# A GEL'FAND TRIPLE APPROACH TO THE SMALL NOISE PROBLEM FOR DISCONTINUOUS ODE'S 

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This Version : May 3, 2011


#### Abstract

In this paper, we develop a variational approach to study perturbation problems of ordinary differential equations (ODE's) with discontinuous coefficients. We propose a mathematical framework which can be used to construct stable (and regular) solution processes of discontinuous ODE's.


Key words and phrases: Malliavin calculus, local time, small random perturbations, strong solutions of SDE's.

MSC2010: $60 \mathrm{H} 05,60 \mathrm{G} 44,60 \mathrm{G} 48$.

## 1. Introduction

In this paper, we aim at analyzing the small noise problem of discontinuous ODE's. More precisely, we want to provide conditions under which the solutions $X_{t}^{n}, n \in \mathbb{N}$, of the stochastic differential equations (SDE's)

$$
\begin{equation*}
d X_{t}^{n}=b\left(t, X_{t}^{n}\right) d t+\frac{1}{n} d B_{t}, 0 \leq t \leq 1, \quad X_{0}^{n}=x \in \mathbb{R}^{d} \tag{1.1}
\end{equation*}
$$

for $n \rightarrow \infty$ converge to a solution (process) $X_{t}$ of the ODE

$$
\begin{equation*}
d X_{t}=b\left(t, X_{t}\right) d t, 0 \leq t \leq 1, \quad X_{0}=x \in \mathbb{R}^{d} \tag{1.2}
\end{equation*}
$$

where the drift term $b:[0, T] \times \mathbb{R}^{d} \rightarrow \mathbb{R}^{d}$ is allowed to be a discontinuous function. Here $\left\{B_{t}\right\}_{0 \leq t \leq 1}$ is a $d$-dimensional $\mathcal{F}_{t}$-Brownian motion on a probability space $(\Omega, \mathcal{F}, \mu)$, where $\left\{\mathcal{F}_{t}\right\}_{0 \leq t \leq 1}$ is a $\mu$-augmented filtration generated by $B$..
In the case of continuous drift coefficients $b$ the small noise problem (1.1), 1.2 has been studied by various authors in the literature. See e.g [2, 3, 4, 8, 13, 25] and [26. The author in [25] introduces the large deviation principle to study the convergence rate of solutions of (1.1) to (1.2) with (Lipschitz-) continuous coefficients. We mention that the authors in [2, 3] and [4] employ the Skorohod embedding in combination with certain boundary value problems to establish criteria for the convergence to solutions processes of (1.2). See also [26]. The work [4] deals with a selection principle based on viscosity solutions to construct Feller solutions of ill-posed degenerate diffusion processes. See also the interesting paper of [13] in the context of (stochastic) superposition solutions of ODE's (SDE's). We shall also refer the reader to [1] and the references therein.

The perturbation problem (1.1), 1.2 for discontinuous or even merely measurable drift terms $b$ is in general challenging and sparsely covered by the current literature. See 7, 9, 15, 16]. In the interesting work [7] the authors use the Skorohod embedding technique
to derive (under fairly general conditions on $b$ ) generalized solutions to $(1.2)$ in the sense of Filippov. Further, the papers [15, 16] are concerned with the convergence rate of the probability densities of $X^{n}$ for some (concrete) non-Lipschitzian drift terms $b$. The method used in the latter papers are based on large deviation techniques and viscosity solutions of Hamilton-Jacobi equations. We also emphasize the work [9, where the authors develop large deviations techniques to treat ODE's for certain discontinuous coefficients $b$. Other techniques for the construction of solutions of discontinuous ODE's can be e.g. found in [6, 24].

Our approach to problem (1.1), (1.2) is different from the above mentioned authors' ones and is based on the use of Gel'fand triples

$$
\begin{equation*}
\mathbb{D}_{1,2} \hookrightarrow L^{2}(\mu) \hookrightarrow \mathbb{D}_{-1,2} \tag{1.3}
\end{equation*}
$$

and

$$
\begin{equation*}
(\mathcal{S}) \hookrightarrow L^{2}(\mu) \hookrightarrow(\mathcal{S})^{*} \tag{1.4}
\end{equation*}
$$

$\mathbb{D}_{1,2}$ denotes the stochastic Sobolev space of Malliavin differentiable square integrable Brownian functionals and $\mathbb{D}_{-1,2}$ is its topological dual. Further, $(\mathcal{S})$ is the Hida test function space and $(\mathcal{S})^{*}$ the Hida distribution space. Here the symbol $\hookrightarrow$ stands for continuous inclusions of spaces. We mention that

$$
\begin{equation*}
(\mathcal{S}) \hookrightarrow \mathbb{D}_{1,2} \hookrightarrow L^{2}(\mu) \hookrightarrow \mathbb{D}_{-1,2} \hookrightarrow(\mathcal{S})^{*} \tag{1.5}
\end{equation*}
$$

For more information about Malliavin calculus the reader may consult [11, 18] or [21]. As for the construction of the triple $(\sqrt{1.4})$ and its applications in white noise analysis, we recommend the books of [17] or [22].

To be more precise, our method to tackle the perturbation problem $\sqrt[1.1]{1},(\sqrt{1.2})$ relies on a compactness criterion in $L^{2}(\mu)$ based on Malliavin calculus (see [10]), a "variational calculus" technique with respect to local time [12], and a compactness criterion for continuous functions with values in $(\mathcal{S})^{*}$. Using these tools, we are able to show (under certain stochastic conditions on $b$ ) that $X^{n}$ in 1.1 converges in $L^{2}(\mu)$ (or even in $\mathbb{D}_{1,2}$ ) for a subsequence to a (possibly Malliavin differentiable) cluster point $X_{t}$, which solves the ODE, almost surely (or on a set with positive probability).

We point out that we obtain solutions of discontinuous ODE's which are stable under random perturbations. This approach also provides a natural selection procedure for solutions of discontinuous ODE's which, as one knows, have no unique solutions in general. See e.g [13] for a general discussion of this topic.

## 2. Main Results

In this section, we want to introduce a new technique to study the behavior of the solutions $X^{n}$ of SDE's (1.1) when $n \rightarrow \infty$. Before we proceed, we shall send ahead some notions and definitions which we will make use of later on in this paper.

