

Random Dynamic Systems with Switching and Applications

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This talk reports some of our recent findings involving joint work with

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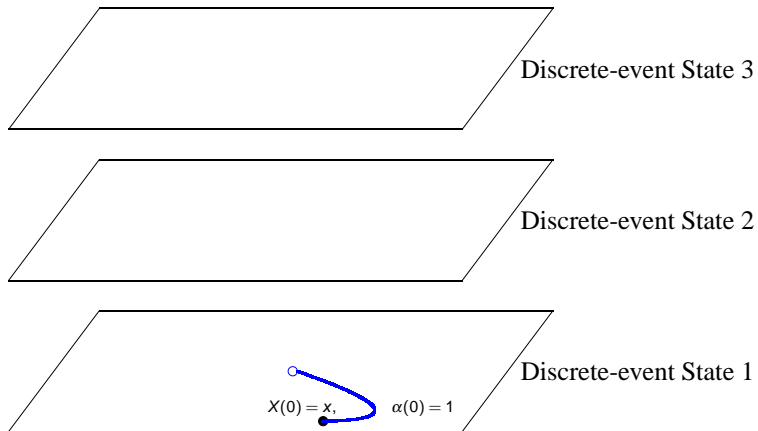
C. Yuan (Univ. of Swansea)

Outline

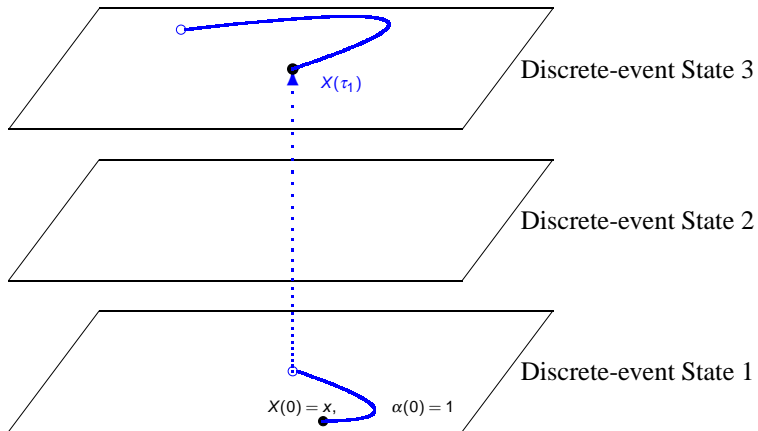
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 - Recurrence
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Switching Random Dynamic Systems

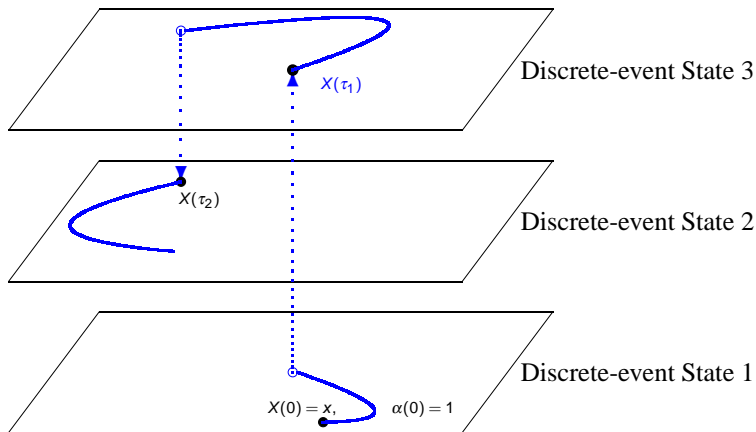
Switched Dynamic System: An Illustration



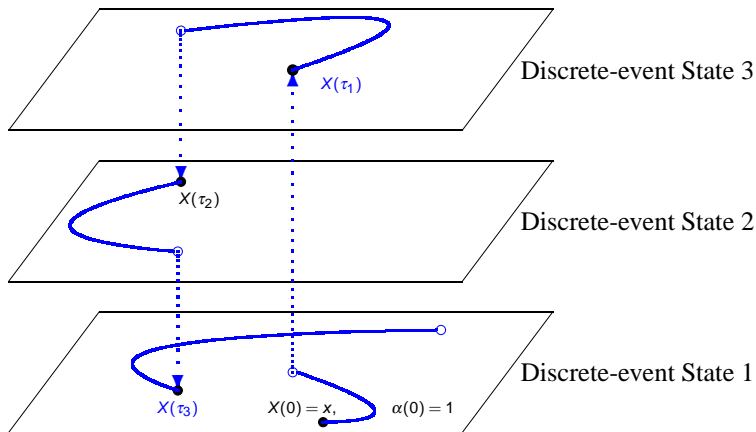
Switched Dynamic System: An Illustration



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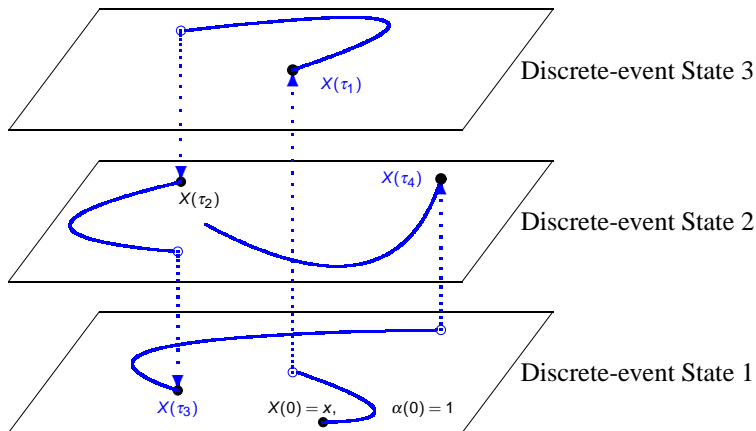


Figure: A “Sample Path” of A Switching Dynamic System $(X(t), \alpha(t))$.

Main Features

- **continuous dynamics** & **discrete events** coexist
- **switching** is used to model **random environment** or other **random factors** that cannot be formulated by the usual differential equations
- problems naturally arise in applications such as distributed, cooperative, and non-cooperative games, wireless communication, target tracking, reconfigurable sensor deployment, autonomous decision making, learning, etc.
- traditional ODE or SDE models are no longer adequate
- **non-Gaussian** distribution

Switching Diffusions

$$\mathcal{M} = \{1, \dots, m\}$$

$\alpha(\cdot)$: taking values in \mathcal{M} .

$w(t)$: d -dimensional standard Brownian motion

$$b(\cdot, \cdot) : \mathbb{R}^r \times \mathcal{M} \mapsto \mathbb{R}^r$$

$$\sigma(\cdot, \cdot) : \mathbb{R}^r \times \mathcal{M} \mapsto \mathbb{R}^r \times \mathbb{R}^d$$

$$\begin{aligned} dX(t) &= b(X(t), \alpha(t))dt + \sigma(X(t), \alpha(t))dw(t), \\ X(0) &= x, \alpha(0) = \alpha, \end{aligned} \tag{1}$$

$$\mathbf{P}\{\alpha(t + \Delta) = j | \alpha(t) = i, (X(s), \alpha(s)), s \leq t\} = q_{ij}(X(t))\Delta + o(\Delta), \quad i \neq j. \tag{2}$$

Formulation (cont.)

