# Toward Understanding State Representation Learning in MuZero: A Case Study in Linear Quadratic Gaussian Control 

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#### Abstract

We study the problem of representation learning for control from partial and potentially high-dimensional observations. We approach this problem via direct latent model learning, where one directly learns a dynamical model in some latent state space by predicting costs. In particular, we establish finite-sample guarantees of finding a near-optimal representation function and a near-optimal controller using the directly learned latent model for infinite-horizon time-invariant Linear Quadratic Gaussian (LQG) control. A part of our approach to latent model learning closely resembles MuZero, a recent breakthrough in empirical reinforcement learning, in that it learns latent dynamics implicitly by predicting cumulative costs. A key technical contribution of this work is to prove persistency of excitation for a new stochastic process that arises from the analysis of quadratic regression in our approach.


## 1 Introduction

Control with a learned latent model is state-of-the-art in several reinforcement learning (RL) benchmarks, including board games, Atari games, and visuomotor control (Schrittwieser et al., 2020; Ye et al., 2021; Hafner et al., 2023). To better understand this modern machinery, we introduce it to a classical optimal control problem, namely Linear Quadratic Gaussian (LQG) control, and study its theoretical, finite-sample performance. Essential to this approach is the learning of two components: a state representation function that maps an observed history to some latent state, and a latent model that predicts the transition and cost in the latent state space. The latent model is usually a Markov decision process, using which we obtain a policy in the latent space or execute online planning.

What is the correct objective to optimize for learning a latent model? One popular choice is to learn a function that reconstructs the observation from the latent state (Hafner et al., 2019a,b,

2020, 2023). A latent model learned this way is agnostic to control tasks and retains all the information about the environment. This class of approaches can achieve satisfactory performance, but are prone to background distraction and control-irrelevant information (Fu et al., 2021). The second class of methods learn an inverse model that infers actions from latent states at different time steps (Pathak et al., 2017; Lamb et al., 2022). A latent model learned with this methodology is also task agnostic but extracts control-relevant information. In contrast, task-relevant representations can be learned by predicting costs in the control task (Oh et al., 2017; Zhang et al., 2020; Schrittwieser et al., 2020). The concept that a good latent state should be able to predict costs is intuitive, and the costs are directly relevant to optimal control. Hence, (Tian et al., 2022) refers to this class of methods as "direct latent model learning", which is the focus of this work.

The direct latent model learning method of particular interest to us is that of MuZero (Schrittwieser et al., 2020). Announced by DeepMind in 2019, MuZero extends the line of works including AlphaGo (Silver et al., 2016), AlphaGo Zero (Silver et al., 2017), and AlphaZero (Silver et al., 2018) by not requiring knowledge of the game rules. MuZero matches the superhuman performance of AlphaZero in Go, shogi and chess, while outperforming model-free RL algorithms in Atari games. MuZero builds upon the powerful planning procedure of Monte Carlo Tree Search, with the major innovation being learning a latent model. The latent model replaces the rule-based simulator during planning, and avoids the burdensome planning in pixel space for Atari games.

MuZero is a milestone algorithm in representation learning for control. Intuitively, the algorithm design makes sense, but its complexity has so far inhibited a formal theoretical study. On the other hand, statistical learning theory for linear dynamical systems and control has evolved rapidly in recent years (Tsiamis et al., 2022); for partially observable linear dynamical systems, much of the work relies on learning Markov parameters, lacking a direct connection to the empirical methods used in practice for possibly nonlinear systems. In this work, we aim to bridge the two areas by studying provable MuZero-style latent model learning in LQG control.

The latent model learning of MuZero features three ingredients: 1) stacking frames, i.e., observations, as input to the representation function; 2) predicting costs, "optimal" values, and "optimal" actions from latent states; and 3) implicit learning of latent dynamics by predicting these quantities from latent states at future time steps. These are the defining characteristics of the MuZero-style algorithm that we shall consider. In MuZero, the "optimal" values and actions are found by the powerful online planning procedure. In this work, we simplify the setup by considering data collected using random actions, which are known to suffice for identifying a partially observable linear dynamical system (Oymak and Ozay, 2019). In this setup, the values become those associated with this trivial policy and we do not predict actions since they are random noises anyway. Note that although our study of the above ingredients is directly motivated by MuZero, previous empirical works have also explored them. For example, frame stacking has been a widely used technique to handle partial observability (Mnih et al., 2013, 2015); predicting values for learning a latent model has been studied in (Oh et al., 2017), which also learns the latent state transition implicitly.

Closely related to our work, (Tian et al., 2022) also considers provable direct latent model learning in LQG, but for the finite-horizon time-varying setting. Our work builds upon it
and complements it in two ways: 1) we extend their algorithm to the time-invariant setting with a stationary representation function and latent model, which is closer to what has been deployed in practice; 2) we present and analyze a new, MuZero-style latent model learning algorithm. Both 1) and 2) introduce new technical challenges to be addressed. We summarize our contributions as follows.

- We show that two direct latent model learning methods provably solve infinite-horizon time-invariant LQG control by establishing finite-sample guarantees. Both methods only need a single trajectory; one resembles the method in (Tian et al., 2022), and the other resembles MuZero.
- By analyzing the MuZero-style algorithm, we notice the potential issue of coordinate misalignment; that is, costs can be invariant to certain transformations of the latent states, and implicit dynamics learning by predicting one-step transition may not recover the latent state coordinates consistently. This insight suggests the need of predicting multi-step transition or other coordinate alignment procedures in implicit dynamics learning.
- Technically, we overcome the difficulty of having correlated data samples in a single trajectory for latent model learning, by proving a new result about the persistency of excitation for a stochastic process that arises from the analysis of the quadratic regression subroutine in both of our methods.

Notation. Random vectors are denoted by lowercase letters; sometimes they also denote their realized values. Uppercase letters denote matrices, some of which can be random. Let $a \wedge b$ denote the minimum between scalars $a$ and $b .1$ denotes either the scalar one or a vector consisting of all ones; $I$ denotes an identity matrix. The dimension, when emphasized, is specified in subscripts, e.g., $1_{d}, I_{d}$. Given vector $v \in \mathbb{R}^{d}$, let $\|v\|$ denote its $\ell_{2}$ norm and $\|v\|_{P}:=\left(v^{\top} P v\right)^{1 / 2}$ for positive semidefinite $P \in \mathbb{R}^{d \times d}$. Given symmetric matrices $P$ and $Q, P \succ Q$ or $Q \prec P$ means $P-Q$ is positive definite, and $P \succcurlyeq Q$ or $Q \preccurlyeq P$ means $P-Q$ is positive semidefinite. Semicolon ";" denotes stacking vectors or matrices vertically. For a collection of $d$-dimensional vectors $\left(v_{t}\right)_{t=i}^{j}$, let $v_{i: j}:=\left[v_{i} ; v_{i+1} ; \ldots ; v_{j}\right] \in \mathbb{R}^{d(j-i+1)}$ denote the concatenation along the column. For random variable $x$, let $\|x\|_{\psi_{\theta}}$ denote its $\theta$-sub-Weibull norm for $\theta>0$, a special case of Orlicz norms (Zhang and Wei, 2022), with $\theta=1,2$ corresponding to subexponential and sub-Gaussian norms. For matrix $A$, let $\sigma_{\min }(A),\|A\|_{2},\|A\|_{F}$, and $\|A\|_{*}$ denote its minimum eigenvalue, minimum singular value, operator norm (induced by vector $\ell_{2}$ norms), Frobenius norm, and nuclear norm, respectively. $\langle\cdot, \cdot\rangle_{F}$ denotes the Frobenius inner product between matrices. For square matrix $A$, let $\lambda_{\text {min }}(A)$ be its minimum eigenvalue and $\rho(A)$ be its spectral radius. Define $\alpha(A):=\sup _{k \geq 0}\left\|A^{k}\right\|_{2} \rho(A)^{-k}$. Let $\operatorname{svec}(\cdot)$ denote the operator of flattening a symmetric matrix by stacking its columns; it does not repeat the off-diagonal elements, but scales them by $\sqrt{2}$ (Schacke, 2004). We adopt the standard use of $\mathcal{O}(\cdot), \Omega(\cdot), \Theta(\cdot)$, where the hidden constants are dimension-free but may depend on system parameters.

## 2 Problem setup

A partially observable linear time-invariant (LTI) dynamical system is described by

$$
\begin{equation*}
x_{t+1}=A^{*} x_{t}+B^{*} u_{t}+w_{t}, \quad y_{t}=C^{*} x_{t}+v_{t} \tag{2.1}
\end{equation*}
$$

with state $x_{t} \in \mathbb{R}^{d_{x}}$, observation $y_{t} \in \mathbb{R}^{d_{y}}$, and control $u_{t} \in \mathbb{R}^{d_{u}}$ for all $t \geq 0$. Process noises $\left(w_{t}\right)_{t \geq 0}$ and observation noises $\left(v_{t}\right)_{t \geq 0}$ are i.i.d. zero-mean Gaussian random vectors with covariance matrices $\Sigma_{w}$ and $\Sigma_{v}$, respectively, and the two sequences are mutually independent. Let initial state $x_{0}$ be sampled from $\mathcal{N}\left(0, \Sigma_{0}\right)$. The quadratic cost function is given by

$$
\begin{equation*}
c(x, u)=\|x\|_{Q^{*}}^{2}+\|u\|_{R^{*}}^{2} \tag{2.2}
\end{equation*}
$$

where $Q^{*} \succcurlyeq 0$ and $R^{*} \succ 0$.
A policy/controller $\pi$ determines an action/control input $u_{t}$ at time step $t$ based on the history $\left[y_{0: t} ; u_{0:(t-1)}\right]$ up to this time step. For $t \geq 0, c_{t}:=c\left(x_{t}, u_{t}\right)$ denotes the cost at time step $t$. Given a policy $\pi$, let

$$
\begin{equation*}
J^{\pi}:=\limsup _{T \rightarrow \infty} \mathbb{E}\left[\frac{1}{T} \sum_{t=0}^{T-1} c_{t}\right] \tag{2.3}
\end{equation*}
$$

denote the time-average expected cost. The objective of LQG control is to find a policy $\pi$ such that $J^{\pi}$ is minimized.

In the fully observable setting, known as the linear quadratic regulator (LQR), $y_{t}=x_{t}$. A linear controller with feedback gain $K \in \mathbb{R}^{d_{u} \times d_{x}}$ determines action $u_{t}=K x_{t}$ at time step $t$. Let $J^{K}\left(A^{*}, B^{*}, Q^{*}, R^{*}\right)$ denote the time-average expected cost (2.3) in the LQR problem $\left(A^{*}, B^{*}, Q^{*}, R^{*}\right)$ under feedback gain $K$ and define $J^{*}\left(A^{*}, B^{*}, Q^{*}, R^{*}\right):=\min _{K} J^{K}\left(A^{*}, B^{*}, Q^{*}, R^{*}\right)$.

We make the following standard assumptions.
Assumption 1. System dynamics (2.1) and cost (2.2) satisfy:

1. The system is stable, that is, $\rho\left(A^{*}\right)<1$.
2. $\left(A^{*}, B^{*}\right)$ is $v$-controllable for some $v>0$, that is, the controllability matrix

$$
\Phi_{c}\left(A^{*}, B^{*}\right):=\left[B^{*}, A^{*} B^{*}, \ldots,\left(A^{*}\right)^{d_{x}-1} B^{*}\right]
$$

has rank $d_{x}$ and $\sigma_{\min }\left(\Phi_{c}\left(A^{*}, B^{*}\right)\right) \geq v$.
3. $\left(A^{*}, C^{*}\right)$ is $\omega$-observable for some $\omega>0$, that is, the observability matrix

$$
\Phi_{o}\left(A^{*}, C^{*}\right):=\left[C^{*} ; C^{*} A^{*} ; \ldots ; C^{*}\left(A^{*}\right)^{d_{x}-1}\right]
$$

has rank $d_{x}$ and $\sigma_{\min }\left(\Phi_{o}\left(A^{*}, C^{*}\right)\right) \geq \omega$.
4. $\left(A^{*}, \Sigma_{w}^{1 / 2}\right)$ is $\eta$-controllable for some $\eta>0$.
5. $\left(A^{*},\left(Q^{*}\right)^{1 / 2}\right)$ is $\mu$-observable for some $\mu>0$.
6. $\Sigma_{v} \succcurlyeq \sigma_{v}^{2} I$ for some $\sigma_{v}>0$; this can always be achieved by inserting Gaussian noises with full-rank covariance matrices to the observations.
7. $R^{*} \succcurlyeq r^{2} I$ for some $r>0$.
8. The operator norms of $A^{*}, B^{*}, C^{*}, Q^{*}, R^{*}, \Sigma_{w}, \Sigma_{v}, \Sigma_{0}$ and $\alpha\left(A^{*}\right), \alpha\left(\bar{A}^{*}\right)$ are $\mathcal{O}(1)$; the singular value lower bounds $v, \omega, v, \eta, \sigma_{v}$, r and spectral radii $\rho\left(A^{*}\right), \rho\left(\bar{A}^{*}\right)$ are $\Omega(1)$, where $\bar{A}^{*}$ is defined in §2.1.

If the system parameters $\left(A^{*}, B^{*}, C^{*}, Q^{*}, R^{*}, \Sigma_{w}, \Sigma_{v}\right)$ are known, the optimal policy is obtained by combining the Kalman filter

$$
\begin{equation*}
z_{t+1}^{*}=A^{*} z_{t}^{*}+B^{*} u_{t}+L^{*}\left(y_{t+1}-C^{*}\left(A^{*} z_{t}^{*}+B^{*} u_{t}\right)\right) \tag{2.4}
\end{equation*}
$$

with the optimal feedback gain $K^{*}$ of the linear quadratic regulator (LQR) such that $u_{t}=K^{*} z_{t}^{*}$, where $L^{*}$ is the Kalman gain, and at the initial time step, we can set, e.g., $z_{0}^{*}=L^{*} y_{0}$. This fact is known as the separation principle, and the Kalman gain and optimal feedback gain are given by

$$
\begin{align*}
L^{*} & =S^{*}\left(C^{*}\right)^{\top}\left(C^{*} S^{*}\left(C^{*}\right)^{\top}+\Sigma_{v}\right)^{-1},  \tag{2.5}\\
K^{*} & =-\left(\left(B^{*}\right)^{\top} P^{*} B^{*}+R\right)^{-1}\left(B^{*}\right)^{\top} P^{*} A^{*}, \tag{2.6}
\end{align*}
$$

where $S^{*}$ and $P^{*}$ are determined by their respective discrete-time algebraic Riccati equations (DAREs):

$$
\begin{align*}
S^{*} & =A^{*}\left(S^{*}-S^{*}\left(C^{*}\right)^{\top}\left(C^{*} S^{*}\left(C^{*}\right)^{\top}+\Sigma_{v}\right)^{-1} C^{*} S^{*}\right)\left(A^{*}\right)^{\top}+\Sigma_{w v}  \tag{2.7}\\
P^{*} & =\left(A^{*}\right)^{\top}\left(P^{*}-P^{*} B^{*}\left(\left(B^{*}\right)^{\top} P^{*} B^{*}+R^{*}\right)^{-1}\left(B^{*}\right)^{\top} P^{*}\right) A^{*}+Q^{*} . \tag{2.8}
\end{align*}
$$

Assumptions 1.2 to 1.7 guarantee the existence and uniqueness of positive definite solutions $S^{*}$ and $P^{*}$; Assumption 1.8 further guarantees that their operator norms are $\mathcal{O}(1)$ and minimum singular values are $\Omega(1)$. The assumption on $\alpha\left(A^{*}\right), \alpha\left(\bar{A}^{*}\right), \rho\left(A^{*}\right), \rho\left(\bar{A}^{*}\right)$ provides guarantees for state estimation from a finite history and is mild (Mania et al., 2019; Oymak and Ozay, 2019). If $\rho\left(A^{*}\right)$ or $\rho\left(\bar{A}^{*}\right)$ equals zero, $\left(A^{*}\right)^{d_{x}}$ or $\left(\bar{A}^{*}\right)^{d_{x}}$ is a zero matrix by the Cayley-Hamilton theorem, so using history length $H \geq d_{x}$ completely eliminates the truncation errors.

We consider the data-driven control setting, where the LQG model ( $\left.A^{*}, B^{*}, C^{*}, Q^{*}, \Sigma_{w}, \Sigma_{v}\right)$ is unknown. For simplicity, we assume $R^{*}$ is known, though our approaches can be readily extended to the case where it is unknown by learning it from predicting costs.

### 2.1 Latent model of LQG

The stationary Kalman filter (2.4) asymptotically produces the optimal state estimation in the sense of minimum mean squared errors. With a finite horizon, however, the optimal state estimator is time-varying, given by

$$
\begin{equation*}
z_{t+1}^{*}=A^{*} z_{t}^{*}+B^{*} u_{t}+L_{t+1}^{*}\left(y_{t+1}-C^{*}\left(A^{*} z_{t}^{*}+B^{*} u_{t}\right)\right), \tag{2.9}
\end{equation*}
$$

where $L_{t}^{*}$ is the time-varying Kalman gain, converging to $L^{*}$ as $t \rightarrow \infty$. This convergence is equivalent to that of error covariance matrix $\mathbb{E}\left[\left(x_{t}-z_{t}^{*}\right)\left(x_{t}-z_{t}^{*}\right)^{\top}\right]$, which happens exponentially fast (Komaroff, 1994). Hence, for simplicity, we assume this error covariance matrix is
stationary at the initial time step by the choice of $z_{0}^{*}$ so that $L_{t}^{*}=L^{*}$ for $t \geq 1$; this assumption is common in the literature (Lale et al., 2020, 2021; Jadbabaie et al., 2021). The innovation term $i_{t+1}:=y_{t+1}-C^{*}\left(A^{*} z_{t}^{*}+B^{*} u_{t}\right)$ is independent of the history $\left(y_{0}, u_{0}, y_{1}, u_{1}, \ldots, y_{t+1}\right)$ and $\left(i_{t}\right)_{t \geq 1}$ are mutually independent. The following proposition taken from (Tian et al., 2022, Proposition 1) represents the system in terms of the state estimates obtained by the Kalman filter, which we refer to as the latent model.

