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NONLINEAR SEMIGROUP FOR CONIROLLED
PARTIALLY OBSERVED DIFFUSIONS ${ }^{+}$
by

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Abstract: In this paper a "separated" control problem associated with controlled, partially observed diffusion processes is considered. The state in the separated problem is an unnormalized conditional distribution measure. The corresponding Nisio nonlinear semigroup associated with the separated problem is found.


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# NONLINEAR SEMIGROUP FOR CONTROLLED <br> PARTIALLY OBSERVED DIFFUSIONS 

Wendell H. Fleming

1. Introduction. In this paper we are concerned with stochastic control problems of the following kind. Let $X_{t}$ denote the state of a process being controlled, $Y_{t}$ the observation process, and $U_{t}$ the control process, $t \geq 0$. The state and observation processes are governed by stochastic differential equations
(a) $d X_{t}=b\left(X_{t}, U_{t}\right) d t+\sigma\left(X_{t}\right) d W_{t}$
(b) $d Y_{t}=h\left(X_{t}\right) d t+d \tilde{W}_{t}$.
$X_{t}$ has values in $N$-dimensional $\mathbb{R}^{N}, Y_{t}$ values in $\mathbb{R}^{M}$, and $U_{t}$ values in $\mathscr{K} \subset \mathbb{R}^{\mathrm{L}}$. $X_{0}$ has given distribution $\mu$, and $Y_{0}=0$. In (1.1), $W$ and $\tilde{W}$ are independent standard Wiener processes, with values in $\mathbb{R}^{\mathrm{D}}, \mathbb{R}^{\mathrm{M}}$ respectively. The problem is to find an admissible control minimizing some criterion $J$.

For instance, we may take $J=E G\left(X_{t_{1}}\right)$ for some fixed time $t_{1}>0$. In case of completely observed, controlled diffusions (with $Y_{t}=X_{t}$ rather than $Y_{t}$ as in (1.1b)), the problem can be treated using dynamic programming. Let $V\left(x, t_{1}\right)$ denote the minimum of $J$, for initial data $X_{0}=x$. Under suitable asssumptions $V(x, t)$ has continuous partial derivatves
$\partial V / \partial t, \partial V / \partial x_{i}, \partial^{2} V / \partial x_{i} \partial x_{j}, i, j:=1, \ldots, N, x=\left(x_{1}, \ldots, x_{N}\right)$. Among these assumptions is the condition that the symmetric matrix $a=\sigma \sigma$ ' has a bounded inverse $a^{-1}$. The function $V$ then satisfies the dynamic programming equation [4, Chap. VI.6]

$$
\begin{equation*}
\frac{\partial V}{\partial t}=L V, \tag{1.2}
\end{equation*}
$$

$$
\begin{equation*}
L V=\min _{u \in \mathscr{U}}\left[\frac{1}{2} \sum_{i, j=1}^{N} a_{i j}(x) \frac{\partial^{2} V}{\partial x_{i} \partial x_{j}}+\sum_{i=1}^{N} b_{i}(x, u) \frac{\partial V}{\partial x_{i}}\right. \tag{1.3}
\end{equation*}
$$

The assumption that $a(x)$ has a bounded inverse can sometimes be weakened, by considering generalized solutions to the dynamic programming equation [4, p. 177].

In [6] Nisio introduced another treatment which is valid under much less restrictive conditions. Let $\quad \mathscr{S}_{t} G(x)=V(x, t)$. Then Nisio showed that $\mathscr{S}_{t}$ is a nonlinear semigroup on the space $C_{b}\left(\mathbb{R}^{N}\right)$ of continuous bounded functions $f$ on $\mathbb{R}^{N}$. Moreover, the operator $L$ in (1.3) agrees with the generator of the semigroup $\mathscr{S}_{t}$ on the space $C_{b}^{2}\left(\mathbb{R}^{N}\right)$ of those $f$ such that $f, f_{x_{i}}, f_{x_{i}} x_{j}$ are in $C_{b}\left(\mathbb{R}^{N}\right)$ for $i, j=1, \ldots, N$. For another treatment of this nonlinear semigroup see [1, Chap. IV.5.1].

In this paper, we find a nonlinear semigroup $\mathscr{T}_{t}$ associated with the partially observed control problem. In this case, one should regard as the true "state" the conditional distribution of $X_{t}$ given past data, or some quantity equivalent to the conditional distribution. For technical reasons, it is more convenient to consider an unnormalized conditional distribution $\Lambda_{t}$ for $X_{t}$.

We have $\Lambda_{t} \in \mathbb{M}$, where $\mathbb{M}$ is the space of finite measures on $\mathbb{R}^{N}$. The problem we consider is to control the measure-valued process $\Lambda_{t}$ such that a criterion of the form $J=E \phi\left(\Lambda_{t_{1}}\right)$ is minimized. The dynamics of the $\Lambda_{t}$-process are governed by the Zakai equation, written in a weak form as (3.1) below.

If one writes $V\left(\mu, t_{1}\right)$ for the minimum of $J$, given initial data $\Lambda_{0}=\mu$, then $V(\mu, t)$ formally satisfies a dynamic programming equation of the form.

$$
\begin{equation*}
\frac{\partial \mathrm{V}}{\partial \mathrm{t}}=\mathscr{L} \mathrm{V} \tag{1.4}
\end{equation*}
$$

where $\mathscr{L} \mathrm{V}=\min _{\mathrm{u} \in \mathscr{U}} \mathscr{L}^{\mathrm{u}} \mathrm{v}$ and $\mathscr{L}^{\mathrm{u}}$ is the generator of the linear semigroup $\mathscr{F}_{t}^{u}$ associated for a constant control $u$ with the process $\Lambda_{t}$ (for constant $u, \Lambda_{t}$ is Markov). Equation (1.4) is called Mortensen's equation. However, (1.4) has been treated rigorously only in very special cases.

Following Nisio, we write $V(\mu, t)=\mathscr{T}_{t} \phi(\mu)$. The purpose of this paper is to show that $\mathscr{F}_{t}$ is a non1inear semigroup, on a space $C(\mathbb{M})$, with $\mathscr{T}_{\mathrm{t}}{ }^{\phi}$ continuous in $t$, and to describe the generator $\mathscr{L}$ on a dense subspace of $C(\mathscr{M})$. We rely heavily on results from [3]. In particular, it was shown in [3] that $\Lambda_{t}$ can be defined pathwise, in such a way that $\Lambda_{t}$ depends continuously on observation and control trajectories $(Y, U)$ and on $\mu=\Lambda_{0}$. This and other results from [ 3 ] needed in this paper are summarized as 3.1-3.4 below.

For the case of a controlled Markov chain $X_{t}$, subject to observations $Y_{t}$ of the form (1.1b) a corresponding nonlinear
semigroup was constructed by Davis [2].
2. The Spaces $C_{K}(\mathcal{M}), C(\mathcal{M})$. Let $C_{b}\left(\mathbb{R}^{N}\right)$ denote the space of bounded, continuous $f$ on $\mathbb{R}^{N}$, and $C_{0}\left(\mathbb{R}^{N}\right)$ the space of continuous $f$ with compact support. Let $C_{b}^{k}\left(\mathbb{R}^{N}\right), C_{0}^{k}\left(\mathbb{R}^{N}\right)$ be the spaces of $f$ such that $f$ together with all partial derivatives of orders $\leq k$ are in $C_{b}\left(\mathbb{R}^{N}\right), C_{0}\left(\mathbb{R}^{N}\right)$ respectively. Similarly, for $\mathbb{R}^{m}$ valued functions on $\mathbb{R}^{N}$ we write $C_{b}^{k}\left(\mathbb{R}^{N} ; \mathbb{R}^{m}\right), C_{0}^{k}\left(\mathbb{R}^{N} ; \mathbb{R}^{m}\right)$. Let $\mathscr{O}\left(\mathbb{R}^{\mathrm{N}}\right)$ denote the Bore $\sigma$-algebra of $\mathbb{R}^{\mathrm{N}}$, and

$$
\begin{equation*}
\mathscr{M}=\left\{\text { measures } \mu \geq 0 \text { on } \mathscr{B}\left(\mathbb{R}^{\mathrm{N}}\right): \mu\left(\mathbb{R}^{\mathrm{N}}\right)<\infty\right\} \tag{2.1}
\end{equation*}
$$

We write

$$
\langle f, \mu\rangle=\int_{\mathbb{R}^{N}} f(x) d \mu(x)
$$

for the scalar product and

$$
\|\mu\|=\langle 1, \mu\rangle=\mu\left(\mathbb{R}^{N}\right) .
$$

By convergence of sequences in $/ \mathbb{w e}$ mean $w^{*}$-convergence: $\mu_{n} \rightarrow \mu$ if and only if $\left\langle f, \mu_{n}\right\rangle \rightarrow\langle f, \mu\rangle$ as $n \rightarrow \infty$ for every $f \in C_{b}\left(\mathbb{R}^{N}\right)$ such that $f(x) \rightarrow 0$ as $|x| \rightarrow \infty$.

We denote real-valued functions on $\mathscr{M}$ by $\phi, \psi, \ldots$. For $K=0,1,2, \ldots 1 \mathrm{et}$

$$
\begin{equation*}
\|\phi\|_{K}=\sup _{\mu \in \mathcal{M}} \frac{|\phi(\mu)|}{1+||\mu||^{K}} \tag{2.2}
\end{equation*}
$$

By $\phi$ continuous on $\mathscr{M}$, we mean of course continuity of $\phi$ under $w^{*}$-sequential convergence. Let

$$
\begin{equation*}
\mathrm{C}_{\mathrm{K}}(\mathscr{M})=\left\{\phi \quad \text { continuous on } \mathscr{M}:\|\phi\|_{K}<\infty\right\} \tag{2.3}
\end{equation*}
$$

Then, $\left\|\|_{K}\right.$ is a norm on $C_{K}(\mathscr{M})$. Let

$$
\begin{equation*}
\mathrm{C}(\mathscr{M})=\cdot \bigcup_{K=0}^{\infty} \mathrm{C}_{\mathrm{K}}(\mathscr{M}) \tag{2.4}
\end{equation*}
$$

For $r<\infty$, let

$$
\begin{equation*}
\mathscr{M}_{\mathrm{r}}=\{\mu \in \mathscr{M}:\|\mu\| \leq \mathrm{r}\} \tag{2.5}
\end{equation*}
$$

We give $C(\mathscr{M})$ the following metric

$$
\begin{equation*}
\mathrm{d}(\phi, \psi)=\sum_{\ell=1}^{\infty} 2^{-\ell}\left(\sup _{\ell}|\phi(\mu)-\psi(\mu)| \wedge 1\right) . \tag{2.6}
\end{equation*}
$$

Thus d-convergence of $\phi_{n}$ to $\phi$ is equivalent to convergence of $\phi_{\mathrm{n}}(\mu)$ to $\phi(\mu)$ uniformly on $\mu_{r}$ for every $r<\infty$. For each $K$, $\left\|\|_{K}\right.$ is a lower semicontinuous function under d-convergence. Moreover, from (2.2), $\phi_{n}, \phi \in C_{K}(\mathscr{M})$ and $\left\|\phi_{n}-\phi\right\|_{K} \rightarrow 0$ imply $\mathrm{d}\left(\phi_{\mathrm{n}}, \phi\right) \rightarrow 0$ as $\mathrm{n} \rightarrow \infty$.

