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STABILITY AND EXISTENCE OF DIFFUSIONS WITH DISCONTINUOUS OR RAPIDLY GROWING DRIFT TERMS+

by

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Abstract

Stochastic differential equations whose drift terms do not satisfy the usual (Itô) Lipschitz or linear growth conditions in the state occur frequently as models in stochastic control theory. Local stability properties are useful for proving global existence for ordinary differential equations whose right hand sides grow too fast or are not Lipschitz in the state. Here, we use a local stochastic stability property to prove global existence, stability, ergodicity, the strong Markov and other properties, for a class of diffusions which occur frequently as models.

1. Introduction

For a vector $x = \{x_i\}$ and matrix $\sigma = \{\sigma_{ij}\}$, define the Euclidean norms $|x|^2 = \sum_i x_i^2$, $|\sigma|^2 = \sum_{i,j} \sigma_{ij}^2$, resp. Consider the homogeneous[+] Itô stochastic differential equation

$$dx = f(x)dt + \sigma(x)dz, \qquad t \geq 0 \tag{1}$$

where $\sigma(\cdot)$ satisfies growth and Lipschitz conditions of the types[++]

$$|\sigma(x)|^2 \leq K(1+|x|^2) \tag{2a}$$

$$|\sigma(x) - \sigma(y)| \leq K(1+|x|), \tag{2b}$$

and $z(t)$ is a normalized vector valued Wiener process. If

$$|f(x)|^2 \leq K(1+|x|^2) \tag{3a}$$

$$|f(x) - f(y)| \leq K|x-y| \tag{3b}$$

then the Ito existence theory is applicable to (1) and the stability properties can be discussed [1]. If (3b) holds locally, but (3a) is violated, a 'local' stability property([1], Theorem 8, Chapter 2) ensures the existence of a solution to (1) for all $t \geq 0$.

[+]The homogeneity condition is not essential, except in Section 4.

[++]K and K_i always denote real numbers; their value may change from usuage to usuage.

Recent investigations [2-5] have studies an important class of equations (1), where $f(\cdot)$ is allowed some discontinuities. Rewrite (1) in the form (x^1 and x^2 are vectors).

$$(4) \qquad dx = \begin{pmatrix} dx^1 \\ dx^2 \end{pmatrix} = \begin{matrix} f^1(x)dt \\ f^2(x)dt + \hat{f}(x)dt + \hat{\sigma}(x)dz \end{matrix}$$

where we assume that the f^i and $\hat{\sigma}$ satisfy (3) and (2), respectively, and $\hat{\sigma}(x)$ has a uniformly bounded inverse. (Thus $\hat{\sigma}^{-1}(x)$ satisfies (2).), but $\hat{f}(\cdot)$ does not necessarily satisfy (3). In the sequel, we prove existence, uniqueness, and other properties of (4), when neither (3a) nor (3b) necessarily holds, but a 'local' stability property obtains, and also treat the problems of asymptotic stability, the existence of a unique invariant measure and the convergence of the measures of (1) to the invariant measure.

Diffusions of the type (4) occur frequently in control applications. Consider, for example, a 'white noise' driven n'th order differential equation where $\hat{f}$ is a 'bang-bang' control taking the values $\{+1,-1\}$, or which may be discontinuous on a smooth 'switching curve', and tend to infinity in certain directions. Also models such as

$$dx = \begin{matrix} dx_1 = x_2 dt \\ dx_2 = -(x_1+x_1^3)dt + \sigma dz \end{matrix}$$

are sometimes used, and the existence, and asymptotic character of the corresponding measures are of interest.

2. Mathematical Preliminaries

Assume

(C1) f^i and $\hat{\sigma}$ satisfy (3) and (2), respectively, and $\hat{\sigma}^{-1}(x)$ is uniformly bounded. $\hat{f}(\cdot)$ is a vector valued Borel function of x which is bounded in any compact set.

(C2) The process (5) has a transition density $p(x; t,y)$.