Proposition 1. Let $\left(z_{t}^{*}\right)_{t \geq 1}$ be state estimates given by the Kalman filter. Then, for $t \geq 0$,

$$
z_{t+1}^{*}=A^{*} z_{t}^{*}+B^{*} u_{t}+L^{*} i_{t+1},
$$

where $L^{*} i_{t+1}$ is independent of $z_{t}^{*}$ and $u_{t}$, i.e., the state estimates follow the same linear dynamics with noises $L^{*} i_{t+1}$. The cost at step $t$ can be reformulated as functions of the state estimates by

$$
c_{t}=\left\|z_{t}^{*}\right\|_{Q^{*}}^{2}+\left\|u_{t}\right\|_{R^{*}}^{2}+b^{*}+\gamma_{t}+\eta_{t},
$$

where $b^{*}>0$, and $\gamma_{t}=\left\|z_{t}^{*}-x_{t}\right\|_{Q^{*}}^{2}-b^{*}, \eta_{t}=\left\langle z_{t}^{*}, x_{t}-z_{t}^{*}\right\rangle_{Q^{*}}$ are both zero-mean subexponential random variables. Moreover, $b^{*}=\mathcal{O}(1)$ and $\left\|\gamma_{t}\right\|_{\psi_{1}}=\mathcal{O}\left(d_{x}^{1 / 2}\right)$; if control $u_{t} \sim \mathcal{N}\left(0, \sigma_{u}^{2} I\right)$ for $t \geq 0$, then we have $\left\|\eta_{t}\right\|_{\psi_{1}}=\mathcal{O}\left(d_{x}^{1 / 2}\right)$.

Proposition 1 shows that the dynamics of the state estimates computed by the time-varying Kalman filter is the same as the original system up to noises; the costs are also the same, up to constants and noises. Hence, a latent model can be parameterized by $\left(A, B, Q, R^{*}\right)$, with the constant $b^{*}$ and noises neglected due to their irrelevance to planning. A stationary latent policy is a linear controller $u_{t}=K z_{t}$ on latent state $z_{t}$, parameterized by feedback gain $K \in \mathbb{R}^{d_{u} \times d_{x}}$.

The latent model enables us to find a good latent policy. To learn such a latent model and to deploy a latent policy in the original partially observable system, we need a representation function. Let $\bar{A}^{*}:=\left(I-L^{*} C^{*}\right) A^{*}$ and $\bar{B}^{*}:=\left(I-L^{*} C^{*}\right) B^{*}$. Then, the Kalman filter can be written as $z_{t+1}^{*}=\bar{A}^{*} z_{t}^{*}+\bar{B}^{*} u_{t}+L^{*} y_{t+1}$. For $t \geq 0$, unrolling the recursion gives

$$
\begin{aligned}
z_{t}^{*} & =\bar{A}^{*}\left(\bar{A}^{*} z_{t-2}^{*}+\bar{B}^{*} u_{t-2}+L^{*} y_{t-1}\right)+\bar{B}^{*} u_{t-1}+L^{*} y_{t} \\
& =\left[\left(\bar{A}^{*}\right)^{t-1} L^{*}, \ldots, L^{*}\right] y_{1: t}+\left[\left(\bar{A}^{*}\right)^{t-1} \bar{B}^{*}, \ldots, \bar{B}^{*}\right] u_{0:(t-1)}+\left(\bar{A}^{*}\right)^{t} z_{0}^{*} \\
& =: M_{t}^{*}\left[y_{1: t} ; u_{0:(t-1)} ; z_{0}^{*}\right],
\end{aligned}
$$

where $M_{t}^{*} \in \mathbb{R}^{d_{x} \times\left(t d_{y}+t d_{u}+d_{x}\right)}$. This means the representation function can be parameterized as linear mappings for full histories (with $y_{0}$ replaced by $z_{0}^{*}$ ). Despite the simplicity, the input dimension of the function grows linearly in time, making it intractable to estimate the state using the full history for large $t$; nor is it necessary, since the impact of old data decreases exponentially. Under Assumption 1, $\rho\left(\bar{A}^{*}\right)<1$ (Bertsekas, 2012, Appendix E.4). With an $H$-step truncated history, the state estimate can be written as

$$
\begin{aligned}
z_{t}^{*} & =\left[\left(\bar{A}^{*}\right)^{H-1} L^{*}, \ldots, L^{*}\right] y_{(t-H+1): t}+\left[\left(\bar{A}^{*}\right)^{H-1} \bar{B}^{*}, \ldots, \bar{B}^{*}\right] u_{(t-H):(t-1)}+\delta_{t} \\
& =: M^{*}\left[y_{(t-H+1)::} ; u_{(t-H):(t-1)}\right]+\delta_{t}
\end{aligned}
$$

where $\delta_{t}=\left(\bar{A}^{*}\right)^{H} z_{t-H}^{*}$, whose impact decays exponentially in $H$ and can be neglected for sufficiently large $H$, since $z_{t-H}^{*}$ converges to a stationary distribution and its norm is bounded
with high probability. Hence, the representation function that we aim to recover is $M^{*} \in$ $\mathbb{R}^{d_{x} \times H\left(d_{y}+d_{u}\right)}$, which takes as input the $H$-step history $h_{t}=\left[y_{(t-H+1): t} ; u_{(t-H):(t-1)}\right]$. Henceforth, we let $d_{h}:=H\left(d_{y}+d_{u}\right)$. Then, a representation function is parameterized by matrix $M \in \mathbb{R}^{d_{x} \times d_{h}}$.

Overall, a policy is a combination of a representation function $M$ and a feedback gain $K$ in the latent model, denoted by $\pi=(M, K)$. Learning to solve LQG control in this framework can thus be achieved by: 1) learning representation function $M$; 2) extracting latent model $\left(A, B, Q, R^{*}\right)$; and 3 ) finding the optimal $K$ by planning in the latent model. Next, we introduce our approach following this pipeline.

## 3 Method

In practice, latent model learning methods collect trajectories by interacting with the system online using some policy; the trajectories are used to improve the learned latent model, which in turn improves the policy. In LQG control, it is known that the simple setup allows us to learn a good latent model from a single trajectory, collected using zero-mean Gaussian inputs; see e.g., (Oymak and Ozay, 2019). This is also how we assume the data are collected. We note that our results also apply to data from multiple independent trajectories using the same zero-mean Gaussian inputs.

In direct latent model learning, state representations are learned by predicting costs. To learn the transition function in the latent model, two approaches are explored in the literature. The first approach explicitly minimizes transition prediction errors (Subramanian et al., 2020; Hafner et al., 2019a). Algorithmically, the overall loss is a combination of cost and transition prediction errors. The second approach, which MuZero takes, learns transition implicitly, by minimizing cost prediction errors at future states generated from the transition function (Oh et al., 2017; Schrittwieser et al., 2020). Algorithmically, the overall loss aggregates the cost prediction errors across multiple time steps. In both approaches, the coupling of different terms in the loss makes finite-sample analysis difficult. As observed in (Tian et al., 2022), the structure of LQG allows us to learn the representation function independently of learning the transition function. This allows us to formulate both approaches under the same direct latent model learning framework (Algorithm 1).

Algorithm 1 consists of three main steps. Lines 3 to 5 correspond to cost-driven representation learning. Lines 6 to 8 correspond to latent model learning, where the system dynamics can be identified either explicitly, by ordinary least squares (SysId), or implicitly, by future cost prediction (CoSysId, Algorithm 2). Line 8 corresponds to latent policy optimization; in LQG this amounts to solving a DARE. Below we elaborate on cost-driven representation learning, SysId, and CoSysId in order.

### 3.1 Cost-driven representation learning

The procedure of cost-driven representation learning is almost identical to that in (Tian et al., 2022). The main idea is to perform quadratic regression (3.2) to the $d_{x}$-step cumulative costs;

```
Algorithm 1 Direct latent model learning for LQG control
    Input: length \(T\), history length \(H\), noise magnitude \(\sigma_{u}\)
    Collect a trajectories of length \(T+H\) using \(u_{t} \sim \mathcal{N}\left(0, \sigma_{u}^{2} I\right)\), for \(t \geq 0\), to obtain \(\mathcal{D}_{\text {raw }}\) :
```

$$
\begin{equation*}
\left(y_{0}, u_{0}, c_{0}, y_{1}, u_{1}, c_{1}, \ldots, y_{T+H}\right) \tag{3.1}
\end{equation*}
$$

3: Estimate the state representation function and cost constants by solving

$$
\begin{equation*}
\hat{N}, \hat{b}_{0} \in \underset{N=N^{\top}, b_{0}}{\operatorname{argmin}} \sum_{t=H}^{T+H-1}\left(\left\|h_{t}\right\|_{N}^{2}+b_{0}-\bar{c}_{t}\right)^{2}, \tag{3.2}
\end{equation*}
$$

where $\bar{c}_{t}:=\sum_{\tau=t}^{t+d_{x}-1}\left(c_{\tau}-\left\|u_{\tau}\right\|_{R^{*}}^{2}\right)$
Find $\hat{M} \in \operatorname{argmin}_{M \in \mathbb{R}^{d_{x} \times H\left(d_{y}+d_{u}\right)}\left\|M^{\top} M-\hat{N}\right\|_{F}, ~}^{\text {and }}$
5: Compute $\hat{z}_{t}=\hat{M}\left[y_{(t-H+1):} ; u_{(t-H):(t-1)}\right]$ for all $t \geq H$, so that the data are converted to $\mathcal{D}_{\text {state }}$ :

$$
\left(\hat{z}_{H}, u_{H}, c_{H}, \ldots, \hat{z}_{T+H-1}, u_{T+H-1}, c_{T+H-1}, \hat{z}_{T+H}\right)
$$

6: Run SysId (3.3) or CoSysId (Algorithm 2) to obtain dynamics matrices $(\hat{A}, \hat{B})$
7: Estimate the cost function by solving

$$
\widetilde{Q}, \hat{b} \in \underset{Q=Q^{\top}, b}{\operatorname{argmin}} \sum_{t=H}^{T+H-1}\left(\left\|\hat{z}_{t}\right\|_{Q}^{2}+b-c_{t}\right)^{2},
$$

8: Truncate negative eigenvalues of $\widetilde{Q}$ to zero to obtain $\hat{Q} \succcurlyeq 0$
9: Find feedback gain $\hat{K}$ from $\left(\hat{A}, \hat{B}, \hat{Q}, R^{*}\right)$ by DARE (2.8) and (2.6)
Return: policy $\hat{\pi}=(\hat{M}, \hat{K})$
these correspond to the value prediction in MuZero. By the $\mu$-observability of ( $A^{*},\left(Q^{*}\right)^{1 / 2}$ ) (Assumption 1.5), the cost observability Gramian

$$
\bar{Q}^{*}:=\sum_{t=0}^{d_{x}-1}\left(\left(A^{*}\right)^{t}\right)^{\top} Q^{*}\left(A^{*}\right)^{t} \succcurlyeq \mu^{2} I .
$$

Under zero control and zero noise, starting from $x$, the $d_{x}$-step cumulative cost is precisely $\|x\|_{\bar{Q}^{*}}^{2}$. Hence, $\hat{N}$ estimates $N^{*}=\left(M^{*}\right)^{\top} \bar{Q}^{*} M^{*}$; up to an orthonormal transformation, $\hat{M}$ recovers $M^{* \prime}:=\left(\bar{Q}^{*}\right)^{1 / 2} M^{*}$, the representation function under an equivalent parameterization, termed as the normalized parameterization in (Tian et al., 2022), where

$$
\begin{gathered}
A^{* \prime}=\left(\bar{Q}^{*}\right)^{1 / 2} A^{*}\left(\bar{Q}^{*}\right)^{-1 / 2}, \quad B^{* \prime}=\left(\bar{Q}^{*}\right)^{1 / 2} B, \quad C^{* \prime}=C^{*}\left(\bar{Q}^{*}\right)^{-1 / 2}, \\
w_{t}^{\prime}=\left(\bar{Q}^{*}\right)^{1 / 2} w_{t}, \quad Q^{* \prime}=\left(\bar{Q}^{*}\right)^{-1 / 2} Q^{*}\left(\bar{Q}^{*}\right)^{-1 / 2} .
\end{gathered}
$$

Due to the following proposition, the algorithm does not need to know the dimension $d_{x}$ of the latent model; it can discover $d_{x}$ from the eigenvalues of $\hat{N}$.

Proposition 2. Under i.i.d. control inputs $u_{t} \sim \mathcal{N}\left(0, \sigma_{u}^{2} I\right)$ for $t \geq 0, \lambda_{\min }\left(\operatorname{Cov}\left(z_{t}^{*}\right)\right)=\Omega\left(v^{2}\right)$ for $t \geq d_{x}$, where $v$ is defined in Assumption 1.3. Recall that for a square matrix $A$, we define $\alpha(A):=$
$\sup _{k \geq 0}\left\|A^{k}\right\|_{2} \rho(A)^{-k}$. As long as $H \geq \frac{\log \left(a \alpha\left(\bar{A}^{*}\right)\right)}{\log \left(\rho\left(\bar{A}^{*}\right)^{-1}\right)}$ for some dimension-free constant $a>0, M^{*}$ has $\operatorname{rank} d_{x}$ and $\sigma_{\min }\left(M^{*}\right) \geq \Omega\left(v H^{-1 / 2}\right)$.

Proof. For $t \geq d_{x}$, unrolling the Kalman filter gives

$$
\begin{aligned}
z_{t}^{*} & =A^{*} z_{t-1}^{*}+B^{*} u_{t-1}+L^{*} i_{t} \\
& =A^{*}\left(A^{*} z_{t-2}^{*}+B^{*} u_{t-2}+L^{*} i_{t-1}\right)+L^{*} i_{t} \\
& =\left[B^{*}, \ldots,\left(A^{*}\right)^{d_{x}-1} B^{*}\right]\left[u_{t-1} ; \ldots ; u_{t-d_{x}}\right]+\left(A^{*}\right)^{d_{x}} z_{t-d_{x}}^{*}+\left[L^{*}, \ldots,\left(A^{*}\right)^{d_{x}-1} L^{*}\right]\left[i_{t} ; \ldots ; i_{t-d_{x}+1}\right],
\end{aligned}
$$

where $\left(u_{\tau}\right)_{\tau=t-d_{x}}^{t-1} z_{t-d_{x}}^{*}$ and $\left(i_{\tau}\right)_{\tau=t-d_{x}+1}^{t}$ are independent. For $H \geq d_{x}$, the matrix multiplied by $\left[u_{t-1} ; \ldots ; u_{t-d_{x}}\right]$ is precisely the controllability matrix $\Phi_{c}\left(A^{*}, B^{*}\right)$. Then,

$$
\begin{aligned}
\operatorname{Cov}\left(z_{t}^{*}\right)=\mathbb{E}\left[z_{t}^{*}\left(z_{t}^{*}\right)^{\top}\right] & \succcurlyeq \Phi_{c}\left(A^{*}, B^{*}\right) \mathbb{E}\left[\left[u_{t-1} ; \ldots ; u_{t-\ell}\right]\left[u_{t-1} ; \ldots ; u_{t-\ell}\right]^{\top}\right] \Phi_{c}^{\top}\left(A^{*}, B^{*}\right) \\
& =\sigma_{u}^{2} \Phi_{c}\left(A^{*}, B^{*}\right) \Phi_{c}^{\top}\left(A^{*}, B^{*}\right) .
\end{aligned}
$$

By the $v$-controllability of $\left(A^{*}, B^{*}\right), \operatorname{Cov}\left(z_{t}^{*}\right)$ is full-rank and $\lambda_{\min }\left(\operatorname{Cov}\left(z_{t}^{*}\right)\right) \geq \sigma_{u}^{2} v^{2}$. Since $z_{t}^{*}=$ $M^{*} h_{t}+\delta_{t}$, we have

$$
\operatorname{Cov}\left(M^{*} h_{t}\right)=\operatorname{Cov}\left(z_{t}^{*}-\delta_{t}\right)=\operatorname{Cov}\left(z_{t}^{*}\right)+\operatorname{Cov}\left(\delta_{t}\right)-\operatorname{Cov}\left(z_{t}^{*}, \delta_{t}\right)-\operatorname{Cov}\left(\delta_{t}, z_{t}^{*}\right)
$$

By Lemma 3,

$$
\left\|\operatorname{Cov}\left(z_{t}^{*}, \delta_{t}\right)\right\|_{2}=\left\|\operatorname{Cov}\left(\delta_{t}, z_{t}^{*}\right)\right\|_{2} \leq\left\|\operatorname{Cov}\left(z_{t}^{*}\right)^{1 / 2}\right\|_{2}\left\|\operatorname{Cov}\left(\delta_{t}\right)^{1 / 2}\right\|_{2} .
$$

Hence, by Weyl's inequality,

$$
\lambda_{\min }\left(\operatorname{Cov}\left(M^{*} h_{t}\right)\right) \geq \lambda_{\min }\left(\operatorname{Cov}\left(z_{t}^{*}\right)\right)-2\left\|\operatorname{Cov}\left(z_{t}^{*}\right)^{1 / 2}\right\|_{2}\left\|\operatorname{Cov}\left(\delta_{t}\right)^{1 / 2}\right\|_{2} .
$$

Since $\left\|\operatorname{Cov}\left(z_{t}^{*}\right)\right\|_{2}=\mathcal{O}(1)$ due to the stability of $A^{*}$ and $\delta_{t}=\left(\bar{A}^{*}\right)^{H} z_{t-H}^{*}$, there exists some dimension-free constant $a>0$ such that as long as $H \geq \frac{\log \left(a \alpha\left(\bar{A}^{*}\right)\right)}{\log \left(\rho\left(\bar{A}^{*}\right)^{-1}\right)}$,

$$
\lambda_{\min }\left(\operatorname{Cov}\left(M^{*} h_{t}\right)\right) \geq \sigma_{u}^{2} v^{2} / 2
$$

On the other hand,

$$
\mathbb{E}\left[M^{*} h_{t} h_{t}^{\top}\left(M^{*}\right)^{\top}\right] \preccurlyeq\left\|\mathbb{E}\left[h_{t} h_{t}^{\top}\right]\right\|_{2} M^{*}\left(M^{*}\right)^{\top} .
$$

Since $h_{t}=\left[y_{(t-H+1): t} ; u_{(t-H):(t-1)}\right]$ and $\left(\operatorname{Cov}\left(y_{t}\right)\right)_{t \geq 0},\left(\operatorname{Cov}\left(u_{t}\right)\right)_{t \geq 0}$ have $\mathcal{O}(1)$ operator norms, by Lemma $4,\left\|\operatorname{Cov}\left(h_{t}\right)\right\|_{2}=\left\|\mathbb{E}\left[h_{t} h_{t}^{\top}\right]\right\|_{2}=\mathcal{O}(H)$. Hence,

$$
0<\sigma_{u}^{2} v^{2} / 2 \leq \lambda_{\min }\left(\operatorname{Cov}\left(M^{*} h_{t}\right)\right)=\mathcal{O}(H) \sigma_{d_{x}}^{2}\left(M^{*}\right)
$$

Since $M^{*} \in \mathbb{R}^{d_{x} \times d_{h}}$, this implies that $\operatorname{rank}\left(M^{*}\right)=d_{x}$ and $\sigma_{\min }\left(M^{*}\right)=\Omega\left(v H^{-1 / 2}\right)$.