Let

$$
\tilde{\mu}=\left\{\mu \geq 0 \text { on } \mathscr{B}\left(\mathbb{R}^{N}\right): \mu(B)<\infty \text { for every compact } B\right\},
$$

with the vague topology: $\mu_{n} \rightarrow \mu$ vaguely means $\left\langle f, \mu_{n}\right\rangle+\langle f, \mu\rangle$
as $n \rightarrow \infty$ for every $f \in C_{0}\left(\mathbb{R}^{N}\right)$. $\tilde{\mu}$ is a Polish space. In fact, one can choose a metric $\delta(\mu, \nu)$ for $\tilde{M}$ of the form

$$
\begin{equation*}
\delta(\mu, v)=\sum_{m=1}^{\infty} 2^{-m}\left(\left|\left\langle f_{m}, \mu\right\rangle-\left\langle f_{m}, v\right\rangle\right| \wedge 1\right) \tag{2.7}
\end{equation*}
$$

for a suitably chosen sequence $f_{m} \in C_{0}\left(\mathbb{R}^{N}\right)$.
For each $r<\infty, \mathcal{M}_{r}$ is a compact subset of $\mathscr{M}$. For sequences in $\mathscr{M}_{r}$, vague convergence is equivalent to $w^{*}$-convergence. Moreover $\mu_{n}, \mu \in \mathcal{M}$ and $\mu_{n} \rightarrow \mu \quad w^{*}$ imply $\left|\left|\mu_{n}\right|\right| \leq r$ for some some r. Thus, we have:

Lemma 2.1. $\Phi$ is continuous on $\mathscr{M}$, under $w^{*}$-sequential convergence, if and only if $\phi \mid \mathscr{M}_{r}$ is vaguely continuous for every $\mathbf{r}<\infty$.

This furnishes an alternate characterization of $C(\mathbb{M})$, in terms of the vague topology rather than in terms of $w^{*}$-sequntiei convergence.

A measure $\mu \in \mathcal{M}$ can be approximated by measures $\rho \mu$ with compact support, as follows. Let $\rho \in C_{0}\left(\mathbb{R}^{N}\right), 0 \leq \rho \leq 1$, and define $\rho \mu$ by $\langle f, \rho \mu\rangle=\langle\rho f, \mu\rangle$ for all $f \in C_{b}\left(\mathbb{R}^{N}\right)$. Define $\phi^{\rho}$ by

$$
\begin{equation*}
\phi^{\rho}(\mu)=\phi(\rho \mu), \quad \mu \in \mathscr{M} \tag{2.8}
\end{equation*}
$$

Then $\phi \in C_{K}(\mathscr{M})$ implies $\phi^{\rho} \in C_{K}(\mathcal{M})$ and $\left\|\phi^{\rho}\right\|_{K} \leq\|\phi\|_{K}$ We write $\mu \mid B$ for the restriction of $\mu$ to a compact set $B:(\mu \mid B)(A)=\mu(A \cap B)$ for all $\left.A \in \mathscr{O} \mathbb{R}^{N}\right)$. Let

$$
\begin{align*}
\mathrm{C}_{\mathrm{K}}^{0}(\mathscr{M})= & \left\{\psi \in \mathrm{C}_{\mathrm{K}}(\mathscr{M}): \text { there exists } \mathrm{B}\right. \text { compact such }  \tag{2.9}\\
& \text { that } \psi(\mu)=\psi(\mu \mid B) \text { for all } \mu \in \mathscr{M}\}
\end{align*}
$$

In particular, $\phi^{\rho} \in C_{K}^{0}(\mathscr{M})$ if $\phi \in C_{K}(\mathscr{M})$ and $\phi^{\rho}$ is defined by (2.8) .

Lemma 2.2. For every $\phi \in C_{K}(\mathbb{M})$ there exists a sequence $\phi_{\mathrm{n}} \in \mathrm{C}_{\mathrm{K}}^{0}(\mathscr{M})$ such that $\left\|\phi_{\mathrm{n}}\right\|_{\mathrm{K}} \leq\|\phi\|_{K}$ and $d\left(\phi_{\mathrm{n}}, \phi\right) \rightarrow 0$ as $\mathrm{n} \rightarrow \infty$.

Proof. Let $\rho_{n} \in C_{0}\left(\mathbb{R}^{N}\right)$ with $0 \leq \rho_{n} \leq 1, \rho_{n}(x)=1$ for $|x| \leq n$ and $\rho_{n}(x)=0$ for $|x| \geq n+1$. Let $\phi_{n}=\phi^{\rho} n$. Then $\left\|\phi_{\mathrm{n}}\right\|_{\mathrm{K}} \leq\|\phi\|_{\mathrm{K}}$. Since $\phi_{\mathrm{n}}(\mu)=\phi\left(\rho_{\mathrm{n}}{ }^{\mu}\right)$ it suffices to show that $\phi\left(\rho_{\mathrm{n}} \mu\right)-\phi(\mu)$ tends to 0 uniformly on $\quad \mu_{r}$ for every $r<\infty$. Let

$$
\eta_{n}=\max _{\boldsymbol{\mu}_{\mathbf{r}}}\left|\phi\left(\rho_{\mathrm{n}} \mu\right)-\phi(\mu)\right|=\left|\phi\left(\rho_{\mathrm{n}} \mu_{\mathrm{n}}\right)-\phi\left(\mu_{\mathrm{n}}\right)\right|
$$

for some $\mu_{n} \in \mathscr{M}_{r}$ (recall that $\mathscr{M}_{r}$ is compact). We have $\rho_{n} \mu_{n} \in \mathscr{M}_{r}$. For each $f \in C_{0}\left(\mathbb{R}^{N}\right),\left\langle f, \rho_{n} \mu_{n}\right\rangle=\left\langle f, \mu_{n}\right\rangle$ for all large enough $n$. Consider any subsequence such that $\mu_{n}$ tends to a limit $\mu$. Then $\rho_{n} \mu_{n}$ also tends to $\mu$ for $n$ in this subsequence. Since $\phi \mid \mu_{r}$ is continuous, both $\phi\left(\rho_{n}{ }_{n}\right)$ and $\phi\left(\mu_{n}\right)$ tend to $\phi(\mu)$. If $1 \mathrm{im} \sup _{n \rightarrow \infty} \eta_{\mathrm{n}}>0$, we could find some such subsequence for which $\left|\phi\left(\rho_{n} \mu_{n}\right)-\phi\left(\mu_{n}\right)\right|$ tends to a positive limit, a contradiction. This proves Lemma 2.2 .

Lemma 2.3. Let $\psi \in C_{K}^{0}(\mathcal{M})$, and $B$ compact such that $\psi(\mu)=\psi(\mu \mid B)$ for all $\mu \in \mathscr{M}$. Then there exists a sequence $\psi_{\mathrm{n}} \in \mathrm{C}_{\mathrm{K}}^{0}(\mathcal{M})$ such that $\left\|\psi_{\mathrm{n}}\right\|_{\mathrm{K}} \leq\|\psi\|_{\mathrm{K}}, \mathrm{d}\left(\psi_{\mathrm{n}}, \psi\right) \rightarrow 0$ as $\mathrm{n} \rightarrow \infty$ and $\psi_{n}(\mu)=0$ whenever $\mu(B) \geq n$.

Proof. Choose $\rho \in C_{0}\left(\mathbb{R}^{N}\right)$ with $0 \leq \rho \leq 1, \rho(x)=1$ for all $x \in B$. Let $g_{n} \in C_{0}\left(\mathbb{R}^{1}\right)$, with $0 \leq g_{n} \leq 1, g_{n}(s)=1$ if $\mathrm{s} \leq \mathrm{n}-1, \mathrm{~g}_{\mathrm{n}}(\mathrm{s})=0$. if $\mathrm{s} \geq \mathrm{n}$. Let

$$
\psi_{n}(\mu)=g_{n}(\langle\rho, \mu\rangle) \psi(\mu)
$$

Since $\left|\psi_{n}(\mu)\right| \leq|\psi(\mu)|,\left\|\psi_{n}\right\|_{K} \leq \| \psi| |_{K}$. For $\mu \in \mu_{r}$, $\langle\rho, \mu\rangle \leq r$. Hence $\psi_{n}(\mu)=\psi(\mu)$ if $n \geq r+1$, which implies $\psi_{n} \rightarrow \psi$ uniformly on $\mu_{r}$. Thus $d\left(\psi_{n}, \psi\right) \rightarrow 0$ as $n \rightarrow \infty$. Finally, $\mu(B) \geq n$ implies $\langle\rho, \mu\rangle \geq n$, and hence $\psi_{n}(\mu)=0$. This proves Lemma 2.3.

The set $\mathscr{D}$ of "test functions". In 55 we shall define a "generator" for the nonlinear semigroup on the following set of functions $\phi$, depending on finitely many scalar products:

$$
\begin{align*}
\mathscr{D}=\{\phi: \phi(\mu) & =F\left(\left\langle f_{1}, \mu\right\rangle, \ldots,\left\langle f_{J}, \mu\right\rangle\right),  \tag{2.10}\\
& \left.F \in C_{b}^{\infty}\left(\mathbb{R}^{J}\right), f_{1}, \ldots, f_{J} \in C_{0}^{\infty}\left(\mathbb{R}^{N}\right), J=1,2, \ldots\right\} .
\end{align*}
$$

In $\S 4$, we shall weaken slightly the conditions on $F, f_{1}, \ldots, f_{J}$, to obtain certain sets $\mathscr{D}_{\mathrm{m}}$ containing $\mathscr{D}$.

Lemma 2.4. For every $\phi \in l_{K}(\mathcal{M})$ therecists a sequence $\psi_{\mathrm{n}} \in \mathscr{D}$ such that $\left\|\psi_{\mathrm{n}}\right\|_{\mathrm{K}} \leqslant\|\psi\|_{K}+\mathrm{n}^{-1}$ and $d\left(\psi_{\mathrm{n}}, \psi\right) \rightarrow 0$ as $n \rightarrow \infty$.