In the following, let $S([0,1]) \subseteq L^{2}([0,1])$ be the Schwartz space on $[0,1]$ as e.g., constructed in [22]. Using the theorem of Bochner-Minlos, we shall denote by $\pi$ the unique probability measure on the Borel sets $\mathcal{B}\left(S^{\prime}([0,1])\right)$ of $S^{\prime}([0,1])$ (topological dual of $\left.S([0,1])\right)$ such that

$$
\int_{S^{\prime}([0,1])} e^{i\langle\omega, \phi\rangle} \pi(d \omega)=e^{-\frac{1}{2}\|\phi\|_{L^{2}([0,1])}^{2}}
$$

61 for all $\phi \in S([0,1])$, where $\langle\omega, \phi\rangle$ is the action of $\omega \in S^{\prime}([0,1])$ on $\phi \in S([0,1])$.

$$
\begin{equation*}
(\Omega, \mathcal{F}, \mu):=\left(\prod_{i=1}^{d} \Omega_{i}, \otimes_{i=1}^{d} \mathcal{F}_{i}, \otimes_{i=1}^{d} \mu_{i}\right) \tag{2.1}
\end{equation*}
$$

where $\Omega_{i}=S^{\prime}([0,1]), \mathcal{F}_{i}=\mathcal{B}\left(S^{\prime}([0,1])\right), \mu_{i}=\pi$ for $i=1, \ldots, d$.
Further, we briefly recall the definition of the $S$-transform, which can be used to characterize elements of the Hida test function and distribution spaces. See [17]. The $S$-transform of a $\Phi \in(\mathcal{S})^{*}$, denoted by $S(\Phi)$ is defined as

$$
\begin{equation*}
S(\Phi)(\phi)=\langle\Phi, \widetilde{e}(\phi, \cdot)\rangle \tag{2.2}
\end{equation*}
$$

for $\phi \in S_{\mathbb{C}}([0,1])^{d}$, where $S_{\mathbb{C}}([0,1])$ is the complexification of $\left.S([0,1])\right)$ and $\widetilde{e}(\phi, \cdot) \in(\mathcal{S})$ is the exponential functional

$$
\widetilde{e}(\phi, \omega):=\exp \left\{\langle\omega, \phi\rangle-\frac{1}{2}\|\phi\|_{L^{2}\left([0,1] ; \mathbb{R}^{d}\right)}^{2}\right\}
$$

68
for $\omega=\left(\omega_{1}, \ldots, \omega_{d}\right) \in \Omega, \quad \Phi=\left(\Phi^{(1)}, \ldots, \Phi^{(d)}\right) \in(S([0,1]))^{d}$, and $\langle\omega, \phi\rangle=\sum_{i=1}^{d}\left\langle\omega_{i}, \phi_{i}\right\rangle$
In what follows, we shall denote by $D$. the Malliavin derivative on $(\Omega, \mathcal{F}, \mu)$, which is a linear operator from $\mathbb{D}_{1,2}$ to $L^{2}(\lambda \otimes \mu)(\lambda$ Lebesgue measure $)$. See e.g [11] or [21] for the definition of $D$.. We Mention that $\mathbb{D}_{1,2}$ in 1.3 is a Hilbert space with a norm $\|\cdot\|_{1,2}$ given by

$$
\begin{equation*}
\|F\|_{1,2}^{2}:=\|F\|_{L^{2}(\mu)}^{2}+\|D \cdot F\|_{L^{2}([0,1] \times \Omega, \lambda \otimes \mu)}^{2} \tag{2.3}
\end{equation*}
$$

(for $d=1$ ). We shall also use the notation $\delta$ for the adjoint operator of $D$., which is referred to as divergence operator.

In this section, we also want to introduce the crucial concept of stochastic integration

$$
\begin{equation*}
\int_{0}^{t} \int_{\mathbb{R}} f(s, x) L(d s, d x) \tag{2.4}
\end{equation*}
$$

over the plane with respect to Brownian local time $L(t, x)$ for integrands $f:[0,1] \times \mathbb{R} \rightarrow \mathbb{R}$ in the Banach space $(\mathcal{H},\|\cdot\|)$ with the norm

$$
\begin{align*}
\|f\|:= & 2\left(\int_{0}^{1} \int_{\mathbb{R}}(f(s, x))^{2} \exp \left(-\frac{x^{2}}{2 s}\right) \frac{d s d x}{\sqrt{2 \pi s}}\right)^{\frac{1}{2}} \\
& +\int_{0}^{1} \int_{\mathbb{R}}|x f(s, x)| \exp \left(-\frac{x^{2}}{2 s}\right) \frac{d s d x}{s \sqrt{2 \pi s}} \tag{2.5}
\end{align*}
$$

77 See [12]. We need the following auxiliary result ([12, Theorem 3.1, Corollary 3.2])
Lemma 2.1. Let $f \in \mathcal{H}$. Suppose that for all $t \in[0,1] f(t, \cdot)$, the derivative $f^{\prime}(t, \cdot)$ (in the generalized sense with respect to the Lebesgue measure) exists and that

$$
\int_{0}^{1} \int_{-A}^{A}\left|f^{\prime}(s, x)\right| \frac{d s}{\sqrt{s}} d x<\infty
$$

78 for all $A \geq 0$. Then

$$
\begin{equation*}
\int_{0}^{t} \int_{\mathbb{R}} f(s, x) L(d s, d x)=-\int_{0}^{t} f^{\prime}\left(s, B_{s}\right) d s \tag{2.6}
\end{equation*}
$$

$$
\begin{align*}
\int_{0}^{t} \int_{\mathbb{R}} f_{i}(s, x) L_{i}(d s, d x)= & \int_{0}^{t} f_{i}\left(s, B_{s}^{(i)}\right) d B_{s}^{(i)}+\int_{1-t}^{1} f_{i}\left(1-s, \widehat{B}_{s}^{(i)}\right) d \widetilde{W}_{s}^{(i)} \\
& +\int_{1-t}^{1} f_{i}\left(1-s, \widehat{B}_{s}^{(i)}\right) \frac{\widehat{B}_{s}^{(i)}}{1-s} d s \tag{2.7}
\end{align*}
$$

$0 \leq t \leq 1$. Further $\widetilde{W}_{t}^{(i)}, 0 \leq t \leq 1$, are independent $\mu_{i}$-Brownian motions (see 2.1) with respect to the filtration $\mathcal{F}_{t}^{\widehat{B}^{(i)}}$ generated by $\widehat{B}_{t}^{(i)}, i=1, \ldots, d$.