Proposition 2 is an adaption of (Tian et al., 2022, Proposition 2) to the infinite-horizon LTI setting. Necessarily, this implies that by our choice of $H, d_{h}=H\left(d_{y}+d_{u}\right) \geq d_{x}$. Moreover, since $\bar{Q}^{*} \succcurlyeq \mu^{2} I, N^{*}=\left(M^{*}\right)^{\top} \bar{Q}^{*} M^{*}$ is a $d_{h} \times d_{h}$ matrix with rank $d_{x}$, and $\lambda_{\text {min }}^{+}\left(N^{*}\right) \geq$ $\lambda_{\min }\left(\bar{Q}^{*}\right) \lambda_{\text {min }}^{2}\left(M^{*}\right)=\Omega\left(\mu^{2} v^{2} H^{-1}\right)$. Hence, if $\hat{N}$ is sufficiently close to $N^{*}$, by setting an appropriate threshold on the eigenvalues of $\hat{N}$, the dimension of the latent model equals the number of eigenvalues above it.

To find an approximate factorization of $\hat{N}$, let $\hat{N}=U \Lambda U^{\top}$ be its eigenvalue decomposition, where the diagonal elements of $\Lambda$ are listed in a descending order, and $U$ is an orthonormal matrix. Let $\Lambda_{d_{x}}$ be the left-top block of $\Lambda$ and $U_{d_{x}}$ be the left $d_{x}$ columns of $U$. By the Eckart-Young-Mirsky theorem, $\hat{M}=\max \left(\Lambda_{d_{x}}, 0\right)^{1 / 2} U_{d_{x}}^{\top}$, where "max" applies elementwise, is the solution to Line 4 of Algorithm 1, that is, the best approximate factorization of $\hat{N}$ among $d_{x} \times d_{h}$ matrices in terms of the Frobenius norm approximation error.

In the next two subsections, we move on to discuss learning latent dynamics, including the explicit approach SysId and the implicit approach CoSysId.

### 3.2 Explicit learning of system dynamics

Explicit learning of the system dynamics simply minimizes the transition prediction error in the latent space (Subramanian et al., 2020), or more generally, the statistical distances between the predicted and estimated distributions of the next latent state, like the KL divergence (Hafner et al., 2019a). In linear systems, it suffices to use the ordinary least squares as the SysId procedure, that is, to solve

$$
\begin{equation*}
(\hat{A}, \hat{B}) \in \underset{A, B}{\operatorname{argmin}} \sum_{t=H}^{T+H-1}\left\|A \hat{z}_{t}+B u_{t}-\hat{z}_{t+1}\right\|^{2} . \tag{3.3}
\end{equation*}
$$

In this linear regression, if $\left(\hat{z}_{t}\right)_{t \geq H}$ are the optimal state estimates $\left(z_{t}^{*}\right)_{t \geq H}(2.9)$, then (Simchowitz et al., 2018) has shown finite-sample guarantees for $(\hat{A}, \hat{B})$. Here, $\hat{z}_{t}$ contains errors resulting from the representation function $\hat{M}$ and the residual error $\delta_{t}$, but as long as $T$ and $H$ are large enough, SysId still has a finite-sample guarantee, as will be shown in Lemma 8. We refer to the algorithm that instantiates Algorithm 1 with SysId as CoReL (Cost-driven state Representation Learning). As the time-varying counterpart in (Tian et al., 2022), it provably solves learning for LQG control, as will be shown in Theorem 2.

### 3.3 Implicit learning of system dynamics (MuZero-style)

An important ingredient of latent model learning in MuZero (Schrittwieser et al., 2020) is to implicitly learn the transition function by minimizing the cost prediction error at future latent states generated from the transition function. Let $z_{t}=M h_{t}$ denote the latent state given by representation function $M$ at step $t$. Let $z_{t, 0}=z_{t}$ and $z_{t, i}=A z_{t, i-1}+B u_{t+i-1}$ for $i \geq 1$ be the future latent state predicted by dynamics $(A, B)$ from $z_{t}$ after $i$ steps of transition. For a trajectory of length $T+H$ like (3.1), the loss that considers $\ell$ steps into the future is given by

$$
\sum_{t=H}^{T+H-K-1} \sum_{i=0}^{\ell}\left(\left\|z_{t, i}\right\|_{Q}^{2}+\left\|u_{t}\right\|_{R^{*}}^{2}+b-c_{t}\right)^{2}
$$

```
Algorithm 2 CoSysId: Cost-driven system identification
    Input: data \(\mathcal{D}_{\text {raw }}\), representation function \(\hat{M}\)
    Estimate the system dynamics by
```

$$
\begin{equation*}
\hat{N}_{1}, \hat{b}_{1} \in \underset{N_{1}=N_{1}^{\top}, b_{1}}{\operatorname{argmin}} \sum_{t=H}^{T+H-1}\left(\left\|\left[h_{t} ; u_{t}\right]\right\|_{N_{1}}^{2}+b_{1}-\bar{c}_{t+1}\right)^{2} \tag{3.6}
\end{equation*}
$$

3: Find $\hat{M}_{1} \in \operatorname{argmin}_{M_{1} \in \mathbb{R}^{d x \times\left(H d_{y}+(H+1) d_{u}\right)}}\left\|M_{1}^{\top} M_{1}-\hat{N}_{1}\right\|_{F}$
4: Split $\hat{M}_{1}$ to $[\widetilde{M}, \widetilde{B}]$ after column $H\left(d_{y}+d_{u}\right)$ and set $\widetilde{A}=\widetilde{M} \hat{M}^{+}$.
5: Find alignment matrix $\hat{S}_{0}$ by

$$
\begin{equation*}
\hat{S}_{0} \in \underset{S_{0} \in \mathbb{R}^{d_{x} \times d_{x}}}{\operatorname{argmin}} \sum_{t=H}^{T+H-1}\left\|S_{0} \hat{M}_{1}\left[h_{t} ; u_{t}\right]-\hat{M} h_{t+1}\right\|^{2} \tag{3.7}
\end{equation*}
$$

6: Return: system dynamics estimate $(\hat{A}, \hat{B})=\left(\hat{S}_{0} \widetilde{A}, \hat{S}_{0} \widetilde{B}\right)$

This loss involves powers of $A$ up to $A^{\ell}$; with the squared norm, the powers double, making the minimization over $A$ hard to solve and analyze for $\ell \geq 2$. In LQG control, our finding is that it suffices to take $\ell=1$. As mentioned in $\S 1$, MuZero also predicts optimal values and optimal actions; in LQG, to handle $Q^{*} \nsucc 0$, like cost-driven representation learning (see §3.1), we adopt the cumulative costs and use the normalized parameterization. Thus, the optimization problem we aim to solve is given by

$$
\begin{equation*}
\min _{M, A, B, b} \sum_{t=H}^{T+H-1}\left(\left(\left\|M h_{t}\right\|^{2}+b-\bar{c}_{t}\right)^{2}+\left(\left\|A M h_{t}+B u_{t}\right\|^{2}+b-\bar{c}_{t+1}\right)^{2}\right) . \tag{3.4}
\end{equation*}
$$

To convexify the optimization problem (3.4), we define $N:=M^{\top} M$ and $N_{1}:=[A M, B]^{\top}[A M, B]$. Then, (3.4) becomes

$$
\begin{equation*}
\min _{N, N_{1}, b} \sum_{t=H}^{T+H-1}\left(\left(\left\|h_{t}\right\|_{N}^{2}+b-\bar{c}_{t}\right)^{2}+\left(\left\|\left[h_{t} ; u_{t}\right]\right\|_{N_{1}}^{2}+b-\bar{c}_{t+1}\right)^{2}\right) . \tag{3.5}
\end{equation*}
$$

This minimization problem is convex in $N, N_{1}$ and $b$, and has a closed-form solution; essentially, it consists of two linear regression problems coupled by $b$. Since constant $b$ is merely a term accounting for the estimation error and not part of the representation function, we can decouple the two regression problems by allowing $b$ to take different values in them. This further simplifies the algorithm: the first regression problem is exactly cost-driven representation learning (§3.1), and the second is cost-driven system identification (CoSysId, Algorithm 2). The algorithm that instantiates Algorithm 1 with CoSysId is called CoReDyL (Cost-driven state Representation and Dynamic Learning). Like CoReL, this MuZero-style latent model learning method provably solves LQG control, as we will show in Theorem 2.

CoSysId has similar steps to cost-driven representation learning (\$3.1), except that in Line 5 of Algorithm 2, it requires fitting a matrix $\hat{S}_{0}$. This is because the cost is invariant to orthonormal transformations of the latent states, and the approximate factorization steps recover $M^{*}$ and $M_{1}^{*}$ up to orthonormal transformations $S$ and $S_{1}$, but there is no guarantee for the two transformations to be the same. MuZero bypasses this problem by predicting multiple steps of costs into
the future, but analyzing such an optimization function involves the additional complexity to deal with powers of $A$. Here, we simply estimate this alignment matrix $S_{0}=S S_{1}^{\top}$. We note that although CoSysId needs the output $\hat{M}$ from cost-driven representation learning, the two quadratic regressions (3.2) and (3.6) are not coupled and can be solved in parallel.

## 4 Theoretical guarantees

The following Theorem 1 shows that both CoReL and CoReDyL are guaranteed to solve unknown LQG control with a finite number of samples.

Theorem 1. Given an unknown LQG problem satisfying Assumption 1, let $M^{* \prime}$ and $\left(A^{* \prime}, B^{* \prime}, Q^{* \prime}, R^{*}\right)$ be the optimal state representation function and the true system parameters under the normalized parameterization. For a given $p \in(0,1)$, if we run CoReL (Algorithm 1 with (3.3)) or CoReDyL (Algorithm 1 with Algorithm 2) for $\left.T \geq \operatorname{poly}\left(d_{x}, d_{y}, d_{u}, \log (T / p)\right)\right), H=\Omega(\log (H T))$, and $\sigma_{u}=\Theta(1)$, then there exists an orthonormal matrix $S \in \mathbb{R}^{d_{x} \times d_{x}}$, such that with probability at least $1-p$, the representation function $\hat{M}$ satisfies

$$
\left\|\hat{M}-S M^{* \prime}\right\|_{2}=\mathcal{O}\left(\operatorname{poly}\left(H, d_{x}, d_{u}, d_{y}, \log (T / p)\right) T^{-1 / 2}\right)
$$

and the feedback gain $\hat{K}$ satisfies

$$
\begin{aligned}
& J^{\hat{K}}\left(S A^{* \prime} S^{\top}, S B^{* \prime}, S Q^{* \prime} S^{\top}, R^{*}\right)-J^{*}\left(S A^{* \prime} S^{\top}, S B^{* \prime}, S Q^{* \prime} S^{\top}, R^{*}\right) \\
= & \mathcal{O}\left(\operatorname{poly}\left(H, d_{x}, d_{u}, d_{y}, \log (T / p)\right) T^{-1}\right) .
\end{aligned}
$$

We defer the proof of Theorem 1 to §F. Compared with the time-varying setting in (Tian et al., 2022), the bounds here do not have a separation for the initial steps and future steps, where the bounds are much worse before the system is fully excited. This is because in the time-invariant setting, the representation function and the latent model are both stationary. On the other hand, to learn such stationary functions across different time steps, we need to aggregate correlated data along a single trajectory, which poses significant challenges for the analysis. A major effort to overcome such difficulties involves proving a new result on the persistency of excitation (Lemma 1) using the small-ball method (Mendelson, 2015; Simchowitz et al., 2018), discussed in $\S 4.1$ in more detail.

Compared with common system identification methods based on learning Markov parameters (Oymak and Ozay, 2019; Simchowitz et al., 2019), the error bounds of the system parameters produced by CoReDyL (or CoReL) have the same dependence on $T$, but worse dependence on system dimensions. Moreover, to establish persistency of excitation, CoReDyL (or CoReL) requires a larger burn-in period. These relative sample inefficiencies are the price we pay for direct latent model learning, which is only supervised by scalar-valued costs that are quadratic in the history, instead of vector-valued observations that are linear in the history. Hence, we have to address the more challenging problem of quadratic regression, which lifts the dimension of the optimization problem. On the other hand, direct latent model learning avoids learning the reconstruction function $C^{*}$ and can learn task-relevant representations in more complex settings, as demonstrated by empirical studies.

### 4.1 Persistency of excitation

Central to the analysis of CoReL and CoReDyL is the finite-sample characterization of the quadratic regression problem. To solve (3.2), notice that

$$
\left\|h_{t}\right\|_{N}^{2}=\left\langle N, h_{t} h_{t}^{\top}\right\rangle_{F}=\left\langle\operatorname{svec}(N), \operatorname{svec}\left(h_{t} h_{t}^{\top}\right)\right\rangle
$$

so this quadratic regression is essentially a linear regression problem in terms of $\left[\operatorname{svec}(N) ; b_{0}\right]$. A major difficulty in the analysis is to establish persistency of excitation for $\left(\left[\operatorname{svec}\left(h_{t} h_{t}^{\top}\right) ; 1\right]\right)_{t \geq H}$, meaning that the minimum eigenvalue of the design matrix $\sum_{t=H}^{T+H-1}\left[\operatorname{svec}\left(h_{t} h_{t}^{\top}\right) ; 1\right]\left[\operatorname{svec}\left(h_{t} h_{t}^{\top}\right)^{\top}, 1\right]$ grows linearly in the size $T$ of the data.

A linear lower bound on $\lambda_{\min }\left(\sum_{t=H}^{T+H-1} h_{t} h_{t}^{\top}\right)$ is a known result for the identification of partially observable linear dynamical systems, see the recent overview in (Tsiamis et al., 2022). In our case, however, elements of $\operatorname{svec}\left(h_{t} h_{t}^{\top}\right)$ are products of Gaussians, making the analysis difficult. If $\left(h_{t}\right)_{t \geq H}$ are independent, which is the case if they are from multiple independent trajectories, the result has been established in (Jadbabaie et al., 2021; Tian et al., 2022). It can also be proved with the matrix Azuma inequality (Tropp, 2012). Here, we need to aggregate correlated data to estimate a set of stationary parameters. In sum, the difficulty we face results from both products of Gaussians and the data dependence.

In principle, given enough burn-in time, state $x_{t}$, and hence observation $y_{t}$ and truncated history $h_{t}$, converges to the steady-state distributions, and samples with an interval of the order of mixing time are approximately independent (Levin and Peres, 2017). Hence, intuitively, a linear lower bound is viable. However, the bound yielded by such an analysis deteriorates as the system is less stable and the mixing time increases, which is qualitatively incorrect for linear systems. To eschew such dependence, (Simchowitz et al., 2018) introduces the small-ball method. We take the same route, while establishing different arguments to handle the products of Gaussians.

Let us first recall the block martingale small-ball condition.
Definition 1 (Block martingale small ball (BMSB) condition (Simchowitz et al., 2018, Definition 2.1)). Let $\left(f_{t}\right)_{t \geq 1}$ be a stochastic process in $\mathbb{R}^{d}$ adapted to filtration $\left(\mathcal{F}_{t}\right)_{t \geq 1}$. We say $\left(f_{t}\right)_{t \geq 1}$ satisfies the $(k, \Gamma, q)$-BMSB condition for $k \in \mathbb{N}^{+}, \Gamma \succcurlyeq 0$ and $q>0$, if for any $t \geq 1$, for any fixed unit vector $v \in \mathbb{R}^{d}, \frac{1}{k} \sum_{i=1}^{k} \mathbb{P}\left(\left|\left\langle f_{t+i}, v\right\rangle\right| \geq\|v\|_{\Gamma} \mid \mathcal{F}_{t}\right) \geq q$ almost surely.

The key Lemma 1 below shows that $\left(\operatorname{svec}\left(h_{t} h_{t}^{\top}\right)\right)_{t \geq H}$ satisfies the BMSB condition.
Lemma 1. Let $h_{t}=\left[y_{(t-H+1): ;} ; u_{(t-H):(t-1)}\right]$ be the $H$-step history at time step $t \geq H$ in system (2.1) with $u_{t} \sim \mathcal{N}\left(0, \sigma_{u}^{2} I\right)$ for $t \geq 0$. Define filtration $\mathcal{F}_{t}:=\sigma\left(x_{0}, v_{0}, u_{0}, w_{0}, v_{1}, \ldots, u_{t-1}, w_{t-1}, v_{t}\right)$. Define $f_{t}:=\operatorname{svec}\left(h_{t} h_{t}^{\top}\right)$, adapted to $\left(\mathcal{F}_{t}\right)_{t \geq H}$. As long as $H \geq \frac{a \log \left(d_{h} \alpha\left(A^{*}\right) \log (T / p)\right)}{\log \left(\rho\left(A^{*}\right)^{-1}\right)},\left(f_{t}\right)_{t \geq H}$ is $\left(k, \gamma^{2} I, q\right)$-BMSB for $k=4 H, \gamma=\Theta\left(1 / d_{h}\right)$, and $q=\Theta\left(1 / d_{h}^{2}\right)$, where $\Theta(\cdot)$ hides the dependence on dimension-free constants.

Then, following the analysis in (Simchowitz et al., 2018, Appendix D), we can show that if additionally, $T$ is large enough, then with high probability,

$$
\lambda_{\min }\left(\sum_{t=H}^{T+H-1} f_{t} f_{t}^{\top}\right)=\Omega\left(\gamma^{2} q^{2} T\right)=\Omega\left(T / d_{h}^{6}\right)
$$

which establishes the persistency of excitation.
The full proof of Lemma 1 is deferred to Appendix D.1. Crucial for its proof is the following Lemma 2, which might be of independent interest.

Lemma 2. Let $x$ be a d-dimensional zero-mean Gaussian random vector with covariance $\Sigma$. For any $d \times d$ symmetric matrix $A$ and constant $b \in \mathbb{R}$ that satisfy $\|A\|_{F}^{2}+b^{2}=1$, there exists an absolute constant $a>0$, such that $\mathbb{E}\left[\left|x^{\top} A x+b\right|\right] \geq a \lambda_{\min }(\Sigma) d^{-3 / 2}$.

### 4.2 Main ideas in proving Theorem 1

Below we sketch the ideas of proving Theorem 1 for CoReDyL. Define $M_{1}^{*}:=\left[A^{*} M^{*}, B^{*}\right]$ as the composition of one-step transition and representation functions and $N_{1}^{*}:=\left(M_{1}^{*}\right)^{\top} M_{1}^{*}$, which is estimated by $\hat{N}_{1}$ in (3.6).