$Q(\mathbf{x}) = (q_{ij}(\mathbf{x}))$: generator associated with $\alpha(t)$ satisfying

$$q_{ij}(\mathbf{x}) \geq 0, \text{ if } j \neq i, \text{ and } \sum_{j=1}^m q_{ij}(\mathbf{x}) = 0, \quad i = 1, 2, \dots, m$$

\mathcal{L} : generator of $(X(t), \alpha(t))$. For each $i \in \mathcal{M}$, and any $g(\cdot, i) \in C^2(\mathbb{R}^r)$,

$$\mathcal{L}g(\mathbf{x}, i) = \frac{1}{2} \text{tr}(a(\mathbf{x}, i) \nabla^2 g(\mathbf{x}, i)) + b'(\mathbf{x}, i) \nabla g(\mathbf{x}, i) + Q(\mathbf{x})g(\mathbf{x}, \cdot)(i) \quad (3)$$

where

$\nabla g(\cdot, i)$ & $\nabla^2 g(\cdot, i)$: gradient & Hessian of $g(\cdot, i)$,

$a(\mathbf{x}, i) = \sigma(\mathbf{x}, i) \sigma'(\mathbf{x}, i)$,

$Q(\mathbf{x})g(\mathbf{x}, \cdot)(i) = \sum_{j=1}^m q_{ij}(\mathbf{x}) g(\mathbf{x}, j)$.

Main Difficulty

- Consider $(X(t), \alpha(t))$ with two different initial data $(X(0), \alpha(0)) = (x, \alpha)$ & $(X(0), \alpha(0)) = (y, \alpha)$, $y \neq x$.
- Since $Q(x)$ depends on x ,
 $\alpha^{x,\alpha}(t) \neq \alpha^{y,\alpha}(t)$ **infinitely often** even though
 $\alpha^{x,\alpha}(0) = \alpha^{y,\alpha}(0) = \alpha$.

Associated Poisson Measure

- $\Delta_{ij}(x)$: left closed, right open intervals of \mathbb{R} , with length $q_{ij}(x)$
- $h: \mathbb{R}^r \times \mathcal{M} \times \mathbb{R} \mapsto \mathbb{R}$:

$$h(x, i, z) = \sum_{j=1}^m (j - i) I_{\{z \in \Delta_{ij}(x)\}}. \quad (4)$$

$$d\alpha(t) = \int_{\mathbb{R}} h(X(t), \alpha(t-), z) p(dt, dz), \quad (5)$$

where

$p(dt, dz)$: a Poisson random measure with intensity $dt \times m(dz)$,
 m : the Lebesgue measure on \mathbb{R} ,
 $p(\cdot, \cdot)$ independent of $w(\cdot)$.

Generalized Itô Lemma

If $V \in C^{1,2}(\mathbb{R}_+ \times \mathbb{R}^r \times \mathcal{M})$, then for any $t \geq 0$:

$$V(t, X(t), \alpha(t)) = V(0, X(0), \alpha(0)) + \int_0^t \left[\frac{\partial}{\partial s} + \mathcal{L} \right] V(s, X(s), \alpha(s)) ds + M_1(t) + M_2(t), \quad (6)$$

where

$$M_1(t) = \int_0^t \langle \nabla V(s-, X(s-), \alpha(s-)), \sigma(X(s-), \alpha(s-)) dw(s) \rangle.$$

$$M_2(t) = \int_0^t \int_{\mathbb{R}} [V(s-, X(s-), \alpha(s-) + h(X(s-), \alpha(s-), z)) - V(s-, X(s-), \alpha(s-))] \mu(ds, dz),$$

$\mu(ds, dz) = p(ds, dz) - ds \times m(dz)$ is a martingale measure.

An Example

Consider

$$\dot{x}(t) = A(\alpha(t))x(t) \quad (7)$$

where $\alpha(t)$ has two states $\{1, 2\}$,

$$A(1) = \begin{bmatrix} 0 & -1 \\ 1 & 0 \end{bmatrix}, \quad A(2) = \begin{bmatrix} -1 & 2 \\ -2 & -1 \end{bmatrix}, \quad Q = \begin{bmatrix} -1 & 1 \\ 2 & -2 \end{bmatrix},$$

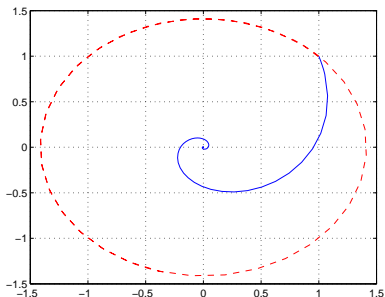
Associated with the hybrid system, there are two ODEs

$$\dot{x}(t) = A(1)x(t), \quad \text{and} \quad (8)$$

$$\dot{x}(t) = A(2)x(t) \quad (9)$$

switching back and forth according to $\alpha(t)$.

Phase Portrait of the Components



Phase portraits of the 'component' with a center (in dashed line) and the 'component' with a stable node (in solid line) with the same initial condition $x_0 = [1, 1]'$

Phase Portrait of Hybrid System

The phase portrait is given below.

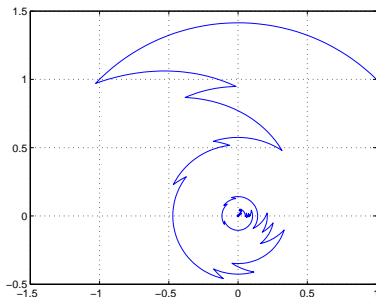


Figure: Switching linear system: Phase portrait of (7) with $x_0 = [1, 1]'$.

Switching ODEs

This example belongs to a more general class of hybrid systems:

$$\dot{x}(t) = f(x(t), \alpha(t)), \quad x(0) = x, \quad \alpha(0) = \alpha, \quad (10)$$

and $Q(x)$ is an x -dependent generator.

By using the Liapunov exponent, we can also obtain necessary and sufficient conditions for stability and instability.

- Some new results different from the usual [Hartman-Grobman theorem](#)

with Zhu & Song, Quarterly Appl. Math (2009)

Seemingly Not Much Different from Diffusions without Switching?

Q: When we have a coupled system with $\mathcal{M} = \{1,2\}$ and two stable linear systems, do we always get a stable system?

Seemingly Not Much Different from Diffusions without Switching?

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Consider $\dot{x} = A(\alpha(t))x + B(\alpha(t))u(t)$, and a state feedback $u(t) = K(\alpha(t))x(t)$. Then one gets

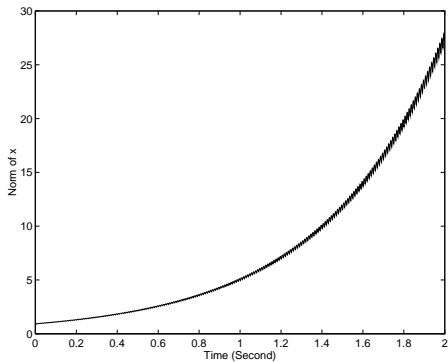
$$\dot{x} = [A(\alpha(t)) - B(\alpha(t))K(\alpha(t))]x.$$

Suppose that $\alpha(t) \in \{1, 2\}$ such that

$$A(1) - B(1)K(1) = \begin{bmatrix} -100 & 20 \\ 200 & -100 \end{bmatrix}, \quad A(2) - B(2)K(2) = \begin{bmatrix} -100 & 200 \\ 20 & -100 \end{bmatrix}.$$

The two feedback systems are stable individually. But if we choose $\alpha(t)$ so that it switches at $k\eta$, where $\eta = 0.01$. Then the resulting system is unstable.