(C3) (A condition on the discontinuities of $\hat{f}$.) Let S_m denote a shpere of radius m, whose center is the origin. Let $N_\epsilon(A)$ denote an ϵ-neighborhood of the set A and $\mu(A)$ the Lebesgue measure of A. Suppose there is a (discontinuity) set D so that

$$\mu(N_\epsilon(D \cap S_m)) \to 0$$

as $\epsilon \to 0$ for each $m < \infty$. For each $\epsilon' > 0$, let there be an $\epsilon > 0$ so that $|x-y| < \epsilon$ implies $|\hat{f}(x,t) - \hat{f}(y,t)| < \epsilon'$ uniformly in x in bounded regions, provided that $x \notin N_\epsilon(D)$.

Assume (C1). Let Ω denote the sample space. We use the notation $(\Omega, z(t), \mathscr{B}_t, P)$ for the Wiener process on $[0,\infty)$, where $\mathscr{B}_t$ measures $z(s)$, $s \leq t$ and $z(r_2) - z(r_1)$ is independent of $\mathscr{B}_t$ for $t \leq r_1 \leq r_2$, and P is the measure on all the $\mathscr{B}_t$. We say that $z(t)$ is a Wiener process on $(\Omega, \mathscr{B}_t, P)$. Let $x(t)$ be the unique solution to the Itô equation (5)

$$(5) \qquad dx = \begin{array}{l} dx^1 = f^1(x)dt \\ dx^2 = f^2(x)dt + \hat{\sigma}(x)dz \end{array}$$

We say that $x(t)$ is an Itô process with respect to $(\Omega, z(t), \mathscr{B}_t, P_x)$, where P_x denotes the probability given that $x(0) = x$ (and E_x denotes the corresponding expectation). E and P denote expectation and probability for functionals of $z(t)$. Define Ω_T as the sample space for $z(t)$, $t \leq T$. Suppose that

$$(6) \qquad \int_0^T |\hat{\sigma}^{-1}(x(t))\hat{f}(x(t))|^2 dt < \infty \qquad \text{w.p.1.}$$

(which is certainly true if $\hat{f}$ is bounded). Define

$$\zeta_0^T(\hat{f}) \equiv \int_0^T \hat{\sigma}^{-1}(x(t))\hat{f}(x(t))dz(t) - \frac{1}{2}\int_0^T |\hat{\sigma}^{-1}(x(t))\hat{f}(x(t))|^2 dt$$

and suppose that

$$(7) \qquad E_x \exp \zeta_0^T(\hat{f}) = 1.$$

((7) holds for all $T < \infty$ if $\hat{f}$ is bounded.) Then the probability

measure $\tilde{P}^T$ defined by[+]

$$\tilde{P}_x^T(A) = \int_A \exp \zeta_o^T(f) \cdot P(d\omega)$$

is a measure on the $\mathscr{B}_t$, $t \leq T$. The process $\tilde{z}(t)$, $t \leq T$

$$\tilde{z}(t) = z(t) - \int_o^t \hat{\sigma}^{-1}(x(s))\hat{f}(x(s))ds$$

is a Wiener process on $(\Omega_T, \mathscr{B}_t, \tilde{P}_x^T)$, and the process

$$(8) \qquad dx = \begin{matrix} f^1(x)dt \\ f^2(x)dt + \hat{f}(x)dt + \hat{\sigma}(x)[dz - \hat{\sigma}^{-1}(x)\hat{f}(x)dt] \end{matrix}$$

$$= \begin{matrix} f^1(x)dt \\ f^2(x)dt + \hat{f}(x)dt + \hat{\sigma}(x)d\tilde{z} \end{matrix}$$

is an Itô process with respect to $(\Omega_T, \tilde{z}(t), \mathscr{B}_t, \tilde{P}_x^T)$. The construction was first done by Girsanov [4], and exploited by Benes [5], Rishel [2] and then Kushner [3], for several control problems. Note the sample space Ω_T, the σ-algebras $\mathscr{B}_t$ and the random variables $x(t)$ for the Wiener process $\tilde{z}(t)$, and Itô process $(\Omega_T, \tilde{z}(t), \mathscr{B}_t, \tilde{P}_x^T)$ are the same as those for the Wiener process $z(t)$ and Itô process (5), for $t \leq T$. Only the measures have been changed. The process (8) is constructed by a transformation of measures on the 'nicer' process (5).

[+]The measure $\tilde{P}_x^T$ depends on the initial condition of (5), as does the Wiener process $\tilde{z}(t)$.