Now consider the SDE's 1.1$]$ with Borel measurable drift $b:[0,1] \times \mathbb{R}^{d} \rightarrow \mathbb{R}^{d}$. For our main result (Theorem 2.2 we will need the existence of a sequence $b_{p}:[0,1] \times \mathbb{R}^{d} \rightarrow \mathbb{R}^{d}, p \in \mathbb{N}$, of approximating drift coefficients which fulfill the following five conditions. For notational convenience we set $b_{0}:=b$.
(C1): The coefficients $b_{p}, p \in \mathbb{N}$, are continuous with compact support such that $b_{p}(t, \cdot)$ is continuously differentiable, $0 \leq t \leq 1$, with bounded derivative on $[0,1] \times \mathbb{R}^{d}$. It is well known that bounded coefficients admit unique strong solutions $X_{t}^{n, p}, n \in \mathbb{N}, p \in \mathbb{N}$, of the SDE's

$$
\begin{equation*}
d X_{t}^{n, p}=b\left(t, X_{t}^{n, p}\right) d t+\frac{1}{n} d B_{t}, 0 \leq t \leq 1, \quad X_{0}^{n, p}=x \in \mathbb{R}^{d} \tag{2.9}
\end{equation*}
$$

93 (C2): Let $\mathcal{M} \subset \mathbb{R}^{d \times d}$ denote the class of continuous matrix valued functions $\mathcal{M}(t):[0,1] \rightarrow$ $94 \mathbb{R}^{d \times d}$ such that $\mathcal{M}(t)$ commutes with $\int_{s}^{t} \mathcal{M}(u) d u$ for all $0 \leq s \leq t \leq 1$. Suppose that $95 b_{p}^{\prime}\left(\cdot, X^{n, p}\right) \in \mathcal{M}$ for all $n \in \mathbb{N}, p \in \mathbb{N}$, where the symbol ' stands for the derivative with 96 respect to the space variable.
(C3): For each $n \in \mathbb{N}$

$$
\sup _{p \geq 0}\left\|\exp \left\{512 \int_{0}^{1} n^{2}\left\|b_{p}\left(s, \frac{B_{s}}{n}+x\right)\right\|^{2} d s\right\}\right\|_{L^{1}(\mu)}<\infty
$$

and the sequence of coefficients $b_{p}, p \in \mathbb{N}$, approximates $b$ in the sense that for each $n \in \mathbb{N}$

$$
E\left[J_{n, p}\right] \underset{p \rightarrow \infty}{\longrightarrow} 0
$$

97 where

$$
\begin{align*}
J_{n, p} & =\sum_{j=1}^{d}\left(2 \int_{0}^{1}\left(n b_{p}^{(j)}\left(s, \frac{B_{s}}{n}+x\right)-n b^{(j)}\left(s, \frac{B_{s}}{n}+x\right)\right)^{2} d s\right. \\
& \left.+\left(\int_{0}^{1}\left|\left(n b_{p}^{(j)}\left(s, \frac{B_{s}}{n}+x\right)\right)^{2}-\left(n b^{(j)}\left(s, \frac{B_{s}}{n}+x\right)\right)^{2}\right| d s\right)^{2}\right) \tag{2.10}
\end{align*}
$$

$$
\begin{align*}
A_{1}\left(n, p, t, t^{\prime}\right)= & \exp \left\{\int_{0}^{1} n b_{p}\left(s, \frac{B_{s}}{n}+x\right) d B_{s}-\frac{1}{2} \int_{0}^{1} n^{2}\left\|b_{p}\left(s, \frac{B_{s}}{n}+x\right)\right\|^{2} d s\right\}  \tag{2.12}\\
A_{2}\left(n, p, t, t^{\prime}\right)= & \| \exp \left\{\left(-\int_{t^{\prime}}^{1} n b_{p}^{(j)}\left(s, \frac{B_{s}}{n}+x\right) d B_{s}^{(i)}-\int_{0}^{1-t^{\prime}} n b_{p}^{(j)}\left(1-s, \frac{\widehat{B}_{s}}{n}+x\right) d \widetilde{W}_{s}^{(i)}\right.\right. \\
& \left.\left.+\int_{0}^{1-t^{\prime}} n b_{p}\left(1-s, \frac{\widehat{B}_{s}}{n}+x\right) \frac{\widehat{B}_{s}^{(i)}}{1-s} d s\right)_{0 \leq i, j \leq d}\right\} \|^{2}  \tag{2.13}\\
A_{3}\left(n, p, t, t^{\prime}\right)= & \sup _{0 \leq \lambda \leq 1} \| \exp \left\{\left(-\lambda \int_{t}^{t^{\prime}} n b_{p}^{(j)}\left(s, \frac{B_{s}}{n}+x\right) d B_{s}^{(i)}-\lambda \int_{1-t^{\prime}}^{1-t} n b_{p}^{(j)}\left(1-s, \frac{\widehat{B}_{s}}{n}+x\right) d \widetilde{W}_{s}^{(i)}\right.\right. \\
& \left.\left.+\lambda \int_{1-t^{\prime}}^{1-t} n b_{p}^{(j)}\left(1-s, \frac{\widehat{B}_{s}}{n}+x\right) \frac{\widehat{B}_{s}^{(i)}}{1-s} d s\right)_{0 \leq i, j \leq d}^{2}\right\} \|^{2}, \tag{2.14}
\end{align*}
$$

$$
\begin{equation*}
A_{4}\left(n, p, t, t^{\prime}\right)=\frac{1}{n^{2}} \frac{\left\|I_{4}\left(n, p, t, t^{\prime}\right)\right\|^{2}}{\left|t-t^{\prime}\right|^{\alpha}}, t \neq t^{\prime} \tag{2.15}
\end{equation*}
$$

where
for some $\alpha>\frac{1}{2}$ with

$$
\begin{align*}
I_{4}\left(n, p, t, t^{\prime}\right)= & \left(\int_{t}^{t^{\prime}} n b_{p}^{(j)}\left(s, \frac{B_{s}}{n}+x\right) d B_{s}^{(i)}-\int_{1-t^{\prime}}^{1-t} n b_{p}^{(j)}\left(1-s, \frac{\widehat{B}_{s}}{n}+x\right) d \widetilde{W}_{s}^{(i)}\right. \\
& \left.+\int_{1-t^{\prime}}^{1-t} n b_{p}^{(j)}\left(1-s, \frac{\widehat{B}_{s}}{n}+x\right) \frac{\widehat{B}_{s}^{(i)}}{1-s} d s\right)_{0 \leq i, j \leq d} \tag{2.16}
\end{align*}
$$