Proof. By Lemmas 2.2 and 2.3 it suffices to suppose that, in addition, there exist compact $B$ and $a>0$ such that $\phi(\mu)=\phi(\mu \mid B)$ for all $\mu$ and $\phi(\mu)=0$ if $\mu(B) \geq$ a. Following a similar construction in [5, §3], given $\varepsilon>0$, we take $g_{1}, \ldots, g_{J}, x_{1}, \ldots, x_{J}$ with the following properties:

$$
\begin{aligned}
& g_{j} \in C_{0}^{\infty}\left(\mathbb{R}^{N}\right), g_{j} \geq 0, \operatorname{diam}\left(\operatorname{spt} g_{j}\right)<\varepsilon \\
& \sum_{j=1}^{J} g_{j}(x) \leq 1 \text { for } x \in \mathbb{R}^{N}, \sum_{j=1}^{J} g_{j}(x)=1, \quad x \in B, \\
& x_{j} \in B \cap \operatorname{spt} g_{j}
\end{aligned}
$$

Let

$$
\begin{aligned}
& \mathbb{R}_{+}^{J}=\left\{z \in \mathbb{R}^{J}: z_{j} \geq 0 \text { for } j=1, \ldots, J\right\}, \\
& \tilde{F}(z)=\phi\left(\sum_{j=1}^{J} z_{j} \delta_{x_{j}}\right)
\end{aligned}
$$

where $\delta_{x}$ denotes the Dirac measure at $x$. Then $\tilde{F} \in C_{0}\left(\mathbb{R}_{+}^{J}\right)$. In fact, $\tilde{F}(z)=0$ whenever

$$
\sum_{j=1}^{J} z_{j} \delta_{x_{j}}(B)=\sum_{j=1}^{J} z_{j} \geq a
$$

By regularizing, there exists $F \in C_{0}^{\infty}\left(\mathbb{R}^{J}\right)$ such that $|F(z)-\tilde{F}(z)| \leq \varepsilon$ for all $z \in \mathbb{R}_{+}^{J}$. Then

$$
\psi(\mu)=F\left(\left\langle g_{1}, \mu\right\rangle, \ldots,\left\langle g_{J}, \mu\right\rangle\right)
$$

is in $\mathscr{D}$. For all $\mu$,

$$
\begin{aligned}
|\psi(\mu)| & \leq\left|\phi\left(\sum_{j=1}^{J}<g_{j}, \mu>\delta_{x_{j}}\right)\right|+\varepsilon \\
& \leq\|\phi\|_{K}\left(1+\left\|\sum_{j=1}^{J}<g_{j}, \mu>\delta_{x_{j}}\right\|^{K}\right)+\varepsilon \\
& \leq\|\phi\|_{K}\left(1+\left.\|\mu\|\right|^{K}\right)+\varepsilon .
\end{aligned}
$$

Therefore, $\left|\mid \psi\left\|_{K} \leq\right\| \phi \|_{K}+\varepsilon\right.$.
We take $\varepsilon=\varepsilon_{n}=n^{-1}$, and corresponding $g_{j n}, x_{j n}, j=1, \ldots, J_{n}$.
The corresponding $\psi_{n}$ obtained from the construction above has the properties required in Lemma 2.4. To show that $d\left(\psi_{n}, \phi\right) \rightarrow 0$, it suffices to show that $\psi_{n}(\mu) \rightarrow \phi(\mu)$ uniformly on $\mu_{r}$ for any $r>0$ as $n \rightarrow \infty$. Now

$$
\begin{aligned}
& \left|\psi_{\mathrm{n}}(\mu)-\phi\left[G_{\mathrm{n}}(\mu)\right]\right| \leq \varepsilon_{\mathrm{n}} \\
& \mathrm{G}_{\mathrm{n}}(\mu)=\sum_{j=1}^{J}\left\langle\mathrm{~g}_{\mathrm{jn}}, \mu>\delta_{x_{j n}} .\right.
\end{aligned}
$$

On $\mathscr{M}_{r}$, both vague and $w^{*}$-convergence of a sequence are equivalent to convergence in the metric $\delta$ in (2.7). For each $m$, $\mid\left\langle f_{m}, G_{n}(\mu)\right\rangle-\left\langle f_{m}, \rho_{n} \mu\right\rangle \rightarrow 0$ as $n \rightarrow \infty$ uniform by for $\mu \in \mathcal{M}_{r}$,
where $\rho_{n}=\sum_{j} g_{j n}$. Therefore, $\delta\left(G_{n}(\mu), \rho_{n}{ }^{\mu}\right) \rightarrow 0$ uniformly on $\quad \mu_{r}$ as $n \rightarrow \infty$. Since $\mathscr{M}_{r}$ is compact and $\phi$ continuous on $\mathscr{M}_{r}, \phi$ is uniformly continuous on $\mu_{r}$. Thus, $\left|\phi\left[G_{n}(\mu)\right]-\phi\left(\rho_{n} \mu\right)\right| \rightarrow 0$ uniformly on $\mu_{r}$. Since $\rho_{n}(x)=1$ on $B, \phi\left(\rho_{n} \mu\right)=\phi(\mu)=\phi(\mu \mid B)$. This proves that $\psi_{n}(\mu) \rightarrow \phi(\mu)$ uniformly on $\mu_{r}$, as required.
3. The Control Problem for $\Lambda_{t}$. We begin with a summary of assumptions and notations, together with a review of concepts from [3]. We make the same assumptions as in [3] about the coefficients in (1.1):
$\left(A_{1}\right) \quad \sigma$ is a bounded, Lipschitz $N \times D$ matrix-valued function on $\mathbb{R}^{N}$.
(A $) b(x, u)=b^{0}(x)+b^{1}(x) u$, where $b^{0}, b^{1}$ are bounded, Lipschitz functions on $\mathbb{R}^{N}$.

Note that $b^{0}$ has values in $\mathbb{R}^{N}$, and $b^{1}$ has $N \times L$ matrices as values. In $\S 5$, we shall impose additional smoothness conditions on $\sigma, b^{0}, b^{1}$.

$$
\left(A_{3}\right) \quad h \in C_{b}^{2}\left(\mathbb{R}^{N} ; \mathbb{R}^{M}\right)
$$

$\left(A_{4}\right) \mathscr{U}$ is a convex, compact subset of $\mathbb{R}^{L}$.
We use $Y$ to denote an $\mathbb{R}^{M}$-valued function, and $U$ a $\mathscr{K}$-valued function, of time $t \geq 0$. Let $Y_{t}, U_{t}$ denote their respective values at time $t$. Let

$$
\begin{gathered}
\Omega=\left\{(Y, U): Y_{0}=0, Y \in C\left([0, \infty): \mathbb{R}^{M}\right), U \in L^{2}([0, T] ; \mathscr{U})\right. \text { for } \\
\text { for each } T<\infty\} .
\end{gathered}
$$

Let $\Omega_{\mathrm{T}}$ denote the set of restrictions to $[0, \mathrm{~T}]$ of functions $(Y, U) \in \Omega$. As in [3], we give $\Omega_{T}$ a metric in which convergence of a sequence $\left(Y_{n}, U_{n}\right)$ means uniform convergence on $[0, T]$ of $Y_{n}$ and weak convergence of $U_{n}$ in $L^{2}([0, T] ; \mathscr{U})$. We give $\Omega$ a metric in which convergence of $\left(Y_{n}, U_{n}\right)$ is equivalent to convergence of $\left(Y_{n}, U_{n}\right)$ restricted to $[0, T]$ for every $T<\infty$. Let

$$
\begin{aligned}
\mathscr{F}_{\mathrm{t}}(\mathrm{Y}) & =. \sigma\left\{\mathrm{Y}_{\mathrm{s}}, 0 \leq \mathrm{s} \leq \mathrm{t}\right\} \\
\mathscr{F}_{\mathrm{t}}(\mathrm{U}) & =\sigma\left\{\mathrm{V}_{\mathrm{s}}, 0 \leq \mathrm{s} \leq \mathrm{t}\right\}, \quad \mathrm{V}_{\mathrm{t}}=\int_{0}^{\mathrm{t}} \mathrm{~V}_{\theta} \mathrm{d} \theta, \\
\mathscr{G}_{\mathrm{t}} & =\mathscr{F}_{\mathrm{t}}(\mathrm{Y}) \times \mathscr{F}_{\mathrm{t}}(\mathrm{U}) .
\end{aligned}
$$

These are $\sigma$-algebras of subsets of $\Omega$. However, if $t \leq T$, they can also be regarded as $\sigma$-algebras of subsets of $\Omega_{\mathrm{T}}$. In [3], $\Omega_{\mathrm{T}}$ was denoted by $\Omega^{2}$ and $\mathscr{G}_{\mathrm{t}}$ by $\mathscr{G}_{\mathrm{t}}^{2}$.

Let $\mathscr{S}_{\infty}$ be the least $\sigma$-algebra containing $\mathscr{G}_{t}$ for all $t \geq 0$.

Definition. An admissible control on $[0, T]$ is a probability measure $\pi_{T}$ on $\left(\Omega_{\mathrm{T}}, \mathscr{G}_{\mathrm{T}}\right)$, such that Y is a $\pi_{\mathrm{T}},\left\{\mathscr{G}_{\mathrm{t}}\right\}$-Wiener process for $0 \leq t \leq T$.

An admissible control is a probability measure $\pi$ on $\left(\Omega, \mathscr{C}_{\infty}\right)$ such that $Y$ is a $\pi,\left\{\mathscr{G}_{\mathrm{t}}\right\}$-Wiener process for $\mathrm{t} \geq 0$.

The definition of admissible control on $[0, T]$ is exactly as in [3]. If $\pi$ is an admissible control, then its restriction $\pi_{T}$ to $\mathscr{G}_{\mathrm{T}}$ is admissible on $[0, T]$.

Let $\mathscr{I}_{\mathrm{T}}$ denote the set of all admissible controls $\pi_{T}$ on $[0, T]$.

Then $\mathscr{A}_{\mathrm{T}}$ is compact under weak sequential convergence of probability measures [3, Lemma 2.3]. Let $\mathscr{A}$ denote the set of all admissible controls with the weak sequential convergence topology. Then $\mathscr{A}$ is a compact metric space under (for instance) the Prokhorov metric. Moreover, $\pi_{n} \rightarrow \pi$ if and only if the restrictions $\pi_{n, T}$ tend to $\pi_{T}$ as $n \rightarrow \infty$ for each $T$ finite.

The unnormalized conditional distribution measure $\Lambda_{t}$. For every $\mu \in \mathscr{M},(Y, U) \in \Omega$, and $t \geq 0$, we define $\Lambda_{t}=\Lambda_{t \mu}^{Y U}$ by formula [3, (3.9)]. (In [3] we wrote $\Lambda_{t}^{Y U}$, but now we wish to emphasize its dependence on the initial value $\mu=\Lambda_{0}$.) From its definition, $\Lambda_{t} \in \mathscr{H}$ and $\Lambda_{t}$ is $\mathscr{S}_{t}$-measurable as a function of $(Y, U) \in \Omega$. In [3,§3] we interpreted $\Lambda_{t}$ as an unnormalized conditional distribution of $X_{t}$ in (1.la) with respect to the $\sigma$-algebra $\mathscr{G}_{\mathrm{t}}$ generated by the observation and control past up to $t$. The normalized conditional distribution of $X_{t}$ is $\left\|\Lambda_{t}\right\|^{-1} \Lambda_{t}$. The intuitive reason for conditioning on $\mathscr{G}_{t}$, rather than on $\mathscr{F}_{\mathrm{t}}(\mathrm{Y})$, is that $\mathrm{U}_{\mathrm{t}}$ is not necessarily $\mathscr{F}_{\mathrm{t}}(\mathrm{Y})$-measurable $\pi$-almost surely, when $\pi \in \mathscr{A}$. For the smaller class of strictsense admissible controls $[3, \S 6]$ one can condition on $\mathscr{F}_{t}(Y)$ instead of $\mathscr{G}_{\mathrm{t}}$.

