(t, s) \mathrm{d} W_{s}\right)^{\top} \mathrm{d} W_{t}$,

$$
\text { (2) } \int_{0}^{T}\left(\int_{0}^{t} g(t, s) \mathrm{d} W_{s}\right)^{\top}\left(\int_{0}^{t} h(t, s) \mathrm{d} W_{s}\right) \mathrm{d} t
$$

are ( $\nu, b$ )-sub-exponential,
Proof. We first prove Item (1) by assuming without loss of generality that $\|g(t, s)\|_{2} \leq L$ for all $t, s \in[0, T]$, and by defining $V^{q}:=\int_{0}^{T}\left(\int_{0}^{t} q(t, s) \mathrm{d} W_{s}\right)^{\top} \mathrm{d} W_{t}$ for any bounded measurable function $q:[0, T] \times[0, T] \rightarrow \mathbb{R}^{n \times n}$. By similar arguments as [37, Lemma 3.2], we have for all $t \in[0, T]$ and $0 \leq \lambda<\frac{1}{2 T}$,

$$
\mathbb{E}\left[\exp \left(2 \lambda V^{\frac{g}{L}}\right)\right] \leq \mathbb{E}\left[\exp \left(2 \lambda V^{I_{n}}\right)\right]=\left(\frac{1}{\sqrt{1-2 \lambda T}} \exp (-\lambda T)\right)^{n}
$$

As $\frac{e^{-\lambda}}{\sqrt{1-2 \lambda}} \leq e^{2 \lambda^{2}}$ for all $|\lambda| \leq 1 / 4$, we see $\mathbb{E}\left[\exp \left(2 \lambda V^{g / L}\right)\right] \leq \exp \left(2 n \lambda^{2} T^{2}\right)$ for all $0 \leq \lambda<\frac{1}{4 T}$. Consequently, for all $0 \leq \lambda<\frac{1}{2 L T}$,

$$
\mathbb{E}\left[\exp \left(\lambda V^{g}\right)\right]=\mathbb{E}\left[\exp \left(2 \frac{\lambda L}{2} V^{\frac{g}{L}}\right)\right] \leq \exp \left(\frac{n L^{2} T^{2} \lambda^{2}}{2}\right)
$$

Replacing $g$ by $-g$ shows the above estimate holds for $|\lambda|<\frac{1}{2 L T}$, which implies the desired sub-exponential property of $V^{g}$.

For Item (2), observe that for each $t \in[0, T]$, the Itô formula allows one to express the product $\left(\int_{0}^{t} g(t, s) \mathrm{d} W_{s}\right)^{\top}\left(\int_{0}^{t} h(t, s) \mathrm{d} W_{s}\right)$ as a linear combination of double iterated Itô integrals and deterministic integrals. Then the desired sub-exponential property follows from the stochastic Fubini theorem (see e.g., [155]) and Item (1).

The following theorem establishes the concentration properties of the random variables involved in the least-squares estimators.

Theorem 2.3.6. Suppose (H.1(1)) holds and let $\Theta$ be a bounded subset of $\mathbb{R}^{(n+d) \times n}$. For each $\theta \in \Theta$ and $N \in \mathbb{N}$, let $\psi^{\theta}$ be defined in (2.2.9), and $\psi^{\theta, \tau}$ be defined in (2.2.20) with stepsize $\tau=T / N$. Then there exist constants $C, \nu, b>0$ such that for all $\theta \in \Theta, N, m \in \mathbb{N}$ and $\epsilon>0$,

$$
\begin{aligned}
& \max \left\{\mathbb{P}\left(\left|V^{\psi^{\theta}, m}-\mathbb{E}\left[V^{\psi^{\theta}}\right]\right| \geq \epsilon\right), \mathbb{P}\left(\left|Y^{\psi^{\theta}, m}-\mathbb{E}\left[Y^{\psi^{\theta}}\right]\right| \geq \epsilon\right)\right. \\
& \left.\quad \mathbb{P}\left(\left|V^{\psi^{\theta, \tau}, \tau, m}-\mathbb{E}\left[V^{\psi^{\theta, \tau}, \tau}\right]\right| \geq \epsilon\right), \mathbb{P}\left(\left|Y^{\psi^{\theta, \tau}, \tau, m}-\mathbb{E}\left[Y^{\psi^{\theta, \tau}, \tau}\right]\right| \geq \epsilon\right)\right\} \leq C \exp \left(-\frac{1}{C} \min \left\{\frac{m \epsilon^{2}}{\nu^{2}}, \frac{m \epsilon}{b}\right\}\right)
\end{aligned}
$$

where $V^{\psi^{\theta}}, Y^{\psi^{\theta}}, V^{\psi^{\theta}, m}, Y^{\psi^{\theta}, m}$ are defined in (2.3.1), and $V^{\psi^{\theta, \tau}, \tau}, Y^{\psi^{\theta, \tau}, \tau}, V^{\psi^{\theta, \tau}, \tau, m}, Y^{\psi^{\theta, \tau}, \tau, m}$ are defined in (2.3.2).
Proof. We first show there exist $\nu, b>0$ such that all entries of $V^{\psi^{\theta}}, Y^{\psi^{\theta}} Y^{\psi^{\theta, \tau}, \tau}, V^{\psi^{\theta, \tau}, \tau}$ are $(\nu, b)$-sub-exponential for all $\theta \in \Theta$ and $N \in \mathbb{N}$. By (2.3.1), we have

$$
V^{\psi^{\theta}}=\int_{0}^{T}\binom{X_{t}^{\psi^{\theta}}}{K_{t}^{\theta} X_{t}^{\psi^{\theta}}}\left(\begin{array}{ll}
\left(X_{t}^{\psi^{\theta}}\right)^{\top} & \left.\left(K_{t}^{\theta} X_{t}^{\psi^{\theta}}\right)^{\top}\right) \mathrm{d} t, \quad Y^{\psi^{\theta}}=V^{\psi^{\theta}} \theta^{\star}+\int_{0}^{T}\binom{X_{t}^{\psi^{\theta}}}{K_{t}^{\theta} X_{t}^{\psi^{\theta}}}\left(\mathrm{d} W_{t}\right)^{\top} . . . . . . . . .
\end{array}\right.
$$

Moreover, applying the variation-of-constants formula (see e.g., [113, Theorem 3.1 p. 96]) to (2.2.11) shows that $X_{t}^{\psi^{\theta}}=\Phi_{t}^{\theta}\left(x_{0}+\int_{0}^{t}\left(\Phi_{s}^{\theta}\right)^{-1} \mathrm{~d} W_{s}\right)$ for all $t \in[0, T]$, where $\Phi^{\theta} \in C\left([0, T] ; \mathbb{R}^{n \times n}\right)$ is the fundamental solution of $\mathrm{d} \Phi_{t}^{\theta}=\left(A^{\star}+B^{\star} K_{t}^{\theta}\right) \Phi_{t}^{\theta} \mathrm{d} t$. The continuity of $\mathbb{R}^{(n+d) \times n} \ni$ $\theta \mapsto K^{\theta} \in C\left([0, T] ; \mathbb{R}^{d \times n}\right)$ (cf. Proposition 2.3.1) and the boundedness of $\Theta$ implies that $K^{\theta}, \Phi^{\theta},\left(\Phi^{\theta}\right)^{-1}$ are uniformly bounded for all $\theta \in \Theta$. Consequently, from Lemma 2.3.5, there exist $\nu, b>0$ such that all entries of $V^{\psi^{\theta}}$ and $Y^{\psi^{\theta}}$ are $(\nu, b)$-sub-exponential.

Similarly, by (2.2.21) and (2.3.2),

$$
\begin{aligned}
& V^{\psi^{\theta, \tau}, \tau}=\int_{0}^{T} \sum_{i=0}^{N-1} \mathbf{1}_{\left[t_{i}, t_{i+1}\right)}(t)\binom{X_{t_{i}}^{\psi^{\theta, \tau}}}{K_{t_{i}}^{\theta, \tau} X_{t_{i}}^{\psi^{\theta, \tau}}}\left(\begin{array}{l}
\left.\left(X_{t_{i}}^{\psi^{\theta, \tau}}\right)^{\top} \quad\left(K_{t_{i}}^{\theta, \tau} X_{t_{i}}^{\psi^{\theta, \tau}}\right)^{\top}\right) \mathrm{d} t,, ~\left(x^{\theta}\right.
\end{array}\right)
\end{aligned}
$$

$$
\begin{aligned}
& +\int_{0}^{T} \sum_{i=0}^{N-1} \mathbf{1}_{\left[t_{i}, t_{i+1}\right)}(t)\binom{X_{t_{i}}^{\psi^{\theta, \tau}}}{K_{t_{i}}^{\theta, \tau} X_{t_{i}}^{\psi^{\theta, \tau}}}\left(\mathrm{d} W_{t}\right)^{\top},
\end{aligned}
$$

where $X_{t}^{\psi^{\theta, \tau}}=\Phi_{t}^{\theta, \tau}\left(x_{0}+\int_{0}^{t}\left(\Phi_{s}^{\theta, \tau}\right)^{-1} \mathrm{~d} W_{s}\right)$ for all $t \in[0, T]$, and $\Phi^{\theta, \tau} \in C\left([0, T] ; \mathbb{R}^{n \times n}\right)$ is the fundamental solution of $\mathrm{d} \Phi_{t}^{\theta, \tau}=\left(A^{\star}+B^{\star} K_{t}^{\theta, \tau}\right) \Phi_{t}^{\theta, \tau} \mathrm{d} t$. By Proposition 2.3.3, $K^{\theta, \tau}, \Phi^{\theta, \tau},\left(\Phi^{\theta, \tau}\right)^{-1}$ are uniformly bounded for all $\theta \in \Theta$ and $N \in \mathbb{N}$, which along with Lemma 2.3.5 leads to the desired sub-exponential properties of $Y^{\psi^{\theta, \tau}, \tau}$ and $V^{\psi^{\theta, \tau}, \tau}$.

Finally, since $\mathbb{P}\left(\left|\sum_{i=1}^{\ell} X_{i}\right| \geq \epsilon\right) \leq \sum_{i=1}^{\ell} \mathbb{P}\left(\left|X_{i}\right| \geq \epsilon / \ell\right)$ for all $\ell \in \mathbb{N}$ and random variables $\left(X_{i}\right)_{i=1}^{\ell}$, we can apply Lemma 2.3.4 to each component of $V^{\psi^{\theta}}, Y^{\psi^{\theta}} Y^{\psi^{\theta, \tau}, \tau}$ and $V^{\psi^{\theta, \tau}, \tau}$, and conclude the desired concentration inequality with a constant $C$ depending polynomially on $n, d$.

### 2.3.3 Regret analysis of continuous-time least-squares algorithm

This section is devoted to the proof of Theorem 2.2.2, which consists of three steps: (1) We first quantify the performance gap between applying feedback controls for an incorrect model and that for the true model; our proof exploits the stability of Riccati equations established in Lemma 2.3.1; (2) We then estimate the parameter estimation error in terms of the number of learning episodes based on the sub-exponential tail behavior of the least-squares estimator (2.2.13); (3) Finally, we estimate the regret for the feedback controls $\left(\psi^{\theta_{\ell}}\right)_{\ell \in \mathbb{N}}$ in Algorithm 8, thus establishing Theorem 2.2.2.

Step 1: Analysis of the performance gap. We start by establishing a quadratic expansion of the cost function at any open-loop control.

Proposition 2.3.7. Suppose (H.1(1)) holds. Let $\psi^{\theta^{\star}}$ be defined in (2.2.3), $X^{\theta^{\star}}$ be the state process associated with $\psi^{\theta^{\star}}(c f .(2.2 .5))$, and $U^{\theta^{\star}} \in \mathcal{H}^{2}\left(\mathbb{R}^{d}\right)$ be such that for all $t \in[0, T]$,
$U_{t}^{\theta^{\star}}=\psi^{\theta^{\star}}\left(t, X_{t}^{\theta^{\star}}\right)$. Then for all $U \in \mathcal{H}^{2}\left(\mathbb{R}^{d}\right)$,

$$
\begin{equation*}
J^{\theta^{\star}}(U)-J^{\theta^{\star}}\left(U^{\theta^{\star}}\right) \leq\|Q\|_{2}\left\|X^{\theta^{\star}, U}-X^{\theta^{\star}}\right\|_{\mathcal{H}^{2}\left(\mathbb{R}^{n}\right)}^{2}+\|R\|_{2}\left\|U-U^{\theta^{*}}\right\|_{\mathcal{H}^{2}\left(\mathbb{R}^{d}\right)}^{2} \tag{2.3.5}
\end{equation*}
$$

where $X^{\theta^{\star}, U}$ is the state process controlled by $U(c f .(2.2 .2))$, and $J^{\theta^{\star}}: \mathcal{H}^{2}\left(\mathbb{R}^{d}\right) \rightarrow \mathbb{R}$ is defined in (2.2.1).
Proof. For notational simplicity, for all $U \in \mathcal{H}^{2}\left(\mathbb{R}^{d}\right)$ and $\epsilon>0$, we write $U^{\epsilon}=U^{\theta^{\star}}+$ $\epsilon\left(U-U^{\theta^{\star}}\right)$, denote by $X^{\epsilon}=X^{\theta^{\star}, U^{\epsilon}}$ the associated state process defined by (2.2.2), and by $X^{U}=X^{0}=X^{\theta^{\star}, U}$. The affineness of (2.2.2) implies that $X^{\epsilon}=(1-\epsilon) X^{\theta^{\star}}+\epsilon X^{U}$ for all $\epsilon>0$. Hence, for all $U \in \mathcal{H}^{2}\left(\mathbb{R}^{d}\right)$,

$$
\begin{aligned}
& \lim _{\epsilon \searrow 0} \frac{1}{\epsilon}\left(J^{\theta^{\star}}\left(U^{\epsilon}\right)-J^{\theta^{\star}}\left(U^{\theta^{\star}}\right)\right) \\
&= \lim _{\epsilon \searrow 0} \frac{1}{\epsilon} \mathbb{E}\left[\int _ { 0 } ^ { T } \left(\left((1-\epsilon) X_{t}^{\theta^{\star}}+\epsilon X_{t}^{U}\right)^{\top} Q\left((1-\epsilon) X_{t}^{\theta^{\star}}+\epsilon X_{t}^{U}\right)-\left(X_{t}^{\theta^{\star}}\right)^{\top} Q X_{t}^{\theta^{\star}}\right.\right. \\
&\left.\left.+\left((1-\epsilon) U_{t}^{\theta^{\star}}+\epsilon U_{t}\right)^{\top} R\left((1-\epsilon) U_{t}^{\theta^{\star}}+\epsilon U_{t}\right)-\left(U_{t}^{\theta^{\star}}\right)^{\top} R U_{t}^{\theta^{\star}}\right) \mathrm{d} t\right] \\
&= \lim _{\epsilon \searrow 0} \epsilon \mathbb{E}\left[\int_{0}^{T}\left(\left(X_{t}^{U}-X_{t}^{\theta^{\star}}\right)^{\top} Q\left(X_{t}^{U}-X_{t}^{\theta^{\star}}\right)+\left(U_{t}-U_{t}^{\theta^{\star}}\right)^{\top} R\left(U_{t}-U_{t}^{\theta^{\star}}\right)\right) \mathrm{d} t\right] \\
&+2 \mathbb{E}\left[\int_{0}^{T}\left(\left(X_{t}^{U}-X_{t}^{\theta^{\star}}\right)^{\top} Q X_{t}^{\theta^{\star}}+\left(U_{t}-U_{t}^{\theta^{\star}}\right)^{\top} R U_{t}^{\theta^{\star}}\right) \mathrm{d} t\right] \\
&= 2 \mathbb{E}\left[\int_{0}^{T}\left(\left(X_{t}^{U}-X_{t}^{\theta^{\star}}\right)^{\top} Q X_{t}^{\theta^{\star}}+\left(U_{t}-U_{t}^{\theta^{\star}}\right)^{\top} R U_{t}^{\theta^{\star}}\right) \mathrm{d} t\right]
\end{aligned}
$$

which is based on the fact that $X^{U}-X^{\theta^{\star}} \in \mathcal{H}^{2}\left(\mathbb{R}^{n}\right)$ and $U-U^{\star} \in \mathcal{H}^{2}\left(\mathbb{R}^{d}\right)$. As $U^{\theta}$ is the optimal control of $J^{\theta^{\star}}, J^{\theta^{\star}}(U) \geq J^{\theta^{\star}}\left(U^{\theta^{\star}}\right)$ for all $U \in \mathcal{H}^{2}\left(\mathbb{R}^{d}\right)$. Hence for all $U \in \mathcal{H}^{2}\left(\mathbb{R}^{d}\right)$,

$$
\begin{equation*}
\mathbb{E}\left[\int_{0}^{T}\left(\left(X_{t}^{U}-X_{t}^{\theta^{\star}}\right)^{\top} Q X_{t}^{\theta^{\star}}+\left(U_{t}-U_{t}^{\theta^{\star}}\right)^{\top} R U_{t}^{\theta^{\star}}\right) \mathrm{d} t\right]=\lim _{\epsilon \searrow 0} \frac{1}{2 \epsilon}\left(J^{\theta^{\star}}\left(U^{\epsilon}\right)-J^{\theta^{\star}}\left(U^{\theta^{\star}}\right)\right) \geq 0 \tag{2.3.6}
\end{equation*}
$$

We now prove that the above quantity is in fact zero for all $U \in \mathcal{H}^{2}\left(\mathbb{R}^{d}\right)$. To this end, let $U \in \mathcal{H}^{2}\left(\mathbb{R}^{d}\right)$ be a given (open-loop) control, and consider $\tilde{U}=U^{\theta^{\star}}-\left(U-U^{\theta^{*}}\right)$. Then by the affineness of (2.2.2), $X^{\tilde{U}}-X^{\theta^{\star}}$ satisfies the following controlled dynamics:

$$
\begin{equation*}
\mathrm{d} X_{t}=\left(A^{\star} X_{t}-B^{\star}\left(U-U^{\theta^{\star}}\right)_{t}\right) \mathrm{d} t, \quad t \in[0, T] ; \quad X_{0}=0 . \tag{2.3.7}
\end{equation*}
$$

Moreover, one can verify by the affineness of (2.2.2) that $-\left(X^{U}-X^{\theta^{\star}}\right)$ also satisfies the dynamics (2.3.7), which along with the uniqueness of solutions to (2.3.7) shows that $X^{\tilde{U}}-$ $X^{\theta^{\star}}=-\left(X^{U}-X^{\theta^{\star}}\right)$. Therefore, applying (2.3.6) with $U=\tilde{U}$ implies that

$$
\begin{aligned}
0 & \leq \mathbb{E}\left[\int_{0}^{T}\left(\left(X_{t}^{\tilde{U}}-X_{t}^{\theta^{\star}}\right)^{\top} Q X_{t}^{\theta^{\star}}+\left(\tilde{U}_{t}-U_{t}^{\theta^{\star}}\right)^{\top} R U_{t}^{\theta^{\star}}\right) \mathrm{d} t\right] \\
& =-\mathbb{E}\left[\int_{0}^{T}\left(\left(X_{t}^{U}-X_{t}^{\theta^{\star}}\right)^{\top} Q X_{t}^{\theta^{\star}}+\left(U_{t}-U_{t}^{\theta^{\star}}\right)^{\top} R U_{t}^{\theta^{\star}}\right) \mathrm{d} t\right] \leq 0
\end{aligned}
$$

Hence for all $U \in \mathcal{H}^{2}\left(\mathbb{R}^{d}\right)$,

$$
\mathbb{E}\left[\int_{0}^{T}\left(\left(X_{t}^{U}-X_{t}^{\theta^{\star}}\right)^{\top} Q X_{t}^{\theta^{\star}}+\left(U_{t}-U_{t}^{\theta^{\star}}\right)^{\top} R U_{t}^{\theta^{\star}}\right) \mathrm{d} t\right]=0,
$$

which leads to the desired result (2.3.5) due to the following identify:

$$
\begin{aligned}
J^{\theta^{\star}}(U)-J^{\theta^{\star}}\left(U^{\theta^{\star}}\right)= & \mathbb{E}\left[\int_{0}^{T}\left(\left(X_{t}^{U}\right)^{\top} Q X_{t}^{U}-\left(X_{t}^{\theta^{\star}}\right)^{\top} Q X_{t}^{\theta^{\star}}+U_{t}^{\top} R U_{t}-\left(U_{t}^{\theta^{\star}}\right)^{\top} R U_{t}^{\theta^{\star}}\right) \mathrm{d} t\right] \\
= & \mathbb{E}\left[\int_{0}^{T}\left(\left(X_{t}^{U}\right)^{\top} Q X_{t}^{U}-\left(X_{t}^{\theta^{\star}}\right)^{\top} Q X_{t}^{\theta^{\star}}+U_{t}^{\top} R U_{t}-\left(U_{t}^{\theta^{\star}}\right)^{\top} R U_{t}^{\theta^{\star}}\right) \mathrm{d} t\right] \\
& -2 \mathbb{E}\left[\int_{0}^{T}\left(\left(X_{t}^{U}-X_{t}^{\theta^{\star}}\right)^{\top} Q X_{t}^{\theta^{\star}}+\left(U_{t}-U_{t}^{\theta^{\star}}\right)^{\top} R U_{t}^{\theta^{\star}}\right) \mathrm{d} t\right] \\
= & \mathbb{E}\left[\int_{0}^{T}\left(\left(X_{t}^{U}-X_{t}^{\theta^{\star}}\right)^{\top} Q\left(X_{t}^{U}-X_{t}^{\theta^{\star}}\right)+\left(U_{t}-U_{t}^{\theta^{\star}}\right)^{\top} R\left(U_{t}-U_{t}^{\theta^{\star}}\right) \mathrm{d} t\right)\right] .
\end{aligned}
$$

Armed with Proposition 2.3.7, the following proposition quantifies the quadratic performance gap of a greedy policy $\psi^{\theta}$.

Proposition 2.3.8. Suppose (H.1(1)) holds and let $\Theta$ be a bounded subset of $\mathbb{R}^{(n+d) \times n}$. For each $\theta \in \Theta$, let $\psi^{\theta}$ be defined in (2.2.9), let $X^{\psi^{\theta}}$ be the state process associated with $\psi^{\theta}$ (cf. (2.2.11)), let $\psi^{\theta^{\star}}$ be defined in (2.2.3), and let $X^{\theta^{\star}}$ be the state process associated with $\psi^{\theta^{\star}}(c f .(2.2 .5))$. Then there exists a constant $C$ such that

$$
\left|J^{\theta^{\star}}\left(U^{\psi^{\theta}}\right)-J^{\theta^{\star}}\left(U^{\theta^{\star}}\right)\right| \leq C\left|\theta-\theta^{\star}\right|^{2}, \quad \forall \theta \in \Theta
$$

where $U_{t}^{\psi^{\theta}}=\psi^{\theta}\left(t, X_{t}^{\psi^{\theta}}\right)$ and $U_{t}^{\theta^{\star}}=\psi^{\theta^{\star}}\left(t, X_{t}^{\theta^{\star}}\right)$ for all $t \in[0, T]$, and $J^{\theta^{\star}}$ is defined in (2.2.1).
Proof. For all $\theta \in \Theta$, applying Proposition 2.3 .7 with $U=U^{\psi^{\theta}}$ gives

$$
\begin{align*}
& J^{\theta^{\star}}\left(U^{\psi^{\theta}}\right)-J^{\theta^{\star}}\left(U^{\theta^{\star}}\right) \\
& \leq\|Q\|_{2}\left\|X^{\theta^{\star}, U^{\psi^{\theta}}}-X^{\theta^{\star}}\right\|_{\mathcal{H}^{2}\left(\mathbb{R}^{n}\right)}^{2}+\|R\|_{2}\left\|U^{\psi^{\theta}}-U^{\theta^{\star}}\right\|_{\mathcal{H}^{2}\left(\mathbb{R}^{d}\right)}^{2}  \tag{2.3.8}\\
& \leq\|Q\|_{2}\left\|X^{\psi^{\theta}}-X^{\psi^{\theta^{\star}}}\right\|_{\mathcal{H}^{2}\left(\mathbb{R}^{n}\right)}^{2}+\|R\|_{2}\left\|\psi^{\theta}\left(\cdot, X_{\psi^{\psi^{\theta}}}\right)-\psi^{\theta^{\star}}\left(\cdot, X^{\psi^{\theta^{\star}}}\right)\right\|_{\mathcal{H}^{2}\left(\mathbb{R}^{d}\right)}^{2}
\end{align*}
$$

where the last inequality used the fact that $X^{\theta^{\star}, U^{\psi^{\theta}}}=X^{\psi^{\theta}}$ (see (2.2.11)), and the definitions of $U^{\psi^{\theta}}$ and $U^{\theta^{\star}}$. It remains to prove

$$
\left\|X^{\psi^{\theta}}-X^{\psi^{\theta^{\star}}}\right\|_{\mathcal{H}^{2}\left(\mathbb{R}^{n}\right)}+\left\|\psi^{\theta}\left(\cdot, X_{.^{\psi^{\theta}}}\right)-\psi^{\theta^{\star}}\left(\cdot, X_{\cdot}^{\psi^{\theta^{\star}}}\right)\right\|_{\mathcal{H}^{2}\left(\mathbb{R}^{d}\right)} \leq C\left|\theta-\theta^{\star}\right|
$$

for a constant $C$ independent of $\theta$.

Observe that by (2.2.9), for all $(t, x) \in[0, T] \times \mathbb{R}^{n}, \psi^{\theta}(t, x)=K_{t}^{\theta} x$ with $K_{t}^{\theta}=-R^{-1} B^{\top} P_{t}^{\theta}$. Now by Lemma 2.3.1 and the boundedness of $\Theta$, there exists a constant $C \geq 0$ such that $\left\|P^{\theta}\right\|_{C\left(\left[0, T ; \mathbb{R}^{n \times n}\right)\right.} \leq C$ and $\left\|P^{\theta}-P^{\theta^{\star}}\right\|_{C\left(\left[0, T ; \mathbb{R}^{n \times n}\right)\right.} \leq C\left|\theta-\theta^{\star}\right|$ for all $\theta \in \Theta \cup\left\{\theta^{\star}\right\}$, which along with $K_{t}^{\theta}=-R^{-1} B^{\top} P_{t}^{\theta}$ implies that $\left\|K^{\theta}\right\|_{C\left(\left[0, T ; \mathbb{R}^{d \times n}\right)\right.} \leq C$ and $\left\|K^{\theta}-K^{\theta^{*}}\right\|_{C\left(\left[0, T ; \mathbb{R}^{d \times n}\right)\right.} \leq$ $C\left|\theta-\theta^{\star}\right|$. Moreover, observe from (2.2.5) and (2.2.11) that $X_{0}^{\theta^{\star}}=X_{0}^{\psi^{\theta}}$ and for all $t \in[0, T]$,

$$
\mathrm{d}\left(X^{\psi^{\theta^{\star}}}-X^{\psi^{\theta}}\right)_{t}=\left(\left(A^{\star}+B^{\star} K_{t}^{\theta^{\star}}\right)\left(X^{\psi^{\theta^{\star}}}-X^{\psi^{\theta}}\right)_{t}+B^{\star}\left(K_{t}^{\theta^{\star}}-K_{t}^{\theta}\right) X_{t}^{\psi^{\theta}}\right) \mathrm{d} t
$$

which combined with the boundedness of $K^{\theta^{\star}}$ and Gronwall's inequality leads to

$$
\begin{aligned}
\left\|X^{\psi^{\theta^{\star}}}-X^{\psi^{\theta}}\right\|_{\mathcal{H}^{2}\left(\mathbb{R}^{n}\right)} & \leq C\left\|X^{\psi^{\theta^{\star}}}-X^{\psi^{\theta}}\right\|_{\mathcal{S}^{2}\left(\mathbb{R}^{n}\right)} \\
& \leq C\left\|\left(K^{\theta^{\star}}-K^{\theta}\right) X^{\psi^{\theta}}\right\|_{\mathcal{H}^{2}\left(\mathbb{R}^{d}\right)} \leq C\left\|K^{\theta^{\star}}-K^{\theta}\right\|_{C\left(\left[0, T ; \mathbb{R}^{d \times n}\right)\right.}\left\|X^{\psi^{\theta}}\right\|_{\mathcal{H}^{2}\left(\mathbb{R}^{n}\right)} \\
& \leq C\left|\theta-\theta^{\star}\right|, \quad \forall \theta \in \Theta,
\end{aligned}
$$

where the last inequality follows from $\left\|X^{\psi^{\theta}}\right\|_{\mathcal{H}^{2}\left(\mathbb{R}^{n}\right)} \leq C$, as $K^{\theta}$ is uniformly bounded. The above inequality further implies

$$
\begin{aligned}
& \left\|\psi^{\theta}\left(\cdot, X_{\cdot}^{\psi^{\theta}}\right)-\psi^{\theta^{\star}}\left(\cdot, X_{\cdot}^{\psi^{\theta^{\star}}}\right)\right\|_{\mathcal{H}^{2}\left(\mathbb{R}^{d}\right)}=\left\|K_{\cdot}^{\theta} X_{\cdot}^{\psi^{\theta}}-K_{\cdot}^{\theta^{\star}} X_{\cdot}^{\psi^{\theta^{\star}}}\right\|_{\mathcal{H}^{2}\left(\mathbb{R}^{d}\right)} \\
& \leq\left\|\left(K_{\cdot}^{\theta}-K_{\cdot}^{\theta^{\star}}\right) X_{\cdot}^{\psi^{\theta}}\right\|_{\mathcal{H}^{2}\left(\mathbb{R}^{d}\right)}+\left\|K_{\cdot}^{\theta^{\star}}\left(X_{\cdot}^{\psi^{\theta}}-X_{\cdot}^{\psi^{\theta^{\star}}}\right)\right\|_{\mathcal{H}^{2}\left(\mathbb{R}^{d}\right)} \\
& \leq\left\|K^{\theta^{\star}}-K^{\theta}\right\|_{C\left(\left[0, T ; \mathbb{R}^{d \times n}\right)\right.}\left\|X^{\psi^{\theta}}\right\|_{\mathcal{H}^{2}\left(\mathbb{R}^{n}\right)}+\left\|K^{\theta^{\star}}\right\|_{C\left(\left[0, T ; \mathbb{R}^{d \times n}\right)\right.}\left\|X^{\psi^{\theta}}-X^{\psi^{\theta^{\star}}}\right\|_{\mathcal{H}^{2}\left(\mathbb{R}^{n}\right)} \\
& \leq C\left|\theta-\theta^{\star}\right|, \quad \forall \theta \in \Theta,
\end{aligned}
$$

which along with (2.3.8) finishes the desired estimate.

## Step 2: Error bound for parameter estimation.

Proposition 2.3.9. Suppose (H.1(1)) holds and let $\Theta \subset \mathbb{R}^{(n+d) \times n}$ such that there exists $C_{1}>0$ satisfying $\left\|\left(\mathbb{E}\left[V^{\psi^{\theta}}\right]\right)^{-1}\right\|_{2} \leq C_{1}$ and $|\theta| \leq C_{1}$ for all $\theta \in \Theta$, with $V^{\psi^{\theta}}$ defined in (2.3.1). Then there exist constants $\bar{C}_{1}, \bar{C}_{2} \geq 0$, such that for all $\theta \in \Theta$ and $\delta \in(0,1 / 2)$, if $m \geq \bar{C}_{1}(-\ln \delta)$, then with probability at least $1-2 \delta$,

$$
\begin{equation*}
\left|\hat{\theta}-\theta^{\star}\right| \leq \bar{C}_{2}\left(\sqrt{\frac{-\ln \delta}{m}}+\frac{-\ln \delta}{m}+\frac{(-\ln \delta)^{2}}{m^{2}}\right) \tag{2.3.9}
\end{equation*}
$$

where $\hat{\theta}$ denotes the right-hand side of (2.2.13) with the control $\psi^{\theta}$.
Proof. Let us fix $\delta \in(0,1 / 2)$ and $\theta \in \Theta$. By (2.2.12) and (2.2.13), we obtain

$$
\begin{aligned}
\left\|\hat{\theta}-\theta^{\star}\right\|_{2} & =\left\|\left(V^{\psi^{\theta}, m}+\frac{1}{m} I\right)^{-1} Y^{\psi^{\theta}, m}-\left(\mathbb{E}\left[V^{\psi^{\theta}}\right]\right)^{-1} \mathbb{E}\left[Y^{\psi^{\theta}}\right]\right\|_{2} \\
& \leq\left\|\left(V^{\psi^{\theta}, m}+\frac{1}{m} I\right)^{-1}-\left(\mathbb{E}\left[V^{\psi^{\theta}}\right]\right)^{-1}\right\|_{2}\left\|Y^{\psi^{\theta}, m}\right\|_{2}+\left\|\left(\mathbb{E}\left[V^{\psi^{\theta}}\right]\right)^{-1}\right\|_{2}\left\|Y^{\psi^{\theta}, m}-\mathbb{E}\left[Y^{\psi^{\theta}}\right]\right\|_{2} .
\end{aligned}
$$

As $E^{-1}-F^{-1}=F^{-1}(F-E) E^{-1}$ for all nonsingular matrices $E$ and $F$, we have

$$
\begin{align*}
& \left\|\hat{\theta}-\theta^{\star}\right\|_{2} \\
& \leq\left\|\left(V^{\psi^{\theta}, m}+\frac{1}{m} I\right)^{-1}\right\|_{2}\left\|\left(\mathbb{E}\left[V^{\psi^{\theta}}\right]\right)^{-1}\right\|_{2}\left\|Y^{\psi^{\theta}, m}\right\|_{2}\left\|V^{\psi^{\theta}, m}-\mathbb{E}\left[V^{\psi^{\theta}}\right]+\frac{1}{m} I\right\|_{2} \\
& \quad+\left\|\left(\mathbb{E}\left[V^{\psi^{\theta}}\right]\right)^{-1}\right\|_{2}\left\|Y^{\psi^{\theta}, m}-\mathbb{E}\left[Y^{\psi^{\theta}}\right]\right\|_{2} \\
& \leq C_{1}\left(\left\|\left(V^{\psi^{\theta}, m}+\frac{1}{m} I\right)^{-1}\right\|_{2}\left\|Y^{\psi^{\theta}, m}\right\|_{2}\left\|V^{\psi^{\theta}, m}-\mathbb{E}\left[V^{\psi^{\theta}}\right]+\frac{1}{m} I\right\|_{2}+\left\|Y^{\psi^{\theta}, m}-\mathbb{E}\left[Y^{\psi^{\theta}}\right]\right\|_{2}\right), \tag{2.3.10}
\end{align*}
$$

where the last inequality follows from the assumption $\left\|\left(\mathbb{E}\left[V^{\psi^{\theta}}\right]\right)^{-1}\right\|_{2} \leq C_{1}$.
We now estimate each term in the right-hand side of (2.3.10), and denote by $C$ a generic constant independent of $\theta \in \Theta, \delta \in(0,1 / 2), m \in \mathbb{N}$. By Theorem 2.3.6, with probability at least $1-2 \delta,\left\|V^{\psi^{\theta}, m}-\mathbb{E}\left[V^{\psi^{\theta}}\right]\right\|_{2} \leq \delta_{m}$ and $\left\|Y^{\psi^{\theta}, m}-\mathbb{E}\left[Y^{\psi^{\theta}}\right]\right\|_{2} \leq \delta_{m}$, with the constant $\delta_{m}$ given by

$$
\begin{equation*}
\delta_{m}:=C \max \left\{\left(\frac{-\ln \delta}{m}\right)^{\frac{1}{2}}, \frac{-\ln \delta}{m}\right\} . \tag{2.3.11}
\end{equation*}
$$

Let $m$ be a sufficiently large constant satisfying $\delta_{m}+1 / m \leq 1 /\left(2 C_{1}\right)$, where $C_{1}$ is the constant such that $\left\|\left(\mathbb{E}\left[V^{\psi^{\theta}}\right]\right)^{-1}\right\|_{2} \leq C_{1}$ for all $\theta \in \Theta$. Then with probability at least $1-2 \delta$, $\left\|V^{\psi^{\theta}, m}-\mathbb{E}\left[V^{\psi^{\theta}}\right]+\frac{1}{m} I\right\|_{2} \leq \delta_{m}+\frac{1}{m} \leq \frac{1}{2 C_{1}}$, which in turn yields

$$
\lambda_{\min }\left(V^{\psi^{\theta}, m}+\frac{1}{m}\right) \geq \lambda_{\min }\left(\mathbb{E}\left[V^{\psi^{\theta}}\right]\right)-\left\|V^{\psi^{\theta}, m}-\mathbb{E}\left[V^{\psi^{\theta}}\right]+\frac{1}{m} I\right\|_{2} \geq \frac{1}{2 C_{1}}
$$

or equivalently $\left\|\left(V^{\psi^{\theta}, m}+\frac{1}{m}\right)^{-1}\right\|_{2} \leq 2 C_{1}$. Moreover, the continuity of $\mathbb{R}^{(n+d) \times n} \ni \theta \mapsto$ $\mathbb{E}\left[Y^{\psi^{\theta}}\right] \in \mathbb{R}$ implies $\left\|Y^{\psi^{\theta}, m}\right\|_{2} \leq\left\|\mathbb{E}\left[Y^{\psi^{\theta}}\right]\right\|_{2}+\left\|Y^{\psi^{\theta}, m}-\mathbb{E}\left[Y^{\psi^{\theta}}\right]\right\|_{2} \leq C+\left\|Y^{\psi^{\theta}, m}-\mathbb{E}\left[Y^{\psi^{\theta}}\right]\right\|_{2}$. Hence, by (2.3.10),

$$
\begin{aligned}
& \left|\hat{\theta}-\theta^{\star}\right| \\
& \leq C\left(\left(1+\left\|Y^{\psi^{\theta}, m}-\mathbb{E}\left[Y^{\psi^{\theta}}\right]\right\|_{2}\right)\left\|V^{\psi^{\theta}, m}-\mathbb{E}\left[V^{\psi^{\theta}}\right]+\frac{1}{m} I\right\|_{2}+\left\|Y^{\psi^{\theta}, m}-\mathbb{E}\left[Y^{\psi^{\theta}}\right]\right\|_{2}\right) \\
& \leq C\left(\left(\delta_{m}+\frac{1}{m}\right)\left(1+\delta_{m}\right)+\delta_{m}\right) \leq C\left(\delta_{m}+\delta_{m}^{2}+\frac{1}{m}\right) .
\end{aligned}
$$

Substituting (3.3.15) into the above estimate yields the desired estimate (3.3.13). As $\delta \in$ $(0,1 / 2)$, it is clear that $\delta_{m}+1 / m \leq 1 /\left(2 C_{1}\right)$ is satisfied for all $m \geq \bar{C}_{1}(-\ln \delta)$, with a sufficiently large $C_{1}$.

Step 3: Proof of Theorem 2.2.2. The following proposition shows that for any given $\theta=(A, B)^{\top} \in \mathbb{R}^{(n+d) \times n}$, the full row rank of $K^{\theta}$ is equivalent to the well-definedness of (2.2.12) for all $\theta^{\prime}$ sufficiently close to $\theta$.

Proposition 2.3.10. Suppose (H.1(1)) holds. For each $\theta \in \mathbb{R}^{(n+d) \times n}$, let $V^{\psi^{\theta}}$ be defined in (2.3.1). Then for any $\theta=(A, B)^{\top} \in \mathbb{R}^{(n+d) \times n}$, the following properties are equivalent:
(1) $\left\{v \in \mathbb{R}^{d} \mid\left(K_{t}^{\theta}\right)^{\top} v=0, \forall t \in[0, T]\right\}=\{0\}$, with $K^{\theta}$ defined in (2.2.9);
(2) $\mathbb{E}\left[V^{\psi^{\theta}}\right] \in \mathbb{S}_{+}^{n+d}$;
(3) there exist $\lambda_{0}, \varepsilon>0$ such that $\lambda_{\min }\left(\mathbb{E}\left[V^{\psi^{\theta^{\prime}}}\right]\right) \geq \lambda_{0}$ for all $\theta^{\prime} \in \Phi_{\varepsilon}:=\left\{\theta^{\prime} \in \mathbb{R}^{(n+d) \times n} \mid\right.$ $\left.\left|\theta^{\prime}-\theta\right| \leq \varepsilon\right\}$, where $\lambda_{\min }(Z)$ is the minimum eigenvalue of $Z \in \mathbb{S}_{0}^{n+d}$.