(C5):

$$
\begin{equation*}
\sup _{n, p \geq 1} \sup _{0 \leq t<t^{\prime} \leq 1}\left\|A_{5}\left(n, p, t, t^{\prime}\right) A_{1}\left(n, p, t, t^{\prime}\right)\right\|_{L^{1}(\mu)}<\infty, \tag{2.17}
\end{equation*}
$$

where

$$
\begin{equation*}
A_{5}\left(n, p, t, t^{\prime}\right)=\frac{\left\|\int_{t}^{t^{\prime}} b_{p}\left(s, \frac{B_{s}}{n}+x\right) d s\right\|^{2}}{\left|t-t^{\prime}\right|^{\beta}}, t \neq t^{\prime} \tag{2.18}
\end{equation*}
$$

for some $\beta>\frac{1}{2}$.
Theorem 2.2. Consider the family of SDE's in (1.1) with Borel measurable drift coefficient $b:[0,1] \times \mathbb{R}^{d} \rightarrow \mathbb{R}^{d}$. Suppose there exists a sequence of approximating coefficients $\left(b_{p}\right)_{p \geq 1}$ such that $\left\{b,\left(b_{p}\right)_{p \geq 1}\right\}$ fulfill conditions (C1)-(C5). Then for all $0 \leq t \leq 1$ the set of solutions $\left(X_{t}^{n, p}\right)_{n \geq 1, p \geq 1}$ of $(2.9)$ is relatively compact in $L^{2}\left(\mu ; \mathbb{R}^{d}\right)$. Further, for all $n \in \mathbb{N}$ there exist
a unique strong solution $X_{t}^{n}$ of (1.1) and the sequence of solutions $X_{t}^{n}$ to (1.1) is relatively compact in $L^{2}\left(\mu ; \mathbb{R}^{d}\right), 0 \leq t \leq 1$, and there exists a cluster point $\left(X_{t}\right)_{0 \leq t \leq 1}$ of $\left(X_{t}^{n}\right)_{0 \leq t \leq 1}$, that is one finds a subsequence $\left(n_{m}\right)_{m \geq 1}$ such that

$$
\begin{equation*}
\lim _{m \rightarrow \infty} X_{t}^{n_{m}}=X_{t} \text { in } L^{2}\left(\mu ; \mathbb{R}^{d}\right) \tag{2.19}
\end{equation*}
$$

for all $0 \leq t \leq 1$. In particular, if $\left\|b\left(t, X_{t}^{n}\right)\right\|_{L^{2}(\mu)} \leq M<\infty, n \geq 1, t$-a.e for some constant $M$, then

$$
\begin{equation*}
X_{t}=x+\int_{0}^{t} \lim _{m \rightarrow \infty} b\left(s, X_{s}^{n_{m}}\right) d s \tag{2.20}
\end{equation*}
$$

in $L^{2}(\mu)$.
Remark 2.3. Note that in case of a bounded drift coefficient b there obviously exists a sequence of approximating coefficients $\left(b_{p}\right)_{p \geq 1}$ that fulfill conditions (C1), (C3), and (C5). In that case, the crucial conditions to check are (C2) and (C4).
Remark 2.4. In the case of dimension $d=1$, the commutativity requirement (C2) is obviously always fulfilled. In the case $d=2$, condition (C2) can be verified, if e.g.,

$$
b(t, x)=\binom{f\left(x_{1}+x_{2}\right)}{f\left(x_{1}+x_{2}\right)},
$$

where $f: \mathbb{R} \rightarrow \mathbb{R}$ is a bounded Borel measurable function. See 19 for other examples and more general criteria.

We postpone the proof of Theorem 2.2 to a later time point. In the sequel, we discuss some consequences of the previous result:

Corollary 2.5. Retain the conditions in Theorem 2.2 and assume additionally that the drift coefficient $b$ in (1.1) is continuous. Then there exists a Malliavin differentiable process $X_{t}$ such that

$$
\begin{equation*}
X_{t}=x+\int_{0}^{t} b\left(s, X_{s}\right) d s \tag{2.21}
\end{equation*}
$$

Proof. Equation (2.21) follows from (2.20) and the continuity of $b$. The Malliavin differentiability of $X_{t}$ follows from a weak compactness argument. See the proof of Theorem 2.2 .

The next two result treats the case of discontinuous ODE's:
Theorem 2.6. Keep the conditions in Theorem 2.2 and assume additionally that the drift coefficient $b$ in (1.1) is bounded. Further require that the process $X_{t}$ in (2.19) doesn't hit the set of points of discontinuity of $b(t, \cdot) \mu$-a.e. for almost all (fixed) $t$. Then $X_{t}$ solves the ODE

$$
\begin{equation*}
X_{t}=x+\int_{0}^{t} b\left(s, X_{s}\right) d s \tag{2.22}
\end{equation*}
$$

Theorem 2.7. Retain the conditions in Theorem 2.2 and require additionally that the drift coefficient $b$ in 1.1) is bounded and time-homogeneous. Then

$$
X_{t}^{(i)} \in \mathbb{D}_{1,2}
$$

for all $i=1, \ldots, d, 0 \leq t \leq 1$. Moreover, if the Malliavin matrix $\sigma_{X_{t}}=\left(\sigma_{X_{t}}^{i, j}\right)_{1 \leq i, j \leq d}$ with

$$
\sigma_{X_{t}}^{i, j}=\left(D \cdot X_{t}^{(i)}, D \cdot X_{t}^{(j)}\right)_{L^{2}([0,1])}
$$

is invertible a.e for each $t$, then $X_{t}$ is a solution of (2.2.2).

The proofs of these two theorems are also put off to a later time point.
The following result will be needed in the proof of Theorem 2.2 .

Lemma 2.8. Suppose that the conditions of Theorem 2.2 hold. Then the double sequence $\left(t \longmapsto X_{t}^{n, p}, \quad n, p \geq 1\right)$ is relatively compact in $C\left([0,1],(\mathcal{S})^{*}\right)$.

Proof. Let $\zeta$ belong to the Hida test function space $(\mathcal{S})$. Denote by $\langle F, \rho\rangle$ the dual pairing for $F \in(\mathcal{S})^{*}, \quad \rho \in(\mathcal{S})$. Using the Cauchy-Schwartz inequality, Girsanov's theorem and (C3), and (C5) we get that

$$
\begin{aligned}
\left|\left\langle X_{t_{1}}^{n, p}-X_{t_{2}}^{n, p}, \zeta\right\rangle\right| & =E\left[\left(X_{t_{1}}^{n, p}-X_{t_{2}}^{n, p}\right) \zeta\right] \leq E\left[\left\|X_{t_{1}}^{n, p}-X_{t_{2}}^{n, p}\right\|^{2}\right]^{\frac{1}{2}} E\left[|\zeta|^{2}\right]^{\frac{1}{2}} \\
& \leq C\left|t_{2}-t_{1}\right|^{\beta} E\left[|\zeta|^{2}\right]^{\frac{1}{2}}
\end{aligned}
$$

for some $\beta>\frac{1}{2}$. On the other hand, we directly see that

$$
\sup _{0 \leq t \leq T}\left\|X_{t}^{n, p}\right\|_{L^{2}(\mu)} \leq M
$$

for all $n, p \geq 1$. The desired result then follows from Mitoma's theorem (see [20]) applied to the conuclear space $(\mathcal{S})^{*}$ and Arzelá-Ascoli's theorem with respect to $C([0,1])$.