For the quadratic regression problems (3.2) and (3.6), by establishing persistency of excitation (Lemma 7) and using sub-Weibull martingale concentration (Lemma 6), we can show that

$$
\begin{aligned}
\left\|\hat{N}-N^{*}\right\|_{F} & =\mathcal{O}\left(\left(H\left(d_{y}+d_{u}\right)\right)^{10} d_{x}^{3 / 2} T^{-1 / 2} \log ^{3}(T / p)\right) \\
\left\|\hat{N}_{1}-N_{1}^{*}\right\|_{F} & =\mathcal{O}\left(\left(H\left(d_{y}+d_{u}\right)+d_{x}\right)^{10} d_{x}^{3 / 2} T^{-1 / 2} \log ^{3}(T / p)\right)
\end{aligned}
$$

our choice of $H$ ensures that the truncation errors are absorbed into these bounds. By the Procrustes-type lemma in (Tu et al., 2016, Lemma 5.4), these two bounds imply bounds of the same order on $\left\|\hat{M}-S M^{* \prime}\right\|_{F}$ and $\left\|\hat{M}_{1}-S_{1} M_{1}^{* \prime}\right\|_{F}$, respectively, for some orthonormal matrices $S$ and $S_{1}$. Thus, $\left\|\widetilde{B}-S_{1} B^{* \prime}\right\|_{2}$ is of the same order as the bound on $\left\|\hat{N}_{1}-N_{1}^{*}\right\|_{2}$, and by the perturbation bounds of the Moore-Penrose inverse (Wedin, 1973), so is $\left\|\widetilde{A}-S_{1} A^{* \prime} S^{\top}\right\|_{2}$.

As explained in $\S 3.3$, we need to fit the matrix $S_{0}=S S_{1}^{\top}$ to align the next states generated by $(\widetilde{A}, \widetilde{B})$ with those by $\hat{M}$. With an analysis of perturbed linear regression (Lemma 8), we show that $\hat{S}_{0}$ in Algorithm 2 satisfies $\left\|\hat{S}_{0}-S_{0}\right\|_{2}=\mathcal{O}\left(\left(H\left(d_{y}+d_{u}\right)+d_{x}\right)^{21 / 2} d_{x}^{2} T^{-1 / 2} \log ^{4}(T / p)\right)$. Therefore, the bounds on $\left\|\hat{A}-S A^{*} S^{\top}\right\|_{2}$ and $\left\|\hat{B}-S B^{*}\right\|_{2}$ are of the same order. Line 7 in Algorithm 1 requires an analysis of perturbed quadratic regression (Lemma 7), which guarantees that

$$
\left\|\hat{Q}-S Q^{*} S^{\top}\right\|_{F}=\mathcal{O}\left(\left(H\left(d_{y}+d_{u}\right) d_{x}\right)^{21 / 2} T^{-1 / 2} \log ^{4}(T / p)\right) .
$$

Hence, $\left\|\hat{A}-S A^{*} S^{\top}\right\|_{2},\left\|\hat{B}-S B^{*}\right\|_{2}$ and $\left\|\hat{Q}-S Q^{*} S^{\top}\right\|_{2}$ are all $\mathcal{O}\left(\left(H\left(d_{y}+d_{u}\right) d_{x}\right)^{21 / 2} T^{-1 / 2} \log ^{4}\left(d_{x} T / p\right)\right)$.
Lastly, we invoke the result on certainty equivalent LQR in (Mania et al., 2019) to certify the suboptimality gap of $\hat{K}$ obtained from $\left(\hat{A}, \hat{B}, \hat{Q}, R^{*}\right)$. This certainty equivalent controller can also be replaced by a robust one (Dean et al., 2020). Let $K^{* \prime}$ denote the optimal controller under the normalized parameterization. Then, for a large enough burn-in period, the certainty-equivalent controller satisfies that $\left\|\hat{K}-K^{* \prime}\right\|_{2}$ is of the same order as the system parameter errors. Note that policy $K^{*} M^{*}$ in the original system is independent of the latent model parameterization, and we have that $\left\|\hat{K} \hat{M}-K^{*} M^{*}\right\|_{2}$ is of the same order.


Figure 1: Optimality of the learned policy versus the length of the trajectory. Each dot averages the results of five independent runs.

## 5 Numerical results

Although this work is theoretical in nature, we conduct preliminary numerical experiments by testing CoReL on a mass-spring-damper system, with mass 1 kg , stiffness $10 \mathrm{~N} / \mathrm{m}$ and damping coefficient $100 \mathrm{~kg} / \mathrm{s}$. We discretize time with the forward Euler method using 0.01 s intervals. The state of the system is the position and the velocity of the mass, with the observation being the position only. The scalar control input is the external force. We take $Q=I$ and $R=1$ in the cost function. We set $H=5$ and measure the performance by the distance between $\hat{K} \hat{M}$ and $K^{*} M^{*}$. In all experiments, we warm up the system by 1000 steps.

The simulation results are summarized in Figure 1. In general, we observe the decrease of the distance to the optimal policy as the length of the trajectory increases. The system has $d_{x}=2$. Besides $d_{z}=2$, we experiment with $d_{z}=1$ and $d_{z}=3$ to examine the impact of dimension mismatch. In this example, we see that $d_{z}=1$ does not produce a meaningful policy and $d_{z}=3$ performs better than the $d_{z}=2$ baseline. We note that the these results are preliminary, and this "blessing of overparameterization" in the latent state space is worth further investigation.

## 6 Conclusion and future work

We studied direct latent model learning for solving unknown infinite-horizon time-invariant LQG control. We established finite-sample guarantees for two methods, which differ in whether the latent dynamics is learned explicitly by minimizing the transition prediction errors, or im-
plicitly by using the transition for future predictions, with the latter being closer to that used in MuZero (Schrittwieser et al., 2020). For MuZero-style latent model learning, our analysis identifies a coordinate misalignment problem in the latent state space, suggesting the value of multi-step future prediction. A limitation of this work is that we only consider state representation based on truncated histories, i.e., frame stacking, as used in MuZero; the recursive form of the representation function, as in the Kalman filter, is also used empirically (Ha and Schmidhuber, 2018; Hafner et al., 2019a), and might be worth further investigation.

Many questions remain to be answered in representation learning for control. Provable generalization of direct latent model learning to nonlinear observation channels or dynamics is a natural consideration. Moreover, with the ubiquity of visual perception in real-world control systems, what if we have a time-varying observation function or multiple observation functions, modeling images taken from different angles? In reality, most of the time we do not have a well-defined cost function; learning task-relevant state representations from demonstrations is another intriguing direction.

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## A Additional related work

(Oymak and Ozay, 2019) studies the identification of partially observable linear dynamical systems from a single trajectory, which presents a finite-sample analysis of identifying the Markov parameter and a perturbation analysis of the Ho-Kalman algorithm (Ho and Kalman, 1966). (Simchowitz et al., 2019) relaxes the stability requirement to marginal stability by using prefiltered least squares to identify the Markov parameter. The method in (Zheng and Li, 2020) applies to unstable systems but requires multiple trajectories. Since the Markov parameter maps control input histories to observations, these methods do not work with costs and use the Markov parameter as an intermediate step to identify the system. By contrast, our methods, entirely driven by the costs and closely connected with empirical methods, directly learn the representation function and the latent model. Directly learning the latent model connects our work to the identification of fully observable linear dynamical systems. (Simchowitz et al., 2018) introduces small-ball conditions to handle correlated data and characterizes the statistical rates for stable and unstable systems, both proving to be useful for our analysis.

Online control of partially observable linear dynamical systems is considered in (Lale et al., 2020, 2021) for stochastic noises and in (Simchowitz et al., 2020) for nonstochastic noises. (Zheng et al., 2021) considers end-to-end sample complexity and is closest to our setup. All these methods rely on the estimation of Markov parameters. For a discussion of the literature in more detail and breadth, we refer the reader to the recent survey (Tsiamis et al., 2022).

## B Additional discussions on CoSysId

In CoSysId (Algorithm 2), the covariates of quadratic regression in (3.6) are $\left(\left[h_{t} ; u_{t}\right]\right)_{t \geq H}$. One may wonder if we can pursue an alternative approach by fixing $M$ to be $\hat{M}$, and using $\left(\left[\hat{z}_{t} ; u_{t}\right]\right)_{t \geq H}$ as covariates, which have a much lower dimension, though the two quadratic regressions cannot be solved in parallel anymore. Specifically, the new quadratic regression we need to solve is given by

$$
\hat{N}_{2}, \hat{b}_{2} \in \underset{N_{2}=N_{2}^{\top}, b_{2}}{\operatorname{argmin}} \sum_{t=H}^{T+H-1}\left(\left\|\left[\hat{z}_{t} ; u_{t}\right]\right\|_{N_{2}}^{2}+b_{2}-\bar{c}_{t+1}\right)^{2}
$$

where $\hat{z}_{t}=\hat{M} h_{t}$ is an approximation of $S z_{t}^{*}$. The ground truth for $\hat{N}_{2}$ is $N_{2}^{*}=\left[S A^{*} S^{\top}, S B^{*}\right]^{\top}\left[S A^{*} S^{\top}, S B^{*}\right]$, so its approximate factorization recovers $\left[S_{2} A^{*} S^{\top}, S_{2} B^{*}\right]$ for some orthonormal matrix $S_{2}$. In a similar way to CoSysId, we still need to fit an alignment matrix $S_{3}=S S_{2}^{\top}$ to align the coordinates. Let $\widetilde{A}, \widetilde{B}$ denote the system parameters recovered from $\hat{N}_{2}$. The linear regression we now need to solve is from $\left([\widetilde{A}, \widetilde{B}]\left[\tilde{z}_{t} ; u_{t}\right]\right)_{t=H}^{T+H-1}$ to $\left(\hat{z}_{t+1}\right)_{t=H}^{T+H-1}$. However, without further assumptions, $\left[A^{*}, B^{*}\right]$ does not necessarily have full row rank, and hence, neither does $[\widetilde{A}, \widetilde{B}]$, in which case recovering the entire $S_{3}$ is impossible.

On the other hand, for CoSysId (Algorithm 2), the ground truth of $\hat{M}_{1}$ is $M_{1}^{*}=\left[A^{*} M^{*}, B^{*}\right]$, which is guaranteed to have full row rank by the same argument as the proof of Proposition 2, since $M_{1}^{*}\left[h_{t} ; u_{t}\right]$ estimates $z_{t+1}^{*}$, which has full-rank covariance. Hence, recovering $S_{0}=S S_{1}^{\top}$ is feasible.

## C Auxiliary results

## C. 1 Basic inequalities involving random vectors

Lemma 3. Let $x$ and $y$ be random vectors defined on the same probability space. Then, $\left\|\mathbb{E}\left[x y^{\top}\right]\right\|_{2}^{2} \leq$ $\left\|\mathbb{E}\left[x x^{\top}\right]\right\|_{2} \cdot\left\|\mathbb{E}\left[y y^{\top}\right]\right\|_{2}$.

Proof. Let $d_{x}, d_{y}$ be the dimensions of the values of $x, y$, respectively. For any vectors $v \in \mathbb{R}^{d_{x}}$, $w \in \mathbb{R}^{d_{y}}$, by the Cauchy-Schwarz inequality,

$$
\begin{aligned}
\left(v^{\top} \mathbb{E}\left[x y^{\top}\right] w\right)^{2} & =\left(\mathbb{E}\left[v^{\top} x y^{\top} w\right]\right)^{2} \\
& \leq \mathbb{E}\left[\left(v^{\top} x\right)^{2}\right] \cdot \mathbb{E}\left[\left(w^{\top} y\right)^{2}\right] \\
& =\mathbb{E}\left[v^{\top} x x^{\top} v\right] \cdot \mathbb{E}\left[w^{\top} y y^{\top} w\right] \\
& =\left(v^{\top} \mathbb{E}_{x x^{\top}} v\right) \cdot\left(w^{\top} \mathbb{E}\left[y y^{\top}\right] w\right) .
\end{aligned}
$$

Taking the maximum over $v, w$ on both sides subject to $\|v\|,\|w\| \leq 1$ gives

$$
\left\|\mathbb{E}\left[x y^{\top}\right]\right\|_{2}^{2} \leq\left\|\mathbb{E}\left[x x^{\top}\right]\right\|_{2} \cdot\left\|\mathbb{E}\left[y y^{\top}\right]\right\|_{2}
$$

which completes the proof.
Lemma 4 ((Tian et al., 2022, Lemma 8)). Let $x, y$ be random vectors of dimensions $d_{x}, d_{y}$, respectively, defined on the same probability space. Then, $\|\operatorname{Cov}([x ; y])\|_{2} \leq\|\operatorname{Cov}(x)\|_{2}+\|\operatorname{Cov}(y)\|_{2}$.

## C. 2 Lower bound about Gaussian quadratic forms

Lemma 5. Let $z_{1}, z_{2}, \ldots, z_{d}$ be independent standard Gaussian random variables. Let $v=\left[v_{1}, v_{2}, \ldots, v_{d+1}\right]^{\top} \in$ $\mathrm{S}^{d}$ be a $(d+1)$-dimensional unit vector. There exists an absolute constant $a>0$, such that

$$
\inf _{v \in S^{d}} \mathbb{E}\left[\left|v_{d+1}+\sum_{i=1}^{d} v_{i} z_{i}^{2}\right|\right] \geq a d^{-3 / 2}
$$

Proof. Let us consider the value of $v_{d+1}$. Since $\mathbb{E}\left[z_{i}^{2}\right]=1$ for all $1 \leq i \leq d$,

$$
\mathbb{E}\left[\left|\sum_{i=1}^{d} v_{i} z_{i}^{2}\right|\right] \leq \sum_{i=1}^{d}\left|v_{i}\right| \leq \sqrt{d \sum_{i=1}^{d} v_{i}^{2}} \leq \sqrt{d\left(1-v_{d+1}^{2}\right)}
$$

Then,

$$
\mathbb{E}\left[\left|v_{d+1}+\sum_{i=1}^{d} v_{i} z_{i}^{2}\right|\right] \geq\left|v_{d+1}\right|-\mathbb{E}\left[\left|\sum_{i=1}^{d} v_{i} z_{i}^{2}\right|\right] \geq\left|v_{d+1}\right|-\sqrt{d\left(1-v_{d+1}^{2}\right)}
$$

Hence, if $\left|v_{d+1}\right| \geq 2 \sqrt{d /(4 d+1)}$, we have $\sqrt{d\left(1-v_{d+1}^{2}\right)} \leq\left|v_{d+1}\right| / 2$. It follows that

$$
\mathbb{E}\left[\left|v_{d+1}+\sum_{i=1}^{d} v_{i} z_{i}^{2}\right|\right] \geq \frac{\left|v_{d+1}\right|}{2} \geq \sqrt{\frac{d}{4 d+1}} \geq \frac{1}{\sqrt{5}} .
$$

Below we consider the case where $\left|v_{d+1}\right|<2 \sqrt{d /(4 d+1)}$. Let $\operatorname{sign}(\cdot)$ denote the sign function. Let $\mathcal{I}^{+}:=\left\{i: \operatorname{sign}\left(v_{i}\right)=1,1 \leq i \leq d\right\}$ and $\mathcal{I}^{-}:=\left\{i: \operatorname{sign}\left(v_{i}\right)=-1,1 \leq i \leq d\right\}$ be the index sets of positive and negative values among $\left(v_{i}\right)_{i=1}^{d}$. Then,

$$
\begin{aligned}
\mathbb{E}\left[\left|v_{d+1}+\sum_{i=1}^{d} v_{i} z_{i}^{2}\right|\right] & =\mathbb{E}\left[\left|v_{d+1}+\sum_{i=1}^{d}\right| v_{i}\left|\operatorname{sign}\left(v_{i}\right) z_{i}^{2}\right|\right] \\
& =\mathbb{E}\left[\left|v_{d+1}+\sum_{i \in \mathcal{I}^{+}}\right| v_{i}\left|z_{i}^{2}-\sum_{j \in \mathcal{I}^{-}}\right| v_{j}\left|z_{j}^{2}\right|\right]
\end{aligned}
$$

For a given $v$, since $\left(z_{i}^{2}\right)_{i=1}^{d}$ have identical distributions, $\mathbb{E}\left[\left|v_{d+1}+\sum_{i \in \mathcal{I}^{+}}\right| v_{i}\left|z_{i}^{2}-\sum_{j \in \mathcal{I}^{-}}\right| v_{j}\left|z_{j}^{2}\right|\right]$ has the same value under permutations of $\left(v_{i}\right)_{i \in \mathcal{I}^{+}}$and $\left(v_{j}\right)_{j \in \mathcal{I}^{-}}$. Summing over all the permutations of $\left(v_{i}\right)_{i \in \mathcal{I}^{+}}$and $\left(v_{j}\right)_{j \in \mathcal{I}^{-}}$gives

$$
d \mathbb{E}\left[\left|v_{d+1}+\sum_{i=1}^{d} v_{i} z_{i}^{2}\right|\right] \geq \mathbb{E}\left[\left|d \cdot v_{d+1}+\left(\sum_{i \in \mathcal{I}^{+}}\left|v_{i}\right|\right) \sum_{i \in \mathcal{I}^{+}} z_{i}^{2}-\left(\sum_{j \in \mathcal{I}^{-}}\left|v_{j}\right|\right) \sum_{j \in \mathcal{I}^{-}} z_{j}^{2}\right|\right]
$$

Hence,

$$
\mathbb{E}\left[\left|v_{d+1}+\sum_{i=1}^{d} v_{i} z_{i}^{2}\right|\right] \geq \frac{1}{d}\left(\sum_{i=1}^{d}\left|v_{i}\right|\right) \mathbb{E}\left[\left|d \cdot v_{d+1}+\sum_{i=1}^{d} \operatorname{sign}\left(v_{i}\right) z_{i}^{2}\right|\right]
$$

Since $\sum_{i=1}^{d}\left|v_{i}\right| \geq\left(\sum_{i=1}^{d} v_{i}^{2}\right)^{1 / 2}=\left(1-v_{d+1}^{2}\right)^{1 / 2}$, we have

$$
\mathbb{E}\left[\left|v_{d+1}+\sum_{i=1}^{d} v_{i} z_{i}^{2}\right|\right] \geq \frac{\left(1-v_{d+1}^{2}\right)^{1 / 2}}{d} \inf _{w \in\{ \pm 1\}^{d}} \mathbb{E}\left[\left|d \cdot v_{d+1}+\sum_{i=1}^{d} w_{i} z_{i}^{2}\right|\right]
$$

It remains to lower bound $\inf _{w \in\{ \pm 1\}^{d}} \mathbb{E}\left[\left|d \cdot v_{d+1}+\sum_{i=1}^{d} w_{i} z_{i}^{2}\right|\right]$. By symmetry, for any pair $w_{i} \neq w_{j}$, the expectation remains the same if we exchange $z_{i}$ and $z_{j}$. Hence, for any random variable $x$,

$$
\mathbb{E}\left[\left|x+z_{i}-z_{j}\right|\right]=\frac{1}{2}\left(\mathbb{E}\left[\left|x+z_{i}-z_{j}\right|\right]+\mathbb{E}\left[\left|x+z_{i}-z_{j}\right|\right]\right) \geq \mathbb{E}[|x|] .
$$

We shall apply this symmetry trick in the following to cancel terms with opposite signs. Let $p$ denote the number of +1 's and $q$ denote the number of -1 's in $w$, such that $p+q=n$. If $p \neq q$, by the symmetry trick,

$$
\mathbb{E}\left[\left|d \cdot v_{d+1}+\sum_{i=1}^{d} w_{i} z_{i}^{2}\right|\right] \geq \mathbb{E}\left[\left|d \cdot v_{d+1}+\sum_{i=1}^{|p-q|} z_{i}^{2}\right|\right] \geq \operatorname{Var}\left(\sum_{i=1}^{|p-q|} z_{i}^{2}\right)=\Omega(1) .
$$

If $p=q$, again, the symmetry trick yields

$$
\mathbb{E}\left[\left|d \cdot v_{d+1}+\sum_{i=1}^{d} w_{i} z_{i}^{2}\right|\right] \geq \mathbb{E}\left[\left|d \cdot v_{d+1}+z_{1}^{2}-z_{2}^{2}\right|\right] \geq \operatorname{Var}\left(z_{1}^{2}-z_{2}^{2}\right)=\Omega(1)
$$

Hence, regardless of $p$ and $q$, we have $\inf _{w \in\{ \pm 1\}^{d}} \mathbb{E}\left[\left|d \cdot v_{d+1}+\sum_{i=1}^{d} w_{i} z_{i}^{2}\right|\right]=\Omega(1)$. Then,

$$
\mathbb{E}\left[\left|v_{d+1}+\sum_{i=1}^{d} v_{i} z_{i}^{2}\right|\right] \geq \Omega(1) \cdot \frac{\left(1-v_{d+1}^{2}\right)^{1 / 2}}{d}
$$

Since $\left|v_{d+1}\right|<2 \sqrt{d /(4 d+1)}$,

$$
\mathbb{E}\left[\left|v_{d+1}+\sum_{i=1}^{d} v_{i} z_{i}^{2}\right|\right]=\Omega(1) \cdot \frac{1}{\sqrt{4 d+1} \cdot d}=\Omega\left(d^{-3 / 2}\right) .
$$

Hence, over all, we have

$$
\inf _{v \in \mathrm{~S}^{d}} \mathbb{E}\left[\left|v_{d+1}+\sum_{i=1}^{d} v_{i} z_{i}^{2}\right|\right] \geq a d^{-3 / 2}
$$

From the proof, we can see that without $v_{d+1}$, we have the improved bound $\inf _{v \in S^{d-1}} \mathbb{E}\left[\left|\sum_{i=1}^{d} v_{i} z_{i}^{2}\right|\right]=$ $\Omega\left(d^{-1}\right)$.