The hybrid system is unstable



[L.Y. Wang, P.P. Khargonecker, and A. Beydoun, 1999, deterministic switching system]

Why is the system unstable?

$$\frac{1}{2}[A(1) - B(1)K(1) + A(2) - B(2)K(2)] = \frac{1}{2} \begin{bmatrix} -200 & 220 \\ 220 & -200 \end{bmatrix}$$

is an unstable matrix.

The **averaging effect** dominates the dynamics.

with Zhang, Springer book, 2nd Ed. 2013

- Consider a system

$$\dot{x}^\varepsilon(t) = b(x^\varepsilon(t), \alpha^\varepsilon(t)), \quad \alpha^\varepsilon(t) \sim Q/\varepsilon \quad (11)$$

- each $\dot{x}(t) = b(x(t), i)$, $i \in \mathcal{M}$ is stable.
- Q irreducible
- $x^\varepsilon(\cdot) \Rightarrow x(\cdot)$ such that

$$\dot{x}(t) = \bar{b}(x(t)), \quad \bar{b}(x) = \sum_{i \in \mathcal{M}} v_i b(x, i). \quad (12)$$

- System (12) is unstable.
- Use perturbed Liapunov function to show that (11) is unstable.

Regime-switching Diffusion Examples

Average Cost Per Unit Time Problem

Consider a controlled switching diffusion $(X(t), \alpha(t))$ (drift and diffusion coefficients also depend on a control u).

Aim: find $u^*(\cdot)$ so

$$\lim_{T \rightarrow \infty} E \frac{1}{T} \int_0^T L(X(t), \alpha(t), u(t)) dt$$

is minimized.

Questions: Does there exist an ergodic measure? If yes, can we replace the instantaneous measure by the ergodic one?

Two-time-scale Markov Chains

- Two-time-scale Markov chain $\alpha(t)$ with $\varepsilon > 0$ small,

$$Q(t) = Q^\varepsilon(t) = \frac{\tilde{Q}(t)}{\varepsilon} + \hat{Q}(t). \quad (13)$$

- ▶ $\tilde{Q}(t), \hat{Q}(t)$ are generators of Markov chains.
 - ▶ $\tilde{Q}(t) = \text{diag}(\tilde{Q}^1(t), \dots, \tilde{Q}^l(t))$ nearly decomposable
 - ▶ $\mathcal{M} = \mathcal{M}_1 \cup \dots \cup \mathcal{M}_l; \mathcal{M}_i = \{s_{i1}, \dots, s_{im_i}\}$
- Consider the scaled sequence

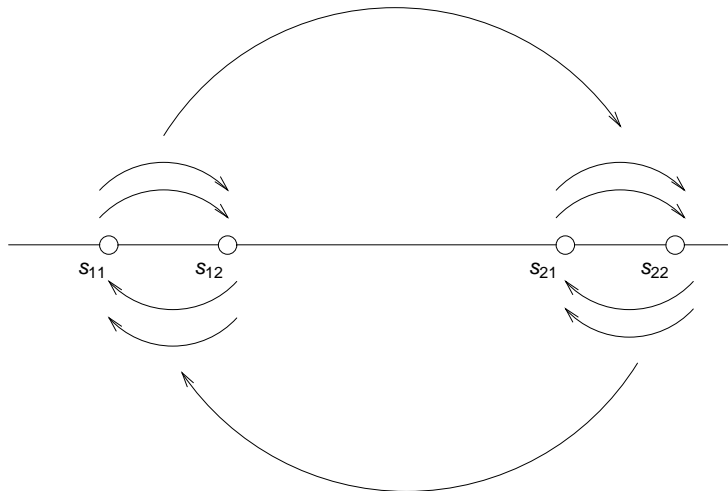
$$\frac{1}{\sqrt{\varepsilon}} \int_0^t I_{\{\alpha^\varepsilon(u)=s_{ij}\}} - v_j^i(u) I_{\{\alpha^\varepsilon(u) \in \mathcal{M}_i\}} du.$$

- Limit: switching diffusion.

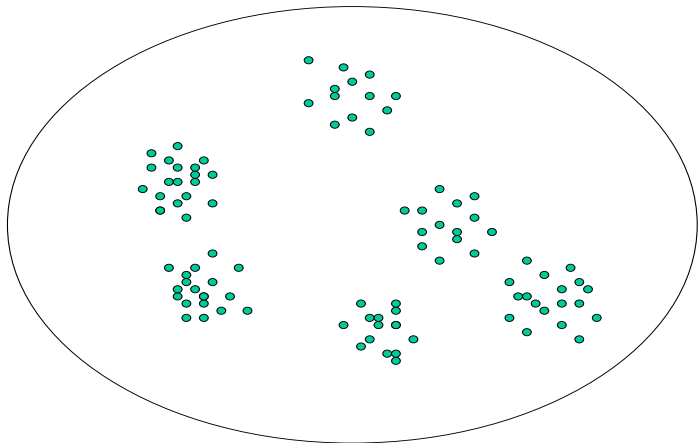
with Badowski, Zhang, Ann. Appl. Probab. (2000)

with Zhang, Ann. Appl. Probab. (2007)

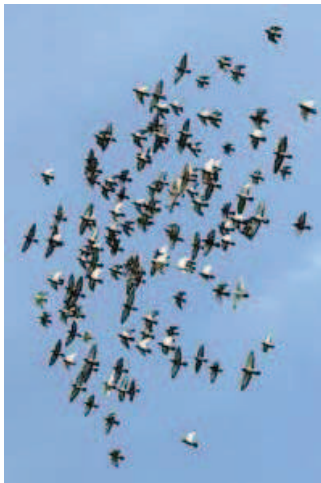
Two-time Scale (a demonstration)



Aggregation for Large-scale Systems



Consensus: Flocking



Consensus: Schooling (Couzin et.al. Nature, 2005)



Consensus Problems: High Way Traffic



Consensus: Honeybee Organization (Visscher, Nature, 2003)



Consensus Issues

- multi-agent coordination
- a group objective (e.g., alignment during motion, UAVs formation)
- to maintain shared information
- some kind of agreement such as objective of operation or a condition for proceeding to further operation

- Our work: Stochastic recursive algorithm, topology switching, multi-scale systems

with L.Y. Wang and Y. Sun, SIAM MMS, Automatica (2011)

Mean-Field Model

- $\alpha(t)$: with $\mathcal{M} = \{1, 2, \dots, m_0\}$.
- Consider an ℓ -body mean-field model For $i = 1, 2, \dots, \ell$,

$$\begin{aligned}dX_i(t) &= [\gamma(\alpha(t))X_i(t) - X_i^3(t) - \beta(\alpha(t))(X_i(t) - \bar{X}(t))] dt \\ &\quad + \sigma_{ii}(X(t), \alpha(t))dw_i(t), \\ \bar{X}(t) &= \frac{1}{\ell} \sum_{j=1}^{\ell} X_j(t), \\ X(t) &= (X_1(t), X_2(t), \dots, X_{\ell}(t))',\end{aligned}\tag{14}$$

$\gamma(i) > 0$ and $\beta(i) > 0$ for $i \in \mathcal{M}$.