Proof. (Theorem 2.2).
We first want to employ a compactness criterion based on Malliavin calculus [10, Theorem 1] to show that $\left(X_{t}^{n, p}\right)_{p \geq 0, n \geq 1}$ is relatively compact in $L^{2}\left(\mu ; \mathbb{R}^{d}\right)$ for all $t \geq 0$. To this end we assume without loss of generality that $t=1$. Our assumptions and the chain rule of the Malliavin derivative $D_{t}$ (see e.g., [21]) imply that

$$
\begin{equation*}
D_{t} X_{1}^{n, p}=\frac{1}{n} \exp \left\{\int_{t}^{1} b_{p}^{\prime}\left(s, \frac{X_{s}^{n, p}}{n}\right) d s\right\} \in \mathbb{R}^{d \times d}, \quad 0 \leq t \leq 1, n, p \geq 1 \tag{2.23}
\end{equation*}
$$

Fix $0 \leq t<t^{\prime} \leq 1$. Then using Girsanov's theorem we find that

$$
\begin{aligned}
& E\left[\left\|D_{t} X_{1}^{n, p}-D_{t^{\prime}} X_{1}^{n, p}\right\|^{2}\right] \\
= & \frac{1}{n^{2}} E\left[\left\|\exp \left\{\int_{t}^{1} b_{p}^{\prime}\left(s, \frac{B_{s}}{n}+x\right) d s\right\}-\exp \left\{\int_{t^{\prime}}^{1} b_{p}^{\prime}\left(s, \frac{B_{s}}{n}+x\right) d s\right\}\right\|^{2} A_{1}\right]
\end{aligned}
$$

where $A_{1}=\exp \left\{\int_{0}^{1} n b_{p}\left(s, \frac{B_{s}}{n}+x\right) d B_{s}-\frac{1}{2} \int_{0}^{1} n^{2}\left\|b_{p}\left(s, \frac{B_{s}}{n}+x\right)\right\|^{2} d s\right\}$.

Applying the properties of evolution operators for linear systems of ODE's and the mean value theorem, we have

$$
\begin{aligned}
& E\left[\left\|D_{t} X_{1}^{n, p}-D_{t^{\prime}} X_{1}^{n, p}\right\|^{2}\right] \\
= & \frac{1}{n^{2}} E\left[\left\|\exp \left\{\int_{t^{\prime}}^{1} b_{p}^{\prime}\left(s, \frac{B_{s}}{n}+x\right) d s\right\}\right\|^{2}\left\|\exp \left\{\int_{t}^{t^{\prime}} b_{p}^{\prime}\left(s, \frac{B_{s}}{n}+x\right) d s\right\}-1\right\|^{2} A_{1}\right] \\
\leq & \frac{1}{n^{2}} C E\left[\left\|\exp \left\{\int_{t^{\prime}}^{1} b_{p}^{\prime}\left(s, \frac{B_{s}}{n}+x\right) d s\right\}\right\|^{2}\left\|\left\{\int_{t}^{t^{\prime}} b_{p}^{\prime}\left(s, \frac{B_{s}}{n}+x\right) d s\right\}\right\|^{2} .\right. \\
& \left.\sup _{0 \leq \lambda \leq 1}\left\|\exp \left\{\lambda \int_{t}^{t^{\prime}} b_{p}^{\prime}\left(s, \frac{B_{s}}{n}+x\right) d s\right\}\right\|^{2} A_{1}\right]
\end{aligned}
$$

Consider the local time-space $L_{i}(d s, d x)$ with respect to $B^{(i)}$ (the i-th component of $B$ ) on $\left(\Omega_{i}, \mu_{i}\right), \quad i=1, \ldots, d$. Using Lemma 2.1 and the decomposition 2.7), we get

$$
\begin{aligned}
& E\left[\left\|D_{t} X_{1}^{n, p}-D_{t^{\prime}} X_{1}^{n, p}\right\|^{2}\right] \\
& \leq C E\left[\left\|\exp \left\{\left(-\int_{t^{\prime}}^{1} \int_{\mathbb{R}} n b_{p}^{(j)}\left(s, \frac{x}{n}\right) L_{i}(d s, d x)\right)_{1 \leq i, j \leq d}\right\}\right\|^{2}\right. \\
& \left\|\left\{\left(-\int_{t}^{t^{\prime}} \int_{\mathbb{R}} n b_{p}^{(j)}\left(s, \frac{x}{n}\right) L_{i}(d s, d x)\right)_{1 \leq i, j \leq d}\right\}\right\|^{2} \\
& \left.\sup _{0 \leq \lambda \leq 1}\left\|\exp \left\{\lambda\left(-\int_{t}^{t^{\prime}} \int_{\mathbb{R}} n b_{p}^{(j)}\left(s, \frac{x}{n}\right) L_{i}(d s, d x)\right)_{1 \leq i, j \leq d}\right\}\right\|^{2} A_{1}\right] \\
& \leq C\left|t^{\prime}-t\right|^{\alpha}\left(\sup _{n, p \geq 1} \sup _{0 \leq t<t^{\prime} \leq 1}\left\|\prod_{i=1}^{4} A_{i}\left(n, p, t, t^{\prime}\right)\right\|_{L^{1}(\mu)}\right)
\end{aligned}
$$

for some constant $C$. In particular, we see that the family $\left(X_{1}^{n, p}\right)_{p \geq 0, n \geq 1}$ is bounded in $\mathbb{D}_{1,2}$. Then the relative compactness of $\left(X_{1}^{n, p}\right)_{p \geq 0, n \geq 1}$ follows from [[10], Lemma 1] in connection with [10, Theorem 1].

In the next step of the proof we aim at constructing a solution process $X_{t}$ to the ODE's 1.2 based on the double sequence $\left(X_{t}^{n, p}\right)_{p \geq 1, n \geq 1}$. Using the condition (C3) in connection with Theorem 4 in [19], we obtain that for all $n \geq 1$ there exists a subsequence ( $p_{k, n}$ ) (independent of $t$ ) such that

$$
X_{t}^{n}=\lim _{k \rightarrow \infty} X_{t}^{n, p_{k, n}} \in L^{2}\left(\mu ; \mathbb{R}^{d}\right)
$$

satisfies the SDE's 1.1). In particular, $\left(X_{t}^{n}\right)_{n \geq 1}$ is relatively compact in $L^{2}\left(\mu ; \mathbb{R}^{d}\right)$ for each $t$. We also mention that $X_{t}^{n}$ is Malliavin differentiable for all $n, t$ by a weak compactness argument (see [19, Lemma 1,2,3]).