Below we use Lemma 5 to prove Lemma 2.
Proof of Lemma 2. Let $y:=\Sigma^{-1 / 2} x$. Then $y$ is a standard Gaussian random vector, and $x^{\top} A x=$ $y^{\top} \Sigma^{1 / 2} A \Sigma^{1 / 2} y$. Let $U^{\top} \Lambda U$ be the eigenvalue decomposition of $\Sigma^{1 / 2} A \Sigma^{1 / 2}$. Then,

$$
x^{\top} A x=y^{\top} U^{\top} \Lambda U y=z^{\top} \Lambda z,
$$

where $z:=U y$ is still a standard Gaussian random vector.
By the unitary invariance of the Frobenius norm,

$$
\|\Lambda\|_{F}=\left\|U^{\top} \Lambda U\right\|_{F}=\left\|\Sigma^{1 / 2} A \Sigma^{1 / 2}\right\|_{F} \geq \lambda_{\min }(\Sigma)\|A\|_{F}
$$

Hence,

$$
\|\Lambda\|_{F}^{2}+b^{2} \geq \lambda_{\min }^{2}(\Sigma)\|A\|_{F}^{2}+b^{2} \geq \lambda_{\min }^{2}(\Sigma) \wedge 1
$$

Therefore, by Lemma 5 , there exists an absolute constant $a>0$, such that

$$
\inf _{\|A\|_{F}^{2}+b^{2}=1} \mathbb{E}\left[\left|x^{\top} A x+b\right|\right] \geq \inf _{\|\Lambda\|_{F}^{2}+b^{2} \geq \lambda_{\min }^{2}(\Sigma) \wedge 1} \mathbb{E}\left[\left|z^{\top} \Lambda z+b\right|\right] \geq a\left(\lambda_{\min }(\Sigma) \wedge 1\right) d^{-3 / 2}
$$

## C. 3 Sum of sub-Weibull martingale difference sequences

To upper bound $\sum_{t=H}^{T+H-1} f_{t} e_{t}$ in the analysis of quadratic regression (see §D), one possible approach (Simchowitz et al., 2018; Oymak and Ozay, 2019; Tsiamis et al., 2022) is to use bounds for self-normalized martingales (Abbasi-Yadkori et al., 2011), but the standard self-normalized martingale lemma assumes the noises $\left(e_{t}\right)_{t \geq H}$ to be sub-Gaussian. (Fan et al., 2012, 2017) study the sum of martingale difference sequences with sub-Weibull distributions, based on which we prove Lemma 6.

Lemma 6. Let $\left(\eta_{t}\right)_{t \geq 1}$ be a martingale difference sequence adapted to filtration $\left(\mathcal{F}_{t}\right)_{t \geq 1}$. Assume $\eta_{t} \mid \mathcal{F}_{t-1}$ is $\theta$-sub-Weibull with $\left\|\eta_{t} \mid \mathcal{F}_{t-1}\right\|_{\psi_{\theta}} \leq K$ for $\theta>0$. Then with probability at least $1-p$, there exist constants $a_{\theta}, a_{\theta}^{\prime}>0$, such that as long as $T \geq a_{\theta}$,

$$
\sum_{t=1}^{T} \eta_{t} \leq a_{\theta}^{\prime} K \sqrt{T}(\log (T / p))^{1+\theta^{-1}}
$$

Proof. By the definition of sub-Weibull distributions, $\mathbb{E}\left[\exp \left(\left|\eta_{t} / K\right|^{\theta}\right) \mid \mathcal{F}_{t-1}\right] \leq 2$. Define $\epsilon_{t}=$ $\eta_{t} / K$. Then by the properties of sub-Weibull distributions, there exists constant $a_{\theta}>0$, such that $\mathbb{E}\left[\epsilon_{t}^{2} \mid \mathcal{F}_{t-1}\right] \leq a_{\theta}$ (Zhang and Wei, 2022). Hence, $\left(\epsilon_{t}\right)_{t \geq 1}$ satisfies the assumptions required in (Fan et al., 2017, Theorem 3.2) for $\alpha=\theta /(\theta+1)$. Taking $\left(\phi_{t}\right)_{t \geq 1}$ in (Fan et al., 2017, Theorem 3.2) to be ones, we have for any $x>0$,

$$
\mathbb{P}\left(\sum_{t=1}^{T} \epsilon_{t} \geq x \sqrt{T}\right) \leq \exp \left(-\frac{x^{2}}{2\left(a_{\theta}+x^{1+\frac{1}{\theta+1}} / 3\right)}\right)+2 T \exp \left(-x^{\frac{\theta}{\theta+1}}\right)
$$

Note that

$$
\frac{x^{2}}{2\left(a_{\theta}+x^{1+\frac{1}{\theta+1}} / 3\right)}=\frac{x^{\frac{\theta}{\theta+1}}}{2 / 3+2 a_{\theta} x^{-1-\frac{1}{\theta+1}}} \geq x^{\frac{\theta}{\theta+1}}
$$

if $2 a_{\theta} x^{-1-\frac{1}{\theta+1}} \leq 1 / 3$, that is, $x \geq\left(6 a_{\theta}\right)^{\frac{\theta+1}{\theta+2}}$. Then, as long as $x \geq 6 a_{\theta}$,

$$
\mathbb{P}\left(\sum_{t=1}^{T} \epsilon_{t} \geq x \sqrt{T}\right) \leq 3 T \exp \left(-x^{\frac{\theta}{\theta+1}}\right)
$$

Hence, for any $p \in(0,1)$,

$$
\mathbb{P}\left(\sum_{t=1}^{T} \epsilon_{t} \geq \sqrt{T}(\log (3 T / p))^{1+\theta^{-1}}\right) \leq p
$$

Therefore, there exists a constant $a_{\theta}^{\prime}>0$, such that as long as $T \geq 3 \exp \left(\left(6 a_{\theta}\right)^{\frac{\theta}{\theta+1}}\right)$, with probability at least $1-p$,

$$
\sum_{t=1}^{T} \eta_{t}=\sum_{t=1}^{T} K \epsilon_{t} \leq a_{\theta}^{\prime} K \sqrt{T}(\log (T / p))^{1+\theta^{-1}}
$$

Taking $3 \exp \left(\left(6 a_{\theta}\right)^{\frac{\theta}{\theta+1}}\right)$ as the new $a_{\theta}$ completes the proof.

## C. 4 Proposition on multi-step cumulative costs

The following proposition adapts (Tian et al., 2022, Proposition 3) to our LTI setting, which is important for analyzing CoReL and CoReDyL.

Proposition 3. Define filtration $\mathcal{F}_{t}:=\sigma\left(x_{0}, v_{0}, u_{0}, w_{0}, v_{1}, \ldots, u_{t-1}, w_{t-1}, v_{t}\right)$. Let $\left(z_{t}^{* \prime}\right)_{t \geq 0}$ be the state estimates generated by the Kalman filter under the normalized parameterization, adapted to $\left(\mathcal{F}_{t}\right)_{t \geq 0}$. If we apply $u_{t} \sim \mathcal{N}\left(0, \sigma_{u}^{2} I\right)$, then for any $t \geq 0$,

$$
\bar{c}_{t}:=\sum_{\tau=t}^{t+d_{x}-1}\left(c_{\tau}-\left\|u_{\tau}\right\|_{R^{*}}^{2}\right)=\left\|z_{t}^{* \prime}\right\|^{2}+b^{* \prime}+e_{t}^{\prime}
$$

where $b^{* \prime}=\mathcal{O}\left(d_{x}\right)$ is a positive constant, and $e_{t}^{\prime} \mid \mathcal{F}_{t-1}$ is a zero-mean subexponential random variable with $\left\|e_{t}^{\prime} \mid \mathcal{F}_{t-1}\right\|_{\psi_{1}}=\mathcal{O}\left(d_{x}^{3 / 2}\right)$.

Proof. By Proposition 1, $z_{t+1}^{* \prime}=A^{* \prime} z_{t}^{* \prime}+B^{* \prime} u_{t}+L^{* \prime} i_{t+1}^{\prime}$, where $L^{* \prime}, i_{t+1}^{\prime}$ are the Kalman gain and the innovation under the normalized parameterization, respectively. Under Assumption 1, $\left(i_{t}^{\prime}\right)_{t \geq 1}$ are Gaussian random vectors whose covariances have $\mathcal{O}(1)$ operator norms, and $L^{* \prime}$
has $\mathcal{O}(1)$ operator norm (Zhang and Zhang, 2021). Hence, The covariance of $L^{*} i_{t+1}^{\prime}$ has $\mathcal{O}(1)$ operator norm. Since $u_{t} \sim \mathcal{N}\left(0, \sigma_{u}^{2} I\right), j_{t}:=B^{* \prime} u_{t}+i_{t}^{\prime}$ can be viewed as a Gaussian noise vector independent of $\mathcal{F}_{t-1}$, whose covariance has $\mathcal{O}(1)$ operator norm. By Proposition 1,

$$
c_{t}=\left\|z_{t}^{* \prime}\right\|_{Q^{* \prime}}^{2}+\left\|u_{t}\right\|_{R^{*}}^{2}+b^{*}+e_{t},
$$

where $b^{*}:=\mathbb{E}\left[\left\|z_{t}^{* \prime}-x_{t}^{\prime}\right\|_{Q^{*+}}^{2}\right]$ and $e_{t}:=\gamma_{t}+\eta_{t}=\left\|z_{t}^{* \prime}-x_{t}^{\prime}\right\|_{Q^{* \prime}}^{2}-b^{*}+\left\langle z_{t}^{* \prime}, x_{t}^{\prime}-z_{t}^{* \prime}\right\rangle_{Q^{* \prime}}$. Since the state estimation error $z_{t}^{* \prime}-x_{t}^{\prime}$ is a zero-mean Gaussian random vector and independent of $\mathcal{F}_{t-1}, e_{t} \mid \mathcal{F}_{t-1}$ is a zero-mean subexponential random variable with $\left\|e_{t}\right\|_{\psi_{1}}=\mathcal{O}\left(d_{x}^{1 / 2}\right)$. Then, for $\tau \geq t+1$,

$$
z_{\tau}^{* \prime}=\left(A^{* \prime}\right)^{\tau-t} z_{t}^{* \prime}+\sum_{s=t}^{\tau-1}\left(A^{* \prime}\right)^{\tau-s} j_{s}:=\left(A^{* \prime}\right)^{\tau-t} z_{t}^{* \prime}+j_{\tau, t}^{\prime},
$$

where $j_{\tau, t}^{\prime}$ is a Gaussian random vector independent of $\mathcal{F}_{t-1}$ and $j_{\tau, t}^{\prime} \mid \mathcal{F}_{t-1}$ has zero mean and $\mathcal{O}(1)$ covariance norm due to stability (Assumption 1.1). Therefore,

$$
\begin{aligned}
\bar{c}_{t}= & \sum_{\tau=t}^{t+d_{x}-1} c_{\tau} \\
= & \sum_{\tau=t}^{t+d_{x}-1}\left(\left\|\left(A^{* \prime}\right)^{\tau-t} z_{t}^{* \prime}+j_{\tau, t}^{\prime}\right\|_{Q_{\tau}^{* \prime *}}^{2}+\left\|u_{\tau}\right\|_{R_{\tau}^{*}}^{2}+b^{*}+e_{\tau}\right) \\
= & \left(z_{t}^{* \prime}\right)^{\top}\left(\sum_{\tau=t}^{t+d_{x}-1}\left(\left(A^{* \prime}\right)^{\tau-t}\right)^{\top} Q_{t}^{* \prime}\left(A^{* \prime}\right)^{\tau-t}\right) z_{t}^{* \prime}+\sum_{\tau=t}^{t+d_{x}-1}\left\|u_{\tau}\right\|_{R_{\tau}^{*}}^{2} \\
& \quad+\sum_{\tau=t}^{t+d_{x}-1}\left(\left\|j_{\tau, t}^{\prime}\right\|_{Q_{\tau}^{* \prime}}^{2}+\left(j_{\tau, t}^{\prime}\right)^{\top} Q_{\tau}^{* \prime}\left(A^{* \prime}\right)^{\tau-t} z_{t}^{* \prime}+b^{*}+e_{\tau}\right) \\
= & \left\|z_{t}^{* \prime}\right\|^{2}+\sum_{\tau=t}^{t+d_{x}-1}\left\|u_{\tau}\right\|_{R_{\tau}^{*}}^{2}+b^{* \prime}+e_{t}^{\prime},
\end{aligned}
$$

where $\sum_{\tau=t}^{t+d_{x}-1}\left(\left(A^{* \prime}\right)^{\tau-t}\right)^{\top} Q_{t}^{* \prime}\left(A^{* \prime}\right)^{\tau-t}=I$ is due to the normalized parameterization, $b^{* \prime}:=$ $\sum_{\tau=t}^{t+d_{x}-1}\left(b^{*}+\mathbb{E}\left[\left\|j_{\tau}^{\prime}\right\|_{Q_{\tau}^{* *}}^{2}\right]\right)=\mathcal{O}\left(d_{x}\right)$, and

$$
e_{t}^{\prime}:=\sum_{\tau=t}^{t+d_{x}-1}\left(\left\|j_{\tau}^{\prime}\right\|_{Q_{\tau}^{* \prime}}^{2}-\mathbb{E}\left[\left\|j_{\tau}^{\prime}\right\|_{Q_{\tau}^{* *}}^{2}\right]+\left(j_{\tau}^{\prime}\right)^{\top} Q_{\tau}^{* \prime}\left(A^{* \prime}\right)^{\tau-t} z_{t}^{* \prime}+e_{\tau}\right)
$$

satisfies that $e_{t}^{\prime} \mid \mathcal{F}_{t-1}$ has zero mean and is subexponential with $\left\|e_{t}^{\prime}\right\|_{\psi_{1}}=\mathcal{O}\left(d_{x}^{3 / 2}\right)$.

## D Quadratic regression bound

The following quadratic regression bound is at the core of proving Theorem 2. Its proof builds on a new persistency of excitation result (Lemma 1) and the concentration of sub-Weibull martingale different sequences (Lemma 6). As a caveat, the $h_{t}, c_{t}, \mathcal{F}_{t}$ in Lemma 7 and its proof are a slight abuse of notation, which denote different variables than the rest of the paper. Hence, the indices start with $t=1$, rather than $t=H$ in the CoReL and CoReDyL algorithms.