- Originated from statistical mechanics, mean-field models are concerned with many-body systems with interactions. To overcome the difficulty of interactions due to the many bodies, one of the main ideas is to **replace all interactions to any one body with an average or effective interaction**.

with F. Xi, J. Appl. Probab. (2009)

Insurance Risk Models

The surplus at time t :

$$S(t, x, i) = x + \int_0^t c(\alpha(s)) ds - \sum_{j=1}^{N(t)} X_j(\alpha(T_j)),$$

- (x, i) : initial (surplus, regime);
- $c(i)$: premium rate;
- $X_j(i)$: claim size;
- T_j : claim time;
- $N(t)$: Poisson process.
- $\alpha(t)$ is used to model:
 - ▶ El Nino/La Nina phenomena in property ins.
 - ▶ economic condition in unemployment policy
 - ▶ certain epidemics in health insurance

Stock Price Models

- Stock market models
 - ▶ $S(t)$: stock price
 - ▶ $w(\cdot)$: stand Brownian motion
 - ▶ μ : return (appreciation) rate
 - ▶ σ : volatility
- traditional GBM model is given by

$$dS(t) = \mu S(t)dt + \sigma S(t)dw.$$

- Regime-switching market models

$$dS(t) = \mu(\alpha(t))S(t)dt + \sigma(\alpha(t))S(t)dw.$$

- ▶ both the return rate & volatility depend on $\alpha(t)$
- ▶ $\alpha(\cdot)$ and $w(\cdot)$ are independent
- ▶ $\alpha(t)$: market mode, investor's mode, & other economic factors (e.g., bull, bear)

with X.Y. Zhou, *SIAM J. Control Optim.* (2003), *IEEE T-AC*, (2004)

with Bensoussan and Yan (2012), *SIAM J. Fin.*

Properties

Regularity & Recurrence

Definition

Regularity. A Markov process $Y^{x,\alpha}(t) = (X^{x,\alpha}(t), \alpha^{x,\alpha}(t))$ is said to be *regular*, if for any $0 < T < \infty$,

$$\mathbf{P}\left\{\sup_{0 \leq t \leq T} |X^{x,\alpha}(t)| = \infty\right\} = 0. \quad (15)$$

Remark

Let $\beta_n := \inf\{t : |X^{x,\alpha}(t)| = n\}$. Then $\{\beta_n\}$ is monotonically increasing and hence has a (finite or infinite) limit. It follows that the process is regular iff

$$\beta_n \rightarrow \infty \text{ almost surely as } n \rightarrow \infty. \quad (16)$$

Definition

- (i) **Recurrence.** For $U := D \times J$, where $J \subset \mathcal{M}$ and $D \subset \mathbb{R}^r$ is an open set with compact closure, let $\sigma_U^{x,\alpha} = \inf\{t : Y^{x,\alpha}(t) \in U\}$. A regular process $Y^{x,\alpha}(\cdot)$ is *recurrent* w.r.t. U if

$$\mathbf{P}\{\sigma_U^{x,\alpha} < \infty\} = 1 \text{ for any } (x, \alpha) \in D^c \times \mathcal{M}.$$

- (ii) **Positive and Null Recurrence.** A recurrent process with finite mean recurrence time for some set $U = D \times J$ is said to be *positive recurrent* w.r.t. U ; otherwise, the process is *null recurrent* w.r.t. U .

Recurrence Is Independent of Sets

- (i) The process $(X(t), \alpha(t))$ is (positive) recurrent w.r.t. $D \times \mathcal{M}$ if and only if it is (positive) recurrent w.r.t. $D \times \{\ell\}$, where $D \subset \mathbb{R}^r$ is a bounded open set with compact closure and $\ell \in \mathcal{M}$.
- (ii) If the process $(X(t), \alpha(t))$ is (positive) recurrent w.r.t. some $U = D \times \mathcal{M}$, where $D \subset \mathbb{R}^r$, then it is (positive) recurrent w.r.t. $\tilde{U} = \tilde{D} \times \mathcal{M}$, where $\tilde{D} \subset \mathbb{R}^r$ is any nonempty open set.

Positive Recurrence (1)

Theorem

A necessary and sufficient condition for positive recurrence with respect to a domain $U = D \times \{\ell\} \subset \mathbb{R}^r \times \mathcal{M}$ is: For each $i \in \mathcal{M}$, there exists a nonnegative function $V(\cdot, i) : D^c \mapsto \mathbb{R}$ s.t. $V(\cdot, i)$ is twice continuously differentiable and that

$$\mathcal{L}V(x, i) = -1, \quad (x, i) \in D^c \times \mathcal{M}. \quad (17)$$

Let $u(x, i) = \mathbf{E}_{x, i} \sigma_D$. It is the smallest positive sol'n to

$$\begin{cases} \mathcal{L}u(x, i) = -1, & (x, i) \in D^c \times \mathcal{M}, \\ u(x, i) = 0, & (x, i) \in \partial D \times \mathcal{M}. \end{cases} \quad (18)$$

Step 1: Positive recurrence. Show the process is positive recurrent if exists $V(\cdot, \cdot) (\geq 0)$ satisfying the conditions of the theorem.

- Fix any $(x, i) \in D^c \times \mathcal{M}$ and set $\sigma_D^{(n)}(t) = \min\{\sigma_D, t, \beta_n\}$. Dynkin's formula implies

$$\begin{aligned} \mathbf{E}_{x,i} V(X(\sigma_D^{(n)}(t)), \alpha(\sigma_D^{(n)}(t))) - V(x, i) &= \mathbf{E}_{x,i} \int_0^{\sigma_D^{(n)}(t)} \mathcal{L} V(X(s), \alpha(s)) ds \\ &= -\mathbf{E}_{x,i} \sigma_D^{(n)}(t). \end{aligned}$$

Since $V(\cdot)$ is nonnegative,

$$\mathbf{E}_{x,i} \sigma_D^{(n)}(t) \leq V(x, i).$$

Letting $n \rightarrow \infty$ and $t \rightarrow \infty$, $\mathbf{E}_{x,i} \sigma_D < \infty$. This is positive recurrence.

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Step 2: Show $u(x, i) := \mathbf{E}_{x,i} \sigma_D < \infty$ is the smallest positive solution of the BVP (18).

- Set $\sigma_D^{(n)} = \min\{\sigma_D, \beta_n\}$ & $u_n(x, i) = \mathbf{E}_{x,i} \sigma_D^{(n)}$. Then $u_n(x, i)$ solves

$$\mathcal{L}u_n(x, i) = -1, \quad u_n(x, i)|_{x \in \partial D} = 0, \quad u_n(x, i)|_{|x|=n} = 0.$$

- $v_n(x, i) := u_{n+1}(x, i) - u_n(x, i)$ is \mathcal{L} -harmonic in $(D^c \cap \{|x| < n\}) \times \mathcal{M}$.
- $\mathbf{E}_{x,i} \sigma_D^{(n)} \nearrow \mathbf{E}_{x,i} \sigma_D$ by regularity and DCT. Hence we can write

$$u(x, i) = u_{n_0}(x, i) + \sum_{k=n_0}^{\infty} v_k(x, i).$$

- Harnack's theorem implies that $u(x, i)$ is a solution of (18).
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- Set $\sigma_D^{(n)} = \min\{\sigma_D, \beta_n\}$ & $u_n(x, i) = \mathbf{E}_{x, i} \sigma_D^{(n)}$. Then $u_n(x, i)$ solves

$$\mathcal{L} u_n(x, i) = -1, \quad u_n(x, i)|_{x \in \partial D} = 0, \quad u_n(x, i)|_{|x|=n} = 0.$$

- $v_n(x, i) := u_{n+1}(x, i) - u_n(x, i)$ is \mathcal{L} -harmonic in $(D^c \cap \{|x| < n\}) \times \mathcal{M}$.
- $\mathbf{E}_{x, i} \sigma_D^{(n)} \nearrow \mathbf{E}_{x, i} \sigma_D$ by regularity and DCT. Hence we can write

$$u(x, i) = u_{n_0}(x, i) + \sum_{k=n_0}^{\infty} v_k(x, i).$$

- Harnack's theorem implies that $u(x, i)$ is a solution of (18).
- Maximum Principle yields $u(x, i)$ is the smallest solution.