Proof. For $(1) \Longrightarrow(2)$ : By (2.3.1), $\mathbb{E}\left[V^{\psi^{\theta}}\right] \in \mathbb{S}_{+}^{n+d}$ if and only if there exists no nonzero $v \in \mathbb{R}^{n+d}$ such that

$$
\mathbb{E}\left[\int_{0}^{T} v^{\top} Z_{t}^{\psi^{\theta}}\left(Z_{t}^{\psi^{\theta}}\right)^{\top} v \mathrm{~d} t\right]=\int_{0}^{T} v^{\top}\binom{I}{K_{t}^{\theta}} \mathbb{E}\left[X_{t}^{\psi^{\theta}}\left(X_{t}^{\psi^{\theta}}\right)^{\top}\right]\left(\begin{array}{ll}
I & \left.\left(K_{t}^{\theta}\right)^{\top}\right) v \mathrm{~d} t=0 \tag{2.3.12}
\end{array}\right.
$$

where we applied Fubini's theorem for the first identity. By (2.2.5), $X_{t}^{\psi^{\theta}}=\Phi_{t}^{\theta}\left(x_{0}+\right.$ $\left.\int_{0}^{t}\left(\Phi_{s}^{\theta}\right)^{-1} \mathrm{~d} W_{s}\right)$ for all $t \in[0, T]$, where $\Phi^{\theta} \in C\left([0, T] ; \mathbb{R}^{n \times n}\right)$ is the fundamental solution of $\mathrm{d} \Phi_{t}^{\theta}=\left(A^{\star}+B^{\star} K_{t}^{\theta}\right) \Phi_{t}^{\theta} \mathrm{d} t, K_{t}^{\theta}=-R^{-1} B^{\top} P_{t}^{\theta}$ for all $t \in[0, T]$, and $P^{\theta}$ satisfies (2.2.10). Hence,

$$
\mathbb{E}\left[X_{t}^{\psi^{\theta}}\left(X_{t}^{\psi^{\theta}}\right)^{\top}\right]=\Phi_{t}^{\theta}\left(x_{0} x_{0}^{\top}+\int_{0}^{t}\left(\Phi_{s}^{\theta}\right)^{-1}\left(\left(\Phi_{s}^{\theta}\right)^{-1}\right)^{\top} \mathrm{d} s\right)\left(\Phi_{t}^{\theta}\right)^{\top} \in \mathbb{S}_{0}^{n}, \quad \forall t \in[0, T] .
$$

Then by (2.3.12) and the continuity of $t \mapsto \mathbb{E}\left[X_{t}^{\psi^{\theta}}\left(X_{t}^{\psi^{\theta}}\right)^{\top}\right]$ and $t \mapsto K_{t}^{\theta}, \mathbb{E}\left[V^{\psi^{\theta}}\right] \in \mathbb{S}_{+}^{n+d}$ if and only if there exists no nonzero $v \in \mathbb{R}^{n+d}$ such that
where $I$ is the $n \times n$ identity matrix. One can easily deduce by the invertibility of $\left(\Phi_{t}^{\theta}\right)^{-1}$ for all $t \in[0, T]$ that $\int_{0}^{t}\left(\Phi_{s}^{\theta}\right)^{-1}\left(\left(\Phi_{s}^{\theta}\right)^{-1}\right)^{\top} \mathrm{d} s \in \mathbb{S}_{+}^{n}$ for all $t>0$, which subsequently shows that $\mathbb{E}\left[V^{\psi^{\theta}}\right] \in \mathbb{S}_{+}^{n+d}$ if and only if there exists no nonzero $\tilde{v} \in \mathbb{R}^{n+d}$ such that $\left(I \quad\left(K_{t}^{\theta}\right)^{\top}\right) \tilde{v}=0$ for all $t \in[0, T]$. Now let us denote without loss of generality that $\tilde{v}=\binom{u}{v}$ for some $u \in \mathbb{R}^{n}$ and $v \in \mathbb{R}^{d}$. Then the above derivation shows that $\mathbb{E}\left[V^{\psi^{\theta}}\right] \in \mathbb{S}_{+}^{n+d}$ is equivalent to the following statement:
if $u \in \mathbb{R}^{n}$ and $v \in \mathbb{R}^{d}$ satisfy $u+\left(K_{t}^{\theta}\right)^{\top} v=0$ for all $t \in[0, T]$, then $u=0$ and $v=0$.
By (2.2.9), $\left(K_{t}^{\theta}\right)^{\top}=-P_{t}^{\theta} B R^{-1}$ for all $t \in[0, T]$ and $P_{T}^{\theta}=0$, implying that $K_{T}^{\theta}=0$. Then (2.3.13) can be rewritten as:

$$
\text { if } v \in \mathbb{R}^{d} \text { satisfies }\left(K_{t}^{\theta}\right)^{\top} v=0 \text { for all } t \in[0, T] \text {, then } v=0
$$

For $(2) \Longleftrightarrow(3)$ : Item (3) clearly implies Item (2). On the other hand, for any given $\theta, \theta^{\prime} \in \mathbb{R}^{(n+d) \times n}$,

$$
\mathrm{d}\left(X^{\psi^{\theta}}-X^{\psi^{\theta^{\prime}}}\right)_{t}=\left(\left(A^{\star}+B^{\star} K_{t}^{\theta}\right)\left(X^{\psi^{\theta}}-X^{\psi^{\theta^{\prime}}}\right)_{t}+B^{\star}\left(K_{t}^{\theta}-K_{t}^{\theta^{\prime}}\right) X_{t}^{\psi^{\theta^{\prime}}}\right) \mathrm{d} t .
$$

Then, we can easily deduce from the continuity of $t \mapsto K^{\theta}$ (see Lemma 2.3.1) that $\mathbb{R}^{(n+d) \times n} \ni$ $\theta \mapsto Z^{\psi^{\theta}} \in \mathcal{H}^{2}\left(\mathbb{R}^{(n+d) \times n}\right)$ is continuous, which implies the continuity of $\mathbb{R}^{(n+d) \times n} \ni \theta \mapsto$ $V^{\psi^{\theta}}=\mathbb{E}\left[\int_{0}^{T} Z_{t}^{\psi^{\theta}}\left(Z_{t}^{\psi^{\theta}}\right)^{\top} \mathrm{d} t\right] \in \mathbb{S}_{0}^{n+d}$. Hence, by the continuity of the minimum eigenvalue function, we can conclude Item (2) from Item (3).

The following proposition provides sufficient conditions for the nondegeneracy of $K^{\theta}$.
Proposition 2.3.11. Let $n, d \in \mathbb{N}, \theta=(A, B)^{\top} \in \mathbb{R}^{(n+d) \times n}, Q \in \mathbb{S}_{0}^{n}$ and $R \in \mathbb{S}_{+}^{d}$.
(1) For all $T>0$, if $B^{\top} Q B \in \mathbb{S}_{+}^{d}$, then $\left\{v \in \mathbb{R}^{d} \mid\left(K_{t}^{\theta}\right)^{\top} v=0, \forall t \in[0, T]\right\}=\{0\}$.
(2) Assume that the algebraic Riccati equation $A^{\top} P+P A-P\left(B R^{-1} B^{\top}\right) P+Q=0$ admits a unique maximal solution $P_{\infty} \in \mathbb{S}_{+}^{n}$. Let $K_{\infty}=-R^{-1} B^{\top} P_{\infty}$, and for each $T>0$, let $P^{(T)} \in C\left([0, T] ; \mathbb{S}_{0}^{n}\right)$ be defined in (2.2.10). Assume that $\lim _{T \rightarrow \infty} P_{0}^{(T)}=P_{\infty}$ and $K_{\infty}\left(K_{\infty}\right)^{\top} \in \mathbb{S}_{+}^{d}$. Then there exists $T_{0}>0$, such that for all $T \geq T_{0},\left\{v \in \mathbb{R}^{d} \mid\right.$ $\left.\left(K_{t}^{\theta}\right)^{\top} v=0, \forall t \in[0, T]\right\}=\{0\}$.

Proof. To prove Item (1), suppose that $B^{\top} Q B \in \mathbb{S}_{+}^{n}$ and $v \in \mathbb{R}^{d}$ such that $\left(K_{t}^{\theta}\right)^{\top} v=$ $-P_{t}^{\theta} B R^{-1} v=0$ for all $t \in[0, T]$, with $P^{\theta}$ defined in (2.2.10). Setting $u=R^{-1} v$, right multiplying (2.2.10) by $B u$, and left multiplying (2.2.10) by $u^{\top} B^{\top}$ shows

$$
u^{\top} B^{\top}\left(\frac{\mathrm{d}}{\mathrm{~d} t} P_{t}^{\theta}\right) B u+A^{\top} P_{t}^{\theta} B u+u^{\top} B^{\top} P_{t}^{\theta} A B u-u^{\top} B^{\top} P_{t}^{\theta} B R^{-1} B^{\top} P_{t}^{\theta} B u+u^{\top} B^{\top} Q B u=0 .
$$

As $P_{t}^{\theta} B u=0$ for all $t \in(0, T), u^{\top} B^{\top}\left(\frac{\mathrm{d}}{\mathrm{d} t} P_{t}^{\theta}\right) B u=u^{\top} B^{\top} P_{t}^{\theta}=0$ for all $t \in(0, T)$, and hence $u^{\top} B^{\top} Q B u=0$. The assumption of $B^{\top} Q B \in \mathbb{S}_{+}^{n}$ then gives $u=R^{-1} v=0$, which along with the invertibility of $R^{-1}$ shows that $v=0$.

To prove Item (2), observe that $\lim _{T \rightarrow \infty}\left(-R^{-1} B^{\top} P_{0}^{(T)}\right)=K_{\infty}$. As $\lambda_{\min }\left(K_{\infty}\left(K_{\infty}\right)^{\top}\right)>0$, there exists $T_{0}>0$ such that for all $T \geq T_{0}, \lambda_{\min }\left(\left(-R^{-1} B^{\top} P_{0}^{(T)}\right)\left(-R^{-1} B^{\top} P_{0}^{(T)}\right)^{\top}\right)>0$. Fix $T \geq T_{0}$ and consider $v \in \mathbb{R}^{d}$ such that $\left(K_{t}^{\theta}\right)^{\top} v=0$ for all $t \in[0, T]$. Then the definitions of $K^{\theta}$ and $P^{(T)}$ imply the invertibility of $K_{0}^{\theta}\left(K_{0}^{\theta}\right)^{\top}$, which yields $v=\left(K_{0}^{\theta}\left(K_{0}^{\theta}\right)^{\top}\right)^{-1} K_{0}^{\theta}\left(K_{0}^{\theta}\right)^{\top} v=$ 0 .

Now we are ready for the proof of Theorem 2.2.2.
Proof of Theorem 2.2.2. As (H.15) holds with $\theta^{\star}$ and $\theta_{0}$, we can obtain from Proposition 2.3.10 that, there exist $C_{1}, \varepsilon>0$ such that for all $\theta \in \Phi_{\varepsilon}:=\left\{\theta\left|\mathbb{R}^{(n+d) \times n}\right|\left|\theta-\theta^{\star}\right| \leq \varepsilon\right\} \cup\left\{\theta_{0}\right\}$, we have $\left\|\left(\mathbb{E}\left[V^{\psi^{\theta}}\right]\right)^{-1}\right\|_{2} \leq C_{1}$. Then by Proposition 3.3.12, there exist constants $\bar{C}_{1}, \bar{C}_{2} \geq 1$,
such that for all $\theta \in \Theta_{\varepsilon}$ and $\delta^{\prime} \in(0,1 / 2)$, if $m \geq \bar{C}_{1}\left(-\ln \delta^{\prime}\right)$, then with probability at least $1-2 \delta^{\prime}$,

$$
\begin{equation*}
\left|\hat{\theta}-\theta^{\star}\right| \leq \bar{C}_{2}\left(\sqrt{\frac{-\ln \delta^{\prime}}{m}}+\frac{-\ln \delta^{\prime}}{m}+\frac{\left(-\ln \delta^{\prime}\right)^{2}}{m^{2}}\right) \tag{2.3.14}
\end{equation*}
$$

where $\hat{\theta}$ denotes the right-hand side of (2.2.13) with the control $\psi^{\theta}$. In the following, we fix $\delta \in\left(0,3 / \pi^{2}\right)$ and $C \geq C_{0}$, with the constant $C_{0} \in(0, \infty)$ satisfying

$$
C_{0} \geq \bar{C}_{1}\left(\sup _{\ell \in \mathbb{N} \cup\{0\}, \delta \in\left(0,3 / \pi^{2}\right)} \frac{-\ln \left(\delta /(\ell+1)^{2}\right)}{2^{\ell}(-\ln \delta)}\right) / \min \left\{\left(\frac{\varepsilon}{3 \bar{C}_{2}}\right)^{2}, 1\right\}
$$

let $m_{0}=C(-\ln \delta)$, and for each $\ell \in \mathbb{N} \cup\{0\}$, let $\delta_{\ell}=\delta /(\ell+1)^{2}$, $m_{\ell}=2^{\ell} m_{0}$, and let $\theta_{\ell+1}$ be generated by (2.2.13) with $m=m_{\ell}$ and $\theta=\theta_{\ell}$. Note that the choices of $C_{0}, m_{\ell}, \delta_{\ell}$ ensure that $m_{\ell} \geq \bar{C}_{1}\left(-\ln \delta_{\ell}\right)$, and

$$
\begin{equation*}
\bar{C}_{2}\left(\sqrt{\frac{-\ln \delta_{\ell}}{m_{\ell}}}+\frac{\left(-\ln \delta_{\ell}\right)}{m_{\ell}}+\frac{\left(-\ln \delta_{\ell}\right)^{2}}{m_{\ell}^{2}}\right) \leq 3 \bar{C}_{2} \sqrt{\frac{-\ln \delta_{\ell}}{m_{\ell}}} \leq \varepsilon, \quad \forall \ell \in \mathbb{N} \cup\{0\} \tag{2.3.15}
\end{equation*}
$$

We now prove with probability at least $1-2 \sum_{\ell=0}^{\infty} \delta_{\ell}=1-\frac{\pi^{2}}{3} \delta$,

$$
\begin{equation*}
\left|\theta_{\ell+1}-\theta^{\star}\right| \leq \bar{C}_{2}\left(\sqrt{\frac{-\ln \delta_{\ell}}{m_{\ell}}}+\frac{\left(-\ln \delta_{\ell}\right)}{m_{\ell}}+\frac{\left(-\ln \delta_{\ell}\right)^{2}}{m_{\ell}^{2}}\right), \quad \forall \ell \in \mathbb{N} \cup\{0\} \tag{2.3.16}
\end{equation*}
$$

Let us consider the induction statement for each $k \in \mathbb{N} \cup\{0\}$ : with probability at least $1-2 \sum_{\ell=0}^{k} \delta_{\ell},(2.3 .16)$ holds for all $\ell=0, \ldots, k$. The fact that $\theta_{0} \in \Theta_{\varepsilon}$ and (2.3.14) yields the induction statement for $k=0$. Now suppose that the induction statement holds for some $k \in \mathbb{N} \cup\{0\}$. Then the induction hypothesis and (3.3.16) ensure that $\left|\theta_{\ell}-\theta^{\star}\right| \leq \varepsilon$ for all $\ell=1, \ldots, k+1$ (and hence $\theta_{k+1} \in \Theta_{\varepsilon}$ ) with probability at least $1-2 \sum_{\ell=0}^{k} \delta_{\ell}$. Conditioning on this event, we can apply (2.3.14) with $\theta=\theta_{k+1}, \delta^{\prime}=\delta_{k+1}<1 / 2$ and $m=m_{k+1} \geq \bar{C}_{1}\left(-\ln \delta_{k+1}\right)$, and deduce with probability at least $1-2 \delta_{k+1}$ that (2.3.16) holds for the index $\ell=k+1$. Combining this with the induction hypothesis yields (2.3.16) holds for the indices $\ell=0, \ldots, k+1$, with probability at least $1-2 \sum_{\ell=0}^{k+1} \delta_{\ell}$.

Observe that for all $i \in \mathbb{N}$, Algorithm 8 generates the $i$-th trajectory with control $\psi^{\theta_{\ell}}$ if $i \in\left(\sum_{j=0}^{\ell-1} m_{j}, \sum_{j=0}^{\ell} m_{j}\right]=\left(m_{0}\left(2^{\ell}-1\right), m_{0}\left(2^{\ell+1}-1\right)\right.$ ] with some $\ell \in \mathbb{N} \cup\{0\}$. Then
conditioning on the event (2.3.16), we can obtain from Proposition 2.3.8 that, for all $M \in \mathbb{N}$,

$$
\begin{align*}
R(M) & \leq \sum_{\ell=0}^{\left\lceil\log _{2}\left(\frac{M}{m_{0}}+1\right)\right\rceil-1} m_{\ell}\left(J^{\theta^{\star}}\left(U^{\psi^{\theta} \ell}\right)-J^{\theta^{\star}}\left(U^{\theta^{\star}}\right)\right) \leq C^{\prime} \sum_{\ell=0}^{\left\lceil\log _{2}\left(\frac{M}{m_{0}}+1\right)\right\rceil-1} m_{\ell}\left|\theta_{\ell}-\theta^{\star}\right|^{2} \\
& \leq C^{\prime} m_{0}+C^{\prime} \sum_{\ell=0}^{\left\lceil\log _{2}\left(\frac{M}{m_{0}}+1\right)\right\rceil-1}\left(-\ln \delta_{\ell}\right)\left(1+\frac{-\ln \delta_{\ell}}{m_{\ell}}+\frac{\left(-\ln \delta_{\ell}\right)^{3}}{m_{\ell}^{3}}\right) \\
& \leq C^{\prime}(-\ln \delta)+C^{\prime} \sum_{\ell=1}^{\left\lceil\log _{2} M\right\rceil}(2 \ln \ell-\ln \delta) \leq C^{\prime}((\ln M)(\ln \ln M)+(\ln M)(-\ln \delta)), \tag{2.3.17}
\end{align*}
$$

with a constant $C^{\prime}$ independent of $M$ and $\delta$, where we have used $\sum_{\ell=1}^{n} \ln \ell=\ln (n!) \leq C^{\prime} n \ln n$ due to Stirling's formula.

### 2.3.4 Regret analysis of discrete-time least-squares algorithm

This section is devoted to the proof of Theorem 2.2.3. The main step is similar to the proof of Theorem 2.2.2 in Section 2.3.3. However, one needs to quantity the precise impact of the piecewise constant policies and the discrete-time observations on the performance gap and the parameter estimation error.

Step 1: Analysis of the performance gap. The following proposition shows the performance gap between applying a piecewise constant feedback control for an incorrect model and a continuous-time feedback control for the true model scales quadratically with respect to the stepsize and the parameter errors.
Proposition 2.3.12. Suppose (H.1(1)) holds and let $\Theta$ be a bounded subset of $\mathbb{R}^{(n+d) \times n}$. For each $\theta \in \Theta$ and $N \in \mathbb{N}$, let $\psi^{\theta, \tau}$ be defined in (2.2.19) with stepsize $\tau=T / N$, let $X^{\psi^{\theta, \tau}}$ be the state process associated with $\psi^{\theta, \tau}$ (cf. (2.2.21)), let $\psi^{\theta^{*}}$ be defined in (2.2.3), and let $X^{\theta^{\star}}$ be the state process associated with $\psi^{\theta^{\star}}(c f .(2.2 .5))$. Then there exists $C>0$ such that

$$
\begin{equation*}
\left|J^{\theta^{\star}}\left(U^{\psi^{\theta, \tau}}\right)-J^{\theta^{\star}}\left(U^{\theta^{\star}}\right)\right| \leq C\left(N^{-2}+\left|\theta-\theta^{\star}\right|^{2}\right), \quad \forall \theta \in \Theta, N \in \mathbb{N} \tag{2.3.18}
\end{equation*}
$$

where $U_{t}^{\psi^{\theta, \tau}}=\psi^{\theta, \tau}\left(t, X_{t}^{\psi^{\theta, \tau}}\right)$ and $U_{t}^{\theta^{\star}}=\psi^{\theta^{\star}}\left(t, X_{t}^{\theta^{\star}}\right)$ for all $t \in[0, T]$, and $J^{\theta^{\star}}$ is defined in (2.2.1).

Proof. Let us fix $\theta \in \Theta$ and $N \in \mathbb{N}$. By applying Proposition 2.3.7 with $U=U^{\psi^{\theta, \tau}}$,

$$
\begin{align*}
& J^{\theta^{\star}}\left(U^{\psi^{\theta, \tau}}\right)-J^{\theta^{\star}}\left(U^{\theta^{\star}}\right) \\
& \leq\|Q\|_{2}\left\|X^{\theta^{\star}, U^{\psi^{\theta, \tau}}}-X^{\theta^{\star}}\right\|_{\mathcal{H}^{2}\left(\mathbb{R}^{n}\right)}^{2}+\|R\|_{2}\left\|U^{\psi^{\theta, \tau}}-U^{\theta^{\star}}\right\|_{\mathcal{H}^{2}\left(\mathbb{R}^{d}\right)}^{2}  \tag{2.3.19}\\
& \leq\|Q\|_{2}\left\|X^{\psi^{\theta, \tau}}-X^{\psi^{\theta^{\star}}}\right\|_{\mathcal{H}^{2}\left(\mathbb{R}^{n}\right)}^{2}+\|R\|_{2}\left\|\psi^{\theta, \tau}\left(\cdot, X^{\psi^{\theta, \tau}}\right)-\psi^{\theta^{\star}}\left(\cdot, X^{\psi^{\theta^{\star}}}\right)\right\|_{\mathcal{H}^{2}\left(\mathbb{R}^{d}\right)}^{2},
\end{align*}
$$

where the last inequality used the fact that $X^{\theta^{\star}, U^{\psi^{\theta, \tau}}}=X^{\psi^{\theta, \tau}}$ (see (2.2.11)), and the definitions of $U^{\psi^{\theta, \tau}}$ and $U^{\theta^{\star}}$.

We then prove that there exists a constant $C$, independent of $\theta, N$, such that

$$
\left\|X^{\psi^{\theta, \tau}}-X^{\psi^{\theta^{\star}}}\right\|_{\mathcal{H}^{2}\left(\mathbb{R}^{n}\right)}+\left\|\psi^{\theta, \tau}\left(\cdot, X^{\psi^{\theta, \tau}}\right)-\psi^{\theta^{\star}}\left(\cdot, X^{\psi^{\theta^{\star}}}\right)\right\|_{\mathcal{H}^{2}\left(\mathbb{R}^{d}\right)} \leq C\left(N^{-1}+\left|\theta-\theta^{\star}\right|\right)
$$

By setting $\delta X=X^{\theta^{\star}}-X^{\psi^{\theta, \tau}}$, we obtain from (2.2.5) and (2.2.21) that

$$
\begin{equation*}
\mathrm{d} \delta X_{t}=\left(A^{\star} \delta X_{t}+B^{\star} K_{t}^{\theta^{\star}} \delta X_{t}+\left(K_{t}^{\theta^{\star}}-K_{t}^{\theta, \tau}\right) X_{t}^{\psi^{\theta, \tau}}\right) \mathrm{d} t, \quad t \in[0, T] ; \quad \delta X_{0}=0 \tag{2.3.20}
\end{equation*}
$$

Since $\left\|P^{\theta^{\star}}\right\|_{C\left([0, T] ; \mathbb{R}^{n \times n}\right)} \leq C$ and $K_{t}^{\theta^{\star}}=-R^{-1} B^{\top} P_{t}^{\theta^{\star}}$ for all $t \in[0, T],\left\|K^{\theta^{\star}}\right\|_{C\left([0, T] ; \mathbb{R}^{d \times n}\right)} \leq C$. Moreover, by $\left\|P_{t_{i}}^{\theta, \tau}\right\|_{2} \leq C$ for all $i=0, \ldots, N$ (see Proposition 2.3.3) and (2.2.20), we have $\left\|K_{t}^{\theta, \tau}\right\|_{2} \leq C$ for all $t \in[0, T]$, which along with a moment estimate of (2.2.21) yields $\left\|X^{\psi^{\theta, \tau}}\right\|_{\mathcal{S}^{2}\left(\mathbb{R}^{n}\right)} \leq C$. Thus, by applying Gronwall's inequality to (2.3.20), Lemma 2.3.1 and Proposition 2.3.3, for all $\theta \in \Theta$ and $N \in \mathbb{N}$,

$$
\begin{align*}
& \left\|X^{\theta^{\star}}-X^{\psi^{\theta, \tau}}\right\|_{\mathcal{H}^{2}\left(\mathbb{R}^{n}\right)} \leq C\left\|X^{\theta^{\star}}-X^{\psi^{\theta, \tau}}\right\|_{\mathcal{S}^{2}\left(\mathbb{R}^{n}\right)} \\
& \leq C\left\|\left(K_{t}^{\theta^{\star}}-K_{t}^{\theta, \tau}\right) X^{\psi^{\theta, \tau}}\right\|_{\mathcal{H}^{2}\left(\mathbb{R}^{d}\right)} \leq C \max _{t \in[0, T]}\left\|K_{t}^{\theta^{\star}}-K_{t}^{\theta, \tau}\right\|_{2}  \tag{2.3.21}\\
& \leq C \max _{t \in[0, T]}\left(\left\|K_{t}^{\theta}-K_{t}^{\theta, \tau}\right\|_{2}+\left\|K_{t}^{\theta^{\star}}-K_{t}^{\theta}\right\|_{2}\right) \leq C\left(N^{-1}+\left|\theta-\theta^{\star}\right|\right)
\end{align*}
$$

The above inequality further implies

$$
\begin{aligned}
& \left\|\psi^{\theta, \tau}\left(\cdot, X_{\cdot}^{\psi^{\theta, \tau}}\right)-\psi^{\theta^{\star}}\left(\cdot, X^{\psi^{\theta^{\star}}}\right)\right\|_{\mathcal{H}^{2}\left(\mathbb{R}^{d}\right)}=\left\|K_{\cdot}^{\theta, \tau} X^{\psi^{\theta, \tau}}-K_{\cdot}^{\theta^{\star}} X_{\cdot}^{\psi^{\theta^{\star}}}\right\|_{\mathcal{H}^{2}\left(\mathbb{R}^{d}\right)} \\
& \leq\left\|\left(K^{\theta, \tau}-K_{\cdot}^{\theta^{\star}}\right) X^{\psi^{\theta, \tau}}\right\|_{\mathcal{H}^{2}\left(\mathbb{R}^{d}\right)}+\left\|K_{\cdot}^{\theta^{\star}}\left(X_{\cdot}^{\psi^{\theta, \tau}}-X_{\psi^{\theta^{\star}}}\right)\right\|_{\mathcal{H}^{2}\left(\mathbb{R}^{d}\right)} \\
& \leq\left\|K^{\theta^{\star}}-K^{\theta, \tau}\right\|_{C\left(\left[0, T ; \mathbb{R}^{d \times n}\right)\right.}\left\|X^{\psi^{\theta, \tau}}\right\|_{\mathcal{H}^{2}\left(\mathbb{R}^{n}\right)}+\left\|K^{\theta^{\star}}\right\|_{C\left(\left[0, T ; \mathbb{R}^{d \times n}\right)\right.}\left\|X^{\psi^{\theta, \tau}}-X^{\psi^{\theta^{\star}}}\right\|_{\mathcal{H}^{2}\left(\mathbb{R}^{n}\right)} \\
& \leq C\left(N^{-1}+\left|\theta-\theta^{\star}\right|\right), \quad \forall \theta \in \Theta, N \in \mathbb{N},
\end{aligned}
$$

which along with (2.3.19) finishes the desired estimate.

Step 2: Error bound for parameter estimation. The following lemma shows that the difference between the expectations of ( $V^{\psi^{\theta, \tau}, \tau}, Y^{\psi^{\theta, \tau}, \tau}$ ) and of ( $V^{\psi^{\theta, \tau}}, Y^{\psi^{\theta, \tau}}$ ) scales linearly with respect to the stepsize.

Lemma 2.3.13. Suppose (H.1(1)) holds and let $\Theta$ be a bounded subset of $\mathbb{R}^{(n+d) \times n}$. For each $\theta \in \Theta$ and $N \in \mathbb{N}$, let $\tau=T / N$, let $\psi^{\theta, \tau}$ be defined in (2.2.19), let $V^{\psi^{\theta, \tau}}, Y^{\psi^{\theta, \tau}}$ be defined in (2.3.1), and let $V^{\psi^{\theta, \tau}, \tau}, Y^{\psi^{\theta, \tau}, \tau}$ be defined in (2.3.2). Then there exists a constant $C$ such that

$$
\left|\mathbb{E}\left[V^{\psi^{\theta, \tau}, \tau}-V^{\psi^{\theta, \tau}}\right]\right|+\left|\mathbb{E}\left[Y^{\psi^{\theta, \tau}, \tau}-Y^{\psi^{\theta, \tau}}\right]\right| \leq C N^{-1}, \quad \forall \theta \in \Theta, N \in \mathbb{N} .
$$

Proof. By (2.2.21), we have for all $i=0, \ldots, N-1, X_{t_{i+1}}^{\psi^{\theta, \tau}}-X_{t_{i}}^{\psi^{\theta, \tau}}=\int_{t_{i}}^{t_{i+1}}\left(\theta^{\star}\right)^{\top} Z_{t}^{\psi^{\theta, \tau}} \mathrm{d} t+$ $W_{t_{i+1}}-W_{t_{i}}$, which implies

$$
\begin{aligned}
& \mathbb{E}\left[V^{\psi^{\theta, \tau}}-V^{\psi^{\theta, \tau}, \tau}\right]=\sum_{i=0}^{N-1} \int_{t_{i}}^{t_{i+1}} \mathbb{E}\left[Z_{t}^{\psi^{\theta, \tau}}\left(Z_{t}^{\psi^{\theta, \tau}}\right)^{\top}-Z_{t_{i}}^{\psi^{\theta, \tau}}\left(Z_{t_{i}}^{\psi^{\theta, \tau}}\right)^{\top}\right] \mathrm{d} t, \\
& \mathbb{E}\left[Y^{\psi^{\theta, \tau}}-Y^{\psi^{\theta, \tau}, \tau}\right]=\sum_{i=0}^{N-1} \int_{t_{i}}^{t_{i+1}} \mathbb{E}\left[\left(Z_{t}^{\psi^{\theta, \tau}}-Z_{t_{i}}^{\psi^{\theta, \tau}}\right)\left(Z_{t}^{\psi^{\theta, \tau}}\right)^{\top} \theta^{\star}\right] \mathrm{d} t .
\end{aligned}
$$

Hence it suffices to prove that $\left|\mathbb{E}\left[Z_{t}^{\psi^{\theta, \tau}}\left(Z_{t}^{\psi^{\theta, \tau}}\right)^{\top}-Z_{t_{i}}^{\psi^{\theta, \tau}}\left(Z_{t_{i}}^{\psi^{\theta, \tau}}\right)^{\top}\right]\right| \leq C N^{-1}$ and $\mid \mathbb{E}\left[\left(Z_{t}^{\psi^{\theta, \tau}}-\right.\right.$ $\left.\left.Z_{t_{i}}^{\psi^{\theta, \tau}}\right)\left(Z_{t}^{\psi^{\theta, \tau}}\right)^{\top}\right] \mid \leq C N^{-1}$ for all $t \in\left[t_{i}, t_{i+1}\right]$ and $i=0, \ldots, N-1$.

Let us fix $i=0, \ldots, N-1$ and $t \in\left[t_{i}, t_{i+1}\right]$. In the following, we shall omit the superscripts of $X^{\psi^{\theta, \tau}}$ and $Z^{\psi^{\theta, \tau}}$ if no confusion occurs. As $t \in\left[t_{i}, t_{i+1}\right]$, by (2.2.21), we have $X_{t}=$ $e^{L t} X_{t_{i}}+\int_{t_{i}}^{t} L^{L(t-s)} \mathrm{d} W_{s}$ with $L:=A^{\star}+B^{\star} K_{t_{i}}^{\theta, \tau}$. Thus,

$$
\begin{aligned}
& X_{t} X_{t}^{\top}-X_{t_{i}} X_{t_{i}}^{\top} \\
& =\left(X_{t}-X_{t_{i}}+X_{t_{i}}\right)\left(X_{t}-X_{t_{i}}+X_{t_{i}}\right)^{\top}-X_{t_{i}} X_{t_{i}}^{\top} \\
& =\left(X_{t}-X_{t_{i}}\right)\left(X_{t}-X_{t_{i}}\right)^{\top}+X_{t_{i}}\left(X_{t}-X_{t_{i}}\right)^{\top}+\left(X_{t}-X_{t_{i}}\right) X_{t_{i}}^{\top} \\
& =\left(\left(e^{L t}-I\right) X_{t_{i}}+\int_{t_{i}}^{t} e^{L(t-s)} \mathrm{d} W_{s}\right)\left(\left(e^{L t}-I\right) X_{t_{i}}+\int_{t_{i}}^{t} e^{L(t-s)} \mathrm{d} W_{s}\right)^{\top} \\
& \quad+X_{t_{i}}\left(\left(e^{L t}-I\right) X_{t_{i}}+\int_{t_{i}}^{t} e^{L(t-s)} \mathrm{d} W_{s}\right)^{\top}+\left(\left(e^{L t}-I\right) X_{t_{i}}+\int_{t_{i}}^{t} e^{L(t-s)} \mathrm{d} W_{s}\right) X_{t_{i}}^{\top} .
\end{aligned}
$$

By taking expectations of both sides of the above identity, the martingale property of the Itô integral, and the Itô isometry,

$$
\begin{aligned}
\mathbb{E}\left[X_{t} X_{t}^{\top}-X_{t_{i}} X_{t_{i}}^{\top}\right]= & \left(e^{L t}-I\right) \mathbb{E}\left[X_{t_{i}} X_{t_{i}}^{\top}\right]\left(e^{L^{\top} t}-I\right)+\int_{t_{i}}^{t} e^{L(t-s)} e^{L^{\top}(t-s)} \mathrm{d} s \\
& +\mathbb{E}\left[X_{t_{i}} X_{t_{i}}^{\top}\right]\left(e^{L^{\top} t}-I\right)+\left(e^{L t}-I\right) \mathbb{E}\left[X_{t_{i}} X_{t_{i}}^{\top}\right] \leq C\left(t-t_{i}\right),
\end{aligned}
$$

where the last inequality follows from $\left\|X^{\psi^{\theta, \tau}}\right\|_{\mathcal{S}^{2}\left(\mathbb{R}^{n}\right)} \leq C$. Since $\psi^{\theta, \tau}\left(t, X_{t}^{\psi^{\theta, \tau}}\right)=K_{t_{i}}^{\theta, \tau} X_{t}^{\psi^{\theta, \tau}}$ and $\left\|K_{t_{i}}^{\theta, \tau}\right\|_{2} \leq C$, one can easily show that $\left|\mathbb{E}\left[Z_{t}^{\psi^{\theta, \tau}}\left(Z_{t}^{\psi^{\theta, \tau}}\right)^{\top}-Z_{t_{i}}^{\psi^{\theta, \tau}}\left(Z_{t_{i}}^{\psi^{\theta, \tau}}\right)^{\top}\right]\right| \leq C N^{-1}$. Furthermore, by $X_{t}^{\psi^{\theta, \tau}}=e^{L t} X_{t_{i}}^{\psi^{\theta, \tau}}+\int_{t_{i}}^{t} e^{L(t-s)} \mathrm{d} W_{s}$ and the identity

$$
Z_{t}^{\psi^{\theta, \tau}}\left(Z_{t}^{\psi^{\theta, \tau}}\right)^{\top}-Z_{t_{i}}^{\psi^{\theta, \tau}}\left(Z_{t_{i}}^{\psi^{\theta, \tau}}\right)^{\top}=\left(Z_{t}^{\psi^{\theta, \tau}}-Z_{t_{i}}^{\psi^{\theta, \tau}}\right)\left(Z_{t}^{\psi^{\theta, \tau}}\right)^{\top}+Z_{t_{i}}^{\psi^{\theta, \tau}}\left(Z_{t}^{\psi^{\theta, \tau}}-Z_{t_{i}}^{\psi^{\theta, \tau}}\right)^{\top},
$$

we can show that

$$
\begin{aligned}
& \left|\mathbb{E}\left[\left(Z_{t}^{\psi^{\theta, \tau}}-Z_{t_{i}}^{\psi^{\theta, \tau}}\right)\left(Z_{t}^{\psi^{\theta, \tau}}\right)^{\top}\right]\right| \\
& \leq\left|\mathbb{E}\left[Z_{t}^{\psi^{\theta, \tau}}\left(Z_{t}^{\psi^{\theta, \tau}}\right)^{\top}-Z_{t_{i}}^{\psi^{\theta, \tau}}\left(Z_{t_{i}}^{\psi^{\theta, \tau}}\right)^{\top}\right]\right|+\left|\mathbb{E}\left[Z_{t_{i}}^{\psi^{\theta, \tau}}\left(Z_{t}^{\psi^{\theta, \tau}}-Z_{t_{i}}^{\psi^{\theta, \tau}}\right)^{\top}\right]\right| \\
& \leq C\left(N^{-1}+\left|\mathbb{E}\left[Z_{t_{i}}^{\psi^{\theta, \tau}}\left(X_{t_{i}}^{\psi^{\theta, \tau}}\right)^{\top}\left(e^{L^{\top} t}-I\right)\left(I \quad\left(K_{t_{i}}^{\theta, \tau}\right)^{\top}\right)\right]\right|\right) \leq C N^{-1},
\end{aligned}
$$

by the uniform boundedness of $\left\|X^{\psi^{\theta, \tau}}\right\|_{\mathcal{S}^{2}\left(\mathbb{R}^{n}\right)}$ and $K^{\theta, \tau}$.
Proposition 2.3.14. Suppose (H.1(1)) holds, and let $\Theta \subset \mathbb{R}^{(n+d) \times n}$ such that there exists $C_{1}>0$ satisfying $\left\|\left(\mathbb{E}\left[V^{\psi^{\theta}}\right]\right)^{-1}\right\|_{2} \leq C_{1}$ and $|\theta| \leq C_{1}$ for all $\theta \in \Theta$, with $V^{\psi^{\theta}}$ defined in (2.3.1). Then there exist constants $\bar{C}_{1}, \bar{C}_{2} \geq 0$ and $n_{0} \in \mathbb{N}$, such that for all $\theta \in \Theta, N \in \mathbb{N} \cap\left[n_{0}, \infty\right)$ and $\delta \in(0,1 / 2)$, if $m \geq \bar{C}_{1}(-\ln \delta)$, then with probability at least $1-2 \delta$,

$$
\begin{equation*}
\left|\hat{\theta}-\theta^{\star}\right| \leq \bar{C}_{2}\left(\sqrt{\frac{-\ln \delta}{m}}+\frac{-\ln \delta}{m}+\frac{(-\ln \delta)^{2}}{m^{2}}+\frac{1}{N}\right) \tag{2.3.22}
\end{equation*}
$$

where $\hat{\theta}$ denotes the right-hand side of (2.2.23) with the control $\psi^{\theta, \tau}$ and stepsize $\tau=T / N$.
Proof. We first prove that there exists $n_{0} \in \mathbb{N}$ such that for all $N \in \mathbb{N} \cap\left[n_{0}, \infty\right)$ and $\theta \in \Theta$, $\left\|\left(\mathbb{E}\left[V^{\psi^{\theta, \tau}}\right]\right)^{-1}\right\|_{2} \leq C$ for a constant $C>0$ independent of $\theta$ and $N$. By (2.2.11) and (2.2.21), we have for all $\theta \in \Theta$ and $N \in \mathbb{N}, X_{0}^{\psi^{\theta}}=X_{0}^{\psi^{\theta}, \tau}$ and

$$
\mathrm{d}\left(X^{\psi^{\theta}}-X^{\psi^{\theta}, \tau}\right)_{t}=\left(\left(A^{\star}+B^{\star} K_{t}^{\theta}\right)\left(X^{\psi^{\theta}}-X^{\psi^{\theta}, \tau}\right)_{t}+B^{\star}\left(K_{t}^{\theta}-K_{t}^{\theta, \tau}\right) X_{t}^{\psi^{\theta}, \tau}\right) \mathrm{d} t, \quad t \in[0, T] .
$$

Proposition 2.3.3 shows $\left\|K_{t}^{\theta}-K_{t}^{\theta, \tau}\right\|_{2} \leq C N^{-1}$ for all $t \in[0, T]$, which along with Gronwall's inequality yields $\left\|X^{\psi^{\theta}}-X^{\psi^{\theta}, \tau}\right\|_{\mathcal{S}^{2}\left(\mathbb{R}^{n}\right)} \leq C N^{-1}$ for all $\theta \in \Theta$ and $N \in \mathbb{N}$. One can further prove that $\left\|U^{\psi^{\theta}}-U^{\psi^{\theta}, \tau}\right\|_{\mathcal{S}^{2}\left(\mathbb{R}^{d}\right)} \leq C N^{-1}$ with $U_{t}^{\psi^{\theta}}=\psi^{\theta}\left(t, X_{t}^{\psi^{\theta}}\right)$ and $U_{t}^{\psi^{\theta, \tau}}=\psi^{\theta, \tau}\left(t, X_{t}^{\psi^{\theta, \tau}}\right)$ for all $t \in[0, T]$. Thus, we have $\left|\mathbb{E}\left[V^{\psi^{\theta}}\right]-\mathbb{E}\left[V^{\psi^{\theta, \tau}}\right]\right| \leq C N^{-1}$, which along with $\left\|\left(\mathbb{E}\left[V^{\psi^{\theta}}\right]\right)^{-1}\right\|_{2} \leq$ $C_{1}$ implies a uniform bound of $\left\|\left(\mathbb{E}\left[V^{\psi^{\theta, \tau}}\right]\right)^{-1}\right\|_{2}$ for all sufficiently large $N$.