We shall need the following properties of $\Lambda_{t}$, proved in [3].
3.1. For each $t \geq 0, r<\infty, \Lambda_{t \mu}^{Y U}$ is continuous on $\mu_{r} \times \Omega$. See [3, Lemma 3.2].
3.2. For each finite $\mathrm{T}, \mathrm{r}, \mathrm{a}$ there exists $\rho=\rho(\mathrm{T}, \mathrm{r}, \mathrm{a})$ such that $0 \leq t \leq T,\|\mu\| \leq r,\|Y\|_{T} \leq a \operatorname{imply}\left\|\Lambda_{t \mu}^{Y U}\right\| \leq \rho$. Here $\|Y\|_{T}=\max _{0 \leq t \leq T}|Y(t)|$. See $[3,(3.6)] ;$ since $\Lambda_{t}$ depends linearly
on $\mu=\Lambda_{0}$ it suffices to consider $\|\mu\|=1$.
3.3. The Zakai equation holds:

$$
\begin{equation*}
\left.d<f, \Lambda_{t}\right\rangle=\left\langle L^{U} t_{f, \Lambda_{t}}>d t+\left\langle h f, \Lambda_{t}>\cdot d Y_{t}, \text { all } f \in C_{b}^{2}\left(\mathbb{R}^{N}\right)\right.\right. \tag{3.1}
\end{equation*}
$$

See [3, Thy. 5.2]. Here, for constant control ut $\mathscr{K}, L^{u}$ is the generator of the diffusion process in $\mathbb{R}^{N}$ corresponding to (1.1a):

$$
\begin{equation*}
L^{u_{f}}=\frac{1}{2} \sum_{i, j=1}^{N} a_{i j}(x) f_{x_{i} x_{j}}+\left(b^{0}(x)+b^{1}(x) u\right) \cdot \nabla f \tag{3.2}
\end{equation*}
$$

with $a=\sigma \sigma^{\prime}$.
3.4. For every $T<\infty, K=1,2, \ldots$, there exists $\gamma_{K l}$ such that

$$
E_{\pi}| | \Lambda_{t}\left\|^{K} \leq \gamma_{K T}| | \mu\right\|^{K}, \quad 0 \leq t \leq T
$$

for all $\pi \in \mathscr{A}^{\prime}$. See $[3$, Thy. 5.3] with $m=0$.
For $t \geq 0, \mu \in \mathscr{M}, \pi \in \mathscr{A}, \phi \in C(M)$ let

$$
\begin{equation*}
J(t, \mu, \pi, \phi)=E_{\pi} \phi\left(\Lambda_{t \mu}^{Y U}\right) \tag{3.3}
\end{equation*}
$$

Since $\phi \in C_{K}(\mathscr{M})$ for some $K$, the expectation exists by 3.1 and 3.4 .

Lemma 3.5. Let $\left\|\phi_{n}\right\|_{K} \leq i$ and $d\left(\phi_{n}, \phi\right) \rightarrow 0$ as $n \rightarrow \infty$. Then

$$
J(t, \mu, \pi, \phi)=\lim _{n \rightarrow \infty} J\left(t, \mu, \pi, \phi_{n}\right)
$$

uniformly on $[0, T] \times M_{r} \times \mathscr{A}$, for any finite $T, r$.

Proof. Consider. $\Gamma \subset \Omega$, and let $\Gamma_{T} \subset{ }^{S_{T}}$ denote the set of restrictions to $[0, T]$ of $(Y, U) \in \Gamma$. Then
(*)

$$
\begin{aligned}
\mid E_{\pi} \phi_{n}\left(\Lambda_{t}\right) & -E_{\pi} \phi\left(\Lambda_{t}\right) \mid \leq \\
& \leq \int_{\Gamma}\left|\phi_{n}\left(\Lambda_{t}\right)-\phi\left(\Lambda_{t}\right)\right| d \pi+\int_{\Gamma}\left|\phi_{n}\left(\Lambda_{t}\right)-\phi\left(\Lambda_{t}\right)\right| d \pi
\end{aligned}
$$

with $\Gamma^{\prime}=\Omega-\Gamma$. If $\Gamma_{T}$ is a compact subset of $\Omega_{T}$, then $\| Y_{T}$ is bounded on $\Gamma$. By $3.2,0 \leq t \leq T,(Y, U) \in \Gamma, \mu \in \mu_{r}$ imply $\Lambda_{t} \in M_{\rho}$ for some $\rho$. Since $d\left(\phi_{n}, \phi\right) \rightarrow 0, \phi_{n} \rightarrow \phi$ uniformly on $\mu_{\rho}$. Therefore, the first term on the right side of (*) tends to 0 as $n \rightarrow \infty$, uniformly with respect to $(t, \mu, \pi) \in[0, T] \times \mu_{r} \times \mathscr{A}$.

It remains to show that, given $\varepsilon>0, \Gamma$ can be chosen such that the last term in (*) is less than $\varepsilon$, uniformly on $[0, T] \times$ $\boldsymbol{M}_{\mathrm{r}} \times$. Now

By Cauchy-Schwartz and 3.4

$$
\begin{aligned}
\int_{\Gamma^{\prime}}\left\|\Lambda_{\mathrm{t}}\right\|^{\mathrm{K}_{\mathrm{d} \pi}} & \leq \pi\left(\Gamma^{\prime}\right)^{\frac{1}{2}}\left(\int_{\Gamma^{\prime}}\left\|\Lambda_{\mathrm{t}}\right\|^{2 K_{d \pi}}\right)^{\frac{1}{2}} \\
& \leq \pi\left(\Gamma^{\prime}\right)^{\frac{1}{2}} \gamma_{2 K, T}^{\frac{1}{2}}\|\mu\|^{K}
\end{aligned}
$$

Under $(Y, U) \rightarrow Y$, $\pi$ projects onto Wiener measure w. Let $A \subset C\left([0, T] ; \mathbb{R}^{N}\right)$ be compact with $Y_{0}=0$ for all $Y \in A$ and

$$
w\left[C\left([0, T] ; \mathbb{R}^{N}\right)-A\right]<\varepsilon^{2}\left[2 C\left(1+\gamma_{2 K, T^{2}}^{\frac{1}{2}} r^{K}\right)\right]^{-2}
$$

We choose $\Gamma$ such that $\Gamma_{T}=A \times L^{2}([0, T] ; \mathscr{Z})$. Since $\mathrm{L}^{2}([0, \mathrm{~T}] ; \mathscr{U})$ is compact (weak topology), $\Gamma_{\mathrm{T}}$ is compact. We have

$$
\int_{\Gamma^{\prime}}\left|\phi_{n}\left(\Lambda_{t}\right)-\phi\left(\Lambda_{t}\right)\right| \mathrm{d}^{\pi} \leq 2 \mathrm{C}\left(\pi\left(1^{\prime}\right)+\pi\left(\Gamma^{\prime}\right)^{\frac{1}{2}} \gamma_{2 \mathrm{~K}, \mathrm{~T}^{\frac{1}{2}}}^{\mathrm{K}}\right)<\varepsilon
$$

as required. This proves Lemma 3.5.

Lemma 3.6. For each $t \geq 0, \phi \in C(M), r<\infty, J(t, \mu, \pi, \phi)$ is continuous on $M_{r} \times \mathscr{A}$.

Proof. Let $g(\mu, Y, U)=\phi\left(\Lambda_{t}^{Y}, U_{\mu}\right) . \quad B y 3.1,3.2, g$ is continuous on $\mathscr{M}_{r} \times \Omega$ (recall that $\left(Y_{n}, U_{n}\right) \rightarrow(Y, U)$ implies $\left\|Y_{n}-Y\right\|_{t} \rightarrow 0$, and hence $\left\|Y_{n}\right\|_{t} \leq a$ for some a.) Moreover, $g(\mu, \cdot, \cdot)$ is $\mathscr{G}_{\mathrm{t}}$-measurable.

Suppose first that $\phi(\mu)$ is bounded on $\mu$. Let $\mu_{n} \rightarrow \mu$, $\pi_{n} \rightarrow \pi$ with $\mu_{n} \in M_{r}$. By definition of weak convergence

$$
\lim _{n \rightarrow \infty} \int_{\Omega} g(\mu, Y, U) d \pi_{n}=\int_{\Omega} g(\mu, Y, U) d \pi
$$

Moreover, $\left|g\left(\mu_{n}, Y, U\right)-g(\mu, Y, U)\right| \rightarrow 0$ as $n \rightarrow \infty$, uniformly on any $\Gamma \subset \Omega$ such that the set $\Gamma_{t}$ of restrictions to $[0, t]$ of $(\mathrm{Y}, \mathrm{U}) \in \Gamma$ is compact. As in the proof of Lemna 3.5, we can choose $\Gamma$ such that $\pi_{n}(\Omega-\Gamma)$ is arbitrarily small, uniformly with respect to $n$. This proves Lemma 3.6 in case $\phi(\mu)$ is bounded on $\mathcal{M}$.

Now take any $\phi \in C_{K}(\mathbb{M})$. By Lemmas 2.2 and 2.3 , there exist $\psi_{\mathrm{n}} \in \mathrm{C}_{\mathrm{K}}(\mathscr{M})$ such that $\left|\phi_{\mathrm{n}}(\mu)\right|$ is bounded on $\mathscr{M}$ for cach $n$, $\left\|\phi_{n}\right\|_{K}$ is bounded, and $d\left(\phi_{n}, \phi\right) \rightarrow 0$ as $n \rightarrow \infty$. Lemma 3.6 now follows from Lemma 3.5 .

The control problem. Given $t, \mu, \phi$, we consider the problem of minimizing $J(t, \mu, \pi, \mu)=E_{\pi} \phi\left(\Lambda_{t}\right)$ on the space of admissible controls $\pi$. We can regard the Zakai equation (3.1) as governing the dynamics of the "state" process $\Lambda_{t}$ for this control problem. Since $\Lambda_{t}$ is an unnormalized conditional distribution measure for $X_{t}$ in the partially-observed control system (1.1), we call the problem of minimizing $E_{\pi} \phi\left(\Lambda_{t}\right)$ a "separated" optimal control problem.

Following Nisio [6] let

$$
\begin{equation*}
\mathscr{\mathscr { T }}_{\mathrm{t}}^{\phi}(\mu)=\min _{\pi \in \mathscr{A}} J(\mathrm{t}, \mu, \pi, \phi) \tag{3.4}
\end{equation*}
$$

The minimum is attained, by Lemma 3.6.

Since $\phi\left(\Lambda_{t}\right)$ is $\mathscr{G}_{t}$-measurable, the minimum is the same taken in the class $\alpha_{t}$ of admissible controls on $[0, t]:$

$$
\begin{equation*}
y_{t}^{\phi(\mu)}=\min _{\pi_{t} \in \mathscr{A}_{t}} J\left(t, \mu, \pi_{t}, \phi\right) \tag{3.5}
\end{equation*}
$$

For the special case $\phi(\mu)=\langle G, \mu\rangle, J=E_{\pi}\left\langle G, \Lambda_{t}\right\rangle$ which is of the form considered in the existence theorem [3, Theorem 4.1]. However, if $\varphi$ has this special linear form, $\dot{y}_{t} \Phi(\mu)$ is not not linear in $\mu$. Hence, we define $\mathscr{S}_{t}$ the bigger space $C(\mathscr{K})$, and not merely on the space of $\phi$ of the form $\phi(\mu)=\langle G, \mu\rangle$.

Theorem 3.1. $\phi \in \mathrm{C}_{K}(\mathcal{M})$ implies $\mathscr{y}_{\mathrm{t}} \boldsymbol{y}^{\phi} \in \mathrm{C}_{K}(\mathcal{M})$.

Proof. By Lemma 3.6 and the fact that $\mathscr{M}_{\mathrm{r}}$ and $\mathscr{A}$ are compact, $\left\|\mu_{n}\right\| \leq r$ and $\mu_{n} \rightarrow \mu$ imply $\dddot{y}_{t}^{\phi}\left(\mu_{n}\right) \rightarrow \mathscr{y}_{t}^{\phi}(\mu)$. Since any $w^{*}$-convergent sequence $\mu_{n}$ has $\left\|\mu_{n}\right\|$ bounded, $y_{t}^{\phi}$ is continuous on $\boldsymbol{M}$. From 3.4,

$$
\begin{aligned}
|J(t, \mu, \pi, \phi)| & \leq\|\phi \mid\|_{K} E \int_{\Omega}\left(1+| | \Lambda_{t} \|^{K}\right) \mathrm{d} \pi \\
& \leq\|\phi \mid\|_{K}\left(1+\gamma_{K t}| | \mu \|^{K}\right) \\
& \leq\left(1+\gamma_{K t}\right)\|\phi\|_{K}\left(1+| | \mu \|^{K}\right) .
\end{aligned}
$$

Thus, $\left\|\psi_{t} \phi\right\|_{K} \leq\left(1+\gamma_{K t}\right)\|\phi\|_{K}$, which proves Theorem 3.1.
In the next section we establish the semigroup property of \%
4. The Semigroup Property. The purpose of this section is to prove the following two theorems.