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Step 3: Show: If $Y(t) = (X(t), \alpha(t))$ is positive recurrent w.r.t. $U = D \times \{\ell\}$, then $\exists V$ satisfying $V \geq 0$ and the conditions of the theorem.

- The positive recurrence implies $\mathbf{E}_{x,i} \sigma_D < \infty$ for all $(x, i) \in D^c \times \mathcal{M}$.
Noting $\sigma_D^{(n)} \leq \sigma_D^{(n+1)}$, Harnack's theorem for \mathcal{L} -elliptic systems implies that the bounded monotone increasing sequence $u_n(x, i)$ converges uniformly on every compact subset of $D^c \times \mathcal{M}$.
Moreover, its limit $u(x, i)$ satisfies $\mathcal{L}u(x, i) = -1$ for every $i \in \mathcal{M}$.
- The function $V(x, i) := u(x, i)$ satisfies the required condition.

Ergodicity

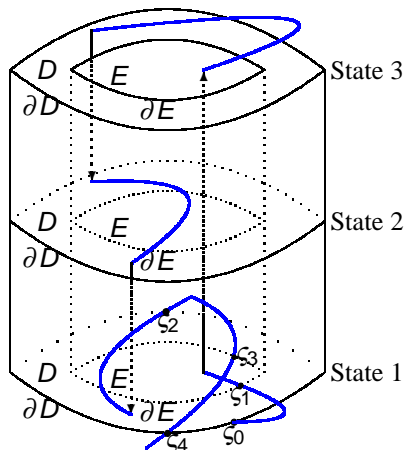


Figure 2: Cycles of $Y(t) = (X(t), \alpha(t))$; $m = 3$ & $\ell = 1$

Cycles

- Assume the process is positive recurrent w.r.t. $U = E \times \{\ell\}$; $E \subset \mathbb{R}^r$ and $\ell \in \mathcal{M}$ are fixed from now on.
- Let ∂E be sufficiently smooth. Let $D \subset \mathbb{R}^r$ be a bdd. ball with suff. smooth ∂D s.t. $E \cup \partial E \subset D$.
- Let $\zeta_0 = 0$ and then define for $n = 0, 1, \dots$

$$\zeta_{2n+1} = \inf\{t \geq \zeta_{2n} : (X(t), \alpha(t)) \in \partial E \times \{\ell\}\},$$

$$\zeta_{2n+2} = \inf\{t \geq \zeta_{2n+1} : (X(t), \alpha(t)) \in \partial D \times \{\ell\}\}.$$

Then we can divide an arbitrary sample path of the process into cycles:

$$[\zeta_0, \zeta_2), [\zeta_2, \zeta_4), \dots, [\zeta_{2n}, \zeta_{2n+2}) \dots \quad (19)$$

- Assume $Y(0) = (X(0), \alpha(0)) = (\mathbf{x}, \ell) \in \partial D \times \{\ell\}$.
- Define $Y_n = Y(\zeta_{2n}) = (X_n, \ell), n = 0, 1, \dots$. It is a MC on $\partial D \times \{\ell\}$ by strong Markov property

Theorem

A positive recurrent process $(X(t), \alpha(t))$ has a unique stationary distribution $\hat{v}(\cdot, \cdot) = (\hat{v}(\cdot, i) : i \in \mathcal{M})$.

Strong Law of Large Numbers

Theorem

Denote by $\mu(\cdot, \cdot)$ the stationary density associated with $\widehat{v}(\cdot, \cdot)$ and $f(\cdot, \cdot) : \mathbb{R}^r \times \mathcal{M} \mapsto \mathbb{R}$ is Borel measurable such that

$$\sum_{i=1}^{m_0} \int_{\mathbb{R}^r} |f(x, i)| \mu(x, i) dx < \infty. \quad (20)$$

Then for any $(x, i) \in \mathbb{R}^r \times \mathcal{M}$

$$\mathbf{P}_{x,i} \left(\frac{1}{T} \int_0^T f(X(t), \alpha(t)) dt \rightarrow \bar{f} \right) = 1, \quad (21)$$

where $\bar{f} = \sum_{i=1}^{m_0} \int_{\mathbb{R}^r} f(x, i) \mu(x, i) dx$.

Cauchy Problem

Let the assumptions of the last theorem be satisfied, and $u(t, \mathbf{x}, i)$ be the solution of the Cauchy problem

$$\begin{cases} \frac{\partial u(t, \mathbf{x}, i)}{\partial t} = \mathcal{L}u(\mathbf{x}, i), & i \in \mathcal{M}, \\ u(0, \mathbf{x}, i) = f(\mathbf{x}, i). \end{cases} \quad (22)$$

Then as $T \rightarrow \infty$,

$$\frac{1}{T} \int_0^T u(t, \mathbf{x}, i) dt \rightarrow \sum_{i=1}^{m_0} \int_{\mathbb{R}^r} f(\mathbf{x}, i) \mu(\mathbf{x}, i) d\mathbf{x}. \quad (23)$$

A key to establish this result is the result of law of large numbers.

Stability

Definitions

The equilibrium point $x = 0$ of system is:

- (i) *stable in probability*, if for any $r > 0$, $\lim_{x \rightarrow 0} \mathbf{P}\{\sup_{t \geq 0} |X^{x,\alpha}(t)| > r\} = 0$;
otherwise $x = 0$ is *unstable in probability*.
- (ii) *asymptotically stable in probability*, if it is stable in prob. &
 $\lim_{x \rightarrow 0} \mathbf{P}\{\lim_{t \rightarrow \infty} X^{x,\alpha}(t) = 0\} = 1$.
- (iii) *p-stable (for $p > 0$)*, if $\lim_{\delta \rightarrow 0} \sup_{|x| \leq \delta, \alpha \in \mathcal{M}, t \geq 0} \mathbf{E}|X^{x,\alpha}(t)|^p = 0$.
- (iv) *asymptotically p-stable*, if it is *p-stable* & $\mathbf{E}|X^{x,\alpha}(t)|^p \rightarrow 0$ as $t \rightarrow \infty$.
- (v) *exponentially p-stable*, if for some $K, k > 0$,
 $\mathbf{E}|X^{x,\alpha}(t)|^p \leq K|x|^p \exp\{-kt\}$, for any $\alpha \in \mathcal{M}$.

