Lawrence Berkeley National Laboratory
Recent Work

Title
LQ-optimal sampled-data control under stochastic delays: Gridding approach for stabilizability and detectability

Permalink
https://escholarship.org/uc/item/7k97b5j3

Journal
SIAM Journal on Control and Optimization, 56(4)

ISSN
0363-0129

Authors
Wakaiki, M
Ogura, M
Hespanha, JP

Publication Date
2018

DOI
10.1137/17M1150608

Peer reviewed
LQ-OPTIMAL SAMPLE-DATA CONTROL UNDER STOCHASTIC DELAYS: GRIDDING APPROACH FOR STABILIZABILITY AND DETECTABILITY∗

MASASHI WAKAIKI†, MASAKI OGURA‡, AND JOÃO P. HESPANHA§

Abstract. We solve a linear quadratic optimal control problem for sampled-data systems with stochastic delays. The delays are stochastically determined by the last few delays. The proposed optimal controller can be efficiently computed by iteratively solving a Riccati difference equation, provided that a discrete-time Markov jump system equivalent to the sampled-data system is stochastic stabilizable and detectable. Sufficient conditions for these notions are provided in the form of linear matrix inequalities, from which stabilizing controllers and state observers can be constructed.

Key words. Time-delay, optimal control, Markov process, stochastic stabilizability, stochastic detectability

AMS subject classifications. 49J15, 93C57, 93E15

1. Introduction. Communication delays occur in networked control systems due to signal processing and congestion in busy channels. Such delays are generally time-varying, and if their range is large, control methods developed for systems with constant delays in [7,12] may not be suitable. With a rapid development of communication technologies, control under time-varying delays has received extensive attention over recent decades, as surveyed in [17,40].

One approach to compensate for time-varying delays is the virtual time-delay method developed, e.g., for bilateral control [24]. In this method, the maximum value of delays is assumed to be known, and control signals are updated when the maximum time of delays has passed. This method keeps the apparent delays constant but may degrade the performance of networked control systems if the maximum delay is quite larger than the average delay. Another approach is to measure delays by time-stamped messages and exploit these measurements in the control algorithms, as in Fig. 1. An example of this scenario can be found in inter-area power systems [31]. Controllers using time-stamp information in this way are delay-dependent, and stabilization by such controllers has been studied in [8,18–20,25,31,33,36,41] and references therein. Time-stamped messages are also used for linear quadratic (LQ) control in [23,26,27,34] and for model predictive control in [35] under stochastically time-varying delays.

In addition to the earlier studies mentioned above, the authors in [11,21,22,38,39] have developed design methods of LQ controllers for scenarios where the measurements of random delays are not available. However, most of those syntheses of optimal controllers require online computation and are based on assumptions that the delays can take only finitely many values or be independent and identically-distributed random variables. In practice, communication delays take a continuum of values and

∗Submitted to the editors DATE. This paper was partially presented at the American Control Conference 2017 [37].

Funding: This material is based upon work supported by JSPS KAKENHI Grant Number 17K14099 and the National Science Foundation under Grants No. CNS-1329650 and EPCN-1608889.

†Graduate School of System Informatics, Kobe University, Hyogo 657-8501, Japan (wakaiki@ruby.kobe.u.ac.jp).
‡Division of Information Science, Graduate School of Science and Technology, Nara Institute of Science and Technology, Nara 630-0192, Japan (oguram@is.naist.jp).
§Center for Control, Dynamical-systems and Computation (CCDC), University of California, Santa Barbara, CA 93106, USA (hespanha@ucsb.edu).
need to be modeled by more general stochastic processes such as Markov processes in [3,15,16]. A notable exception can be found in [23], in which the authors have presented an offline computation method of delay-dependent LQ controllers for systems with continuous-valued Markovian delays. The formulation in [23] requires a solution of a nonlinear vector integral equation called the Riccati integral equation and ignores the intersample behavior of the closed-loop system.

In this paper, we study delay-dependent LQ control for sampled-data linear systems. The advantages of the proposed method are twofold. First, our delay model is more general than that in the above previous studies. Indeed, in the model we consider, the present delay is determined by the last few delays like in an autoregressive models (see, e.g., Chapter 9 of [29] for autoregressive models), and hence our delay model belongs to a class of higher-order Markov models. Second, we can efficiently compute an LQ control law that takes into account the intersample behavior. This controller is obtained by iteratively solving a Riccati difference equation.

A key step in the construction of the controller consists of reducing the original sampled-data problem into an LQ problem for a discrete-time Markov jump system whose jumps are modeled by a Markov chain taking values in a general Borel space. In [5], the reduced LQ problem has been solved under the assumption that the plant and the LQ criterion satisfy appropriately defined stochastic stabilizability and detectability notions. However, there has been relatively little work on the test of these properties. To use the results on stochastic stability in [4], we obtain novel sufficient conditions for stochastic stabilizability and detectability in terms of linear matrix inequalities (LMIs). From these results, we can also construct stabilizing controllers and state observers. The proposed method is inspired by the gridding methods for establishing the stability of networked control systems with aperiodic sampling and time-varying delays in [9,10,13,18,28]. Moreover, we show that the sufficient conditions can be arbitrarily tight under certain assumptions.

The remainder of the paper is organized as follows. We provide the problem statement in Section 2. Section 3 is devoted to reducing our optimal control problem to an LQ problem for discrete-time Markov jump systems. In Section 4, we recall the general results in [5] on the equivalent discrete-time LQ problem. Section 5 addresses the derivation of sufficient conditions for stochastic stabilizability and detectability. In Section 6, we illustrate the proposed method with a numerical simulation of a batch reactor.

**Notation.** Let $\mathbb{Z}_+, \mathbb{R}^{n \times m}$, and $\mathbb{C}^{n \times m}$ denote the set of nonnegative integers and the sets of real and complex matrices with size $n \times m$, respectively. For a real matrix $M$, let us denote its transpose by $M^\top$. The Euclidean norm of $v \in \mathbb{R}^n$ is denoted by $\|v\| := (v^\top v)^{1/2}$ and the corresponding induced norm of a matrix $M \in \mathbb{R}^{m \times n}$ by $\|M\| := \sup\{\|Mv\| : v \in \mathbb{R}^n, \|v\| = 1\}$. For simplicity, we write a partitioned
Let \((\mathcal{M}, \mathcal{B}(\mathcal{M}))\) be a Borel space, that is, \(\mathcal{M}\) be a Borel subset of a complete and separable metric space and \(\mathcal{B}(\mathcal{M})\) be its Borel \(\sigma\)-algebra. In this article, \(\mathcal{M}\) is a compact subset of \(\mathbb{R}^p\) except in Section 4. For a \(\sigma\)-finite measure \(\mu\) on \(\mathcal{M}\), we denote by \(\mathbb{H}^{1,n}_{\mathcal{M}}\) the space of matrix-valued functions \(P(\bullet) : \mathcal{M} \to \mathbb{R}^{n \times m}\) that are measurable and integrable in \(\mathcal{M}\), and similarly, by \(\mathbb{H}^{n+}_{\mathcal{M}}\) the space of matrix-valued functions \(P(\bullet) : \mathcal{M} \to \mathbb{R}^{n \times m}\) that are measurable and essentially bounded in \(\mathcal{M}\). For \(P \in \mathbb{H}^{n+}_{\mathcal{M}}\), we define a norm \(\|P\|_{\infty}\) by the essential supremum of the function \(\|P\| : \mathcal{M} \to [0, \infty)\). For simplicity, we will write \(\mathbb{H}^1 := \mathbb{H}^{1,n}_{\mathcal{M}}\), and \(\mathbb{H}^n := \mathbb{H}^{n+}_{\mathcal{M}}\).

Additionally, we denote by \(\mathbb{H}_{\mathcal{M},\mathbb{C}}^{1,n}\) the space of matrix-valued functions \(P(\bullet) : \mathcal{M} \to \mathbb{C}^{n \times m}\) that are measurable and integrable in \(\mathcal{M}\) and by \(\mathbb{H}^{n+}_{\mathcal{M},\mathbb{C}}\) the space of matrix-valued functions \(P(\bullet) : \mathcal{M} \to \mathbb{C}^{n \times m}\) that are measurable and essentially bounded in \(\mathcal{M}\) and satisfy \(P(\phi) \geq 0\) for \(\mu\)-almost every \(\phi \in \mathcal{M}\). For a bounded linear operator \(T\) on a Banach space, let \(\tau_c(T)\) denote the spectral radius of \(T\).

2. Problem Statement. Consider the following linear continuous-time plant:

\[
\begin{align*}
\dot{x}(t) &= A_x x(t) + B_x u(t), & x(0) &= x_0
\end{align*}
\]

where \(x(t) \in \mathbb{R}^n\) and \(u(t) \in \mathbb{R}^m\) are the state and the input of the plant. This plant is connected to a controller through a time-driven sampler with period \(h > 0\) and an event-driven zero-order hold as shown in Fig. 1.

The state \(x\) is measured at each sampling time \(t = kh\) \((k \in \mathbb{Z}_+)\), and the controller receives the sampled state \(x(kh)\) at time \(t = kh + \tau_k\), where \(\tau_k > 0\) is a sensor-to-controller delay. We assume that the delay \(\tau_k\) becomes known to the controller at the time \(t = kh + \tau_k\) when the sampled state \(x(kh)\) arrives. One way to measure the delays is to mark every output of the sampler with a time-stamp and then to compute the difference between the value of the time-stamp and the present time of a clock in the controller. Through the zero-order hold, the discrete-time signal \(u_k\) generated from the controller is transformed to the continuous-time signal

\[
\begin{align*}
u(t) = \begin{cases} u_{-1} & 0 \leq t < \tau_0 \\ u_k & kh + \tau_k \leq t < (k + 1)h + \tau_{k+1}, \text{ } k \in \mathbb{Z}_+. \end{cases} \end{align*}
\]

where \(u_{-1}\) is an initial state of the zero-order hold.

Throughout this paper, we fix the probability space \((\Omega, \mathcal{F}, P)\). We assume that the delays \(\{\tau_k : k \in \mathbb{Z}_+\}\) is smaller than one sampling period and that the latest delay is stochastically determined by the last few delays. We specifically assume that the delay sequence \(\{\tau_k : k \in \mathbb{Z}_+\}\) is a higher-order Markov chain. For some known \(p \in \mathbb{N}\), define a delay vector \(\phi_k\) by

\[
\begin{align*}
\phi_k := \begin{bmatrix} \tau_k \\ \vdots \\ \tau_{k-p+1} \end{bmatrix} & \quad \forall k \in \mathbb{Z}_+,
\end{align*}
\]

where \(\tau_{-p+1}, \ldots, \tau_1 < h\) are the time delays associated with the sampling instants \(t = (-p + 1)h, \ldots, -h\).
Assumption 2.1 (Higher-order Markovian delays). The sequence \( \{\phi_k : k \in \mathbb{Z}_+\} \) in (3) is a time-homogeneous Markov chain taking values in \( \mathcal{M} := [\tau_{\min}, \tau_{\max}]^p \subset [0, h]^p \) and having transition probability kernel \( \mathcal{G}(\cdot, \cdot) \) with a density \( g(\cdot, \cdot) \) with respect to a \( \sigma \)-finite measure \( \mu \) on \( \mathcal{M} \), so that for every \( k \in \mathbb{Z}_+ \) and every Borel set \( \mathcal{B} \) of \( \mathcal{M} \),

\[
\mathcal{G}(\phi, \mathcal{B}) := P(\phi_{k+1} \in \mathcal{B} | \phi_k = \phi) = \int_{\mathcal{B}} g(\phi, \ell) \mu(d\ell).
\]

The choice of the dimension \( p \) depends on accuracy of delay models and computational cost. As the dimension \( p \) increases, we may obtain more accurate models of delays. However, a large \( p \) requires substantial computational resources for optimal controllers. Moreover, the gridding method presented in Section 5 suffers from the curse of dimensionality in the case of large \( p \).

Define

\[
(4)\quad \xi_0 := \begin{bmatrix} x_0 \\ u_{-1} \end{bmatrix}.
\]

Let \( \hat{\mu} \) be a probability measure on \( \mathbb{R}^{n+m} \times \mathcal{M} \). We assume that the pair of the initial state and delay \( (\xi_0, \phi_0) \) has a distribution \( \hat{\mu} \). Define \( \hat{\mu}_M \) by \( \hat{\mu}_M(\mathcal{B}) := \int (\mathbb{R}^{n+m} \times \mathcal{B}) = P(\hat{\theta}_0 \in \mathcal{B}) \) for all Borel sets \( \mathcal{B} \) of \( \mathcal{M} \). We place the following mild assumption on the initial distribution \( \hat{\mu} \):

**Assumption 2.2** (Initial distribution). The initial distribution \( \hat{\mu} \) of \( (\xi_0, \phi_0) \) satisfies A1) \( E(\|\xi_0\|^2) < \infty \) and A2) \( \hat{\mu}_M \) is absolutely continuous with respect to \( \mu \).

The assumption of absolute continuity guarantees the existence of the Radon-Nikodym derivative of \( \hat{\mu}_M \).