Let us fix $N \in \mathbb{N} \cap\left[n_{0}, \infty\right)$ and $\theta \in \Theta$ for the subsequent analysis. The invertibility of $\mathbb{E}\left[V^{\psi^{\theta, \tau}}\right]$ implies that $\left.\theta^{\star}=\left(\mathbb{E}\left[V^{\psi^{\theta, \tau}}\right]\right)^{-1} \mathbb{E}\left[Y^{\psi^{\theta, \tau}}\right]\right)(c f .(2.2 .12))$. Then by (2.2.23), we can derive the following analogues of (2.3.10):

$$
\begin{aligned}
& \left\|\hat{\theta}-\theta^{\star}\right\|_{2} \\
& =\left\|\left(V^{\psi^{\theta, \tau}, \tau, m}+\frac{1}{m} I\right)^{-1} Y^{\psi^{\theta, \tau}, \tau, m}-\left(\mathbb{E}\left[V^{\psi^{\theta, \tau}}\right]\right)^{-1} \mathbb{E}\left[Y^{\psi^{\theta, \tau}}\right]\right\|_{2} \\
& \leq\left\|\left(V^{\psi^{\theta, \tau}, \tau, m}+\frac{1}{m} I\right)^{-1}-\left(\mathbb{E}\left[V^{\psi^{\theta, \tau}}\right]\right)^{-1}\right\|_{2}\left\|Y^{\psi^{\theta, \tau}, m, \tau}\right\|_{2}+\left\|\left(\mathbb{E}\left[V^{\psi^{\theta, \tau}}\right]\right)^{-1}\right\|_{2}\left\|Y^{\psi^{\theta, \tau, \tau}, m}-\mathbb{E}\left[Y^{\psi^{\theta, \tau}}\right]\right\|_{2} \\
& \leq\left\|\left(V^{\psi^{\theta, \tau}, \tau, m}+\frac{1}{m} I\right)^{-1}\right\|\left\|\left(\mathbb{E}\left[V^{\psi^{\theta, \tau}}\right]\right)^{-1}\right\|_{2}\left\|Y^{\psi^{\theta, \tau}, m, \tau}\right\|_{2}\left\|V^{\psi^{\theta, \tau}, \tau, m}-\mathbb{E}\left[V^{\psi^{\theta, \tau}}\right]+\frac{1}{m} I\right\|_{2} \\
& \quad+\left\|\left(\mathbb{E}\left[V^{\psi^{\theta, \tau}}\right]\right)^{-1}\right\|_{2}\left\|Y^{\psi^{\theta, \tau, \tau}, m}-\mathbb{E}\left[Y^{\psi^{\theta, \tau}}\right]\right\|_{2},
\end{aligned}
$$

where $V^{\psi^{\theta, \tau}, \tau, m}$ and $Y^{\psi^{\theta, \tau}, \tau, m}$ are defined in (2.3.2). Note that

$$
\begin{aligned}
\left\|V^{\psi^{\theta, \tau}, \tau, m}-\mathbb{E}\left[V^{\psi^{\theta, \tau}}\right]+\frac{1}{m} I\right\|_{2} & \leq\left\|V^{\psi^{\theta, \tau}, \tau, m}-\mathbb{E}\left[V^{\psi^{\theta, \tau}, \tau}\right]\right\|_{2}+\left\|\mathbb{E}\left[V^{\psi^{\theta, \tau}, \tau}\right]-\mathbb{E}\left[V^{\psi^{\theta, \tau}}\right]\right\|_{2}+\frac{1}{m}, \\
\left\|Y^{\psi^{\theta, \tau, \tau}, m}-\mathbb{E}\left[Y^{\psi^{\theta, \tau}}\right]\right\|_{2} & \leq\left\|Y^{\psi^{\theta, \tau, \tau}, m}-\mathbb{E}\left[Y^{\psi^{\theta, \tau, \tau}}\right]\right\|_{2}+\left\|\mathbb{E}\left[Y^{\psi^{\theta, \tau, \tau}}\right]-\mathbb{E}\left[Y^{\psi^{\theta, \tau}}\right]\right\|_{2},
\end{aligned}
$$

where for both inequalities, the first term on the right-hand side can be estimated by Theorem 2.3.6 (uniformly in $N$ ), and the second term is of the magnitude $\mathcal{O}\left(N^{-1}\right)$ due to Lemma 2.3.13. Hence, proceeding along the lines of the proof of Proposition 3.3.12 leads to the desired result.

Step 3: Proof of Theorem 2.2.3. The proof follows from similar arguments as that of Theorem 2.2.2, and we only present the main steps here. As (H.15) holds with $\theta_{0}$ and $\theta^{\star}$, we can obtain from Propositions 2.3.10 and 2.3.14 that, there exists a bounded set $\Phi_{\varepsilon} \subset \mathbb{R}^{(n+d) \times n}$ and constants $\bar{C}_{1}, \bar{C}_{2} \geq 1, n_{0} \in \mathbb{N}$ that for all $\theta \in \Phi_{\varepsilon}, N \in \mathbb{N} \cap\left[n_{0}, \infty\right)$ and $\delta^{\prime} \in(0,1 / 2)$, if $m \geq \bar{C}_{1}(-\ln \delta)$, then with probability at least $1-2 \delta^{\prime}$,

$$
\begin{equation*}
\left|\hat{\theta}-\theta^{\star}\right| \leq \bar{C}_{2}\left(\sqrt{\frac{-\ln \delta^{\prime}}{m}}+\frac{-\ln \delta^{\prime}}{m}+\frac{\left(-\ln \delta^{\prime}\right)^{2}}{m^{2}}+\frac{1}{N}\right) \tag{2.3.23}
\end{equation*}
$$

where $\hat{\theta}$ denotes the right-hand side of (2.2.23) with the control $\psi^{\theta, \tau}$ and stepsize $\tau=T / N$. Then by proceeding along the lines of the proof of Theorem 2.2.2, there exists $C_{0}>0$ and $n_{0} \in \mathbb{N}$, such that for any given $\delta \in\left(0, \frac{3}{\pi^{2}}\right)$, if $m_{0}=C(-\ln \delta)$ with $C \geq C_{0}$ and $N_{\ell} \geq n_{0}$ for all $\ell \in \mathbb{N} \cup\{0\}$, then with probability at least $1-\frac{\pi^{2}}{3} \delta$,

$$
\begin{equation*}
\left|\theta_{\ell+1}-\theta^{\star}\right| \leq \bar{C}_{2}\left(\sqrt{\frac{-\ln \delta_{\ell}}{m_{\ell}}}+\frac{\left(-\ln \delta_{\ell}\right)}{m_{\ell}}+\frac{\left(-\ln \delta_{\ell}\right)^{2}}{m_{\ell}^{2}}+\frac{1}{N_{\ell}}\right), \quad \forall \ell \in \mathbb{N} \cup\{0\} \tag{2.3.24}
\end{equation*}
$$

where $\delta_{\ell}=\delta /(\ell+1)^{2}$ and $m_{\ell}=2^{\ell} m_{0}$ for all $\ell$. Consequently, we can conclude the desired regret bound from Proposition 2.3 .12 (cf. (2.3.17)), with an additional term $\sum_{\ell=0}^{\ln M} m_{\ell} N_{\ell}^{-2}$ due to the time discretization errors in (2.3.18) and (2.3.24).

## Chapter 3

## Reinforcement learning for linear-convex models with jumps via stability analysis of feedback controls

### 3.1 Introduction

Reinforcement learning (RL) seeks optimal strategies to control an unknown dynamical system by interacting with the random environment through exploration and exploitation [147]. This paper studies a reinforcement learning problem for controlled linear-convex models with unknown drift parameters. The controlled dynamics are with possible jumps, the objectives are extended real-valued nonsmooth convex functions, and the learning is in an episodic setting for a finite-time horizon.

Regret analysis of RL algorithm and stability of controls. RL algorithms are in general characterized by iterations of exploitation and exploration (see e.g. [1, 112, 17]). In the model-based approach, for instance, the agent interacts with the environment via policies based on the present estimation of the unknown model parameters, and then incorporates the responses of these interactions to improve their knowledge of the system. One of the main performance criteria for RL algorithm, called regret, is to measure its deviation from the optimality over the learning process.

One key component in regret analysis is the Lipschitz stability of feedback controls which quantifies the mismatch between the assumed and actual models, or the stability of controls with respect to model perturbations. It is to analyze the precise derivation of a pre-computed feedback control from the optimal one, and is also known as the robustness of control policies in the learning community $[112,17,94]$ ).

Despite the long history of stability of controls in the control literature, its main focus in classical control theory has been restricted to the continuity of value functions and optimal open-loop controls (see e.g. [6, 169, 15, 18, 94]). Studies of high-order stability of controls
such as the Lipschitz stability, has only attracted attention very recently, largely due to its crucial importance in characterizing the precise regret order of learning algorithms (see [112, 17, 135]). Analyzing Lipschitz stability of feedback control is technically more challenging. It requires analyzing the derivatives of the value function in a suitable function space, as optimal feedback controls are usually characterized via the derivatives of the value function.

Due to this technical difficulty, most existing works on regret analysis of RL algorithms concentrate on the linear-quadratic (LQ) control framework. In this special setting, the optimal feedback control is an affine function of state variables, whose coefficients satisfy an associated algebraic or ordinary Riccati equation. Consequently, the Lipschitz stability of feedback controls is simplified by analyzing the robustness of the Riccati equation (see e.g. $[1,112,17]$ ). Unfortunately, these techniques developed specifically for Riccati equations in LQ-RL problems are clearly not applicable for general RL problems (see e.g. [21, 39, 54, 107]). In particular, optimal policies are typically nonlinear in the state variable, especially with the inclusion of entropy regularization for the exploration strategy in the optimization objective (see e.g. [159, 77, 143, 135]).

Our work. This paper consists of three parts.

- It first establishes the Lipschitz stability for finite-time horizon linear-convex control problems, whose dynamics are linear jump-diffusion processes with controlled drifts and possibly degenerate additive noises, and objectives are extended real-valued lower semicontinuous convex functions. Such control problems include as special cases LQ problems with convex control constraints, sparse and switching control of linear systems, and entropyregularized relaxed control problems (see Examples 3.2.1 and 3.2.2). It shows that these control problems admit Lipchitz continuous optimal feedback controls with linear growth in the spatial variables (Theorem 3.2.5). It further proves that the performance gap between applying feedback controls for an incorrect model and for the true model depends Lipschitz-continuously on the magnitude of perturbations in the model coefficients, even with lower semicontinuous cost functions (Theorem 3.2.7). The Lipschitz stability of feedback controls is extended to entropy-regularized control problems with controlled diffusion in Proposition 3.4.1.
- It then proposes a greedy least-squares (GLS) algorithm for a class of continuous-time linear-convex RL problems in an episodic setting. At each iteration, the GLS algorithm estimates the unknown drift parameters by a regularized least-squares estimator based on observed trajectories, and then designs a feedback control for the estimated model. It establishes that the regret of this GLS algorithm is sublinear, i.e., of the magnitude $\mathcal{O}(\sqrt{N \ln N})$ with $N$ being the number of learning episodes, provided that the least-squares estimator satisfies a general concentration inequality (Theorem 3.3.2). It further characterizes the explicit concentration behaviour of the least-squares estimator (and hence the precise regret bound of the GLS algorithm), depending on tail behaviours of the random
jumps in the state dynamics (Theorem 3.3.3). In the pure diffusion case, a sharper regret bound has been obtained (Theorem 3.3.4).
- It finally verifies the theoretical properties of the proposed GLS algorithm through numerical experiment on a three-dimensional LQ RL problem. It shows the convergence of the least-squares estimations to the true parameters as the number of episodes increases, as well as a sublinear regret as indicated in theoretical results. It also demonstrates the GLS algorithm is robust with respect to initializations.

Our approaches and related works. Optimal control of stochastic systems with parametric uncertainty has been studied in the classical adaptive control literature (see [50, 139, 89, 12]), where stationary policy is constructed to minimize the long term average cost and where the asymptotic stability and convergence of an adaptive control law is analyzed when the time horizon goes to infinity. However, research on rate of convergence is virtually non-existent. The problem studied here is different. The main objective is to construct optimal (and time-dependent) policies for finite-horizon problems, with the finitesample regret analysis for the learning algorithm. Compared with the classical adaptive control literature, the regret analysis in this work, also known as the non-asymptotic performance analysis, requires novel techniques, consisting of a precise performance estimate of a greedy policy (namely the Lipschitz stability of feedback controls) and a finite-sample analysis of the parameter estimation scheme.

Analyzing the Lipschitz stability of feedback controls in a continuous-time setting requires quantifying the impact of parameter uncertainty on the derivatives of the value functions. [135] studies the so-called exit time problem and the Lipschitz stability of regularized relaxed controls of diffusion processes via a partial differential equation (PDE) approach, which assumes that the diffusion coefficients are non-degenerate and the state process takes values in a compact set. In contrast, we consider (see Section 3.2) unconstrained jump-diffusion process with unbounded drift and (uncontrolled) degenerate noise, and the cost functions are nonsmooth and unbounded. Consequently, the PDE approach requires to deal with a degenerate nonlocal PDE with non-Lipschitz nonlinearity, whose solution (i.e., the value function) is unbounded and may be nonsmooth due to the lack of regularization from the Laplacian operator. Here the Lipschitz stability of feedback controls is established by analyzing the stability of the associated coupled forward-backward stochastic differential equations (FBSDEs). This is possible by a) first exploiting the linear-convex structure of the control problem, which enables constructing a Lipschitz continuous feedback control via solutions of coupled FBSDEs, and then b) by extending the stochastic maximum principle in [151] to feedback controls with nonsmooth costs. To the best of our knowledge, this is the first time FBSDE has been used to study stability of feedback controls.

Analyzing the (finite-sample) accuracy of the least-squares estimator for jump-diffusion models involves integrations of the state and control processes with respect to Brownian motions and Poisson random measures. Now, the nonlinearity of feedback controls renders it impossible to analyze the tail behaviour of these stochastic integrals as [77] does for

LQ problems with analytical solutions; Additionally, the presence of random jumps implies that the state process is no longer sub-Gaussian, and hence the stochastic integrals in the least-squares estimator no longer sub-exponential. To overcome these difficulties, a convex concentration inequality is employed for SDEs with jumps [111], along with Burkholder's inequality and the Girsanov theorem to characterize precisely the sub-Weibull behaviour of the required stochastic integrals in terms of their Orlicz norms (Lemmas 3.3.6 and 3.3.7). Leveraging recent developments in the theory of sub-Weibull random variables, the precise parameter estimation error of the least-squares estimator is quantified in terms of the sample size.

It is worth pointing out that the stability analysis of feedback controls can be extended (see Section 3.4) to entropy-regularized control problems with controlled diffusion and without the linear-convex structure. Instead of the maximum principle for the linear-convex setting, regularity analysis of the associated fully-nonlinear parabolic PDEs may be needed for nondegenerate noise with regular (such as bounded and high-order differentiable) coefficients. (See the discussion after Proposition 3.4.1 for more details).

Notation. For each $T>0$, filtered probability space $\left(\Omega, \mathcal{F}, \mathbb{F}=\left\{\mathcal{F}_{t}\right\}_{t \in[0, T]}, \mathbb{P}\right)$ satisfying the usual condition and Euclidean space $(E,|\cdot|)$, we introduce the following spaces:

- $L^{p}(0, T ; E), p \in[2, \infty]$, is the space of (Borel) measurable functions $\phi:[0, T] \rightarrow E$ satisfying $\|\phi\|_{L^{p}}=\left(\int_{0}^{T}\left|\phi_{t}\right|^{p} \mathrm{~d} t\right)^{1 / p}<\infty$ if $p \in[2, \infty)$ and $\|\phi\|_{L^{\infty}}=\operatorname{ess}_{\sup }^{t \in[0, T]}|~| \phi_{t} \mid<\infty$ if $p=\infty$;
- $L^{2}(\Omega ; E)$ is the space of $E$-valued $\mathcal{F}$-measurable random variables $X$ satisfying $\|X\|_{L^{2}}=$ $\mathbb{E}\left[|X|^{2}\right]^{1 / 2}<\infty$;
- $\mathcal{S}^{2}(t, T ; E), t \in[0, T]$, is the space of $E$-valued $\mathbb{F}$-progressively measurable càdlàg processes $Y: \Omega \times[t, T] \rightarrow E$ satisfying $\|Y\|_{\mathcal{S}^{2}}=\mathbb{E}\left[\sup _{s \in[t, T]}\left|Y_{s}\right|^{2}\right]^{1 / 2}<\infty ;$
- $\mathcal{H}^{2}(t, T ; E), t \in[0, T]$, is the space of $E$-valued $\mathbb{F}$-progressively measurable processes $Z$ : $\Omega \times[t, T] \rightarrow E$ satisfying $\|Z\|_{\mathcal{H}^{2}}=\mathbb{E}\left[\int_{t}^{T}\left|Z_{s}\right|^{2} \mathrm{~d} s\right]^{1 / 2}<\infty ;$
- $\mathcal{H}_{\nu}^{2}(t, T ; E), t \in[0, T]$, is the space of $E$-valued $\mathbb{F}$-progressively measurable processes $M: \Omega \times[t, T] \times \mathbb{R}_{0}^{p} \rightarrow E$ satisfying $\|M\|_{\mathcal{H}_{\nu}^{2}}=\mathbb{E}\left[\int_{t}^{T} \int_{\mathbb{R}_{0}^{p}}\left|M_{s}(u)\right|^{2} \nu(\mathrm{~d} u) \mathrm{d} s\right]^{1 / 2}<\infty$, where $\mathbb{R}_{0}^{p}:=\mathbb{R}^{p} \backslash\{0\}$ and $\nu$ is a $\sigma$-finite measure on $\mathbb{R}_{0}^{p}$.
For notational simplicity, we denote $\mathcal{S}^{2}(E)=\mathcal{S}^{2}(0, T ; E), \mathcal{H}^{2}(E)=\mathcal{H}^{2}(0, T ; E)$ and $\mathcal{H}_{\nu}^{2}(E)=$ $\mathcal{H}_{\nu}^{2}(0, T ; E)$. We shall also denote by $\langle\cdot, \cdot\rangle$ the usual inner product in a given Euclidean space, by $|\cdot|$ the norm induced by $\langle\cdot, \cdot\rangle$, by $A^{\top}$ the transpose of a matrix $A$, and by $C \in[0, \infty)$ a generic constant, which depends only on the constants appearing in the assumptions and may take a different value at each occurrence.


### 3.2 Lipschitz stability of linear-convex control problems

### 3.2.1 Problem formulation with nonsmooth costs

In this section, we introduce the linear-convex control problems with nonsmooth costs.
Let $T>0$ be a given terminal time and $(\Omega, \mathcal{F}, \mathbb{P})$ be a complete probability space, in which two mutually independent processes, a $d$-dimensional Brownian motion $W$ and a Poisson random measure $N(\mathrm{~d} t, \mathrm{~d} u)$ with compensator $\nu(\mathrm{d} u) \mathrm{d} t$, are defined. We assume that $\nu$ is a $\sigma$ finite measure on $\mathbb{R}_{0}^{p}$ equipped with its Borel field $\mathcal{B}\left(\mathbb{R}_{0}^{p}\right)$ and satisfies $\int_{\mathbb{R}_{0}^{p}} \min \left(1,|u|^{2}\right) \nu(\mathrm{d} u)<$ $\infty$. We denote by $\tilde{N}(\mathrm{~d} t, \mathrm{~d} u)=N(\mathrm{~d} t, \mathrm{~d} u)-\nu(\mathrm{d} u) \mathrm{d} t$ the compensated process of $N$ and by $\mathbb{F}=\left(\mathcal{F}_{t}\right)_{t \in[0, T]}$ the filtration generated by $W$ and $N$ and augmented by the $\mathbb{P}$-null sets.

For any given initial state $x_{0} \in \mathbb{R}^{n}$, we consider the following minimization problem

$$
\begin{equation*}
V\left(x_{0}\right)=\inf _{\alpha \in \mathcal{H}^{2}\left(\mathbb{R}^{k}\right)} J\left(\alpha ; x_{0}\right), \quad \text { with } \quad J\left(\alpha ; x_{0}\right)=\mathbb{E}\left[\int_{0}^{T} f\left(t, X_{t}^{x_{0}, \alpha}, \alpha_{t}\right) \mathrm{d} t+g\left(X_{T}^{x_{0}, \alpha}\right)\right], \tag{3.2.1}
\end{equation*}
$$

where for each $\alpha \in \mathcal{H}^{2}\left(\mathbb{R}^{k}\right)$, the process $X^{x_{0}, \alpha}$ satisfies the following controlled dynamics:

$$
\begin{equation*}
\mathrm{d} X_{t}=b\left(t, X_{t}, \alpha_{t}\right) \mathrm{d} t+\sigma(t) \mathrm{d} W_{t}+\int_{\mathbb{R}_{0}^{p}} \gamma(t, u) \tilde{N}(\mathrm{~d} t, \mathrm{~d} u), \quad t \in[0, T], \quad X_{0}=x_{0} \tag{3.2.2}
\end{equation*}
$$

where $b, \sigma, \gamma, f$ and $g$ are given functions satisfying the following conditions:
H.2. $b:[0, T] \times \mathbb{R}^{n} \times \mathbb{R}^{k} \rightarrow \mathbb{R}^{n}, \sigma:[0, T] \rightarrow \mathbb{R}^{n \times d}, \gamma:[0, T] \times \mathbb{R}_{0}^{p} \rightarrow \mathbb{R}^{n}, f:[0, T] \times \mathbb{R}^{n} \times \mathbb{R}^{k} \rightarrow$ $\mathbb{R} \cup\{\infty\}, g: \mathbb{R}^{n} \rightarrow \mathbb{R}$ are measurable functions such that for some $L \geq 0$ and $\lambda>0$,
(1) there exist measurable functions $\left(b_{0}, b_{1}, b_{2}\right):[0, T] \rightarrow \mathbb{R}^{n} \times \mathbb{R}^{n \times n} \times \mathbb{R}^{n \times k}$ such that $b(t, x, a)=b_{0}(t)+b_{1}(t) x+b_{2}(t)$ a for all $(t, x, a) \in[0, T] \times \mathbb{R}^{n} \times \mathbb{R}^{k}$, with $\left\|b_{0}\right\|_{L^{2}}+$ $\left\|b_{1}\right\|_{L^{\infty}}+\left\|b_{2}\right\|_{L^{\infty}}+\|\sigma\|_{L^{2}}+\left(\int_{0}^{T} \int_{\mathbb{R}_{0}^{p}}|\gamma(t, u)|^{2} \nu(\mathrm{~d} u) \mathrm{d} t\right)^{1 / 2} \leq L$.
(2) $g$ is convex and differentiable with an L-Lipschitz derivative such that $|\nabla g(0)| \leq L$.
(3) there exist functions $f_{0}:[0, T] \times \mathbb{R}^{n} \times \mathbb{R}^{k} \rightarrow \mathbb{R}$ and $\mathcal{R}: \mathbb{R}^{k} \rightarrow \mathbb{R} \cup\{\infty\}$ such that

$$
f(t, x, a)=f_{0}(t, x, a)+\mathcal{R}(a), \quad \forall(t, x, a) \in[0, T] \times \mathbb{R}^{n} \times \mathbb{R}^{k}
$$

For all $(t, x) \in[0, T] \times \mathbb{R}^{n}, f_{0}(t, x, \cdot)$ is convex in $\mathbb{R}^{k}, f_{0}(t, \cdot, \cdot)$ is differentiable in $\mathbb{R}^{n} \times \mathbb{R}^{k}$ with an L-Lipschitz derivative, and $\left|f_{0}(t, 0,0)\right|+\left|\partial_{(x, a)} f_{0}(t, 0,0)\right| \leq L$. Moreover, $\mathcal{R}$ is proper, lower semicontinuous, and convex. ${ }^{1}$

[^6](4) for all $t \in[0, T],(x, a),\left(x^{\prime}, a^{\prime}\right) \in \mathbb{R}^{n} \times \mathbb{R}^{k}$, and $\eta \in[0,1]$,
\[

$$
\begin{equation*}
\eta f(t, x, a)+(1-\eta) f\left(t, x^{\prime}, a^{\prime}\right) \geq f\left(t, \eta x+(1-\eta) x^{\prime}, \eta a+(1-\eta) a^{\prime}\right)+\eta(1-\eta) \frac{\lambda}{2}\left|a-a^{\prime}\right|^{2} . \tag{3.2.3}
\end{equation*}
$$

\]

Remark 3.2.1. Throughout this paper, let $\operatorname{dom} \mathcal{R}=\left\{a \in \mathbb{R}^{k} \mid \mathcal{R}(a)<\infty\right\}$ be the effective domain of $\mathcal{R}$ (or equivalently the effective domain of $f$ ). Under (H.2), we can show that both the function $f$ and its conjugate function

$$
\begin{equation*}
[0, T] \times \mathbb{R}^{n} \times \mathbb{R}^{k} \ni(t, x, z) \mapsto f^{*}(t, x, z):=\sup \left\{\langle a, z\rangle-f(t, x, a) \mid a \in \mathbb{R}^{k}\right\} \in \mathbb{R} \cup\{\infty\} \tag{3.2.4}
\end{equation*}
$$

are normal convex integrands in the sense of [137, Section 14] and hence measurable, which are crucial for the well-definedness of the control problem (3.2.1) and the characterization of optimal controls. Furthermore, the strong convexity condition (H.2(4)) enables us to establish the Lipschitz stability of feedback controls to (3.2.1), which is essential for the analysis of learning algorithms.

Our analysis and results can be extended to control problems with time-space dependent nonsmooth cost function $\mathcal{R}:[0, T] \times \mathbb{R}^{n} \times \mathbb{R}^{k} \rightarrow \mathbb{R} \cup\{\infty\}$ by assuming $\mathcal{R}$ is a normal convex integrand and satisfies suitable subdifferentability conditions. For notational simplicity and clarity, we choose to refrain from further generalization.

Note that (H.2) allows the diffusion coefficient $\sigma$ to be degenerate, hence the stability results in Section 3.2.3 apply to deterministic control problems. Moreover, (H.2) requires neither the effective domain $\operatorname{dom} \mathcal{R}$ to be closed nor the function $\mathcal{R}$ to be bounded or continuous on $\operatorname{dom} \mathcal{R}$, which is important for problems in engineering and machine learning, as shown in the following examples.

Example 3.2.1 (Sparse and switching controls). Let $\boldsymbol{A} \subset \mathbb{R}^{k}$ be a nonempty closed convex set, $\delta_{\boldsymbol{A}}$ be the indicator of $\boldsymbol{A}$ satisfying $\delta_{\boldsymbol{A}}(x)=0$ for $x \in \boldsymbol{A}$ and $\delta_{\boldsymbol{A}}(x)=\infty$ for $x \in \mathbb{R}^{k} \backslash \boldsymbol{A}$, and $\ell: \mathbb{R}^{k} \rightarrow \mathbb{R}$ be a lower semicontinuous and convex function. Then $\mathcal{R}:=\ell+\delta_{\boldsymbol{A}}$ satisfies (H.2(3)). In particular, by setting $\ell \equiv 0$, we can consider the linear-convex control problems with smooth running costs and control constraints (see e.g. [21] and [169, Theorem 5.2 on p. 137]), which include the most commonly used linear-quadratic models as special cases.

More importantly, it is well-known in optimal control literature (see e.g. [39, 54, 107] and references therein) that, one can employ a nonsmooth function $\ell$ involving $L^{1}$-norm of controls to enhance the sparsity and switching property of optimal controls, which are practically important for mimimum fuel problems and optimal device placement problems. Here by sparsity we refer to the situation where the whole vector $\alpha_{t}$ is zero, while by switching control we refer to the phenomena where at most one coordinate of $\alpha_{t}$ is non-zero at each $t$.

Example 3.2.2 (Regularized relaxed controls). Consider a regularized control problem arising from reinforcement learning (see e.g. [159, 77, 143, 135]), whose cost function $f$ is of the following form:

$$
\begin{equation*}
f(t, x, a)=f_{0}(t, x)+\left\langle f_{1}(t, x), a\right\rangle+\rho D_{\mathfrak{f}}(a \| \mu) \quad \forall(t, x, a) \in[0, T] \times \mathbb{R}^{n} \times \mathbb{R}^{k} \tag{3.2.5}
\end{equation*}
$$

where $f_{0}:[0, T] \times \mathbb{R}^{n} \rightarrow \mathbb{R}, f_{1}:[0, T] \times \mathbb{R}^{n} \rightarrow \mathbb{R}^{k}$ are given functions, $\rho>0$ is a regularization parameter, and $D_{\mathfrak{f}}(\cdot \| \mu): \mathbb{R}^{k} \rightarrow \mathbb{R} \cup\{\infty\}$ is an $\mathfrak{f}$-divergence defined as follows. Let $\Delta_{k}:=$ $\left\{a \in[0,1]^{k} \mid \sum_{i=1}^{k} a_{i}=1\right\}, \mu=\left(\mu_{i}\right)_{i=1}^{k} \in \Delta_{k} \cap(0,1)^{k}$, and $\mathfrak{f}:[0, \infty) \rightarrow \mathbb{R} \cup\{\infty\}$ be a lower semicontinuous function which satisfies $\mathfrak{f}(0)=\lim _{x \rightarrow 0} \mathfrak{f}(x), \mathfrak{f}(1)=0$ and $\mathfrak{f}$ is $\kappa_{\mu}$-strongly convex on $\left[0, \frac{1}{\min _{i} \mu_{i}}\right]$ with a constant $\kappa_{\mu}>0$. Then, the $\mathfrak{f}$-divergence $D_{\mathfrak{f}}(\cdot \| \mu): \mathbb{R}^{k} \rightarrow \mathbb{R} \cup\{\infty\}$ satisfies $D_{\mathfrak{f}}(a \| \mu)=\infty$ for $a \notin \Delta_{k}$ and

$$
D_{\mathfrak{f}}(a \| \mu):=\sum_{i=1}^{k} \mu_{i} f\left(\frac{a_{i}}{\mu_{i}}\right) \in \mathbb{R} \cup\{\infty\} \quad \forall a \in \Delta_{k}
$$

One can easily see from $\mathfrak{f}(1)=0$ and the lower semicontinuity of $\mathfrak{f}$ that $D_{\mathfrak{f}}(\cdot \| \mu)$ is a proper, lower semicontinuous function with effective domain dom $D_{\mathfrak{f}}(\cdot \| \mu) \subset \Delta_{k}$. Moreover, by the $\kappa_{\mu}$-strong convexity of $\mathfrak{f}$, we have for all a, $\tilde{a} \in \Delta_{k}, \eta \in[0,1]$ that

$$
\begin{aligned}
& \eta D_{\mathfrak{f}}(a \| \mu)+(1-\eta) D_{\mathfrak{f}}(\tilde{a} \| \mu) \\
& =\sum_{i=1}^{k} \mu_{i}\left(\eta \mathfrak{f}\left(\frac{a_{i}}{\mu_{i}}\right)+(1-\eta) \mathfrak{f}\left(\frac{\tilde{a}_{i}}{\mu_{i}}\right)\right) \geq \sum_{i=1}^{k} \mu_{i}\left(\mathfrak{f}\left(\frac{\eta a_{i}+(1-\eta) \tilde{a}_{i}}{\mu_{i}}\right)+\eta(1-\eta) \frac{\kappa_{\mu}}{2}\left|\frac{a_{i}-\tilde{a}_{i}}{\mu_{i}}\right|^{2}\right) \\
& \geq D_{\mathfrak{f}}(\eta a+(1-\eta) \tilde{a} \| \mu)+\eta(1-\eta) \frac{\kappa_{\mu}}{2 \max _{i} \mu_{i}}|a-\tilde{a}|^{2},
\end{aligned}
$$

which implies the $\frac{\kappa_{\mu}}{\max _{i} \mu_{i}}$-strong convexity of $D_{\mathfrak{f}}(\cdot \| \mu)$ in $\mathbb{R}^{k}$. It is clear that for suitable choices of $f_{0}, f_{1}$, the function $f$ in (3.2.5) satisfies (H.2(3)).

It is important to notice that an $\mathfrak{f}$-divergence $D_{\mathfrak{f}}(\cdot \| \mu)$ is in general non-differentiable and unbounded on its effective domain. For example, one may consider the relative entropy (with $\mathfrak{f}(s)=s \log s$ ) and the squared Hellinger divergence (with $\mathfrak{f}(s)=2(1-\sqrt{s})$ ), which are not subdifferentiable at the boundary of $\Delta_{k}$. Moreover, the reverse relative entropy (with $\mathfrak{f}(s)=-\log s)$ and the Neyman's $\chi^{2}$ divergence (with $\mathfrak{f}(s)=\frac{1}{s}-1$ ) are unbounded near the boundary of $\Delta_{k}$.

### 3.2.2 Construction of optimal feedback controls

In this section, we apply the maximum principle to (3.2.1) and explicitly construct optimal feedback controls of (3.2.1) based on the associated coupled FBSDE.

The following proposition shows that under (H.2), the control problem (3.2.1) admits a unique optimal open-loop control.

Proposition 3.2.1. Suppose (H.2) holds and let $x_{0} \in \mathbb{R}^{n}$. Then the cost functional $J\left(\alpha ; x_{0}\right)$ : $\mathcal{H}^{2}\left(\mathbb{R}^{k}\right) \rightarrow \mathbb{R} \cup\{\infty\}$ is proper, lower semicontinuous, and $\lambda$-strongly convex. Consequently, $J\left(\cdot ; x_{0}\right)$ admits a unique minimizer $\alpha^{x_{0}}$ in $\mathcal{H}^{2}\left(\mathbb{R}^{k}\right)$.

Proof. The desired properties of $J$ follow directly from the corresponding properties of $f, g$ in (H.2) and the fact that (3.2.2) has affine coefficients. The well-posedness of minimizers
then follows from the standard theory of strongly convex minimization problems on Hilbert spaces (see e.g. [26, Lemma 2.33 (ii)]).

We then proceed to study optimal feedback controls of (3.2.1). The classical control theory shows that under suitable coercivity and convexity conditions, the optimal openloop control of (3.2.1) can be expressed in a feedback form, i.e., there exists a measurable function $\psi:[0, T] \times \mathbb{R}^{n} \rightarrow \mathbb{R}^{k}$ such that $\alpha^{x_{0}}=\psi\left(t, X_{t}^{x_{0}, \alpha^{x_{0}}}\right)$ for $\mathrm{d} \mathbb{P} \otimes \mathrm{d} t$ a.e. (see [123] for the case with controlled jump-diffusions and smooth costs and [81] for the case with controlled diffusions and nonsmooth costs). However, since these non-constructive proofs are based on a measurable selection theorem, the resulting feedback policy $\psi$ may not be unique, and may be unstable with respect to perturbations of the state dynamics.

In the subsequent analysis, we give a constructive proof of the existence of Lipschitz continuous feedback controls by exploiting the linear-convex structure of the control problem (3.2.1)-(3.2.2). Such a feedback control can be explicitly represented as solutions of a suitable FBSDE, and hence is Lipschitz stable with respect to perturbations of underlying models (see Theorem 3.2.6).

We first present the precise definitions of feedback controls and the associated state processes.

Definition 3.2.1. Let $\mathcal{V}$ be the following space of feedback controls:

$$
\mathcal{V}:=\left\{\begin{array}{l|l}
\psi:[0, T] \times \mathbb{R}^{n} \rightarrow \mathbb{R}^{k} & \begin{array}{l}
\psi \text { is measurable and there exists } C \geq 0 \text { such that } \\
\text { for all }(t, x, y) \in[0, T] \times \mathbb{R}^{n} \times \mathbb{R}^{n},|\psi(t, 0)| \leq C \\
\text { and }|\psi(t, x)-\psi(t, y)| \leq C|x-y|
\end{array} \tag{3.2.6}
\end{array}\right\}
$$

For any given $x_{0} \in \mathbb{R}^{n}$ and $\psi \in \mathcal{V}$, we say $X^{x_{0}, \psi} \in \mathcal{S}^{2}\left(\mathbb{R}^{n}\right)$ is the state process associated with $\psi$ if it satisfies the following dynamics:

$$
\begin{equation*}
\mathrm{d} X_{t}=b\left(t, X_{t}, \psi\left(t, X_{t}\right)\right) \mathrm{d} t+\sigma(t) \mathrm{d} W_{t}+\int_{\mathbb{R}_{0}^{p}} \gamma(t, u) \tilde{N}(\mathrm{~d} t, \mathrm{~d} u), \quad t \in[0, T], \quad X_{0}=x_{0} \tag{3.2.7}
\end{equation*}
$$

We say $\psi \in \mathcal{V}$ is an optimal feedback control of (3.2.1) if it holds for $\mathrm{d} \mathbb{P} \otimes \mathrm{d} t$ a.e. that $\alpha_{t}^{x_{0}}=\psi\left(t, X_{t}^{x_{0}, \psi}\right)$, where $\alpha^{x_{0}} \in \mathcal{H}^{2}\left(\mathbb{R}^{n}\right)$ is the optimal control of (3.2.1).

We then proceed to establish a maximum principle for feedback controls of the control problem (3.2.1) with non-smooth costs. Let $H:[0, T] \times \mathbb{R}^{n} \times \mathbb{R}^{k} \times \mathbb{R}^{n} \rightarrow \mathbb{R} \cup\{\infty\}$ and $\phi:[0, T] \times \mathbb{R}^{n} \times \mathbb{R}^{n} \rightarrow \mathbb{R}^{k}$ such that for all $(t, x, a, y) \in[0, T] \times \mathbb{R}^{n} \times \mathbb{R}^{k} \times \mathbb{R}^{n}$,

$$
\begin{equation*}
H(t, x, a, y):=\langle b(t, x, a), y\rangle+f(t, x, a), \quad \phi(t, x, y):=\underset{a \in \mathbb{R}^{k}}{\arg \min } H(t, x, a, y) \in \operatorname{dom} \mathcal{R} . \tag{3.2.8}
\end{equation*}
$$

The following lemma shows that the function $\phi$ is well-defined and measurable.

Lemma 3.2.2. Suppose (H.2) holds. Then the function $\phi:[0, T] \times \mathbb{R}^{n} \times \mathbb{R}^{n} \rightarrow \mathbb{R}^{k}$ defined in (3.2.8) is measurable and satisfies for all $t \in[0, T], x, y \in \mathbb{R}^{n}$ that

$$
\begin{equation*}
\phi(t, x, y)=\partial_{z} f^{*}\left(t, x,-b_{2}(t)^{\top} y\right) \tag{3.2.9}
\end{equation*}
$$

where the function $f^{*}$ is defined in (3.2.4).
Proof. Let $f^{*}:[0, T] \times \mathbb{R}^{n} \times \mathbb{R}^{k} \rightarrow \mathbb{R} \cup\{\infty\}$ be the function defined in (3.2.4). Recall that for each $(t, x) \in[0, T] \times \mathbb{R}^{n}, f(t, x, \cdot)$ is $\lambda$-strongly convex and lower semicontinuous. Hence by [137, Theorems 11.3 and 11.8], $f^{*}(t, x, \cdot)$ is finite and differentiable on $\mathbb{R}^{k}$, and $\partial_{z} f^{*}(t, x, z)=\arg \max _{a \in \operatorname{dom} \mathcal{R}}(\langle a, z\rangle-f(t, x, a))$ for all $z \in \mathbb{R}^{k}$. Moreover, by [83, Theorem E4.2.1], $\mathbb{R}^{k} \ni z \mapsto \partial_{z} f^{*}(t, x, z) \in \mathbb{R}^{k}$ is $1 / \lambda$-Lipschitz continuous. Hence, from the definition of $\phi$ and (H.2(1)), for all $t \in[0, T], x, y \in \mathbb{R}^{n}$,

$$
\begin{align*}
\phi(t, x, y) & =\underset{a \in \mathbb{R}^{k}}{\arg \min }(\langle b(t, x, a), y\rangle+f(t, x, a))=\underset{a \in \mathbb{R}^{k}}{\arg \max }\left(\left\langle a,-b_{2}(t)^{\top} y\right\rangle-f(t, x, a)\right) \\
& =\partial_{z} f^{*}\left(t, x,-b_{2}(t)^{\top} y\right) \tag{3.2.10}
\end{align*}
$$

Note that the measurability of $f^{*}$ (see Remark 3.2.1) implies that the derivative $\partial_{z} f^{*}$ is measurable, which along with the continuity of $z \mapsto \partial_{z} f^{*}(t, x, z)$ leads to the measurability of $\phi$.

With the measurable function $\phi$ in hand, for each $(t, x) \in[0, T] \times \mathbb{R}^{n}$, let us consider the following coupled FBSDE on $[t, T]$ : for all $s \in[t, T]$,

$$
\begin{align*}
\mathrm{d} X_{s} & =b\left(s, X_{s}, \phi\left(s, X_{s}, Y_{s}\right)\right) \mathrm{d} s+\sigma(s) \mathrm{d} W_{s}+\int_{\mathbb{R}_{0}^{p}} \gamma(s, u) \tilde{N}(\mathrm{~d} s, \mathrm{~d} u),  \tag{3.2.11a}\\
\mathrm{d} Y_{s} & =-\partial_{x} H\left(s, X_{s}, \phi\left(s, X_{s}, Y_{s}\right), Y_{s}\right) \mathrm{d} s+Z_{s} \mathrm{~d} W_{s}+\int_{\mathbb{R}_{0}^{p}} M_{s} \tilde{N}(\mathrm{~d} s, \mathrm{~d} u), \tag{3.2.11b}
\end{align*} \quad Y_{T}=\nabla g\left(X_{T}\right) .
$$

We say a tuple of processes $\left(X^{t, x}, Y^{t, x}, Z^{t, x}, M^{t, x}\right) \in \mathbb{S}(t, T):=\mathcal{S}^{2}\left(t, T ; \mathbb{R}^{n}\right) \times \mathcal{S}^{2}\left(t, T ; \mathbb{R}^{n}\right) \times$ $\mathcal{H}^{2}\left(t, T ; \mathbb{R}^{n \times d}\right) \times \mathcal{H}_{\nu}^{2}\left(t, T ; \mathbb{R}^{n}\right)$ is a solution to (3.2.11) (on $[t, T]$ with initial condition $\left.X_{t}^{t, x}=x\right)$ if it satisfies (3.2.11) $\mathbb{P}$-almost surely.

The next lemma presents several important properties of the Hamiltonian $H$ and the function $\phi$ defined in (3.2.8), which are essential for the well-posedness and stability of (3.2.11).