On the other hand, it follows from Lemma 2.8 that there exists a subsequence $\left(n_{k}\right)$ such that

$$
X_{t}^{n_{k}} \underset{k \rightarrow \infty}{\longrightarrow} X_{t} \text { in }(\mathcal{S})^{*}
$$

uniformly in $t$. The latter and the uniqueness of chaos decompositions in $(\mathcal{S})^{*}$ entail that

$$
X_{t}^{n_{k}} \underset{k \rightarrow \infty}{\longrightarrow} X_{t} \text { in } L^{2}\left(\mu ; \mathbb{R}^{d}\right)
$$

for all $t$.
Finally, if the drift coefficient is bounded, we can apply dominated convergence for functions from $[0,1]$ to $L^{2}\left(\mu ; \mathbb{R}^{d}\right)$ and obtain 2.20 .

Proof. (Theorem 2.6).
We shall argue by contradiction. Assume that $b\left(t, X_{t}^{n}\right)$ does not converge to $b\left(t, X_{t}\right)$ in $L^{2}(\mu)$ for some $t$ for which the points of discontinuity cannot be reached. Then there exists a $\epsilon>0$ and a subsequence $\left(n_{k}\right)$ such that

$$
\begin{equation*}
\left\|b\left(t, X_{t}^{n_{k}}\right)-b\left(t, X_{t}\right)\right\|_{L^{2}(\mu)}>\epsilon \tag{2.24}
\end{equation*}
$$

We know that

$$
X_{t}^{n_{\tilde{n}_{l}(t)}} \longrightarrow X_{t} \text { a.e. }
$$

for some subsequence $\left(\tilde{n}_{l}(t)\right)$. Using the fact that $X_{t}$ doesn't hit the points of discontinuity of $b(t, \cdot)$ a.e., we see that

$$
b\left(t, X_{t}^{n_{\tilde{n}_{l}(t)}}\right) \longrightarrow b\left(t, X_{t}\right) \text { a.e. }
$$

Since $b$ is bounded, it follows from the dominated convergence theorem that

$$
\left\|b\left(t, X_{t}^{n_{\tilde{n}_{l}(t)}}\right)-b\left(t, X_{t}\right)\right\|_{L^{2}(\mu)} \underset{l \rightarrow \infty}{\longrightarrow} 0
$$

For $k=\tilde{n}_{l}(t)$, this leads to a contradiction to 2.24 . Therefore

$$
\lim _{n \rightarrow \infty} b\left(t, X_{t}^{n}\right)=b\left(t, \overline{\left.X_{t}\right) \text { in }} L^{2}(\mu), t\right. \text {-a.e. }
$$

Proof. (Theorem 2.7).
We recall that each $X_{s}^{n}$ is Malliavin differentiable (see [19]). We want to justify that we may set $b_{p}=b$ for all $p \geq 1$ in the proof of Theorem 2.2 . To this end we shall derive a certain representation for $D_{t} X_{s}^{n}$ by employing the $S$-transform (see $(2.2)$ ). Without loss of generality, we assume that $s=1$ and $d=1$ (one-dimensional case). Let us evaluate

$$
S\left(D_{t} X_{1}^{n, p}\right)(\phi), \quad \phi \in S_{\mathbb{C}}(\mathbb{R}), \quad n \geq 1
$$

Then, using Girsanov's theorem and the local time-space decomposition (2.7), we find that

$$
\begin{align*}
& S\left(D_{t} X_{1}^{n, p}\right)(\phi) \\
&= E\left[\frac { 1 } { n } \operatorname { e x p } \left\{\int_{t}^{1} n b_{p}\left(\frac{1}{n} B_{s}+x\right) d B_{s}-\int_{0}^{1-t} n b_{p}\left(\frac{1}{n} B_{s}+x\right) d \widetilde{W}_{s}\right.\right. \\
&\left.+\int_{0}^{1-t} n b_{p}\left(\frac{1}{n} \widehat{B}_{s}+x\right) \frac{\widehat{B}_{s}}{1-s} d s\right\}  \tag{2.25}\\
&\left.\exp \left\{\int_{0}^{1}\left(n b_{p}\left(\frac{1}{n} B_{s}+x\right)+\phi(x)\right) d B_{s}-\frac{1}{2} \int_{0}^{1}\left(n b_{p}\left(\frac{1}{n} B_{s}+x\right)+\phi(x)\right)^{2} d s\right\}\right]
\end{align*}
$$

Hilbert spaces, we deduce that

$$
\begin{align*}
& S\left(\int_{0}^{1} D_{t} X_{1}^{n} \cdot h(t) d t\right)(\phi) \\
& =E\left[\int _ { 0 } ^ { 1 } \left(\frac { 1 } { n } \operatorname { e x p } \left\{-\int_{t}^{1} n b\left(\frac{B_{s}}{n}+x\right) d B_{s}-\int_{0}^{1-t} n b\left(\frac{B_{s}}{n}+x\right) d \widetilde{W}_{s}\right.\right.\right. \\
& \left.\quad+\int_{0}^{1-t} n b\left(\frac{\widehat{B}_{s}}{n}+x\right) \frac{\widehat{B}_{s}}{1-s} d s\right\}  \tag{2.26}\\
& \left.\left.\exp \left\{\int_{0}^{1}\left(n b\left(\frac{B_{s}}{n}+x\right)+\phi(s)\right) d B_{s}-\frac{1}{2} \int_{0}^{1}\left(n b\left(\frac{B_{s}}{n}+x\right)+\phi(s)\right)^{2} d s\right\}\right) h(t) d t\right]
\end{align*}
$$

for all bounded Borel-measurable functions $h$ on $[0,1], \phi \in S_{\mathbb{C}}(\mathbb{R})$ and $n \geq 1$. Repeated use of the local time-space decomposition (2.7), Girsanov's theorem and the Itô-Tanaka formula for continuous semimartingales in [23, p.220] give that

$$
S\left(D_{t} X_{1}^{n}\right)(\phi)=S\left(\Psi_{t}^{n}\right)(\phi)
$$

for all $\phi \in S_{\mathbb{C}}(\mathbb{R})$, where

$$
\Psi_{t}^{n}=\frac{1}{n} \exp \left\{\int_{t}^{1} \int_{\mathbb{R}} n b\left(\frac{y}{n}+x\right) L^{n\left(X^{n}-x\right)}(d s, d y)\right\}
$$

where $L^{n\left(X^{n}-x\right)}(s, y)$ denotes the local time at $y$ of $n\left(X^{n}-x\right)$. Thus

$$
\begin{equation*}
D \cdot X_{1}^{n}=\Psi^{n} \tag{2.27}
\end{equation*}
$$

for all $n$.
Using this representation and the line of reasoning in the proof of Theorem 2.2 in connection with the weak compactness in $\mathbb{D}_{1,2}$, we conclude that $X_{t}$ is Malliavin differentiable for all $t$.