The following facts (drawn from [2-4]) about (8) will be needed. Assume that $\hat{f}$ is bounded and that (C1-3) hold.

(01) ([3], Theorem 5). The multivariate distributions of (8) are continuous with respect to the initial condition $x(0)$, (in the sense that the characteristic functions are continuous in $x(0)$).

(02) ([3], Theorem 2). The solution to (8) is unique in the sense that any two solutions to (8) have the same multivariate distributions.

$$(03) \quad \tilde{E}_x^T \sup_{t \geq s \geq 0} |x(s) - x|^2 \leq K_1 t(1+|x|^2), \quad t \leq T$$

$$\tilde{E}_x^T \sup_{t \geq s \geq 0} |x(s) - x|^4 \leq K_1 t^2(1+|x|^4), \quad t \leq T$$

where $\tilde{E}_x^T$ is the expectation given $x(0) = x$, and K_1 depends on the bound on $\hat{f}$. The proof of (03) is close to that of (27) - (28) of [3] Theorem 6. K_1 depends on the bound on $\hat{f}$.

(04) If the process (5) has a density $p(x;\ t,y)$, then so does (8) and the density of (8) is any version of ([2], Lemma 1), (boundedness of $\hat{f}$ is not required if (6) - (7) hold) for $t \leq T$

$$q(x;\ t,y) = \tilde{E}_x^T[\exp \zeta_0^t(\hat{f}) | x(t) = y]p(x;\ t,y).$$

Also ($\hat{f}$ is not required to be bounded in (05)).

(05) ([4], Corollary to Lemma 3). Let $g(\omega)$ be $\mathscr{B}_t$ measurable with $\tilde{E}_x^T|g(\omega)| < \infty$, and $t \leq T$. Then, for $s \leq t \leq T$, w.p.1.

$$\tilde{E}_x^T[g(\omega)|\mathscr{B}_s] = E[g(\omega)\exp\zeta_s^t(\hat{f})|\mathscr{B}_s].$$

(The equation also holds if $\mathscr{B}_s$ is replaced by any sub σ-algebra of $\mathscr{B}_s$.)

Fix T, and define $\tilde{z}(t)$ and $\tilde{P}_x^T$ by the Girsanov transformation. Write $\tilde{z}(t)$ as $\tilde{z}^T(t)$. Suppose that (6) - (7) hold for a time $T_1 > T$, and define the corresponding Ω_{T_1}, $\tilde{z}^{T_1}(t)$, $\tilde{P}_x^{T_1}$. Then $\tilde{z}^{T_1}(t) = \tilde{z}^T(t)$ for $t \leq T$, and on sets B of $\mathscr{B}_T$ we have $\tilde{P}_x^T(B) = \tilde{P}_x^{T_1}(B)$. This follows from (05) since (χ_B is the characteristic function of the set B)

$$\begin{aligned}\tilde{P}_x^{T_1}(B) &= E_x[E_x(\chi_B\exp\zeta_o^{T_1}(\hat{f})|\mathscr{B}_T)] \\ &= E_x\chi_B\exp\zeta_o^T(\hat{f})[E_x(\exp\zeta_T^{T_1}(\hat{f})|\mathscr{B}_T)] \\ &= E_x\chi_B\exp\zeta_o^T(\hat{f}) = \tilde{P}_x^T(B).\end{aligned}$$

Thus $\tilde{P}_x^{T_1}$ is an extension of $\tilde{P}_x^T$. If (6) - (7) hold for each $T < \infty$, we can replace Ω_T by Ω and define a unique measure $\tilde{P}_x$ on all the $\mathscr{B}_t$, $t < \infty$, which will be consistent with the $\tilde{P}_x^T$ on $\mathscr{B}_T$. Then $\tilde{z}(t)$ will be an Itô process with respect to $(\Omega, \mathscr{B}_t, \tilde{P}_x)$, and $(\Omega, \tilde{z}(t), \mathscr{B}_t, \tilde{P}_x)$ an Itô process (for all $t < \infty$). Both (6) - (7) hold for all $T < \infty$ if $\hat{f}$ is bounded. Let $\mathscr{B} = \bigcup_{t \geq 0} \mathscr{B}_t$.