Let \( \{\mathcal{F}_k : k \in \mathbb{Z}_+\} \) denote a filtration, where \( \mathcal{F}_k \) represents the \( \sigma \)-field generated by

\[
\{u_{-1}, x(0), \phi_0, \ldots, x(kh), \phi_k\} = \{\tau_{-p+1}, \ldots, \tau_{-1}, u_{-1}, x(0), \tau_0, \ldots, x(kh), \tau_k\}.
\]

Set \( \mathcal{U}_c \) as the class of control inputs \( u = \{u_k : k \in \mathbb{Z}_+\} \) such that \( u_k \) is \( \mathcal{F}_k \) measurable and the controlled system (1) and (2) satisfies \( E(\|x(t)\|^2) \to 0 \) as \( t \to \infty \) and \( E(\|u_k\|^2) \to 0 \) as \( k \to \infty \) for every initial distribution \( \hat{\mu} \) satisfying Assumption 2.2. For all \( u \in \mathcal{U}_c \), we consider the infinite-horizon continuous-time quadratic cost functional \( \mathcal{J}_c \) defined by

\[
(5)\quad \mathcal{J}_c(\hat{\mu}, u) := E \left( \int_0^\infty x(t)\top Q_c x(t) + u(t)\top R_c u(t) dt \right),
\]

where \( Q_c \geq 0 \) and \( R_c > 0 \) are weighting matrices with appropriate dimensions.

In this paper, we study the following LQ problem:

**Problem 2.3.** Consider a sampled-data system (1) and (2), and let Assumptions 2.1 and 2.2 hold. Find an optimal control law \( u^{\text{opt}} \in \mathcal{U}_c \) that achieves \( \mathcal{J}_c(\hat{\mu}, u^{\text{opt}}) = \inf_{u \in \mathcal{U}_c} \mathcal{J}_c(\hat{\mu}, u) \) for every initial distribution \( \hat{\mu} \) satisfying Assumption 2.2.

**Remark 2.4.** In this paper, we impose the following two assumptions on delays: (i) Delays are smaller than one sampling period \( h \). If the controller-to-actuator delays \( \tau_{ca} \) are deterministic (see, e.g., [40, Section 2.3.2] for this situation), then we can also apply the proposed method in the presence of the controller-to-actuator delays by using the total delays \( \tau_k + \tau_{ca} \) instead of the sensor-to-controller delays \( \tau_k \). To deal
with delays larger than one sampling period \( h \), we can employ the technique presented in Section II of [2], but a stochastic model on delays becomes complicated. Therefore, we here assume that delays are smaller than one sampling period \( h \).

3. Reduction to Discrete-time LQ Problem. In this section, we transform Problem 2.3 to an LQ problem of discrete-time Markov jump linear systems.

Consider the sampled-data system (1) and (2). We define

\[
\xi_k := \begin{bmatrix} x(kh) \\ u_{k-1} \end{bmatrix}.
\]

Then the dynamics of \( \xi \) can be described by the following discrete-time Markov jump linear system

\[
\xi_{k+1} = A(\phi_k)\xi_k + B(\phi_k)u_k,
\]

where, for every vector \( \phi \in \mathcal{M} \) whose the first element is given by \( \tau \), we define

\[
A(\phi) := \begin{bmatrix} A_d & B_d - \Gamma(\tau) \\ 0 & 0 \end{bmatrix}, \quad B(\phi) := \begin{bmatrix} \Gamma(\tau) \\ I \end{bmatrix}
\]

\[
A_d := e^{A_h}, \quad B_d := \int_0^h e^{A_\tau}B_\tau ds, \quad \Gamma(\tau) := \int_0^{h-\tau} e^{A_\tau}B_\tau ds.
\]

By definition, the matrices \( A \) and \( B \) satisfy \( A \in \mathbb{H}^{n}_{\sup} \) and \( B \in \mathbb{H}^{n \times m}_{\sup} \). This delay-dependent discrete-time system is widely used for the analysis of time-delay systems, e.g., in [9, 10, 18, 23, 26, 27, 34].

Let \( \mathcal{F}_k^d \) denote a filtration, where \( \mathcal{F}_k^d \) represents the \( \sigma \)-field generated by

\[
\{ \xi_0, \phi_0, \ldots, \xi_k, \phi_k \}.
\]

We denote by \( \mathcal{U}_d \) the discrete-time counterpart of \( \mathcal{U}_c \), defined as the class of control inputs \( \{ u_k : k \in \mathbb{Z}_+ \} \) such that \( u_k \) is \( \mathcal{F}_k^d \) measurable and \( E(\| \xi_k \|^2) \to 0 \) as \( k \to \infty \) for every initial distribution \( \hat{\mu} \). The following result establishes that these classes of control inputs are equal.

**Lemma 3.1.** For the sampled-data system (1), (2) and its discretization (7), we obtain \( \mathcal{U}_c = \mathcal{U}_d \).

**Proof.** Since the filtrations \( \{ \mathcal{F}_k \} \) and \( \{ \mathcal{F}_k^d \} \) are equal by definition, it is enough to prove that the following two conditions are equivalent:

1. \( \lim_{k \to \infty} E(\| \xi_k \|^2) = 0 \).
2. \( \lim_{k \to \infty} E(\| x(t) \|^2) = 0 \) and \( \lim_{k \to \infty} E(\| u_k \|^2) = 0 \).

The statement 2 \( \Rightarrow \) 1 follows directly from the definition of \( \xi_k \). To prove the converse statement, we note that for the system dynamics (1) and (2), there exist constants \( M_1, M_2, M_3 > 0 \) such that

\[
\| x(kh + \theta) \| \leq M_1 \| x(kh) \| + M_2 \| u_{k-1} \| + M_3 \| u_k \| \quad \forall \theta \in [0, h).
\]

Therefore

\[
\| x(kh + \theta) \|^2 \leq M_1^2 \| x(kh) \|^2 + M_2^2 \| u_{k-1} \|^2 + M_3^2 \| u_k \|^2 + 2M_1M_2\| x(kh) \| \cdot \| u_{k-1} \| + 2M_2M_3\| u_{k-1} \| \cdot \| u_k \| + 2M_3M_1 \| u_k \| \cdot \| x(kh) \| \quad \forall \theta \in [0, h).
\]
For every vector \( \phi \in \mathcal{M} \) whose first element is given by \( \tau \), we define the matrices \( Q, W, \) and \( R \) by

\begin{align}
Q(\phi) &:= \begin{bmatrix} Q_{11}(\tau) & Q_{12}(\tau) \\ * & Q_{22}(\tau) \end{bmatrix} \\
W(\phi) &:= \int_{\tau}^{h} \begin{bmatrix} \alpha(\theta)^\top Q_c \gamma(\tau, \theta) \\ \beta(\theta)^\top Q_c \gamma(\tau, \theta) - \gamma(\tau, \theta)^\top Q_c \gamma(\tau, \theta) \end{bmatrix} d\theta \\
R(\phi) &:= (h - \tau)R_c + \int_{\tau}^{h} \gamma(\tau, \theta)^\top Q_c \gamma(\tau, \theta) d\theta \\
Q_{11}(\tau) &:= \int_{0}^{h} \alpha(\theta)^\top Q_c \alpha(\theta) d\theta \\
Q_{12}(\tau) &:= \int_{0}^{h} \alpha(\theta)^\top Q_c \beta(\theta) d\theta - \int_{\tau}^{h} \alpha(\theta)^\top Q_c \gamma(\tau, \theta) d\theta \\
Q_{22}(\tau) &:= \tau R_c + \int_{0}^{h} \beta(\theta)^\top Q_c \beta(\theta) d\theta + \int_{\tau}^{h} \gamma(\tau, \theta)^\top Q_c \gamma(\tau, \theta) d\theta \\
&\quad - \int_{\tau}^{h} (\beta(\theta)^\top Q_c \gamma(\tau, \theta) + \gamma(\tau, \theta)^\top Q_c \beta(\theta)) d\theta,
\end{align}

where the functions \( \alpha, \beta, \) and \( \gamma \) are given by

\begin{align}
\alpha(\theta) &:= e^{A_c \theta}, \quad \beta(\theta) := \int_{0}^{\theta} e^{A_c s} B_c ds \quad \forall \theta \in [0, h] \\
\gamma(\tau, \theta) &:= \int_{0}^{\theta - \tau} e^{A_c s} B_c ds \quad \forall \tau \in [\tau_{\min}, \tau_{\max}], \forall \theta \in [0, h].
\end{align}

The next lemma shows that each \( J_k \) is a quadratic form on the state \( \xi_k \) and the input \( u_k \) of the discrete-time system \( \Sigma_d \) in (7).
Lemma 3.2. Let \( x \) and \( \xi \) be the solutions of the sampled-data system (1) and (2) and of the discrete-time system (7) with the initial state \( \xi_0 \) defined by (4), receptively. Then \( J_k \) defined by (10) satisfies

\[
J_k = \begin{bmatrix} \xi_k \\ u_k \end{bmatrix}^\top \begin{bmatrix} Q(\phi_k) & W(\phi_k) \\ \ast & R(\phi_k) \end{bmatrix} \begin{bmatrix} \xi_k \\ u_k \end{bmatrix} \quad \forall k \in \mathbb{Z}_+,
\]

where the matrices \( Q, W, \) and \( R \) are defined as in (11).

**Proof.** If \( 0 < \theta < \tau_k \), then

\[
x(kh + \theta) = \alpha(\theta)x_k + \beta(\theta)u_{k-1},
\]

\[
u(kh + \theta) = u_{k-1},
\]

and we have

\[
x(kh + \theta)^\top Q_c x(kh + \theta) + u(kh + \theta)^\top R_c u(kh + \theta)
\]

\[= x_k^\top \alpha(\theta) Q_c \alpha(\theta)x_k + 2x_k^\top \alpha(\theta)^\top Q_c \beta(\theta)u_{k-1} + u_{k-1}^\top \beta(\theta)^\top Q_c \beta(\theta)u_{k-1}
\]

\[+ u_{k-1}^\top R_c u_{k-1}.
\]

On the other hand, if \( \tau_k \leq \theta < h \), then

\[
x(kh + \theta) = \alpha(\theta)x_k + (\beta(\theta) - \gamma(\tau_k, \theta))u_{k-1} + \gamma(\tau_k, \theta)u_k
\]

\[
u(kh + \theta) = u_k.
\]

Hence

\[
x(kh + \theta)^\top Q_c x(kh + \theta) + u(kh + \theta)^\top R_c u(kh + \theta)
\]

\[= x_k^\top \alpha(\theta) Q_c \alpha(\theta)x_k + 2x_k^\top \alpha(\theta)^\top Q_c \beta(\theta) - \gamma(\tau_k, \theta)u_{k-1}
\]

\[+ 2x_k^\top \alpha(\theta)^\top Q_c \gamma(\tau_k, \theta)u_k + u_{k-1}^\top (\beta(\theta) - \gamma(\tau_k, \theta))^\top Q_c (\beta(\theta) - \gamma(\tau_k, \theta))u_{k-1}
\]

\[+ 2u_{k-1}^\top (\beta(\theta) - \gamma(\tau_k, \theta))^\top Q_c (\gamma(\tau_k, \theta)u_k + u_k^\top (\gamma(\tau_k, \theta))^\top Q_c (\gamma(\tau_k, \theta)u_k + u_k^\top R_c u_k.
\]

Substituting these equations into

\[
J_k = \int_{kh}^{kh+\tau_k} (x(t)^\top Q_c x(t) + u(t)^\top R_c u(t)) \, dt
\]

\[+ \int_{(k+1)h}^{(k+1)h+\tau_k} (x(t)^\top Q_c x(t) + u(t)^\top R_c u(t)) \, dt,
\]

we obtain the desired expression (12).

For the discrete-time Markov jump system (7), we define the infinite-horizon discrete-time quadratic cost functional \( J_{d1} \) by

\[
J_{d1}(\hat{\mu}, u) := \sum_{k=0}^{\infty} E \left( \begin{bmatrix} \xi_k \\ u_k \end{bmatrix}^\top \begin{bmatrix} Q(\phi_k) & W(\phi_k) \\ \ast & R(\phi_k) \end{bmatrix} \begin{bmatrix} \xi_k \\ u_k \end{bmatrix} \right).
\]

Then a solution to Problem 2.3 for sampled-data systems with stochastic delays can be obtained as a solution to the following problem for discrete-time Markov jump systems:
PROBLEM 3.3. Consider a discrete-time Markov jump system (7), and let Assumptions 2.1 and 2.2 hold. Find an optimal control law $u^{opt} \in \mathcal{U}_d$ that achieves $J_{d1}(\mu, u^{opt}) = \inf_{u \in \mathcal{U}_d} J_{d1}(\mu, u)$ for every initial distribution $\mu$ satisfying Assumption 2.2.

LEMMA 3.4. A control input $u^{opt}$ is a solution to Problem 2.3 if and only if $u^{opt}$ is also a solution to Problem 3.3 where the system matrices $A$, $B$ and the weighting matrices $Q$, $W$, $R$ are defined by (8) and (11).