Theorem 4.1. For every $\phi \in C(\mathbb{K}), \mathrm{s}, \mathrm{t} \geq 0$,

$$
y_{s}+t^{\phi}=\mathscr{y}_{s} y_{t}^{\phi} .
$$

Theorems 3.1 and 4.1 imply that $\%_{t}$ is a (nonlinear) semigroup on $C(\mathscr{M})$. Let $C_{b}(\mathscr{M})$ denote the space of bounded continuous functions on $\mathscr{M}$ (it is the same as $C_{K}(\mathscr{M})$ when $K=0$. ) From (3.3), (3.4) $\left\|\mathscr{T}_{t}^{\phi}-\mathscr{T}_{t} \psi\right\|_{0} \leq\|\phi-\psi\|_{0}$. Hence, when restricted to $\mathrm{C}_{\mathrm{b}}(\mathscr{M}), \mathscr{T}_{\mathrm{t}}$ is a contracting semigroup on $\mathrm{C}_{\mathrm{b}}(\mathbb{M})$.

Theorem 4.2. For every $\phi \in C(\mathbb{M}), \mathscr{T}_{\mathrm{t}}{ }^{\phi}$ is a continuous function of $t \in[0, \infty)$ in the d-metric on $C(\mathbb{M})$.

The proof of Theorem 4.1 will be based on a series of three lemmas. We begin by temporarily imposing rather stringent conditions on the coefficients in (1.1), and on $\mathrm{Y}, \mathrm{U}, \mu$. We say that the coefficients are regular if $\sigma, b^{0}, b^{1}, g$ are of class $C_{b}^{\infty}\left(\mathbb{R}^{N} ; \mathbb{R}^{\ell}\right)$ for the appropriate $\ell=N D, N, N L, M$, respectively. Let us denote by $C_{e}^{1,2}$ the class of functions $q$ on $[0, \infty) \times \mathbb{R}^{N}$ with the following properties:
(i) $q$ and the partial derivatives $q_{t}, q_{x_{i}} ; q_{x_{i}} x_{j}$ are continuous, $\mathrm{i}, \mathrm{j}=1, \ldots, \mathrm{~N}$.
(ii) For each $T>0$, there exist $C, k>0$ (depending perhaps on $T$ ) such that

$$
|r(x, t)| \leq C \exp (-k|x|), 0 \leq t \leq T,
$$

where $r$ denotes any of the functions $q, q_{x_{i}}, q_{x_{i}} x_{j}$.
For brevity, we write $q(t)=q(t, \cdot)$.

Lemma 4.1. Assume that the coefficients in (1.1) are regular, and that $Y \in C^{1}\left([0, \infty) ; \mathbb{R}^{N}\right), U \in C([0, \infty) ; \mathscr{K})$. Then:
(a) If $\mu$ has a density $p_{0} \in C_{0}^{\infty}\left(\mathbb{R}^{N}\right)$, then $\Lambda_{t}\left(=\Lambda_{t \mu}^{Y, U}\right)$ has a density $q \in C_{e}^{1,2}$, satisfying the partial differential equation

$$
\begin{align*}
\frac{d q}{d t} & =\left(L^{U}\right)^{*} q+h q \cdot \dot{Y}_{t}-\frac{1}{2}|h|^{2} q, \quad t \geq 0  \tag{4.1}\\
q(0) & =p_{0} .
\end{align*}
$$

(b) If $q \in C_{e}^{1,2}$ is a solution of (4.1) with $q(0)$ the density of $\mu$, then $q(t)$ is the density of $\Lambda_{t}$ for all $t \geq 0$.

Here ( $L^{u}$ ) * denotes the formal adjoint of the operator $L^{u}$ in (3.2), and $\dot{Y}_{t}=d Y / d t$. Note that part (a) of the Lemma, but not part (b), requires that $q(0)$ has compact support.

Proof of Lemma 4.1. To prove (a), we recall from [3, §5] that

$$
\begin{equation*}
p(t)=q(t) \exp \left(-Y_{t} \cdot h\right) \tag{4.2}
\end{equation*}
$$

is a solution in $C_{e}^{1,2}$ to the partial differential equation

$$
\begin{align*}
\frac{d p}{d t} & =V_{t}^{*} p+c(t) p, \quad \text { where }  \tag{4.3}\\
V_{t} & =L_{t}-\left(a Y_{t} \cdot \nabla h, \nabla\right), L_{t}=L_{t} \\
e(t) & =\frac{1}{2}\left(a Y_{t} \cdot \nabla h, Y_{t} \cdot \nabla h\right)-Y_{t} \cdot L_{t} h-\frac{1}{2}|h|^{2}
\end{align*}
$$

where $(a \xi, \eta)=\sum_{\substack{i, j=1 \\, N}}^{N}{ }_{i j} \xi_{i} \eta_{j}$ and denotes the product in $\mathbb{R}^{M}$. The operators $L_{t}^{*}, Y_{t}^{*}$ are related by

$$
\begin{equation*}
\left(\stackrel{V}{L}_{t}^{*} p\right) \exp \left(Y_{t} \cdot h\right)=L_{t}^{*} q-e q-\frac{1}{2}|h|^{2} q . \tag{4.4}
\end{equation*}
$$

Equation (4.4) follows upon multiplying both sides of (4.4) by $f \in C_{0}\left(\mathbb{R}^{N}\right)$, integrating by parts, and using the relation

$$
\begin{aligned}
\exp \left(Y_{t} \cdot h\right) L_{t} f & =\check{L}_{t}^{v}\left[f \exp \left(Y_{t} \cdot h\right)\right]+e(t) f \exp \left(Y_{t} \cdot h\right) \\
& +\frac{1}{2}|h|^{2} f \exp \left(Y_{t} \cdot h\right) .
\end{aligned}
$$

Then equation (4.1) follows from (4.3), (4.4) and the product rule applied to $\frac{d}{d t}\left[p \exp \left(Y_{t} \cdot h\right)\right]$.

To prove (b), if $a \in C_{e}^{1,2}$ satisfies (4.1), then the above calculation shows that $p(t)$ defined by (4.2) is a solution in $C_{e}^{1,2}$ to (4.3). It follows from [3, (5.5)] that $q(t)$ is the density of $\Lambda_{t}$. (In the derivation of $[3,(5.5)]$ it was stated that $q(0) \in C_{0}\left(\mathbb{R}^{N}\right)$. However, the proof there is based on integrations by parts, and is the same if $q \in C_{e}^{1,2}$.) This proves Lemma 4.1.

For $s \geq 0$, let us introduce the notation

$$
Y_{\tau}^{S}=Y_{S+\tau}-Y_{S}, \quad U_{\tau}^{S}=U_{S+\tau}, \quad \tau \geq 0
$$

In particular, $Y_{0}^{S}=0 ;$ and $(Y, U) \in \Omega$ implies $\left(Y^{S}, U^{S}\right) \in \Omega$.

Lemma 4.2. For every $(Y, U) \in \Omega, \mu \in \mathscr{M}, \mathrm{s}, \mathrm{t} \geq 0$,

$$
\begin{equation*}
\Lambda_{\mathrm{S}+\mathrm{t}, \mu}^{\mathrm{YU}}=\Lambda_{\mathrm{t}}^{\mathrm{Y} \Lambda_{\mathrm{S}}}{ }^{\mathrm{S}} \text {, where } \Lambda_{\mathrm{s}}=\Lambda_{\mathrm{s} \mu}^{\mathrm{YU}} \tag{4.5}
\end{equation*}
$$

Proof. Step 1. First assume the conditions of Lemma 4.1 on $b^{\ell}, \sigma, h, Y, U$, and that $\mu=\Lambda_{0}$ has a density $p_{0} \in C_{0}^{\infty}\left(\mathbb{R}^{N}\right)$. By Lemma $4.1(a), \Lambda_{\tau}$ has a density $q(\tau) \in C_{e}^{1,2}$ satisfying (4.1) for $\tau \geq 0$. Let $q^{s}(\tau)=q(s+\tau)$. Then $q^{s}$ is a solution in $C_{e}^{1,2}$ of (4.1), with (Y,U) replaced by $\left(Y^{S}, U^{S}\right)$; note that $\dot{Y}_{s+\tau}=\dot{Y}_{\tau}^{S}$ and $q^{s}(0)=q(s)$. By Lemma $4.1(b), q^{s}(t)$ is the density of $\Lambda_{t \Lambda_{S}}^{S_{U}^{S}}$. This proves (4.5) under these conditions.

Step 2. Again assume regular coefficients $b^{\ell}, \sigma, h, \ell=0,1$.
Let $(Y, U) \in \Omega, \mu \in \mathscr{M}$. Let $\left(Y_{n}, U_{n}\right) \rightarrow(Y, U), \mu_{n} \rightarrow \mu$, where $Y_{n}, U_{n}, \mu_{n}$ satisfy the conditions in Step 1 for each $n$. Write $\Lambda_{\mathrm{s}}^{\mathrm{n}}=\Lambda_{\mathrm{s} \mu_{\mathrm{n}}}^{\mathrm{Y}_{\mathrm{n}}} \mathrm{U}_{\mathrm{n}}$. By property 3.1 , as $\mathrm{n} \rightarrow \infty$

$$
\begin{aligned}
& \Lambda_{\mathrm{s}+\mathrm{t}, \mu_{\mathrm{n}}}^{\mathrm{Y}_{\mathrm{n}} \mathrm{U}_{\mathrm{n}}} \quad \Lambda_{\mathrm{s}+\mathrm{t}, \mu}^{\mathrm{YU}} \\
& \Lambda_{\mathrm{s}}^{\mathrm{n}} \rightarrow \Lambda_{\mathrm{s}}, \\
& \Lambda_{\mathrm{t}}^{\mathrm{Y}} \Lambda_{\mathrm{s}}^{\mathrm{S}} \mathrm{U}_{\mathrm{n}}^{\mathrm{s}} \rightarrow \Lambda_{\mathrm{t}}^{Y^{s} U_{\mathrm{s}}^{\mathrm{s}}}
\end{aligned}
$$

At the last step we used the fact that $\left(Y_{n}^{S}, U_{n}^{s}\right) \rightarrow\left(Y^{s}, U^{s}\right)$. This implies (4.5).
 each $n$, uniformly bounded together with their first order partial derivatives and tending uniformly to $\sigma, b^{\ell}, h$ as $n \rightarrow \infty, \ell=0, l$. Write $\Lambda_{t \mu}^{n}=\Lambda_{t \mu}^{n Y U}$ to indicate that the coefficients depend on $n$. The proof of [3, Theorem 5.1] shows the following: $v_{n} \rightarrow v$, $\nu_{n} \in \mathscr{M}_{r}$, implies $\Lambda_{\tau \nu}^{n} \rightarrow \Lambda_{\tau \nu}$ for any $\tau \geq 0$. We then have as $n \rightarrow \infty$

$$
\Lambda_{\mathrm{s}+\mathrm{t}, \mu}^{\mathrm{n}} \rightarrow \Lambda_{\mathrm{s}+\mathrm{t}, \mu}, \quad \Lambda_{\mathrm{s} \mu}^{\mathrm{n}} \rightarrow \Lambda_{\mathrm{s} \mu} .
$$

Similarly, if we write $\Lambda_{s \mu}^{n}=\Lambda_{s}^{n}$, then

$$
\underset{t \Lambda_{S}^{n}}{\Lambda^{S} U^{S}} \rightarrow \Lambda_{t \Lambda}^{Y^{S} U^{S}}
$$

This implies (4.5), and hence Lemma 4.2.
As in $\S 3$ let $\pi_{s}$ denote the restriction to $\mathscr{G}_{S}$ of $\pi \in \mathbb{Z}$. Let $\pi_{5}{ }_{S}$ be a regular conditional distribution for ( $Y^{s}, U^{s}$ ) given $\mathscr{S}_{S}$.