Lemma 3.2.3. Suppose (H.2) holds. Let $\phi:[0, T] \times \mathbb{R}^{n} \times \mathbb{R}^{n} \rightarrow \mathbb{R}^{k}$ be the function defined in (3.2.8). Then there exists a constant $C$ such that for all $t \in[0, T]$ and $(x, y),\left(x^{\prime}, y^{\prime}\right) \in$ $\mathbb{R}^{n} \times \mathbb{R}^{n},|\phi(t, 0,0)| \leq C,\left|\phi(t, x, y)-\phi\left(t, x^{\prime}, y^{\prime}\right)\right| \leq C\left(\left|x-x^{\prime}\right|+\left|y-y^{\prime}\right|\right)$ and

$$
\begin{align*}
& \left\langle b(t, x, \phi(t, x, y))-b\left(t, x^{\prime}, \phi\left(t, x^{\prime}, y^{\prime}\right)\right), y-y^{\prime}\right\rangle \\
& \quad+\left\langle-\partial_{x} H(t, x, \phi(t, x, y), y)+\partial_{x} H\left(t, x^{\prime}, \phi\left(t, x^{\prime}, y^{\prime}\right), y^{\prime}\right), x-x^{\prime}\right\rangle  \tag{3.2.12}\\
& \leq-\lambda\left|\phi(t, x, y)-\phi\left(t, x^{\prime}, y^{\prime}\right)\right|^{2}
\end{align*}
$$

with the constant $\lambda$ in (H.2).
Proof. We start by showing the boundedness of $\phi(\cdot, 0,0)$ by considering $a(t):=\left(\partial_{z} f^{*}\right)(t, 0,0)$ for each $t \in[0, T]$, where $f^{*}:[0, T] \times \mathbb{R}^{n} \times \mathbb{R}^{k} \rightarrow \mathbb{R}$ is defined as in (3.2.4). The fact that $f(t, 0, \cdot)$ is proper, lower semicontinuous and convex implies that $0 \in \widehat{\partial}_{a} f(t, 0, a(t))$ for all $t \in[0, T]$, where $\widehat{\partial}_{a} f(t, 0, a(t))$ is the subdifferential of $f(t, 0, \cdot)$ at $a(t)$. Note that $f_{0}(t, 0, \cdot)$ and $\mathcal{R}$ are proper, lower semicontinuous, and convex, and $\operatorname{dom} \mathcal{R} \subset \operatorname{dom} f_{0}(t, 0, \cdot)=\mathbb{R}^{k}$. Hence by [137, Corollary 10.9], $\widehat{\partial}_{a} f(t, 0, a)=\partial_{a} f_{0}(t, 0, a)+\widehat{\partial} \mathcal{R}(a)$ for all $(t, a) \in[0, T] \times \mathbb{R}^{k}$, where $\widehat{\partial} \mathcal{R}(a)$ is the subdifferential of $\mathcal{R}$ at $a$. Now fix an arbitrary $t_{0} \in[0, T]$ and set $a_{0}=a\left(t_{0}\right)$. The fact that $0 \in \widehat{\partial}_{a} f\left(t_{0}, 0, a_{0}\right)$ implies that $-\partial_{a} f_{0}\left(t_{0}, 0, a_{0}\right) \in \widehat{\partial} \mathcal{R}\left(a_{0}\right)$ and hence $\partial_{a} f_{0}\left(t, 0, a_{0}\right)-\partial_{a} f_{0}\left(t_{0}, 0, a_{0}\right) \in \widehat{\partial}_{a} f\left(t, 0, a_{0}\right)$ for all $t \in[0, T]$. By the strong convexity condition (3.2.3), for all $t \in[0, T], \xi_{1} \in \widehat{\partial}_{a} f\left(t, 0, a_{0}\right)$ and $\xi_{2} \in \widehat{\partial}_{a} f(t, 0, a(t))$,

$$
\lambda\left|a_{0}-a(t)\right|^{2} \leq\left\langle\xi_{1}-\xi_{2}, a_{0}-a(t)\right\rangle \leq\left|\xi_{1}-\xi_{2}\right|\left|a_{0}-a(t)\right| .
$$

Taking $\xi_{1}=\partial_{a} f_{0}\left(t, 0, a_{0}\right)-\partial_{a} f_{0}\left(t_{0}, 0, a_{0}\right)$ and $\xi_{2}=0$ in the above inequality yields

$$
\left|a_{0}-a(t)\right| \leq\left|\partial_{a} f_{0}\left(t, 0, a_{0}\right)-\partial_{a} f_{0}\left(t_{0}, 0, a_{0}\right)\right| / \lambda \leq C,
$$

by the linear growth of $\partial_{a} f_{0}(t, 0, \cdot)$. This implies that $\left|\left(\partial_{z} f^{*}\right)(t, 0,0)\right| \leq C$ for all $t \in[0, T]$, which along with (3.2.10) leads to the desired uniform boundedness of $\phi(\cdot, 0,0)$.

We proceed to establish the Lipschitz continuity of $\phi$ with respect to $(x, y)$. The $1 / \lambda$ Lipschitz continuity of $\partial_{z} f^{*}(t, x, \cdot)$ and the boundedness of $b_{2}$ imply that $\phi$ is Lipschitz continuous in $y$, uniformly with respect to $(t, x)$. It remains to show the Lipschitz continuity of $\partial_{z} f^{*}$ with respect to $x$, which along with (3.2.10) leads to the desired Lipchitz continuity of $\phi$. For any given $(t, z) \in[0, T] \times \mathbb{R}^{k}$ and $x, x^{\prime} \in \mathbb{R}^{n}$, let $a=\partial_{z} f^{*}(t, x, z)$ and $a^{\prime}=$ $\partial_{z} f^{*}\left(t, x^{\prime}, z\right)$. Then we have $z \in \widehat{\partial}_{a} f(t, x, a)$ and $z \in \widehat{\partial}_{a} f\left(t, x^{\prime}, a^{\prime}\right)$. Moreover, by the convexity of $f(t, x, \cdot)$ for all $(t, x) \in[0, T] \times \mathbb{R}^{n}$ and similar arguments as above, we can show that $z-\partial_{a} f_{0}\left(t, x^{\prime}, a^{\prime}\right)+\partial_{a} f_{0}\left(t, x, a^{\prime}\right) \in \widehat{\partial}_{a} f\left(t, x, a^{\prime}\right)$, which together with the convexity condition (3.2.3) and $z \in \widehat{\partial}_{a} f(t, x, a)$ leads to

$$
\lambda\left|a^{\prime}-a\right| \leq\left|z-\partial_{a} f_{0}\left(t, x^{\prime}, a^{\prime}\right)+\partial_{a} f_{0}\left(t, x, a^{\prime}\right)-z\right| \leq L\left|x-x^{\prime}\right|
$$

where we have used the $L$-Lipchitz continuity of $\partial_{a} f_{0}(t, \cdot, \cdot)$. This finishes the proof of the Lipschitz continuity of $\partial_{z} f^{*}(t, \cdot, \cdot)$ and $\phi(t, \cdot, \cdot)$.

Finally, we establish the monotonicity condition (3.2.12). By (H.2(4)), for all $t \in[0, T]$, $x, x^{\prime} \in \mathbb{R}^{n}, a, a^{\prime} \in \mathbb{R}^{k}, y \in \mathbb{R}^{n}$, the function $\mathbb{R}^{n} \times \mathbb{R}^{k} \ni(x, a) \mapsto H(t, x, a, y) \in \mathbb{R}^{n} \cup\{\infty\}$ satisfies the same convexity condition (3.2.3) as the function $f$, and hence

$$
\begin{equation*}
H\left(t, x^{\prime}, a^{\prime}, y\right)-H(t, x, a, y) \geq\left\langle\xi, x^{\prime}-x, a^{\prime}-a\right\rangle+\frac{\lambda}{2}\left|a^{\prime}-a\right|^{2} \quad \forall \xi \in \widehat{\partial}_{(x, a)} H(t, x, a, y) \tag{3.2.13}
\end{equation*}
$$

where $\widehat{\partial}_{(x, a)} H(t, x, a, y)$ denotes the subdifferential of the function $H(t, \cdot, y)$ at $(x, a)$. Moreover, for any given $t \in[0, T]$ and $x, y \in \mathbb{R}^{n}$, the definition of $\phi$ in (3.2.8) implies that
$0 \in \widehat{\partial}_{a} H(t, x, \phi(t, x, y), y)$, where $\widehat{\partial}_{a} H(t, x, \phi(t, x, y), y)$ denotes the subdifferential of the function $H(t, x, \cdot, y)$ at $\phi(t, x, y)$. Now recall that for any Euclidean space $E$, convex function $F: E \rightarrow \mathbb{R} \cup\{\infty\}$ and $x \in \operatorname{dom} F, v \in \widehat{\partial} F(x)$ if and only if $\liminf _{\tau \rightarrow 0, \tilde{w} \rightarrow w} \frac{F(x+\tau \tilde{w})-F(x)}{\tau} \geq$ $\langle v, w\rangle$ for all $w \in E$ (see e.g., Exercise 8.4 and Proposition 8.12 in [137]). Thus, for any $t \in[0, T]$ and $x, y \in \mathbb{R}^{n}, 0 \in \widehat{\partial}_{a} H(t, x, \phi(t, x, y), y)$ yields for all $z \in \mathbb{R}^{k}$,

$$
\begin{equation*}
\liminf _{\tau \rightarrow 0, \tilde{z} \rightarrow z} \frac{H(t, x, \phi(t, x, y)+\tau \tilde{z}, y)-H(t, x, \phi(t, x, y), y)}{\tau} \geq\langle 0, z\rangle=0 \tag{3.2.14}
\end{equation*}
$$

Moreover, by the convexity of $H$ and the continuity of $\partial_{x} H$ in $(x, a)$, for any $t \in[0, T]$ and $x, y, w \in \mathbb{R}^{n}$ and $z \in \mathbb{R}^{k}$,

$$
\begin{align*}
& \quad \liminf _{\tau \rightarrow 0,(\tilde{w}, \tilde{z}) \rightarrow(w, z)} \frac{H(t, x+\tau \tilde{w}, \phi(t, x, y)+\tau \tilde{z}, y)-H(t, x, \phi(t, x, y)+\tau \tilde{z}, y)}{\tau} \\
& \geq \liminf _{\tau \rightarrow 0,(\tilde{w}, \tilde{z}) \rightarrow(w, z)} \frac{\left\langle\partial_{x} H(t, x, \phi(t, x, y)+\tau \tilde{z}, y), \tau \tilde{w}\right\rangle}{\tau} \geq\left\langle\partial_{x} H(t, x, \phi(t, x, y), y), w\right\rangle \tag{3.2.15}
\end{align*}
$$

provided that $\phi(t, x, y)+\tau \tilde{z} \in \operatorname{dom} \mathcal{R}$ (cf. (3.2.8)). Then for any $t \in[0, T]$ and $x, y \in \mathbb{R}^{n}$, adding up (3.2.14) and (3.2.15) and using the fact that $\phi(t, x, y) \in \operatorname{dom} \mathcal{R}$ give for all $(w, z) \in \mathbb{R}^{n} \times \mathbb{R}^{k}$,

$$
\begin{aligned}
& \liminf _{\tau \rightarrow 0,(\tilde{w}, \tilde{z}) \rightarrow(w, z)} \frac{H(t, x+\tau \tilde{w}, \phi(t, x, y)+\tau \tilde{z}, y)-H(t, x, \phi(t, x, y), y)}{\tau} \\
& \geq\left\langle\partial_{x} H(t, x, \phi(t, x, y), y), w\right\rangle+\langle 0, z\rangle
\end{aligned}
$$

which implies

$$
\begin{equation*}
\left(\partial_{x} H(t, x, \phi(t, x, y), y), 0\right) \subset \widehat{\partial}_{(x, a)} H(t, x, \phi(t, x, y), y) \tag{3.2.16}
\end{equation*}
$$

Hence for all $t \in[0, T],\left(x_{1}, y_{1}\right),\left(x_{2}, y_{2}\right) \in \mathbb{R}^{n} \times \mathbb{R}^{n}$, we can define $a_{1}=\phi\left(t, x_{1}, y_{1}\right), a_{2}=$ $\phi\left(t, x_{2}, y_{2}\right)$ and deduce that

$$
\begin{aligned}
& \left\langle b\left(t, x_{1}, a_{1}\right)-b\left(t, x_{2}, a_{2}\right), y_{1}-y_{2}\right\rangle+\left\langle-\partial_{x} H\left(t, x_{1}, a_{1}, y_{1}\right)+\partial_{x} H\left(t, x_{2}, a_{2}, y_{2}\right), x_{1}-x_{2}\right\rangle \\
& =H\left(t, x_{1}, a_{1}, y_{1}\right)-H\left(t, x_{2}, a_{2}, y_{1}\right)-\left\langle\partial_{x} H\left(t, x_{1}, a_{1}, y_{1}\right), x_{1}-x_{2}\right\rangle \\
& \quad-\left(H\left(t, x_{1}, a_{1}, y_{2}\right)-H\left(t, x_{2}, a_{2}, y_{2}\right)-\left\langle\partial_{x} H\left(t, x_{2}, a_{2}, y_{2}\right), x_{1}-x_{2}\right\rangle\right) \\
& \leq-\lambda\left|a_{1}-a_{2}\right|^{2}
\end{aligned}
$$

which finishes the proof of the desired monotonicity condition.
The following proposition shows that (3.2.11) admits a unique solution, which is Lipschitz continuous with respect to the initial state. The proof is based on the stability of (3.2.11) under the generalized monotonicity condition (3.2.12) (see Lemma 3.6.1), and follows [134, Corollary 2.4] for the case without jumps.

Proposition 3.2.4. Suppose (H.2) holds. Then for any given $(t, x) \in[0, T] \times \mathbb{R}^{n}$, the $F B S D E$ (3.2.11) admits a unique solution $\left(X^{t, x}, Y^{t, x}, Z^{t, x}, M^{t, x}\right) \in \mathbb{S}(t, T)$. Moreover, there exists a constant $C$ such that for all $t \in[0, T]$ and $x, x^{\prime} \in \mathbb{R}^{n},\left\|\left(X^{t, x}, Y^{t, x}, Z^{t, x}, M^{t, x}\right)\right\|_{\mathbb{S}(t, T)} \leq$ $C(1+|x|)$ and $\left\|\left(X^{t, x}-X^{t, x^{\prime}}, Y^{t, x}-Y^{t, x^{\prime}}, Z^{t, x}-Z^{t, x^{\prime}}, M^{t, x}-M^{t, x^{\prime}}\right)\right\|_{\mathbb{S}(t, T)} \leq C\left|x-x^{\prime}\right|$.

Now we are ready to present the main result of this section, which constructs an optimal feedback control of (3.2.1) based on the Hamiltonian (3.2.8) and the solutions to the FBSDE (3.2.11).

Theorem 3.2.5. Suppose (H.2) holds. Let $\psi:[0, T] \times \mathbb{R}^{n} \rightarrow \mathbb{R}^{k}$ be the function defined as

$$
\begin{equation*}
\psi(t, x):=\phi\left(t, x, Y_{t}^{t, x}\right), \quad(t, x) \in[0, T] \times \mathbb{R}^{n} \tag{3.2.17}
\end{equation*}
$$

where the function $\phi$ is defined in (3.2.8). Then there exists a constant $C$ such that $|\psi(t, 0)| \leq$ $C$ and $\left|\psi(t, x)-\psi\left(t, x^{\prime}\right)\right| \leq C\left|x-x^{\prime}\right|$ for all $t \in[0, T], x, x^{\prime} \in \mathbb{R}^{n}$. Moreover, for all $x_{0} \in \mathbb{R}^{n}$, $\psi$ is an optimal feedback control of (3.2.1).

Proof. We first analyze the mapping $[0, T] \times \mathbb{R}^{n} \ni(t, x) \mapsto v(t, x):=Y_{t}^{t, x} \in \mathbb{R}^{n}$. Note by Proposition 3.2.4, for any given $(t, x) \in[0, T] \times \mathbb{R}^{n}$, the solution to (3.2.11) (with initial time $t$ and initial state $x$ ) is pathwise unique and Lipschitz continuous with respect to the initial state $x \in \mathbb{R}^{n}$. Hence, it is well-known that (see e.g., Theorem 3.1 and Remarks 3.2-3.3 in [108]) that the map $v$ can be identified with a deterministic function in the space $\mathcal{V}$ and it holds for all $(t, x) \in[0, T] \times \mathbb{R}^{n}$ that $\mathbb{P}\left(\forall s \in[t, T], Y_{s}^{t, x}=v\left(s, X_{s}^{t, x}\right)\right)=1$. Thus, from the regularity of $\phi$ and $v,|\psi(t, 0)| \leq C$ and $\left|\psi(t, x)-\psi\left(t, x^{\prime}\right)\right| \leq C\left|x-x^{\prime}\right|$ for all $x, x^{\prime} \in \mathbb{R}^{n}$, i.e., $\psi$ is in the space $\mathcal{V}$.

Now let $x_{0} \in \mathbb{R}^{n}$ be a given initial state and $\tilde{\alpha} \in \mathcal{A}$ satisfy for $\mathrm{d} \mathbb{P} \otimes \mathrm{d} t$ a.e. that $\tilde{\alpha}_{t}=$ $\phi\left(t, X_{t}^{0, x_{0}}, Y_{t}^{0, x_{0}}\right)$. Then for $\mathrm{d} \mathbb{P} \otimes \mathrm{d} t$ a.e., $\tilde{\alpha}_{t}=\phi\left(t, X_{t}^{0, x_{0}}, v\left(t, X_{t}^{0, x_{0}}\right)\right)=\psi\left(t, X_{t}^{0, x_{0}}\right)$, and $X^{0, x_{0}}$ is the solution to (3.2.2) controlled by $\tilde{\alpha}$, because ( $X^{0, x_{0}}, Y^{0, x_{0}}$ ) satisfy (3.2.11a). Since the control problem (3.2.1) admits an unique optimal control in $\mathcal{H}^{2}\left(\mathbb{R}^{k}\right)$, it suffices to show that $\tilde{\alpha}$ is optimal. By (3.2.16), for $\mathrm{d} \mathbb{P} \otimes \mathrm{d} t$ a.e.,

$$
\left(\partial_{x} H\left(t, X_{t}^{0, x_{0}}, \phi\left(t, X_{t}^{0, x_{0}}, Y_{t}^{0, x_{0}}\right), Y_{t}^{0, x_{0}}\right), 0\right) \subset \widehat{\partial}_{(x, a)} H\left(t, X_{t}^{0, x_{0}}, \phi\left(t, X_{t}^{0, x_{0}}, Y_{t}^{0, x_{0}}\right), Y_{t}^{0, x_{0}}\right)
$$

Then for any given $\alpha \in \mathcal{H}^{2}\left(\mathbb{R}^{n}\right)$ with the state process $X^{x_{0}, \alpha}$ satisfying the controlled dynamics (3.2.2), by the definition of $H$ in (3.2.8), (H.2(2)) and (3.2.13),

$$
\begin{aligned}
& J\left(\alpha ; x_{0}\right)-J\left(\tilde{\alpha} ; x_{0}\right) \\
& =\mathbb{E}\left[g\left(X_{T}^{x_{0}, \alpha}\right)-g\left(X_{T}^{0, x_{0}}\right)+\int_{0}^{T}\left(H\left(t, X_{t}^{x_{0}, \alpha}, \alpha_{t}, Y_{t}^{0, x_{0}}\right)-H\left(t, X_{t}^{0, x_{0}}, \tilde{\alpha}_{t}, Y_{t}^{0, x_{0}}\right)\right) \mathrm{d} t\right] \\
& \left.\quad-\int_{0}^{T}\left\langle b\left(t, X_{t}^{x_{0}, \alpha}, \alpha_{t}\right)-b\left(t, X_{t}^{0, x_{0}}, \tilde{\alpha}_{t}\right), Y_{t}^{0, x_{0}}\right\rangle \mathrm{d} t\right] \\
& \geq \mathbb{E}\left[\left\langle\nabla g\left(X_{T}^{0, x_{0}}\right), X_{T}^{x_{0}, \alpha}-X_{T}^{0, x_{0}}\right\rangle+\int_{0}^{T}\left\langle\partial_{x} H\left(t, X_{t}^{0, x_{0}}, \phi\left(t, X_{t}^{0, x_{0}}, Y_{t}^{0, x_{0}}\right), Y_{t}^{0, x_{0}}\right), X_{t}^{x_{0}, \alpha}-X_{t}^{0, x_{0}}\right\rangle \mathrm{d} t\right. \\
& \left.\quad-\int_{0}^{T}\left\langle b\left(t, X_{t}^{x_{0}, \alpha}, \alpha_{t}\right)-b\left(t, X_{t}^{0, x_{0}}, \tilde{\alpha}_{t}\right), Y_{t}^{0, x_{0}}\right\rangle \mathrm{d} t\right]=0,
\end{aligned}
$$

where the last equality is by applying Itô's formula to the process $\left(\left\langle X_{t}^{x_{0}, \alpha}-X_{t}^{0, x_{0}}, Y_{t}^{0, x_{0}}\right\rangle\right)_{t \geq 0}$ and by the FBSDE (3.2.11). That is, $\psi \in \mathcal{V}$ is an optimal feedback control of (3.2.1).

### 3.2.3 Lipschitz stability of optimal feedback controls and associated costs

In this section, we establish the Lipschitz stability of the optimal feedback controls constructed in Theorem 3.2.5 and their associated costs: that is, they are are Lipschitz continuous with respect to the perturbation in the coefficients of (3.2.2). Such a Lipschitz stability property is crucial for the subsequent analysis of learning algorithms.

More precisely, for any given $x_{0} \in \mathbb{R}^{n}$, we consider a perturbed control problem where the cost functions $f, g$ are the same as those in (3.2.1), and for each $\alpha \in \mathcal{H}^{2}\left(\mathbb{R}^{n}\right)$, the corresponding state dynamics satisfies the following perturbed dynamics:

$$
\begin{equation*}
\mathrm{d} X_{t}=\tilde{b}\left(t, X_{t}, \alpha_{t}\right) \mathrm{d} t+\tilde{\sigma}(t) \mathrm{d} W_{t}+\int_{\mathbb{R}_{0}^{p}} \tilde{\gamma}(t, u) \tilde{N}(\mathrm{~d} t, \mathrm{~d} u), \quad t \in[0, T], \quad X_{0}=x_{0} \tag{3.2.18}
\end{equation*}
$$

whose coefficients satisfy the following assumption:
H.3. $\tilde{b}:[0, T] \times \mathbb{R}^{n} \times \mathbb{R}^{k} \rightarrow \mathbb{R}^{n}, \tilde{\sigma}:[0, T] \rightarrow \mathbb{R}^{n \times d}$ and $\tilde{\gamma}:[0, T] \times \mathbb{R}_{0}^{p} \rightarrow \mathbb{R}^{n}$ satisfy (H.2(1)) with the same constant $L$, i.e., there exist measurable functions $\left(\tilde{b}_{0}, \tilde{b}_{1}, \tilde{b}_{2}\right):[0, T] \rightarrow$ $\mathbb{R}^{n} \times \mathbb{R}^{n \times n} \times \mathbb{R}^{n \times k}$ such that $\tilde{b}(t, x, a)=\tilde{b}_{0}(t)+\tilde{b}_{1}(t) x+\tilde{b}_{2}(t)$ a for all $(t, x, a) \in[0, T] \times \mathbb{R}^{n} \times \mathbb{R}^{k}$ and $\left\|\tilde{b}_{0}\right\|_{L^{2}}+\left\|\tilde{b}_{1}\right\|_{L^{\infty}}+\left\|\tilde{b}_{2}\right\|_{L^{\infty}}+\|\tilde{\sigma}\|_{L^{2}}+\left(\int_{0}^{T} \int_{\mathbb{R}_{0}^{p}}|\tilde{\gamma}(t, u)|^{2} \nu(\mathrm{~d} u) \mathrm{d} t\right)^{1 / 2} \leq L$.

Under (H.2) and (H.3), Theorem 3.2.5 ensures that an optimal feedback control of the perturbed control problem can be obtained by

$$
\begin{equation*}
[0, T] \times \mathbb{R}^{n} \ni(t, x) \mapsto \tilde{\psi}(t, x):=\tilde{\phi}\left(t, x, \tilde{Y}_{t}^{t, x}\right) \in \mathbb{R}^{k} \tag{3.2.19}
\end{equation*}
$$

where $\tilde{\phi}:[0, T] \times \mathbb{R}^{n} \times \mathbb{R}^{n} \rightarrow \mathbb{R}^{k}$ satisfies for all $(t, x, a, y) \in[0, T] \times \mathbb{R}^{n} \times \mathbb{R}^{k} \times \mathbb{R}^{n}$ that

$$
\begin{equation*}
\tilde{\phi}(t, x, y):=\underset{a \in \mathbb{R}^{k}}{\arg \min } \tilde{H}(t, x, a, y), \quad \tilde{H}(t, x, a, y):=\langle\tilde{b}(t, x, a), y\rangle+f(t, x, a), \tag{3.2.20}
\end{equation*}
$$

and for each $(t, x) \in[0, T] \times \mathbb{R}^{n},\left(\tilde{X}^{t, x}, \tilde{Y}^{t, x}, \tilde{Z}^{t, x}, \tilde{M}^{t, x}\right) \in \mathbb{S}(t, T)$ is the solution to the following perturbed FBSDE: for all $s \in[t, T]$,

$$
\begin{align*}
\mathrm{d} X_{s} & =\tilde{b}\left(s, X_{s}, \tilde{\phi}\left(s, X_{s}, Y_{s}\right)\right) \mathrm{d} s+\tilde{\sigma}(s) \mathrm{d} W_{s}+\int_{\mathbb{R}_{0}^{p}} \tilde{\gamma}(s, u) \tilde{N}(\mathrm{~d} s, \mathrm{~d} u), \\
\mathrm{d} Y_{s} & =-\partial_{x} \tilde{H}\left(s, X_{s}, \tilde{\phi}\left(s, X_{s}, Y_{s}\right), Y_{s}\right) \mathrm{d} s+Z_{s} \mathrm{~d} W_{s}+\int_{\mathbb{R}_{0}^{p}} M_{s} \tilde{N}(\mathrm{~d} s, \mathrm{~d} u), \tag{3.2.21}
\end{align*} \quad Y_{T}=\nabla g\left(X_{T}\right) .
$$

The following theorem quantifies the difference of optimal feedback controls in terms of the magnitude of perturbations in the coefficients.

Theorem 3.2.6. Suppose (H.Q) and (H.3) hold. Let $\psi, \tilde{\psi}:[0, T] \times \mathbb{R}^{n} \times \mathbb{R}^{n} \rightarrow \mathbb{R}^{k}$ be the functions defined in (3.2.17) and (3.2.19), respectively. Then there exists a constant $C$ such that $|\psi(t, x)-\tilde{\psi}(t, x)| \leq C(1+|x|) \mathcal{E}_{\text {per }}$ for all $(t, x) \in[0, T] \times \mathbb{R}^{n}$, with the constant $\mathcal{E}_{\text {per }}$ defined by

$$
\begin{align*}
\mathcal{E}_{\text {per }}:= & \left\|b_{0}-\tilde{b}_{0}\right\|_{L^{2}}+\left\|b_{1}-\tilde{b}_{1}\right\|_{L^{\infty}}+\left\|b_{2}-\tilde{b}_{2}\right\|_{L^{\infty}}+\|\sigma-\tilde{\sigma}\|_{L^{2}} \\
& +\left(\int_{0}^{T} \int_{\mathbb{R}_{0}^{p}}|\gamma(t, u)-\tilde{\gamma}(t, u)|^{2} \nu(\mathrm{~d} u) \mathrm{d} t\right)^{1 / 2} . \tag{3.2.22}
\end{align*}
$$

Proof. Throughout this proof, for each $(t, x) \in[0, T] \times \mathbb{R}^{n}$, let $\left(X^{t, x}, Y^{t, x}, Z^{t, x}, M^{t, x}\right) \in \mathbb{S}(t, T)$ and $\left(\tilde{X}^{t, x}, \tilde{Y}^{t, x}, \tilde{Z}^{t, x}, \tilde{M}^{t, x}\right) \in \mathbb{S}(t, T)$ be the solutions to (3.2.11) and (3.2.21), respectively, and let $C$ be a generic constant which is independent of $(t, x) \in[0, T] \times \mathbb{R}^{n}$. Then by Proposition 3.2.4, there exists $C \geq 0$ such that for all $(t, x) \in[0, T] \times \mathbb{R}^{n},\left\|\left(X^{t, x}, Y^{t, x}, Z^{t, x}, M^{t, x}\right)\right\|_{\mathbb{S}(t, T)} \leq$ $C(1+|x|)$ and $\left\|\left(\tilde{X}^{t, x}, \tilde{Y}^{t, x}, \tilde{Z}^{t, x}, \tilde{M}^{t, x}\right)\right\|_{\mathbb{S}(t, T)} \leq C(1+|x|)$.

We first estimate the difference between ( $X^{t, x}, Y^{t, x}, Z^{t, x}, M^{t, x}$ ) and ( $\tilde{X}^{t, x}, \tilde{Y}^{t, x}, \tilde{Z}^{t, x}, \tilde{M}^{t, x}$ ) for a given $(t, x) \in[0, T] \times \mathbb{R}^{n}$. By Lemmas 3.2.3 and 3.6.1,

$$
\begin{aligned}
& \|\left(X^{t, x}-\tilde{X}^{t, x}, Y^{t, x}-\tilde{Y}^{t, x}, Z^{t, x}-\tilde{Z}^{t, x}, M^{t, x}-\tilde{M}^{t, x}\right) \|_{\mathbb{S}(t, T)} \\
& \leq C\left\{\left\|b\left(\cdot, \tilde{X}^{t, x}, \phi\left(\cdot, \tilde{X}^{t, x}, \tilde{Y}^{t, x}\right)\right)-\tilde{b}\left(\cdot, \tilde{X}^{t, x}, \tilde{\phi}\left(\cdot, \tilde{X}^{t, x}, \tilde{Y}^{t, x}\right)\right)\right\|_{\mathcal{H}^{2}}\right. \\
& \quad+\left\|\partial_{x} H\left(\cdot, \tilde{X}^{t, x}, \phi\left(\cdot, \tilde{X}^{t, x}, \tilde{Y}^{t, x}\right), \tilde{Y}^{t, x}\right)-\partial_{x} \tilde{H}\left(\cdot, \tilde{X}^{t, x}, \tilde{\phi}\left(\cdot, \tilde{X}^{t, x}, \tilde{Y}^{t, x}\right), \tilde{Y}^{t, x}\right)\right\|_{\mathcal{H}^{2}} \\
&\left.+\|\sigma-\tilde{\sigma}\|_{L^{2}}+\left(\int_{0}^{T} \int_{\mathbb{R}_{0}^{p}}|\gamma(t, u)-\tilde{\gamma}(t, u)|^{2} \nu(\mathrm{~d} u) \mathrm{d} t\right)^{1 / 2}\right\} .
\end{aligned}
$$

It remains to estimate the first two terms on the right-hand side of the above inequality. By (H.2(1)),

$$
\begin{aligned}
& \left\|b\left(\cdot, \tilde{X}^{t, x}, \phi\left(\cdot, \tilde{X}^{t, x}, \tilde{Y}^{t, x}\right)\right)-\tilde{b}\left(\cdot, \tilde{X}^{t, x}, \tilde{\phi}\left(\cdot, \tilde{X}^{t, x}, \tilde{Y}^{t, x}\right)\right)\right\|_{\mathcal{H}^{2}} \\
& \leq\left\|b\left(\cdot, \tilde{X}^{t, x}, \phi\left(\cdot, \tilde{X}^{t, x}, \tilde{Y}^{t, x}\right)\right)-b\left(\cdot, \tilde{X}^{t, x}, \tilde{\phi}\left(\cdot, \tilde{X}^{t, x}, \tilde{Y}^{t, x}\right)\right)\right\|_{\mathcal{H}^{2}} \\
& \quad+\left\|b\left(\cdot, \tilde{X}^{t, x}, \tilde{\phi}\left(\cdot, \tilde{X}^{t, x}, \tilde{Y}^{t, x}\right)\right)-\tilde{b}\left(\cdot, \tilde{X}^{t, x}, \tilde{\phi}\left(\cdot, \tilde{X}^{t, x}, \tilde{Y}^{t, x}\right)\right)\right\|_{\mathcal{H}^{2}} \\
& \leq\left\|b_{2}\right\|_{L^{\infty}}\left\|\phi\left(\cdot, \tilde{X}^{t, x}, \tilde{Y}^{t, x}\right)-\tilde{\phi}\left(\cdot, \tilde{X}^{t, x}, \tilde{Y}^{t, x}\right)\right\|_{\mathcal{H}^{2}}+\left\|b_{0}-\tilde{b}_{0}\right\|_{L^{2}}+\left\|b_{1}-\tilde{b}_{1}\right\|_{L^{\infty}}\left\|\tilde{X}^{t, x}\right\|_{\mathcal{H}^{2}} \\
& \quad+\left\|b_{2}-\tilde{b}_{2}\right\|_{L^{\infty}}\left\|\tilde{\phi}\left(\cdot, \tilde{X}^{t, x}, \tilde{Y}^{t, x}\right)\right\|_{\mathcal{H}^{2}} .
\end{aligned}
$$

Note that by (3.2.10), for all $t \in[0, T]$ and $x, y \in \mathbb{R}^{n}, \phi(t, x, y)=\left(\partial_{z} f^{*}\right)\left(t, x,-b_{2}(t)^{\top} y\right)$ and $\tilde{\phi}(t, x, y)=\left(\partial_{z} f^{*}\right)\left(t, x,-\tilde{b}_{2}(t)^{\top} y\right)$, where $f^{*}$ is the function defined in (3.2.4). Hence, from the $1 / \lambda$-Lipschitz continuity of $\partial_{z} f^{*}(t, x, \cdot)$ (see the proof of Lemma 3.2.2),

$$
\begin{equation*}
\left\|\phi\left(\cdot, \tilde{X}^{t, x}, \tilde{Y}^{t, x}\right)-\tilde{\phi}\left(\cdot, \tilde{X}^{t, x}, \tilde{Y}^{t, x}\right)\right\|_{\mathcal{H}^{2}} \leq C\left\|b_{2}-\tilde{b}_{2}\right\|_{L^{\infty}}\left\|\tilde{Y}^{t, x}\right\|_{\mathcal{H}^{2}} \leq C(1+|x|) \mathcal{E}_{\mathrm{per}}, \tag{3.2.23}
\end{equation*}
$$

where the last inequality follows from the moment estimate of $\tilde{Y}^{t, x}$. Moreover, the regularity of $\tilde{\phi}$ (see Lemma 3.2.3) and the moment estimate of ( $\left.\tilde{X}^{t, x}, \tilde{Y}^{t, x}\right)$ imply that $\left\|\tilde{\phi}\left(\cdot, \tilde{X}^{t, x}, \tilde{Y}^{t, x}\right)\right\|_{\mathcal{H}^{2}} \leq$
$C(1+|x|)$, which shows that $\left\|b\left(\cdot, \tilde{X}^{t, x}, \phi\left(\cdot, \tilde{X}^{t, x}, \tilde{Y}^{t, x}\right)\right)-b\left(\cdot, \tilde{X}^{t, x}, \tilde{\phi}\left(\cdot, \tilde{X}^{t, x}, \tilde{Y}^{t, x}\right)\right)\right\|_{\mathcal{H}^{2}} \leq C(1+$ $|x|) \mathcal{E}_{\text {per }}$. By the definitions of $H$ and $\tilde{H}$, the Lipschitz continuity of $\partial_{x} f_{0}$ in (H.2(3)) and (3.2.23),

$$
\begin{aligned}
& \left\|\partial_{x} H\left(\cdot, \tilde{X}^{t, x}, \phi\left(\cdot, \tilde{X}^{t, x}, \tilde{Y}^{t, x}\right), \tilde{Y}^{t, x}\right)-\partial_{x} \tilde{H}\left(\cdot, \tilde{X}^{t, x}, \tilde{\phi}\left(\cdot, \tilde{X}^{t, x}, \tilde{Y}^{t, x}\right), \tilde{Y}^{t, x}\right)\right\|_{\mathcal{H}^{2}} \\
& \leq\left\|\left(b_{1}-\tilde{b}_{1}\right)^{\top} \tilde{Y}^{t, x}\right\|_{\mathcal{H}^{2}}+\left\|\partial_{x} f_{0}\left(\cdot, \tilde{X}^{t, x}, \phi\left(\cdot, \tilde{X}^{t, x}, \tilde{Y}^{t, x}\right)\right)-\partial_{x} f_{0}\left(\cdot, \tilde{X}^{t, x}, \tilde{\phi}\left(\cdot, \tilde{X}^{t, x}, \tilde{Y}^{t, x}\right)\right)\right\|_{\mathcal{H}^{2}} \\
& \leq C(1+|x|) \mathcal{E}_{\text {per }} .
\end{aligned}
$$

Thus, we have proved the stability estimate that $\|\left(X^{t, x}-\tilde{X}^{t, x}, Y^{t, x}-\tilde{Y}^{t, x}, Z^{t, x}-\tilde{Z}^{t, x}, M^{t, x}-\right.$ $\left.\tilde{M}^{t, x}\right) \|_{\mathbb{S}(t, T)} \leq C(1+|x|) \mathcal{E}_{\text {per }}$.

We now establish the stability of feedback controls. By (3.2.10) and the $1 / \lambda$-Lipschitz continuity of $\partial_{z} f^{*}(t, x, \cdot)$, for all $(t, x) \in[0, T] \times \mathbb{R}^{n}$,

$$
\begin{aligned}
& \left.|\psi(t, x)-\tilde{\psi}(t, x)|=\mid\left(\partial_{z} f^{*}\right)\left(t, x,-b_{2}(t)^{\top} Y_{t}^{t, x}\right)-\left(\partial_{z} f^{*}\right)\left(t, x,-\tilde{b}_{2}(t)^{\top} \tilde{Y}_{t}^{t, x}\right)\right) \mid \\
& \leq\left|b_{2}(t)^{\top} Y_{t}^{t, x}-\tilde{b}_{2}(t)^{\top} \tilde{Y}_{t}^{t, x}\right| / \lambda \leq C\left(\left\|b_{2}-\tilde{b}_{2}\right\|_{L^{\infty}}\left|Y_{t}^{t, x}\right|+\left|Y_{t}^{t, x}-\tilde{Y}_{t}^{t, x}\right|\right) \\
& \leq C\left(\left\|b_{2}-\tilde{b}_{2}\right\|_{L^{\infty}}\left\|Y^{t, x}\right\|_{\mathcal{S}^{2}}+\left\|Y_{t}^{t, x}-\tilde{Y}_{t}^{t, x}\right\|_{\mathcal{S}^{2}}\right) \leq C(1+|x|) \mathcal{E}_{\text {per }} .
\end{aligned}
$$

An important application of the Lipschitz stability of feedback controls (Theorem 3.2.6) is the analysis of model misspecification error of a given learning algorithm. One essential component is to examine the performance of the feedback control $\tilde{\psi}$, computed based on the control problem (3.2.1) with the perturbed coefficients ( $\tilde{b}, \tilde{\sigma}, \tilde{\gamma}, f, g$ ), on the true model with coefficients $(b, \sigma, \gamma, f, g)$. For any given $x_{0} \in \mathbb{R}^{n}$, implementing the feedback control $\tilde{\psi}$ on the original system (3.2.2) will lead to the sub-optimal cost:

$$
\begin{equation*}
J\left(\tilde{\psi} ; x_{0}\right):=\mathbb{E}\left[\int_{0}^{T} f\left(t, X_{t}^{x_{0}, \tilde{\psi}}, \tilde{\psi}\left(t, X_{t}^{x_{0}, \tilde{\psi}}\right)\right) \mathrm{d} t+g\left(X_{T}^{x_{0}, \tilde{\psi}}\right)\right] \tag{3.2.24}
\end{equation*}
$$

where $X^{x_{0}, \tilde{\psi}} \in \mathcal{S}^{2}\left(\mathbb{R}^{n}\right)$ is the state process (with coefficients $b, \sigma$ and $\gamma$ ) associated with $\tilde{\psi}$ (see Definition 3.2.1). The following theorem shows that the difference between this suboptimal cost $J\left(\tilde{\psi} ; x_{0}\right)$ and the optimal cost $V$ in (3.2.1) depends Lipschitz-continuously on the magnitude of perturbations in the coefficients.

Theorem 3.2.7. Suppose (H.2) and (H.3) hold. Let $\psi \in \mathcal{V}$ (resp. $\tilde{\psi} \in \mathcal{V}$ ) be defined in (3.2.17) (resp. (3.2.19)), and for each $x_{0} \in \mathbb{R}^{n}$, let $X^{x_{0}, \psi} \in \mathcal{S}^{2}\left(\mathbb{R}^{n}\right)$ (resp. $X^{x_{0}, \tilde{\psi}} \in \mathcal{S}^{2}\left(\mathbb{R}^{n}\right)$ ) be the state process (3.2.2) associated with $\psi$ (resp. $\tilde{\psi}$ ), and let $V\left(x_{0}\right)$ (resp. $J\left(\tilde{\psi} ; x_{0}\right)$ ) be defined in (3.2.1) (resp. (3.2.24)). Then there exists a constant $C$ such that for all $x_{0} \in \mathbb{R}^{n}$, $\left\|X^{x_{0}, \psi}-X^{x_{0}, \tilde{\psi}}\right\|_{\mathcal{S}^{2}} \leq C\left(1+\left|x_{0}\right|\right) \mathcal{E}_{\text {per }}$ and $\left|V\left(x_{0}\right)-J\left(\tilde{\psi} ; x_{0}\right)\right| \leq C\left(1+\left|x_{0}\right|^{2}\right) \mathcal{E}_{\text {per }}$, with the constant $\mathcal{E}_{\text {per }}$ defined in (3.2.22).

To prove Theorem 3.2.7, we first establish that the composition of $f$ and the optimal feedback control is Lipschitz continuous, even though the cost function $f$ is merely lower
semicontinuous in the control variable (cf. (H.2(3))). The proof is based on the FenchelYoung identity:

$$
f\left(t, x, \partial_{z} f^{*}(t, x, z)\right)=\left\langle z, \partial_{z} f^{*}(t, x, z)\right\rangle-f^{*}(t, x, z) \in \mathbb{R}, \quad \forall(t, x, z) \in[0, T] \times \mathbb{R}^{n} \times \mathbb{R}^{k}
$$

the regularity of $f^{*}$ and Theorem 3.2.6, and has been given in Section 3.6.2.
Lemma 3.2.8. Suppose (H.2) and (H.3) hold. Let $\psi, \tilde{\psi}:[0, T] \times \mathbb{R}^{n} \rightarrow \mathbb{R}^{k}$ be the functions defined in (3.2.17) and (3.2.19), respectively. Then there exists a constant $C$ such that for all $t \in[0, T], x, x^{\prime} \in \mathbb{R}^{n},\left|f(t, x, \psi(t, x))-f\left(t, x^{\prime}, \tilde{\psi}\left(t, x^{\prime}\right)\right)\right| \leq C\left(\left(1+|x|+\left|x^{\prime}\right|\right)\left|x-x^{\prime}\right|+(1+\right.$ $\left.|x|^{2}+\left|x^{\prime}\right|^{2}\right) \mathcal{E}_{\text {per }}$ ), where the constant $\mathcal{E}_{\text {per }}$ is defined in (3.2.22).