Proof. Since $\mathcal{U}_c = \mathcal{U}_d$ from Lemma 3.1, it follows that $u^{opt} \in \mathcal{U}_c$ if and only if $u^{opt} \in \mathcal{U}_d$. Let $x$ and $\xi$ be the solutions of the sampled-data system (1) and (2) and of the discrete-time system (7) with the initial state $\xi_0$ defined by (4), receptively. Since

$$x(t)^\top Q_c x(t) + u(t)^\top R_c u(t) \geq 0 \quad \forall t \geq 0,$$

we obtain

$$J_c(\mu, u) = E \left( \int_0^\infty x(t)^\top Q_c x(t) + u(t)^\top R_c u(t) dt \right)$$

$$= E \sum_{k=0}^{\infty} \left( \int_{kh}^{(k+1)h} x(t)^\top Q_c x(t) + u(t)^\top R_c u(t) dt \right)$$

$$= \sum_{k=0}^{\infty} E \left( \int_{kh}^{(k+1)h} x(t)^\top Q_c x(t) + u(t)^\top R_c u(t) dt \right).$$

It also follows from Lemma 3.2 that

$$\sum_{k=0}^{\infty} E \left( \int_{kh}^{(k+1)h} x(t)^\top Q_c x(t) + u(t)^\top R_c u(t) dt \right)$$

$$= \sum_{k=0}^{\infty} E \left( \begin{bmatrix} \xi_k \\ u_k \end{bmatrix}^\top \begin{bmatrix} Q(\phi_k) & W(\phi_k) \\ \ast & R(\phi_k) \end{bmatrix} \begin{bmatrix} \xi_k \\ u_k \end{bmatrix} \right) = J_{d1}(\mu, u).$$

Thus if a control input $u^{opt} \in \mathcal{U}_c$ is a solution to Problem 2.3, then $u^{opt}$ satisfies $u^{opt} \in \mathcal{U}_d$ and is a solution to Problem 3.3 where the system matrices $A$, $B$ and the weighting matrices $Q$, $W$, $R$ are defined by (8), (11), and vice versa.

Let us next remove the cross term of the cost function $J_{d1}$. To this end, as in the deterministic case [1, Section 3.4], we transform $u_k$ into $\bar{u}_k$ in the following way:

(14) $$\bar{u}_k = u_k + R(\phi_k)^{-1} W(\phi_k)^\top \xi_k \quad \forall k \in \mathbb{Z}_+.$$ 

Since $Q_c \geq 0$, $R_c > 0$, and $h - \tau_{\text{max}} > 0$, it follows that $R(\phi)$ in (11) is invertible for all $\phi \in \mathcal{M}$. Therefore the right-hand side of (14) is well-defined for all $\phi_k \in \mathcal{M}$.

Define

(15) $$\bar{A}(\phi) := A(\phi) - B(\phi) R(\phi)^{-1} W(\phi)^\top \quad \forall \phi \in \mathcal{M}$$

and let $C(\phi)$ and $D(\phi)$ be the matrices obtained from the following Cholesky decompositions:

(16a) $$C(\phi)^\top C(\phi) = Q(\phi) - W(\phi) R(\phi)^{-1} W(\phi)^\top \quad \forall \phi \in \mathcal{M}$$

(16b) $$D(\phi)^\top D(\phi) = R(\phi) \quad \forall \phi \in \mathcal{M}.$$
These Cholesky decompositions are possible if the weighting matrix in Lemma 3.2 satisfies
\begin{equation}
\begin{bmatrix}
Q(\phi) & W(\phi) \\
* & R(\phi)
\end{bmatrix} \geq 0 \quad \forall \phi \in \mathcal{M}.
\end{equation}
This is because \( R(\phi) > 0 \) for every \( \phi \in \mathcal{M} \) and the Schur complement formula leads to
\[ Q(\phi) - W(\phi)R(\phi)^{-1}W(\phi)^T \geq 0 \quad \forall \phi \in \mathcal{M}. \]
Under the transformation (14), we obtain the following result:

**Lemma 3.5.** Assume that the weighting matrices \( Q, W, R \) in (11) satisfy the inequality (17). A control input \( u^{\text{opt}} \) is a solution to Problem 3.3 with the system (7) and the LQ cost (13) if and only if \( \bar{u}^{\text{opt}} = u^{\text{opt}}_k + R(\phi_k)^{-1}W(\phi_k)^T \xi_k \) is a solution to Problem 3.3 where the Markov jump system is given by
\begin{equation}
\xi_{k+1} = \bar{A}(\phi_k)\xi_k + B(\phi_k)\bar{u}_k
\end{equation}
and the LQ cost by
\begin{equation}
J_{d2}(\bar{\mu}, \bar{u}) := \sum_{k=0}^{\infty} E \left( ||C(\phi_k)\xi_k||^2 + ||D(\phi_k)\bar{u}_k||^2 \right).
\end{equation}
Here the matrices \( \bar{A}, C, D \) are defined as in (15) and (16).

**Proof.** Let \( \xi \) and \( \bar{\xi} \) be the solutions of the difference equations (7) and (18) with the same Markov parameter \( \phi \). If \( \bar{u}_k \) satisfies (14) and if \( \xi_k = \bar{\xi}_k \), then
\begin{align*}
\xi_{k+1} &= \bar{A}(\phi_k)\xi_k + B(\phi_k)\bar{u}_k \\
&= \bar{A}(\phi_k)\bar{\xi}_k + B(\phi_k)(\bar{u}_k - R(\phi_k)^{-1}W(\phi_k)^T \xi_k) \\
&= \bar{A}(\phi_k)\bar{\xi}_k + B(\phi_k)\bar{u}_k = \bar{\xi}_{k+1}.
\end{align*}
Therefore, if \( \xi_0 = \bar{\xi}_0 \), then \( \xi_k = \bar{\xi}_k \) for every \( k \in \mathbb{Z}_+ \). Thus \( u \in \mathcal{U}_d \) for the system (7) if and only if \( \bar{u} \in \mathcal{U}_d \) for the system (18). Moreover, if (14) holds, then
\begin{align*}
J_{d2}(\bar{\mu}, \bar{u}) &= \sum_{k=0}^{\infty} E \left( \begin{bmatrix} \xi_k \\ \bar{u}_k \end{bmatrix}^T \begin{bmatrix}
C(\phi_k) & 0 \\
D(\phi_k) & 0
\end{bmatrix} \begin{bmatrix} \xi_k \\ \bar{u}_k \end{bmatrix} \right) \\
&= \sum_{k=0}^{\infty} E \left( \begin{bmatrix} \xi_k \\ \bar{u}_k \end{bmatrix}^T \begin{bmatrix}
Q(\phi_k) - W(\phi_k)R(\phi_k)^{-1}W(\phi_k)^T & 0 \\
0 & R(\phi_k)
\end{bmatrix} \begin{bmatrix} \xi_k \\ \bar{u}_k \end{bmatrix} \right) \\
&= \sum_{k=0}^{\infty} E \left( \begin{bmatrix} \xi_k \\ \bar{u}_k \end{bmatrix}^T \begin{bmatrix}
Q(\phi_k) & W(\phi_k) \\
* & R(\phi_k)
\end{bmatrix} \begin{bmatrix} \xi_k \\ \bar{u}_k \end{bmatrix} \right) = J_{d1}(\bar{\mu}, \bar{u}).
\end{align*}
Hence \( u^{\text{opt}} \) is a solution to Problem 3.3 with the system (7) and the LQ cost (13) if and only if \( \bar{u}^{\text{opt}} = u^{\text{opt}}_k + R(\phi_k)^{-1}W(\phi_k)^T \xi_k \) is a solution to Problem 3.3 with the system (18) and the LQ cost (19).

Finally we can reduce the LQ problem 2.3 for sampled-data systems with stochastic delays to the LQ problem 3.3 for discrete-time Markov jump systems and LQ costs in the form (19).
Theorem 3.6. A control input $u_{\text{opt}}$ is a solution to Problem 2.3 if and only if $\bar{u}_{\text{opt}} = u_{\text{opt}} + R(\phi_k)^{-1}W(\phi_k)^\top \xi_k$ is a solution to Problem 3.3 for the Markov jump system (18) and the LQ cost (19), where the matrices $\bar{A}, C,$ and $D$ are obtained in (15) and (16) from the matrices $A, B, Q, W,$ and $R$ defined by (8) and (11).

Proof. This is an immediate consequence of Lemmas 3.4 and 3.5.

Remark 3.7. If, in (17), we have strict positive definiteness

\begin{align}
\left[\begin{array}{cc}
Q(\phi) & W(\phi) \\
* & R(\phi)
\end{array}\right] > 0 & \forall \phi \in \mathcal{M}
\end{align}

instead of the semidefiniteness, then

\begin{align}
Q(\phi) - W(\phi)R(\phi)^{-1}W(\phi)^\top > 0 & \forall \phi \in \mathcal{M}
\end{align}

by the Schur complement formula. In this case, $C(\phi)$ and $D(\phi)$, derived from the Cholesky decompositions (16), are unique in the following sense: For all $\phi \in \mathcal{M}$, there exist unique upper triangular matrices $C(\phi)$ and $D(\phi)$ with strictly positive diagonal entries such that (16) holds. Moreover, $C(\phi), D(\phi)$ are continuous with respect to $\phi$. See, e.g., Chapters 9 and 12 of [32]. Thus $C$ and $D$ satisfies $C \in \mathbb{H}_{n\times n}^\sup$ and $D \in \mathbb{H}_{m \times m}^\sup$.

4. LQ control for Discrete-time Markov Jump Systems. In the previous section, we reduced the LQ problem 2.3 for sampled-data systems with stochastic delays into the LQ problem 3.3 for discrete-time Markov jump systems and LQ costs in the form (19). In this section, we recall results from [5] on such an LQ problem for Markov jump systems.

First we define stochastic stability for discrete-time Markov jump linear systems. On a probability space $(\Omega, \mathcal{F}, P)$, consider the following autonomous system

\begin{align}
\xi_{k+1} = A(\phi_k)\xi_k,
\end{align}

where $A \in \mathbb{H}_n^\sup$ and the sequence \{\phi_k : k \in \mathbb{Z}_+\} is a time-homogeneous Markov chain in a Borel space $\mathcal{M}$. Throughout this section, we assume that the initial distribution $\hat{\mu}$ of $(\xi_0, \phi_0)$, which is a probability measure on $\mathbb{R}^n \times \mathcal{M}$, satisfies the following conditions analogous to the ones in Assumption 2.2: A1') $E(\|\xi_0\|^2) < \infty$ and A2') $\hat{\mu}_{\mathcal{M}}(\bullet) = \hat{\mu}(\mathbb{R}^n \times \bullet)$ is absolutely continuous with respect to $\mu$.

Definition 4.1 (Stochastic stability, [4]). The autonomous Markov jump linear system (21) is said to be stochastically stable if $\sum_{k=0}^{\infty} E(\|\xi_k\|^2) < \infty$ for any initial distribution $\hat{\mu}$ satisfying A1') and A2').

Let $g(\bullet, \bullet)$ be the density function with respect to a $\sigma$-finite measure $\mu$ on $\mathcal{M}$ for the transition of the Markov chain \{\phi_k : k \in \mathbb{Z}_+\} as in Assumption 2.1. For every $A \in \mathbb{H}_n^\sup$, define an operator $L_A : \mathbb{H}_{n,C}^\sup \to \mathbb{H}_{1,C}^\sup$ by

\begin{align}
L_A(V)(\bullet) := \int_\mathcal{M} g(\ell, \bullet)A(\ell)V(\ell)A(\ell)^\top \mu(d\ell).
\end{align}

We recall a relationship among stochastic stability, the spectral radius $r_\sigma(L_A)$, and a Lyapunov inequality condition.

Theorem 4.2 ([4]). The following assertions are equivalent:

1. The system (21) is stochastically stable.
2. The spectral radius \( r_\sigma(\mathcal{L}_A) < 1 \), where \( \mathcal{L}_A \) is defined as in (22).

3. There exist \( S \in \mathbb{H}_{\sup}^{n+} \) and \( \epsilon > 0 \) such that the following Lyapunov inequality holds for \( \mu \)-almost every \( \phi \in \mathcal{M} \):

\[
S(\phi) - A(\phi)^\top \left( \int_{\mathcal{M}} g(\phi, \ell) S(\ell) \mu(\ell) d\ell \right) A(\phi) \geq \epsilon I.
\]  

(23)

Note that the matrix \( S \) in the statement 3 is a complex-valued function, because the system matrix \( A \) is a complex-valued function in [4]. However, for a real-valued function \( A \), it is enough to find a real-valued \( S \) in the statement 3.

**Proposition 4.3.** For a real-valued function \( A \in \mathbb{H}_{\sup}^n \), the statement 3 in Theorem 4.2 is equivalent to

3’. There exist \( S \in \mathbb{H}_{\sup}^{n+} \) and \( \epsilon > 0 \) such that the Lyapunov inequality (23) holds for \( \mu \)-almost every \( \phi \in \mathcal{M} \).

**Proof.** The statement 3’ \( \Rightarrow \) 3 is trivial because \( \mathbb{H}_{\sup}^{n+} \subset \mathbb{H}_{\sup}^{n+} \). To prove 3 \( \Rightarrow \) 3’, we let \( S \in \mathbb{H}_{\sup}^{n+} \) and \( \epsilon > 0 \) satisfy the Lyapunov inequality (23). Let \( S_R(\phi) \in \mathbb{R}^{n \times n} \) and \( S_I(\phi) \in \mathbb{R}^{n \times n} \) be the real and imaginary part of \( S(\phi) \), that is,

\[
S(\phi) = S_R(\phi) + iS_I(\phi).
\]

Since \( S(\phi) \) is Hermitian, it follows that \( S_R(\phi) = S_R(\phi)^\top \) and \( S_I(\phi) = -S_I(\phi)^\top \).

Therefore,

\[
0 \leq x^\top S(\phi)x = x^\top (S_R(\phi) + iS_I(\phi))x = x^\top S_R(\phi)x \quad \forall x \in \mathbb{R}^n.
\]

Thus we obtain \( S_R \in \mathbb{H}_{\sup}^{n+} \). Similarly, since \( A(\phi) \) and \( g(\phi, \ell) \) are real-valued, it follows that

\[
\epsilon \|x\|^2 \leq x^\top \left( S(\phi) - A(\phi)^\top \left( \int_{\mathcal{M}} g(\phi, \ell) S(\ell) \mu(\ell) d\ell \right) A(\phi) \right) x
\]

\[
= x^\top \left( S_R(\phi) - A(\phi)^\top \left( \int_{\mathcal{M}} g(\phi, \ell) S_R(\ell) \mu(\ell) d\ell \right) A(\phi) \right) x
\]

\[
\forall x \in \mathbb{R}^n.
\]

Hence \( S_R \) also satisfies the Lyapunov inequality (23) for \( \mu \)-almost every \( \phi \in \mathcal{M} \). This completes the proof. \( \square \)

We next provide the definition of stochastic stabilizability and stochastic detectability.