Lemma 7. Let $\left(h_{t}^{*}\right)_{t \geq 1}$ be a sequence of d-dimensional Gaussian random vectors adapting to filtration $\left(\mathcal{F}_{t}\right)_{t \geq 1}$ with $\left\|\mathbb{E}\left[h_{t}^{*}\left(h_{t}^{*}\right)^{\top}\right]\right\|_{2}^{1 / 2} \leq \sigma$. Define random variable $c_{t}=\left(h_{t}^{*}\right)^{\top} N^{*} h_{t}^{*}+b^{*}+e_{t}$, where $N^{*} \in \mathbb{R}^{d \times d}$ is a positive semidefinite matrix, $b^{*} \in \mathbb{R}$ is a constant, and $e_{t} \mid \mathcal{F}_{t-1}$ is zero-mean, subexponential with $\left\|e_{t} \mid \mathcal{F}_{t-1}\right\|_{\psi_{1}} \leq E$. Assume $\sigma$ and $\left\|N^{*}\right\|_{2}$ are $\mathcal{O}(1)$. Define $h_{t}=h_{t}^{*}+\delta_{t}$, where the noise vector $\delta_{t}$ can be correlated with $h_{t}^{*}$ and its $\ell_{2}$ norm is sub-Gaussian with $\mathbb{E}\left[\left\|\delta_{t}\right\|\right] \leq \epsilon$,
$\left\|\left\|\delta_{t}\right\|_{\psi_{2}} \leq \epsilon\right.$. Assume that $\left(\left[\operatorname{svec}\left(h_{t}^{*}\left(h_{t}^{*}\right)^{\top}\right) ; 1\right]\right)_{t \geq 1}$ satisfies $\left(k, \gamma^{2} I, q\right)-B M S B$ condition and $\epsilon \leq$ $\min \left(\sigma d^{1 / 2}, a_{0} \gamma \sigma^{-1} d^{-1} \log ^{-2}(T / p)\right)$ for some absolute constant $a_{0}>0$. Consider

$$
\begin{equation*}
(\hat{N}, \hat{b}) \in \underset{N=N^{\top}, b}{\operatorname{argmin}} \sum_{t=1}^{T}\left(c_{t}-\left\|h_{t}\right\|_{N}^{2}-b\right)^{2} . \tag{D.1}
\end{equation*}
$$

Then, as long as $T \geq a_{1} k d^{5} \log (d / p)$ for some dimension-free constant $a_{1}>0$, we have that with probability at least $1-p$,

$$
\left\|\hat{N}-N^{*}\right\|_{F}=\mathcal{O}\left(\epsilon(\gamma q)^{-1} d^{1 / 2} \log (T / p)+(\gamma q)^{-2} d E T^{-1 / 2} \log ^{3}(T / p)\right),
$$

where $\sigma, E$ and $\left\|N^{*}\right\|_{2}$ are problem-dependent constants hidden in $\mathcal{O}(\cdot)$.
Proof. Regression (D.1) can be written as

$$
\underset{\operatorname{svec}(N), b}{\operatorname{argmin}} \sum_{t=1}^{T}\left(c_{t}-\operatorname{svec}\left(h_{t} h_{t}^{\top}\right)^{\top} \operatorname{svec}(N)-b\right)^{2} .
$$

Let $f_{t}:=\operatorname{svec}\left(h_{t} h_{t}^{\top}\right)$ denote the covariates and $\bar{f}_{t}:=\left[f_{t} ; 1\right]$ denote the extended covariates. Define $f_{t}^{*}$ and $\bar{f}_{t}^{*}$ similarly by replacing $h_{t}$ with $h_{t}^{*}$. Then, regression (D.1) can be written as

$$
\begin{equation*}
\underset{\operatorname{svec}(N), b}{\operatorname{argmin}} \sum_{t=1}^{T}\left(c_{t}-f_{t}^{\top} \operatorname{svec}(N)-b\right)^{2} . \tag{D.2}
\end{equation*}
$$

Let $\bar{F}:=\left[\bar{f}_{1}, \bar{f}_{2}, \ldots, \bar{f}_{T}\right]^{\top}$ be the $T \times \frac{d(d+1)}{2}$ matrix whose $t$ th row is $f_{t}^{\top}$. Define $\bar{F}^{*}$ similarly by replacing $\bar{f}_{t}$ by $\bar{f}_{t}^{*}$. Solving linear regression (D.2) gives

$$
\bar{F}^{\top} \bar{F}[\operatorname{svec}(\hat{N}) ; \hat{b}]=\sum_{t=1}^{T} \bar{f}_{t} c_{t} .
$$

Substituting $c_{t}=\left(\bar{f}_{t}^{*}\right)^{\top}\left[\operatorname{svec} N^{*} ; b^{*}\right]+e_{t}$ into the above equation yields

$$
\bar{F}^{\top} \bar{F}[\operatorname{svec}(\hat{N}) ; \hat{b}]=\bar{F}^{\top} \bar{F}^{*}\left[\operatorname{svec}\left(N^{*}\right) ; b^{*}\right]+\bar{F}^{\top} \xi
$$

where $\xi$ denotes the vector whose $t$ th element is $e_{t}$. Rearranging the terms, we have

$$
\begin{equation*}
\bar{F}^{\top} \bar{F}\left[\operatorname{svec}\left(\hat{N}-N^{*}\right) ; \hat{b}-b^{*}\right]=\bar{F}^{\top}\left(\bar{F}^{*}-\bar{F}\right)\left[\operatorname{svec}\left(N^{*}\right) ; b^{*}\right]+\bar{F}^{\top} \xi . \tag{D.3}
\end{equation*}
$$

Next, we show that $\lambda_{\text {min }}\left(\bar{F}^{\top} \bar{F}\right)=\Omega\left(\gamma^{2} q^{2} T\right)$, which we achieve by showing $\left(\bar{f}_{t}\right)_{t \geq 1}$ satisfies the BMSB condition. By our assumption, $\left(\bar{f}_{t}^{*}\right)_{t \geq 1}$ satisfies $\left(k, \gamma^{2} I, q\right)$-BMSB condition, meaning that for any fixed unit vector $v \in \mathbb{R}^{\frac{d(d+1)}{2}+1}$, it holds almost surely that

$$
\frac{1}{k} \sum_{i=1}^{k} \mathbb{P}\left(\left|\left\langle\bar{f}_{t+i}^{*}, v\right\rangle\right| \geq \gamma \mid \mathcal{F}_{t}\right) \geq q .
$$

For any fixed unit vector $v \in \mathbb{R}^{\frac{d(d+1)}{2}+1}$, we have

$$
\left|\left\langle\bar{f}_{t}, v\right\rangle\right|=\left|\left\langle\bar{f}_{t}^{*}, v\right\rangle+\left\langle\bar{f}_{t}-\bar{f}_{t}^{*}, v\right\rangle\right| \geq\left|\left\langle\bar{f}_{t}^{*}, v\right\rangle\right|-\left|\left\langle\bar{f}_{t}-\bar{f}_{t}^{*}, v\right\rangle\right| \geq\left|\left\langle\bar{f}_{t}^{*}, v\right\rangle\right|-\left\|\bar{f}_{t}-\bar{f}_{t}^{*}\right\| .
$$

For all $1 \leq t \leq T$, since $\left\|\bar{f}_{t}-\bar{f}_{t}^{*}\right\|=\left\|f_{t}-f_{t}^{*}\right\|=\left\|h_{t} h_{t}^{\top}-h_{t}^{*}\left(h_{t}^{*}\right)^{\top}\right\|_{F}=\mathcal{O}\left(\epsilon \sigma d \log ^{2}(T / p)\right)$, there exists an absolute constant $a_{0}>0$, such that as long as $\epsilon \leq \frac{a_{0} \gamma}{\sigma d \log ^{2}(T / p)},\left\|\bar{f}_{t}-\bar{f}_{t}^{*}\right\| \leq \gamma / 2$. It follows that

$$
\frac{1}{k} \sum_{i=1}^{k} \mathbb{P}\left(\left|\left\langle\bar{f}_{t+i}, v\right\rangle\right| \geq \gamma / 2 \mid \mathcal{F}_{t}\right) \geq \frac{1}{k} \sum_{i=1}^{k} \mathbb{P}\left(\left|\left\langle\bar{f}_{t+i}^{*}, v\right\rangle\right| \geq \gamma \mid \mathcal{F}_{t}\right) \geq q
$$

which means that $\left(\bar{f}_{t}\right)_{1 \leq t \leq T}$ is $\left(k, \gamma^{2} I / 4, q\right)$-BMSB. Following the analysis in (Simchowitz et al., 2018, Appendix D), by lower bounding $\inf _{v:\|v\|=1} \sum_{t=1}^{T}\left\langle v, f_{t}\right\rangle^{2}$ using a covering argument (Simchowitz et al., 2018, Lemma 4.1), we can show that for a given $p \in(0,1)$, as long as $T \geq$ $a_{1} k d^{5} \log (d / p)$ for some dimension-free constant $a_{1}>0$, then with probability at least $1-p$, we have

$$
\lambda_{\min }\left(\sum_{t=1}^{T} \bar{f}_{t} \bar{f}_{t}^{\top}\right)=\Omega\left(\gamma^{2} q^{2} T\right)
$$

Hence, we have $\lambda_{\min }\left(\bar{F}^{\top} \bar{F}\right)=\Omega\left(\gamma^{2} q^{2} T\right)$.
Now we return to (D.3). By inverting $\bar{F}^{\top} \bar{F}$, we obtain

$$
\begin{align*}
\left\|\left[\operatorname{svec}\left(\hat{N}-N^{*}\right) ; \hat{b}-b^{*}\right]\right\| & =\left\|\bar{F}^{\dagger}\left(\bar{F}^{*}-\bar{F}\right)\left[\operatorname{svec}\left(N^{*}\right) ; b^{*}\right]+\bar{F}^{\dagger} \xi\right\| \\
& \leq \underbrace{\left\|\bar{F}^{\dagger}\left(\bar{F}^{*}-\bar{F}\right)\left[\operatorname{svec}\left(N^{*}\right) ; b^{*}\right]\right\|}_{(a)}+\underbrace{\left\|\bar{F}^{\dagger} \tilde{\zeta}\right\|}_{(b)} . \tag{D.4}
\end{align*}
$$

Term (a) is upper bounded by

$$
\begin{aligned}
& \sigma_{\min }(\bar{F})^{-1}\left\|\left(\bar{F}^{*}-\bar{F}\right)\left[\operatorname{svec}\left(N^{*}\right) ; b^{*}\right]\right\| \\
= & \mathcal{O}\left(\sigma_{\min }(\bar{F})^{-1}\right)\left\|\left(F^{*}-F\right) \operatorname{svec}\left(N^{*}\right)\right\| \\
= & \mathcal{O}\left((\gamma q)^{-1} T^{-1 / 2}\right)\left\|\left(F^{*}-F\right) \operatorname{svec}\left(N^{*}\right)\right\| .
\end{aligned}
$$

Using arguments similar to those in (Mhammedi et al., 2020, Section B.2.13), we have

$$
\begin{aligned}
\left\|\left(F^{*}-F\right) \operatorname{svec}\left(N^{*}\right)\right\|^{2} & =\sum_{t=1}^{T}\left\langle\operatorname{svec}\left(h_{t}^{*}\left(h_{t}^{*}\right)^{\top}\right)-\operatorname{svec}\left(h_{t} h_{t}^{\top}\right), \operatorname{svec}\left(N^{*}\right)\right\rangle^{2} \\
& =\sum_{t=1}^{T}\left\langle h_{t}^{*}\left(h_{t}^{*}\right)^{\top}-h_{t} h_{t}^{\top}, N^{*}\right\rangle_{F}^{2} \\
& \leq\left\|N^{*}\right\|_{2}^{2} \sum_{t=1}^{T}\left\|h_{t}^{*}\left(h_{t}^{*}\right)^{\top}-h_{t} h_{t}^{\top}\right\|_{*}^{2} \\
& \stackrel{(i)}{\leq} 4\left\|N^{*}\right\|_{2}^{2} \sum_{t=1}^{T}\left\|h_{t}^{*}\left(h_{t}^{*}\right)^{\top}-h_{t} h_{t}^{\top}\right\|_{2}^{2}
\end{aligned}
$$

where (i) follows from the fact that the matrix $h_{t}^{*}\left(h_{t}^{*}\right)^{\top}-h_{t} h_{t}^{\top}$ has at most rank two. Note that

$$
\left\|h_{t}^{*}\left(h_{t}^{*}\right)^{\top}-h_{t} h_{t}^{\top}\right\|_{2}=\left\|h_{t}^{*}\left(h_{t}^{*}-h_{t}\right)^{\top}+\left(h_{t}^{*}-h_{t}\right) h_{t}^{\top}\right\|_{2} \leq\left(\left\|h_{t}^{*}\right\|+\left\|h_{t}\right\|\right)\|\delta\| .
$$

Since $h_{t}^{*}$ is Gaussian with $\left\|\mathbb{E}\left[h_{t}^{*}\left(h_{t}^{*}\right)^{\top}\right]\right\|_{2}^{1 / 2} \leq \sigma,\left\|h^{*}\right\|$ is sub-Gaussian with its mean and subGaussian norm bounded by $\mathcal{O}\left(\sigma d^{1 / 2}\right)$. Since $\|\delta\|$ is sub-Gaussian with its mean and subGaussian norm bounded by $\epsilon \leq \sigma d^{1 / 2}$, we conclude that $\left\|h^{*}\left(h^{*}\right)^{\top}-h h^{\top}\right\|_{2}$ is subexponential
with its mean and subexponential norm bounded by $\mathcal{O}\left(\epsilon \sigma d^{1 / 2}\right)$. Hence, with probability at least $1-p$,

$$
\left\|h^{*}\left(h^{*}\right)^{\top}-h h^{\top}\right\|_{2}^{2}=\mathcal{O}\left(\epsilon^{2} \sigma^{2} d \| \log ^{2}(T / p)\right)
$$

Therefore, term (a) in (D.4) is bounded by

$$
\mathcal{O}\left((\gamma q)^{-1} T^{-1 / 2} \epsilon \sigma\left\|N^{*}\right\|_{2} d^{1 / 2} T^{1 / 2} \log (T / p)\right)=\mathcal{O}\left((\gamma q)^{-1} d^{1 / 2} \epsilon \log (T / p)\right)
$$

Now we consider term (b) in (D.4):

$$
(b)=\left\|\bar{F}^{\dagger} \xi\right\| \leq \lambda_{\min }\left(\bar{F}^{\top} \bar{F}\right)^{-1}\left\|\bar{F}^{\top} \xi\right\|=\mathcal{O}\left((\gamma q)^{-2} T^{-1}\right)\left\|\sum_{t=1}^{T} \bar{f}_{t} e_{t}\right\| .
$$

Let $\left[\bar{f}_{t}\right]_{i}$ denote the $i$ th element of vector $\bar{f}_{t}$. Since $\left\|\left[\bar{f}_{t}\right]_{i} \mid \mathcal{F}_{t-1}\right\|_{\psi_{1}}=\mathcal{O}(1)$ and $\left\|e_{t} \mid \mathcal{F}_{t-1}\right\|_{\psi_{1}} \leq$ $E$, the product $\left[\bar{f}_{t}\right]_{i} e_{t} \mid \mathcal{F}_{t-1}$ is $\frac{1}{2}$-sub-Weibull, with the sub-Weibull norm being $\mathcal{O}(E)$. By Lemma 6 , with probability at least $1-p$,

$$
\sum_{t=1}^{T}\left[\bar{f}_{t}\right]_{i} e_{t}=\mathcal{O}\left(E T^{1 / 2} \log ^{3}(T / p)\right)
$$

Then, since $\sum_{t=1}^{T} \bar{f}_{t} e_{t}$ has $d(d+1) / 2+1$ components,

$$
\left\|\sum_{t=1}^{T} \bar{f}_{t} e_{t}\right\|=\mathcal{O}\left(d E T^{1 / 2} \log ^{3}(T / p)\right)
$$

Hence, term (b) in (D.4) is bounded by

$$
\mathcal{O}\left((\gamma q)^{-2} d E T^{-1 / 2} \log ^{3}(T / p)\right)
$$

Finally, combining the bounds on $(a)$ and (b), we show that as long as $T \geq a_{1} k d^{5} \log (d / p)$, with probability at least $1-p$,

$$
\left\|\left[\operatorname{svec}\left(\hat{N}-N^{*}\right) ; \hat{b}-b^{*}\right]\right\|=\mathcal{O}\left(\epsilon(\gamma q)^{-1} d^{1 / 2} \log (T / p)+(\gamma q)^{-2} d E T^{-1 / 2} \log ^{3}(T / p)\right) .
$$

## D. 1 Persistency of excitation

Below we prove Lemma 1, which claims that $\left(\left[f_{t} ; 1\right]\right)_{t \geq H}$ satisfies the $\left(k, \gamma^{2} I, q\right)$-BMSB condition. With some additional arguments (Simchowitz et al., 2018; Matni and Tu, 2019), this implies that $\lambda_{\text {min }}\left(\sum_{t=H}^{T+H-1}\left[f_{t} ; 1\right]\left[f_{t}^{\top}, 1\right]\right)=\Omega\left(\gamma^{2} q^{2} T\right)$, establishing the persistency of excitation.

Proof of Lemma 1. Since svec is a bijection, every vector $w \in \mathbb{R}^{d_{h}\left(d_{h}+1\right) / 2}$ corresponds to a symmetric matrix $D \in \mathbb{R}^{d_{h} \times d_{h}}$ with Frobenius norm $\|w\|$. Then, for any unit vector $v=[w ; s]$ with $w \in \mathbb{R}^{d_{h}\left(d_{h}+1\right) / 2}$ and $s \in \mathbb{R}$,

$$
\left\langle\bar{f}_{t+i}, v\right\rangle=\left\langle f_{t+i}, w\right\rangle+s=\left\langle\operatorname{svec}\left(h_{t+i} h_{t+i}^{\top}\right), \operatorname{svec}(D)\right\rangle+s=h_{t+i}^{\top} D h_{t+i}+s
$$

Take $\Gamma=\gamma^{2} I$ for some $\gamma>0$ to be specified later. Then, $\|v\|_{\Gamma}=\gamma$. It suffices to show that for $i>G$ for some $G>0$,

$$
\mathbb{P}\left(\left|h_{t+i}^{\top} D h_{t+i}+s\right| \geq \gamma \mid \mathcal{F}_{t}\right) \geq q,
$$

since if so, we have

$$
\frac{1}{2 G} \sum_{i=1}^{2 G} \mathbb{P}\left(\left|h_{t+i}^{\top} D h_{t+i}+s\right| \geq \gamma \mid \mathcal{F}_{t}\right) \geq \frac{1}{2 G} \sum_{i=G+1}^{2 G} \mathbb{P}\left(\left|h_{t+i}^{\top} D h_{t+i}+s\right| \geq \gamma \mid \mathcal{F}_{t}\right) \geq q / 2
$$

which means $\left(\bar{f}_{t}\right)_{t \geq H}$ is $\left(2 G, \gamma^{2} I, q / 2\right)$-BMSB.
Now let us take a close look at

$$
h_{t+i}=\left[y_{(t+i-H+1):(t+i)} ; u_{(t+i-H):(t+i-1)}\right] .
$$

Since

$$
y_{t+i}=C^{*}\left(A^{*}\right)^{i} x_{t}+\sum_{j=1}^{i} C^{*}\left(A^{*}\right)^{j}\left(B^{*} u_{t+i-j}+w_{t+i-j}\right)+v_{t+i}
$$

$y_{t+i} \mid \mathcal{F}_{t}$ is Gaussian with mean $C^{*}\left(A^{*}\right)^{i} x_{t}$ and covariance determined by $\sum_{j=1}^{i} C^{*}\left(A^{*}\right)^{j}\left(B^{*} u_{t+i-j}+\right.$ $\left.w_{t+i-j}\right)+v_{t+i}$, where we note that $v_{t+i}$ is independent of all other random variables and has full-rank covariance. Hence, for $i>H, h_{t+i} \mid \mathcal{F}_{t}$ is Gaussian and has full-rank covariance. Then intuitively, since $\|D\|_{F}=1,\left|h_{t+i}^{\top} D h_{t+i}\right| \mid \mathcal{F}_{t}$ is a well-behaved random variable that can exceed some $\gamma>0$ with a positive probability $q$. Formally, let $\mu_{t, i}:=\mathbb{E}\left[h_{t+i} \mid \mathcal{F}_{t}\right]$. By Lemma 2 , for $i>H$, there exists some absolute constant $a>0$, such that

$$
\mathbb{E}\left[\left|\left(h_{t+i}-\mu_{t, i}\right)^{\top} D\left(h_{t+i}-\mu_{t, i}\right)+s\right| \mid \mathcal{F}_{t}\right] \geq a \min \left\{\sigma_{u}, \sigma_{v}\right\} d_{h}^{-3 / 2}
$$