3. Existence of a Solution to (8) for Unbounded $\hat{f}$

Let $V(x)$ denote a non-negative twice continuously differentiable function which tends to infinity as $|x| \to \infty$. Define $Q_N = \{x: V(x) < N\}$ and let $\hat{f}_N(x) = \hat{f}(x)$ for $x \in Q_N$ and $\hat{f}_N(x) = 0$, $x \notin Q_N$. Define $C_N^T = \{\omega: x(t) \in Q_N, t \in [0,T]\}$. Let $\tilde{\mathcal{L}}$ denote the differential generator of the process (8) and write $\tilde{\mathcal{L}}^N$ for the differential generator when $\hat{f}$ is replaced by $\hat{f}^N$ in (8). Theorem 1 uses a stability idea to prove existence for (8), for all $t < \infty$.

Theorem 1. *Assume* (C1) *and the above conditions on* $V(x)$. *Let* $\tilde{\mathcal{L}}V(x) \leq 0$ *for* x *not in some* Q_a, $a < \infty$. *Then*

$$E_x \exp \zeta_0^T(\hat{f}) = 1 \tag{9}$$

for all $T < \infty$, *and*

$$\tilde{z}(t) = z(t) - \int_0^t \hat{\sigma}^{-1}(x(s))\hat{f}(x(s))ds$$

is a Wiener process, for all $t < \infty$ *with respect to* $(\Omega, \mathscr{B}_t, \tilde{P}_x)$. *The solution to* (8) *exists for all* $t < \infty$. *It is an Itô process with respect to* $(\Omega, \tilde{z}(t), \mathscr{B}_t, \tilde{P}_x)$ *and, under the additional assumptions* (C2-3), *it is unique* (*in the sense that the multivariate distributions of any two solutions are equal*) *for all* $t < \infty$.

Remark. Let $f(y), \sigma(y)$ satisfy (3), (2) locally, and let $\mathcal{L}_1$ denote the differential generator, with coefficients determined by $f(y), \sigma(y)$. If $V(x)$ and $\mathcal{L}_1 V(x)$ have the properties required in Theorem 1, then the proof can be altered to yield existence and uniqueness for the process $dy = f(y)dt + \sigma(y)dz$.

Proof. Let $\hat{f}^N$ replace $\hat{f}$, in (8), where $N > a$. Let $\tilde{P}_x^{N,T}$ denote the transformed measure with $\tilde{P}_x^{N,T}(A) = \int_A \exp \zeta_0^T(\hat{f}^N)dP$ and $\tilde{P}_x^N$ the extension of $\tilde{P}_x^{N,T}$ to the σ-algebra $\mathscr{B}$ on Ω. Write the Wiener process corresponding to $\tilde{P}_x^N$ as $\tilde{z}^N(t)$ (instead of $\tilde{z}(t)$). Then (8) is an Itô process with respect to $(\Omega, \tilde{z}^N(t), \mathscr{B}_t, \tilde{P}_x^N)$. By virtue of (03) (for $x = x(0)$)

$$\tilde{P}_x^N\{\sup_{t \geq s \geq 0} |x(s) - x| \geq \epsilon > 0\} \to 0 \tag{10}$$

as $t \to 0$, uniformly for x in compact intervals. Also $\tilde{\mathcal{L}}^N V(x) \leq 0$ in[+] $Q_N - Q_a - \partial Q_a \equiv Q_{N,a}$. Let τ denote the first exist[++] time of the path $x(t)$ from $Q_N - Q_a - \partial Q_a$, and $t \cap \tau \equiv \min(t, \tau)$. Then, by Itô's Lemma $\tilde{E}_x^N V(x(t \cap \tau)) - V(x) \leq 0$ for $x \in Q_N - Q_a$. Since

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[+] ∂Q_N is the boundary of the set Q_N.