**Definition 4.4 (Stochastic stabilizability, [5]).** Let \( A \in \mathbb{H}_{\sup}^n \) and \( B \in \mathbb{H}_{\sup}^{n \times m} \). We say that \((A, B)\) is stochastically stabilizable if there exists \( F \in \mathbb{H}_{\sup}^{n \times n} \) such that \( r_\sigma(\mathcal{L}_{A+BF}) < 1 \), where \( \mathcal{L}_{A+BF} \) is defined as in (22). In this case, \( F \) is said to stochastically stabilize \((A, B)\).

**Definition 4.5 (Stochastic detectability, [5]).** Let \( A \in \mathbb{H}_{\sup}^n \) and \( C \in \mathbb{H}_{\sup}^n \). We say that \((C, A)\) is stochastically detectable if there exists \( L \in \mathbb{H}_{\sup}^{n \times n} \) such that \( r_\sigma(\mathcal{L}_{A+LC}) < 1 \), where \( \mathcal{L}_{A+LC} \) is defined as in (22).

Consider the controlled system

\[
\xi_{k+1} = A(\phi_k)\xi_k + B(\phi_k)u_k
\]

and the LQ cost

\[
\mathcal{J}_{d}(\hat{\mu}, u) = \sum_{k=0}^{\infty} E \left( \| C(\phi_k)\xi_k \|^2 + \| D(\phi_k)u_k \|^2 \right).
\]
where $A \in \mathbb{H}^{n \times n}_{\text{sup}}, B \in \mathbb{H}^{m \times n}_{\text{sup}}, C \in \mathbb{H}^{n \times r}_{\text{sup}},$ and $D \in \mathbb{H}^{m \times q}_{\text{sup}}$. We assume that there exists $\epsilon_D > 0$ such that $D(\phi)^TD(\phi) > \epsilon_D I$ for $\mu$-almost every $\phi \in \mathcal{M}$. As in Section 3, let $\{ \mathcal{F}^d_k : k \in \mathbb{Z}_+ \}$ denote a filtration, where $\mathcal{F}^d_k$ represents the $\sigma$-field generated by $\{\xi_0, \phi_0, \ldots, \xi_k, \phi_k\}$, and set $\mathcal{U}_d$ as the class of control inputs $u = \{u_k : k \in \mathbb{Z}_+\}$ such that $u_k$ is $\mathcal{F}^d_k$ measurable and $E(\|\xi_k\|^2) \to 0$ as $k \to \infty$ for every initial distribution $\mu$ satisfying $A1'$ and $A2'$.

Define operators $\mathcal{E} : \mathbb{H}^{n+}_{\text{sup}} \to \mathbb{H}^{n+}_{\text{sup}}, \mathcal{V} : \mathbb{H}^{n+}_{\text{sup}} \to \mathbb{H}^{m+}_{\text{sup}},$ and $\mathcal{R} : \mathbb{H}^{n+}_{\text{sup}} \to \mathbb{H}^{n+}_{\text{sup}},$ as follows:

$$
\mathcal{E}(Z)(\bullet) := \int_\mathcal{M} Z(\ell)g(\bullet, \ell)\mu(d\ell) \\
\mathcal{V}(Z) := D^TD + B^T\mathcal{E}(Z)B \\
\mathcal{R}(Z) := C^TC + A^T(\mathcal{E}(Z) - \mathcal{E}(Z)BV(Z)^{-1}B^T\mathcal{E}(Z))A.
$$

Using these operators, we can obtain a solution to the LQ problem for discrete-time Markov jump linear systems from the iterative computation of $\mathcal{R}(Z)$.

**Theorem 4.6 ([5]).** Consider the Markov jump system (18) with the LQ cost $J_d$ in (19). If $(A, B)$ is stochastically stabilizable and $(C, A)$ is stochastically detectable, then there exists a function $S \in \mathbb{H}^{n+}_{\text{sup}}$ such that $S$ is the unique solution in $\mathbb{H}^{n+}_{\text{sup}}$ of the $\mathcal{M}$-coupled algebraic Riccati equation

$$(24) \\
S(\phi) = \mathcal{R}(S)(\phi) \quad \mu\text{-almost every } \phi \in \mathcal{M}$$

and such that

$$
K := -\mathcal{V}(S)^{-1}B^T\mathcal{E}(S)A \in \mathbb{H}^{m \times n}_{\text{sup}}
$$

stochastically stabilizes $(A, B)$. The control input $u^{\text{opt}} \in \mathcal{U}_d$ defined by $u_k^{\text{opt}} := K(\phi_k)\xi_k$ achieves

$$
J_d(\hat{\mu}, u^{\text{opt}}) = \inf_{u \in \mathcal{U}_d} J_d(\hat{\mu}, u) = E(\xi_0^T S(\phi_0)\xi_0)
$$

for every initial distribution $\hat{\mu}$ satisfying $A1'$ and $A2'$. Moreover, we can compute the solution $S$ of the Riccati equation (24) in the following way: For any $\Xi \in \mathbb{H}^{n+}_{\text{sup}},$ the sequence $\{Y^\eta_k\}_{k=0}^{\infty}$ that is calculated by solving a (backward recursive) Riccati difference equation $Y^\eta_k = \mathcal{R}(Y^\eta_{k+1})$ with the initial value $Y^\eta_0 = \Xi$ satisfies $Y^\eta_0(\phi) \to S(\phi)$ as $\eta \to \infty$ for $\mu$-almost every $\phi \in \mathcal{M}$.

Let us go back to the reduced LQ problem 3.3 for the Markov jump system (18) and the LQ cost (19), where the matrices $\hat{A}, C,$ and $D$ are defined by (15) and (16). The Cholesky decompositions for all $\phi \in \mathcal{M}$ in (16) requires heavy computational cost, but the weighting functions $C$ and $D$ appear only in $\mathcal{R}$ and $\mathcal{V}$ in the forms $C^TC$ and $D^TD$. Hence, to compute an optimal control input $u^{\text{opt}}$, we do not need the Cholesky decompositions in (16). Although we still need $C$ to check the stochastic detectability of $(C, \hat{A})$, we see from Proposition 4.7 below that it is enough to test the stochastic detectability of $(C^TC, \hat{A})$ if $C^TC$ is positive definite.

**Proposition 4.7.** Define $\hat{A}$ and $C$ as in (15) and (16), respectively. The pair $(A, B)$ is stochastically stabilizable if and only if $(\hat{A}, B)$ is stochastically stabilizable. Moreover, under the positive definiteness of the weighting matrix in (20), $(C, \hat{A})$ is stochastically detectable if and only if $(C^TC, \hat{A})$ is stochastically detectable.
Proof. By the definition of $\hat{A}$, $K \in \mathbb{H}^{m \times n}_{\text{st}}$ stochastically stabilizes $(A, B)$ if and only if $K - R^{-1}W \in \mathbb{H}^{m \times n}_{\text{st}}$ stochastically stabilizes $(\hat{A}, B)$.

From the discussion in Remark 3.7, $C^{-1} \in \mathbb{H}^{n \times m}_{\text{sup}}$ if (20) holds. Hence if $(C, \hat{A})$ is stochastically detectable with an observer gain $L \in \mathbb{H}^{n \times r}_{\text{sup}}$, then $(C^\top C, \hat{A})$ is also stochastically detectable with an observer gain $L(C^\top)^{-1} \in \mathbb{H}^{n \times n}_{\text{sup}}$, and vice versa. $rown$

5. Sufficient conditions for Stochastic Stabilizability and Detectability.

From the results in Sections 3 and 4, we can obtain an optimal controller under the assumption that $(A, B)$ defined by (8) is stochastically stabilizable and $(C, \hat{A})$ defined by (15) and (16) is stochastically detectable. This assumption does not hold in general (and hence the solution of the Riccati difference equation may diverge) even if $(A_c, B_c)$ is controllable and $(Q_c, A_c)$ is observable. The major difficulty in this controller design is to check the stochastic stabilizability and detectability, namely, to show the existence of $F \in \mathbb{H}^{m \times n}_{\text{sup}}$ and $L \in \mathbb{H}^{n \times r}_{\text{sup}}$ such that the spectral radii of the operators $L_{A_c}BF$ and $L_{A_c}LC$ are less than one. In this section, we provide sufficient conditions for these properties in terms of LMIs. To this end, we use the following technical result:

**Lemma 5.1 ([14]).** For every square matrix $U$ and every positive definite matrix $S$,

$$US^{-1}U^\top \geq U + U^\top - S.$$  

We here assume that $\{\phi_k : k \in \mathbb{Z}_+\}$ is a time-homogeneous Markov chain taking values in the box $\mathcal{M} = [\tau_{\min}, \tau_{\max}]^r$. As in Assumption 2.1, let $g(\bullet, \bullet)$ be the density function with respect to a $\sigma$-finite measure $\mu$ on $\mathcal{M}$ for the transition of the Markov chain $\{\phi_k : k \in \mathbb{Z}_+\}$.

5.1. Stochastic Stabilizability. We first study the stochastic stabilizability of the pair $A \in \mathbb{H}^n_{\text{sup}}$ and $B \in \mathbb{H}^{m \times n}$. In our LQ problem, we need to check the stochastic stabilizability of the pair $(A, B)$ defined by (8).

Divide $\mathcal{M} = [\tau_{\min}, \tau_{\max}]^r$ into $N$ disjoint boxes $\{\mathcal{B}_i\}_{i=1}^N$ (whose union is $\mathcal{M}$), e.g., by splitting each interval $[\tau_{\min}, \tau_{\max}]$ into $r$ intervals $[s_i, s_{i+1})$ $(i = 1, \ldots, r - 1)$ and $[s_r, s_{r+1})$ such that

$$\tau_{\min} := s_1 < s_2 < \cdots < s_{r+1} := \tau_{\max}.$$  

For each $i = 1, \ldots, N$, let $c_i \in \mathcal{B}_i$, e.g., the center of $\mathcal{B}_i$. Consider a piecewise-constant feedback gain $F \in \mathbb{H}^{m \times n}$ defined by

$$F(\phi) := F_i \in \mathbb{R}^{m \times n} \quad \forall \phi \in \mathcal{B}_i.$$  

We provide a sufficient condition for the feedback gain $F$ in (25) to stochastically stabilize $(A, B)$, inspired by the gridding approach developed, e.g., in [9,10,13,18,28].

**Theorem 5.2.** Let $A \in \mathbb{H}^n_{\text{sup}}$ and $B \in \mathbb{H}^{m \times n}$. For each $i = 1, \ldots, N$, define

$$w_i(\phi) := \int_{\mathcal{B}_i} g(\phi, \ell) \mu(d\ell) \geq 0 \quad \forall \phi \in \mathcal{M},$$  

and

$$[\Gamma_{A,i} \Gamma_{B,i}] := \begin{bmatrix} \sqrt{w_1(c_i)A(c_i)} & \sqrt{w_1(c_i)B(c_i)} \\ \vdots & \vdots \\ \sqrt{w_N(c_i)A(c_i)} & \sqrt{w_N(c_i)B(c_i)} \end{bmatrix}.$$  

Assume that, for every \( i = 1, \ldots, N \), a scalar \( \kappa_i > 0 \) satisfies

\[
\begin{bmatrix}
\sqrt{w_1(\phi)}A(\phi) & \sqrt{w_1(\phi)}B(\phi) \\
\vdots & \vdots \\
\sqrt{w_N(\phi)}A(\phi) & \sqrt{w_N(\phi)}B(\phi)
\end{bmatrix} - \begin{bmatrix}
\Gamma_{A,i} & \Gamma_{B,i}
\end{bmatrix} \leq \kappa_i
\]

for \( \mu \)-almost every \( \phi \in B_i \). If there exist positive definite matrices \( R_i \in \mathbb{R}^{n \times n} \), (not necessarily symmetric) matrices \( U_i \in \mathbb{R}^{n \times n}, \bar{F}_i \in \mathbb{R}^{m \times n} \), and scalars \( \lambda_i > 0 \) such that the following LMIs hold for all \( i = 1, \ldots, N \):

\[
\begin{bmatrix}
U_i + U_i^T - R_i & 0 & U_i^T \Gamma_{A,i} + \bar{F}_i^T \Gamma_{B,i} & \kappa_i [U_i^T \bar{F}_i^T] \\
\star & \lambda_i I & \lambda_i I & 0 \\
\star & \star & \bar{R} & 0 \\
\star & \star & \star & \lambda_i I
\end{bmatrix} > 0,
\]

where \( \bar{R} := \text{diag}(R_1, \ldots, R_N) \), then the pair \((A, B)\) is stochastically stabilizable by the controller (25) with \( F_i := \bar{F}_i U_i^{-1} \).