Proof of Theorem 3.2.7. According to Theorems 3.2.5 and 3.2.6, there exists a constant $C$ such that for all $t \in[0, T], x, x^{\prime} \in \mathbb{R}^{n},|\psi(t, 0)|+|\tilde{\psi}(t, 0)| \leq C,\left|\psi(t, x)-\psi\left(t, x^{\prime}\right)\right|+\mid \tilde{\psi}(t, x)-$ $\tilde{\psi}\left(t, x^{\prime}\right)|\leq C| x-x^{\prime} \mid$, and $|\psi(t, x)-\tilde{\psi}(t, x)| \leq C(1+|x|) \mathcal{E}_{\text {per }}$. Then, for any given $x_{0} \in \mathbb{R}^{n}$, standard moment and stability estimates of (3.2.7) yield $\left\|X^{x_{0}, \psi}\right\|_{\mathcal{S}^{2}}+\left\|X^{x_{0}, \dot{\psi}}\right\|_{\mathcal{S}^{2}} \leq C\left(1+\left|x_{0}\right|\right)$ and

$$
\begin{aligned}
\left\|X^{x_{0}, \psi}-X^{x_{0}, \tilde{\psi}}\right\|_{\mathcal{S}^{2}} & \leq C\left\|b\left(\cdot, X^{x_{0}, \tilde{\psi}}, \psi\left(\cdot, X^{x_{0}, \tilde{\psi}}\right)\right)-b\left(\cdot, X^{x_{0}, \tilde{\psi}}, \tilde{\psi}\left(\cdot, X^{x_{0}, \tilde{\psi}}\right)\right)\right\|_{\mathcal{H}^{2}} \\
& \leq C\left\|\psi\left(\cdot, X^{x_{0}, \tilde{\psi}}\right)-\tilde{\psi}\left(\cdot, X^{x_{0}, \tilde{\psi}}\right)\right\|_{\mathcal{H}^{2}} \leq C\left(1+\left\|X^{x_{0}, \tilde{\psi}}\right\|_{\mathcal{H}^{2}}\right) \mathcal{E}_{\text {per }} \leq C\left(1+\left|x_{0}\right|\right) \mathcal{E}_{\text {per }}
\end{aligned}
$$

We now proceed to estimate $\left|V\left(x_{0}\right)-\tilde{V}\left(x_{0}\right)\right|$ for any given $x_{0} \in \mathbb{R}^{n}$. By the mean value theorem, (H.2(2)) and the Cauchy-Schwarz inequality,

$$
\begin{aligned}
\mathbb{E}\left[\left|g\left(X_{T}^{x_{0}, \psi}\right)-g\left(X_{T}^{x_{0}, \tilde{\psi}}\right)\right|\right] & \leq C \mathbb{E}\left[\left|\left(1+\left|X_{T}^{x_{0}, \psi}\right|+\left|X_{T}^{x_{0}, \tilde{\psi}}\right|\right)\right| X_{T}^{x_{0}, \psi}-X_{T}^{x_{0}, \tilde{\psi}} \mid\right] \\
& \leq C\left(1+\left\|X_{T}^{x_{0}, \psi}\right\|_{L^{2}}+\left\|X_{T}^{x_{0}, \tilde{\psi}^{2}}\right\|_{L^{2}}\right)\left\|X_{T}^{x_{0}, \psi}-X_{T}^{x_{0}, \tilde{\psi}}\right\|_{L^{2}} \\
& \leq C\left(1+\left|x_{0}\right|^{2}\right) \mathcal{E}_{\text {per }} .
\end{aligned}
$$

Moreover, from Lemma 3.2.8 and the Cauchy-Schwarz inequality,

$$
\begin{aligned}
& \mathbb{E}\left[\int_{0}^{T}\left|f\left(t, X_{t}^{x_{0}, \psi}, \psi\left(t, X_{t}^{x_{0}, \psi}\right)\right)-f\left(t, X_{t}^{x_{0}, \tilde{\psi}}, \tilde{\psi}\left(t, X_{t}^{x_{0}, \tilde{\psi}}\right)\right)\right| \mathrm{d} t\right] \\
& \leq C \mathbb{E}\left[\int_{0}^{T}\left(\left(1+\left|X_{t}^{x_{0}, \psi}\right|+\left|X_{t}^{x_{0}, \tilde{\psi}}\right|\right)\left|X_{t}^{x_{0}, \psi}-X_{t}^{x_{0}, \tilde{\psi}}\right|+\left(1+\left|X_{t}^{x_{0}, \psi}\right|^{2}+\left|X_{t}^{x_{0}, \tilde{\psi}}\right|^{2}\right) \mathcal{E}_{\text {per }}\right) \mathrm{d} t\right] \\
& \leq C\left(\left(1+\left\|X^{x_{0}, \psi}\right\|_{\mathcal{H}^{2}}+\left\|X^{x_{0}, \tilde{\psi}}\right\|_{\mathcal{H}^{2}}\right)\left\|X^{x_{0}, \psi}-X^{x_{0}, \tilde{\psi}}\right\|_{\mathcal{H}^{2}}+\left(1+\left\|X^{x_{0}, \psi}\right\|_{\mathcal{H}^{2}}^{2}+\left\|X^{x_{0}, \tilde{\psi}}\right\|_{\mathcal{H}^{2}}^{2}\right) \mathcal{E}_{\text {per }}\right) \\
& \leq C\left(1+\left|x_{0}\right|^{2}\right) \mathcal{E}_{\text {per }} .
\end{aligned}
$$

Since $\psi$ is an optimal feedback control of (3.2.1) with the initial state $x_{0} \in \mathbb{R}^{n}$, the desired estimate $\left|V\left(x_{0}\right)-J\left(\tilde{\psi} ; x_{0}\right)\right| \leq C\left(1+\left|x_{0}\right|^{2}\right) \mathcal{E}_{\text {per }}$ follows.

### 3.3 Regret analysis for linear-convex reinforcement learning

The focus of this section is the linear-convex reinforcement learning (RL) problem, where the drift coefficient of the state dynamics (3.2.2) is unknown to the controller, and the objective is to control the system optimally while simultaneously learning the dynamics. We shall propose a greedy least-squares algorithm to solve such problems, and show that the algorithm provides a sublinear regret with high probability guarantees. The analysis of the regret bounds for the algorithm relies on the Lipschitz stability of feedback controls established in Section 3.2.1.

### 3.3.1 Reinforcement learning problem and least-squares algorithm

The RL problem goes as follows. Let $x_{0} \in \mathbb{R}^{n}$ be a given initial state and $\theta^{\star}=\left(A^{\star}, B^{\star}\right) \in$ $\mathbb{R}^{n \times(n+k)}$ be fixed but unknown constants, consider the following problem:
$V\left(x_{0} ; \theta^{\star}\right)=\inf _{\alpha \in \mathcal{H}^{2}\left(\mathbb{R}^{k}\right)} J^{\theta^{\star}}\left(\alpha ; x_{0}\right), \quad$ with $\quad J^{\theta^{\star}}\left(\alpha ; x_{0}\right)=\mathbb{E}\left[\int_{0}^{T} f\left(t, X_{t}^{x_{0}, \theta^{\star}, \alpha}, \alpha_{t}\right) \mathrm{d} t+g\left(X_{T}^{x_{0}, \theta^{\star}, \alpha}\right)\right]$,
where for each $\alpha \in \mathcal{H}^{2}\left(\mathbb{R}^{k}\right)$, the process $X^{x_{0}, \theta^{\star}, \alpha}$ satisfies the following controlled dynamics associated with the parameter $\theta^{\star}$ :

$$
\begin{equation*}
\mathrm{d} X_{t}=\left(A^{\star} X_{t}+B^{\star} \alpha_{t}\right) \mathrm{d} t+\sigma \mathrm{d} W_{t}+\int_{\mathbb{R}_{0}^{p}} \gamma(u) \tilde{N}(\mathrm{~d} t, \mathrm{~d} u), \quad t \in[0, T], \quad X_{0}=x_{0} \tag{3.3.2}
\end{equation*}
$$

with a given constant $\sigma \in \mathbb{R}^{n \times d}$ and given functions $\gamma: \mathbb{R}_{0}^{p} \rightarrow \mathbb{R}^{n}, f:[0, T] \times \mathbb{R}^{n} \times \mathbb{R}^{k} \rightarrow$ $\mathbb{R} \cup\{\infty\}$ and $g: \mathbb{R}^{n} \rightarrow \mathbb{R}$. If $\theta^{\star}=\left(A^{\star}, B^{\star}\right)$ were known, then (3.3.1) is a control problem.

It is clear that (3.3.1)-(3.3.2) is a special case of (3.2.1)-(3.2.2) with $b(t, x, a)=A^{\star} x+B^{\star} a$, $\sigma(t)=\sigma$ and $\gamma(t, u)=\gamma(u)$ for all $(t, x, a, u) \in[0, T] \times \mathbb{R}^{n} \times \mathbb{R}^{k} \times \mathbb{R}_{0}^{p}$. Hence, if $f$ and $g$ satisfy (H.2) with for some $L \geq 0$ and $\lambda>0$, then (3.3.1)-(3.3.2) admits an optimal feedback control $\psi^{\theta^{\star}} \in \mathcal{V}$ as shown in Theorem 3.2.5. Note that to simplify the presentation, we assume that (3.3.2) has time homogenous coefficients as in [1, 112, 17], but similar analysis can be performed if the drift is a linear combination of given time-and-space-dependent basis functions with unknown weights or the diffusion/jump coefficients are also unknown.

To solve (3.3.1)-(3.3.2) with unknown $\theta^{\star}$, in an episodic reinforcement learning framework, the controller improves their knowledge of the parameter $\theta^{\star}$ through successive learning episodes. In particular, for each episode $i \in \mathbb{N}$, based on her observations in the past episodes, the controller executes a suitable control policy in $\psi_{i} \in \mathcal{V}$, whose associated state dynamics (3.3.2) leads to an expected cost $J^{\theta^{\star}}\left(\psi_{i} ; x_{0}\right)$. To measure the performance of an learning algorithm in this setting, one widely adopted criteria is the (expected) regret of the
algorithm defined as follows (see e.g. [43, 17]):

$$
\begin{equation*}
R(N)=\sum_{i=1}^{N}\left(J^{\theta^{\star}}\left(\psi_{i} ; x_{0}\right)-V\left(x_{0} ; \theta^{\star}\right)\right), \quad \forall N \in \mathbb{N} \tag{3.3.3}
\end{equation*}
$$

where $N$ denotes the total number of learning episodes. Intuitively, this regret characterizes the cumulative loss from taking sub-optimal policies in all episodes.

To start, let us consider a greedy algorithm, which chooses the optimal feedback control based on the current estimation of the parameter, and provides a sublinear regret with respect to the number of episodes $N$. More precisely, let $\theta=(A, B) \in \mathbb{R}^{n \times(n+k)}$ be the current estimate of $\theta^{\star}$, then the controller would exercise the optimal feedback control $\psi^{\theta} \in \mathcal{V}$ defined in Theorem 3.2.5 for the control problem (3.3.1)-(3.3.2) with $\theta^{\star}$ replaced by $\theta$, which leads to the state process $X^{x_{0}, \theta} \in \mathcal{S}^{2}\left(\mathbb{R}^{n}\right)$ satisfying:

$$
\begin{equation*}
\mathrm{d} X_{t}=\left(A^{\star} X_{t}+B^{\star} \psi^{\theta}\left(t, X_{t}\right)\right) \mathrm{d} t+\sigma \mathrm{d} W_{t}+\int_{\mathbb{R}_{0}^{p}} \gamma(u) \tilde{N}(\mathrm{~d} t, \mathrm{~d} u), \quad t \in[0, T], \quad X_{0}=x_{0} . \tag{3.3.4}
\end{equation*}
$$

By the martingale properties of stochastic integrals, we can then estimate $\theta^{\star}$ based on the process $Z_{t}^{x_{0}, \theta}:=\binom{X_{t}^{x_{0}, \theta}}{\psi^{\theta}\left(t, X_{t}^{x_{0}, \theta}\right)}, t \in[0, T]$, as follows:

$$
\begin{equation*}
\left(\theta^{\star}\right)^{\top}=\left(\mathbb{E}\left[\int_{0}^{T} Z_{t}^{x_{0}, \theta}\left(Z_{t}^{x_{0}, \theta}\right)^{\top} \mathrm{d} t\right]\right)^{-1} \mathbb{E}\left[\int_{0}^{T} Z_{t}^{x_{0}, \theta}\left(\mathrm{~d} X_{t}^{x_{0}, \theta}\right)^{\top}\right], \tag{3.3.5}
\end{equation*}
$$

provided that $\mathbb{E}\left[\int_{0}^{T} Z_{t}^{x_{0}, \theta}\left(Z_{t}^{x_{0}, \theta}\right)^{\top} \mathrm{d} t\right] \in \mathbb{R}^{(n+k) \times(n+k)}$ is invertible. This motivates us to introduce an iterative procedure to estimate $\theta^{\star}$, where the expectations in (3.3.5) are replaced by empirical averages over independent realizations. More precisely, let $m \in \mathbb{N}$ and $\left(X_{t}^{x_{0}, \theta, i}, \psi^{\theta}\left(t, X_{t}^{x_{0}, \theta, i}\right)\right)_{t \in[0, T]}, i=1, \ldots, m$, be trajectories of $m$ independent realizations of the state and control processes, we shall update the estimate $\theta$, denoted by $\hat{\theta}$, according to (3.3.5):

$$
\begin{equation*}
\hat{\theta}^{\top}:=\left(\frac{1}{m} \sum_{i=1}^{m} \int_{0}^{T} Z_{t}^{x_{0}, \theta, i}\left(Z_{t}^{x_{0}, \theta, i}\right)^{\top} \mathrm{d} t+\frac{1}{m} \mathbb{I}_{n+k}\right)^{-1}\left(\frac{1}{m} \sum_{i=1}^{m} \int_{0}^{T} Z_{t}^{x_{0}, \theta, i}\left(\mathrm{~d} X_{t}^{x_{0}, \theta, i}\right)^{\top}\right) \tag{3.3.6}
\end{equation*}
$$

where $Z_{t}^{x_{0}, \theta, i}:=\binom{X_{t}^{x_{0}, \theta, i}}{\psi^{\theta}\left(t, X_{t}^{x_{0}, \theta, i}\right)}$ for all $t \in[0, T]$ and $i=1, \ldots, m$, and $\mathbb{I}$ is the $(n+k) \times(n+k)$ identity matrix used to ensure the existence of the required matrix inverse. This leads to the following greedy least-squares (GLS) algorithm:

```
Algorithm 10 Greedy least-squares (GLS) algorithm
    Input: Choose an initial estimation \(\theta_{0}\) of \(\theta^{\star}\) and numbers of learning episodes
    \(\left\{m_{\ell}\right\}_{\ell \in \mathbb{N} \cup\{0\}}\).
    for \(\ell=0,1, \cdots\) do
        Obtain the optimal feedback control \(\psi^{\theta_{\ell}}\) for (3.3.1)-(3.3.2) with \(\theta^{\star}=\theta_{\ell}\) as in Theorem
        3.2.5.
        Execute the feedback control \(\psi^{\theta_{\ell}}\) for \(m_{\ell}\) indepenent episodes, and collect the trajectory
        data \(\left(X_{t}^{x_{0}, \theta_{\ell}, i}, \psi^{\theta_{\ell}}\left(t, X_{t}^{x_{0}, \theta_{\ell}, i}\right)\right)_{t \in[0, T]}, i=1, \ldots, m_{\ell}\).
        Obtain an updated estimation \(\theta_{\ell+1}\) by using (3.3.6) and the \(m_{\ell}\) trajectories collected
        above.
    end for
```


### 3.3.2 Structural assumptions for learning problems

In this section, we analyze the regret of Algorithm 10 based on the following assumptions of the learning problem (3.3.1)-(3.3.2).
H.4. (1) Let $x_{0} \in \mathbb{R}^{n}$, $\theta^{\star}=\left(A^{\star}, B^{\star}\right) \in \mathbb{R}^{n \times(n+k)}, \sigma \in \mathbb{R}^{n \times d}, \gamma: \mathbb{R}_{0}^{p} \rightarrow \mathbb{R}^{n}, f:[0, T] \times \mathbb{R}^{n} \times$ $\mathbb{R}^{k} \rightarrow \mathbb{R} \cup\{\infty\}$ and $g: \mathbb{R}^{n} \rightarrow \mathbb{R}$ satisfy (H.2) with some constants $L \geq 0$ and $\lambda>0$.
(2) There exist $\gamma_{\max } \geq 0$ and $\vartheta \in[0,1]$ such that $\sup _{q \geq 2} q^{-\vartheta}\left(\int_{\mathbb{R}_{0}^{p}}|\gamma(u)|^{q} \nu(\mathrm{~d} u)\right)^{1 / q} \leq \gamma_{\max }$.

Remark 3.3.1. Condition (H.4(1)) implies that for each $\theta=(A, B)$, the control problem of (3.3.1)-(3.3.2) with $\theta^{\star}$ replaced by $\theta$ is a nonsmooth linear-convex control problem studied in Section 3.2.

Condition (H.4(2)) describes the large jumps of the pure jump process

$$
L_{t}:=\int_{0}^{t} \int_{\mathbb{R}_{0}^{p}} \gamma(u) \tilde{N}(\mathrm{~d} s, \mathrm{~d} u), \quad t \in[0, T],
$$

which enables estimating the tail behaviour of the state process $X^{\theta}$, and subsequently quantifying the parameter estimation error of the least-squares estimator (3.3.6) (see Section 3.3.4.2). If the jump coefficient $\gamma$ is bounded, then one can easily see from $\int_{\mathbb{R}_{0}^{p}}|\gamma(u)|^{2} \nu(\mathrm{~d} u)<\infty$ that (H.4(2)) holds with $\vartheta=0$. Another important case is when $\gamma(u)=u$ for all $u \in \mathbb{R}_{0}^{p}$, under which the process $\left(L_{t}\right)_{t \in[0, T]}$ is a Lévy process of pure jumps with Lévy measure $\nu(\mathrm{d} u)$. In this case, (H.4(2)) holds with $\vartheta \in(0,1]$ if and only if $\left(\int_{\mathbb{R}_{0}^{p}}|u|^{q} \nu(\mathrm{~d} u)\right)^{1 / q} \leq \mathcal{O}\left(q^{\theta}\right)$ as $q \rightarrow \infty$.
H.5. $\theta^{\star}$ is identifiable, i.e., the optimal control $\alpha^{x_{0}, \star} \in \mathcal{H}^{2}\left(\mathbb{R}^{k}\right)$ and the optimal state process $X^{x_{0}, \theta^{\star}, \alpha^{\star}} \in \mathcal{S}^{2}\left(\mathbb{R}^{n}\right)$ of (3.3.1)-(3.3.2) (with initial state $x_{0}$ and parameter $\theta^{\star}$ ) satisfy the following linear independence condition: if $u_{1} \in \mathbb{R}^{n}$ and $u_{2} \in \mathbb{R}^{k}$ satisfy $u_{1}^{\top} X_{t}^{x_{0}, \theta^{\star}, \alpha^{\star}}+$ $u_{2}^{\top} \alpha_{t}^{x_{0}, \star}=0$ for $\mathrm{d} \mathbb{P} \otimes \mathrm{d} t$ a.e., then $u_{1}$ and $u_{2}$ are zero vectors.

Condition (H.5) implies that the true parameter $\theta^{\star}$ can be uniquely identified if we observe sufficiently many trajectories of the optimal state and control processes of (3.3.1)-(3.3.2). Such a self-exploration property allows us to design exploration-free learning algorithms for (3.3.1)-(3.3.2).

The following proposition shows that if the laws of the state processes are supported on the whole space, then (H.5) is equivalent to a self-exploration property of the optimal feedback control. The proof essentially follows the argument of [148, Lemma 6.1], and hence is omitted.

Proposition 3.3.1. Assume (H.4(1)). Let $\psi \in \mathcal{V}$. Assume that for all $t \in(0, T]$, and any open set $O \subset \mathbb{R}^{n}$ with positive Lebesgue measure, the state process $X^{\theta^{*}, \psi}$ (defined by (3.3.4) with $\left.\psi^{\theta}=\psi\right)$ satisfies that $\mathbb{P}\left(\left\{\omega \in \Omega \mid X_{t}^{\theta^{*}, \psi}(\omega) \in O\right\}\right)>0$. Then the following two statements are equivalent:
(a) if $u_{1} \in \mathbb{R}^{n}$ and $u_{2} \in \mathbb{R}^{k}$ satisfy $u_{1}^{\top} X_{t}^{\theta^{\star}, \psi}+u_{2}^{\top} \psi\left(t, X_{t}^{\theta^{\star}, \psi}\right)=0$ for $\mathrm{d} \mathbb{P} \otimes \mathrm{d} t$ a.e., then $u_{1}$ and $u_{2}$ are zero vectors;
(b) if $u_{1} \in \mathbb{R}^{n}$ and $u_{2} \in \mathbb{R}^{k}$ satisfy $u_{1}^{\top} x+u_{2}^{\top} \psi(t, x)=0$ for almost every $(t, x) \in[0, T] \times \mathbb{R}^{n}$, then $u_{1}$ and $u_{2}$ are zero vectors.

Consequently, suppose that (H.4(1)) holds and $\sigma \sigma^{\top}$ is positive definite, then (H.5) holds if and only if the optimal feedback control $\psi^{\theta^{\star}}$ of (3.3.1) satisfies Item (b).

Proposition 3.3.1 allows for more explicit expressions of (H.5). For instance, as shown in [17, Proposition 3.9], for quadratic cost functions $g=0$ and $f(t, x, a)=x^{\top} Q x+a^{\top} R a$ with positive definite matrices $Q$ and $R$, (H.5) holds if and only if $B^{\star}$ in (3.3.2) is full column rank. Alternatively, by [148, Proposition 3.3], if (3.3.1)-(3.3.2) has a bounded action set, i.e., $\mathcal{R}$ in (H.2(3)) has a bounded domain $\operatorname{dom} \mathcal{R}$ (cf. Example 3.2.1), then (H.5) holds if and only if the range of $\psi^{\theta^{*}}$ contains $k$ linearly independent vectors.

We remark that for general linear-convex learning problems without (H.5), an explicit exploration is necessary for learning [148]. Instead of merely employing greedy polices as in Algorithm 10, they dedicate certain episodes to actively explore the environment with some exploration policy $\psi^{e}$ satisfying Proposition 3.3.1 Item (b). The numbers of exploration and exploitation episodes are then balanced based on the performance gap in Theorem 3.2.7 and the finite-sample accuracy of the parameter estimator. Note, however, this explicit exploration may yield larger regrets for algorithm in [148] than that in Theorem 3.3.2.

### 3.3.3 Main results on sublinear regret bounds

We now state the main result which shows that the regret of Algorithm 10 grows at most sublinearly with respect to the number of episodes, provided that the hyper-parameters $\theta_{0}$ and $\left\{m_{j}\right\}_{j \in \mathbb{N} \cup\{0\}}$ are chosen properly. In particular, we shall choose an initial guess $\theta_{0}$ of $\theta^{\star}$ which satisfies the identifiability condition in (H.5) and we shall also double the number of learning episodes between two successive updates of the estimation of $\theta^{\star}$, which is a commonly
used strategy (the so-called doubling trick) in the design of online learning algorithms (see e.g. [17]). The proof of this theorem is given in Section 3.3.4.3.

To simplify the notation, we introduce the following quantifies for each $x_{0} \in \mathbb{R}^{n}, \theta=$ $(A, B) \in \mathbb{R}^{n \times(n+k)}$ and $m \in \mathbb{N}$ :

$$
\begin{align*}
\bar{U}^{x_{0}, \theta} & :=\mathbb{E}\left[\int_{0}^{T} Z_{t}^{x_{0}, \theta}\left(Z_{t}^{x_{0}, \theta}\right)^{\top} \mathrm{d} t\right], & \bar{V}^{x_{0}, \theta}:=\mathbb{E}\left[\int_{0}^{T} Z_{t}^{x_{0}, \theta}\left(\mathrm{~d} X_{t}^{x_{0}, \theta}\right)^{\top}\right] \\
U^{x_{0}, \theta, m} & :=\frac{1}{m} \sum_{i=1}^{m} \int_{0}^{T} Z_{t}^{x_{0}, \theta, i}\left(Z_{t}^{x_{0}, \theta, i}\right)^{\top} \mathrm{d} t, & V^{x_{0}, \theta, m}:=\frac{1}{m} \sum_{i=1}^{m} \int_{0}^{T} Z_{t}^{x_{0}, \theta, i}\left(\mathrm{~d} X_{t}^{x_{0}, \theta, i}\right)^{\top} \tag{3.3.7}
\end{align*}
$$

where $X^{x_{0}, \theta} \in \mathcal{S}^{2}\left(\mathbb{R}^{n}\right)$ is the solution of (3.3.4), $\left(X^{x_{0}, \theta, i}\right)_{i=1}^{m}$ are independent copies of $X^{x_{0}, \theta}$, and $Z^{x_{0}, \theta}$ and $\left(Z^{x_{0}, \theta, i}\right)_{i=1}^{m}$ are defined as in (3.3.5) and (3.3.6), respectively. For any given symmetric matrix $A$, we denote by $\lambda_{\min }(A)$ the smallest eigenvalue of $A$.

Theorem 3.3.2. Suppose (H.4(1)) and (H.5) hold. Assume further that $\lambda_{\min }\left(\bar{U}^{x_{0}, \theta_{0}}\right)>0$, and for any given bounded set $\mathcal{K} \subset \mathbb{R}^{n \times(n+k)}$, there exist constants $C_{1}, C_{2}>0$ and $\beta \geq 1$, such that the following concentration inequality holds for all $\varepsilon \geq 0, m \in \mathbb{N}$ and $\theta \in \mathcal{K}$,

$$
\begin{align*}
& \max \left\{\mathbb{P}\left(\left|U^{x_{0}, \theta, m}-\bar{U}^{x_{0}, \theta}\right| \geq \varepsilon\right), \mathbb{P}\left(\left|V^{x_{0}, \theta, m}-\bar{V}^{x_{0}, \theta}\right| \geq \varepsilon\right)\right\} \\
& \quad \leq C_{2} \exp \left(-C_{1} \min \left\{\frac{m \varepsilon^{2}}{C_{2}^{2}},\left(\frac{m \varepsilon}{C_{2}}\right)^{\frac{1}{\beta}}\right\}\right) \tag{3.3.8}
\end{align*}
$$

Then there exists a constant $C_{0}>0$, such that for all $C \geq C_{0}$ and $\delta \in(0,1 / 4)$, if we set $m_{0}=C(-\ln \delta)^{\beta}$ and $m_{\ell}=2^{\ell} m_{0}$ for all $\ell \in \mathbb{N}$, then the regret of Algorithm 10 (cf. (3.3.3)) satisfies the following properties:
(1) It holds with probability at least $1-4 \delta$ that $R(N) \leq C^{\prime}(\sqrt{N} \sqrt{\ln N}+\sqrt{-\ln \delta} \sqrt{N}+$ $\left.(-\ln \delta)^{\beta} \ln N\right)$ for all $N \in \mathbb{N}$, where $C^{\prime}$ is a constant independent of $N$ and $\delta$.
(2) It holds with probability 1 that $R(N)=\mathcal{O}(\sqrt{N \ln N})$ as $N \rightarrow \infty$.

The following theorem presents a precise sublinear regret bound of Algorithm 10 for the jump-diffusion model (3.3.2), depending on the jump sizes of the Poisson random measure. The proof follows from Theorem 3.3.2 and Proposition 3.3.9.

Theorem 3.3.3. Suppose (H.4) and (H.5) hold, and $\lambda_{\min }\left(\bar{U}^{x_{0}, \theta_{0}}\right)>0$. Then there exists a constant $C_{0}>0$, such that for all $C \geq C_{0}$ and $\delta \in(0,1 / 4)$, if we set $m_{0}=C(-\ln \delta)^{3+\vartheta}$ and $m_{\ell}=2^{\ell} m_{0}$ for all $\ell \in \mathbb{N}$, then the regret of Algorithm 10 (cf. (3.3.3)) satisfies the following properties:
(1) It holds with probability at least $1-4 \delta$ that $R(N) \leq C^{\prime}(\sqrt{N} \sqrt{\ln N}+\sqrt{-\ln \delta} \sqrt{N}+$ $\left.(-\ln \delta)^{3+\vartheta} \ln N\right)$ for all $N \in \mathbb{N}$, where $\vartheta$ is the constant in (H.4(2)) and $C^{\prime}$ is a constant independent of $\vartheta, N$ and $\delta$.
(2) It holds with probability 1 that $R(N)=\mathcal{O}(\sqrt{N \ln N})$ as $N \rightarrow \infty$.

In the case where (3.3.2) is only driven by the Brownian motion, we can exploit the sub-Gaussianity of the state process and obtain a shaper regret bound based on Theorem 3.3.2 and Proposition 3.3.10.

Theorem 3.3.4. Suppose (H.4) and (H.5) hold with $\gamma_{\max }=0$, and $\lambda_{\min }\left(\bar{U}^{x_{0}, \theta_{0}}\right)>0$. Then there exists a constant $C_{0}>0$, such that for all $C \geq C_{0}$ and $\delta \in(0,1 / 4)$, if we set $m_{0}=$ $C(-\ln \delta)$ and $m_{\ell}=2^{\ell} m_{0}$ for all $\ell \in \mathbb{N}$, then the regret of Algorithm 10 (cf. (3.3.3)) satisfies the following properties:
(1) It holds with probability at least $1-4 \delta$ that $R(N) \leq C^{\prime}(\sqrt{N} \sqrt{\ln N}+\sqrt{-\ln \delta} \sqrt{N}+$ $(-\ln \delta) \ln N)$ for all $N \in \mathbb{N}$, where $C^{\prime}$ is a constant independent of $N$ and $\delta$.
(2) It holds with probability 1 that $R(N)=\mathcal{O}(\sqrt{N \ln N})$ as $N \rightarrow \infty$.

Remark 3.3.2. The condition $\lambda_{\min }\left(\bar{U}^{x_{0}, \theta_{0}}\right)>0$ in Theorems 3.3.3 and 3.3.4 ensures that the greedy policy $\psi^{\theta_{0}}$ explores the parameter space and improves the accuracy of parameter estimation. By Proposition 3.3.1, if (3.3.4) has nondegenerate Brownian noises, then it suffices to choose $\theta_{0}$ such that the corresponding greedy policy $\psi^{\theta_{0}}$ enjoys the exploration property stated in Item (b).

The choice of $m_{0}=C_{0}(-\ln \delta)^{\beta}$ along with (H.5) ensures that $\left(\theta_{\ell}\right)_{\ell \in \mathbb{N}}$ stays close to $\theta^{\star}$ so that (3.3.8) is applicable. Here $\delta$ is an arbitrarily small constant indicating the agent's confidence of the regret bound, and $C_{0}$ is a constant depending on the exploration strength of $\psi^{\theta^{\star}}$, namely the constant $\lambda_{\min }\left(\bar{U}^{x_{0}, \theta^{\star}}\right)>0$ (see Section 3.3.4.3). Note that to analyze algorithm regrets, it is common to assume some a-priori information on the true parameter and the algorithm being initialized with sufficiently many learning episodes (see e.g., [45]). Obtaining an explicit dependence of $C_{0}$ on model parameters, however, could be challenging. A practical strategy for validating (H.5) and for choosing the initial episode $m_{0}$ is to ensure that the obtained estimations $\left(\theta_{\ell}\right)_{\ell \in \mathbb{N}}$ remain bounded and that the resulting greedy policies $\left(\psi^{\theta_{\ell}}\right)_{\ell \in \mathbb{N}}$ satisfy Proposition 3.3.1 Item (b). Our numerical experiments in Section 3.5 demonstrate that the performance of Algorithm 10 is stable with respect to $m_{0}$, and that a small $m_{0}$ in general suffices to guarantee a sublinear regret.

### 3.3.4 Proofs of sublinear regret bounds

This section is devoted to the proofs of Theorem 3.3.2, 3.3.3 and 3.3.4.
As we have seen in Theorems 3.3.3-3.3.4, an essential step for estimating the regret of Algorithm 10 is to establish the concentration inequality (3.3.8) for the least-squares estimator (3.3.6). Compared to the classical learning problems with Brownian-motion-driven state dynamics (see e.g. [17]), the presence of jumps in the state dynamics creates a crucial
difficulty in quantifying the precise value of $\beta$ in (3.3.8), since the state variable $X^{\theta}$ is in general not sub-Gaussian, and hence (3.3.8) does not hold with $\beta=1$.

In the subsequent analysis, we overcome the above difficulty by introducing a notation of sub-Weibull random variables as in [98] and establishing that both deterministic and stochastic integrals preserve sub-Weibull random variables in Section 3.3.4.1. We then show in Section 3.3.4.2 that (3.3.6) behaves like sub-Weibull random variables and (3.3.8) holds with some $\beta \geq 1$, provided that the jumps of the state dynamics are sub-exponential. Finally, we prove the general regret result Theorem 3.3.2 for Algorithm 10 in Section 3.3.4.3.

### 3.3.4.1 Step 1: Analysis of sub-Weibull random variables

The first step is to analyze integrals of sub-Weibull random variables. We start by recalling the precise definition of sub-Weibull random variables in terms of their Orlicz norms (see [98]).

Definition 3.3.1. For every $\alpha>0$, let $\Psi_{\alpha}:[0, \infty) \rightarrow \mathbb{R}$ such that $\Psi_{\alpha}(x)=e^{x^{\alpha}}-1$ for all $x \geq 0$, and let $\|\cdot\|_{\Psi_{\alpha}}$ be the corresponding $\Psi_{\alpha}$-Orlicz (quasi-)norm such that for any given random variable $X$,

$$
\|X\|_{\Psi_{\alpha}}:=\inf \left\{t>0 \left\lvert\, \mathbb{E}\left[\Psi_{\alpha}\left(\frac{|X|}{t}\right)\right] \leq 1\right.\right\}
$$

Then a random variable $X$ is said to be sub-Weibull of order $\alpha>0$, denoted by $X \in$ $\operatorname{subW}(\alpha)$, if $\|X\|_{\Psi_{\alpha}}<\infty$.

Note that $\|\cdot\|_{\Psi_{\alpha}}$ is a norm if and only if $\alpha \geq 1$, as otherwise the triangle inequality does not hold. Examples of sub-Weibull random variables include sub-Gaussian and subexponential random variables, which correspond to subW (2) and subW (1), respectively. We point out that the class of sub-Weibull random variables is closed under multiplication and addition, and for all $\alpha>0$, there exists a constant $C_{\alpha}$, depending only on $\alpha$, such that

$$
\begin{equation*}
C_{\alpha}^{-1} \sup _{q \geq 1} q^{-1 / \alpha}\|X\|_{L^{q}} \leq\|X\|_{\Psi_{\alpha}} \leq C_{\alpha} \sup _{q \geq 1} q^{-1 / \alpha}\|X\|_{L^{q}} \tag{3.3.9}
\end{equation*}
$$

for all random variables $X$ (see [70, Appendix A] for a proof of these properties).
We now present several important lemmas regarding the behavior of integrals of subWeibull random variables. The first lemma shows that deterministic integral of a product of sub-Weibull random variables is still sub-Weibull. The proof is based on Definition 3.3.1 and Hölder's inequality, and is given in Section 3.6.2.

Lemma 3.3.5. For all $\alpha>0$ and every stochastic process $X, Y: \Omega \times[0, T] \rightarrow \mathbb{R}$,

$$
\left\|\int_{0}^{T} X Y \mathrm{~d} t\right\|_{\Psi_{\alpha / 2}} \leq\left\|\left(\int_{0}^{T}|X|^{2} \mathrm{~d} t\right)^{\frac{1}{2}}\right\|_{\Psi_{\alpha}}\left\|\left(\int_{0}^{T}|Y|^{2} \mathrm{~d} t\right)^{\frac{1}{2}}\right\|_{\Psi_{\alpha}}
$$

The second lemma shows that stochastic integrals preserve the property of being subWeibull random variables. The proof is based on the equivalent characterization (3.3.9) of sub-Weibull random variables and Burkholder's inequality, whose details are given in Section 3.6.2.

Lemma 3.3.6. There exists $C \geq 0$ such that for all $\sigma \in \mathbb{R}^{d}, X \in \mathcal{S}^{2}(\mathbb{R})$ and every measurable function $\gamma: \mathbb{R}_{0}^{p} \rightarrow \mathbb{R}$ satisfying (H.4(2)), $\left\|\int_{0}^{T} X_{t} \sigma^{\top} \mathrm{d} W_{t}\right\|_{\Psi_{1 / 2}} \leq C|\sigma|\left\|\left(\int_{0}^{T}|X|^{2} \mathrm{~d} t\right)^{\frac{1}{2}}\right\|_{\Psi_{1}}$ and

$$
\left\|\int_{0}^{T} \int_{\mathbb{R}_{0}^{p}} X_{t} \gamma(u) \tilde{N}(\mathrm{~d} t, \mathrm{~d} u)\right\|_{\Psi_{1 /(3+\vartheta)}} \leq C \gamma_{\max }\left(\sup _{p \geq 2}\left\|\left(\int_{0}^{T}\left|X_{t}\right|^{q} \mathrm{~d} t\right)^{\frac{1}{q}}\right\|_{\Psi_{1}}\right)
$$

with the constants $\gamma_{\max }$ and $\vartheta$ in (H.4(2)).
Lemma 3.3.6 focuses on the case where $\left(\int_{0}^{T}|X|^{2} \mathrm{~d} t\right)^{1 / 2} \in \operatorname{subW}(1) \backslash \operatorname{subW}(2)$, which is important for control problems whose state dynamics is driven by a Poisson random measure. Hence we establish the sub-Weibull properties of the stochastic integrals by applying the Burkholder's inequality to estimate the growth of their $L^{q}$-norms, precise order of which depends on the constants $C_{q}$ and $\tilde{C}_{q}$ in the inequalities (3.6.7) and (3.6.8).

In the case where $\left(\int_{0}^{T}|X|^{2} \mathrm{~d} t\right)^{1 / 2} \in \operatorname{subW}(2)$, we can establish the optimal sub-Weibull order $\int_{0}^{T} X_{t} \sigma^{\top} \mathrm{d} W_{t} \in \operatorname{subW}(1)$. Such a characeterization is essential for obtaining a sharper regret bound of Algorithm 10 when the state dynamics is only driven by the Brownian motion. The proof is based on the Girsanov theorem and is given in Section 3.6.2.

Lemma 3.3.7. There exists $C \geq 0$ such that for all $\sigma \in \mathbb{R}^{d}$ and $X \in \mathcal{S}^{2}(\mathbb{R})$,

$$
\left\|\int_{0}^{T} X_{t} \sigma^{\top} \mathrm{d} W_{t}\right\|_{\Psi_{1}} \leq C|\sigma|\left\|\left(\int_{0}^{T}|X|^{2} \mathrm{~d} t\right)^{\frac{1}{2}}\right\|_{\Psi_{2}}
$$

### 3.3.4.2 Step 2: Concentration inequalities for the least-squares estimator

Based on the fact that sub-Weibull properties are preserved under algebraic and integral operations as shown in Section 3.3.4.1, we now quantify the precise tail behavior of the least-squares estimator (3.3.6), namely the constant $\beta$ in (3.3.8), for the jump-diffusion model (3.3.2).

We start by establishing the sub-exponential properties of Lipschitz functionals of the state process $X^{\theta}$ driven by both Brownian motions and Poisson random measures as in (3.3.4). The proof follows as a special case of [111] and is given in Section 3.6.2.

Lemma 3.3.8. Suppose (H.4) holds. Let $K \in \mathbb{R}$ and $\theta=(A, B) \in \mathbb{R}^{n \times(n+k)}$ satisfy $|\theta| \leq K$. Then there exists $C \geq 0$, depending only on $K, T$ and the constants in (H.4), such that for all $x_{0} \in \mathbb{R}^{n}$ and for every Lipschitz continuous function $\mathfrak{f}:\left(\mathbb{D}\left([0, T] ; \mathbb{R}^{n}\right), d_{\infty}\right) \rightarrow \mathbb{R}$, the solution $X^{x_{0}, \theta}$ of (3.3.4) satisfies $\left\|\mathfrak{f}\left(X^{x_{0}, \theta}\right)\right\|_{\Psi_{1}} \leq C\left(\|\mathfrak{f}\|_{\text {Lip }}+\left|\mathbb{E}\left[\mathfrak{f}\left(X^{x_{0}, \theta}\right)\right]\right|\right)$, where $\mathbb{D}\left([0, T] ; \mathbb{R}^{n}\right)$ is the space of $\mathbb{R}^{n}$-valued càdlàg functions on $[0, T]$ endowed with the uniform metric $d_{\infty}$, and $\|\mathfrak{f}\|_{\text {Lip }}$ is the Lipschitz constant of $\mathfrak{f}$. (cf. Lemma 3.6.3).

We now characterize the parameter $\beta$ in the concentration inequality (3.3.8) based on Lemmas 3.3.5, 3.3.6 and 3.3.8.

Proposition 3.3.9. Suppose (H.4) holds and let $\mathcal{K} \subset \mathbb{R}^{n \times(n+k)}$ be a bounded set. Then there exist constants $C_{1}, C_{2} \geq 0$ such that (3.3.8) holds for all $\varepsilon \geq 0, m \in \mathbb{N}$ and $\theta \in \mathcal{K}$ with $\beta=3+\vartheta$, where $\vartheta$ is the constant in (H.4(2)).

Proof. Throughout this proof, let $\theta$ be a given constant satisfying $|\theta| \leq K$ for some $K \geq 0$. For notational simplicity, we shall omit the dependence on $\left(x_{0}, \theta\right)$ in the subscripts of all random variables, and denote by $C_{2}$ a generic constant, which is independent of $m$ and the precise value of $\theta$, and depends possibly on $K, x_{0}$, the constants in (H.4) and the dimensions.