[++] If τ is undefined for some path, set $\tau = +\infty$. Note that the exit time $\tau(\omega)$ (as a path function) for $x(t)$ and $\tilde{x}^N(t)$ are the same; but their distributions may differ.

$$\tilde{E}_x^N V(x(t \cap \tau)) - a \leq (N-a)\tilde{P}_x^N\{x(s) \text{ hits } \partial Q_N \text{ before } \partial Q_a \text{ and leaves } Q_{N,a} \text{ in } [0,t]\},$$

we can conclude that

$$(11) \quad \tilde{P}_x^N\{x^N(t) \text{ hits } \partial Q_N \text{ before } \partial Q_a \text{ and leaves } Q_{N,a} \text{ in } [0,T]\} \leq \frac{V(x)-a}{N-a} \equiv \epsilon_3.$$

We will show that for each $\epsilon > 0$, there is an $N < \infty$ so that

$$(12) \qquad \tilde{P}_x^N\{C_N^T\} \geq 1 - \epsilon.$$

Fix $a_1 > a$. Let $x \in \partial Q_a$. There is a $\delta_o > 0$ so that

$$\min_{x \in \partial Q_a, y \in \partial Q_{a_1}} |x-y| \geq \delta_o.$$

Let $A \in C_N^T$. Then, since $\hat{f}^N(x(t)) = \hat{f}^M(x(t))$ on $[0,T]$ for $M \geq N$ and $\omega \in C_N^T$, we have

$$(13) \qquad \tilde{P}_x^M(A) = E_x \exp \zeta_o^T(\hat{f}^M)\chi_A = E_x \exp \zeta_o^T(\hat{f}^N)\chi_A = \tilde{P}_x^N(A).$$

(03) implies that

$$\sup_{x\in\partial Q_a} \tilde{P}_x^N\{\max_{\delta_1\geq t\geq 0} |x(t) - x| \geq \delta_o\} \leq K_2 \frac{\delta_1^2}{\delta_o^4} = \epsilon_2.$$

But (13) implies that the constant K_2 depends only on the number a_1 and does not depend on N, for $N > a_1$. Thus, we can assume that K_2 does not depend on N.

Let G_N^T denote the event that $x(t)$ goes to ∂Q_a before ∂Q_N (<u>or never leaves</u> $Q_{N,a}$), then takes more time than $T/n \equiv \delta_1$ to reach ∂Q_{a_1}, then returns to ∂Q_a <u>no fewer than</u> $n - 1$ additional times and after each return takes no less than δ_1 to reach ∂Q_{a_1}, before leaving Q_N for the first time. Then $\tilde{P}_x^N\{C_N^T\} \geq \tilde{P}_x^N\{G_N^T\}$ and $\tilde{P}_x^N\{G_N^T\} \geq 1 - n(\epsilon_1+\epsilon_2) - \epsilon_3$, where

$$\epsilon_1 = \max_{x\in\partial Q_{a_1}} \tilde{P}_x^N\{x(t) \text{ reaches } \partial Q_N \text{ before } \partial Q_a\} \leq \frac{(a_1-a)}{(N-a)}.$$

Thus, using $\delta_1 = T/n$,

$$\tilde{P}_x^N\{G_N^T\} \geq 1 - n\left(\frac{a_1-a}{N-a} + \frac{K_2T^2}{n^2\delta_o^4}\right) - \frac{V(x)-a}{N-a}$$

and N and n can be chosen so that $\tilde{P}_x^N\{G_N^T\} \geq 1 - \epsilon$.

There is a unique measure $\tilde{P}_x^T$ on $\mathscr{B}_T$ which is consistent with the $\tilde{P}_x^N$ on the sets C_N^T. Furthermore, (the left hand inequality is [4], Lemma 2)

$$1 \geq \tilde{P}_x^T(\Omega_T) = E_x \exp \zeta_o^T(\hat{f}) \geq E_x \exp \zeta_o^T(\hat{f}^N)\chi_{C_N^T} \geq 1 - \epsilon.$$

Since ϵ is arbitrary, (9) holds, $\tilde{z}(t)$, $t \leq T$, is a Brownian motion with respect to $(\Omega_T, \mathscr{B}_t, \tilde{P}_x^T)$ and $x(t)$, $t \leq T$, an Itô process with respect to $(\Omega_T, \tilde{z}(t), \mathscr{B}_t, \tilde{P}_x^T)$. Furtheremore, since T is arbitrary, we can replace $t \leq T$ by $t < \infty$ and $\tilde{P}_x^T$ and Ω_T by $\tilde{P}_x$ and Ω.