Necessary Conditions

Theorem

$x = 0$ is exponentially p -stable iff b and σ have continuous bdd. derivatives w.r.t. x up to the 2nd order. Then for $i \in \mathcal{M}$, $\exists V(\cdot, i) : \mathbb{R}^r \mapsto \mathbb{R}$ s.t.

$$\begin{aligned} k_1|x|^p &\leq V(x, i) \leq k_2|x|^p, \quad x \in N, \\ \mathcal{L}V(x, i) &\leq -k_3|x|^p \quad \text{for all } x \in N - \{0\}, \\ \left| \frac{\partial V}{\partial x_j}(x, i) \right| &< k_4|x|^{p-1}, \quad \left| \frac{\partial^2 V}{\partial x_j \partial x_k}(x, i) \right| < k_4|x|^{p-2}, \end{aligned} \tag{24}$$

for all $1 \leq j, k \leq n$, $x \in N - \{0\}$, and for some $k_i > 0$ ($i = 1, 2, 3, 4$), where N is a neighborhood of 0.

with Khasminskii & Zhu, Stochastic Proc. Appl. (2007)
with Mao & Yuan, Automatica (2007); (Markov switching)
with Xi SIAM J. Control Optim. (2010)

Linearized Systems

- For each $i \in \mathcal{M}$, $\exists b(i)$, $\sigma_j(i) \in \mathbb{R}^{r \times r}$, $j = 1, \dots, d$, and $\widehat{Q} = (\widehat{q}_{ij})$ s.t.

$$\left. \begin{aligned} b(x, i) &= b(i)x + o(|x|), \\ \sigma(x, i) &= (\sigma_1(i)x, \dots, \sigma_d(i)x) + o(|x|), \\ Q(x) &= \widehat{Q} + o(1), \end{aligned} \right\} \text{as } x \rightarrow 0 \quad (25)$$

\widehat{Q} is an irreducible generator of a Markov chain $\widehat{\alpha}(t)$. Use $\pi = (\pi_1, \dots, \pi_m) \in \mathbb{R}^{1 \times m}$ to denote the stationary dist. associated with \widehat{Q} .

Easily Verifiable Conditions

Theorem

Suppose that $\frac{\sigma_j'(i) + \sigma_j(i)}{2} \geq 0$.

(a) Then $x = 0$ (i) is asymptotically stable in prob. if

$$\sum_{i=1}^m \pi_i \left(\Lambda_{b(i)} + \frac{1}{2} \sum_{j=1}^d [\Lambda_{a_j(i)} - 2(\lambda_{\sigma_j(i)})^2] \right) < 0 \quad (26)$$

and (ii) is unstable if

$$\sum_{i=1}^m \pi_i \left(\lambda_{b(i)} + \frac{1}{2} \sum_{j=1}^d [\lambda_{a_j(i)} - 2(\Lambda_{\sigma_j(i)})^2] \right) > 0, \quad (27)$$

where Λ_A and λ_A denote the max. and min. eigenvalue of $\frac{1}{2}(A + A')$, resp.

(b) If $X(t)$ is 1-d, then $x = 0$ is (i) asymptotically stable in probab if

$\sum_{i=1}^m \pi_i \left(b_i - \frac{\sigma_i^2}{2} \right) < 0$, & (ii) unstable in probab if $\sum_{i=1}^m \pi_i \left(b_i - \frac{\sigma_i^2}{2} \right) > 0$.

Idea of Proof

We only consider (i).

$\mu = (\mu_1, \dots, \mu_m)' \in \mathbb{R}^m$ with $\mu_i = \Lambda_{b(i)} + \frac{1}{2} \sum_{j=1}^d \Lambda_{a_j(i)}$. Let $\beta := -\pi\mu > 0$. Then that $\widehat{Q}c = \mu + \beta \mathbf{1}$ has a soln. $c = (c_1, \dots, c_m)' \in \mathbb{R}^m$.

Consider $V(x, i) = (1 - \gamma c_i) |x|^\gamma$, where $0 < \gamma < 1$ is suff. small s.t. $1 - \gamma c_i > 0$, $i \in \mathcal{M}$. $V(\cdot, i)$ is continuous, nonnegative, & vanishes only at $x = 0$. $\mathcal{L}V(x, i) < 0$ for any $(x, i) \in (N - \{0\}) \times \mathcal{M}$, where $N \subset \mathbb{R}^r$ is a small neighborhood of 0. Then we can show that 0 is asymptotically stable.

Closing the "Gap"

Consider 'linear system' with $Q(x) \equiv Q$, and define $Y(t) = X(t)/|X(t)|$. Itô's formula implies that

$$dY(t) = \Phi(Y(t), \alpha(t))dt + \Psi(Y(t), \alpha(t))dw(t),$$

where $w(t) = (w_1(t), \dots, w_d(t))' \in \mathbb{R}^d$ with $w_k(t), k = 1, \dots, d$ being indep. 1-dim. Brownian motions and Φ, Ψ are appropriate functions. We can represent $\ln|X(t)|$ in terms of $(Y(t), \alpha(t))$.

Denote the stationary density of $(Y(t), \alpha(t))$ by $\mu(y, i), i \in \mathcal{M}$

$$\rho_0 = \sum_{i=1}^m \int_{\mathbb{S}} \left[y' b(i) y + \frac{1}{2} \sum_{k=1}^d (|\sigma_k(i) y|^2 - 2|y' \sigma_k(i) y|^2) \right] \mu(y, i) dy.$$

Then asymptotically stable if $\rho_0 < 0$ and unstable if $\rho_0 > 0$.

Explosion Suppression & Stabilization

Regularity Criterion (cont.)

Theorem

Suppose that $b(\cdot, \cdot) : \mathbb{R}^r \times \mathcal{M} \mapsto \mathbb{R}^r$ and that $\sigma(\cdot, \cdot) : \mathbb{R}^r \times \mathcal{M} \mapsto \mathbb{R}^{r \times d}$,

$$\begin{aligned} dX(t) &= b(X(t), \alpha(t))dt + \sigma(X(t), \alpha(t))dw(t), \quad (X(0), \alpha(0)) = (x, \alpha), \\ P\{\alpha(t + \delta) = j | \alpha(t) = i, X(s), \alpha(s), s \leq t\} &= q_{ij}(X(t))\delta + o(\delta), \quad i \neq j. \end{aligned} \quad (28)$$

Suppose that for each $i \in \mathcal{M}$, both $b(\cdot, i)$ and $\sigma(\cdot, i)$ are local linear growth and local Lipschitzian and that \exists a nonnegative $V(\cdot, \cdot) : \mathbb{R}^r \times \mathcal{M} \mapsto \mathbb{R}^+$ that is C^2 in $x \in \mathbb{R}^r$ for each $i \in \mathcal{M}$ s.t. $\exists \gamma_0 > 0$

$$\begin{aligned} \mathcal{L}V(x, i) &\leq \gamma_0 V(x, i), \quad \text{for all } (x, i) \in \mathbb{R}^r \times \mathcal{M}, \\ V_R &:= \inf_{|x| \geq R, i \in \mathcal{M}} V(x, i) \rightarrow \infty \quad \text{as } R \rightarrow \infty. \end{aligned} \quad (29)$$

Then the process $(X(t), \alpha(t))$ is regular.

Explosion Suppression

$$x \in \mathbb{R}^r$$

$$f(\cdot, \cdot) : \mathbb{R}^r \times \mathcal{M} \mapsto \mathbb{R}^r$$

$$\alpha(t) \in \mathcal{M} = \{1, \dots, m\}$$

$$\frac{dX(t)}{dt} = f(X(t), \alpha(t)) \quad (30)$$

$f(\cdot, i)$ continuous but the growth rate is faster than linear

We wish to stabilize (30).

Motivational Example

- Consider an even simpler problem: the logistic system

$$\dot{x}(t) = x(t)(1 + x(t)), \quad x(0) = 1.$$

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- Question: How can we get a global soln; how can we stabilize this?

Motivational Example

- Consider an even simpler problem: the logistic system

$$\dot{x}(t) = x(t)(1 - x(t)), \quad x(0) = 1.$$

- solution:

$$x(t) = \frac{1}{-1 + 2e^{-t}}.$$

- It will blow up and the explosion time $\tau = \log 2$.
- Question: How can we get a global soln; how can we stabilize this?

Two things are needed:

- 1) extend to a global solution;
- 2) stabilization.

What have been done?