Note that for each $i=1, \ldots, m$, the entries of $\int_{0}^{T} Z_{t}^{i}\left(Z_{t}^{i}\right)^{\top} \mathrm{d} t$ are one of the three cases:

$$
\begin{equation*}
\int_{0}^{T} X_{\ell, t}^{i} X_{j, t}^{i} \mathrm{~d} t, \quad \int_{0}^{T} X_{\ell, t}^{i} \psi^{\theta}\left(t, X_{t}^{i}\right)_{j} \mathrm{~d} t, \quad \int_{0}^{T} \psi^{\theta}\left(t, X_{t}^{i}\right)_{\ell} \psi^{\theta}\left(t, X_{t}^{i}\right)_{j} \mathrm{~d} t \tag{3.3.10}
\end{equation*}
$$

where $X_{\ell, t}^{i}$ and $\psi^{\theta}\left(t, X_{t}^{i}\right)_{\ell}$ are the $\ell$ th-entry of $X_{\ell, t}^{i}$ and $\psi^{\theta}\left(t, X_{t}^{i}\right)_{\ell}$, respectively. Similarly, the entries of $\int_{0}^{T} Z_{t}^{i}\left(\mathrm{~d} X_{t}^{i}\right)^{\top}$ are one of the two cases:

$$
\begin{align*}
& \int_{0}^{T} X_{\ell, t}^{i}\left(A^{\star} X_{t}^{i}\right)_{j} \mathrm{~d} t+\int_{0}^{T} X_{\ell, t}^{i}\left(B^{\star} \psi^{\theta}\left(t, X_{t}^{i}\right)\right)_{j} \mathrm{~d} t+\int_{0}^{T} X_{\ell, t}^{i} \sigma_{j} \mathrm{~d} W_{t}^{i}+\int_{0}^{T} \int_{\mathbb{R}_{0}^{p}} X_{\ell, t}^{i} \gamma(u)_{j} \tilde{N}^{i}(\mathrm{~d} t, \mathrm{~d} u), \\
& \int_{0}^{T} \psi^{\theta}\left(t, X_{t}^{i}\right)_{\ell}\left(A^{\star} X_{t}^{i}\right)_{j} \mathrm{~d} t+\int_{0}^{T} \psi^{\theta}\left(t, X_{t}^{i}\right)_{\ell}\left(B^{\star} \psi^{\theta}\left(t, X_{t}^{i}\right)\right)_{j} \mathrm{~d} t+\int_{0}^{T} \psi^{\theta}\left(t, X_{t}^{i}\right)_{\ell} \sigma_{j} \mathrm{~d} W_{t}^{i}  \tag{3.3.11}\\
& \quad+\int_{0}^{T} \int_{\mathbb{R}_{0}^{p}} \psi^{\theta}\left(t, X_{t}^{i}\right)_{\ell} \gamma(u)_{j} \tilde{N}^{i}(\mathrm{~d} t, \mathrm{~d} u),
\end{align*}
$$

where $\sigma_{j}$ is the $j$-th row of $\sigma \in \mathbb{R}^{n \times d}, \gamma_{j}$ is the $j$-th entry of the function $\gamma: \mathbb{R}_{0}^{p} \rightarrow \mathbb{R}^{n}$, $\left(W^{i}\right)_{i=1}^{m}$ are $m$-independent $d$-dimensional Brownian motion, and $\left(\tilde{N}^{i}\right)_{i=1}^{m}$ are $m$-independent compensated Poisson random measures. By the definitions of $U^{x_{0}, \theta, m}, V^{x_{0}, \theta, m}$ in (3.3.7), and the inequality that $\mathbb{P}\left(\left|\sum_{i=1}^{\ell} X_{i}\right| \geq \varepsilon\right) \leq \sum_{i=1}^{\ell} \mathbb{P}\left(\left|X_{i}\right| \geq \varepsilon / \ell\right)$ for all $\ell \in \mathbb{N}$ and random variables $\left(X_{i}\right)_{i=1}^{\ell}$, it suffices to obtain a concentration inequality for each term in (3.3.10) and (3.3.11).

Since $|\theta| \leq K$, by Theorem 3.2.5, there exists $C_{2} \geq 0$ such that $\left|\psi^{\theta}(t, 0)\right| \leq C_{2}$ and $\left|\psi^{\theta}(t, x)-\psi^{\theta}\left(t, x^{\prime}\right)\right| \leq C_{2}\left|x-x^{\prime}\right|$ for all $t \in[0, T], x, x^{\prime} \in \mathbb{R}^{n}$. Then standard moment estimates of (3.3.4) (with the initial condition $x_{0}$ ) shows that $\left\|X^{i}\right\|_{\mathcal{S}^{2}\left(\mathbb{R}^{n}\right)} \leq C_{2}$ for all $i=$ $1, \ldots, m$, with a constant $C_{2}$ depending on $x_{0}$. Then, for each $q \geq 2, \ell=1, \ldots, n$ and $j=1, \ldots k$, we consider the functions $\mathfrak{f}_{\ell}^{(q)}, \overline{\mathfrak{f}}_{j}^{(q)}:\left(\mathbb{D}\left([0, T] ; \mathbb{R}^{n}\right), d_{\infty}\right) \rightarrow \mathbb{R}$ satisfying for all $\rho \in \mathbb{D}\left([0, T] ; \mathbb{R}^{n}\right)$ that $\mathfrak{f}_{\ell}^{(q)}(\rho)=\left(\int_{0}^{T}\left|\rho_{\ell, t}\right|^{q} \mathrm{~d} t\right)^{\frac{1}{q}}$ and $\overline{\mathfrak{f}}_{j}^{(q)}(\rho)=\left(\int_{0}^{T}\left|\psi^{\theta}\left(t, \rho_{t}\right)_{j}\right|^{q} \mathrm{~d} t\right)^{\frac{1}{q}}$, where $\rho_{\ell, t}$ is the $\ell$ th component of $\rho_{t}$ and $\psi^{\theta}\left(t, \rho_{t}\right)_{j}$ is the $j$ th component of $\psi^{\theta}\left(t, \rho_{t}\right)$. One can easily show that $\mathfrak{f}_{\ell}^{(q)}(0)=0$ and $\left|\mathfrak{f}_{j}^{(q)}(0)\right|,\left\|\mathfrak{f}_{\ell}^{(q)}\right\|_{\text {Lip }},\left\|\overline{\mathfrak{f}}_{j}^{(q)}\right\|_{\text {Lip }} \leq C$, which along with Lemma 3.3.8 implies that $\left\|\left(\int_{0}^{T}\left|X_{\ell, t}^{i}\right|^{q} \mathrm{~d} t\right)^{\frac{1}{q}}\right\|_{\Psi_{1}} \leq C$ and $\left\|\left(\int_{0}^{T}\left|\psi^{\theta}\left(t, X_{t}^{i}\right)_{j}\right|^{q} \mathrm{~d} t\right)^{\frac{1}{q}}\right\|_{\Psi_{1}} \leq C$, uniformly with
respect to $i, \ell, j, q, \theta$. Hence, we can obtain from Lemmas 3.3.5 and 3.3.6 a uniform bound for the $\|\cdot\|_{\Psi_{1 /(3+\vartheta)}}$-norms of all the terms in (3.3.10) and (3.3.11).

Consequently, we can deduce the desired concentration inequality by applying Lemma 3.6.4 (with $\alpha=1 /(3+\vartheta), N=m$ and $\varepsilon^{\prime}=m \varepsilon$ ) to each component of the zero-mean random varables $\left(\int_{0}^{T} Z_{t}^{i}\left(Z_{t}^{i}\right)^{\top} \mathrm{d} t-\bar{U}\right)_{i=1}^{m}$ and $\left(\int_{0}^{T} Z_{t}^{i}\left(\mathrm{~d} X_{t}^{i}\right)^{\top}-\bar{V}\right)_{i=1}^{m}$.

The following proposition improves the concentration inequality in Proposition 3.3.9 for the case without jumps.

Proposition 3.3.10. Suppose (H.4) holds with $\gamma_{\max }=0$ and let $\mathcal{K} \subset \mathbb{R}^{n \times(n+k)}$ be a bounded set. Then there exist constants $C_{1}, C_{2} \geq 0$ such that (3.3.8) holds for all $\varepsilon \geq 0, m \in \mathbb{N}$ and $\theta \in \mathcal{K}$ with $\beta=1$.

Proof. We first refine the result of Lemma 3.3.8 and prove Lipschitz functionals of the state process $X^{x_{0}, \theta}$ is sub-Gaussian. By [49, Theorem 1.1 and Corollary 4.1], there exists $C \geq 0$ such that for all $x_{0} \in \mathbb{R}^{n}$ and for every Lipschitz continuous function $\mathfrak{f}:\left(\mathbb{D}\left([0, T] ; \mathbb{R}^{n}\right), d_{\infty}\right) \rightarrow$ $\mathbb{R}$ with $\|\mathfrak{f}\|_{\text {Lip }} \leq 1, \mathbb{E}\left[\exp \left(\lambda\left(\mathfrak{f}\left(X^{x_{0}, \theta}\right)-\mathbb{E}\left[\mathfrak{f}\left(X^{x_{0}, \theta}\right)\right]\right)\right)\right] \leq \exp \left(C^{2} \lambda^{2}\right)$ for all $\lambda>0$, which along with $\left[156\right.$, Proposition 2.5.2 (v)] implies that $\left\|\mathfrak{f}\left(X^{x_{0}, \theta}\right)-\mathbb{E}\left[\mathfrak{f}\left(X^{x_{0}, \theta}\right)\right]\right\|_{\Psi_{2}} \leq C$ for some constant $C$, uniformly with respect to $x_{0} \in \mathbb{R}^{n}, \theta \in \mathcal{K}$ and $\mathfrak{f}:\left(\mathbb{D}\left([0, T] ; \mathbb{R}^{n}\right), d_{\infty}\right) \rightarrow \mathbb{R}$ satisfying $\|\mathfrak{f}\|_{\text {Lip }} \leq 1$. Then, we can deduce from the fact that $\|\cdot\|_{\Psi_{2}}$ is a norm that $\left\|\mathfrak{f}\left(X^{x_{0}, \theta}\right)\right\|_{\Psi_{2}} \leq C\left(\|\mathfrak{f}\|_{\text {Lip }}+\left|\mathbb{E}\left[\mathfrak{f}\left(X^{x_{0}, \theta}\right)\right]\right|\right)$ for all $x_{0} \in \mathbb{R}^{n}, \theta \in \mathcal{K}$ and Lipschitz continuous functions $\mathfrak{f}$.

We then proceed along the proof of Proposition 3.3.9. For each $i=1, \ldots, m$, all entries of $\int_{0}^{T} Z_{t}^{i}\left(Z_{t}^{i}\right)^{\top} \mathrm{d} t$ are given in (3.3.10), and all entries of $\int_{0}^{T} Z_{t}^{i}\left(\mathrm{~d} X_{t}^{i}\right)^{\top}$ are given by (cf. (3.3.11)):

$$
\begin{align*}
& \int_{0}^{T} X_{\ell, t}^{i}\left(A^{\star} X_{t}^{i}\right)_{j} \mathrm{~d} t+\int_{0}^{T} X_{\ell, t}^{i}\left(B^{\star} \psi^{\theta}\left(t, X_{t}^{i}\right)\right)_{j} \mathrm{~d} t+\int_{0}^{T} X_{\ell, t}^{i} \sigma_{j} \mathrm{~d} W_{t}^{i} \\
& \int_{0}^{T} \psi^{\theta}\left(t, X_{t}^{i}\right)_{\ell}\left(A^{\star} X_{t}^{i}\right)_{j} \mathrm{~d} t+\int_{0}^{T} \psi^{\theta}\left(t, X_{t}^{i}\right)_{\ell}\left(B^{\star} \psi^{\theta}\left(t, X_{t}^{i}\right)\right)_{j} \mathrm{~d} t+\int_{0}^{T} \psi^{\theta}\left(t, X_{t}^{i}\right)_{\ell} \sigma_{j} \mathrm{~d} W_{t}^{i}, \tag{3.3.12}
\end{align*}
$$

for all $\ell=1, \ldots, n$ and $j=1, \ldots, k$, where we have omitted the dependence on $\left(x_{0}, \theta\right)$ in the subscripts for notational simplicity. Hence, by following the same argument as in Proposition 3.3.9, we can show there exists a constant $C$, such that for all $i=1, \ldots, m, \ell=1, \ldots, n$, $j=1, \ldots, k$ and $\theta \in \mathcal{K}$, we have $\left\|\left(\int_{0}^{T}\left|X_{\ell, t}^{i}\right|^{2} \mathrm{~d} t\right)^{\frac{1}{2}}\right\|_{\Psi_{2}} \leq C$ and $\left\|\left(\int_{0}^{T}\left|\psi^{\theta}\left(t, X_{t}^{i}\right)_{j}\right|^{2} \mathrm{~d} t\right)^{\frac{1}{2}}\right\|_{\Psi_{2}} \leq$ $C$. Then, we can obtain from Lemmas 3.3.5 and 3.3.7 a uniform bound for the $\|\cdot\|_{\Psi_{1}-}$ norms of all entires of $\int_{0}^{T} Z_{t}^{i}\left(Z_{t}^{i}\right)^{\top} \mathrm{d} t$ and $\int_{0}^{T} Z_{t}^{i}\left(\mathrm{~d} X_{t}^{i}\right)^{\top}$. Consequently, we can apply Lemma 3.6.4 (with $\alpha=1, N=m$ and $\varepsilon^{\prime}=m \varepsilon$ ) to each entry of $\left(\int_{0}^{T} Z_{t}^{i}\left(Z_{t}^{i}\right)^{\top} \mathrm{d} t-\bar{U}\right)_{i=1}^{m}$ and $\left(\int_{0}^{T} Z_{t}^{i}\left(\mathrm{~d} X_{t}^{i}\right)^{\top}-\bar{V}\right)_{i=1}^{m}$, and deduce the desired concentration inequality.

### 3.3.4.3 Step 3: Proof of general regret bounds

After demonstrating how to verify (3.3.8) based on the precise jump sizes in the state dynamics, it remains to establish the general regret result in Theorem 3.3.2 under the assumption that (3.3.8) holds for some $\beta \geq 1$.

We start by showing that under (H.4(1)) and (H.5), the expression (3.3.5) is well-defined if $\theta$ is a sufficiently accurate estimation of the true parameter $\theta^{\star}$.

Lemma 3.3.11. Suppose (H.4(1)) and (H.5) hold. Then there exist constants $\varepsilon_{0}>0$ and $\tau_{0}>0$, such that for all $\theta \in \mathcal{K}_{0}:=\left\{\theta \in \mathbb{R}^{n \times(n+k)}| | \theta-\theta^{\star} \mid \leq \varepsilon_{0}\right\}$, we have $\lambda_{\min }\left(\bar{U}^{x_{0}, \theta}\right) \geq \tau_{0}$, where $\bar{U}^{x_{0}, \theta}$ is defined as in (3.3.7) and $\lambda_{\min }(A)$ is the smallest eigenvalue of a symmetric matrix $A$.
Proof. Since $\bar{U}^{x_{0}, \theta^{\star}}$ is positive semidefinite, we shall prove $\lambda_{\min }\left(\bar{U}^{x_{0}, \theta^{\star}}\right)>0$ by assuming that $\lambda_{\min }\left(\bar{U}^{x_{0}, \theta^{\star}}\right)=0$. Then we see there exists a non-zero vector $u=\binom{u_{1}}{u_{2}} \in \mathbb{R}^{n+k}$ with $u_{1} \in \mathbb{R}^{n}$ and $u_{2} \in \mathbb{R}^{k}$, such that $u^{\top} \bar{U}^{x_{0}, \theta^{\star}} u=0$. By the definition of $\bar{U}^{x_{0}, \theta^{\star}}$ in (3.3.7), we can deduce that $\mathbb{E}\left[\int_{0}^{T}\left|u^{\top} Z_{t}^{x_{0}, \theta^{\star}}\right|^{2} \mathrm{~d} t\right]=0$, which along with the definition of $Z_{t}^{x_{0}, \theta^{\star}}$ in (3.3.5) implies for $\mathrm{d} \mathbb{P} \otimes \mathrm{d} t$ a.e. that $u_{1}^{\top} X_{t}^{x_{0}, \theta^{\star}, \alpha^{\star}}+u_{2}^{\top} \alpha_{t}^{x_{0, \star}}=0$. This contradicts to (H.5), which leads to the desired inequality that $\lambda_{\text {min }}\left(\bar{U}^{x_{0}, \theta^{\star}}\right)>0$.

We then show that the map $\mathbb{R}^{n \times(n+k)} \ni \theta \mapsto \bar{U}^{x_{0}, \theta} \in \mathbb{R}$ is continuous. Theorem 3.2.7 shows that the map $\mathbb{R}^{n \times(n+k)} \ni \theta \mapsto X^{x_{0}, \theta} \in \mathcal{H}^{2}\left(\mathbb{R}^{n}\right)$ is continuous. Moreover, Theorems 3.2.5 and 3.2.6 imply that there exists a constant $C \geq 0$, such that for all $\theta \in \mathbb{R}^{n \times(n+k)}$ satisfying $\left|\theta-\theta^{\star}\right| \leq 1, t \in[0, T]$ and $x, x^{\prime} \in \mathbb{R}^{n}$, we have that $\left|\psi^{\theta}(t, 0)\right| \leq C, \mid \psi^{\theta}(t, x)-$ $\psi^{\theta}\left(t, x^{\prime}\right)|\leq C| x-x^{\prime} \mid$ and $\left|\psi^{\theta}(t, x)-\psi^{\theta^{\star}}(t, x)\right| \leq C(1+|x|)\left|\theta-\theta^{\star}\right|$, from which we can deduce that

$$
\begin{aligned}
\left|\psi^{\theta}(t, x)-\psi^{\theta^{\star}}\left(t, x^{\prime}\right)\right| & \leq\left|\psi^{\theta}(t, x)-\psi^{\theta^{\star}}(t, x)\right|+\left|\psi^{\theta^{\star}}(t, x)-\psi^{\theta^{\star}}\left(t, x^{\prime}\right)\right| \\
& \leq C(1+|x|)\left|\theta-\theta^{\star}\right|+C\left|x-x^{\prime}\right|
\end{aligned}
$$

Hence, for all $\theta \in \mathbb{R}^{n \times(n+k)}$ with $\left|\theta-\theta^{\star}\right| \leq 1$,

$$
\left\|\psi^{\theta}\left(\cdot, X^{x_{0}, \theta}\right)-\psi^{\theta^{\star}}\left(\cdot, X^{x_{0}, \theta^{\star}}\right)\right\|_{\mathcal{H}^{2}} \leq C\left(1+\left\|X^{x_{0}, \theta}\right\|_{\mathcal{H}^{2}}\right)\left|\theta-\theta^{\star}\right|+C\left\|X^{x_{0}, \theta}-X^{x_{0}, \theta^{\star}}\right\|_{\mathcal{H}^{2}}
$$

which along with the continuity of the map $\mathbb{R}^{n \times(n+k)} \ni \theta \mapsto X^{x_{0}, \theta} \in \mathcal{H}^{2}\left(\mathbb{R}^{n}\right)$ implies that the map $\mathbb{R}^{n \times(n+k)} \ni \theta \mapsto \psi^{\theta}\left(\cdot, X^{x_{0}, \theta}\right) \in \mathcal{H}^{2}\left(\mathbb{R}^{k}\right)$ is continuous. Since the entires of $\bar{U}^{x_{0}, \theta}$ involve only the expectations of products of $X^{x_{0}, \theta}$ and $\psi^{\theta}\left(\cdot, X^{x_{0}, \theta}\right)$, the desired continuity of the map $\mathbb{R}^{n \times(n+k)} \ni \theta \mapsto \bar{U}^{x_{0}, \theta} \in \mathbb{R}$ follows.

Finally, by the continuity of the minimum eigenvalue function, clearly $\mathbb{R}^{n \times(n+k)} \ni \theta \mapsto$ $\lambda_{\text {min }}\left(\bar{U}^{x_{0}, \theta}\right) \in \mathbb{R}$ is continuous, which along with the fact that $\lambda_{\text {min }}\left(\bar{U}^{x_{0}, \theta^{*}}\right)>0$ leads to the desired result.

We then quantify the estimation error of the least-squares estimator (3.3.6) by assuming the concentration inequality (3.3.8) holds for the compact set $\mathcal{K}_{0}$ in Lemma 3.3.11.

Proposition 3.3.12. Suppose (H.4(1)) and (H.5) hold. Let $\mathcal{K}_{0}$ be the set in Lemma 3.3.11. Assume further that there exist constants $C_{1}, C_{2}>0$ and $\beta \geq 1$ such that (3.3.8) holds for all $\varepsilon \geq 0, m \in \mathbb{N}$ and $\theta \in \mathcal{K}_{0}$. Then there exist constants $\bar{C}_{1}, \bar{C}_{2} \geq 0$, such that for all $\theta \in \mathcal{K}_{0}$ and $\delta \in(0,1 / 2)$, if $m \geq \bar{C}_{1}(-\ln \delta)^{\beta}$, then we have with probability at least $1-2 \delta$ that

$$
\begin{equation*}
\left|\hat{\theta}-\theta^{\star}\right| \leq \bar{C}_{2}\left(\sqrt{\frac{-\ln \delta}{m}}+\frac{(-\ln \delta)^{\beta}}{m}+\frac{(-\ln \delta)^{2 \beta}}{m^{2}}\right) \tag{3.3.13}
\end{equation*}
$$

where $\hat{\theta}$ denotes the transpose of the left-hand side of (3.3.6) associated with $\theta$.
Proof. Throughout the proof, let $\delta \in(0,1 / 2)$ and $\theta \in \mathcal{K}_{0}$ be fixed and let $\|\cdot\|_{2}$ be the matrix norm induced by Euclidean norms. The invertibility of $\bar{U}^{x_{0}, \theta}$ (see Lemma 3.3.11) implies that (3.3.5) is well-defined, which along with (3.3.6) leads to

$$
\begin{align*}
\left\|\hat{\theta}-\theta^{\star}\right\|_{2}= & \left\|\left(U^{x_{0}, \theta, m}+\frac{1}{m} \mathbb{I}\right)^{-1} V^{x_{0}, \theta, m}-\left(\bar{U}^{x_{0}, \theta}\right)^{-1} \bar{V}^{x_{0}, \theta}\right\|_{2} \\
\leq & \left\|\left(U^{x_{0}, \theta, m}+\frac{1}{m} \mathbb{I}\right)^{-1}-\left(\bar{U}^{x_{0}, \theta}\right)^{-1}\right\|_{2}\left\|V^{x_{0}, \theta, m}\right\|_{2}  \tag{3.3.14}\\
& +\left\|\left(\bar{U}^{x_{0}, \theta}\right)^{-1}\right\|_{2}\left\|V^{x_{0}, \theta, m}-\bar{V}^{x_{0}, \theta}\right\|_{2} .
\end{align*}
$$

We now estimate each term in the right-hand side of (3.3.14). By Lemma 3.3.11, $\lambda_{\text {min }}\left(\bar{U}^{x_{0}, \theta}\right) \geq \tau_{0}$ for some $\tau_{0}>0$, which implies that $\left\|\left(\bar{U}^{x_{0}, \theta}\right)^{-1}\right\|_{2} \leq 1 / \tau_{0}$. Moreover, by setting the right-hand side of (3.3.8) to be $\delta$, we can deduce with probability at least $1-2 \delta$ that $\left|U^{x_{0}, \theta, m}-\bar{U}^{x_{0}, \theta}\right| \leq \delta_{m}$ and $\left|V^{x_{0}, \theta, m}-\bar{V}^{x_{0}, \theta}\right| \leq \delta_{m}$ with the constant $\delta_{m}$ given by

$$
\begin{equation*}
\delta_{m}:=\max \left\{\left(\frac{C_{2}^{2}}{C_{1} m} \ln \left(\frac{C_{2}}{\delta}\right)\right)^{\frac{1}{2}}, \frac{C_{2}}{m}\left(\frac{1}{C_{1}} \ln \left(\frac{C_{2}}{\delta}\right)\right)^{\beta}\right\} \tag{3.3.15}
\end{equation*}
$$

where we have assumed without loss of generality that $C_{2} \geq 1$.
Let $m$ be a sufficiently large constant satisfying $\delta_{m}+1 / m \leq \tau_{0} / 2$. The fact that $\|\cdot\|_{2} \leq|\cdot|$ indicates with probability at least $1-2 \delta$ that $\left\|U^{x_{0}, \theta, m}+\frac{1}{m} \mathbb{I}-\bar{U}^{x_{0}, \theta}\right\|_{2} \leq \frac{1}{m}+\delta_{m} \leq \frac{\tau_{0}}{2}$, which in turn yields

$$
\lambda_{\min }\left(U^{x_{0}, \theta, m}+\frac{1}{m} \mathbb{I}\right) \geq \lambda_{\min }\left(\bar{U}^{x_{0}, \theta}\right)-\left\|U^{x_{0}, \theta, m}+\frac{1}{m} \mathbb{I}-\bar{U}^{x_{0}, \theta}\right\|_{2} \geq \frac{\tau_{0}}{2},
$$

or equivalently $\left\|\left(U^{x_{0}, \theta, m}+\frac{1}{m} \mathbb{I}\right)^{-1}\right\|_{2} \leq 2 / \tau_{0}$. Then, since $A^{-1}-(A+B)^{-1}=(A+B)^{-1} B A^{-1}$ for all nonsingular matrices $A$ and $A+B$, we have with probability at least $1-2 \delta$ that,

$$
\begin{aligned}
\left\|\left(U^{x_{0}, \theta, m}+\frac{1}{m} \mathbb{I}\right)^{-1}-\left(\bar{U}^{x_{0}, \theta}\right)^{-1}\right\|_{2} & =\left\|\left(\bar{U}^{x_{0}, \theta}+U^{x_{0}, \theta, m}+\frac{1}{m} \mathbb{I}-\bar{U}^{x_{0}, \theta}\right)^{-1}-\left(\bar{U}^{x_{0}, \theta}\right)^{-1}\right\|_{2} \\
& \leq\left\|\left(U^{x_{0}, \theta, m}+\frac{1}{m} \mathbb{I}\right)^{-1}\right\|_{2}\left\|\left(\bar{U}^{x_{0}, \theta}\right)^{-1}\right\|_{2}\left\|\left(U^{x_{0}, \theta, m}+\frac{1}{m} \mathbb{I}\right)-\bar{U}^{x_{0}, \theta}\right\|_{2} \\
& \leq \frac{2}{\tau_{0}^{2}}\left(\frac{1}{m}+\delta_{m}\right),
\end{aligned}
$$

which along with the inequality that $\left\|V^{x_{0}, \theta, m}\right\|_{2} \leq\left\|\bar{V}^{x_{0}, \theta}\right\|_{2}+\left|V^{x_{0}, \theta, m}-\bar{V}^{x_{0}, \theta}\right|$ allows us to derive the following estimate from (3.3.14):

$$
\left\|\hat{\theta}-\theta^{\star}\right\|_{2} \leq \frac{2}{\tau_{0}^{2}}\left(\frac{1}{m}+\delta_{m}\right)\left(\left\|\bar{V}^{x_{0}, \theta}\right\|_{2}+\delta_{m}\right)+\frac{\delta_{m}}{\tau_{0}} .
$$

Note that $\left\|\bar{V}^{x_{0}, \theta}\right\|_{2}$ is uniformly bounded for all $\theta \in \mathcal{K}_{0}$ by the compactness of $\mathcal{K}_{0}$ and the continuity of the map $\theta \mapsto \bar{V}^{x_{0}, \theta}$ (cf. Lemma 3.3.11). Thus, by the condition that $\beta \geq 1$ and the definition of $\delta_{m}$ in (3.3.15), we see that there exists a constant $\bar{C}_{2}$, depending only on $C_{1}, C_{2}, \beta, \tau_{0}$, and the constants in (H.4(1)), such that the desired estimate (3.3.13) holds with probability at least $1-2 \delta$, provided that $m$ satisfies $\delta_{m}+1 / m \leq \tau_{0} / 2$. Since $\beta \geq 1$ and $\delta \leq 1 / 2$, we see there exists $\bar{C}_{1} \geq 0$, independent of $m, \delta$ and $\theta$, such that the inequality (3.3.13) holds for all $m$ satisfying $m \geq \bar{C}_{1}(-\ln \delta)^{\beta}$.

Now we are ready to present the proof of Theorem 3.3.2.
Proof of Theorem 3.3.2. We start by proving Item (1). Then, by the assumptions that $\lambda_{\text {min }}\left(\bar{U}^{x_{0}, \theta_{0}}\right)>0$ and (3.3.8) holds for $\mathcal{K}=\overline{\mathcal{K}}_{0}:=\mathcal{K}_{0} \cup\left\{\theta_{0}\right\}$ with $\mathcal{K}_{0}$ from Lemma 3.3.11, we can extend Proposition 3.3.12 to show that (3.3.13) holds for all $\theta \in \overline{\mathcal{K}}_{0}, \delta \in(0,1 / 2)$ and $m \geq \bar{C}_{1}(-\ln \delta)^{\beta}$, with some constants $\bar{C}_{1}, \bar{C}_{2} \geq 1$ depending on $\overline{\mathcal{K}}_{0}$. In the subsequent analysis, we fix $\delta \in(0,1 / 4)$ and for all $\ell \in \mathbb{N} \cup\{0\}$, we define $\delta_{\ell}=2^{-\ell} \delta$, and let $\theta_{\ell+1}$ be generated by using (3.3.6) with $m=m_{\ell}$ and $\theta=\theta_{\ell}$. We shall specify the precise choice of $m_{0}$ later.

In the sequel, we assume without loss of generality that $\varepsilon_{0} /\left(3 \bar{C}_{2}\right) \leq 1$ and $\bar{C}_{2} / \varepsilon_{0} \geq \bar{C}_{1}$, where $\varepsilon_{0}>0$ is the constant in the definition of $\mathcal{K}_{0}$ (see Lemma 3.3.11). We first show that there exists $\hat{C}_{0}>0$, independent of $\delta$, such that if $m_{0} \geq \hat{C}_{0}(-\ln \delta)^{\beta}$, then for all $\ell \in \mathbb{N} \cup\{0\}$,

$$
\begin{equation*}
\bar{C}_{2}\left(\sqrt{\frac{-\ln \delta_{\ell}}{m_{\ell}}}+\frac{\left(-\ln \delta_{\ell}\right)^{\beta}}{m_{\ell}}+\frac{\left(-\ln \delta_{\ell}\right)^{2 \beta}}{m_{\ell}^{2}}\right) \leq \varepsilon_{0} . \tag{3.3.16}
\end{equation*}
$$

By the assumption that $\varepsilon_{0} /\left(3 \bar{C}_{2}\right) \leq 1$, it suffices to show that for all $\ell \in \mathbb{N} \cup\{0\},-\ln \delta_{\ell} / m_{\ell} \leq$ $\left(\varepsilon_{0} /\left(3 \bar{C}_{2}\right)\right)^{2}$ and $\left(-\ln \delta_{\ell}\right)^{\beta} / m_{\ell} \leq \varepsilon_{0} /\left(3 \bar{C}_{2}\right)$. Given $\beta \geq 1$ and $\delta_{\ell}<1 / 4$, it suffices to ensure $m_{\ell} \geq C\left(-\ln \delta_{\ell}\right)^{\beta}$ for all $\ell \in \mathbb{N} \cup\{0\}$, where $C$ is a sufficiently large constant independent of $\delta$ and $\ell$. By the definitions of $\left(\delta_{\ell}\right)_{\ell \in \mathbb{N}}$ and $\left(m_{\ell}\right)_{\ell \in \mathbb{N}}$ and the fact that $\delta<1 / 4$, the desired condition can be achieved by choosing $m_{0} \geq \hat{C}_{0}(-\ln \delta)^{\beta}$, for a sufficiently large constant $\hat{C}_{0}$ satisfying
$\sup _{\ell \in \mathbb{N} \cup\{0\}, \delta \in\left(0, \frac{1}{4}\right)} \frac{\left(-\ln \left(2^{-\ell} \delta\right)\right)^{\beta}}{2^{\ell}(-\ln \delta)^{\beta}}=\sup _{\ell \in \mathbb{N} \cup\{0\}, \delta \in\left(0, \frac{1}{4}\right)} 2^{-\ell}\left(\frac{\ell \ln 2}{-\ln \delta}+1\right)^{\beta} \leq \sup _{\ell \in \mathbb{N} \cup\{0\}} 2^{-\ell}\left(\frac{\ell}{2}+1\right)^{\beta} \leq \hat{C}_{0}<\infty$.
Now we choose $m_{0} \geq \max \left(\hat{C}_{0}, \bar{C}_{1}\right)(-\ln \delta)^{\beta}$, and show by induction that for all $k \in \mathbb{N} \cup\{0\}$, it holds with probability at least $1-2 \sum_{\ell=0}^{k-1} \delta_{\ell}$ that $\theta_{\ell} \in \overline{\mathcal{K}}_{0}$ for all $\ell=0, \ldots, k$ and

$$
\left|\theta_{k}-\theta^{\star}\right|_{2} \leq \begin{cases}\left|\theta_{0}-\theta^{\star}\right|_{2}  \tag{3.3.17}\\ \bar{C}_{2}\left(\sqrt{\frac{-\ln \delta_{k}}{m_{k}}}+\frac{\left(-\ln \delta_{k}\right)^{\beta}}{m_{k}}+\frac{\left(-\ln \delta_{k}\right)^{2 \beta}}{m_{k}^{2}}\right), & k \in \mathbb{N}\end{cases}
$$

The statement clearly holds for $k=0$. Now suppose that the induction statement holds for some $k \in \mathbb{N} \cup\{0\}$. Conditioning on $\theta_{k} \in \overline{\mathcal{K}}_{0}$, we can apply (3.3.13) with $\theta=\theta_{k}, \delta=\delta_{k}<1 / 2$ and $m=m_{k} \geq \bar{C}_{1}\left(-\ln \delta_{k}\right)^{\beta}$ (see (3.3.16) and $\bar{C}_{2} / \varepsilon_{0} \geq \bar{C}_{1}$ ), and deduce with probability at least $1-2 \delta_{k}$ that (3.3.17) holds for the index $k+1$, which along with (3.3.16) shows that $\theta_{k+1} \in \mathcal{K}_{0} \subset \overline{\mathcal{K}}_{0}$. Since the induction hypothesis implies that $\theta_{k} \in \overline{\mathcal{K}}_{0}$ holds with probability at least $1-2 \sum_{\ell=0}^{k-1} \delta_{\ell}$, one can deduce that the induction statement also holds $k+1$.

The above induction argument shows that if $m_{0}=C(-\ln \delta)^{\beta}$ for any constant $C \geq C_{0}:=$ $\max \left(\hat{C}_{0}, \bar{C}_{1}\right)$, then with probability at least $1-2 \sum_{\ell=0}^{\infty} \delta_{\ell}=1-4 \delta, \theta_{k} \in \overline{\mathcal{K}}_{0}$ and (3.3.17) holds for all $k \in \mathbb{N} \cup\{0\}$. Now let us assume such a setting, and observe that the $i$-th trajectory is generated with control $\psi^{\theta_{\ell}}$ if $i \in\left(\sum_{j=0}^{\ell-1} m_{j}, \sum_{j=0}^{\ell} m_{j}\right]=\left(m_{0}\left(2^{\ell}-1\right), m_{0}\left(2^{\ell+1}-1\right)\right]$ for $\ell \in \mathbb{N} \cup\{0\}$ (cf. Algorithm 10). Then we can apply Theorem 3.2.7 and deduce for all $N \in \mathbb{N}$ that

$$
\begin{align*}
R(N) & \leq \sum_{\ell=0}^{\left\lceil\log _{2}\left(\frac{N}{m_{0}}+1\right)\right\rceil-1} m_{\ell}\left(J^{\theta^{\star}}\left(\psi^{\theta_{\ell}} ; x_{0}\right)-V\left(x_{0} ; \theta^{\star}\right)\right) \leq C^{\prime} \sum_{\ell=0}^{\left\lceil\log _{2}\left(\frac{N}{m_{0}}+1\right)\right\rceil-1} m_{\ell}\left|\theta_{\ell}-\theta^{\star}\right| \\
& \leq C^{\prime} m_{0}+C^{\prime} \sum_{\ell=1}^{\left\lceil\log _{2}\left(\frac{N}{m_{0}}+1\right)\right\rceil-1}\left(\sqrt{\left(-\ln \delta_{\ell}\right) m_{\ell}}+\left(-\ln \delta_{\ell}\right)^{\beta}\left(1+\frac{\left(-\ln \delta_{\ell}\right)^{\beta}}{m_{\ell}}\right)\right) \\
& \leq C^{\prime}(-\ln \delta)^{\beta}+C^{\prime} \sum_{\ell=1}^{\left\lceil\log _{2}\left(\frac{N}{m_{0}}+1\right)\right\rceil-1}\left(\sqrt{\left(-\ln \delta_{\ell}\right) m_{\ell}}+\left(-\ln \delta_{\ell}\right)^{\beta}\right), \tag{3.3.18}
\end{align*}
$$

where we have denoted by $C^{\prime}$ a generic constant independent of $\ell, N, \delta$, and used the fact that $\left(-\ln \delta_{\ell}\right)^{\beta} / m_{\ell} \leq C^{\prime}$ for the last inequality (cf. the choice of $\hat{C}_{0}$ ). We then derive an upper bound of (3.3.18). By virtue of the inequality that $\sqrt{\left(-\ln \delta_{\ell}\right) m_{\ell}}=\sqrt{(\ell \ln 2-\ln \delta) 2^{\ell} m_{0}} \leq$ $C^{\prime} \sqrt{(\ell-\ln \delta) m_{0}} \sqrt{2}$ for all $\ell \in \mathbb{N}$, we have

$$
\begin{aligned}
\sum_{\ell=1}^{\left\lceil\log _{2}\left(\frac{N}{m_{0}}+1\right)\right\rceil-1} \sqrt{\left(-\ln \delta_{\ell}\right) m_{\ell}} & \leq C^{\prime} \sqrt{(\ln N-\ln \delta) m_{0}} \sqrt{2} \log _{2}\left(\frac{N}{m_{0}}+1\right) \\
& \leq C^{\prime} \sqrt{(\ln N-\ln \delta)\left(N+(-\ln \delta)^{\beta}\right)}
\end{aligned}
$$

Moreover, by $\ln \delta_{\ell}=-\ell \ln 2+\ln \delta$ and Hölder's inequality,

$$
\sum_{\ell=1}^{\left\lceil\log _{2}\left(\frac{N}{m_{0}}+1\right)\right\rceil-1}\left(-\ln \delta_{\ell}\right)^{\beta} \leq \sum_{\ell=1}^{C^{\prime} \ln N} C^{\prime}\left((\ell \ln 2)^{\beta}+(-\ln \delta)^{\beta}\right) \leq C^{\prime}\left((\ln N)^{\beta+1}+\ln N(-\ln \delta)^{\beta}\right)
$$

Consequently, from (3.3.18), $\beta \geq 1$ and the inequality $\sqrt{x+y} \leq \sqrt{x}+\sqrt{y}$ for all $x, y \geq 0$, it is clear for all $N \in \mathbb{N}, R(N) \leq C^{\prime}\left(\sqrt{N} \sqrt{\ln N}+\sqrt{-\ln \delta} \sqrt{N}+(-\ln \delta)^{\beta} \ln N\right)$ for some constant $C^{\prime}$ independent of $\beta$ and $N$, which finishes the proof of Item (1).

We are ready to show Item (2). For each $N \in \mathbb{N} \cap[3, \infty)$, we define $\delta_{N}=1 / N^{2}$ and the event $A_{N}=\left\{R(N)>C^{\prime}\left(\sqrt{N} \sqrt{\ln N}+\sqrt{-\ln \delta_{N}} \sqrt{N}+\left(-\ln \delta_{N}\right)^{\beta} \ln N\right)\right\}$. Item (1) shows that $\sum_{N=3}^{\infty} \mathbb{P}\left(A_{N}\right) \leq 4 \sum_{N=3}^{\infty} \delta_{N}<\infty$. Hence, from the Borel-Cantelli lemma, $\mathbb{P}\left(\lim \sup _{N \rightarrow \infty} A_{N}\right)=0$, which along with the definition of $\delta_{N}$ implies the desired conclusion.

### 3.4 Extension: RL problems with controlled diffusion

In this section, we extend our framework to analyze the regret order of learning algorithms for general continuous-time RL problems, whose state dynamics involves controlled diffusion. To simplify the presentation, we focus on entropy-regularized problems studied in [159, 143, 135,152 ] and outline the essential steps of the argument.