**Proof.** From Theorem 4.2 and Proposition 4.3, \((A, B)\) is stochastically stabilizable if and only if there exist \( S \in \mathbb{H}^{n \times n}_{\text{sup}}, F \in \mathbb{H}^{m \times n}_{\text{sup}} \), and \( \epsilon > 0 \) such that the following Lyapunov inequality holds for \( \mu \)-almost every \( \phi \in \mathcal{M} \):

\[
S(\phi) - (A(\phi) + B(\phi)F(\phi))^T \left( \int_{\mathcal{M}} g(\phi, \ell)S(\ell)\mu(d\ell) \right) (A(\phi) + B(\phi)F(\phi)) \geq \epsilon I.
\]

We employ a piecewise-constant matrix function \( S \) for the Lyapunov inequality (29). Define

\[
S(\phi) := S_i \quad \forall \phi \in B_i
\]

with \( S_i > 0 \). In what follow, we prove that if the LMIs in (28) are feasible for all \( i = 1, \ldots, N \), then the Lyapunov inequality (29) holds with \( S \in \mathbb{H}^{n \times n}_{\text{sup}} \) defined by (30).

By construction, we have that

\[
\int_{\mathcal{M}} g(\phi, \ell)S(\ell)\mu(d\ell) = \sum_{i=1}^{N} w_i(\phi)S_i \quad \forall \phi \in \mathcal{M}.
\]

Note that

\[
\sum_{i=1}^{N} w_i(\phi)S_i = \begin{bmatrix}
\sqrt{w_1(\phi)}I \\
\vdots \\
\sqrt{w_N(\phi)}I
\end{bmatrix}^T S \begin{bmatrix}
\sqrt{w_1(\phi)}I \\
\vdots \\
\sqrt{w_N(\phi)}I
\end{bmatrix},
\]

where \( S := \text{diag}(S_1, \ldots, S_N) \). Substituting (31) and (32) into (29), we see that (29) can be transformed into the matrix inequality

\[
S(\phi) - (A(\phi) + B(\phi)F(\phi))^T \begin{bmatrix}
\sqrt{w_1(\phi)}I \\
\vdots \\
\sqrt{w_N(\phi)}I
\end{bmatrix}^T S \begin{bmatrix}
\sqrt{w_1(\phi)}I \\
\vdots \\
\sqrt{w_N(\phi)}I
\end{bmatrix} (A(\phi) + B(\phi)F(\phi)) \geq \epsilon I.
\]
Moreover, by the Schur complement formula, the matrix inequality (33) is equivalent to
\[
\begin{bmatrix}
S(\phi) - \epsilon I & (A(\phi) + B(\phi)F(\phi))^T \\
\ast & S
\end{bmatrix} \geq 0.
\]

Using the inequality (27), we next discretize (34), more specifically, show that if the LMIs (28) are feasible for every \( i = 1, \ldots, N \), then the matrix inequality (34) holds for \( \mu \)-almost every \( \phi \in \mathcal{M} \). In terms of the upper-right part of the matrix in (34), we obtain
\[
(A(\phi) + B(\phi)F(\phi))^T \begin{bmatrix}
\sqrt{w_1(\phi)}I \\
\vdots \\
\sqrt{w_N(\phi)}I
\end{bmatrix} = [I \ F(\phi)^T] \begin{bmatrix}
\sqrt{w_1(\phi)}A(\phi)^T \\
\sqrt{w_1(\phi)}B(\phi)^T \\
\sqrt{w_N(\phi)}A(\phi)^T \\
\sqrt{w_N(\phi)}B(\phi)^T
\end{bmatrix}.
\]
We also have from the inequality (27) that, for \( \mu \)-almost every \( \phi \in \mathcal{B}_i \), there exists
\[
[F_A \ F_B] \in \mathbb{R}^{nN \times (n+m)} \text{ with } \|[F_A \ F_B]\| < 1
\]
such that
\[
\begin{bmatrix}
\sqrt{w_1(\phi)}A(\phi)^T \\
\sqrt{w_1(\phi)}B(\phi)^T \\
\vdots \\
\sqrt{w_N(\phi)}A(\phi)^T \\
\sqrt{w_N(\phi)}B(\phi)^T
\end{bmatrix} = [\Gamma^T_{A,i} \ \Gamma^T_{B,i}] + \kappa_i [F_A^T \ F_B^T].
\]
Hence the matrix inequality (34) holds with some \( \epsilon > 0 \) for \( \mu \)-almost every \( \phi \in \mathcal{M} \) if
\[
\begin{bmatrix}
S_i & [I \ F_i^T] \left( \begin{bmatrix} \Gamma^T_{A,i} \\ \Gamma^T_{B,i} \end{bmatrix} + \kappa_i [F_A^T \ F_B^T] \right) S_i \\
\ast & S
\end{bmatrix} > 0
\]
for every \( i = 1, \ldots, N \) and for every \( [F_A \ F_B] \in \mathbb{R}^{nN \times (n+m)} \) with \( \|[F_A \ F_B]\| < 1 \).

The resulting matrix inequality (35) has the product term of the variables \( F_i \) and \( S \). To remove this product term, we employ Lemma 5.1. Let \( U_i \in \mathbb{R}^{n \times n} \) be a nonsingular matrix. Defining \( R_i := S_i^{-1} \) and \( F_i := F_i U_i \), we have from Lemma 5.1 that
\[
\begin{bmatrix}
U_i^T & 0 \\
0 & S^{-1}
\end{bmatrix} \begin{bmatrix}
S_i & [I \ F_i^T] \left( \begin{bmatrix} \Gamma^T_{A,i} \\ \Gamma^T_{B,i} \end{bmatrix} + \kappa_i [F_A^T \ F_B^T] \right) S_i \\
\ast & S
\end{bmatrix} \begin{bmatrix}
U_i & 0 \\
0 & S^{-1}
\end{bmatrix} \geq \begin{bmatrix}
U_i + U_i^T - R_i & G_i^T + H_i^T \Phi^T \\
\ast & R
\end{bmatrix},
\]
where the matrices \( R, \Phi, G_i, \) and \( H_i \) are defined by
\[
R := \text{diag}(R_1, \ldots, R_n), \quad \Phi := [F_A \ F_B]
\]
\[
G_i := [\Gamma_{A,i} \ \Gamma_{B,i}] U_i, \quad H_i := \kappa_i F_i.
\]
Finally, we obtain the sufficient LMI condition (28) by removing $\Phi$ from the matrix in the right-hand side of the inequality (36). Since $\|\Phi\| < 1$, it follows that for every $\rho_i > 0$, we have $\rho_i H_i^T (I - \Phi^T \Phi) H_i \geq 0$. Moreover,

$$
\begin{bmatrix}
U_i + U_i^T - R_i & G_i^T + H_i^T \Phi^T \\
* & R
\end{bmatrix} - \begin{bmatrix}
\rho_i H_i^T (I - \Phi^T \Phi) H_i & 0 \\
* & \rho_i I
\end{bmatrix} = V_i^T \Omega_i V_i,
$$

where

$$
V_i := \begin{bmatrix}
I & 0 \\
\Phi H_i & 0 \\
0 & I \\
-H_i & 0
\end{bmatrix}, \quad \Omega_i := \begin{bmatrix}
U_i + U_i^T - R_i & 0 & G_i^T & \rho_i H_i^T \\
* & \rho_i I & I & 0 \\
* & * & R & 0 \\
* & * & * & \rho_i I
\end{bmatrix}.
$$

Since $V_i$ is full column rank, it follows that if $\Omega_i > 0$, then $V_i^T \Omega_i V_i > 0$ and hence the matrix inequality (35) holds. Note that $\rho_i H_i^T$ in $\Omega_i$ has the product of the variables $\rho_i$ and $[U_i^T, F_i^T]$. However, using the similarity transformation with $\text{diag}(I, 1/\rho_i I, I, 1/\rho_i I)$, we see that $\Omega_i$ is similar to

$$
\begin{bmatrix}
U_i + U_i^T - R_i & 0 & G_i^T & H_i^T \\
* & \lambda_i I & \lambda_i I & 0 \\
* & * & R & 0 \\
* & * & * & \lambda_i I
\end{bmatrix},
$$

in which $\lambda_i := 1/\rho_i$ and the variables appear in a linear form, and the matrix in (37) is the one in the left-hand side of the LMIs (28). Thus if the LMIs (28) hold for all $i = 1, \ldots, N$, then the controller (25) with $F_i := F_i U_i^{-1}$ stochastically stabilizes $(A, B)$. \hfill \Box

The controller obtained in Theorem 5.2 is assumed to know to which box $B_i$ the parameter $\phi_k$ belongs for each $k \in \mathbb{Z}_+$. The following result can be used to test stabilizability and to obtain a stabilizing controller when no information about the delays is available.

**Corollary 5.3.** Under the same hypothesis as in Theorem 5.2, if there exist positive definite matrices $R_i \in \mathbb{R}^{n \times n}$, (not necessarily symmetric) matrices $U \in \mathbb{R}^{n \times n}$, $F \in \mathbb{R}^{m \times n}$, and scalars $\lambda_i > 0$ such that the following LMIs hold for all $i = 1, \ldots, N$:

$$
\begin{bmatrix}
U + U^T - R_i & 0 & U^T \Gamma_{A,i} + F^T \Gamma_{B,i} & \kappa_i [U^T \hat{F}^T] \\
* & \lambda_i I & * & \lambda_i I \\
* & * & R & 0 \\
* & * & * & \lambda_i I
\end{bmatrix} > 0,
$$

where $R := \text{diag}(R_1, \ldots, R_N)$, then the delay-independent controller $F := \hat{F} U^{-1}$ stochastically stabilizes $(A, B)$. \hfill \Box

**Proof.** This is an immediate consequence of Theorem 5.2 with $U_1 = \cdots = U_N = U$ and $\bar{F}_1 = \cdots = \bar{F}_N = \bar{F}$. \hfill \Box

We next see how conservative the proposed gridding method is. We impose the following three assumptions on discrete-time Markov jump systems.

**Assumption 5.4.** For all $c \in \mathcal{M}$ and $\delta > 0$, the $\sigma$-finite measure $\mu$ on $\mathcal{M}$ satisfies

$$
\mu(\{\phi \in \mathcal{M} : \|\phi - c\| < \delta\}) > 0.
$$
Assumption 5.5. The functions $A \in \mathbb{H}^{n\times n}_{\text{sup}}$ and $B \in \mathbb{H}^{m\times m}_{\text{sup}}$ are continuous.

Assumption 5.6. There exist $S \in \mathbb{H}^{n\times n}_{\text{sup}}$, $F \in \mathbb{H}^{m\times n}_{\text{sup}}$, and $\epsilon > 0$ such that the Lyapunov inequality (29) holds for $\mu$-almost every $\phi \in \mathcal{M}$. Moreover, for every $\epsilon_a, \epsilon_b > 0$, there exist disjoint boxes $\{B_i\}_{i=1}^N$ whose union is $\mathcal{M}$, points $c_i \in B_i$ ($i = 1, \ldots, N$), and piecewise constant functions $S_a \in \mathbb{H}^{n\times n}_{\text{sup}}$ and $F_a \in \mathbb{H}^{m\times n}_{\text{sup}}$ defined as in (30) and (25) such that the following three conditions holds:

1. $\|S - S_a\|_{\infty} < \epsilon_a$, $\|F - F_a\|_{\infty} < \epsilon_a$.
2. For all $i, j = 1, \ldots, N$, $w_j(\phi)$ defined by (26) is continuous at $\phi = c_i$.
3. For $\mu$-almost every $\phi \in B_i$ and for every $i = 1, \ldots, N$,

$$
\begin{align*}
&\|\begin{bmatrix} A(\phi) & B(\phi) \end{bmatrix} - \begin{bmatrix} A(c_i) & B(c_i) \end{bmatrix}\| < \epsilon_b \\
&\|\begin{bmatrix} \sqrt{w_1(\phi)} & \cdots & \sqrt{w_N(\phi)} \end{bmatrix} - \begin{bmatrix} \sqrt{w_1(c_i)} & \cdots & \sqrt{w_N(c_i)} \end{bmatrix}\| < \epsilon_b.
\end{align*}
$$

Assumption 5.4 holds for the standard Borel measure. The functions $A$ and $B$ defined by (8) satisfy Assumption 5.5. The first statement of Assumption 5.6, together with Theorem 4.2 and Proposition 4.3, implies that the pair $(A, B)$ is stochastically stabilizable. Since we approximate a solution of the Lyapunov inequality (29) by piecewise-constant functions in Theorem 5.2, we need 1 of Assumption 5.6. We use 2 and 3 of Assumption 5.6 together with Assumptions 5.4 and 5.5 to show that a certain inequality holds at $\phi = c_i$ for every $i = 1, \ldots, N$. These assumptions on non-zero measure and continuity are required because the Lyapunov inequality (29) is assumed to be satisfied only at $\mu$-almost everywhere.

The following proposition shows that if Assumptions 5.4–5.6 hold, then the presented gridding method will guarantee stochastic stabilizability given that approximation errors are sufficiently small.

Proposition 5.7. If Assumptions 5.4–5.6 hold, then there exist disjoint boxes $\{B_i\}_{i=1}^N$ whose union is $\mathcal{M}$ and points $c_i \in B_i$ ($i = 1, \ldots, N$) such that the LMIs in (28) are feasible.

The proof of Proposition 5.7 can be found in the Appendix.

5.2. Stochastic Detectability. Next we study the stochastic detectability of the pair $A \in \mathbb{H}^{n\times n}_{\text{sup}}$ and $C \in \mathbb{H}^{m\times n}_{\text{sup}}$. In our LQ problem, we need to check the stochastic detectability of $(C, A)$ or $(Q - WR^{-1}W^T, \hat{A})$ in (15) and (16).