Lemma 4.3. If $\pi \in \mathscr{A}$, then:
(a) $\pi_{\mathrm{S}}^{\mathrm{YU}} \in \mathscr{A}, \pi_{\mathrm{s}}$-almost surely;
(b) $J(s+t, \mu, \pi, \phi)=\int_{\Omega} J\left(s, \Lambda \frac{Y U}{S \mu}, \pi_{s}^{Y U}, \phi\right) d \pi_{s}$,
for any $\phi \in C(\mathscr{M})$.

Proof. To prove (a) it suffices to verify that, for any $\mathscr{G}_{\mathrm{s}}$-measurable $\Phi \in \mathrm{C}_{\mathrm{b}}(\Omega), \mathscr{S}_{\mathrm{t}}$-measurable $\Psi \in \mathrm{C}_{\mathrm{b}}(\Omega), \mathrm{F} \in \mathrm{C}_{\mathrm{b}}\left(\mathbb{R}^{M}\right)$, and $r>t$

$$
E_{\pi}\left[\Psi(Y, U) \Phi\left(Y^{S}, U^{S}\right) F\left(Y_{r}^{S}-Y_{t}^{S}\right)\right]=E_{\pi}\left[\Psi(Y, U) \Phi\left(Y^{S}, U^{S}\right)\right] E_{\pi} F\left(Y_{r}^{S}-Y_{t}^{S}\right)
$$

But this follows from independence under $\pi$ of the random variables $\Psi(Y, U) \Phi\left(Y^{S-}, U^{S}\right)$ and $F\left(Y_{r}^{S}-Y_{t}^{S}\right)$.

Part (b) is immediate from (3.3), Lemma 4.2 and properties of conditional expectations.

Proof of Theorem 4.1. For every $\pi \in \mathscr{A}$, Lemma 4.3, the definition (3.4) of $\mathscr{S}_{\mathrm{t}}{ }^{\phi}$, and $(3,5)$ imply

$$
\begin{aligned}
\mathrm{J}(\mathrm{~s}+\mathrm{t}, \mu, \pi, \phi) & =\int_{\Omega} \mathrm{J}\left(\mathrm{~s}, \Lambda_{\mathrm{s} \mu}^{\mathrm{YU}}, \pi \mathrm{~s}^{\mathrm{YU}}, \phi\right) \mathrm{d}_{\mathrm{s}} \\
& \geq \int_{\Omega} \mathscr{G}_{\mathrm{t}} \phi\left(\Lambda_{\mathrm{s} \mu}^{\mathrm{YU}}\right) \mathrm{d}_{\mathrm{s}}=\mathrm{E}_{\pi_{\mathrm{s}}} \mathscr{T}_{\mathrm{t}} \phi\left(\Lambda_{\mathrm{s} \mu}^{\mathrm{YU}}\right) \\
& \geq \pi_{s} \mathscr{G}_{\mathrm{t}} \phi(\mu) .
\end{aligned}
$$

Since this is true for every $\pi \in \mathscr{A}$,

$$
\mathscr{I}_{s}+t{ }^{\phi(\mu)} \geq \mathscr{S}_{s} \mathscr{G}_{t} \phi(\mu) .
$$

To prove the opposite inequality, we make the following construction. Let $\rho>0, \delta>0$ to be chosen later. Let $A_{0}=M-M_{\rho}$ and $A_{1}, \ldots, A_{m}$ disjoint Borel subsets of $\mathcal{N}_{\rho}$,
such that

$$
\mathscr{\mu}_{\rho}=A_{1} \cup \ldots \cup A_{m}
$$

and for $\nu, \nu^{\prime} \in A_{i}, i=1, \ldots, m, \pi \in \mathbb{A}$,

$$
\left|J(t, v, \pi, \phi)-J\left(t, v^{\prime}, \pi, \phi\right)\right|<\delta .
$$

This is possible by Lemma 3.6. Choose $\mu_{i} \in A_{i}$ and $\pi_{i} \in \mathscr{A}$ such that

$$
J\left(t, \mu_{i}, \pi_{i}, \phi\right)<\mathscr{S}_{t} \phi\left(\mu_{i}\right)+\delta .
$$

For all $v \in A_{i}$,
(*)

$$
J\left(t, v, \pi_{i}, \phi\right)<\mathscr{F}_{t} \Phi(v)+3 \delta .
$$

Let $\pi_{0} \in \mathscr{A}$ be arbitrary. Let

$$
\pi_{S}^{Y U}=\pi_{i} \quad \text { if } \quad \Lambda_{s \mu}^{Y U} \in A_{i}
$$

Given $\pi_{s} \epsilon \mathscr{A}_{s}$, this defines $\pi \in \mathscr{A}$ such that $\pi_{S}^{Y U}$ is a regular conditional distribution for $Y^{s}, U^{s}$ given $\mathscr{G}_{S}$ and $\pi \mid \mathscr{G}_{S}=\pi_{S}$. By Lemma 4.3 and (*), with $v=\Lambda_{s \mu}^{Y U}=\Lambda_{s}$

$$
\begin{aligned}
J(s+t, \mu, \pi, \phi) & =\int_{\Omega} J\left(s, \Lambda_{s}, \pi_{s}^{Y U} \phi\right) d \pi_{s} \\
& \leq \int_{\mathscr{H}_{\rho}} \mathscr{F}_{t} \phi\left(\Lambda_{s}\right) d \pi_{s}+\int_{\Lambda_{0}} J\left(s, \Lambda_{s}, \pi_{0}, \phi\right) d \pi_{s}+3 \delta .
\end{aligned}
$$

Since

$$
\mathscr{\mathscr { M }}_{s+t} \phi(\mu) \leq J(s+t, \mu, \pi, \phi), \text { we have }
$$

$$
\begin{aligned}
\mathscr{S}_{s}+t^{\phi}(\mu) & \leq E_{\pi_{s}} \mathscr{C}_{t} \phi\left(\Lambda_{s}\right)+\int_{A_{0}} J\left(s, \Lambda_{s}, \pi_{0}, \phi\right) d \pi_{s} \\
& +\int_{A_{0}}\left|\mathscr{C}_{t} \phi\left(\Lambda_{s}\right)\right| d \pi_{s}+3 \delta .
\end{aligned}
$$

Now $\phi \in C_{K}(\mathbb{M})$. for some $K$. We have, for some $C_{1}$,

$$
\begin{aligned}
& \left|J\left(s, \Lambda_{s}, \pi_{0}, \phi\right)\right| \leq C_{1}\left(1+\left\|\Lambda_{s}\right\|^{K}\right) \\
& \left|\mathscr{T}_{t}{ }^{\phi}\left(\Lambda_{s}\right)\right| \leq C_{1}\left(1+\left\|\Lambda_{s}\right\|^{K}\right)
\end{aligned}
$$

while by 3.4 and the fact that $A_{0}=\{\nu:\|\nu\|>\rho\}$

$$
\begin{aligned}
\int_{A_{0}}\left(1+\left\|\Lambda_{s}\right\|^{K}\right) d \pi_{s} & \leq \rho^{-K_{E_{\pi_{s}}}\left[\left\|\Lambda_{s}\right\|^{K}+\left\|\Lambda_{s}\right\|^{2 K}\right]} \\
& \leq C_{2} \rho^{-K}\left(1+r^{2 K}\right)
\end{aligned}
$$

for $\mu \in M_{r}$. Therefore, given $\varepsilon>0$ we can choose $\rho$ large enough and $\delta$ small enough that

$$
\mathscr{G}_{s+t^{\phi}(\mu) \leq}^{E_{\pi_{s}} \mathscr{T}_{t}^{\phi}\left(\Lambda_{s}\right)+\varepsilon, ~}
$$

for all $\mu \in \mathscr{M}_{r}$ and $\pi_{s} \in \mathscr{A}_{s}$. Upon taking the inf over $\pi_{s}$ (recall (3.5))

$$
\mathscr{T}_{S}+t^{\phi(\mu)} \leq \mathscr{F}_{\mathrm{t}} \mathscr{T}^{\phi(\mu)+\varepsilon .}
$$

Since $\varepsilon$ is arbitrary, we obtain Theorem 4.1.
In preparation for the proof of Theorem 4.2, and for 55 , let us introduce the following family of operators $\mathscr{L}^{\mathrm{u}}$, for constant controls $u \in \mathscr{U}$. Let

$$
\begin{aligned}
\tilde{\mathscr{D}}= & \left\{\phi: \phi(\mu)=F\left(\left\langlef_{1}, \mu>, \ldots,\left\langle f_{J}, \mu>\right),\right.\right.\right. \\
& \left.F \in C^{2}\left(\mathbb{R}^{J}\right), f_{1}, \ldots, f_{J} \in C_{0}^{2}\left(\mathbb{R}^{J}\right), J=1,2, \ldots\right\},
\end{aligned}
$$

and for each integer $m \geq 0$

$$
\begin{gather*}
\mathscr{D}_{\mathrm{m}}=\left\{\phi \in \tilde{\mathscr{D}}:\left|\mathrm{F}_{z_{j}}(z)\right| \leq \mathrm{C}\left(1+|z|^{\mathrm{m}+1}\right),\left|\mathrm{F}_{z_{j} z_{k}}(z)\right|\right.  \tag{4.6}\\
\left.\leq \mathrm{C}\left(1+|z|^{m}\right), \quad j, k=1, \ldots, \mathrm{~J}\right\} .
\end{gather*}
$$

We have the inclusions $\mathscr{D} \subset \mathscr{D}_{\mathrm{m}} \subset \mathrm{C}_{\mathrm{m}+2}(\mathscr{M})$.
For $\phi \in \mathscr{D}_{\mathrm{m}}$ and $u \in \mathscr{U}$, let

$$
\begin{align*}
\mathscr{L}^{u}{ }_{\phi(\mu)} & =\sum_{j=1}^{J} F_{z_{j}}(\ldots)<L{ }^{u_{f}}, \mu>  \tag{4.7}\\
& +\sum_{j, k=1}^{J} F_{z_{j}} z_{k}(\ldots)<h f_{j}, \mu>\cdot\left\langle h f_{k}, \mu>\right.
\end{align*}
$$

where ... denotes that the partial derivatives $F_{z_{j}}, F_{z_{j}} z_{k}$ are evaluated at the vector $z=\left(\left\langle f_{1}, \mu\right\rangle, \ldots,\left\langle f_{J}, \mu\right\rangle\right)$. It might seem that $\mathscr{L}^{u_{\phi}}$ depends not just on $\phi$, but also on $F, f_{1}, \ldots, f_{J}$.

However, it follows from (4.13) below that this difficulty does not occur.