For each $\theta=(A, B) \in \mathbb{R}^{n \times(n+k)}$, define $V(\cdot ; \theta):[0, T] \times \mathbb{R}^{n} \rightarrow \mathbb{R}$ by

$$
\begin{equation*}
V(t, x ; \theta):=\inf _{\alpha \in \mathcal{H}^{2}\left(\mathbb{R}^{k}\right)} \mathbb{E}\left[\int_{t}^{T} f\left(s, X_{s}^{t, x, \alpha}, \alpha_{s}\right) \mathrm{d} s+g\left(X_{T}^{t, x, \alpha}\right)\right], \quad \forall(t, x) \in[0, T] \times \mathbb{R}^{n} \tag{3.4.1}
\end{equation*}
$$

where for each $\alpha \in \mathcal{H}^{2}\left(\mathbb{R}^{k}\right), X^{t, x, \alpha} \in \mathcal{S}^{2}\left(\mathbb{R}^{n}\right)$ satisfies the controlled dynamics:

$$
\begin{equation*}
\mathrm{d} X_{s}=\left(A X_{s}+B \alpha_{s}\right) \mathrm{d} s+\sigma\left(s, X_{s}, \alpha_{s}\right) \mathrm{d} W_{s}, \quad s \in[t, T], \quad X_{t}=x \tag{3.4.2}
\end{equation*}
$$

The functions $f:[0, T] \times \mathbb{R}^{n} \times \mathbb{R}^{k} \rightarrow \mathbb{R} \cup\{\infty\}$ and $\sigma:[0, T] \times \mathbb{R}^{n} \times \mathbb{R}^{k} \rightarrow \mathbb{R}^{n \times d}$ are such that for all $(t, x) \in[0, T] \times \times \mathbb{R}^{n}$ and $a=\left(a_{i}\right)_{i=1}^{k} \in \mathbb{R}^{k}$,

$$
\begin{equation*}
f(t, x, a)=\sum_{i=1}^{k} \bar{f}_{i}(t, x) a_{i}+\mathcal{R}_{\mathrm{en}}(a), \quad \sigma(t, x, a) \sigma(t, x, a)^{\mathrm{\top}}=\sum_{i=1}^{k} \bar{\sigma}_{i}(t, x) \bar{\sigma}_{i}(t, x)^{\mathrm{\top}} a_{i} \tag{3.4.3}
\end{equation*}
$$

where for each $i=1, \ldots k, \bar{f}_{i}:[0, T] \times \mathbb{R}^{n} \rightarrow \mathbb{R}, \bar{\sigma}_{i}:[0, T] \times \mathbb{R}^{n} \rightarrow \mathbb{R}^{n \times d}$ are some given functions and $\mathcal{R}_{\text {en }}: \mathbb{R}^{k} \rightarrow \mathbb{R} \cup\{\infty\}$ is Shannon's entropy function (cf. Example 3.2.2) such that

$$
\mathcal{R}_{\mathrm{en}}(a)= \begin{cases}\sum_{i=1}^{k} a_{i} \ln \left(a_{i}\right), & a \in \Delta_{k}:=\left\{a \in[0,1]^{k} \mid \sum_{i=1}^{k} a_{i}=1\right\}  \tag{3.4.4}\\ \infty, & a \in \mathbb{R}^{k} \backslash \Delta_{k}\end{cases}
$$

To avoid needless technicalities, we assume $\left(\bar{f}_{i}\right)_{i=1}^{k},\left(\bar{\sigma}_{i}\right)_{i=1}^{k}$ and $g$ to be bounded and sufficiently regular as in Proposition 3.4.1.

Note that (3.4.4) restricts control processes to those taking values in $\Delta_{k}$. Hence, if $\bar{\sigma}_{\ell}(t, x) \equiv \bar{\sigma}$ for some $\bar{\sigma} \in \mathbb{R}^{n \times d}$, then (3.4.1)-(3.4.2) is a special case of the linear-convex model studied in Sections 3.2-3.3. Consequently, Theorem 3.3.2 can be applied to study the regret order of GLS algorithms for (3.4.1)-(3.4.2) with given initial time and state $(t, x) \in$ $[0, T] \times \mathbb{R}^{n}$ but with unknown parameter $\theta$.

To analyze the regret order of learning algorithms with general $\sigma$, a crucial step is to extend Theorem 3.2.7 and quantify the performance of a greedy policy from an incorrect model.

The fact that control affects the diffusion coefficients complicates the stability analysis of optimal feedback controls (i.e., Theorem 3.2.6) for (3.4.1)-(3.4.2). The following proposition proves a linear performance gap under the condition that the value function $V(t, x ; \theta)$ in (3.4.1) is sufficiently regular in $t, x$ and $\theta$. Recall that the first and second-order derivatives of a sufficiently regular value function can be represented by solutions to the associated FBSDE (3.2.11) (see e.g., [169, Theorem 4.1, p. 250]). Hence the linear performance gap can also be established by assuming sufficient regularity of the solution process ( $Y^{t, x}, Z^{t, x}$ ) to (3.2.11), whose details are omitted here.

Proposition 3.4.1. For each $\theta \in \mathbb{R}^{n \times(n+k)}$, let $V(\cdot ; \theta):[0, T] \times \mathbb{R}^{n} \rightarrow \mathbb{R}$ be defined by (3.4.1). Suppose that $\left(\bar{f}_{i}\right)_{i=1}^{k},\left(\bar{\sigma}_{i}\right)_{i=1}^{k} \subset C^{0,1}\left([0, T] \times \mathbb{R}^{n}\right), g \in C^{1}\left(\mathbb{R}^{n}\right)$, and there exists $\mathfrak{m}:[0, \infty) \rightarrow[0, \infty)$ such that for all $(t, x) \in[0, T] \times \mathbb{R}^{n}$ and $\theta, \theta^{\prime} \in \mathbb{R}^{n \times(n+k)}, \frac{\partial}{\partial t} V(\cdot ; \theta)$ is continuous, $\|V(\cdot ; \theta)\|_{C^{0,3}\left([0, T] \times \mathbb{R}^{n}\right)} \leq \mathfrak{m}(|\theta|)$,

$$
\begin{aligned}
\left|\nabla_{x} V(t, x ; \theta)-\nabla_{x} V\left(t, x ; \theta^{\prime}\right)\right| & +\left|\operatorname{Hess}_{x} V(t, x ; \theta)-\operatorname{Hess}_{x} V\left(t, x ; \theta^{\prime}\right)\right| \\
& \leq\left(\mathfrak{m}(|\theta|)+\mathfrak{m}\left(|\theta|^{\prime}\right)\right)\left|\theta-\theta^{\prime}\right|(1+|x|)
\end{aligned}
$$

Then for all $\theta \in \mathbb{R}^{n \times(n+k)}$, there exists $\psi^{\theta} \in \mathcal{V}$ such that
(1) $\psi^{\theta}$ is an optimal feedback control of (3.4.1)-(3.4.2) satisfying for all $x_{0} \in \mathbb{R}^{n}$ and $\theta \in \mathbb{R}^{n \times(n+k)}, V\left(0, x_{0} ; \theta\right)=J\left(\psi^{\theta} ; x_{0}, \theta\right)$, where for each $\psi \in \mathcal{V}$,

$$
J\left(\psi ; x_{0}, \theta\right):=\mathbb{E}\left[\int_{0}^{T} f\left(t, X_{t}^{x_{0}, \theta, \psi}, \psi\left(t, X_{t}^{x_{0}, \theta, \psi}\right)\right) \mathrm{d} t+g\left(X_{T}^{x_{0}, \theta, \psi}\right)\right]
$$

and $X^{x_{0}, \theta, \psi} \in \mathcal{S}^{2}\left(\mathbb{R}^{n}\right)$ satisfies the following dynamics:

$$
\mathrm{d} X_{t}=\theta\binom{X_{t}}{\psi\left(t, X_{t}\right)} \mathrm{d} t+\sigma\left(t, X_{t}, \psi\left(t, X_{t}\right)\right) \mathrm{d} W_{t}, \quad t \in[0, T], \quad X_{0}=x_{0},
$$

(2) for all $x_{0} \in \mathbb{R}^{n}$ and $R \geq 0$ there exists a constant $C$ such that for all $\theta, \theta^{\prime} \in \mathbb{R}^{n \times(n+k)}$ with $|\theta|,\left|\theta^{\prime}\right| \leq R$,

$$
\left|J\left(\psi^{\theta^{\prime}} ; x_{0}, \theta\right)-J\left(\psi^{\theta} ; x_{0}, \theta\right)\right| \leq C\left|\theta^{\prime}-\theta\right|
$$

Proposition 3.4.1 relies on the regularity and Lipschitz stability of the value function $V$. For instance, if all coefficients are bounded and sufficiently smooth, and $\sigma$ satisfies the uniform parabolicity condition, then $C^{2+\alpha}$ regularity results for fully nonlinear parabolic PDEs (see e.g., the Evan-Kryolv theorem in [97, Theorems 6.4.3 and 6.4.4, p. 301]) and a bootstrap argument would ensure that for any given $\theta$, the function $V(\cdot, \theta)$ is continuously differentiable in $t$ and three-time continuously differentiable in $x$. Due to the unbounded drift coefficient of (3.4.2), the boundedness in the $C^{0,3}\left([0, T] \times \mathbb{R}^{n}\right)$-norm and the locally Lipschitz continuity of $V$ in $\theta$ follow from an extension of the Schauder estimate (see e.g., [96]) to nonlinear parabolic equations with unbounded coefficients in the whole space.

With Proposition 3.4.1, we can then quantify the regrets of GLS algorithms (see Algorithm 10) for (3.4.1)-(3.4.2) with unknown drift parameter $\theta$ and known diffusion coefficient $\sigma$. By the boundedness of $\sigma$ and the regularity of $\psi^{\theta}$, one can prove Proposition 3.3.10 in the present setting. Hence, Theorem 3.3.2 (with $\beta=1$ in (3.3.8)) shows that Algorithm 10 enjoys a sublinear regret as shown in Theorem 3.3.4.

Proof of Proposition 3.4.1. For any given $\theta=(A, B) \in \mathbb{R}^{n \times(n+k)}$, the regularity of $V(\cdot ; \theta)$ and [169, Proposition 3.5, p. 182] imply that $V(\cdot ; \theta)$ is the unique classical solution to the associated HJB equation. That is, for all $(t, x) \in[0, T) \times \mathbb{R}^{d}$,
$\frac{\partial}{\partial t} V(t, x)+\inf _{a \in \Delta_{k}}\left(\frac{1}{2} \operatorname{tr}\left(\sigma(t, x, a) \sigma(t, x, a)^{\top} \operatorname{Hess}_{x} V(t, x)\right)+\left\langle A x+B a, \nabla_{x} V(t, x)\right\rangle+f(t, x, a)\right)=0$, and $V(T, x)=g(x)$ for all $x \in \mathbb{R}^{d}$. By (3.4.3), for all $(t, x) \in[0, T] \times \mathbb{R}^{n}$,

$$
\begin{aligned}
\psi^{\theta}(t, x) & :=\arg \min _{a \in \Delta_{k}}\left(\frac{1}{2} \operatorname{tr}\left(\sigma(t, x, a) \sigma(t, x, a)^{\top} \operatorname{Hess}_{x} V(t, x ; \theta)\right)+\left\langle B a, \nabla_{x} V(t, x ; \theta)\right\rangle+f(t, x, a)\right) \\
& =\nabla \mathcal{R}_{\mathrm{en}}^{*}\left(-\frac{1}{2} \operatorname{tr}\left(\bar{\sigma}(t, x) \bar{\sigma}(t, x)^{\top} \operatorname{Hess}_{x} V(t, x ; \theta)\right)-B^{\top} \nabla_{x} V(t, x ; \theta)-\bar{f}(t, x)\right),
\end{aligned}
$$

where for all $z \in \mathbb{R}^{k}, \mathcal{R}_{\text {en }}^{*}(z)=\sup _{a \in \Delta_{k}}\left(\langle a, z\rangle-\mathcal{R}_{\text {en }}(a)\right)=\ln \sum_{i=1}^{k} \exp \left(z_{i}\right)$, and

$$
\begin{aligned}
\operatorname{tr}\left(\bar{\sigma}(t, x) \bar{\sigma}(t, x)^{\top} \operatorname{Hess}_{x} V(t, x ; \theta)\right) & =\left(\begin{array}{c}
\operatorname{tr}\left(\bar{\sigma}_{1}(t, x) \bar{\sigma}_{1}(t, x)^{\top} \operatorname{Hess}_{x} V(t, x ; \theta)\right) \\
\vdots \\
\operatorname{tr}\left(\bar{\sigma}_{k}(t, x) \bar{\sigma}_{k}(t, x)^{\top} \operatorname{Hess}_{x} V(t, x ; \theta)\right)
\end{array}\right), \\
\bar{f}(t, x) & =\left(\begin{array}{c}
\bar{f}_{1}(t, x) \\
\vdots \\
\bar{f}_{k}(t, x)
\end{array}\right) .
\end{aligned}
$$

The Lipschitz continuity of $\nabla \mathcal{R}_{\text {en }}^{*}$ and the regularity assumptions imply that $\psi^{\theta} \in \mathcal{V}$ and the corresponding state process $X^{x_{0}, \theta, \psi^{\theta}}$ is well defined. Then a standard verification argument (see e.g., [169, Theorem 6.6, p. 278]) shows $\psi^{\theta}$ is an optimal feedback control and finishes the proof of Item (1).

To prove Item (2), Fix $x_{0} \in \mathbb{R}^{n}$ and $R \geq 0$ and let $C$ be a generic constant independent of $\theta$. Note that the Fenchel-Young identity gives that $\mathcal{R}_{\text {en }}^{*}\left(\nabla \mathcal{R}_{\text {en }}^{*}(z)\right)=\left\langle z, \nabla \mathcal{R}_{\text {en }}^{*}(z)\right\rangle-\mathcal{R}_{\text {en }}^{*}(z)$ for all $z \in \mathbb{R}^{k}$, which along with (3.4.3) implies that for all $(t, x, \theta) \in[0, T] \times \mathbb{R}^{n} \times \mathbb{R}^{n \times(n+k)}$,

$$
\begin{aligned}
f\left(t, x, \psi^{\theta}(t, x)\right)= & -\left\langle\frac{1}{2} \operatorname{tr}\left(\bar{\sigma}(t, x) \bar{\sigma}(t, x)^{\top} \operatorname{Hess}_{x} V(t, x ; \theta)\right)+B^{\top} \nabla_{x} V(t, x ; \theta), \psi^{\theta}(t, x)\right\rangle \\
& -\mathcal{R}_{\text {en }}^{*}\left(-\frac{1}{2} \operatorname{tr}\left(\bar{\sigma}(t, x) \bar{\sigma}(t, x)^{\top} \operatorname{Hess}_{x} V(t, x ; \theta)\right)-B^{\top} \nabla_{x} V(t, x ; \theta)-\bar{f}(t, x)\right) .
\end{aligned}
$$

By the regularity assumptions of the coefficients and the function $V$, for all $t \in[0, T]$, $x, x^{\prime} \in \mathbb{R}^{n}$ and $\theta, \theta^{\prime} \in \mathbb{R}^{n \times(n+k)}$ with $|\theta|,\left|\theta^{\prime}\right| \leq R$, there exists $C \geq 0$ such that

$$
\begin{aligned}
\left|\psi^{\theta^{\prime}}\left(t, x^{\prime}\right)-\psi^{\theta}(t, x)\right|+\mid f\left(t, x^{\prime}, \psi^{\theta^{\prime}}\left(t, x^{\prime}\right)\right) & -f\left(t, x, \psi^{\theta}(t, x)\right) \mid \\
& \leq C\left(\left|x-x^{\prime}\right|+\left(1+\left|x^{\prime}\right|+|x|\right)\left|\theta-\theta^{\prime}\right|\right)
\end{aligned}
$$

Proceeding along the lines of the proof of Theorem 3.2.7 leads to the desired estimate in Item (2).

### 3.5 Numerical experiments

In this section, we test the theoretical findings and Algorithm 10 through numerical experiment on a three-dimensional LQ RL problem considered in [45, 46]. Our experiments show the convergence of the least-squares estimations to the true parameters as the number of episodes increases, as well as the sublinear cumulative regret as indicated in Theorem 3.3.4. Moreover, it confirms that the state coefficient $A^{\star}$ is easier to learn than the control coefficient $B^{\star}$, consistent with the observations in [46]. Our numerical result shows that a rough estimation of the control parameter $B^{\star}$ is often sufficient to design a nearly optimal feedback control, and that the Algorithm 10 is robust with respect to the initial batch size $m_{0}$.

Problem setup. We consider a three-dimensional LQ RL problems over the time horizon $[0, T]$ with $T=1.5$, where the linear state dynamics (3.3.2) has the initial state $x_{0}$ and unknown coefficients $\theta^{\star}=\left(A^{\star}, B^{\star}\right) \in \mathbb{R}^{3 \times(3+3)}$ chosen as in [45, 46]:

$$
A^{\star}=\left[\begin{array}{ccc}
1.01 & 0.01 & 0 \\
0.01 & 1.01 & 0.01 \\
0 & 0.01 & 1.01
\end{array}\right], B^{\star}=\mathbb{I}_{3}, \sigma=\mathbb{I}_{3}, \gamma \equiv 0, x_{0}=0
$$

with $\mathbb{I}_{3}$ being the $3 \times 3$ identity matrix, and the cost functional (3.3.1) involves quadratic functions $g \equiv 0$ and $f(t, x, a)=\left(x^{\top} Q x+a^{\top} R a\right) / 2$, with $Q=0.1 \mathbb{I}_{3}$ and $R=\mathbb{I}_{3}$. As mentioned in [45, 46], this state dynamics corresponds to a marginally unstable graph Laplacian system where adjacent nodes are weakly connected, which arises naturally from consensus and distributed averaging problems. Since the cost penalizes the control inputs more than the states, it is essential to learn the unstable components of $A^{\star}$ and perform control on these components in order to achieve an optimal cost. Note that this LQ RL problem satisfies (H.4); see the last paragraph of Remark 3.3.1.

The numerical experiments are coded using Python. Algorithm 10 is initialized with $m_{0}=4$ and the initial guess $A_{0}=\left[\begin{array}{ccc}1.6243 & -0.6118 & -0.5282 \\ -1.0730 & 0.8654 & -2.3015 \\ 1.7418 & -0.7612 & 0.3190\end{array}\right]$ and $B_{0}=\left[\begin{array}{ccc}-0.2444 & 1.4621 & -2.0601 \\ -0.324 & -0.3841 \\ -1.0999 & -0.1382 & 1.1338 \\ -0.12979\end{array}\right]$, whose entries are sampled independently from the standard normal distribution. For each $\ell \in \mathbb{N} \cup\{0\}$, given the current estimate $\theta_{\ell}=\left(A_{\ell}, B_{\ell}\right)$ of $\theta^{\star}$, classical LQ control theory (see e.g., [169]) shows that solutions to (3.2.11) can be found analytically via Riccati equations, and the greedy policy $\psi^{\theta_{\ell}}$ is given by $\psi^{\theta_{\ell}}(t, x)=-R^{-1} B^{\top} P_{t}^{\theta_{\ell}} x$, where $P^{\theta_{\ell}}$ is the unique positive semidefinite solution to

$$
\begin{equation*}
\frac{\mathrm{d}}{\mathrm{~d} t} P_{t}+A_{\ell}^{\top} P_{t}+P_{t} A_{\ell}-P_{t}\left(B_{\ell} R^{-1} B_{\ell}^{\top}\right) P_{t}+Q=0, t \in(0, T) ; \quad P_{T}=0 \tag{3.5.1}
\end{equation*}
$$

We solve (3.5.1) numerically via a high-order Runge-Kutta method on a uniform time grid with stepsize $T / 100$, and then simulate $m_{\ell}=2^{\ell} m_{0}$ independent trajectories of the state dynamics (3.3.4) (controlled by $\psi^{\theta_{\ell}}$ ) using the Euler-Maruyama method on the same time grid. To estimate statistical properties of the algorithm regret (3.3.3), we execute Algorithm 10 for 100 independent runs, where among different executions, the observed state trajectories are simulated based on independent Brownian motion increments.

Performance with $m_{0}=4$. Figure 3.1 exhibits the performance of Algorithm 10 for this LQ-RL problem, where the solid lines and the shallow areas indicate the sample mean and the $95 \%$ confidence interval over 100 repeated experiments. The numerical results indicate that algorithm 10 manages to learn the parameters over time while incurring a desirable sublinear regret, which is consistent with our theoretical result in Theorem 3.3.4. More precisely,

- Figure 3.1a presents the logarithmic relative error of the estimate $\left(A_{\ell}, B_{\ell}\right)$ (in the Frobenius norm) after the $\ell$-th update for $\ell \in\{0, \ldots, 10\}$. One can observe that the estimate $\left(A_{\ell}, B_{\ell}\right)_{\ell}$ converge to the true parameter $\left(A^{\star}, B^{\star}\right)$ as the number of episodes increases. Our experiment shows that it is much easier to learn the state coefficient $A^{\star}$ than the control coefficient $B^{\star}$, which is consistent with the observation in [46] for other adaptive control schemes.
- Figure 3.1 b presents the relative error between the expected $\operatorname{cost} J^{\theta^{\star}}\left(\psi^{\theta_{\ell}} ; x_{0}\right)$ and the optimal expected cost $J^{\theta^{\star}}\left(\psi^{\theta^{\star}} ; x_{0}\right)$. One can see that a rough estimate of the control parameter $B^{\star}$ is often sufficient to design a nearly optimal feedback control. In particular, after the 10 -th update $(\ell=10)$, although the relative approximation errors of $A_{\ell}$ and $B_{\ell}$ are $2.7 \%$ and $24.9 \%$, respectively, the cost of $\psi^{\theta_{\ell}}$ approximates the optimal cost accurately with a relative error $0.6 \%$.
- Figure 3.1c presents the cumulative regret over episodes. One can see that the small performance gap results in a slowly growing algorithm regret. In fact, performing a linear regression for logarithms of expected regret and episode shows that the regret after the $N$-th episode is of the magnitude $\mathcal{O}\left(N^{0.34}\right)$, which is slightly better than the theoretical upper bound in Theorem 3.3.4.


Figure 3.1: Performance of Algorithm 10 for the LQ-RL problem ( $m_{0}=4$ ).

Robustness with respect to the initial batch size $m_{0}$. We next demonstrate the robustness of Algorithm 10 by performing computations with $m_{0}=1$ and fixing other settings as above. The results are shown in Figure 3.2. Note that the smaller initial batch size $m_{0}$ makes the learning more challenging. By comparing the results against those with $m_{0}=4$, one can see that our algorithm is robust and performs well with the small $m_{0}$. In particular, we see that

- Estimating parameters with fewer sample trajectories leads to larger parameter estimation errors with suboptimality gaps, especially for the first few iterations. It also leads to a wider range of $\left(A_{\ell}, B_{\ell}\right)_{\ell}$ among different algorithm executions and hence a larger variance of the algorithm regret.
- As the number of episodes increases, the estimate $\left(A_{\ell}, B_{\ell}\right)_{\ell}$ converge to the true parameter $\left(A^{\star}, B^{\star}\right)$ and the suboptimality gap quickly converges to 0 , see Figures 3.2 a and 3.2 b . The algorithm regret grows sublinearly (see Figure 3.2c), and the regret after the $N$-th episode is of the magnitude $\mathcal{O}\left(N^{0.51}\right)$. This confirms the theoretical results in Theorem 3.3.4 even for a small $m_{0}$.


Figure 3.2: Performance of Algorithm 10 for the LQ-RL problem ( $m_{0}=1$ ).

### 3.6 Appendix

### 3.6.1 Preliminaries

Here, we collect some fundamental results which are used for our analysis.
We start with a stability result for coupled FBSDEs under a generalized monotonicity condition, which is crucial for our stability analysis of feedback controls. For any given $t \in[0, T]$ and $\lambda \in[0,1]$, we consider the following FBSDE defined on $[t, T]$ : for $s \in[t, T]$,

$$
\begin{align*}
\mathrm{d} X_{s} & =\left(\lambda \bar{b}\left(s, X_{s}, Y_{s}\right)+\mathcal{I}_{s}^{b}\right) \mathrm{d} s+\bar{\sigma}(s) \mathrm{d} W_{s}+\int_{\mathbb{R}_{0}^{p}} \bar{\gamma}(s, u) \tilde{N}(\mathrm{~d} s, \mathrm{~d} u),  \tag{3.6.1a}\\
\mathrm{d} Y_{s} & =-\left(\lambda \bar{f}\left(s, X_{s}, Y_{s}\right)+\mathcal{I}_{s}^{f}\right) \mathrm{d} t+Z_{s} \mathrm{~d} W_{s}+\int_{\mathbb{R}_{0}^{p}} M_{s} \tilde{N}(\mathrm{~d} s, \mathrm{~d} u), \tag{3.6.1b}
\end{align*} \quad Y_{T}=\lambda \bar{g}\left(X_{T}\right)+\mathcal{I}^{g},
$$

with given $\xi \in L^{2}\left(\mathcal{F}_{t} ; \mathbb{R}^{n}\right)$, $\left(\mathcal{I}^{b}, \mathcal{I}^{f}\right) \in \mathcal{H}^{2}\left(\mathbb{R}^{n} \times \mathbb{R}^{n}\right), \mathcal{I}^{g} \in L^{2}\left(\mathcal{F}_{T} ; \mathbb{R}^{m}\right)$ and measurable functions $\bar{\sigma}:[0, T] \rightarrow \mathbb{R}^{n \times d}, \bar{\gamma}:[0, T] \times \mathbb{R}_{0}^{p} \rightarrow \mathbb{R}^{n}, \bar{b}, \bar{f}:[0, T] \times \mathbb{R}^{n} \times \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ and $\bar{g}: \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$.

Lemma 3.6.1. Let $K \geq 0$, for each $i \in\{1,2\}$, let $\bar{b}_{i}, \bar{f}_{i}:[0, T] \times \mathbb{R}^{n} \times \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$, $\bar{g}_{i}$ : $\mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ satisfy $\int_{0}^{T}\left(\left|\bar{b}_{i}(t, 0,0)\right|^{2}+\left|\bar{f}_{i}(t, 0,0)\right|^{2}\right) \mathrm{d} t<\infty$ and for all $t \in[0, T]$ that $\bar{f}_{i}(t, \cdot), \bar{g}_{i}$ are $K$-Lipschitz continuous, let $\bar{\sigma}_{i}:[0, T] \rightarrow \mathbb{R}^{n \times d}$ satisfy $\int_{0}^{T}\left|\bar{\sigma}_{i}(t)\right|^{2} \mathrm{~d} t<\infty$ and let $\bar{\gamma}_{i}$ : $[0, T] \times \mathbb{R}_{0}^{p} \rightarrow \mathbb{R}^{n}$ satisfy $\int_{0}^{T} \int_{\mathbb{R}_{0}^{p}}\left|\bar{\gamma}_{i}(t, u)\right|^{2} \nu(\mathrm{~d} u) \mathrm{d} t<\infty$. Assume further that there exists $\tau>0$ and a measurable function $\eta:[0, T] \times \mathbb{R}^{n} \times \mathbb{R}^{n} \times \mathbb{R}^{n} \times \mathbb{R}^{n} \rightarrow[0, \infty)$ such that for all $t \in[0, T],(x, y),\left(x^{\prime}, y^{\prime}\right) \in \mathbb{R}^{n} \times \mathbb{R}^{n}$,

$$
\begin{align*}
\left\langle\bar{b}_{1}(t, x, y)-\bar{b}_{1}\left(t, x^{\prime}, y^{\prime}\right), y-y^{\prime}\right\rangle & +\left\langle-\bar{f}_{1}(t, x, y)+\bar{f}_{1}\left(t, x^{\prime}, y^{\prime}\right), x-x^{\prime}\right\rangle \leq-\tau \eta\left(t, x, y, x^{\prime}, y^{\prime}\right), \\
\left|\bar{b}_{1}(t, x, y)-\bar{b}_{1}\left(t, x^{\prime}, y^{\prime}\right)\right| & \leq K\left(\left|x-x^{\prime}\right|+\eta\left(t, x, y, x^{\prime}, y^{\prime}\right)\right),  \tag{3.6.2}\\
\left\langle\bar{g}(x)-\bar{g}\left(x^{\prime}\right), x-x^{\prime}\right\rangle & \geq 0 . \tag{3.6.4}
\end{align*}
$$

Then there exists $C>0$, depending only on $T, K, \lambda$ and the dimensions, such that for all $t \in[0, T], \lambda_{0} \in[0,1], i \in\{1,2\}$, for every $\left(X_{i}, Y_{i}, Z_{i}, M_{i}\right) \in \mathcal{S}^{2}\left(t, T ; \mathbb{R}^{n}\right) \times \mathcal{S}^{2}\left(t, T ; \mathbb{R}^{n}\right) \times$ $\mathcal{H}^{2}\left(t, T ; \mathbb{R}^{n \times d}\right) \times \mathcal{H}_{\nu}^{2}\left(t, T ; \mathbb{R}^{n}\right)$ satisfying (3.6.1) with $\lambda=\lambda_{0}$, $(\bar{b}, \bar{\sigma}, \bar{\gamma}, \bar{f}, \bar{g})=\left(\bar{b}_{i}, \bar{\sigma}_{i}, \bar{\gamma}_{i}, \bar{f}_{i}, \bar{g}_{i}\right)$, $\xi=\xi_{i} \in L^{2}\left(\mathcal{F}_{t} ; \mathbb{R}^{n}\right), \mathcal{I}^{g}=\mathcal{I}_{i}^{g} \in L^{2}\left(\mathcal{F}_{T} ; \mathbb{R}^{n}\right)$ and $\left(\mathcal{I}^{b}, \mathcal{I}^{f}\right)=\left(\mathcal{I}_{i}^{b}, \mathcal{I}_{i}^{f}\right) \in \mathcal{H}^{2}\left(\mathbb{R}^{n} \times \mathbb{R}^{n}\right)$, we have that

$$
\begin{align*}
& \left\|X_{1}-X_{2}\right\|_{\mathcal{S}^{2}}^{2}+\left\|Y_{1}-Y_{2}\right\|_{\mathcal{S}^{2}}^{2}+\left\|Z_{1}-Z_{2}\right\|_{\mathcal{H}^{2}}^{2}+\left\|M_{1}-M_{2}\right\|_{\mathcal{H}_{\nu}^{2}}^{2} \\
& \leq C\left\{\left\|\xi_{1}-\xi_{2}\right\|_{L^{2}}^{2}+\left\|\lambda_{0}\left(\bar{g}_{1}\left(X_{2, T}\right)-\bar{g}_{2}\left(X_{2, T}\right)\right)+\mathcal{I}_{1}^{g}-\mathcal{I}_{2}^{g}\right\|_{L^{2}}^{2}+\left\|\bar{\sigma}_{1}-\bar{\sigma}_{2}\right\|_{\mathcal{H}^{2}}^{2}+\left\|\bar{\gamma}_{1}-\bar{\gamma}_{2}\right\|_{\mathcal{H}_{\nu}^{2}}^{2}\right. \\
& \quad+\left\|\lambda_{0}\left(\bar{b}_{1}\left(\cdot, X_{2}, Y_{2}\right)-\bar{b}_{2}\left(\cdot, X_{2}, Y_{2}\right)\right)+\mathcal{I}_{1}^{b}-\mathcal{I}_{2}^{b}\right\|_{\mathcal{H}^{2}}^{2} \\
& \left.\quad+\left\|\lambda_{0}\left(\bar{f}_{1}\left(\cdot, X_{2}, Y_{2}\right)-\bar{f}_{2}\left(\cdot, X_{2}, Y_{2}\right)\right)+\mathcal{I}_{1}^{f}-\mathcal{I}_{2}^{f}\right\|_{\mathcal{H}^{2}}^{2}\right\} . \tag{3.6.5}
\end{align*}
$$

Proof. Throughout this proof, let $C$ be a generic constant depending only on $T, K, \lambda$ and the dimensions, let $t \in[0, T], \lambda_{0} \in[0,1]$, let $(\delta X, \delta Y, \delta Z, \delta M)=\left(X_{1}-X_{2}, Y_{1}-Y_{2}, Z_{1}-Z_{2}, M_{1}-\right.$ $\left.M_{2}\right), \delta \xi=\xi_{1}-\xi_{2}, \delta \sigma=\bar{\sigma}_{1}-\bar{\sigma}_{2}, \delta \gamma=\bar{\gamma}_{1}-\bar{\gamma}_{2}, \delta \mathcal{I}^{g}=\mathcal{I}_{1}^{g}-\mathcal{I}_{2}^{g}$, and for each $s \in[t, T]$ let $\delta \mathcal{I}_{s}^{b}=\mathcal{I}_{1, s}^{b}-\mathcal{I}_{2, s}^{b}, \delta \mathcal{I}_{s}^{f}=\mathcal{I}_{1, s}^{f}-\mathcal{I}_{2, s}^{f}, \bar{b}_{1}\left(\Theta_{1, s}\right)=\bar{b}_{1}\left(t, X_{1, s}, Y_{1, s}\right), \bar{b}_{1}\left(\Theta_{2, s}\right)=\bar{b}_{1}\left(t, X_{2, s}, Y_{2, s}\right)$ and $\bar{b}_{2}\left(\Theta_{2, s}\right)=\bar{b}_{2}\left(t, X_{2, s}, Y_{2, s}\right)$. Similarly, we introduce the notation $\bar{f}_{1}\left(\Theta_{1, s}\right), \bar{f}_{1}\left(\Theta_{2, s}\right), \bar{f}_{2}\left(\Theta_{2, s}\right)$ for $s \in[t, T]$.

By applying Itô's formula to $\left(\left\langle Y_{1, s}-Y_{2, s}, X_{1, s}-X_{2, s}\right\rangle\right)_{s \in[t, T]}$, we can obtain from (3.6.1) that

$$
\begin{aligned}
\mathbb{E} & {\left[\left\langle\lambda_{0}\left(\bar{g}_{1}\left(X_{1, T}\right)-\bar{g}_{2}\left(X_{2, T}\right)\right)+\delta \mathcal{I}^{g}, \delta X_{T}\right\rangle-\left\langle\delta Y_{t}, \delta \xi\right\rangle\right] } \\
= & \mathbb{E}\left[\int _ { t } ^ { T } \left(\left\langle\lambda_{0}\left(\bar{b}_{1}\left(\Theta_{1, s}\right)-\bar{b}_{2}\left(\Theta_{2, s}\right)\right)+\delta \mathcal{I}_{s}^{b}, \delta Y_{s}\right\rangle-\left\langle\lambda_{0}\left(\bar{f}_{1}\left(\Theta_{1, s}\right)-\bar{f}_{2}\left(\Theta_{2, s}\right)\right)+\delta \mathcal{I}_{s}^{f}, \delta X_{s}\right\rangle\right.\right. \\
& \left.\left.+\left\langle\delta \sigma(s), \delta Z_{s}\right\rangle+\int_{\mathbb{R}_{0}^{p}}\left\langle\delta \gamma(s, u), \delta M_{s}\right\rangle \nu(\mathrm{d} u)\right) \mathrm{d} s\right] \\
\leq & \mathbb{E}\left[\int _ { t } ^ { T } \left(-\lambda_{0} \tau \eta\left(s, X_{1, s}, Y_{1, s}, X_{2, s}, Y_{2, s}\right)+\left\langle\lambda_{0}\left(\bar{b}_{1}\left(\Theta_{2, s}\right)-\bar{b}_{2}\left(\Theta_{2, s}\right)\right)+\delta \mathcal{I}_{s}^{b}, \delta Y_{s}\right\rangle\right.\right. \\
& \left.\left.-\left\langle\lambda_{0}\left(\bar{f}_{1}\left(\Theta_{2, s}\right)-\bar{f}_{2}\left(\Theta_{2, s}\right)\right)+\delta \mathcal{I}_{s}^{f}, \delta X_{s}\right\rangle+\left\langle\delta \sigma(s), \delta Z_{s}\right\rangle+\int_{\mathbb{R}_{0}^{p}}\left\langle\delta \gamma(s, u), \delta M_{s}\right\rangle \nu(\mathrm{d} u)\right) \mathrm{d} s\right]
\end{aligned}
$$

where for the last inequality, we have added and subtracted the terms $\lambda_{0} \bar{b}_{1}\left(\Theta_{2, s}\right)$ and $-\lambda_{0} \bar{f}_{1}\left(\Theta_{2, s}\right)$, and applied (3.6.2). Then, we can further deduce from (3.6.4) that

$$
\begin{aligned}
& \lambda_{0} \tau \mathbb{E}\left[\int_{t}^{T} \eta\left(s, X_{1, s}, Y_{1, s}, X_{2, s}, Y_{2, s}\right) \mathrm{d} s\right] \\
& \leq-\mathbb{E}\left[\left\langle\lambda_{0}\left(\bar{g}_{1}\left(X_{2, T}\right)-\bar{g}_{2}\left(X_{2, T}\right)\right)+\delta \mathcal{I}^{g}, \delta X_{T}\right\rangle-\left\langle\delta Y_{t}, \delta \xi\right\rangle\right] \\
&+\mathbb{E}\left[\int _ { t } ^ { T } \left(\left\langle\lambda_{0}\left(\bar{b}_{1}\left(\Theta_{2, s}\right)-\bar{b}_{2}\left(\Theta_{2, s}\right)\right)+\delta \mathcal{I}_{s}^{b}, \delta Y_{s}\right\rangle-\left\langle\lambda_{0}\left(\bar{f}_{1}\left(\Theta_{2, s}\right)-\bar{f}_{2}\left(\Theta_{2, s}\right)\right)+\delta \mathcal{I}_{s}^{f}, \delta X_{s}\right\rangle\right.\right. \\
&\left.\left.+\left\langle\delta \sigma(s), \delta Z_{s}\right\rangle+\int_{\mathbb{R}_{0}^{p}}\left\langle\delta \gamma(s, u), \delta M_{s}\right\rangle \nu(\mathrm{d} u)\right) \mathrm{d} s\right]
\end{aligned}
$$

from which we can apply Young's inequality and obtain for all $\varepsilon>0$ that

$$
\begin{aligned}
& \lambda_{0} \mathbb{E}\left[\int_{t}^{T} \eta\left(s, X_{1, s}, Y_{1, s}, X_{2, s}, Y_{2, s}\right) \mathrm{d} s\right] \\
& \leq \varepsilon\left(\left\|\delta X_{T}\right\|_{L^{2}}^{2}+\left\|\delta Y_{t}\right\|_{L^{2}}^{2}+\|\delta X\|_{\mathcal{H}^{2}}^{2}+\|\delta Y\|_{\mathcal{H}^{2}}^{2}+\|\delta Z\|_{\mathcal{H}^{2}}^{2}+\|\delta M\|_{\mathcal{H}_{\nu}^{2}}^{2}\right)+C \mathrm{RHS} / \varepsilon
\end{aligned}
$$

where RHS denotes the terms at the right-hand side of (3.6.5).
By (3.6.3) and a standard stability estimate of (3.6.1a), we can deduce that

$$
\begin{align*}
\|\delta X\|_{\mathcal{S}^{2}}^{2} & \leq C\left(\lambda_{0} \mathbb{E}\left[\int_{t}^{T} \eta\left(s, X_{1, s}, Y_{1, s}, X_{2, s}, Y_{2, s}\right) \mathrm{d} s\right]+\mathrm{RHS}\right)  \tag{3.6.6}\\
& \leq \varepsilon C\left(\left\|\delta Y_{t}\right\|_{L^{2}}^{2}+\|\delta Y\|_{\mathcal{H}^{2}}^{2}+\|\delta Z\|_{\mathcal{H}^{2}}^{2}+\|\delta M\|_{\mathcal{H}_{\nu}^{2}}^{2}\right)+C \mathrm{RHS} / \varepsilon
\end{align*}
$$

for all small enough $\varepsilon>0$. Moreover, by the Lipschitz continuity of $\bar{f}_{1}, \bar{g}_{1}$ and the stability estimate of (3.6.1b) (see e.g. [133, Proposition A4]), we have that

$$
\begin{aligned}
& \|\delta Y\|_{\mathcal{S}^{2}}^{2}+\|\delta Z\|_{\mathcal{H}^{2}}^{2}+\|\delta M\|_{\mathcal{H}_{\nu}^{2}}^{2} \\
& \leq C\left(\left\|\lambda_{0}\left(\bar{g}_{1}\left(X_{1, T}\right)-\bar{g}_{2}\left(X_{2, T}\right)\right)+\delta \mathcal{I}^{g}\right\|_{L^{2}}^{2}+\left\|\lambda_{0}\left(\bar{f}_{1}\left(\cdot, X_{1}, Y_{2}\right)-\bar{f}_{2}\left(\cdot, X_{2}, Y_{2}\right)\right)+\delta \mathcal{I}^{f}\right\|_{\mathcal{H}^{2}}^{2}\right) \\
& \leq C\left(\|\delta X\|_{\mathcal{S}^{2}}^{2}+\left\|\lambda_{0}\left(\bar{g}_{1}\left(X_{2, T}\right)-\bar{g}_{2}\left(X_{2, T}\right)\right)+\delta \mathcal{I}^{g}\right\|_{L^{2}}^{2}+\left\|\lambda_{0}\left(\bar{f}_{1}\left(\cdot, X_{2}, Y_{2}\right)-\bar{f}_{2}\left(\cdot, X_{2}, Y_{2}\right)\right)+\delta \mathcal{I}^{f}\right\|_{\mathcal{H}^{2}}^{2}\right) \\
& \leq C \text { RHS }
\end{aligned}
$$

where we have applied (3.6.6) with a sufficiently small $\varepsilon$ for the last inequality. This completes the desired stability estimate.

We then present a version of Burkholder's inequality for the $\|\cdot\|_{L^{q}}$-norm of stochastic integrals, which not only extends [27, Corollary 2.2] to stochastic integrals with respect to general Poisson random measures on $[0, T] \times \mathbb{R}_{0}^{p}$, but also improves the bounding constants there with a sharper dependence on the index $q$.

Lemma 3.6.2. For all $v \in \mathcal{H}^{2}\left(0, T ; \mathbb{R}^{d}\right)$, $w \in \mathcal{H}_{\nu}^{2}(0, T ; \mathbb{R})$ and $q \geq 2$, we have

$$
\begin{align*}
\mathbb{E}\left[\left|\int_{0}^{T} v_{t}^{\top} \mathrm{d} W_{t}\right|^{q}\right] & \leq C_{q} \mathbb{E}\left[\left(\int_{0}^{T}\left|v_{t}\right|^{2} \mathrm{~d} t\right)^{q / 2}\right]  \tag{3.6.7}\\
\mathbb{E}\left[\left|\int_{0}^{T} \int_{\mathbb{R}_{0}^{p}} w(t, u) \tilde{N}(\mathrm{~d} t, \mathrm{~d} u)\right|^{q}\right] & \leq \tilde{C}_{q}\left(\mathbb{E}\left[\int_{0}^{T} \int_{\mathbb{R}_{0}^{p}}|w(t, u)|^{q} \nu(\mathrm{~d} u) \mathrm{d} t\right]\right. \\
& \left.+\mathbb{E}\left[\left(\int_{0}^{T} \int_{\mathbb{R}_{0}^{p}}|w(t, u)|^{2} \nu(\mathrm{~d} u) \mathrm{d} t\right)^{q / 2}\right]\right) \tag{3.6.8}
\end{align*}
$$

where $C_{q}=(\sqrt{e / 2} q)^{q}$ and $\tilde{C}_{q}=21 e^{q} q^{2 q}$.
Proof. Recall that Burkholder's inequality in [23, Theorem 4.2.12] shows for all $1 \leq q<\infty$ and for every local martingale $\left(M_{t}\right)_{t \in[0, T]}$ that $\mathbb{E}\left[\left|M_{T}^{*}\right|^{q}\right] \leq C_{q} \mathbb{E}\left[[M, M]_{T}^{q / 2}\right]$, where $M_{T}^{*}=$ $\sup _{t \in[0, T]}\left|M_{t}\right|,[M, M]$ is the quadratic variation of $M$, and $C_{q}=(\sqrt{10 q})^{q}$ for $q \in[1,2)$ and $C_{q}=(\sqrt{e / 2} q)^{q}$ for $q \in[2, \infty)$. Hence we can obtain (3.6.7) by setting $M_{t}=\int_{0}^{t} v_{s}^{\top} \mathrm{d} W_{s}$ for all $t \in[0, T]$, whose quadratic variation process is given by $[M, M]_{T}=\int_{0}^{T}|v|^{2} \mathrm{~d} t$.