The process (8) is unique in the following sense. Suppose that both $x^i(t)$, $i = 1,2$ satisfy (8). Let $x^{i,N}(t)$ denote the processes which result when $\hat{f}^N$ replaces $\hat{f}$. Suppose that if $x^{i,N}(t) \in Q_N$ for all $t \in [0,T]$, then $x^i(t)$ coincides with $x^{i,N}(t)$ on $[0,T]$. Then the uniqueness of the $x^{i,N}(t)$ (in the sense of multivariate distributions) and the fact that $\tilde{P}_x^N\{C_N^T\} = \tilde{P}_x^M\{C_N^T\} \geq 1 - \epsilon$ for $M > N$ (the $\tilde{P}_x^N$ do not depend on i) imply uniqueness of the $x^i(t)$ in the sense of multivariate distributions. Q.E.D.

<u>Remark</u>. Lemma 7 of [4] would appear to yield existence for a large class of unbounded $\hat{f}$. But an examination of the proof shows that its content is the following. Let processes (5) and (8) exist with respect to some Wiener process, with (5) being unique, and $\int_o^T |\hat{\sigma}^{-1}(x(t))\hat{f}(x(t))|^2 dt < \infty$ w.p.1, where $x(t)$ is the solution to (5). Under some minor subsidiary condition, it is proved that

$$E_x \exp \zeta_o^T(\hat{f}) = 1$$

where the expectation corresponds to (5). Then (8) can be obtained by a Girsanov transformation from (5). But both the square integrability property and existence for (8) must be established first. But these properties are essentially the desired result.

3. Markov Properties of (8)

Write (C4): In each compact x set, there is an $\alpha > 1$ and $M < \infty$ so that

$$\int p^{\alpha}(x;\ t,y) \leq M < \infty.$$

Theorem 2. Assume (C1) - (C3) and the condition on V and $\tilde{\mathcal{L}}V$ of Theorem 1. Then the process (8) is a strong Markov process.

If (C4) holds, for some $\alpha > 1$, (8) is a strong Feller process.

Proof. T he terminology of Theorem 1 will be used. By Theorem 1, the process is defined on the time interval $[0,\infty)$, and has continuous paths w.p.1.

First, we prove that (8) is a Markov process. Let $\mathscr{B}_t^x \subset \mathscr{B}_t$ measure $x(s)$, $s \leq t$. Define the transition function $\tilde{P}_x(x;\ t,A) = \tilde{P}_x\{x(t) \in A\}$. Since the right hand term of

$$\tilde{P}_x\{x(t) \in A\} = E_x \chi_{\{x(t) \in A\}} \exp \zeta_0^t(\hat{f})$$

is a Borel measurable function of x, so is $\tilde{P}(x; t,A)$ for each $A \in \mathscr{B}_t^x$. Now assume that $\hat{f}^N$ replaces $\hat{f}$. The Chapman-Kolmogorov equation holds since, by (05) and the fact that (5) is a Markov process,

$$\tilde{E}_x^N[\chi_{\{x(t+s)\in A\}} | \mathscr{B}_s^x] = E_x[\chi_{\{x(t+s)\in A\}} \exp \zeta_s^{s+t}(\hat{f}^N) | \mathscr{B}_s^x] .$$

$$= E_{x(s)}[\chi_{\{x(t)\in A\}} \exp \zeta_o^t(\hat{f}^N)] = \tilde{P}^N(x(s);t,A)$$

w.p.1. Thus by the definition Dynkin [6, Chapter 3], $x^N(t)$ (the Itô process on $(\Omega, \tilde{z}^N(t), \mathscr{B}_t, \tilde{P}_x^N)$ corresponding to the use of $\hat{f}^N$) is a Markov process.

The σ-algebras $\mathscr{B}_t^x$ also measure (8). The measure $\tilde{P}_x$ for the unbounded $\hat{f}$, has the correct conditioning properties since, by (05) and the dominated convergence theorem,

$$\tilde{E}_x[\chi_{\{x(t+s)\in A\}} \chi_{\{C_{t+s}^N\}} | \mathscr{B}_s^x]$$

$$= E_x[\chi_{\{x(t+s)\in A\}} \chi_{\{C_{t+s}^N\}} \exp \zeta_s^{t+s}(\hat{f}) | \mathscr{B}_s^x]$$

$$\to E_x[\chi_{\{x(t+s)\in A\}} \exp \zeta_s^{t+s}(\hat{f}) | \mathscr{B}_s^x] =$$

$$= E_{x(s)}[\chi_{\{x(t)\in A\}} \exp \zeta_o^t(\hat{f})] = \tilde{P}(x(s);t,A)$$

w.p.1. Then, by the definition [6, Chapter 3], (8) is a Markov process.