Lemma 4.4. Let $\phi \in \mathscr{D}_{\mathrm{m}}$. Then there exists c such that:
(a) $\mathscr{L}^{\mathrm{u}_{\Phi} \in \mathrm{C}_{\mathrm{m}+2}(\mathscr{M}),\left\|\mathscr{L}^{\mathrm{u}_{\Phi}}\right\|_{\mathrm{m}+2} \leq \mathrm{c} \text { for all} u \in \mathscr{U} . ~}$
(b) The mapping $(u, \mu) \rightarrow \mathscr{L}^{u_{\phi}(\mu)}$ is continuous from $\mathscr{K} \times \mathcal{M}_{r}$ into $\mathbb{R}^{1}$ for every $r<\infty$.

This follows at once from (4.7).
Let us next apply the It $\hat{o}$ differential rule to $\phi\left(\Lambda_{t}\right)$; for $\varphi \in \mathscr{D}_{\mathrm{m}}$,

$$
\varphi\left(\Lambda_{t}\right)=F\left(<f_{1}, \Lambda_{t}>, \ldots,<f_{J}, \Lambda_{t}>\right)
$$

We get, using the zakai equation (3.1),

$$
\begin{equation*}
d \phi\left(\Lambda_{t}\right)=\mathscr{L}^{U} t_{\phi\left(\Lambda_{t}\right) d t}+\sum_{j=1}^{J} F_{z_{j}}(\ldots)<h f_{j}, \Lambda_{t}>\cdot d Y_{t}, \tag{4.8}
\end{equation*}
$$

where.. denotes $\left(\left\langle f_{1}, \Lambda_{t}\right\rangle, \ldots,\left\langle f_{J}, \Lambda_{t}\right\rangle\right)$. Since $\left|F_{z_{j}}\right| \leq C\left(1+|z|^{m+1}\right)$, the components of $i_{z_{j}}\left(\left\langle f_{1}, \mu\right\rangle, \ldots,\left\langle f_{J}, \mu\right\rangle\right)\left\langle h f_{j}, \mu\right\rangle$ are in $C_{m+2}(\mathcal{M})$. From 3.4, the integral on $[0, t]$ of the last term in (4.8) is a square integrable $\pi,\left\{\mathscr{S}_{t}\right\}$ martingale for any $\pi \in \mathscr{A}$. By taking $\mathrm{E}_{\pi} \int_{0}^{\mathrm{t}}$ in (4.8) and using Lemma 4.4(a) we get

$$
\begin{equation*}
\left.E_{\pi^{\phi}\left(\Lambda_{t}\right.}\right)=\phi(\mu)+E_{\pi} \int_{0}^{t} \mathscr{L}^{U_{\theta}} \phi\left(\Lambda_{\theta}\right) d \theta \tag{4.9}
\end{equation*}
$$

for any $\phi \in \mathscr{O}_{\mathrm{m}}, \| \in \mathscr{A}$, and any initial data $\mu=\Lambda_{0}$.

Lemma 4.5. Let $\psi \in \mathscr{D}_{\mathrm{m}}, 0 \leq \mathrm{s} \leq \mathrm{t} \leq \mathrm{T}$. Then there exists $\alpha$ (depending on $\psi$ and $T$ ) such that

$$
\left|\left|\mathscr{I}_{\mathrm{t}} \psi-\mathscr{S}_{\mathrm{s}} \psi\right|_{\mathrm{m}+2} \leq \alpha(\mathrm{t}-\mathrm{s}) .\right.
$$

Proof. Consider any $\pi \in \mathscr{A}$. By (4.9)

$$
\begin{aligned}
\mid E_{\pi} \phi\left(\Lambda_{t}\right) & -E_{\pi} \phi\left(\Lambda_{s}\right)\left|\leq E \int_{s}^{t}\right| \mathscr{L}^{U_{\theta}} \phi\left(\Lambda_{\theta}\right) \mid d^{\theta} \\
& \leq \max _{u \in \mathscr{K}}| | \mathscr{L}^{u_{\phi}}| |_{m+2} \int_{s}^{t}\left(1+E| | \Lambda_{\theta}| |^{m+2}\right) d \theta .
\end{aligned}
$$

By Lemma 4.4 a and 3.4

$$
\left|E_{\pi} \phi\left(\Lambda_{t}\right)-E_{\pi} \phi\left(\Lambda_{s}\right)\right| \leq c\left(1+\gamma_{m+2, T}\right)(t-s)\left(1+||\mu||^{m+2}\right)
$$

Since this holds for all $\pi \in \mathscr{A}$, we get Lemma 4.5 with $\alpha=c\left(1+\gamma_{m+2, T}\right)$.

Proof of Theorem 4.2. For some $K, \phi \in C_{K}(\mathbb{M})$. By Lemma 2.4 there exists $\psi_{n} \in \mathscr{D}, \mathrm{n}=1,2, \ldots$, such that $\left\|\psi_{n} \mid\right\|_{K}$ is bounded and $\mathrm{d}\left(\psi_{n}, \phi\right) \rightarrow 0$. Fix $T>0$. For $0 \leq s<t \leq T$, we write

$$
\mathscr{F}_{\mathbf{t}} \phi-\mathscr{T}_{\mathbf{s}} \phi=\left[\mathscr{J}_{\mathbf{t}} \phi-\mathscr{T}_{\mathbf{t}} \psi_{\mathrm{n}}\right]+\left[\mathscr{T}_{\mathbf{t}} \psi_{\mathbf{n}}-\mathscr{F}_{\mathbf{s}} \psi_{\mathrm{n}}\right]+\left[\mathscr{T}_{\mathbf{s}} \psi_{\mathrm{n}}-\mathscr{S}_{\mathbf{s}} \psi\right]
$$

Lemma 3.5 implies that the first and third terms on the right side
tend to 0 as $n \rightarrow \infty$, uniformly for $0 \leq s<t \leq T$ and $\mu \in \mathscr{M}_{r}$. Lemma 4.5 with $m=0$ implies, for $\mu \in \mathscr{H}_{r}$,

$$
\left|\mathscr{G}_{t} \psi_{n}(\mu)-\mathscr{S}_{s} \psi_{n}(\mu)\right| \leq \alpha_{n}\left(1+r^{2}\right)(t-s)
$$

where $\alpha_{n}$ is some constant. Let

$$
\eta(\varepsilon, r)=\sup \left\{\left|\mathscr{T}_{t}^{\Phi}(\mu)-\mathscr{G}_{\mathrm{s}} \Phi(\mu)\right|: \mu \in \mu_{\mathrm{r}}, 0 \leq \mathrm{s}<\mathrm{t} \leq \mathrm{T}, \mathrm{t}-\mathrm{s}<\varepsilon\right\} .
$$

For each $r, \eta(\varepsilon, r) \rightarrow 0$ as $\varepsilon \rightarrow 0$. This implies $d\left(\mathscr{G}_{t} \phi, \mathscr{V}_{s} \phi\right) \rightarrow 0$ as $t-s \rightarrow 0$, as required.

This proves Theorem 4.2.

Constant controls. In particular, let us consider a constant control $u$. In our formulation, this corresponds to taking $\pi=\pi^{u}=w \times \delta_{u}$, where $w$ is Wiener measure on $C\left([0, \infty) ; \mathbb{R}^{m}\right)$ and $\delta_{u}$ is the Dirac measure on $L_{1 o c}^{2}([0, \infty) ; \mathscr{W})$ concentrated on the constant trajectory $U_{t} \equiv u$. We can then write $E\left(=E_{w}\right)$ instead of $E_{\pi} u$, and obtain from

$$
\begin{equation*}
E \phi\left(\Lambda_{t}\right)=\phi(\mu)+E \int_{0}^{t} \mathscr{L}^{u} \phi\left(\Lambda_{\theta}\right) d \theta, \quad \phi \in \mathscr{D}_{m} \tag{4.10}
\end{equation*}
$$

For constant $u$, we may regard $\Lambda_{t}$ as defined on the sample space $C\left([0, \infty) ; \mathbb{R}^{N}\right)$ of $Y$-trajectories, endowed with the family $\left\{\mathscr{F}_{t}(Y)\right\}$ of $\sigma$-algebras and with Wiener measure $w$. It follows from Lemma 4.2 that $\Lambda_{t}=\Lambda_{t \mu}^{Y u}$ is a Markov process (u fixed), with which is associated the linear semigroup $\mathscr{T}_{\mathrm{t}}^{\mathrm{u}}$ on $\mathrm{C}(\mathbb{M})$ :

$$
\begin{equation*}
y_{t}^{u_{t}} u_{t}(\mu)=E\left(\Lambda_{t}\right), \tag{4.11}
\end{equation*}
$$

where $E=E_{w}$.
From (4.10) we have, for $\phi \in \mathscr{D}_{\mathrm{m}}$,
 $\mathscr{F}_{\theta}^{\mathrm{u}}\left(\mathscr{L}^{\mathrm{u}_{\phi}}\right) \rightarrow \mathscr{L}^{\mathrm{u}_{\phi}}$ as $\theta \rightarrow 0^{+}$, uniformly on $\mu_{\mathrm{r}}$ for each $r<\infty$ (alternatively we could apply Theorem 4.2 with the control space $\mathscr{U}$ replaced by a new one-element control space \{u\}.) Hence the left side of (4.12) tends to 0 as $t \rightarrow 0^{+}$uniformly on $\mu_{r}$, which implies

$$
\begin{equation*}
\mathscr{L}^{u_{\phi}}=\underset{\mathrm{t} \rightarrow 0^{+}}{\mathrm{d}-1 \mathrm{im}} \mathrm{t}^{-1}\left[\mathscr{T}_{\mathrm{t}}^{\left.\mathrm{u}_{\phi-\phi}\right]}, \quad \phi \in \mathscr{U}_{\mathrm{m}}\right. \tag{4.13}
\end{equation*}
$$

This shows that for each $m=0,1,2, \ldots, \mathscr{O}_{m}$ is contained in the domain of the generator of the linear semigroup $\mathscr{T}_{\mathrm{t}}^{\mathrm{u}}$ and that $\mathscr{L}^{\mathrm{u}}$ agrees on $\mathscr{D}_{\mathrm{m}}$ with the generator.
5. The Generator of the Semigroup $\mathscr{G}_{t}$. We define the operator $\mathscr{L}$ on the dense subset $\mathscr{O}$ of $C(\mathbb{M})$ by

$$
\begin{equation*}
\mathscr{L} \phi(\mu)=\min _{u \in \mathscr{U}} \mathscr{L}^{u_{\phi}}(\mu), \quad \Phi \in \mathscr{D} \tag{5.1}
\end{equation*}
$$

Lemma 4.4 implies that $\mathscr{L} \phi \in C_{2}(\mathbb{M})$ for every $\phi \in \mathscr{D}$.

We need slightly stronger hypotheses on $\sigma, b^{0}, b^{1}$ than $\left(A_{1}\right)$, $\left(\mathrm{A}_{2}\right)$ in 53 :
( $A_{i}$ ) Condition ( $A_{1}$ ) holds and, in addition, a $\in C_{b}^{2}\left(\mathbb{R}^{N} ; \mathbb{R}^{N^{2}}\right)$, where $a=\sigma \sigma^{\prime}$.
( $\left.A_{2}^{\prime}\right) b(x, u)=b^{0}(x)+b^{1}(x) u$, where $b^{0} \in C_{b}^{2}\left(\mathbb{R}^{N} ; \mathbb{R}^{N}\right)$ and $b^{1} \in C_{b}^{2}\left(\mathbb{R}^{N} ; \mathbb{R}^{N L}\right)$.