We proceed to establish (3.6.8) by following the arguments of [27, Lemma 2.1]. For all $r \geq 1$ and $t \in[0, T]$, let $K_{t}^{(r)}=\int_{0}^{t} \int_{\mathbb{R}_{0}^{p}} w(s, u)^{r} \tilde{N}(\mathrm{~d} s, \mathrm{~d} u)$. For any given $r \geq 1$ and $q \geq 2$, we
can obtain from Burkholder's inequality and $\tilde{N}(\mathrm{~d} t, \mathrm{~d} u)=N(\mathrm{~d} t, \mathrm{~d} u)-\nu(\mathrm{d} u) \mathrm{d} t$ that

$$
\begin{aligned}
\mathbb{E}\left[\left|\left(K^{(r)}\right)_{T}^{*}\right|^{q}\right] & \leq C_{q} \mathbb{E}\left[\left(\int_{0}^{T} \int_{\mathbb{R}_{0}^{p}}|w(t, u)|^{2 r} N(\mathrm{~d} t, \mathrm{~d} u)\right)^{q / 2}\right] \\
& =C_{q} \mathbb{E}\left[\left(\int_{0}^{T} \int_{\mathbb{R}_{0}^{p}}|w(t, u)|^{2 r} \tilde{N}(\mathrm{~d} t, \mathrm{~d} u)+\int_{0}^{T} \int_{\mathbb{R}_{0}^{p}}|w(t, u)|^{2 r} \nu(\mathrm{~d} u) \mathrm{d} t\right)^{q / 2}\right] \\
& \leq 2^{\frac{q}{2}-1} C_{q} \mathbb{E}\left[\left|\left(K^{(2 r)}\right)_{T}^{*}\right|^{q / 2}\right]+2^{\frac{q}{2}-1} C_{q} \mathbb{E}\left[\left.\left.\left|\int_{0}^{T} \int_{\mathbb{R}_{0}^{p}}\right| w(t, u)\right|^{2 r} \nu(\mathrm{~d} u) \mathrm{d} t\right|^{q / 2}\right] .
\end{aligned}
$$

Hence, recursively applying the above estimate yields for all $q \geq 2$ and $n \in \mathbb{N}$ with $q / 2^{n-1} \geq 2$ that

$$
\begin{align*}
\mathbb{E}\left[\left|\int_{0}^{T} \int_{\mathbb{R}_{0}^{p}} w(t, u) \tilde{N}(\mathrm{~d} t, \mathrm{~d} u)\right|^{q}\right] \leq & \left(\prod_{j=1}^{n} 2^{\frac{q}{2^{j}-1}} C_{\frac{q}{2^{j-1}}}\right) \mathbb{E}\left[\left|\left(K^{\left(2^{n}\right)}\right)_{T}^{*}\right|^{q / 2^{n}}\right] \\
& +\sum_{k=1}^{n}\left(\prod_{j=1}^{k} 2^{\left.\frac{q}{2^{j}-1} C_{\frac{q}{2^{j-1}}}\right) \mathbb{E}\left[\left.\left.\left|\int_{0}^{T} \int_{\mathbb{R}_{0}^{p}}\right| w(t, u)\right|^{2^{k}} \nu(\mathrm{~d} u) \mathrm{d} t\right|^{q / 2^{k}}\right]} .\right. \tag{3.6.9}
\end{align*}
$$

Now let $q \geq 2$ be fixed and set $n=\left\lfloor\log _{2} q\right\rfloor$ such that $q \in\left[2^{n}, 2^{n+1}\right)$. Since $q / 2^{n} \in[1,2)$, the constant $C_{q / 2^{n}}$ in Burkholder's inequality satisfies $C_{q / 2^{n}} \leq 20$, from which we can show that (see [27, Lemma 2.1]):

$$
\mathbb{E}\left[\left|\left(K^{\left(2^{n}\right)}\right)_{T}^{*}\right|^{q / 2^{n}}\right] \leq 20 \mathbb{E}\left[\int_{0}^{T} \int_{\mathbb{R}_{0}^{p}}|w(t, u)|^{q} \nu(\mathrm{~d} u) \mathrm{d} t\right] .
$$

Moreover, by proceeding along the lines of [27, Corollary 2.2], we obtain for all $k=1, \ldots, n$ that

$$
\begin{aligned}
& \mathbb{E}\left[\left.\left.\left|\int_{0}^{T} \int_{\mathbb{R}_{0}^{p}}\right| w(t, u)\right|^{2^{k}} \nu(\mathrm{~d} u) \mathrm{d} t\right|^{q / 2^{k}}\right] \\
& \leq \mathbb{E}\left[\int_{0}^{T} \int_{\mathbb{R}_{0}^{p}}|w(t, u)|^{q} \nu(\mathrm{~d} u) \mathrm{d} t\right]+\mathbb{E}\left[\left.\left.\left|\int_{0}^{T} \int_{\mathbb{R}_{0}^{p}}\right| w(t, u)\right|^{2} \nu(\mathrm{~d} u) \mathrm{d} t\right|^{q / 2}\right] .
\end{aligned}
$$

Hence, we can deduce from (3.6.9) that

$$
\begin{aligned}
\mathbb{E}\left[\left|\int_{0}^{T} \int_{\mathbb{R}_{0}^{p}} w(t, u) \tilde{N}(\mathrm{~d} t, \mathrm{~d} u)\right|^{q}\right] \leq & 21 \sum_{k=1}^{\left\lceil\log _{2} q\right\rceil}\left(\prod_{j=1}^{k} 2^{\frac{q}{2 j}-1} C_{\frac{q}{2 j-1}}\right)\left(\mathbb{E}\left[\int_{0}^{T} \int_{\mathbb{R}_{0}^{p}}|w(t, u)|^{q} \nu(\mathrm{~d} u) \mathrm{d} t\right]\right. \\
& \left.+\mathbb{E}\left[\left.\left.\left|\int_{0}^{T} \int_{\mathbb{R}_{0}^{p}}\right| w(t, u)\right|^{2} \nu(\mathrm{~d} u) \mathrm{d} t\right|^{q / 2}\right]\right) .
\end{aligned}
$$

We now obtain an upper bound of the constant $21 \sum_{k=1}^{\left.\log _{2} q\right\rfloor}\left(\prod_{j=1}^{k} 2^{\frac{q}{2 j}-1} C_{\frac{q}{2^{j-1}}}\right)$ as follows:

$$
\begin{aligned}
& 21 \sum_{k=1}^{\left\lfloor\log _{2} q\right\rfloor}\left(\prod_{j=1}^{k} 2^{\frac{q}{2^{j}-1}} C_{\frac{q}{2 j-1}}\right)=21 \sum_{k=1}^{\left\lfloor\log _{2} q\right\rfloor} \prod_{j=1}^{k} 2^{\frac{q}{2^{j}-1}}\left(\sqrt{\frac{e}{2}} \frac{q}{2^{j-1}}\right)^{\frac{q}{2 j-1}} \\
& \leq 21\left(\sum_{k=1}^{\left\lfloor\log _{2} q\right\rfloor} 2^{-k}\right) e^{\sum_{j=1}^{\left\lfloor\log _{2} q\right\rfloor} \frac{q}{2^{j}}} \prod_{j=1}^{\left\lfloor\log _{2} q\right\rfloor}\left(\frac{q}{2^{j-1}}\right)^{\frac{q}{2 j-1}} \leq 21 e^{q} 2^{\sum_{j=1}^{\left\lfloor\log _{2} q\right\rfloor} \frac{q}{2^{j-1}} \log _{2}\left(\frac{q}{2 j-1}\right)} \\
& \leq 21 e^{q} 2^{\sum_{j=1}^{\left\lfloor\log _{2} q\right\rfloor} \frac{q}{2^{j-1}} \log _{2} q} \leq 21 e^{q} 2^{2 q \log _{2} q}=21 e^{q} q^{2 q}:=\tilde{C}_{q},
\end{aligned}
$$

which leads us to the desired conclusion.
The following lemma estimates the tail behaviors of solutions to SDEs with jumps. The result has been established in Lemma 2.1 and Theorem 2.8 of [111] for SDEs with time homogenous coefficients and bounded Lipschitz continuous functions $\mathfrak{f}$ via Malliavin Calculus, which can be extended to SDEs with time inhomogeneous coefficients and unbounded $\mathfrak{f}$ (via Fatou's lemma) in a straightforward manner.

Lemma 3.6.3. Let $T \geq 0$ and $b:[0, T] \times \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$, $\sigma:[0, T] \rightarrow \mathbb{R}^{n \times d}, \gamma:[0, T] \times \mathbb{R}_{0}^{p} \rightarrow \mathbb{R}^{n}$ be measurable functions such that there exist $K, \sigma_{\max } \geq 0$ and a measurable function $\bar{\gamma}: \mathbb{R}_{0}^{p} \rightarrow \mathbb{R}$ satisfying for all $(t, u) \in[0, T] \times \mathbb{R}_{0}^{p}$, $x, x^{\prime} \in \mathbb{R}^{n}$ that $|b(t, 0)| \leq K,\left|b(t, x)-b\left(t, x^{\prime}\right)\right| \leq$ $K\left|x-x^{\prime}\right|,|\sigma(t)| \leq \sigma_{\max }$ and $|\gamma(t, u)| \leq \bar{\gamma}(u)$, $\nu$-a.e.. Let $\beta:[0, \infty) \rightarrow[0, \infty]$ be defined by $\beta(\lambda):=\int_{\mathbb{R}_{0}^{p}}\left(e^{\lambda \bar{\gamma}(u)}-\lambda \bar{\gamma}(u)-1\right) \nu(\mathrm{d} u)$ for any $\lambda \geq 0$. Assume that $\beta(\lambda)<\infty$ for some $\lambda>0$.

Then there exists a constant $C>0$, depending only on $K$ and $T$, such that for all $x \in \mathbb{R}^{n}$, the unique solution $X^{x} \in \mathcal{S}^{2}\left(\mathbb{R}^{n}\right)$ to the following $S D E$

$$
\mathrm{d} X_{t}=b\left(t, X_{t}\right) \mathrm{d} t+\sigma(t) \mathrm{d} W_{t}+\int_{\mathbb{R}_{0}^{p}} \gamma(t, u) \tilde{N}(\mathrm{~d} t, \mathrm{~d} u), \quad t \in[0, T], \quad X_{0}=x
$$

satisfies for every Lipschitz continuous function $\mathfrak{f}:\left(\mathbb{D}\left([0, T] ; \mathbb{R}^{n}\right), d_{\infty}\right) \rightarrow \mathbb{R}$ that

$$
\begin{equation*}
\mathbb{E}\left[e^{\lambda\left(f\left(X^{x}\right)-\mathbb{E}\left[f\left(X^{x}\right)\right]\right)}\right] \leq e^{C \eta\left(C \lambda\|f\|_{\text {Lip }}\right)} \quad \forall \lambda>0 \tag{3.6.10}
\end{equation*}
$$

where $\mathbb{D}\left([0, T] ; \mathbb{R}^{n}\right)$ is the space of $\mathbb{R}^{n}$-valued càdlàg functions on $[0, T], d_{\infty}$ is the uniform metric defined by $d_{\infty}\left(\rho_{1}, \rho_{2}\right):=\sup _{t \in[0, T]}\left|\rho_{1}(t)-\rho_{2}(t)\right|$ for any $\rho_{1}, \rho_{2} \in \mathbb{D}\left([0, T] ; \mathbb{R}^{n}\right)$, $\|\mathfrak{f}\|_{\text {Lip }}$ is the constant defined by $\|\mathfrak{f}\|_{\text {Lip }}:=\sup _{\rho_{1} \neq \rho_{2}} \frac{\left|\mathfrak{f}\left(\rho_{1}\right)-\mathrm{f}\left(\rho_{2}\right)\right|}{d_{\infty}\left(\rho_{1}, \rho_{2}\right)}$, and $\eta:[0, \infty) \rightarrow[0, \infty]$ is the function defined by $\eta(\lambda):=\beta(\lambda)+\sigma_{\max }^{2} \lambda^{2} / 2$ for any $\lambda \geq 0$.

The next lemma presents a concentration inequality for the sum of independent subWeibull random variables, which follows directly from Theorem 3.1 and Proposition A3 in [98].

Lemma 3.6.4. Let $\alpha \in(0,1], N \in \mathbb{N}$ and $X_{1}, \ldots, X_{N} \in \operatorname{subW}(\alpha)$ be independent random variables satisfying $\mathbb{E}\left[X_{i}\right]=0$ for all $i=1, \ldots, N$. Then there exists a constant $C \geq 0$, depending only on $\alpha$, such that

$$
\mathbb{P}\left(\left|\sum_{i=1}^{N} X_{i}\right| \geq \varepsilon^{\prime}\right) \leq 2 \exp \left(-C \min \left\{\frac{\left(\varepsilon^{\prime}\right)^{2}}{\sum_{i=1}^{N}\left\|X_{i}\right\|_{\Psi_{\alpha}}^{2}},\left(\frac{\varepsilon^{\prime}}{\max _{i}\left\|X_{i}\right\|_{\Psi_{\alpha}}}\right)^{\alpha}\right\}\right), \quad \forall \varepsilon^{\prime} \geq 0
$$

### 3.6.2 Proofs of Lemmas 3.2.8, 3.3.5, 3.3.6, 3.3.7, 3.3.8

Proof of Lemma 3.2.8. We start by establishing the regularity of $\mathbb{R}^{n} \times \mathbb{R}^{k} \ni(x, z) \mapsto$ $f\left(t, x, \partial_{z} f^{*}(t, x, z)\right) \in \mathbb{R} \cup\{\infty\}$ for a given $t \in[0, T]$. Observe that for all $(t, x) \in[0, T] \times \mathbb{R}^{n}$, $f(t, x, \cdot)$ is proper, convex, and lower semicontinuous, which along with the Fenchel-Young identity implies

$$
\begin{equation*}
f\left(t, x, \partial_{z} f^{*}(t, x, z)\right)=\left\langle z, \partial_{z} f^{*}(t, x, z)\right\rangle-f^{*}(t, x, z) \in \mathbb{R}, \quad \forall(t, x, z) \in[0, T] \times \mathbb{R}^{n} \times \mathbb{R}^{k} \tag{3.6.11}
\end{equation*}
$$

Given $t \in[0, T]$ and $\left(x_{1}, z_{1}\right),\left(x_{2}, z_{2}\right) \in \mathbb{R}^{n} \times \mathbb{R}^{k}$, by (3.6.11),

$$
\begin{aligned}
& \left|f\left(t, x_{1}, \partial_{z} f^{*}\left(t, x_{1}, z_{1}\right)\right)-f\left(t, x_{2}, \partial_{z} f^{*}\left(t, x_{2}, z_{2}\right)\right)\right| \\
& \quad \leq\left|\left\langle z_{1}, \partial_{z} f^{*}\left(t, x_{1}, z_{1}\right)\right\rangle-\left\langle z_{2}, \partial_{z} f^{*}\left(t, x_{2}, z_{2}\right)\right\rangle\right|+\left|f^{*}\left(t, x_{1}, z_{1}\right)-f^{*}\left(t, x_{2}, z_{1}\right)\right| \\
& \quad+\left|f^{*}\left(t, x_{2}, z_{1}\right)-f^{*}\left(t, x_{2}, z_{2}\right)\right|
\end{aligned}
$$

We now estimate all the terms on the right hand side of the above inequality. By the Lipchitz continuity and local boundedness of $\partial_{z} f^{*}(t, \cdot)$ (see the proof of Lemma 3.2.3), we can obtain the following upper bound for the first and third terms:

$$
\begin{align*}
& \left|\left\langle z_{1}, \partial_{z} f^{*}\left(t, x_{1}, z_{1}\right)\right\rangle-\left\langle z_{2}, \partial_{z} f^{*}\left(t, x_{2}, z_{2}\right)\right\rangle\right|+\left|f^{*}\left(t, x_{2}, z_{1}\right)-f^{*}\left(t, x_{2}, z_{2}\right)\right| \\
& \leq\left|\left\langle z_{1}-z_{2}, \partial_{z} f^{*}\left(t, x_{1}, z_{1}\right)\right\rangle\right|+\left|\left\langle z_{2}, \partial_{z} f^{*}\left(t, x_{1}, z_{1}\right)-\partial_{z} f^{*}\left(t, x_{2}, z_{2}\right)\right\rangle\right|  \tag{3.6.12}\\
& \quad+\left|f^{*}\left(t, x_{2}, z_{1}\right)-f^{*}\left(t, x_{2}, z_{2}\right)\right| \\
& \leq C\left(1+\left|x_{1}\right|+\left|x_{2}\right|+\left|z_{1}\right|+\left|z_{2}\right|\right)\left(\left|x_{1}-x_{2}\right|+\left|z_{1}-z_{2}\right|\right),
\end{align*}
$$

where the last inequality is by the mean value theorem. Moreover, by applying (3.6.11) to $f^{*}\left(t, x_{1}, z_{1}\right)$ and by the definition of $f^{*}\left(t, x_{2}, z_{1}\right)$ in (3.2.4), (H.2(3)), the linear growth of $\partial_{z} f^{*}(t, \cdot)$,

$$
\begin{align*}
f^{*}\left(t, x_{1}, z_{1}\right)-f^{*}\left(t, x_{2}, z_{1}\right) \leq & \left\langle z_{1}, \partial_{z} f^{*}\left(t, x_{1}, z_{1}\right)\right\rangle-f\left(t, x_{1}, \partial_{z} f^{*}\left(t, x_{1}, z_{1}\right)\right) \\
& -\left(\left\langle z_{1}, \partial_{z} f^{*}\left(t, x_{1}, z_{1}\right)\right\rangle-f\left(t, x_{2}, \partial_{z} f^{*}\left(t, x_{1}, z_{1}\right)\right)\right) \\
= & -f_{0}\left(t, x_{1}, \partial_{z} f^{*}\left(t, x_{1}, z_{1}\right)\right)+f_{0}\left(t, x_{2}, \partial_{z} f^{*}\left(t, x_{1}, z_{1}\right)\right)  \tag{3.6.13}\\
\leq & C\left(1+\left|x_{1}\right|+\left|x_{2}\right|+\left|z_{1}\right|\right)\left|x_{1}-x_{2}\right| .
\end{align*}
$$

Then, by interchanging the roles of $x_{1}, x_{2}$ in (3.6.13) and taking account of (3.6.12), we can obtain the following estimate for all $t \in[0, T],\left(x_{1}, z_{1}\right),\left(x_{2}, z_{2}\right) \in \mathbb{R}^{n} \times \mathbb{R}^{k}$ :

$$
\begin{aligned}
& \left|f\left(t, x_{1}, \partial_{z} f^{*}\left(t, x_{1}, z_{1}\right)\right)-f\left(t, x_{2}, \partial_{z} f^{*}\left(t, x_{2}, z_{2}\right)\right)\right| \\
& \quad \leq C\left(1+\left|x_{1}\right|+\left|x_{2}\right|+\left|z_{1}\right|+\left|z_{2}\right|\right)\left(\left|x_{1}-x_{2}\right|+\left|z_{1}-z_{2}\right|\right)
\end{aligned}
$$

Therefore, by (3.2.10), (3.2.17) and (3.2.19), for all $t \in[0, T]$ and $x, x^{\prime} \in \mathbb{R}^{n}$,

$$
\begin{aligned}
& \left|f(t, x, \psi(t, x))-f\left(t, x^{\prime}, \tilde{\psi}\left(t, x^{\prime}\right)\right)\right|=\left|f\left(t, x, \phi\left(t, x, Y_{t}^{t . x}\right)\right)-f\left(t, x^{\prime}, \tilde{\phi}\left(t, x^{\prime}, \tilde{Y}_{t}^{t . x^{\prime}}\right)\right)\right| \\
& =\left|f\left(t, x, \partial_{z} f^{*}\left(t, x,-b_{2}(t)^{\top} Y_{t}^{t . x}\right)\right)-f\left(t, x^{\prime}, \partial_{z} f^{*}\left(t, x^{\prime},-\tilde{b}_{2}(t)^{\top} \tilde{Y}_{t}^{t . x^{\prime}}\right)\right)\right| \\
& \leq C\left(1+|x|+\left|x^{\prime}\right|+\left|b_{2}(t)^{\top} Y_{t}^{t . x}\right|+\left|\tilde{b}_{2}(t)^{\top} \tilde{Y}_{t}^{t \cdot x^{\prime}}\right|\right)\left(\left|x-x^{\prime}\right|+\left|b_{2}(t)^{\top} Y_{t}^{t . x}-\tilde{b}_{2}(t)^{\top} \tilde{Y}_{t}^{t . x^{\prime}}\right|\right) \\
& \leq C\left(1+|x|+\left|x^{\prime}\right|+\left\|Y^{t . x}\right\|_{\mathcal{S}^{2}}+\left\|\tilde{Y}^{t . x^{\prime}}\right\|_{\mathcal{S}^{2}}\right) \\
& \quad \times\left(\left|x-x^{\prime}\right|+\left\|b_{2}-b_{2}\right\|_{L^{\infty}}\left\|Y^{t . x}\right\|_{\mathcal{S}^{2}}+\left\|Y^{t . x}-\tilde{Y}^{t . x^{\prime}}\right\|_{\mathcal{S}^{2}}\right) \\
& \leq C\left(1+|x|+\left|x^{\prime}\right|\right)\left(\left|x-x^{\prime}\right|+\mathcal{E}_{\text {per }}(1+|x|)\right),
\end{aligned}
$$

which along with Young's inequality leads to the desired conclusion.
Proof of Lemma 3.3.5. It suffices to show the statement for processes $X, Y$ such that

$$
\|X\|_{L^{2}(0, T)},\|Y\|_{L^{2}(0, T)} \in \operatorname{subW}(\alpha)
$$

with $\|X\|_{L^{2}(0, T)}:=\left(\int_{0}^{T}|X|^{2} \mathrm{~d} t\right)^{\frac{1}{2}}$ and $\|Y\|_{L^{2}(0, T)}:=\left(\int_{0}^{T}|Y|^{2} \mathrm{~d} t\right)^{\frac{1}{2}}$, as otherwise the right-hand side of the inequality would be infinity. Since $\|\cdot\|_{\Psi_{\alpha}}$ is a quasi-norm for any $\alpha>0$, we shall assume without loss of generality that $\left\|\|X\|_{L^{2}(0, T)}\right\|_{\Psi_{\alpha}}=\| \| Y\left\|_{L^{2}(0, T)}\right\|_{\Psi_{\alpha}}=1$. Then, we can deduce from Hölder's inequality and Young's inequality that

$$
\begin{aligned}
& \mathbb{E}\left[\exp \left(\left|\int_{0}^{T} X Y \mathrm{~d} t\right|^{\frac{\alpha}{2}}\right)\right] \leq \mathbb{E}\left[\exp \left(\left|\|X\|_{L^{2}(0, T)}\|Y\|_{L^{2}(0, T)}\right|^{\frac{\alpha}{2}}\right)\right] \\
& \leq \mathbb{E}\left[\exp \left(\frac{1}{2}\|X\|_{L^{2}(0, T)}^{\alpha}+\frac{1}{2}\|Y\|_{L^{2}(0, T)}^{\alpha}\right)\right]=\mathbb{E}\left[\exp \left(\frac{1}{2}\|X\|_{L^{2}(0, T)}^{\alpha}\right) \exp \left(\frac{1}{2}\|Y\|_{L^{2}(0, T)}^{\alpha}\right)\right] \\
& \leq\left(\mathbb{E}\left[\exp \left(\|X\|_{L^{2}(0, T)}^{\alpha}\right)\right]\right)^{\frac{1}{2}}\left(\mathbb{E}\left[\exp \left(\|Y\|_{L^{2}(0, T)}^{\alpha}\right)\right]\right)^{\frac{1}{2}} \leq 2,
\end{aligned}
$$

which implies that $\left\|\int_{0}^{T} X Y \mathrm{~d} t\right\|_{\Psi_{\alpha / 2}} \leq 1$ and finishes the proof.
Proof of Lemma 3.3.6. Note that (3.3.9) and Hölder's inequality suggest that it suffices to estimate the growth of $\|\cdot\|_{L^{q}}$-norms of the stochastic integrals for $q \geq 2$. Hence, by (3.6.7), there exists a constant $C$ such that for all $q \geq 2$,

$$
\begin{aligned}
q^{-2}\left\|\int_{0}^{T} X_{t} \sigma^{\top} \mathrm{d} W_{t}\right\|_{L^{q}} & \leq q^{-2} C q\left\|\left(\int_{0}^{T}\left|X_{t} \sigma\right|^{2} \mathrm{~d} t\right)^{\frac{1}{2}}\right\|_{L^{q}} \leq C|\sigma| \sup _{q \geq 1}\left(q^{-1}\left\|\left(\int_{0}^{T}\left|X_{t}\right|^{2} \mathrm{~d} t\right)^{\frac{1}{2}}\right\|_{L^{q}}\right) \\
& \leq C|\sigma|\left\|\left(\int_{0}^{T}\left|X_{t}\right|^{2} \mathrm{~d} t\right)^{\frac{1}{2}}\right\|_{\Psi_{1}}
\end{aligned}
$$

which along with (3.3.9) leads to the desired estimate for $\left\|\int_{0}^{T} X_{t} \sigma^{\top} \mathrm{d} W_{t}\right\|_{\Psi_{1 / 2}}$.

Similarly, by (3.6.8), there exists a constant $C$ satisfying for all $q \geq 2$ that

$$
\begin{aligned}
& \left\|\int_{0}^{T} \int_{\mathbb{R}_{0}^{p}} X_{t} \gamma(u) \tilde{N}(\mathrm{~d} t, \mathrm{~d} u)\right\|_{L^{q}} \\
& \leq C q^{2}\left\{\left(\mathbb{E}\left[\int_{0}^{T} \int_{\mathbb{R}_{0}^{p}}\left|X_{t} \gamma(u)\right|^{q} \nu(\mathrm{~d} u) \mathrm{d} t\right]\right)^{\frac{1}{q}}+\left(\mathbb{E}\left[\left(\int_{0}^{T} \int_{\mathbb{R}_{0}^{p}}\left|X_{t} \gamma(u)\right|^{2} \nu(\mathrm{~d} u) \mathrm{d} t\right)^{\frac{q}{2}}\right]\right)^{\frac{1}{q}}\right\} \\
& \leq C q^{2}\left\{\left(\int_{\mathbb{R}_{0}^{p}}|\gamma(u)|^{q} \nu(\mathrm{~d} u) \mathbb{E}\left[\int_{0}^{T}\left|X_{t}\right|^{q} \mathrm{~d} t\right]\right)^{\frac{1}{q}}+\left(\int_{\mathbb{R}_{0}^{p}}|\gamma(u)|^{2} \nu(\mathrm{~d} u)\right)^{\frac{1}{2}}\left(\mathbb{E}\left[\left(\int_{0}^{T}\left|X_{t}\right|^{2} \mathrm{~d} t\right)^{\frac{q}{2}}\right]\right)^{\frac{1}{q}}\right\} \\
& \leq C q^{2}\left\{\left(\int_{\mathbb{R}_{0}^{p}}|\gamma(u)|^{q} \nu(\mathrm{~d} u)\right)^{\frac{1}{q}}\left\|\left(\int_{0}^{T}\left|X_{t}\right|^{q} \mathrm{~d} t\right)^{\frac{1}{q}}\right\|_{L^{q}}+\left(\int_{\mathbb{R}_{0}^{p}}|\gamma(u)|^{2} \nu(\mathrm{~d} u)\right)^{\frac{1}{2}}\left\|\left(\int_{0}^{T}\left|X_{t}\right|^{2} \mathrm{~d} t\right)^{\frac{1}{2}}\right\|_{L^{q}}\right\}
\end{aligned}
$$

Hence by (H.4(2)) and (3.3.9), for all $q \geq 2$,

$$
\begin{aligned}
& q^{-(3+\vartheta)}\left\|\int_{0}^{T} \int_{\mathbb{R}_{0}^{p}} X_{t} \gamma(u) \tilde{N}(\mathrm{~d} t, \mathrm{~d} u)\right\|_{L^{q}} \\
& \leq C\left(\sup _{q \geq 2} q^{-\vartheta}\left(\int_{\mathbb{R}_{0}^{p}}|\gamma(u)|^{q} \nu(\mathrm{~d} u)\right)^{\frac{1}{q}}\right)\left\{q^{-1}\left\|\left(\int_{0}^{T}\left|X_{t}\right|^{q} \mathrm{~d} t\right)^{\frac{1}{q}}\right\|_{L^{q}}+q^{-1}\left\|\left(\int_{0}^{T}\left|X_{t}\right|^{2} \mathrm{~d} t\right)^{\frac{1}{2}}\right\|_{L^{q}}\right\} \\
& \leq C \gamma_{\max }\left(\left\|\left(\int_{0}^{T}\left|X_{t}\right|^{q} \mathrm{~d} t\right)^{\frac{1}{q}}\right\|_{\Psi_{1}}+\left\|\left(\int_{0}^{T}\left|X_{t}\right|^{2} \mathrm{~d} t\right)^{\frac{1}{2}}\right\|_{\Psi_{1}}\right)
\end{aligned}
$$

Therefore, taking the supremum over $q \geq 2$ in the above inequality leads to the desired estimate of $\left\|\int_{0}^{T} \int_{\mathbb{R}_{0}^{p}} X_{t} \gamma(u) \tilde{N}(\mathrm{~d} t, \mathrm{~d} u)\right\|_{\Psi_{1 /(3+\vartheta)}}$ from (3.3.9).

Proof of Lemma 3.3.7. Let us assume without loss of generality that $|\sigma|>0$ and $\tau:=$ $\left\|\left(\int_{0}^{T}\left|X_{t}\right|^{2} \mathrm{~d} t\right)^{1 / 2}\right\|_{\Psi_{2}}<\infty$, which implies that $\left\|\left(\int_{0}^{T} 2\left|X_{t} \sigma\right|^{2} \mathrm{~d} t\right)^{1 / 2}\right\|_{\Psi_{2}} \leq \sqrt{2}|\sigma| \tau$. Then, by the characterization of sub-Gaussian random variable in [156, Proposition 2.5.2(iii)], there exists $C \geq 0$ such that

$$
\mathbb{E}\left[\exp \left(2 \lambda^{2} \int_{0}^{T}\left|X_{t} \sigma\right|^{2} \mathrm{~d} t\right)\right] \leq \exp \left(2 C^{2} \lambda^{2}|\sigma|^{2} \tau^{2}\right)<\infty \quad \forall|\lambda| \leq \frac{1}{\sqrt{2} C|\sigma| \tau}
$$

Hence, it holds for all $|\lambda| \leq 1 /(\sqrt{2} C|\sigma| \tau)$ that the process $\left(M_{\lambda, t}\right)_{t \in[0, T]}$ defined by:

$$
M_{\lambda, t}:=\exp \left(\int_{0}^{t} 2 \lambda X_{s} \sigma^{\top} \mathrm{d} W_{s}-\frac{1}{2} \int_{0}^{t} 4 \lambda^{2}\left|X_{s} \sigma\right|^{2} \mathrm{~d} s\right) \quad \forall t \in[0, T]
$$

is a martingale, since Novikov's condition is satisfied, which implies that $\mathbb{E}\left[M_{\lambda, T}\right]=1$. Thus, for any given $|\lambda| \leq 1 /(\sqrt{2} C|\sigma| \tau)$, by the Cauchy-Schwarz inequality,

$$
\begin{aligned}
& \mathbb{E}\left[\exp \left(\lambda \int_{0}^{T} X_{t} \sigma^{\top} \mathrm{d} W_{t}\right)\right] \\
& =\mathbb{E}\left[\exp \left(\int_{0}^{T} \lambda X_{t} \sigma^{\top} \mathrm{d} W_{t}-\frac{(2 \lambda)^{2}}{4} \int_{0}^{T}\left|X_{t} \sigma\right|^{2} \mathrm{~d} t\right) \exp \left(\frac{(2 \lambda)^{2}}{4} \int_{0}^{T}\left|X_{t} \sigma\right|^{2} \mathrm{~d} t\right)\right] \\
& \leq \mathbb{E}\left[M_{\lambda, T}\right]^{1 / 2} \mathbb{E}\left[\exp \left(2 \lambda^{2} \int_{0}^{T}\left|X_{t} \sigma\right|^{2} \mathrm{~d} t\right)\right]^{1 / 2} \leq \exp \left(C^{2} \lambda^{2}|\sigma|^{2} \tau^{2}\right),
\end{aligned}
$$

which along with the fact that $\mathbb{E}\left[\int_{0}^{T} X_{t} \sigma^{\top} \mathrm{d} W_{t}\right]=0$ and the characterization of sub-exponential random variable [156, Proposition 2.7.1(v)] yields that $\left\|\int_{0}^{T} X_{t} \sigma^{\top} \mathrm{d} W_{t}\right\|_{\Psi_{1}} \leq C|\sigma| \tau$.

Proof of Lemma 3.3.8. Throughout this proof, let $x_{0} \in \mathbb{R}^{n}$ and $\theta \in \mathbb{R}^{n \times(n+k)}$ be given constants satisfying $|\theta| \leq K$, and let $C$ be a generic constant depend on $K, T$ and the constants in (H.4), but independent of $x_{0}$ and $\theta$.

By (3.3.4), we see that the process $X^{x_{0}, \theta}$ satisfies the SDE:

$$
\mathrm{d} X_{t}=b^{\theta}\left(t, X_{t}\right) \mathrm{d} t+\sigma \mathrm{d} W_{t}+\int_{\mathbb{R}_{0}^{p}} \gamma(u) \tilde{N}(\mathrm{~d} t, \mathrm{~d} u), \quad t \in[0, T], \quad X_{0}=x_{0}
$$

where $b^{\theta}(t, x)=A^{\star} x+B^{\star} \psi^{\theta}(t, x)$ for all $(t, x) \in[0, T] \times \mathbb{R}^{n}$. The definition of the feedback control $\psi^{\theta}$, (H.4(1)) and Theorem 3.2.5 show that there exists $C \geq 0$ such that $\left|\psi^{\theta}(t, 0)\right| \leq C$ and $\left|\psi^{\theta}(t, x)-\psi^{\theta}\left(t, x^{\prime}\right)\right| \leq C\left|x-x^{\prime}\right|$ for all $t \in[0, T], x, x^{\prime} \in \mathbb{R}^{n}$, which implies the same properties for the function $b^{\theta}$. Then, by Lemma 3.6.3, for every Lipschitz continuous function $\mathfrak{f}:\left(\mathbb{D}\left([0, T] ; \mathbb{R}^{n}\right), d_{\infty}\right) \rightarrow \mathbb{R}, \mathbb{E}\left[\exp \left(\lambda\left(\mathfrak{f}\left(X^{x_{0}, \theta}\right)-\mathbb{E}\left[\mathfrak{f}\left(X^{x_{0}, \theta}\right)\right]\right)\right)\right] \leq \exp \left(C \eta\left(C \lambda\|\mathfrak{f}\|_{\text {Lip }}\right)\right)$ for all $\lambda>0$, with the function $\eta:[0, \infty] \rightarrow[0, \infty]$ defined by:

$$
\begin{equation*}
\eta(\lambda):=\int_{\mathbb{R}_{0}^{p}}\left(e^{\lambda \gamma(u)}-\lambda \gamma(u)-1\right) \nu(\mathrm{d} u)+\frac{\sigma^{2}}{2} \lambda^{2} \quad \forall \lambda>0 . \tag{3.6.14}
\end{equation*}
$$

By (H.4(2)) and Stirling's approximation $q!\geq(q / e)^{q}$ for all $q \geq 2$, we have for each $\lambda \in$ $\left[0,1 /\left(2 \gamma_{\max } e\right)\right)$,

$$
\begin{aligned}
& \int_{\mathbb{R}_{0}^{p}}\left(e^{\lambda \gamma(u)}-\lambda \gamma(u)-1\right) \nu(\mathrm{d} u) \\
& =\int_{\mathbb{R}_{0}^{p}} \sum_{q=2}^{\infty} \frac{|\lambda \gamma(u)|^{q}}{q!} \nu(\mathrm{d} u)=\sum_{q=2}^{\infty} \frac{\lambda^{q}}{q!} \int_{\mathbb{R}_{0}^{p}}|\gamma(u)|^{q} \nu(\mathrm{~d} u) \leq \sum_{q=2}^{\infty} \frac{\lambda^{q}}{q!} \gamma_{\max }^{q} q^{\vartheta q} \\
& \leq \sum_{q=2}^{\infty} \frac{\left(\lambda \gamma_{\max } e\right)^{q}}{q^{(1-\vartheta) q}} \leq \frac{\left(\lambda \gamma_{\max } e\right)^{2}}{1-\lambda \gamma_{\max } e} \leq 2\left(\lambda \gamma_{\max } e\right)^{2},
\end{aligned}
$$

which implies for all $0 \leq \lambda \leq 1 / C$ and $\mathfrak{f}:\left(\mathbb{D}\left([0, T] ; \mathbb{R}^{n}\right), d_{\infty}\right) \rightarrow \mathbb{R}$ satisfying $\|\mathfrak{f}\|_{\text {Lip }} \leq$ 1 that $\mathbb{E}\left[\exp \left(\lambda\left(\mathfrak{f}\left(X^{x_{0}, \theta}\right)-\mathbb{E}\left[\mathfrak{f}\left(X^{x_{0}, \theta}\right)\right]\right)\right)\right] \leq \exp \left(C^{2} \lambda^{2}\right)$. Replacing $\mathfrak{f}$ with $-\mathfrak{f}$ shows that the same estimate holds for all for all $|\lambda| \leq 1 / C$, which, along with the characterization of sub-exponential random variable in [156, Proposition 2.7.1(v)], leads to $\| \mathfrak{f}\left(X^{x_{0}, \theta}\right)$ $\mathbb{E}\left[\mathfrak{f}\left(X^{x_{0}, \theta}\right)\right] \|_{\Psi_{1}} \leq C$ for some constant $C$, uniformly with respect to $x_{0} \in \mathbb{R}^{n},|\theta| \leq K$ and $\mathfrak{f}:\left(\mathbb{D}\left([0, T] ; \mathbb{R}^{n}\right), d_{\infty}\right) \rightarrow \mathbb{R}$ satisfying $\|\mathfrak{f}\|_{\text {Lip }} \leq 1$.

Since $\|\cdot\|_{\Psi_{1}}$ is a norm and $\left\|\mathbb{E}\left[\mathfrak{f}\left(X^{x}\right)\right]\right\|_{\Psi_{1}} \leq\left|\mathbb{E}\left[\mathfrak{f}\left(X^{x}\right)\right]\right| / \ln 2,\left\|\mathfrak{f}\left(X^{x_{0}, \theta}\right)\right\|_{\Psi_{1}} \leq C(1+$ $\left.\left|\mathbb{E}\left[\mathfrak{f}\left(X^{x_{0}, \theta}\right)\right]\right|\right)$ for all $\mathfrak{f}$ with $\|\mathfrak{f}\|_{\text {Lip }} \leq 1$. The estimate for a general Lipschitz continuous function $\mathfrak{f}$ follows by considering $\mathfrak{f} /\|\mathfrak{f}\|_{\text {Lip }}$ and by using the fact that $\|\cdot\|_{\Psi_{1}}$ is a norm.

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[^0]:    ${ }^{1}$ The numerical experiments on the application of GMF-V-Q to the motivating Ad auction problem can be found in the conference version of our paper [76].

[^1]:    ${ }^{2}$ Here the indicator function $\mathbf{I}_{s_{t}^{j}=s}=1$ if $s_{t}^{j}=s$ and 0 otherwise.

[^2]:    ${ }^{3}$ Note that MF-Q is designed for global states and coupled local actions, while in our equilibrium price example we have coupled local (private) states and decoupled local actions. To suit this setting, we adapt MF-Q by replacing the mean-field action term with the mean-field state term.
    ${ }^{4}$ For the sake of brevity, we omit the algorithm frame for GMF-PW-TRPO.

[^3]:    ${ }^{5}$ Lemma 1.5.2 indicates that the learning rate should be inversely proportional to the current visitation number of a given state-action pair, we observe that constant learning rate works well in practice which is easier to implement.

[^4]:    ${ }^{6}$ Here we actually only need $c^{\prime}=\Omega\left(\frac{\log (1 / \epsilon)}{\phi_{N}}\right)$ and $c=O\left(\frac{\log (1 / \epsilon)}{\phi_{N}}\right)$, and the corresponding result will differ only in some absolute constants.

[^5]:    ${ }^{7}$ Here again we actually only need $c^{\prime}=\Omega\left(\frac{\log (1 / \epsilon)}{\phi_{N}}\right)$ and $c=O\left(\frac{\log (1 / \epsilon)}{\phi_{N}}\right)$, and the corresponding result will differ only in some absolute constants.

[^6]:    ${ }^{1}$ We say a function $\mathcal{R}: \mathbb{R}^{k} \rightarrow \mathbb{R} \cup\{\infty\}$ is proper if it has a nonempty effective domain $\operatorname{dom} \mathcal{R}:=\{a \in$ $\left.\mathbb{R}^{k} \mid \mathcal{R}(a)<\infty\right\}$.