(8) is a Feller[+] process, hence a strong Markov process [6, Theorem 3.10]. The proof is omitted. The proof of the stronger 'strong' Feller[++] property will be given next, under the additional condition (C4). Let (C4) hold.

Supposing that (8) is a strong Feller process if $\hat{f}$ is bounded, we show that it is a strong Feller process for unbounded $\hat{f}$. Let $g(\cdot)$ be bounded and measurable. Then $\tilde{E}_x^N g(x(t)) \equiv G^N(x)$ is continuous in x, for $t > 0$. Write $G(x) = \tilde{E}_x g(x(t))$. Then

$$|G(x) - G^N(x)| \leq \max_x |g(x)| \cdot [\tilde{P}_x\{\Omega - C_N^T\} + \tilde{P}_x^N\{\Omega - C_N^T\}] \to 0 \quad \text{as} \quad N \to \infty,$$

uniformly in any compact x set. Thus, $G(x)$, being the uniform limit of continuous functions, is continuous.

Finally, <u>suppose</u> $\hat{f}$ <u>is bounded and</u> (C4) <u>holds</u>. Reproducing an argument of Rishel [2], we show for each compact x set, there is a $\beta > 1$ and $M < \infty$ so that (q is the density of (8) - see (O4))

$$\int q^\beta(x;\ t,y)\, dy \leq M_1 < \infty. \tag{14}$$

[+] A process $x(t)$ is a Feller process if $E_x f(x(t))$ is a continuous function of x, if $f(x)$ is continuous and bounded.

[++] $x(t)$ is a strong Feller process if $E_x f(x(t))$ is continuous in x for any bounded Borel function $f(x)$ and $t > 0$.

Define $r(x;\ t,y) \equiv \tilde{E}_x[\exp \zeta_c^t(\hat{f}) | x(t) = y]$. Let $m^{-1} + n^{-1} = 1$, and note that, for any $\rho > 1$ and compact x set, there is an $N_\rho < \infty$ so that $\tilde{E}_x \exp \rho\ \zeta_o^t(\hat{f}) \leq N_\rho$ ([4], Lemma 1). Let $\beta > \beta_1$, $\beta > 1$. By Holder's inequality

$$\int p^\beta(x;\ t,y) r^\beta(x;\ t,y) = \int p^{\beta_1}(x;\ t,y) r^\beta(x;\ t,y) p^{\beta-\beta_1}(x;\ t,y)dy$$

$$\leq [\int p^{\beta_1 n}(x;\ t,y) r^{\beta n}(x;\ t,y)dy]^{1/n} [\int p^{(\beta-\beta_1)m}(x;\ t,y)dy]^{1/m}.$$

We can choose $\beta > 1, \beta > \beta_1$, m, n and $\rho > 1$ so that $(\beta-\beta_1)m = \alpha$, $\beta n = \rho$, $\beta_1 n = 1$, which , together with (C4), proves (14). Equation (14) implies that, as x varies in any compact set, the family $q(x;\ t,y)$ of functions of y is uniformly integrable. This, together with the continuity (in x) of $\tilde{P}(x;\ t,(-\infty,a))$ for any vector a (recall that there is a density) implies that $\tilde{P}(x;\ t,A)$ is continuous in x for any Borel set A, which implies, in turn, the strong Feller property. For more detail, note that the boundary of any rectangle in the range space of $x(t)$ has zero probability, and that $\tilde{P}(x;\ t,A)$ is continuous in x if A is the sum of rectangles (open or closed). Let $\tilde{P}(x;\ t,A_j)$ be continuous in x for a collection of sets A_j, which increase monotonically to A

$$\tilde{P}(x;\ t,A) = \int_{A_j} q(x;\ t,y)dy + \int_{A-A_j} q(x;\ t,y)dy.$$

The second integral goes to zero as $j \to \infty$ uniformly in x in any

compact set, by the uniform integrability of $q(x; t,y)$. Since the first integral is continuous, so is the uniform limit $\tilde{P}(x; t,A)$. Q.E.D.