By triangle inequality, we have

$$
\begin{aligned}
\left|\left(h_{t+i}-\mu_{t, i}\right)^{\top} D\left(h_{t+i}-\mu_{t, i}\right)+s\right| & =\left|h_{t+i}^{\top} D h_{t+i}+\mu_{t, i}^{\top} D \mu_{t, i}-2 h_{t+i}^{\top} D \mu_{t, i}+s\right| \\
& \leq\left|h_{t+i}^{\top} D h_{t+i}+s\right|+\left|\mu_{t, i}^{\top} D \mu_{t, i}\right|+2\left|h_{t+i}^{\top} D \mu_{t, i}\right| .
\end{aligned}
$$

Hence,

$$
\mathbb{E}\left[\left|h_{t+i}^{\top} D h_{t+i}+s\right| \mid \mathcal{F}_{t}\right] \geq a \min \left\{\sigma_{u}, \sigma_{v}\right\} d_{h}^{-3 / 2}-\mathbb{E}\left[\left|\mu_{t, i}^{\top} D \mu_{t, i}\right|+2\left|h_{t+i}^{\top} D \mu_{t, i}\right| \mid \mathcal{F}_{t}\right] .
$$

Now we argue that for large enough $i, \mathbb{E}\left[\left|\mu_{t, i}^{\top} D \mu_{t, i}\right|+2\left|h_{t+i}^{\top} D \mu_{t, i}\right|\right]$ is negligible. Since matrix $A^{*}$ is stable, with probability at least $1-p,\left\|x_{t}\right\|=\mathcal{O}\left(d_{x}^{1 / 2} \log (T / p)\right)$ for all $t \geq 0$. Hence,

$$
\left\|C^{*}\left(A^{*}\right)^{i} x_{t}\right\|=\mathcal{O}\left(\alpha\left(A^{*}\right) \rho\left(A^{*}\right)^{i} d_{x}^{1 / 2} \log (T / p)\right)
$$

where we recall that $\alpha\left(A^{*}\right):=\sup _{k \geq 0}\left\|\left(A^{*}\right)^{k}\right\|_{2} \rho\left(A^{*}\right)^{-k}$ and $\left\|C^{*}\right\|_{2},\left\|A^{*}\right\|_{2}$ are hidden in $\mathcal{O}(\cdot)$. Then, for $i>H$,

$$
\begin{aligned}
\mathbb{E}\left[\left|\mu_{t, i}^{\top} D \mu_{t, i}\right|+2\left|h_{t+i}^{\top} D \mu_{t, i}\right| \mid \mathcal{F}_{t}\right] & =\left|\left\langle\mu_{t, i} i_{t, i}^{\top} D\right\rangle_{F}\right|+2 \mathbb{E}\left[\left|\left\langle\mu_{t, i} h_{t+i}^{\top}, D\right\rangle_{F}\right| \mid \mathcal{F}_{t}\right] \\
& \leq\left\|\mu_{t, i} \mu_{t, i}^{\top}\right\|_{F} \cdot\|D\|_{F}+2 \mathbb{E}\left[\left\|\mu_{t, i} i_{t+i}^{\top}\right\|_{F} \cdot\|D\|_{F} \mid \mathcal{F}_{t}\right] \\
& =\left\|\mu_{t, i}\right\|^{2}+2\left\|\mu_{t, i}\right\| \cdot \mathbb{E}\left[\left\|h_{t+i}\right\| \mid \mathcal{F}_{t}\right] .
\end{aligned}
$$

By definition, $\mu_{t, i}$ is the concatenation of $\left(C^{*}\left(A^{*}\right)^{j} x_{t}\right)_{i-H+1 \leq j \leq i}$ and zero vectors. Hence,

$$
\left\|\mu_{t, i}\right\|=\mathcal{O}\left(d_{h}^{1 / 2} \alpha\left(A^{*}\right) \rho\left(A^{*}\right)^{i} \log (T / p)\right) .
$$

Choose $H \geq \frac{a_{1} \log \left(d_{h} \alpha\left(A^{*}\right) \log (T / p)\right)}{\log \left(\rho\left(A^{*}\right)^{-1}\right)}$ for some dimension-free constant $a_{1}>0$, such that for $i>2 H$, we have

$$
\left\|\mu_{t, i}\right\|^{2}+2\left\|\mu_{t, i}\right\| \cdot \mathbb{E}\left[\left\|h_{t+i}\right\| \mid \mathcal{F}_{t}\right] \leq a \min \left\{\sigma_{u}, \sigma_{v}\right\} d_{h}^{-3 / 2} / 2
$$

Then, we have the desired lower bound that

$$
\mathbb{E}\left[\left|h_{t+i}^{\top} D h_{t+i}+s\right| \mid \mathcal{F}_{t}\right] \geq a \min \left\{\sigma_{u}, \sigma_{v}\right\} d_{h}^{-3 / 2} / 2
$$

On the other hand, since

$$
\left|h_{t+i}^{\top} D h_{t+i}+s\right|=\left|\left\langle D, h_{t+i} h_{t+i}^{\top}\right\rangle_{F}+s\right| \leq\|D\|_{F}\left\|h_{t+i} h_{t+i}^{\top}\right\|_{F}+|s| \leq h_{t+i}^{\top} h_{t+i}+|s|,
$$

we have $\mathbb{E}\left[\left|h_{t+i}^{\top} D h_{t+i}+s\right|^{2} \mid \mathcal{F}_{t}\right] \leq 2 \mathbb{E}\left[\left\|h_{t+i}\right\|^{4} \mid \mathcal{F}_{t}\right]+2 s^{2}$. Since $\left\|h_{t+i}\right\| \mid \mathcal{F}_{t}$ is sub-Gaussian with

$$
\left\|\left\|h_{t+i}\right\| \mid \mathcal{F}_{t}\right\|_{\psi_{2}}=\mathcal{O}\left(\left\|\mathbb{E}\left[h_{t+i} h_{t+i}^{\top} \mid \mathcal{F}_{t}\right]\right\|_{2}^{1 / 2}\right)=\mathcal{O}(1),
$$

$\mathbb{E}\left[\left|h_{t+i}^{\top} D h_{t+i}+s\right|^{2} \mid \mathcal{F}_{t}\right]=\mathcal{O}(1)$. By the Paley-Zygmund inequality, for $\beta \in[0,1]$ we have

$$
\mathbb{P}\left(\left|h_{t+i}^{\top} D h_{t+i}+s\right| \geq \beta a \min \left\{\sigma_{u}, \sigma_{v}\right\} d_{h}^{-3 / 2} / 2 \mid \mathcal{F}_{t}\right)=\Omega\left((1-\beta)^{2} a^{2} d_{h}^{-3}\right),
$$

where the dependence on $\sigma_{u}, \sigma_{v}$ is hidden in $\Omega(\cdot)$. By taking $\beta=1 / 2$, we can see that $\left(f_{t}\right)_{t>H}$ satisfies the $\left(k, \gamma^{2} I, q\right)$-BMSB condition for $k=4 H, \gamma=\Theta\left(d_{h}^{-3 / 2}\right)$ and $q=\Theta\left(d_{h}^{-3}\right)$.

## E Linear system identification with noisy measurements

Identifying the time-invariant latent dynamics involves linear regression with correlated data and noisy measurements. The following Lemma 8 extends the previous linear system identification result in (Simchowitz et al., 2018) to the case with noises in both input and output variables. In Lemma $8, \gamma$ and $q$ are treated as dimension-free constants (in contrast to Lemma 7), which is indeed the case in our application of Lemma 8 to $\left(z_{t}^{*}\right)_{t \geq H}$ in analyzing SysId (3.3) for CoReL in §F.1, and in analyzing alignment matrix estimation (3.7) in Algorithm 2 for CoReDyL in §F.2. Note that the bound in Lemma 8 is worse than that in the time-varying setting in (Tian et al., 2022), due to the treatment of correlated data.

Lemma 8. Let $\left(x_{t}^{*}\right)_{t \geq 1}$ be a sequence of $d_{1}$-dimensional Gaussian random vectors adapted to a filtration $\left(\mathcal{F}_{t}\right)_{t \geq 1}$ with $\left\|\mathbb{E}\left[x_{t}^{*}\left(x_{t}^{*}\right)^{\top}\right]\right\|_{2}^{1 / 2} \leq \sigma$. Define $y_{t}^{*}=A^{*} x_{t}^{*}+e_{t}$, where $A^{*} \in \mathbb{R}^{d_{2} \times d_{1}}$ and $e_{t} \mid \mathcal{F}_{t}$ is Gaussian with zero mean and $\left\|\mathbb{E}\left[e_{t} e_{t}^{\top}\right]\right\|_{2}^{1 / 2} \leq \epsilon$. Define $y_{t}=y_{t}^{*}+\delta_{t}^{y}$ and $x_{t}=x_{t}^{*}+\delta_{t}^{x}$, where the noise vectors $\delta_{t}^{x}$ and $\delta_{t}^{y}$ can be correlated with $x_{t}^{*}$ and $y_{t}^{*}$, and their $\ell_{2}$ norms are sub-Gaussian with $\mathbb{E}\left[\left\|\delta_{t}^{x}\right\|\right] \leq \epsilon_{x},\| \| \delta_{t}^{x}\| \|_{\psi_{2}} \leq \epsilon_{x}$ and $\mathbb{E}\left[\left\|\delta_{t}^{y}\right\|\right] \leq \epsilon_{y},\| \| \delta_{t}^{y}\| \|_{\psi_{2}} \leq \epsilon_{y}$. Assume that $\left(x_{t}^{*}\right)_{t \geq 1}$ satisfies the
$\left(k, \gamma^{2} I, q\right)$-BMSB condition, $\epsilon_{x} \leq a_{0} \gamma^{2} q^{2} / \sigma$ for some absolute constant $a_{0}, a_{1}>0$ and both $\epsilon_{x}$ and $\epsilon_{y}$ are $\mathcal{O}(1)$. Consider

$$
\begin{equation*}
\hat{A} \in \underset{A \in \mathbb{R}^{d_{2} \times d_{1}}}{\operatorname{argmin}} \sum_{t=1}^{T}\left\|y_{t}-A x_{t}\right\|^{2} . \tag{E.1}
\end{equation*}
$$

Then, as long as $T \geq a_{1} k q^{-1}\left(\log (1 / p)+d_{1} \log (10 / q)+d_{1} \log \left(\sigma \gamma^{-1} d_{1} \log (T / p)\right)\right)$ for some absolute constant $a_{1}>0$, we have that with probability at least $1-p$,

$$
\left\|\hat{A}-A^{*}\right\|_{2}=\mathcal{O}\left(\left(\epsilon_{x}+\epsilon_{y}\right) d_{1}^{1 / 2} \log (T / p)+\left(d_{2}+d_{1} \log \left(d_{1} \log (T / p)\right)+\log (1 / p)\right)^{1 / 2} T^{-1 / 2}\right) .
$$

Proof. Let $X \in \mathbb{R}^{T \times d_{1}}$ denote the matrix whose $t$ th row is $x_{t}^{\top}$. Define $X^{*}, Y, E, \Delta_{x}, \Delta_{y}$ similarly. To solve the regression problem, we set its gradient to be zero and substitute in $Y=X^{*}\left(A^{*}\right)^{\top}+$ $E+\Delta_{y}$ to obtain

$$
\begin{equation*}
\hat{A}\left(X^{\top} X\right)=A^{*}\left(X^{*}\right)^{\top} X+E^{\top} X+\Delta_{y}^{\top} X \tag{E.2}
\end{equation*}
$$

Substituting in $X=X^{*}+\Delta_{x}$ gives

$$
\begin{gather*}
\left(\hat{A}-A^{*}\right)\left(\left(X^{*}\right)^{\top} X^{*}\right)=A^{*}\left(X^{*}\right)^{\top} \Delta_{x}-\hat{A}\left(\Delta_{x}^{\top} \Delta_{x}+\Delta_{x}^{\top} X^{*}+\left(X^{*}\right)^{\top} \Delta_{x}\right)  \tag{E.3}\\
+E^{\top} X^{*}+E^{\top} \Delta_{x}+\Delta_{y}^{\top} X^{*}+\Delta_{y}^{\top} \Delta_{x} .
\end{gather*}
$$

Now we deal with each term on the right-hand side. Since $\left(X^{*}\right)^{\top} \Delta_{x}=\sum_{t=1}^{T} x_{t}^{*}\left(\delta_{t}^{x}\right)^{\top}$, by the triangle inequality,

$$
\left\|\left(X^{*}\right)^{\top} \Delta_{x}\right\|_{2} \leq \sum_{t=1}^{T}\left\|x_{t}^{*}\left(\delta_{t}^{x}\right)^{\top}\right\|_{2} \leq \sum_{t=1}^{T}\left\|x_{t}^{*}\right\| \cdot\left\|\delta_{t}^{x}\right\| .
$$

Since $\left(x_{t}^{*}\right)_{t \geq 1}$ are Gaussian, with probability at least $1-p,\left\|x_{t}^{*}\right\|=\mathcal{O}\left(\sigma d_{1}^{1 / 2} \log ^{1 / 2}(T / p)\right)$. Since $\left(\left\|\delta_{t}^{x}\right\|\right)_{t \geq 1}$ are sub-Gaussian with $\mathbb{E}\left[\left\|\delta_{t}^{x}\right\|\right] \leq \epsilon_{x}$ and $\left\|\left\|\delta_{t}^{x}\right\|\right\|_{\psi_{2}} \leq \epsilon_{x}$, with probability at least $1-p,\left\|\delta_{t}^{x}\right\|=\mathcal{O}\left(\epsilon_{x} \log ^{1 / 2}(T / p)\right)$. Hence,

$$
\left\|\left(X^{*}\right)^{\top} \Delta_{x}\right\|_{2}=\mathcal{O}\left(\epsilon_{x} \sigma d_{1}^{1 / 2} T \log (T / p)\right) .
$$

Similarly, with probability at least $1-p$,

$$
\begin{gathered}
\left\|\Delta_{x}^{\top} \Delta_{x}\right\|_{2}=\mathcal{O}\left(\epsilon_{x}^{2} T \log (T / p)\right), \quad\left\|E^{\top} \Delta_{x}\right\|_{2}=\mathcal{O}\left(\epsilon \epsilon_{x} d_{2}^{1 / 2} T \log (T / p)\right), \\
\left\|\Delta_{y}^{\top} X^{*}\right\|_{2}=\mathcal{O}\left(\epsilon_{y} \sigma d_{1}^{1 / 2} T \log (T / p)\right), \quad\left\|\Delta_{y}^{\top} \Delta_{x}\right\|_{2}=\mathcal{O}\left(\epsilon_{x} \epsilon_{y} T \log (T / p)\right) .
\end{gathered}
$$

It remains to bound $\|\hat{A}\|_{2}$. Notice that with probability at least $1-p$,

$$
\left\|\left(X^{*}\right)^{\top} X^{*}\right\|_{2} \leq \sum_{t=1}^{T}\left\|x_{t}^{*}\right\|^{2}=\mathcal{O}\left(\sigma^{2} d_{1} \log (T / p) T\right)
$$

Let $T_{0}:=a_{1} k q^{-1}\left(\log (1 / p)+d_{1} \log (10 / q)+d_{1} \log \left(\sigma \gamma^{-1} d_{1} \log (T / p)\right)\right)$ for some absolute constant $a_{1}>0$. Then, by (Simchowitz et al., 2018, Appendix D), as long as $T \geq T_{0}$, with probability at least $1-p, \lambda_{\min }\left(\left(X^{*}\right)^{\top} X^{*}\right)=\gamma^{2} q^{2} T / 32$. Since $X^{\top} X=\left(X^{*}\right)^{\top} X^{*}+\Delta_{x}^{\top} X^{*}+\left(X^{*}\right)^{\top} \Delta_{x}+\Delta_{x}^{\top} \Delta_{x}$,

$$
\lambda_{\min }\left(X^{\top} X\right) \geq \lambda_{\min }\left(\left(X^{*}\right)^{\top} X^{*}\right)-\left\|\Delta_{x}^{\top} X^{*}+\left(X^{*}\right)^{\top} \Delta_{x}+\Delta_{x}^{\top} \Delta_{x}\right\|_{2} .
$$

Hence, there exists an absolute constant $a_{0}>0$, such that as long as $\epsilon_{x} \leq a_{0} \gamma^{2} q^{2} / \sigma, \lambda_{\min }\left(X^{\top} X\right)=$ $\Omega\left(\gamma^{2} q^{2} T\right)$, which implies $\left\|X^{\dagger}\right\|_{2}=\mathcal{O}\left(\gamma^{-1} q^{-1} T^{-1 / 2}\right)$. From (E.2), we have

$$
\|\hat{A}\|_{2}=\left(\left\|A^{*}\right\|_{2}\left\|X^{*}\right\|_{2}+\|E\|_{2}+\left\|\Delta_{y}\right\|_{2}\right)\left\|X^{\dagger}\right\|_{2}=\mathcal{O}\left(\gamma^{-1} q^{-1}\left(\sigma\left\|A^{*}\right\|_{2}+\epsilon+\epsilon_{y}\right)\right)=\mathcal{O}(1),
$$

where in the last equality we treat $\gamma, q, \sigma, \epsilon, \epsilon_{y}$ as problem-dependent constants.
Finally, by (Simchowitz et al., 2018, Theorem 2.4), as long as $T \geq T_{0}$,

$$
\left\|E^{\top}\left(X^{*}\right)^{\dagger}\right\|_{2}=\mathcal{O}\left(\left(d_{2}+d_{1} \log \left(d_{1} \log (T / p)\right)+\log (1 / p)\right)^{1 / 2} T^{-1 / 2}\right)
$$

Combining all the above individual bounds with the terms on the right-hand side of (E.3), we have

$$
\left\|\hat{A}-A^{*}\right\|_{2}=\mathcal{O}\left(\left(\epsilon_{x}+\epsilon_{y}\right) d_{1}^{1 / 2} \log (T / p)+\left(d_{2}+d_{1} \log \left(d_{1} \log (T / p)\right)+\log (1 / p)\right)^{1 / 2} T^{-1 / 2}\right)
$$

which completes the proof.

## F Proof of the main results

In this section, we prove the sample complexity bounds for CoReL and CoReDyL in Theorem 1, which is a simplified version of the following Theorem 2. As we shall see, the proofs for the two algorithms share similar ideas and tools.