The last statement of Theorem 2.7 is a direct consequence of [21, Theorem 2.1.2]
Remark 2.9. Assume $b: \mathbb{R} \rightarrow \mathbb{R}$ satisfies the assumptions of Theorem 2.6. Consider the case, when

$$
\begin{equation*}
D \cdot X_{u}=0 \tag{2.28}
\end{equation*}
$$

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on a measurable set $A$ such that $(\lambda \otimes \mu)(A)>0$ for some $0<u \leq 1$. Then using relation (2.27) in the proof of Theorem 2.6 in connection with Girsanov's theorem shows that there is a subsequence $n_{k}$ such that

$$
\begin{equation*}
-\log n_{k}+\mathcal{L}_{1}\left(n_{k}, t, u\right)+\mathcal{L}_{2}\left(n_{k}, u\right) \underset{k \rightarrow \infty}{\longrightarrow}-\infty \tag{2.29}
\end{equation*}
$$

on $A(t, \omega)$-a.e., where

$$
\begin{aligned}
\mathcal{L}_{1}(n, t, u)= & \int_{t}^{u} n b_{p}\left(\frac{1}{n} B_{s}+x\right) d B_{s}-\int_{1-u}^{1-t} n b_{p}\left(\frac{1}{n} B_{s}+x\right) d \widetilde{W}_{s} \\
& +\int_{1-u}^{1-t} n b_{p}\left(\frac{1}{n} \widehat{B}_{s}+x\right) \frac{\widehat{B}_{s}}{1-s} d s
\end{aligned}
$$

$$
\begin{equation*}
\rho_{s}^{n, p}(y)=E\left[\chi_{(y, \infty)}\left(X_{s}^{n, p}\right) \delta\left(\frac{D \cdot X_{s}^{n, p}}{\left\|D \cdot X_{s}^{n, p}\right\|_{L^{2}[0,1]}^{2}}\right)\right], y \in \mathbb{R}, n, p \geq 1 \tag{2.32}
\end{equation*}
$$

See [21, Proposition 2.1] or [11]. Consider now the sequence of Lipschitz continuous functions $0 \leq \varrho_{m} \leq \chi_{(x, \infty)}$ with $\varrho_{m}(z) \rightarrow \chi_{(y, \infty)}(z), \quad z \in \mathbb{R}$ given by

$$
\varrho_{m}(z)= \begin{cases}m z-m y & , \quad y<z<y+\frac{1}{m} \\ 0 & , \quad z \leq y \\ 1 & , \quad z \geq y+\frac{1}{m}\end{cases}
$$

Then the functions $\rho_{s}^{m, n, p}$ defined as

$$
\rho_{s}^{m, n, p}(y)=E\left[\varrho_{m}\left(X_{s}^{n, p}\right) \delta\left(\frac{D \cdot X_{s}^{n, p}}{\left\|D \cdot X_{s}^{n, p}\right\|_{L^{2}[0,1]}^{2}}\right)\right]
$$

converge to $\rho_{s}^{n, p}$, pointwisely for all $s, n, p$. On the other hand one infers from the duality relation and the chain rule of the Malliavin derivative (see e.g [21, 11]) that

$$
\rho_{s}^{m, n, p}(y)=E\left[\int_{0}^{s} \chi_{\left(y, y+\frac{1}{m}\right)}\left(X_{s}^{n, p}\right) \frac{\left(D_{u} X_{s}^{n, p}\right)^{2}}{\left\|D \cdot X_{s}^{n, p}\right\|_{L^{2}[0,1]}^{2}} d u\right]
$$

Then we obtain from $(2.31)$ in connection with the Girsanov's theorem and the decomposition (2.7) that

$$
\left\|\rho_{s}^{m, n, p}\right\|_{L^{2}(\mathbb{R})}^{2} \leq M<\infty \quad \text { for all } m, n, p
$$

Using weak compactness of $\rho_{s}^{m, n, p}, m, n, p$ in $L^{2}(\mathbb{R})$, pointwise convergence of $\rho_{s}^{m, n, p}$ with respect to $m$ and the fact that $X_{s}^{n, p}$ converges to $X_{s}^{n}$ in $L^{2}(\mu)$ (for a subsequence), we observe that $X_{s}^{n}$ has a probability density $\rho_{s}^{n}$ and that $\rho_{s}^{n}$ is weakly compact in $L^{2}(\mathbb{R})$. Repeated use of weak compactness and $L^{2}(\mu)$-convergence shows that the cluster point $X_{s}$ in Theorem 2.6 has a density $\rho_{s}, 0<s \leq 1$. So the result follows.

Finally, we give an application of Theorem 2.6 in the case of a discontinuous ODE.
Example 2.11. Consider the ODE (1.2) with initial value $x$ and the drift coefficient $b$ given by the sign function, that is the special case of a step function

$$
b(t, y)=\operatorname{sign}(y)=\left\{\begin{array}{lll}
1 & , y \geq 0 \\
-1 & , y<0
\end{array}\right.
$$

We want to show that there exists a subsequence $\left(n_{k}\right)$ such that the solutions $X_{s}^{n}$ converge in $\mathbb{D}_{1,2}$ to a deterministic process $X_{s}, 0 \leq s \leq 1$ (for certain $x \neq 0$ ).