When $\left(A_{j}^{\prime}\right),\left(A_{2}^{\prime}\right),\left(A_{3}\right)$ hold, $f \in C_{0}^{\infty}\left(\mathbb{R}^{N}\right)$ implies $L_{f} f \in C_{0}^{2}\left(\mathbb{R}^{N}\right)$ and $h f \in C_{0}^{2}\left(\mathbb{R}^{N} ; \mathbb{R}^{M}\right)$. From (4.7), $\varnothing \in \mathscr{C}$ implies $\mathscr{f}^{\prime} u_{\phi} \in \mathscr{U}_{2}$.

Theorem 5.1. For every $\phi \in \mathscr{C}$

$$
\begin{equation*}
\mathscr{L} \phi=\underset{t \rightarrow 0^{+}}{\mathrm{d}-\lim ^{-1}\left(\mathscr{G}_{\mathrm{t}}{ }^{\phi-\phi}\right) . . . . . .} \tag{5.2}
\end{equation*}
$$

This theorem justifies our calling $\mathscr{L}$ the generator of the nonlinear semigroup $\mathscr{F}_{t}$. Our proof of Theorem 5.1 follows the same general line of reasoning as Nisio [6].

The proof of Theorem 5.1 depends on the following estimates for the semigroups $\mathscr{J}_{t}^{u}$, for any constant control $u \in \mathscr{K}$. By the same calculation used in the proof of Theorem 3.1

$$
\begin{equation*}
\| \mathcal{y}_{\mathrm{t}}^{-u_{\phi}\left\|_{K} \leq\left(1+\gamma_{K t}\right)\right\| \phi \|_{K}, \quad \phi \in C_{K}(\mathcal{M}) . . . . .} \tag{5.3}
\end{equation*}
$$

For $\phi \in \mathscr{D}_{\mathrm{m}}, 3.4$ and (4.10) imply

$$
\begin{equation*}
\| \mathscr{F}_{\mathrm{t}}^{\mathrm{u}_{\phi-\phi}\left\|_{\mathrm{m}+2} \leq\right\| \mathscr{L}^{u_{\phi}} \|_{\mathrm{m}+2}\left(1+\gamma_{\mathrm{m}+2, \mathrm{t}}\right) \mathrm{t}} \tag{5.4}
\end{equation*}
$$

Lemma $4.4(a)$ gives a bound for $\left\|\mathscr{S}^{\mathrm{u}} \phi\right\|_{\mathrm{m}+2}$. Now consider $\phi \in \mathscr{L}$,

$$
\phi(\mu)=F\left(\left\langle\mathrm{f}_{1}, \mu\right\rangle, \ldots,\left\langle\mathrm{f}_{\mathrm{J}}, \mu\right\rangle\right)
$$

with $F \in C_{b}^{\infty}\left(\mathbb{R}^{J}\right), f_{j} \in C_{0}^{\infty}\left(\mathbb{R}^{N}\right)$. Then

$$
\begin{aligned}
& \mathscr{C}^{u_{\phi}}=-\sum_{j=1}^{J} \phi_{j}+\sum_{j, k=1}^{J} \phi_{j k}, \\
& \phi_{j}(\mu)=F_{z_{j}}\left(\left\langle f_{1}, \mu\right\rangle, \ldots,\left\langle f_{J}, \mu\right\rangle\left\langle L{ }^{u_{f}}, \mu\right\rangle\right. \\
& \Phi_{j k}(\mu)=F_{z_{j} z_{k}}\left(\left\langle f_{1}, \mu\right\rangle, \ldots,\left\langle f_{J}, \mu\right\rangle\right)\left\langle h f_{j}, \mu\right\rangle \cdot\left\langle h f_{k}, \mu\right\rangle \\
& \mathscr{y}_{\mathrm{t}}^{\mathrm{u}_{\phi}}-\phi-\mathrm{t} \mathscr{L}^{\mathrm{u}_{\phi}}=\int_{0}^{\mathrm{t}}\left[\mathscr{F}_{\theta}^{\mathrm{u}}\left(\mathscr{L}^{\mathrm{u}_{\phi}}\right)-\mathscr{L}^{\mathrm{u}_{\phi}}\right] \mathrm{d} \theta \\
& =\sum_{j} \int_{0}^{t}\left[\mathscr{F}_{\theta}^{-u_{\phi}}{ }_{j}-\Phi_{j}\right] d \theta+\sum_{j, k} \int_{0}^{t}\left[\mathscr{C}_{\theta}^{-u_{\phi}}{ }_{j k}-\phi_{j k}\right] d \theta .
\end{aligned}
$$

Since $\phi_{j}, \phi_{j k} \in \mathscr{D}_{2}$, we can apply (5.4) to $\phi_{j}, \phi_{j k}$ to get, for $0 \leq t \leq 1$,

$$
\begin{equation*}
\left\|\mathscr{T}_{t^{p}}^{u_{p}}-\phi-t \mathscr{L}^{u_{\phi}}\right\|_{4} \leq \beta t^{2} \tag{5.5}
\end{equation*}
$$

where the constant $\beta$ depends on $\phi$ but not on $u \in \mathscr{X}$.

Lemma 5.1. For $\phi \in \mathscr{D}$

$$
\mathscr{T}_{t} \phi-\phi \geq \int_{0}^{\mathrm{t}} \mathscr{S}_{\theta}(\mathscr{L} \phi) \mathrm{d} \theta
$$

Proof. By (4.9), for any $u \in \mathscr{A}$

$$
\begin{aligned}
\mathrm{E}_{\pi} \phi\left(\Lambda_{t}\right) & -\phi(\mu)=\mathrm{E}_{\pi} \int_{0}^{\mathrm{t}} \mathscr{L}^{\mathrm{U}_{\theta}} \phi\left(\Lambda_{\theta}\right) \mathrm{d} \theta \\
& \geq \mathrm{E}_{\pi} \int_{0}^{\mathrm{t}} \mathscr{L} \phi\left(\Lambda_{\theta}\right) \mathrm{d}^{\theta}=\int_{0}^{\mathrm{t}} \mathrm{E}_{\pi} \mathscr{L} \phi\left(\Lambda_{\theta}\right) \mathrm{d}^{\theta} \\
& \geq \int_{0}^{\mathrm{t}} \mathscr{F}_{\theta}(\mathscr{L} \phi)(\mu) \mathrm{d} \theta .
\end{aligned}
$$

The minimum over of the left side is $T_{t} \phi(\mu)-\phi(\mu)$. This proves Lemma 5.1.

Proof of Theorem 5.1. Observe that $\mathscr{T}_{t}{ }^{\phi} \leq \mathscr{T}_{t}^{u}{ }^{\mathbf{u}}$ for all $u \in \mathscr{Z}$ (constant controls are suboptimal) Then, for $\phi \in \mathscr{O}$, $0<t \leq 1$,

$$
t^{-1}\left[\mathscr{T}_{t}^{\phi}-\phi-t \mathscr{L}^{u_{\phi}}\right] \leq t^{-1}\left[\mathscr{T}_{t}^{u_{\phi}}-\phi-t \mathscr{L}^{u_{\phi}}\right] .
$$

In particular, given $\mu$ we take $u$ such that $\mathscr{L}^{u} \phi(\mu)=\mathscr{L} \phi(\mu)$ [recall (5.1)]. By (5.5), when $0<t \leq 1$,

$$
t^{-1}\left[\mathscr{T}_{t} \phi(\mu)-\phi(\mu)-t \mathscr{L} \phi(\mu)\right] \leq \beta t\left(1+||\mu||^{4}\right) .
$$

Therefore, uniformly for $\mu \in \mathcal{M}_{r}$,

$$
\lim \sup _{t \rightarrow 0^{+}} t^{-1}\left[\mathscr{T}_{t} \phi(\mu)-\phi(\mu)\right] \leq \mathscr{L} \phi(\mu) .
$$

On the other hand, by Lemma 5.1

$$
\liminf _{t \rightarrow 0^{+}} t^{-1}\left[\mathscr{Y}_{\mathrm{t}} \phi(\mu)-\phi(\mu)\right] \geq \lim \inf _{\mathrm{t} \rightarrow 0} \mathrm{t}^{-1} \int_{0}^{\mathrm{t}} \mathscr{S}_{\mathrm{J}}(\mathscr{L} \phi)(\mu) \mathrm{d} \theta .
$$

Since $\mathscr{L} \phi \in \mathrm{C}(\mathcal{M})$, Theorem 4.2 implies that $\mathscr{T}_{\theta}(\mathscr{L} \phi)(\mu) \rightarrow \mathscr{L} \phi(\mu)$ as $\theta \rightarrow 0^{+}$, uniformly on $\mu_{r}$. Hence,

$$
\underset{t \rightarrow 0^{+}}{\operatorname{im}_{t}}{ }^{-1}\left[\mathscr{T}_{t} \phi(\mu)-\phi(\mu)\right)=\mathscr{L} \phi(\mu)
$$

uniformly on $\mathscr{M}_{r}$, for each $r$. This proves Theorem 5.1.

Remark. The nonlinear semigroup $\mathscr{F}_{\mathrm{t}}$ can be obtained from the family of linear semigroups $\mathscr{T}_{t}^{u}, u \in \mathscr{U}$, by the following procedure used in [6]. For $\Delta>0$, let $\mathscr{J}_{\Delta} \phi(\mu)=\min _{u \in \mathscr{U}} \mathscr{S}_{\Delta}^{u}(\mu)$. For $n=1,2, \ldots$ and dyadic rational $t=m 2^{-n}(m=1,2, \ldots)$ let

$$
\mathscr{T}_{\mathrm{t}}^{\mathrm{n}_{\phi}}=\sigma_{\Delta}^{\mathrm{m}_{\phi}}, \quad \Delta_{\mathrm{n}}=2^{-\mathrm{n}}, \quad \phi \in \mathrm{C}(\mathscr{M})
$$

It is easy to show that, for dyadic rational $t=m 2^{-n}$.

$$
\mathscr{T}_{\mathrm{t}}^{\mathrm{n}_{\phi}} \geq \mathscr{T}_{\mathrm{t}}^{\mathrm{n}+1_{\phi}} \geq \ldots \geq \mathscr{T}_{\mathrm{t}}{ }^{\phi}
$$

By considering controls piecewise constant in time, one can show that $\mathscr{T}_{t}^{-} n_{t} \rightarrow \mathscr{T}^{\phi}$ as $n \rightarrow \infty$, if $t$ is dyadic rational. Choose $n$ large enough such that $t=m 2^{-n}$. Let $\tau_{k}=k 2^{-n}$ and

$$
\mathscr{\Phi}_{n t}=\left\{\pi \in \mathscr{A}_{t}: \pi\left[U_{\tau}=U_{\tau_{k}} \text { for } \tau \in\left[\tau_{k}, \tau_{k+1}\right), \quad k=0,1, \ldots, m-1\right]=1\right\} .
$$

By induction on $m$ (for fixed $n$ ) and a construction like that in
the proof of Theorem 4.1, it can be shown that

$$
\mathscr{y}_{t}^{n_{t}} \phi(\mu)=\min _{\pi \in \mathscr{A}_{n t}} J(t, \mu, \pi, \phi)
$$

By [3, Corollary 6.1], every $\pi \in \mathscr{A}_{\mathrm{t}}$ is the limit of $\pi_{n t}$ as $n \rightarrow \infty$, with $\pi_{n t} \in \mathscr{A}_{n t}$. Lemma 3.6 then implies that $\bar{y}_{t}^{n^{\prime}}(\mu) \rightarrow \mathscr{T}_{t} \phi(\mu)$ as. $\mathrm{n} \rightarrow \infty$.

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## DTIC


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