4. The Invariant Measure, and the Asymptotic Properties of the Measures of (8)

In [8], under the conditions (D1) - (D5), Khasminskii proved the existence of a unique σ-finite invariant measure for a process $x(t)$ with a stationary transition function $\tilde{P}(x; t,A)$ under the conditions (D1-5).

(D1) For any ϵ neighborhood $N_\epsilon(x)$ of x, $1 - P(x; t,N_\epsilon(x)) = o(t)$ uniformly in x in any compact set.

(D2) The process is a strong Markov and strong Feller process.

(D3) $\tilde{P}(x; t,U) > 0$ for all open sets U and $t > 0$.

(D4) The paths are continuous w.p.1.

(D5) The process is recurrent. (There is some compact set K and a random time $\tau < \infty$ w.p.1. so that $x(\tau) \in K$ w.p.1., for each initial condition.)

In [9], Kushner applied the result in [8] to obtain a sufficient condition for the convergence of the measures of class of diffusions to a unique invariant measure. Theorem 3 includes the prior result as a special case. Zakai [10] has treated the invariant measure problem for a class of diffusions satisfying (2) - (3), using

a general method of Benes [11]. A similar problem is treated in Elliot [12]. Elliot's method involves a condition on a Lie algebra generated by certain functions of the diffusion coefficients, which is hard to check in special cases. The result of Benes [11] (concerning only existence of an invariant measure) uses the condition that $\lim_{|x|\to\infty} P(x;\ t,K) \to 0$ for all compact sets K. This would not always hold under our conditions. E.g., the solution to $\dot{x} + x^3 = 0$, reaches $x = 1$ in a time that is bounded as $x(0) \to \infty$, and we would expect a pect a similar result for $dx = -x^3 dt + \sigma dz$.

Theorem 3. Assume (C1) - (C4), and the conditions on $V(\cdot)$ in Theorem 1, except let $\tilde{\mathcal{L}}V(x) \leq -\epsilon < 0$ outside of Q_a. Let (5) have a nowhere-zero density, for each initial condition x. Then (8) has a unique invariant measure $Q(\cdot)$ and $\tilde{P}(x;\ t,A) \to Q(A)$ as $t \to \infty$ for any x. Both $\tilde{P}(x;\ t,A)$ and $Q(A)$ have nowhere zero densities.

Remark. Theorem 3 only deals with invariant measures, but almost all of stability results in [1] can be carried over to the problem with discontinuous drift terms.

Proof. The second inequality of (03) implies (D1) for bounded $\hat{f}$, and, hence, for the processes $x^N(t)$. But, if (D1) holds for each $x^N(t)$, it holds for (8). (D2) is proved in Theorem 2. Since $\tilde{E}_x[\exp \zeta_0^t(\hat{f}) | x(t) = y] > 0$ w.p.1. and $p(x;\ t,y) > 0$ for

y by assumption, $q(x; t,y)$ (the density for $\tilde{P}(x; t,A)$) is positive for almost all y (Lebesgue measure). This implies (D3). (D4) is a consequence of Theorem 1. (D5) is a consequence of $\tilde{\mathcal{L}}V(x) \leq -\epsilon < 0$ for all large x. (See Theorem 4 in [9]). Indeed, the average time to leave the set $Q_N - Q_a - \partial Q_a$ (for $x(0) = x$) is bounded above by $(V(x) - a)/\epsilon < \infty$. This together with (11) gives (D5). Thus all (D1-5) hold.

$Q(A)$ satisfies

$$Q(A) = \int Q(dx)\tilde{P}(x; t,A)$$
$$= \int_A du \int Q(dx)q(x; t,u).$$

Thus $Q(A) > 0$ for all sets A of positive Lebesgue measure and has density $\int Q(dx)q(x; t,u)$, which must be positive almost everywhere.

For a process with a transition density and a unique invariant measure $Q(\cdot)$ with a nowhere zero density, Doob [7, Theorem 5] proves that $\tilde{P}(x; t,A) \to Q(A)$ as $t \to \infty$ for any x. Q.E.D.

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