Without loss of generality, let $s=1$. Since the sign function is bounded, we know from the proof of Theorem 2.6 that

$$
D_{t} X_{1}^{n}=\frac{1}{n} \exp \left\{-\int_{t}^{1} \int_{\mathbb{R}} n \operatorname{sign}\left(\frac{1}{n} y+x\right) L^{n\left(X^{n}-x\right)}(d s, d y)\right\}
$$

where $L^{n\left(X^{n}-x\right)}(s, y)$ is the local time at $y$ of $n\left(X^{n}-x\right)$. Using the latter representation, we may replace the coefficient $b_{p}, p \geq 1$ in Theorem 2.2 by the sign function itself. In order to verify condition (C4) we apply Girsanov's theorem and Hölder's inequality and find that it is sufficient to show that

$$
\begin{equation*}
I_{1}\left(n, t, t^{\prime}\right) \cdot I_{2}\left(n, t, t^{\prime}\right) \leq C \cdot\left|t-t^{\prime}\right|^{\alpha}, \quad 0 \leq t \leq t^{\prime} \leq 1 \tag{2.33}
\end{equation*}
$$

for some $\alpha>\frac{1}{2}$ and a constant $C$ (independent of $n$ ), where

$$
\begin{equation*}
I_{1}\left(n, t, t^{\prime}\right):=\frac{1}{n^{2}} E\left[\left(-\int_{t}^{t^{\prime}} \int_{\mathbb{R}} n \operatorname{sign}\left(\frac{1}{n} y+x\right) L(d s, d y)\right)^{4}\right]^{\frac{1}{2}} \tag{2.34}
\end{equation*}
$$

and

$$
\begin{align*}
I_{2}\left(n, t, t^{\prime}\right): & =E\left[\exp \left\{-4 \int_{t}^{t^{\prime}} \int_{\mathbb{R}} n \operatorname{sign}\left(\frac{1}{n} y+x\right) L(d s, d y)\right\}\right. \\
& \left.\cdot \exp \left\{\int_{0}^{1} n \operatorname{sign}\left(\frac{1}{n} B_{s}+x\right) d B_{s}-\frac{1}{2} \int_{0}^{1} n^{2} d s\right\}\right]^{\frac{1}{2}} \tag{2.35}
\end{align*}
$$

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Using the Itô-Tanaka formula and Burkholder's inequality we find that

$$
I_{1}\left(n, t, t^{\prime}\right)
$$

$$
=\frac{1}{n^{2}} E\left[\left(\int_{t}^{t^{\prime}} n^{2}\left(\operatorname{sign}\left(X_{u}^{n}\right)\right)^{2} d u+\int_{t}^{t^{\prime}} n \operatorname{sign}\left(X_{u}^{n}\right) d B_{u}-\left(\left|n X_{t^{\prime}}^{n}-n x\right|-\mid n X_{t}^{n}-n x\right)\right)^{4}\right]^{\frac{1}{2}}
$$

$$
\begin{equation*}
\leq C n^{4}\left|t-t^{\prime}\right| \tag{2.36}
\end{equation*}
$$

for some constant $C$.
On the other hand, by applying [12, Corollary 3.2] we get that

$$
\begin{aligned}
-\int_{t}^{1} \int_{\mathbb{R}} n \operatorname{sign}\left(\frac{1}{n} y+x\right) L(d s, d y) & =\int_{\mathbb{R}}\left(\int_{t}^{1} 2 n d_{s} L_{s}^{y}\right) \delta_{\{-n x\}}(d y) \\
& =2 n(L(1,-n x)-L(t,-n x))
\end{aligned}
$$

where $\delta_{\{-n x\}}$ is the Dirac measure in $-n x$.
Repeated use of Girsanov's theorem and the formula of Itô-Tanaka gives that

$$
\begin{aligned}
I_{2}\left(n, t, t^{\prime}\right)= & E[\exp \{8 n(L(1,-n x)-L(t,-n x))\} \\
& \left.\exp \left\{n\left(\left|B_{1}+n x\right|-n|x|-2 L(1,-n x)-\frac{1}{2} n\right)\right\}\right]^{\frac{1}{2}} \\
\leq & E\left[\exp \left\{6 n L(1,-n x)+n\left|B_{1}+n x\right|-n^{2}\left(|x|+\frac{1}{2}\right)\right\}\right]^{\frac{1}{2}}
\end{aligned}
$$

we obtain that

$$
\begin{aligned}
I_{2}^{2}\left(n, t, t^{\prime}\right) & \leq \int_{0}^{\infty} \int_{\mathbb{R}} \exp \left\{6 n y+n|z+n x|-n^{2}\left(|x|+\frac{1}{2}\right)\right\} \\
& \frac{1}{\sqrt{2 \pi}}(y+|z+n x|+|n x|) \exp \left\{-\frac{(y+|z+n x|+|n x|)^{2}}{2}\right\} d z d y
\end{aligned}
$$

Using substitution and the fact that

$$
\frac{2}{\sqrt{\pi}} \int_{r}^{\infty} e^{-v^{2}} d v \cong \frac{1}{\sqrt{\pi} r} e^{-r^{2}}
$$

for $r \rightarrow \infty$ (see e.g. [5]). We conclude that

$$
\begin{align*}
& I_{2}^{2}\left(n, t, t^{\prime}\right) \\
& \leq \frac{2}{(n(|x|-1)-1)(n(|x|-11)-1)} \exp \left\{72 n^{2}-7 n^{2}|x|-(n(|x|-11)-1)^{2}\right\} \tag{2.37}
\end{align*}
$$

for $n \geq n_{0}$ and $|x|>11$.
Combining this with the estimate in (2.36) we see that ( $\boldsymbol{C 4}$ ) is fulfilled for initial values with $|x|>11$. On the other hand the boundedness of the sign function implies the validity of the conditions (C3) and (C5) for $|x|>11$. So it follows from Theorem 2.2 that the solutions
$X_{s}^{n}, 0 \leq s \leq 1$ converge to $X_{s}, 0 \leq s \leq 1$ in $L^{2}(\mu)$ for a subsequence if $|x|>11$. Moreover, by weak compactness and the estimates in (2.36) and (2.37) we can even deduce that this convergence is in $\mathbb{D}_{1,2}$ and that

$$
D \cdot X_{s}=0, \quad 0 \leq s \leq 1 .
$$

Hence, $X_{s}, 0 \leq s \leq 1$ is a deterministic process. On the other hand, since $|x|>11$ we get that

$$
\left|X_{s}\right| \geq||x|-s| \geq 10 \text { for all } 0 \leq s \leq 1 \text {, a.e., }
$$

that is $X_{s}$ cannot hit the discontinuity point zero.
So $X$. must be a deterministic solution (i.e., $x \pm t$ ) of the ODE (1.2).
Remark 2.12. The arguments in Example 2.11 show that we may also consider drift coefficients $b$ given by e.g. step functions of the form

$$
b(x)=\sum_{i=1}^{n} \xi_{i} \chi_{\left[0, b_{i}\right)}
$$

$\xi_{i} \geq 0, \quad b_{i} \in[0, \infty], \quad i=1, \ldots, n$.

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The research of this author was supported by the European Research Council under the European Community's Seventh Framework Programme (FP7/2007-2013) / ERC grant agreement NO [228087].

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