- Khasminskii's book (1981): stabilize 2-d system with two white noise
- Arnold (1972): $\dot{x} = Ax$ can be stabilized by zero mean stationary process iff $\text{tr}(A) < 0$
- Mao (1994) established a general stabilization results of Brownian noise under linear growth condition.
- Wu & Hu (2009) treated one-sided growth condition
- Mao, Yin, and Yuan (2007): showed that both Brownian motion and Markov Chain can be used to stabilize systems.

Motivation (diffusion case)

$$dx = \mu x dt + \sigma x dw, \quad x(0) = x_0.$$

$$x(t) = x_0 \exp \left(\left(\mu - \frac{\sigma^2}{2} \right) t + \sigma w(t) \right).$$

$$\text{when } \sigma^2 > 2\mu,$$

$$\limsup_t \frac{\log |x(t)|}{t} \leq \left(\mu - \frac{\sigma^2}{2} \right) < 0.$$

This implies exponential stability.

How to Get a Global Solution? Stabilization?

- add a diffusion perturbation

$$dX(t) = f(X(t), \alpha(t))dt + a_1(\alpha(t))|X(t)|^\beta X(t)dw_1(t)$$

such that $2\beta - \beta_1 > 0$, where $w_1(\cdot)$ is scalar Brownian motion.

- add another diffusion to get stability

$$dX(t) = f(X(t), \alpha(t))dt + a_1(\alpha(t))|X(t)|^\beta X(t)dw_1(t) + a_2(\alpha(t))X(t)dw_2(t), \quad (31)$$

where $w_2(\cdot)$ is a scalar Brownian motion independent of $w_1(\cdot)$.

- More general,

$$dX(t) = f(X(t), \alpha(t))dt + \sigma_1(X(t), \alpha(t))dw_1 + \sigma_2(X(t), \alpha(t))dw_2. \quad (32)$$

Results

- With proper choice of the perturbations, we get a global solution
- $\limsup_{t \rightarrow \infty} P(|X(t)| \geq K_\delta) \leq \delta$
- The resulting system is stable w.p.1. In fact, $\limsup_t \log |X(t)|/t < 0$ w.p.1.

with Wu and Zhao, SIAM J. Appl. Math (2012)

Example

Begin with (30) together with initial condition $X(0) = 1$. Suppose that $\alpha(t)$ is a Markov chain with two states $\mathcal{M} = \{1, 2\}$ and

$$Q = \begin{pmatrix} -0.1 & 0.1 \\ 1 & -1 \end{pmatrix}, f(x, 1) = x(x + 1) \text{ and } f(x, 2) = x(2x + 1).$$

Corresponding to the states, we have two equations

$$\begin{aligned} \frac{d}{dt}X(t) &= X(t)(X(t) + 1), \\ \frac{d}{dt}X(t) &= X(t)(2X(t) + 1). \end{aligned} \tag{33}$$

Neither equation has a global soln. For the 1st equation, we have $X(t) = e^t/(2 - e^t)$ that will blow up at time $\ln 2$; for the second equation, $X(t) = e^t/(3 - 2e^t)$ that will blow up at time $\ln(3/2)$. We plot the trajectories of the switched system as well as each individual system.

To regularize the system, use a feedback control $a_1(\alpha(t))X^2(t)dw_1(t)$, where $w_1(t)$ is a 1-d Brownian motion. The resulting eq is

$$dX(t) = f(X(t), \alpha(t))dt + a_1(\alpha(t))X^2(t)dw_1(t), \quad (34)$$

$a_1(i) = 2$ for $i = 1, 2$.

Although the system has a global solution, it is not asymptotically stable. To stabilize the system, we add another feedback control $a_2(\alpha(t))X(t)dw_2(t)$, $w_2(t)$ is 1-d standard Brownian motion independent of $w_1(t)$ and $a_2(1) = 19$ and $a_2(2) = 24$.

$$dX(t) = f(X(t), \alpha(t))dt + a_1(\alpha(t))X^2(t)dw_1(t) + a_2(\alpha(t))X(t)dw_2(t). \quad (35)$$

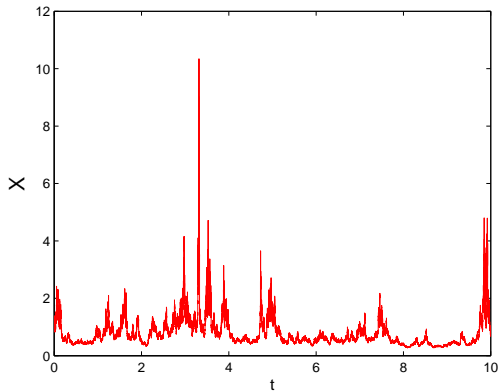


Figure: Trajectory of system (34) with stepsize $\Delta t = 10^{-4}$.

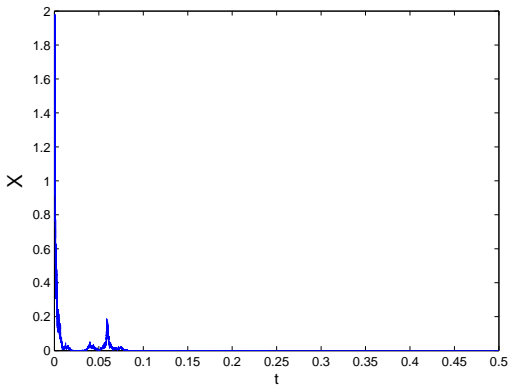
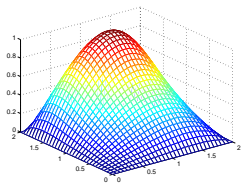


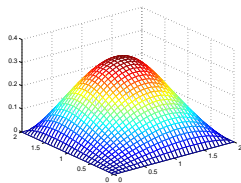
Figure: Trajectory of system (35) with stepsize $\Delta t = 10^{-6}$.

Numerical Approximations, Controlled Switching Diffusions, Games

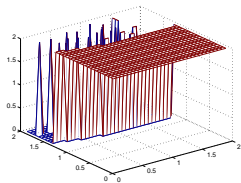
Numerical Methods for SDE, Controls, and Games



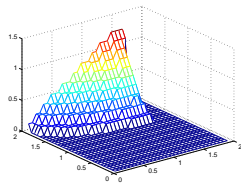
(a) $V^{h,+}(\cdot, \cdot, 1)$



(b) $V^{h,+}(\cdot, \cdot, 2)$



(c) $U_1(\cdot, \cdot, 1)$: player1 1st



(d) $U_1(\cdot, \cdot, 2)$ player1 1st

Numerics for Controlled Switching Diffusions

$$\begin{cases} X(t) = x + \int_0^t b(X(s), \alpha(s), u(s)) ds + \int_0^t \sigma(X(s), \alpha(s)) dw, \\ \alpha(t) \text{ continuous-time MC } \alpha(0) = i, \end{cases} \quad (36)$$

where $w(t)$ is a standard Brownian motion independent of the Markov chain $\alpha(t)$.

- Kushner & Dupuis, Springer, Markov chain approximation
- with Song & Zhang (2006), regime-switching & jump diffusion

Controlled Switching Diffusions (cont.)

Given $B > 0$, define a stopping time as

$$\tau_B^{x,i,u} = \inf\{t : X^{x,i,u}(t) \notin (-B, B)\}.$$

Objective: choose control u . to minimize the expected cost function

$$\begin{cases} J_i^B(x, u) = \mathbf{E} \int_0^{\tau_B^{x,i,u}} f(X(s), \alpha(s), u(s)) ds, \\ \quad \forall x \in (-B, B), i \in \mathcal{M}, \\ J_i^B(x, u) = 0, \forall x \notin (-B, B), i \in \mathcal{M}, \end{cases} \quad (37)$$

where for each $i \in \mathcal{M}$, $f(\cdot, i, \cdot)$ is an appropriate function representing the running cost function.