Theorem 2. Given an unknown LQG problem satisfying Assumption 1, let $M^{* \prime}$ and $\left(A^{* \prime}, B^{* \prime}, Q^{* \prime}, R^{*}\right)$ be the optimal state representation function and the true system parameters under the normalized parameterization. Recall that for a square matrix $A$, we define $\alpha(A):=\sup _{k>0}\left\|A^{k}\right\|_{2} \rho(A)^{-k}$. For a given $p \in(0,1)$, if we run CoReL or CoReDyL for $T \geq \operatorname{poly}\left(d_{x}, d_{y}, d_{u}, \log (T / p)\right), H=$ $\Omega\left(\frac{\log \left(\alpha\left(\bar{A}^{*}\right) T\right)}{\log \left(\rho\left(\bar{A}^{*}\right)^{-1}\right)}+\frac{\log \left(\alpha\left(A^{*}\right) H\left(d_{y}+d_{u}\right) \log (T / p)\right)}{\log \left(\rho\left(A^{*}\right)^{-1}\right)}\right)$, and $\sigma_{u}=\Theta(1)$, then there exists an orthonormal matrix $S \in \mathbb{R}^{d_{x} \times d_{x}}$, such that with probability at least $1-p$, the representation function $\hat{M}$ satisfies

$$
\left\|\hat{M}-S M^{* \prime}\right\|_{2}=\mathcal{O}\left(\left(H\left(d_{y}+d_{u}\right) d_{x}\right)^{21 / 2} T^{-1 / 2} \log ^{4}(T / p)\right),
$$

and feedback gain $\hat{K}$ satisfies

$$
\begin{aligned}
& J^{\hat{K}}\left(S A^{* \prime} S^{\top}, S B^{* \prime}, S Q^{* \prime} S^{\top}, R^{*}\right)-J^{*}\left(S A^{* \prime} S^{\top}, S B^{* \prime}, S Q^{* \prime} S^{\top}, R^{*}\right) \\
= & \mathcal{O}\left(\left(H\left(d_{y}+d_{u}\right) d_{x}\right)^{21}\left(d_{x} \wedge d_{u}\right) T^{-1} \log ^{8}(T / p)\right) .
\end{aligned}
$$

## F. 1 Proof of Theorem 2 for CoReL

Proof. Define filtration

$$
\mathcal{F}_{t}:=\sigma\left(x_{0}, v_{0}, u_{0}, w_{0}, v_{1}, \ldots, u_{t-1}, w_{t-1}, v_{t}\right)
$$

By definition, $x_{t}, y_{t} \in \mathcal{F}_{t}$. Hence, $h_{t}=\left[y_{(t-H+1): t} ; u_{(t-H):(t-1)}\right] \in \mathcal{F}_{t}$ and $z_{t}^{*} \in \mathcal{F}_{t}$.

Recall that we define $f_{t}:=\operatorname{svec}\left(h_{t} h_{t}^{\top}\right)$ and $d_{h}=H\left(d_{y}+d_{u}\right)$. By Lemma 1, $f_{t}$ is $\left(k, \gamma^{2} I, q\right)$ BMSB for $k=4 H, \gamma=\Theta\left(d_{h}^{-3 / 2}\right)$ and $q=\Theta\left(d_{h}^{-3}\right)$. By Proposition 3 and Lemma 7, there exists some absolute constant $a_{0}>0$, such that as long as $T \geq a_{0} H^{6}\left(d_{y}+d_{u}\right)^{5} \log \left(H\left(d_{y}+d_{u}\right) / p\right)$, with probability at least $1-p$,

$$
\left\|\hat{N}-N^{*}\right\|_{F}=\mathcal{O}\left((\gamma q)^{-2} d_{h} E T^{-1 / 2} \log ^{3}(T / p)\right)=\mathcal{O}\left(\left(H\left(d_{y}+d_{u}\right)\right)^{10} d_{x}^{3 / 2} T^{-1 / 2} \log ^{3}(T / p)\right)
$$

By (Tu et al., 2016, Lemma 5.4), there exists an orthogonal matrix $S$, such that $\left\|\hat{M}-S M^{*}\right\|_{F}$ is of the same order of $\left\|\hat{N}-N^{*}\right\|_{F}$.

Recall that $z_{t}^{*}=M^{*} h_{t}+\delta_{t}$, where $\delta_{t}=\left(\bar{A}^{*}\right)^{H} z_{t-H}^{*}$. Then,

$$
\left\|\hat{z}_{t}-S z_{t}^{*}\right\|=\left\|\left(\hat{M}-S M^{*}\right) h_{t}-S \delta_{t}\right\| \leq\left\|\hat{M}-S M^{*}\right\|_{2}\left\|h_{t}\right\|+\left\|\delta_{t}\right\| .
$$

Since $\left\|h_{t}\right\|$ is sub-Gaussian with $\mathbb{E}\left[\left\|h_{t}\right\|\right]=\mathcal{O}\left(d_{h}^{1 / 2}\right),\| \| h_{t}\| \|_{\psi_{2}}=\mathcal{O}\left(d_{h}^{1 / 2}\right)$, we have $\| \hat{M}-$ $S M^{*}\left\|_{2}\right\| h_{t} \|$ is sub-Gaussian with its mean and sub-Gaussian norm bounded by

$$
\begin{equation*}
\mathcal{O}\left(\left(H\left(d_{y}+d_{u}\right)\right)^{21 / 2} d_{x}^{3 / 2} T^{-1 / 2} \log ^{3}(T / p)\right) . \tag{F.1}
\end{equation*}
$$

Notice that $\left\|\delta_{t}\right\|$ is sub-Gaussian with mean and sub-Gaussian norm bounded by

$$
\mathcal{O}\left(\alpha\left(\bar{A}^{*}\right) \rho\left(\bar{A}^{*}\right)^{H} d_{x}^{1 / 2}\right)
$$

which, by our choice of $H$, is dominated by (F.1). Hence, for all $t \geq H,\left\|\hat{z}_{t}-S z_{t}^{*}\right\|$ is subGaussian with its mean and sub-Gaussian norm bounded by (F.1).

Then, by the noisy linear regression bound (Lemma 8), for $T$ greater than a constant polynomial in the problem parameters, we have

$$
\begin{aligned}
\left\|[\hat{A}, \hat{B}]-S\left[A^{*}, B^{*}\right]\right\|_{2} & =\mathcal{O}\left(\left(H\left(d_{y}+d_{u}\right)\right)^{21 / 2} d_{x}^{3 / 2} T^{-1 / 2} \log ^{3}(T / p) \cdot d_{x}^{1 / 2} \log (T / p)\right) \\
& =\mathcal{O}\left(\left(H\left(d_{y}+d_{u}\right)\right)^{21 / 2} d_{x}^{2} T^{-1 / 2} \log ^{4}(T / p)\right) .
\end{aligned}
$$

By the noisy quadratic regression bound (Lemma 7) and $\operatorname{svec}\left(z_{t}^{*}\left(z_{t}^{*}\right)^{\top}\right)$ being $\left(\Theta(1), \Theta\left(d_{x}^{-3 / 2}\right), \Theta\left(d_{x}^{-3}\right)\right)$ BMSB,

$$
\begin{aligned}
\left\|\hat{Q}-S Q^{*} S^{\top}\right\|_{F}= & \mathcal{O}\left(\left(H\left(d_{y}+d_{u}\right)\right)^{21 / 2} d_{x}^{3 / 2} T^{-1 / 2} \log ^{3}(T / p) \cdot d_{x}^{9 / 2} d_{x}^{1 / 2} \log (T / p)\right. \\
& \left.+d_{x}^{9} d_{x} d_{x}^{1 / 2} T^{-1 / 2} \log ^{3}(T / p)\right) \\
= & \mathcal{O}\left(\left(H\left(d_{y}+d_{u}\right) d_{x}\right)^{21 / 2} T^{-1 / 2} \log ^{4}(T / p)\right)
\end{aligned}
$$

Hence, $\left\|\hat{A}-S A^{*} S^{\top}\right\|_{2},\left\|\hat{B}-S B^{*}\right\|_{2}$ and $\left\|\hat{Q}-S Q^{*} S^{\top}\right\|_{2}$ are all bounded by

$$
\mathcal{O}\left(\left(H\left(d_{y}+d_{u}\right) d_{x}\right)^{21 / 2} T^{-1 / 2} \log ^{4}(T / p)\right)
$$

By (Mania et al., 2019), for $T$ greater than a constant that is polynomial in the problem parameters,

$$
\left\|\hat{K}-K^{*} S^{\top}\right\|_{2}=\mathcal{O}\left(\left(H\left(d_{y}+d_{u}\right) d_{x}\right)^{21 / 2} T^{-1 / 2} \log ^{4}(T / p)\right)
$$

is of the same order, and $\hat{K}$ satisfies

$$
\begin{aligned}
& J^{\hat{K}}\left(S A^{* \prime} S^{\top}, S B^{* \prime}, S Q^{* \prime} S^{\top}, R^{*}\right)-J^{*}\left(S A^{* \prime} S^{\top}, S B^{* \prime}, S Q^{* \prime} S^{\top}, R^{*}\right) \\
= & \mathcal{O}\left(\left(H\left(d_{y}+d_{u}\right) d_{x}\right)^{21}\left(d_{x} \wedge d_{u}\right) T^{-1} \log ^{8}\left(d_{x} T / p\right)\right) .
\end{aligned}
$$

## F. 2 Proof of Theorem 2 for CoReDyL

Proof. Define filtration

$$
\mathcal{F}_{t}:=\sigma\left(x_{0}, v_{0}, u_{0}, w_{0}, v_{1}, \ldots, u_{t-1}, w_{t-1}, v_{t}\right)
$$

By definition, $x_{t}, y_{t} \in \mathcal{F}_{t}$. Hence, $h_{t}=\left[y_{(t-H+1): t} ; u_{(t-H):(t-1)}\right] \in \mathcal{F}_{t}$ and $z_{t}^{*} \in \mathcal{F}_{t}$.
Recall that we define $f_{t}:=\operatorname{svec}\left(h_{t} h_{t}^{\top}\right)$. By Lemma $1, f_{t}$ is $\left(k, \gamma^{2} I, q\right)$-BMSB for $k=4 H$, $\gamma=\Theta\left(d_{h}^{-3 / 2}\right)$ and $q=\Theta\left(d_{h}^{-3}\right)$. By Proposition 3 and Lemma 7, there exists some absolute constant $a_{0}>0$, such that as long as $T \geq a_{0} H^{6}\left(d_{y}+d_{u}\right)^{5} \log \left(H\left(d_{y}+d_{u}\right) / p\right)$, with probability at least $1-p$,

$$
\begin{gathered}
\left\|\hat{N}-N^{*}\right\|_{F}=\mathcal{O}\left(\left(H\left(d_{y}+d_{u}\right)\right)^{10} d_{x}^{3 / 2} T^{-1 / 2} \log ^{3}(T / p)\right), \\
\left\|\hat{N}_{1}-N_{1}^{*}\right\|_{F}=\mathcal{O}\left(\left(H\left(d_{y}+d_{u}\right)+d_{x}\right)^{10} d_{x}^{3 / 2} T^{-1 / 2} \log ^{3}(T / p)\right) .
\end{gathered}
$$

By (Tu et al., 2016, Lemma 5.4), there exists orthogonal matrices $S, S_{1}$, such that $\left\|\hat{M}-S M^{*}\right\|_{F}$ and $\left\|\hat{M}_{1}-S_{1} M_{1}^{*}\right\|_{F}$ are on the same order of $\left\|\hat{N}-N^{*}\right\|_{F}$ and $\left\|\hat{N}_{1}-N_{1}^{*}\right\|_{F}$, respectively. Then,

$$
\left\|\widetilde{B}-S_{1} B^{*}\right\|_{2}=\mathcal{O}\left(\left(H\left(d_{y}+d_{u}\right)+d_{x}\right)^{10} d_{x}^{3 / 2} T^{-1 / 2} \log ^{3}(T / p)\right) .
$$

Since $\left[S_{1} A^{*} S^{\top} S M^{*}, S_{1} B^{*}\right]=S_{1} M_{1}^{*}$, by the perturbation bounds of the Moore-Penrose inverse (Wedin, 1973), we have

$$
\left\|\widetilde{A}-S_{1} A^{*} S^{\top}\right\|_{2}=\mathcal{O}\left(\left(H\left(d_{y}+d_{u}\right)+d_{x}\right)^{10} d_{x}^{3 / 2} T^{-1 / 2} \log ^{3}(T / p)\right)
$$

To align $\widetilde{A}$ with $S A^{*} S^{\top}$, we compute another matrix $\hat{S}_{0}$ by solving the regression (3.7) from $\hat{M}_{1}\left[h_{t} ; u_{t}\right]$ to $\hat{M} h_{t+1}$. Since $\hat{M_{1}}\left[h_{t} ; u_{t}\right]$ and $\hat{M} h_{t+1}$ approximate $S_{1} z_{t+1}^{*}$ and $S z_{t+1}^{*}$, respectively, (3.7) is essentially a linear regression that estimates the alignment matrix $S S_{1}^{\top}$ with noisy variables $\hat{M}_{1}\left[h_{t} ; u_{t}\right]$ and $\hat{M} h_{t+1}$. Below we understand the level of the noise terms.

Recall that $z_{t}^{*}=M^{*} h_{t}+\delta_{t}$, where $\delta_{t}=\left(\bar{A}^{*}\right)^{H} z_{t-H}^{*}$. Then,

$$
\left\|\hat{z}_{t}-S z_{t}^{*}\right\|=\left\|\left(\hat{M}-S M^{*}\right) h_{t}-S \delta_{t}\right\| \leq\left\|\hat{M}-S M^{*}\right\|_{2}\left\|h_{t}\right\|+\left\|\delta_{t}\right\| .
$$

Since $\left\|h_{t}\right\|$ is sub-Gaussian with $\mathbb{E}\left[\left\|h_{t}\right\|\right]=\mathcal{O}\left(d_{h}^{1 / 2}\right),\| \| h_{t}\| \|_{\psi_{2}}=\mathcal{O}\left(d_{h}^{1 / 2}\right)$, we have $\| \hat{M}-$ $S M^{*}\left\|_{2}\right\| h_{t} \|$ is sub-Gaussian with its mean and sub-Gaussian norm bounded by

$$
\begin{equation*}
\mathcal{O}\left(\left(H\left(d_{y}+d_{u}\right)\right)^{21 / 2} d_{x}^{3 / 2} T^{-1 / 2} \log ^{3}(T / p)\right) . \tag{F.2}
\end{equation*}
$$

Notice that $\left\|\delta_{t}\right\|$ is sub-Gaussian with mean and sub-Gaussian norm bounded by

$$
\mathcal{O}\left(\alpha\left(\bar{A}^{*}\right) \rho\left(\bar{A}^{*}\right)^{H} d_{x}^{1 / 2}\right),
$$

which, by our choice of $H$, is dominated by (F.2). Hence, for all $t \geq H,\left\|\hat{z}_{t}-S z_{t}^{*}\right\|$ is subGaussian with its mean and sub-Gaussian norm of the order in (F.2). Similarly, $\| \hat{M}_{1}\left[h_{t} ; u_{t}\right]-$ $S_{1} z_{t+1}^{*} \|$ is sub-Gaussian with its mean and sub-Gaussian norm bounded by (F.2).

Hence, by the noisy linear regression bound (Lemma 8), for $T$ greater than a constant polynomial in the problem parameters, we have

$$
\begin{aligned}
\left\|\hat{S}_{0}-S S_{1}^{\top}\right\|_{2} & =\mathcal{O}\left(\left(H\left(d_{y}+d_{u}\right)+d_{x}\right)^{21 / 2} d_{x}^{3 / 2} T^{-1 / 2} \log ^{3}(T / p) \cdot d_{x}^{1 / 2} \log (T / p)\right) \\
& =\mathcal{O}\left(\left(H\left(d_{y}+d_{u}\right)+d_{x}\right)^{21 / 2} d_{x}^{2} T^{-1 / 2} \log ^{4}(T / p)\right)
\end{aligned}
$$

Hence,

$$
\begin{aligned}
\left\|\hat{A}-S A^{*} S^{\top}\right\|_{2} & =\left\|\hat{S}_{0} \widetilde{A}-S S_{1}^{\top} S_{1} A^{*} S^{\top}\right\|_{2} \\
& =\left\|\left(\hat{S}_{0}-S S_{1}^{\top}\right) \widetilde{A}\right\|_{2}+\left\|S S_{1}^{\top}\left(\widetilde{A}-S_{1} A^{*} S^{\top}\right)\right\|_{2} \\
& =\mathcal{O}\left(\left(H\left(d_{y}+d_{u}\right)+d_{x}\right)^{21 / 2} d_{x}^{2} T^{-1 / 2} \log ^{4}(T / p)\right),
\end{aligned}
$$

and $\left\|\hat{B}-S B^{*}\right\|_{2}$ has the same order. By the noisy quadratic regression bound (Lemma 7) and $\operatorname{svec}\left(z_{t}^{*}\left(z_{t}^{*}\right)^{\top}\right)$ being $\left(\Theta(1), \Theta\left(d_{x}^{-3 / 2}\right), \Theta\left(d_{x}^{-3}\right)\right)$-BMSB,

$$
\begin{aligned}
& \left\|\hat{Q}-S Q^{*} S^{\top}\right\|_{F} \\
= & \mathcal{O}\left(\left(H\left(d_{y}+d_{u}\right)\right)^{21 / 2} d_{x}^{3 / 2} T^{-1 / 2} \log ^{3}(T / p) \cdot d_{x}^{9 / 2} d_{x}^{1 / 2} \log (T / p)+d_{x}^{9} d_{x} d_{x}^{1 / 2} T^{-1 / 2} \log ^{4}(T / p)\right) \\
= & \mathcal{O}\left(\left(H\left(d_{y}+d_{u}\right) d_{x}\right)^{21 / 2} T^{-1 / 2} \log ^{4}(T / p)\right)
\end{aligned}
$$

Hence, $\left\|\hat{A}-S A^{*} S^{\top}\right\|_{2},\left\|\hat{B}-S B^{*}\right\|_{2}$ and $\left\|\hat{Q}-S Q^{*} S^{\top}\right\|_{2}$ are all bounded by

$$
\mathcal{O}\left(\left(H\left(d_{y}+d_{u}\right) d_{x}\right)^{21 / 2} T^{-1 / 2} \log ^{4}(T / p)\right)
$$

By (Mania et al., 2019), for $T$ greater than a constant polynomial in the problem parameters,

$$
\left\|\hat{K}-K^{*} S^{\top}\right\|_{2}=\mathcal{O}\left(\left(H\left(d_{y}+d_{u}\right) d_{x}\right)^{21 / 2} T^{-1 / 2} \log ^{4}(T / p)\right)
$$

is of the same order, and $\hat{K}$ satisfies

$$
\begin{aligned}
& J^{\hat{K}}\left(S A^{* \prime} S^{\top}, S B^{* \prime}, S Q^{* \prime} S^{\top}, R^{*}\right)-J^{*}\left(S A^{* \prime} S^{\top}, S B^{* \prime}, S Q^{* \prime} S^{\top}, R^{*}\right) \\
= & \mathcal{O}\left(\left(H\left(d_{y}+d_{u}\right) d_{x}\right)^{21}\left(d_{x} \wedge d_{u}\right) T^{-1} \log ^{8}\left(d_{x} T / p\right)\right),
\end{aligned}
$$

which completes the proof.