Define an observer gain $L \in \mathbb{H}^{m\times n}_{\text{sup}}$ as the piecewise-constant function:

$$
L(\phi) := L_i \in \mathbb{R}^{m\times n} \quad \forall \phi \in B_i,
$$

where the disjoint boxes $\{B_i\}_{i=1}^N$ are chosen as in the previous subsection. Note that the positions of the variables $K, L$ are different between $A + BK$ (stabilizability) and $A + LC$ (detectability). Moreover, unlike the case of countable-state Markov chains (see, e.g., [6]), the duality of stochastic stabilizability and stochastic detectability is not proved yet for the case of continuous-state Markov chains. Hence we cannot use Theorem 5.2 directly, but the gridding method still provides a sufficient condition for stochastic detectability in terms of LMIs.

Theorem 5.8. Let $A \in \mathbb{H}^{n\times n}_{\text{sup}}$ and $C \in \mathbb{H}^{m\times n}_{\text{sup}}$. For each $i = 1, \ldots, N$, define $w_i$ as in Theorem 5.2 and

$$
\begin{bmatrix}
\Upsilon_{A,i} \\
\Upsilon_{C,i}
\end{bmatrix} := \begin{bmatrix} A(c_i) \\
C(c_i) \end{bmatrix}, \quad \Upsilon_{w,i} := \begin{bmatrix} \sqrt{w_1(c_i)}I & \cdots & \sqrt{w_N(c_i)}I \end{bmatrix}.
$$
Assume that, for all $i = 1, \ldots, N$, scalars $\kappa_{A,i}, \kappa_{w,i} > 0$ satisfy

\begin{align}
\| [A(\phi)] - \begin{bmatrix} \Upsilon_{A,i} \\ \Upsilon_{C,i} \end{bmatrix} \| \leq \kappa_{A,i} \tag{40a} \\
\| [\sqrt{w_1(\phi)} I \cdots \sqrt{w_N(\phi)}] - \Upsilon_{w,i} \| \leq \kappa_{w,i} \tag{40b}
\end{align}

for $\mu$-almost every $\phi \in \mathcal{B}_i$. If there exists positive definite matrices $S_i \in \mathbb{R}^{n \times n}$, (not necessarily symmetric) matrices $U_i \in \mathbb{R}^{n \times n}$, $\bar{L}_i \in \mathbb{R}^{n \times r}$, and scalars $\lambda_i, \rho_i > 0$ such that the following LMIs hold for all $i = 1, \ldots, N$:

\begin{align}
\begin{bmatrix}
U_i + U_i^\top & 0 & U_i \Upsilon_{A,i} + \bar{L}_i \Upsilon_{C,i} & \kappa_{A,i} [U_i \bar{L}_i] & 0 & \Upsilon_{w,i} S & \rho_i I \\
* & \lambda_i I & \lambda_i I & 0 & 0 & 0 & 0 \\
* & * & S_i & 0 & 0 & 0 & 0 \\
* & * & * & \lambda_i I & 0 & 0 & 0 \\
* & * & * & * & \rho_i I & \kappa_{w,i} S & 0 \\
* & * & * & * & * & S & 0 \\
* & * & * & * & * & * & \rho_i I 
\end{bmatrix} > 0, \tag{41}
\end{align}

where $S := \text{diag}(S_1, \ldots, S_N)$, then the pair $(C, A)$ is stochastically detectable by the observer gain $L$ in (39) with $L_i := U_i^{-1} \bar{L}_i$.

Proof. As in the proof of Theorem 5.2, we see from Theorem 4.2 and Proposition 4.3, that $(C, A)$ is stochastically detectable if and only if there exist $S \in \mathbb{H}^{n^+}_{\sup}$, $L \in \mathbb{H}^{n \times r}_{\sup}$, and $\epsilon > 0$ such that the following Lyapunov inequality holds for $\mu$-almost every $\phi \in \mathcal{M}$:

\begin{align}
S(\phi) - (A(\phi) + L(\phi) C(\phi))^\top \left( \int_{\mathcal{M}} g(\phi, \ell) S(\ell) \mu(d\ell) \right) (A(\phi) + L(\phi) C(\phi)) \geq \epsilon I. \tag{42}
\end{align}

We prove that if the LMIs in (41) are feasible for all $i = 1, \ldots, N$, then the Lyapunov inequality (42) holds with $S \in \mathbb{H}^{n^+}_{\sup}$ defined as a piecewise-constant matrix function $S$ in (30).

By the definitions of $L$ and $S$, we have from the Schur complement formula and (31) that the Lyapunov inequality (42) can be transformed into

\begin{align}
\left( \sum_{j=1}^N w_j(\phi) S_j \right)^{-1} A(\phi) + L_i C(\phi) \\
* \\
S_i - \epsilon I 
\end{align}

for $\mu$-almost every $\phi \in \mathcal{B}_i$ and for every $i = 1, \ldots, N$. First we remove the nonlinear term $(\sum_{j=1}^N w_j(\phi) S_j)^{-1}$ in (43), by using Lemma 5.1. If $U_i \in \mathbb{R}^{n \times n}$ is nonsingular, then the matrix in the left-hand side of this inequality is similar to

\begin{align}
\begin{bmatrix}
U_i & 0 \\
0 & I 
\end{bmatrix} \left( \sum_{j=1}^N w_j(\phi) S_j \right)^{-1} A(\phi) + L_i C(\phi) \\
* \\
S_i - \epsilon I 
\end{bmatrix} \begin{bmatrix}
U_i^\top & 0 \\
0 & I 
\end{bmatrix}.
\end{align}
Using Lemma 5.1, we see that
\[
\begin{bmatrix}
U_i & 0 \\
0 & I
\end{bmatrix}
\begin{bmatrix}
(\sum_{j=1}^{N} w_j(\phi)S_j)^{-1} & A(\phi) + L_iC(\phi) \\
* & S_i - \epsilon I
\end{bmatrix}
\begin{bmatrix}
U_i^T & 0 \\
0 & I
\end{bmatrix}
\geq
\begin{bmatrix}
U_i + U_i^T - \sum_{j=1}^{N} w_j(\phi)S_j & U_i(A(\phi) + L_iC(\phi)) \\
* & S_i - \epsilon I
\end{bmatrix}
\begin{bmatrix}
\sqrt{w_1(\phi)}I \\
\vdots \\
\sqrt{w_N(\phi)}I
\end{bmatrix}^T
S
\begin{bmatrix}
\sqrt{w_1(\phi)}I \\
\vdots \\
\sqrt{w_N(\phi)}I
\end{bmatrix},
\]
where \(S := \text{diag}(S_1, \ldots, S_N)\). Therefore, from the Schur complement formula, the matrix inequality (43) holds if
\[
\begin{bmatrix}
U_i + U_i^T & U_i(A(\phi) + L_iC(\phi)) & S_w(\phi) \\
* & S_i & 0 \\
* & * & S
\end{bmatrix}
> 0,
\]
where \(S_w(\phi) := \begin{bmatrix}
\sqrt{w_1(\phi)}I \\
\vdots \\
\sqrt{w_N(\phi)}I
\end{bmatrix} S \forall \phi \in \mathcal{M}.
\]

Let us next discretize the matrix inequality (44), by using the inequalities (40). In other words, we show that if the LMIs (41) are feasible for every \(i = 1, \ldots, N\), then the matrix inequality (44) holds for \(\mu\)-almost every \(\phi \in B_i\) and for every \(i = 1, \ldots, N\). From the inequalities (40), we see that for \(\mu\)-almost every \(\phi \in B_i\), there exist
\[
\begin{bmatrix}
\Phi_A \\
\Phi_C
\end{bmatrix} \in \mathbb{R}^{(n+r) \times n}, \ \Phi_w \in \mathbb{R}^{n \times n N} \ \text{with} \ \left\| \begin{bmatrix}
\Phi_A \\
\Phi_C
\end{bmatrix} \right\| < 1, \ \|\Phi_w\| < 1
\]
such that
\[
\begin{bmatrix}
A(\phi) \\
C(\phi)
\end{bmatrix} = \begin{bmatrix}
\Upsilon_{A,i} \\
\Upsilon_{C,i}
\end{bmatrix} + \kappa_{A,i} \begin{bmatrix}
\Phi_A \\
\Phi_C
\end{bmatrix}
\]
\[
\begin{bmatrix}
\sqrt{w_1(\phi)}I \\
\vdots \\
\sqrt{w_N(\phi)}I
\end{bmatrix} = \Upsilon_{w,i} + \kappa_{w,i} \Phi_w.
\]
Then we obtain
\[
\begin{bmatrix}
U_i + U_i^T & U_i(A(\phi) + L_iC(\phi)) & S_w(\phi) \\
* & S_i & 0 \\
* & * & S
\end{bmatrix}
= \begin{bmatrix}
U_i + U_i^T & G_{A,i} + H_{A,i} \Phi & G_{w,i} + \Phi_w H_{w,i} \\
* & S_i & 0 \\
* & * & S
\end{bmatrix},
\]
where, using \(\bar{L}_i := U_i \bar{L}_i\), we define the matrices \(\Phi, G_{A,i}, H_{A,i}, G_{w,i}\), and \(H_{w,i}\) by
\[
\begin{align*}
\Phi & := \begin{bmatrix}
\Phi_A \\
\Phi_C
\end{bmatrix}, \quad G_{A,i} := \begin{bmatrix}
U_i \\
\bar{L}_i
\end{bmatrix} \begin{bmatrix}
\Upsilon_{A,i} \\
\Upsilon_{C,i}
\end{bmatrix} \\
H_{A,i} & := \kappa_{A,i} \begin{bmatrix}
U_i \\
\bar{L}_i
\end{bmatrix}, \quad G_{w,i} := \Upsilon_{w,i} S, \quad H_{w,i} := \kappa_{w,i} S.
\end{align*}
\]
Since \(\|\Phi\|, \|\Phi_w\| < 1\), it follows that for all \(\rho_i, \rho_{w,i} > 0\), we have

\[
\rho_i H_{A,i}(I - \Phi \Phi^T)H_{A,i}^T + \rho_{w,i}(I - \Phi_w \Phi_w^T) \geq 0.
\]

Moreover,

\[
\begin{bmatrix}
U_i + U_i^T & G_{A,i} + H_{A,i} \Phi & G_{w,i} + \Phi_w H_{w,i}
\end{bmatrix}
\begin{bmatrix}
I & 0 & 0
\end{bmatrix}
- \rho_i H_{A,i}(I - \Phi \Phi^T)H_{A,i}^T + \rho_{w,i}(I - \Phi_w \Phi_w^T)
\begin{bmatrix}
0 & 0 & 0
\end{bmatrix}
= V_i^T \Omega_i V_i,
\]

where

\[
V_i := \begin{bmatrix}
I & 0 & 0 \\
\Phi^T H_{A,i} & I & 0 \\
0 & I & 0 \\
-\Phi_w & 0 & 0 \\
0 & 0 & I \\
-J & 0 & 0
\end{bmatrix},
\quad
\Omega_i := \begin{bmatrix}
I & 0 & 0 & \rho_i I & 0 & 0 & 0 & \rho_{w,i} I \\
0 & \rho_i I & 0 & 0 & 0 & 0 & 0 & 0 \\
0 & 0 & \rho_i I & 0 & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & \rho_{w,i} I & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & \rho_{w,i} I & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & 0 & \rho_{w,i} I & 0 & 0
\end{bmatrix}.
\]

Since \(V_i\) is full column rank, \(\Omega_i > 0\) leads to \(V_i \Omega_i V_i > 0\), which implies that the matrix inequality (44) is satisfied for \(\mu\)-almost every \(\phi \in B_i\). Note that \(\rho_i H_{A,i}\) has the product of the variables \(\rho_i\) and \(U_i L_i\). However, applying the similarity transformation with \(\text{diag}(I, 1/\rho_i I, I, 1/\rho_i I, I, I, I)\), we see that \(\Omega_i\) is similar to the following matrix:

\[
\begin{bmatrix}
U_i + U_i^T & 0 & G_{A,i} & H_{A,i} & 0 & 0 & G_{w,i} & \rho_{w,i} I \\
\lambda_i I & \lambda_i I & 0 & 0 & 0 & 0 & 0 & 0 \\
0 & S_i & 0 & 0 & 0 & 0 & 0 & 0 \\
0 & 0 & \lambda_i I & 0 & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & \rho_w I & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & \rho_{w,i} I & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & 0 & \rho_{w,i} I & 0 & 0
\end{bmatrix},
\]

where \(\lambda_i := 1/\rho_i\), and this matrix is the one in the left-hand side of (41). Thus, if the LMIs (41) are feasible, then \((C, A)\) is stochastically detectable. \(\square\)

As in the case of stochastic stabilizability, we see that the proposed gridding method does not introduce conservatism if approximation errors are sufficiently small.

**Assumption 5.9.** The functions \(A \in \mathbb{H}^n_{\text{sup}}\) and \(C \in \mathbb{H}^{r \times n}_{\text{sup}}\) are continuous.