For each $i \in \mathcal{M}$, the value function is given by

$$V^B(x, i) = \inf_{u \in \mathcal{U}} J^B(x, i, u), \quad (38)$$

where \mathcal{U} is the space of all \mathcal{F}_t -adapted controls taking values on a compact set U .

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Formally, the value functions satisfy Hamilton-Jacobi-Bellman (HJB) equations,

$$\begin{cases} \inf_{u \in U} \{L^u V^B(x, i) + f(x, i, u)\} = 0, & \forall x \in (-B, B), i \in \mathcal{M}, \\ V^B(x, i) = 0, & \forall x \notin (-B, B), i \in \mathcal{M}, \end{cases} \quad (39)$$

where

$$L^u \varphi(x, i) = \frac{1}{2} \sigma^2(x, i) \frac{d^2 \varphi(x, i)}{dx^2} + b(x, i, u) \frac{d\varphi(x, i)}{dx} + \sum_{j \in \mathcal{M}} q_{ij} \varphi(x, j).$$

Algorithm

- $h > 0$: discretization parameter.
- $S_h = \{x : x = kh, k = 0, \pm 1, \pm 2, \dots\}$. Let $\{(\xi_n^h, \alpha_n^h), n < \infty\}$ be a controlled discrete-time Markov chain on a discrete state space $S_h \times \mathcal{M}$
- $p^h((x, i), (y, j)|u)$: transition probabilities from $(x, i) \in S_h \times \mathcal{M}$ to $(y, j) \in S_h \times \mathcal{M}$, for $u \in U$.

Then, $\bar{V}^{B,h}(x, i)$, the discretization of $V^B(x, i)$ with step size $h > 0$, is the solution of

$$\begin{cases} \inf_{u \in U} \{L_h^u \bar{V}^{B,h}(x, i) + f(x, i, u)\} = 0, & \forall x \in (-B, B)_h, i \in \mathcal{M}, \\ \bar{V}^{B,h}(x, i) = 0, & \forall x \notin (-B, B)_h, i \in \mathcal{M}, \end{cases} \quad (40)$$

where

$$(-B, B)_h = (-B, B) \cap S_h, \quad [-B, B]_h = (-B, B)_h \cup \{B, -B\}. \quad (41)$$

$$\begin{aligned} \bar{V}^{B,h}(x, i) = \inf_{u \in U} \left\{ \bar{p}_i^{h,+}(x, u) \bar{V}^{B,h}(x+h, i) + \bar{p}_i^{h,-}(x, u) \bar{V}^{B,h}(x-h, i) \right. \\ \left. + \sum_{j \neq i} \bar{p}_{ij}^h(x) \bar{V}^{B,h}(x, j) + f(x, i, u) \Delta \bar{t}_i^h(x) \right\} \end{aligned} \quad (42)$$

Rates of Convergence

Theorem

Under suitable conditions, $\exists \gamma \in (2, 3]$ and $\rho \in (0, 1]$ s.t. the Markov chain approximation algorithm converges at the rate $(\gamma - 2) \wedge \rho \wedge \frac{1}{2}$. That is,

$$|\bar{V}_i^{B,h}(x) - V_i^B(x)| \leq Kh^{\frac{1}{2} \wedge \rho \wedge (\gamma - 2)}, \quad \forall (i, x) \in \mathcal{M} \times \mathbf{G}.$$

- Note that $\gamma \in (2, 3]$ comes from Markov chain \approx , ρ is the Hölder exponent of the cost function.
- PDE approach for controlled diffusions (finite difference approx of PDEs)
 - ▶ Menaldi, SIAM J. Control Optim. (1989)
 - ▶ Krylov, Probab. Theory Related Fields, (2000)
 - ▶ Dong & N.V. Krylov, Appl. Math Optim.
- we use probabilistic approach for controlled switching diffusions
 - ▶ with Q.S. Song, SIAM J. Control Optim. (2009)

Main Ideas

- Use relaxed controls (measures)
- Construct strong approximation
- Consider boundary perturbations
 - ▶ usual notion of cost $J_i(x, \tilde{m})$;
 - ▶ ours $J_i^B(x, \tilde{m})$

Main Ideas

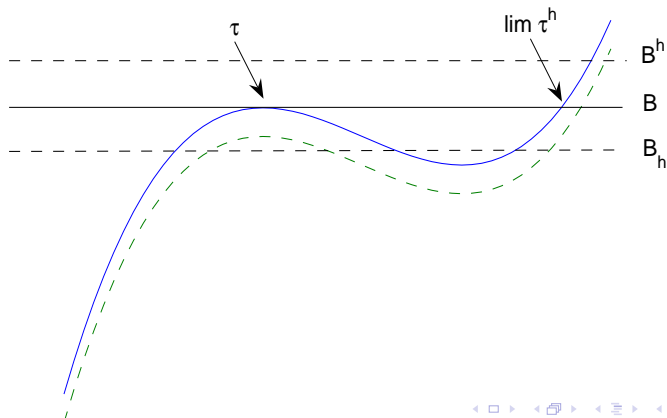
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Tangency Problem

- τ and τ^h : the first hitting time of $X(t)$ and $x^h(t)$ to the boundary.
- Objective: $\approx \mathbf{E}\tau$ by $\mathbf{E}\tau^h$
- In the Figure, $\tau^h \not\rightarrow \tau$, even though $x^h(\cdot)$ converges to $X(\cdot)$.
- Q: extra conditions needed?



Concluding Remarks

In this talk, we

- presented several switching diffusion examples
- considered such properties as recurrence, ergodicity, stability etc.
- developed numerical algorithms for control and game problems
- ascertained rates of convergence and treated tangency problems

Further work:

- rates of convergence for games
- large deviations
- null-recurrent switching diffusion systems ...

- Past dependent switching and countable switching space:

$$\begin{aligned} \mathbf{P}\{\alpha(t+\Delta) = j | \alpha(t) = i, X_s, \alpha(s), s \leq t\} &= q_{ij}(X_t)\Delta + o(\Delta) \text{ if } i \neq j \\ \mathbf{P}\{\alpha(t+\Delta) = i | \alpha(t) = i, X_s, \alpha(s), s \leq t\} &= 1 - q_i(X_t)\Delta + o(\Delta), \\ q_i(\phi) &= \sum_{j \neq i} q_{ij}(\phi) \text{ for any } (\phi, i) \in \mathcal{C} \times \mathbb{Z}_+. \end{aligned}$$

With D. Nguyen, SIAM J. Control Optim. (2016), Potential Anal. (2017)

- Switching jump diffusion:

$$\begin{aligned} \mathcal{L}f(x) &= \sum_{k,l=1}^r a_{kl}(x) \frac{\partial^2 f(x)}{\partial x_k \partial x_l} + \sum_{k=1}^r b_k(x) \frac{\partial f(x)}{\partial x_k} \\ &+ \int_{\mathbb{R}^r} (f(x+z) - f(x) - \nabla f(x) \cdot \mathbf{z} \mathbf{1}_{\{|z|<1\}}) \pi(x, dz). \end{aligned}$$

with Chen, Chen, Tran, Bernoulli (2018), Appl. Math Optim. (2018)

Thank you