**Assumption 5.10.** There exist \(S \in \mathbb{H}^n_{\text{sup}}^+, L \in \mathbb{H}^{n \times r}_{\text{sup}}\), and \(\epsilon > 0\) such that the Lyapunov inequality (42) holds for \(\mu\)-almost every \(\phi \in \mathcal{M}\). Moreover, for every \(\epsilon_a, \epsilon_b > 0\), there exist disjoint boxes \(\{B_i\}_{i=1}^N\) whose union is \(\mathcal{M}\), points \(c_i \in B_i\) (\(i = 1, \ldots, N\)), and piecewise constant functions \(S_a \in \mathbb{H}^n_{\text{sup}}^+\) and \(L_a \in \mathbb{H}^{n \times r}_{\text{sup}}\) defined by (30) and (39) such that the following three conditions hold:

1. \(\|S - S_a\|_{\text{sup}} < \epsilon_a, \quad \|L - L_a\|_{\text{sup}} < \epsilon_a\).
2. For all \(i, j = 1, \ldots, N\), \(w_j(\phi)\) defined by (26) is continuous at \(\phi = c_i\).
3. For \(\mu\)-almost every \(\phi \in \mathcal{B}_i\), and for every \(i = 1, \ldots, N\), the following inequality and (38b) are satisfied:

\[
\left\|\begin{bmatrix}
A(\phi) \\
C(\phi)
\end{bmatrix} - \begin{bmatrix}
A(c_i) \\
C(c_i)
\end{bmatrix}\right\| < \epsilon_b.
\]
Fig. 2. Sample path of $\tau(t)$.

**Proposition 5.11.** If Assumptions 5.4, 5.9, and 5.10 hold, then there exist disjoint boxes $\{B_i\}_{i=1}^N$ whose union is $\mathcal{M}$ and points $c_i \in B_i$ ($i = 1, \ldots, N$) such that the LMI's in (41) are feasible.

The proof of Proposition 5.11 can be found in the Appendix.

### 6. Numerical Example.

Consider the unstable batch reactor studied in [30], where the system matrices $A_c$ and $B_c$ in the continuous-time plant (1) are given by

$$A_c := \begin{bmatrix} 1.38 & -0.2077 & 6.715 & -5.676 \\ -0.5814 & -4.29 & 0 & 0.675 \\ 1.067 & 4.273 & -6.654 & 5.893 \\ 0.048 & 4.273 & 1.343 & -2.104 \end{bmatrix}, \quad B_c := \begin{bmatrix} 0 & 0 \\ 5.679 & 0 \\ 1.136 & -3.146 \\ 1.136 & 0 \end{bmatrix}.$$

This model is widely used as a benchmark example. We take the sampling period $h = 0.2$ and the delay interval $[\tau_{\text{min}}, \tau_{\text{max}}] = [0, 0.03]$.

We consider that the latest delay $\tau_k$ is stochastically determined by the average of the last two delays $\tau_{k-1}$ and $\tau_{k-2}$. More precisely, the sequence $\{\phi_k : k \in \mathbb{Z}_+\}$, where $\phi_k := \begin{bmatrix} \tau_k \\ \tau_{k-1} \end{bmatrix}$, is a Markov chain, and its transition probability kernel $\mathcal{G}$ is given in the following way: For every box $B = [b_{11}, b_{12}] \times [b_{21}, b_{22}]$,

$$\mathcal{G} \left( \begin{bmatrix} d_1 \\ d_2 \end{bmatrix}, B \right) = \begin{cases} \Phi_{d\text{ave}}(b_{12}) - \Phi_{d\text{ave}}(b_{11}) \\ \Phi_{d\text{ave}}(\tau_{\text{max}}) - \Phi_{d\text{ave}}(\tau_{\text{min}}) \end{cases} \quad \text{if } d_1 \in [b_{21}, b_{22}],$$

where $d_{\text{ave}} := (d_1 + d_2)/2$ and $\Phi_d(x)$ is the probability distribution function of the normal distribution with mean $d$ and standard deviation $\sigma$. Fig. 2 illustrates a sample path of the delay $\tau(t)$ with the initial data $\tau_0 = \tau_{-1} = 0.02$ and the standard derivation $\sigma = 1/100$, where $\tau(t)$ is defined by

$$\tau(t) := \tau_k \quad \forall t \in [kh + \tau_k, (k+1)h + \tau_{k+1}).$$

The weighting matrices $Q_c, R_c$ for the state and the input in (5) are the identity matrices with compatible dimensions. Using Theorem 5.2, we can confirm that $(A, B)$ in (8) is stochastically stabilizable. Additionally, $(Q - WR^{-1}W^T, A)$ in Lemma 3.2 is stochastically detectable by Theorem 5.8. Hence, by Theorems 3.6 and 4.6, we can derive an optimal controller $u^{\text{opt}}$ from the iteration of a Riccati difference equation.
Fig. 3. Ten sample paths of $\|x(t)\|^2 + \|u(t)\|^2$: The solid lines are the time responses with stochastic delays. The dotted line is the time response with no delays, for which we used the (conventional) discrete-time LQ regulator computed with the same weighting matrices.

Time responses are computed for an deterministic initial state

$$x(0) = \begin{bmatrix} 1 \\ -1 \\ 2 \\ -2 \end{bmatrix}, \quad u_{-1} = \begin{bmatrix} 0 \\ 0 \end{bmatrix}.$$  

Fig. 3 depicts ten sample paths of the performance function $\|x(t)\|^2 + \|u(t)\|^2$, where initial delays $\tau_0, \tau_{-1}$ are uniformly distributed in the interval $[\tau_{\text{min}}, \tau_{\text{max}}]$ and the standard deviation $\sigma$ of the probability distribution function $\Phi_d$ in (46) is given by $\sigma = 1/100$.

We observe that the time responses with small initial delays are similar to the response with no delays by the conventional discrete-time LQ regulator with the same weighting matrices. Although larger delays degrade the control performance, the optimal controller achieves almost the same decrease rate for every initial delay.

7. Concluding Remarks. We provided the design of delay-dependent optimal controllers for sampled-data systems whose sensor-to-controller delays are stochastically determined by the last few delays. Our optimal control problem was reduced to the LQ problem of discrete-time Markov jump systems. We can efficiently compute an optimal controller by iteratively solving a Riccati difference equation. Moreover, we derived the sufficient conditions for stochastic stabilizability and detectability in terms of LMIs via the gridding approach. From these conditions, we can also construct stabilizing controllers and state observers. Future work will focus on addressing more general systems by incorporating packet losses and output feedback.

To solve the stabilization problem of discrete-time systems with continuous-valued Markovian jumping parameters, we approximated the function $S$ in the Lyapunov inequality (23) by a piecewise constant function. Another interesting study is to approximate the system

$$\{(A(\phi), B(\phi)) : \phi \in \mathcal{M}\}$$

by

$$\{(A(c_i), B(c_i)) : c_i \text{ is the center of } \mathcal{B}_i, \quad i = 1, \ldots, N\}$$
and then to consider the stabilization problem of systems having discrete-valued Markovian jumping parameters but time-varying uncertainty in the coefficient matrices.

**Appendix A. Proof of Proposition 5.7.** First, we prove that if \( \epsilon_a > 0 \) in Assumption 5.6 is sufficiently small, then there exists \( \epsilon_1 > 0 \) such that, for \( \mu \)-almost every \( \phi \in \mathcal{M} \),

\[
\begin{pmatrix}
S_a(\phi) & (A(\phi) + B(\phi)F_a(\phi))^\top \\
* & S^{-1}
\end{pmatrix}
\begin{bmatrix}
\sqrt{w_1(\phi)} I \\
\vdots \\
\sqrt{w_N(\phi)} I
\end{bmatrix}^\top > \epsilon_1 I,
\]

where \( S := \text{diag}(S_1, \ldots, S_N) \).

Since

\[
\int_{\mathcal{M}} g(\phi, \ell) \mu(d\ell) = 1
\]

for all \( \phi \in \mathcal{M} \), it follows from 1 of Assumption 5.6 that

\[
\left\| \int_{\mathcal{M}} g(\phi, \ell)(S(\ell) - S_a(\ell)) \mu(d\ell) \right\| \leq \int_{\mathcal{M}} g(\phi, \ell)\|S(\ell) - S_a(\ell)\| \mu(d\ell) \leq \epsilon_a
\]

for \( \mu \)-almost every \( \phi \in \mathcal{M} \). We obtain

\[
\left\| \begin{pmatrix}
S(\phi) - (A(\phi) + B(\phi)F(\phi))^\top \\
\int_{\mathcal{M}} g(\phi, \ell)S(\ell) \mu(d\ell)
\end{pmatrix}
\begin{pmatrix}
A(\phi) + B(\phi)F(\phi)
\end{pmatrix}
\right\|
\leq (1 + \|A + BF\|_2^2) \epsilon_a
\]

for \( \mu \)-almost every \( \phi \in \mathcal{M} \). Therefore, if we set \( \tau_a > 0 \) to be a value with

\[
\tau_a < \frac{\epsilon}{1 + \|A + BF\|_\infty^2},
\]

then \( \epsilon_2 := \epsilon - (1 + \|A + BF\|_\infty^2)\tau_a > 0 \) satisfies

\[
(49) \quad S_a(\phi) - (A(\phi) + B(\phi)F(\phi))^\top \left( \int_{\mathcal{M}} g(\phi, \ell)S_a(\ell) \mu(d\ell) \right) (A(\phi) + B(\phi)F(\phi)) > \epsilon_2 I
\]

for \( \mu \)-almost every \( \phi \in \mathcal{M} \) and for every \( \epsilon_a \in (0, \tau_a) \). Substituting (31) and (32) into (49) and applying the Schur complement formula, we have that there exists \( \epsilon_3 > 0 \) such that

\[
\begin{pmatrix}
S_a(\phi) & (A(\phi) + B(\phi)F(\phi))^\top w(\phi) \\
* & S^{-1}
\end{pmatrix} > \epsilon_3 I
\]

for \( \mu \)-almost every \( \phi \in \mathcal{M} \), where

\[
w(\phi) := \begin{bmatrix}
\sqrt{w_1(\phi)} I \\
\vdots \\
\sqrt{w_N(\phi)} I
\end{bmatrix}.
\]
We obtain
\[
\begin{bmatrix}
S_a(\phi) & (A(\phi) + B(\phi)F_a(\phi))^Tw(\phi) \\
* & S^{-1}
\end{bmatrix}
= \begin{bmatrix}
S_a(\phi) & (A(\phi) + B(\phi)F(\phi))^Tw(\phi) \\
* & S^{-1}
\end{bmatrix} - \begin{bmatrix}
0 & (F(\phi) - F_a(\phi))^TB(\phi)^Tw(\phi) \\
0 & 0
\end{bmatrix}.
\]

By (48), \( \{w_i\}_{i=1}^N \) satisfies
\[
\sum_{i=1}^N w_i(\phi) = \int_{B_i} g(\phi, \ell)\mu(d\ell) = \int_M g(\phi, \ell)\mu(d\ell) = 1 \quad \forall \phi \in \mathcal{M}.
\]

It follows that \( \|w(\phi)\| = 1 \) for every choice of disjoint boxes \( \{B_i\}_{i=1}^N \) and for every \( \phi \in \mathcal{M} \). Hence, if \( \|F - F_a\|_\infty < \epsilon_a \), then
\[
\begin{bmatrix}
0 & (F(\phi) - F_a(\phi))^TB(\phi)^Tw(\phi) \\
0 & 0
\end{bmatrix} \leq \|B\|_\infty \epsilon_a
\]
for \( \mu \)-almost every \( \phi \in \mathcal{M} \). If \( \epsilon_a \in (0, \tau_a) \) satisfies \( \epsilon_a < \epsilon_3/\|B\|_\infty \), then the desired inequality (47) holds with \( \epsilon_1 := \epsilon_3 - \|B\|_\infty \epsilon_a \).

Let us next derive the feasibility of the LMIs in (28) from the inequality (47). Using the similarity transformation \( \text{diag}(S^{-1}_a, I) \), we find that there exists \( \epsilon_4 > 0 \) such that
\[
\begin{bmatrix}
R_i & \tilde{F}_i^T \\
* & R
\end{bmatrix} \begin{bmatrix}
A(\phi)^T \\
B(\phi)^T
\end{bmatrix}w(\phi) > \epsilon_4 I
\]
for \( \mu \)-almost every \( \phi \in B_i \) and for every \( i = 1, \ldots, N \), where \( R_i := S^{-1}_i, R := S^{-1} \) and \( \tilde{F}_i := F_iR_i \). Assumption 5.4 on the non-zero property of \( \mu \) and Assumptions 5.5 and 5.6 on the continuity of \( A, B, w \) at \( \phi = \phi_i \) imply that for every \( \epsilon_5 \in (0, \epsilon_4) \),
\[
\begin{bmatrix}
R_i & R_i^\top \Gamma_{A,i} + \tilde{F}_i^\top \Gamma_{B,i} \\
* & R
\end{bmatrix} > \epsilon_5 I \quad \forall i = 1, \ldots, N,
\]
where \( \Gamma_{A,i} \) and \( \Gamma_{B,i} \) are defined as in Theorem 5.2.

By 3 of Assumption 5.6, we see that for \( \mu \)-almost every \( \phi \in B_i \) and for every \( i = 1, \ldots, N \),
\[
\|w(\phi)^T[A(\phi) - B(\phi)] - [\Gamma_{A,i} - \Gamma_{B,i}]\| \leq (1 + \|A\|_\infty)\epsilon_b =: \epsilon_\kappa.
\]

On the other hand, the the LMI in (28) with \( U_i = R_i \) and \( \kappa_i = \epsilon_\kappa \) is equivalent to
\[
\begin{bmatrix}
R_i & R_i^\top \Gamma_{A,i} + \tilde{F}_i^\top \Gamma_{B,i} \\
* & R
\end{bmatrix} > 0.
\]

By the Schur complement formula, if \( \tilde{F}_i = F_iR_i \), then the above inequality is equivalent to
\[
\begin{bmatrix}
R_i & R_i^\top \Gamma_{A,i} + \tilde{F}_i^\top \Gamma_{B,i} \\
* & R
\end{bmatrix} - \begin{bmatrix}
(\epsilon_\kappa/\lambda_i) \cdot R_i^\top (I + F_i^\top F_i)R_i & 0 \\
0 & \lambda_i I
\end{bmatrix} > 0.
\]
From (49), we know that
\[
\|R_i\| = \|S_i^{-1}\| < \frac{1}{\epsilon_2} \quad \forall i = 1, \ldots, N.
\]
Therefore, if \(\lambda_i > 0\) and \(\epsilon_\kappa > 0\) satisfy
\[
\lambda_i < \epsilon_5, \quad \epsilon_\kappa < \epsilon_2 \sqrt{\frac{\lambda_i \epsilon_5}{1 + (\|F\|_\infty + \epsilon_a)^2}} \quad \forall i = 1, \ldots, N,
\]
then the inequality (51) leads to the desired conclusion (53). If we choose sufficiently small \(\epsilon_b > 0\), then \(\epsilon_\kappa > 0\) satisfies the above inequality. This completes the proof.

**Appendix B. Proof of Proposition 5.11.** By the same discussion as in the proof of Proposition 5.7, there exists \(\epsilon_1 > 0\) such that
\[
\left(\sum_{j=1}^{N} w_j(\phi)S_j\right)^{-1} A(\phi) + L_i C(\phi) \succ \epsilon_1 I
\]
for \(\mu\)-almost every \(\phi \in B_i\) and for every \(i = 1, \ldots, N\). Using the similarity transformation \(\text{diag}\left(\sum_{j=1}^{N} w_j(\phi)S_j, I\right)\) and applying the Schur complement formula, we find that there exists \(\epsilon_2 > 0\) such that
\[
\begin{bmatrix}
2U_1(\phi) & U_1(\phi)(A(\phi) + L_i C(\phi)) & S_w(\phi) \\
* & S_i & 0 \\
* & * & S
\end{bmatrix} \succ \epsilon_2 I
\]
for \(\mu\)-almost every \(\phi \in B_i\) and for every \(i = 1, \ldots, N\), where \(U_1(\phi) := \sum_{j=1}^{N} w_j(\phi)S_j\), \(S := \text{diag}(S_1, \ldots, S_N)\), and \(S_w\) is defined as in (45). Assumption 5.4 on the non-zero property of \(\mu\) and Assumptions 5.9 and 5.10 on the continuity of \(A, C, w\) at \(\phi = c_i\) lead to
\[
(54) \quad \left[\begin{array}{cc}
2U_i & U_i \Upsilon_{A,i} + \bar{L}_i \Upsilon_{C,i} \\
* & S_i \\
* & * \\
* & S
\end{array}\right] \succ \epsilon_3 I \quad \forall \epsilon_3 \in (0, \epsilon_2), \; \forall i = 1, \ldots, N,
\]
where \(\Upsilon_{A,i}, \Upsilon_{C,i}, \Upsilon_w\) are defined as in Theorem 5.8 and
\[
(55) \quad U_i := U(c_i) = \sum_{j=1}^{N} w_j(c_i)S_j, \quad \bar{L}_i := U_i L_i.
\]

The LMIs (41) with \(\kappa_{A,i} = \epsilon_b = \kappa_{w,i}\) are equivalent to
\[
\begin{bmatrix}
U_i + U_i^T & U_i \Upsilon_{A,i} + \bar{L}_i \Upsilon_{C,i} & S_w & \epsilon_b \left[ U_i \bar{L}_i \right] & \lambda_i I & 0 & 0 & \rho_i I \\
* & S_i & 0 & 0 & \lambda_i I & 0 \\
* & * & S & 0 & \epsilon_b S & 0 \\
* & * & * & \lambda_i I & 0 & 0 & 0 \\
* & * & * & * & \rho_i I & 0 \\
* & * & * & * & * & \lambda_i I & 0 \\
* & * & * & * & * & * & \rho_i I
\end{bmatrix} \succ 0
\]
and hence to
\[
\begin{bmatrix}
U_i + U_i^T \\
* \\
\end{bmatrix}
\begin{bmatrix}
U_i \Gamma_{A,i} + \bar{L}_i \Gamma_{C,i} \\
\Sigma_i \\
\end{bmatrix}
\begin{bmatrix}
Y_{w,i} S_i \\
0 \\
S \\
\end{bmatrix}
- \begin{bmatrix}
\frac{\epsilon^2}{2} (U_i U_i^T + \bar{L}_i \bar{L}_i^T) + \rho_i I \\
0 \\
0 \\
\frac{\epsilon^2}{\rho_i} S^2 \\
\end{bmatrix} > 0
\]

by the Schur complement formula. From (50), we know that $U_i$ and $\bar{L}_i$ defined by (55) satisfy
\[
\|U_i\| \leq \|U\|_{\infty} \leq \sum_{j=1}^{N} w_j(\phi) \|S_j\| \leq \|S\|_{\infty} \leq \|S\|_{\infty} + \epsilon_a
\]
\[
\|\bar{L}_i\| \leq \|U_i\| \cdot \|L_i\| \leq (\|S\|_{\infty} + \epsilon_a) \cdot (\|L\|_{\infty} + \epsilon_a)
\]
for every $i = 1, \ldots, N$. Moreover, $S = \text{diag}(S_1, \ldots, S_N)$ satisfies $\|S\| \leq \|S\|_{\infty} + \epsilon_a$. Hence, for the matrices $U_i$ and $\bar{L}_i$ in (55) and $S = \text{diag}(S_1, \ldots, S_N)$, there exist $\lambda_i > 0$, $\rho_i > 0$, $\epsilon_b > 0$ such that
\[
\begin{bmatrix}
\frac{\epsilon^2}{2} (U_i U_i^T + \bar{L}_i \bar{L}_i^T) + \rho_i I \\
0 \\
0 \\
\frac{\epsilon^2}{\rho_i} S^2 \\
\end{bmatrix} < \frac{\epsilon_3}{2} I.
\]
By (54), the desired LMIs in (41) are satisfied.

REFERENCES

[1] B. D. O. Anderson and J. Moore, Optimal Control-Linear Quadratic Methods, Englewood Cliffs: Prentice Hall, 1990.
[2] M. B. G. Cloosterman, N. van de Wouw, W. P. M. H. Heemels, and H. Nijmeijer, Controller synthesis for networked control systems, IEEE Trans. Automat. Control, 54 (2009), pp. 1575–1580.
[3] S. Cong, Y. Ge, Q. Chen, M. Jiang, and W. Shang, DTHMM based delay modeling and prediction for networked control systems, J. Systems Eng. Electron, 21 (2010), pp. 1014–1024.
[4] O. L. V. Costa and D. Z. Figueiredo, Stochastic stability of jump discrete-time linear systems with Markov chain in a general Borel space, IEEE Trans. Automat. Control, 59 (2014), pp. 223–227.
[5] O. L. V. Costa and D. Z. Figueiredo, LQ control of discrete-time jump systems with Markov chain in a General Borel space, IEEE Trans. Automat. Control, 60 (2015), pp. 2530–2535.
[6] O. L. V. Costa, M. D. Fragoso, and R. P. Marques, Discrete-Time Markovian Jump Linear Systems, London: Springer, 2005.
[7] R. F. Curtain and H. J. Zwart, An Introduction to Infinite-Dimensional Linear Systems Theory, New York: Springer, 1995.
[8] B. Demirel, C. Briat, and K. Johansson, Deterministic and stochastic approaches to supervisory control design for networked systems with time-varying communication delays, Nonlinear Anal.: Hybrid Systems, 10 (2013), pp. 94–110.
[9] M. C. F. Donkers, W. P. M. H. Heemels, D. Bernardini, A. Bemporad, and V. Shneer, Stability analysis of stochastic networked control systems, Automatica, 48 (2012), pp. 917–925.
[10] M. C. F. Donkers, W. P. M. H. Heemels, N. van de Wouw, and L. Hetel, Stability analysis of networked control systems using a switched linear systems approach, IEEE Trans. Automat. Control, 56 (2011), pp. 2101–2115.
[11] S. Elvira-Ceja, E. N. Sanchez, and S. Jagannathan, Stochastic inverse optimal control of unknown linear networked control system in the presence of random delays and packet losses, in Proc. ACC’15, 2015.
[12] C. Foiaş, H. Özbay, and A. Tannenbaum, Robust Control of Infinite Dimensional Systems: Frequency Domain Methods, London: Springer, 1996.
[13] H. Fujioka, A discrete-time approach to stability analysis of systems with aperiodic sample-and-hold devices, IEEE Trans. Automat. Control, 54 (2009), pp. 2440–2445.

[14] J. C. Geromel, R. H. Korogui, and J. Bernussou, $\mathcal{H}_2$ and $\mathcal{H}_\infty$ robust output feedback control for continuous time polytopic systems, IET Control Theory Appl., 1 (2007), pp. 1541–1549.

[15] A. Ghanaim and G. Frey, Markov modeling of delays in networked automation and control systems using colored Petri net models simulation, in Proc. 18th IFAC WC, 2011.

[16] C. Hermann, R. Kissling, and A. Donner, A delay model for satellite constellation networks with inter-satellite links, in Proc. IWSSC’09, 2009.

[17] J. P. Hespanha, N. Nachtsheimrizi, and Y. Xu, A survey of recent results in networked control systems, Proc. IEEE, 95 (2007), pp. 138–162.

[18] L. Hetel, J. Daafouz, J.-P. Richard, and M. Jungers, Delay-dependent sampled-data control based on delay estimates, Systems Control Lett., 60 (2011), pp. 146–150.

[19] S. Hirche, C. Chen, and M. Buß, Performance oriented control over networks: Switching controllers and switched time delay, Asian J. Control, 10 (2008), pp. 24–33.

[20] D. Huang and S.-K. Nguang, State feedback control of uncertain networked control systems with random time delays, IEEE Trans. Automat. Control, 53 (2008), pp. 829–834.

[21] K. Kobayashi and K. Hiraishi, Modeling and design of networked control systems using a stochastic switching systems approach, IEEEJ Trans. Electr. Electron. Eng., 9 (2014), pp. 56–61.

[22] I. V. Kolmanovsky and T. L. Maizenberg, Optimal control of continuous-time linear systems with a time-varying, random delay, Systems Control Lett., 44 (2001), pp. 119–126.

[23] I. Kordonis and G. P. Papavassilopoulos, On stability and LQ control of MJLSS with a Markov chain with general state space, IEEE Trans. Automat. Control, 59 (2014), pp. 535–540.

[24] K. Kosuge, H. Murayama, and K. Takeo, Bilateral feedback control of telesmanipulators via computer network, in Proc. IROS 96, 1996.

[25] A. Kruszewski, W.-J. Jiang, E. Fridman, and A. Richard, J.-P. Toguyeni, A switched system approach to exponential stabilization through communication network, IEEE Trans. Control Systems Tech., 20 (2012), pp. 887–900.

[26] J. Nilsson and B. Bernhardsson, LQG control over a Markov communication network, in Proc. 36th IEEE CDC, 1997.

[27] J. Nilsson, B. Bernhardsson, and B. Wittenmark, Stochastic analysis and control of real-time systems with random time delays, Automatica, 34 (1998), pp. 57–64.

[28] H. Oishi, H. Fujioka, Stability and stabilization of aperiodic sampled-data control systems using robust linear matrix inequalities, Automatica, 46 (2010), pp. 1327–1333.

[29] D. B. Percival and A. T. Walden, Spectral Analysis for Physical Application, Cambridge: Cambridge Univ. Press, 1993.

[30] H. H. Rosenbrock, Computer-Aided Control System Design, New York: Academic Press, 1974.

[31] F. R. P. Safaei, S. G. Ghiocel, J. P. Hespanha, and J. H. Chow, Stability of an adaptive switched controller for power system oscillation damping using remote synchrophasor signals, in Proc. 53rd IEEE CDC, 2014.

[32] M. Schatzman, Numerical Analysis: A Mathematical Introduction, Oxford : Oxford Univ. Press, 2002.

[33] Y. Shi and B. Yu, Output feedback stabilization of networked control systems with random delays modeled by Markov chains, IEEE Trans. Automat. Control, 54 (2009), pp. 1668–1674.

[34] H. Shousong and Z. Qixin, Stochastic optimal control and analysis of stability of networked control systems with long delay, Automatica, 39 (2003), pp. 1877–1884.

[35] D. Srinivasagupta, H. Schättler, and B. Joseph, Time-stamped model predictive control: an algorithm for control of processes with random delays, Comput. Chem. Eng., 28 (2004), pp. 1337–1346.

[36] Y. Sun and S. Qin, Stability analysis and controller design for networked control systems with random time delay, Int. J. Systems Science, 42 (2011), pp. 359–367.

[37] M. Wakaiki, M. Ogura, and J. P. Hespanha, Linear quadratic control for sampled-data systems with stochastic delays, in Proc. ACC’17, 2017.

[38] Z. Wang, X. Wang, and L. Liu, Stochastic optimal linear control of wireless networked control systems with delays and packet losses, IET Control Theory Appl., 10 (2016), pp. 742–751.

[39] H. Xu, S. Jagannathanam, and F. L. Lewis, Stochastic optimal control of unknown linear networked control system in the presence of random delays and packet losses, Automatica, 48 (2012), pp. 1017–1030.

[40] T. C. Yang, Networked control system: a brief survey, IEE Proc. Control Theory Appl., 153
[41] L. Zhang, Y. Shi, T. Chen, and B. Huang, *A new method for stabilization of networked control systems with random delays*, IEEE Trans. Automat. Control, 50 (2005), pp. 1177–1181.