# Time-periodic measures, random periodic orbits, and the linear response for dissipative non-autonomous stochastic differential equations 



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#### Abstract

We consider a class of dissipative stochastic differential equations (SDE's) with time-periodic coefficients in finite dimension, and the response of time-asymptotic probability measures induced by such SDE's to sufficiently regular, small perturbations of the underlying dynamics. Understanding such a response provides a systematic way to study changes of statistical observables in response to perturbations, and it is often very useful for sensitivity analysis, uncertainty quantification, and improving probabilistic predictions of nonlinear dynamical systems, especially in high dimensions. Here, we are concerned with the linear response to small perturbations in the case when the time-asymptotic probability measures are time-periodic. First, we establish sufficient conditions for the existence of stable random time-periodic orbits generated by the underlying SDE. Ergodicity of time-periodic probability measures supported on these random periodic orbits is subsequently discussed. Then, we derive the so-called fluctuation-dissipation relations which allow to describe the linear response of statistical observables to small perturbations away from the time-periodic ergodic regime in a manner which only exploits the unperturbed dynamics. The results are formulated in an abstract setting, but they apply to problems ranging from aspects of climate modelling, to molecular dynamics, to the study of approximation capacity of neural networks and robustness of their estimates.


## 1 Introduction

In many scientific applications, a systematic determination of the response of a complex nonlinear dynamical system to time-dependent perturbations is of key importance; topical examples in high-dimensional, non-autonomous and/or stochastic settings include climate models (e.g. [1,16,34,36,59,62]), statistical physics and non-equilibrium thermodynamics (e.g. [47,50,75, 82, 83]), and even neural networks (e.g. [18,22,76]). The sought response is usually quantified in terms of a change in an 'observable' expressed as a statistical/ensemble average of some functional defined on the trajectories of the underlying
dynamical system. The classical theory of linear response (e.g. [59,81]) is concerned with capturing changes in observables to sufficiently small perturbations of the original dynamics close to its statistical equilibrium. It turns out that in such a setting the response can be expressed, with some caveats, through formulas linking the external perturbations to spontaneous fluctuations and dissipation in the unperturbed time-asymptotic dynamics (e.g. [56,82,83]). The classical fluctuation-dissipation theorem (FDT) is of fundamental importance in statistical physics (e.g. [5,28,51]), and it roughly states that for systems of identical particles in statistical equilibrium, the average response to small external perturbations can be calculated through the knowledge of suitable correlation functions of the unperturbed time-asymptotic dynamics; see, for example, $[14,52]$ for some of the many applications of the FDT in the statistical physics setting.
The validity of the linear response and fluctuation-dissipation relationships for more general dynamical systems encountered, for example, in climate modelling (e.g. [59]) is an important topic which is particularly relevant for uncertainty quantification in reducedorder predictions and reduced model tuning (e.g. [16,34,61,63]). In an early influential work, Leith [55] suggested that if the climate dynamics satisfied a suitable FDT, the climate response to small external forcing could be calculated by estimating suitable statistics in the unperturbed climate. ${ }^{1}$ Climate dynamics is modelled as a forced dissipative chaotic or stochastic dynamical system which is arguably rather far from the statistical physics' setting for FDT. Nevertheless, Leith's conjecture stimulated a lot of activity in generating new theoretical formulations (e.g. [40,62]) and in designing approximate algorithms for FDT to study the climate response (e.g. [1-3,34,36-39,59,62,64]). However, despite numerous applications in autonomous and non-autonomous settings, there is little rigorous evidence supporting the validity of the linear response and FDT in the non-autonomous setting beyond the formal derivation of FDT for time-dependent stochastic systems [62].
The goal here is to provide a more rigorous justification of the linear response theory for a class of forced dissipative stochastic differential equations (SDE's) with time-periodic coefficients which induce time-periodic probability measures. Our objective is twofold:
(i) Establish sufficient conditions for the existence and ergodicity (in an appropriate sense) of time-periodic measures associated with time-asymptotic dynamics for a class of 'dissipative'SDE's (defined later in (4.13)) with time-periodic coefficients in finite dimensions.
(ii) Analyse the linear response of such SDE's in the time-periodic regime to small perturbations, and express the change in the statistical observables based on time-periodic ergodic measures via fluctuation-dissipation-type relations.

The results derived in the sequel will concern SDE's whose time-periodic measures are supported on certain stable random periodic solutions. In principle, the results discussed in the context of the linear response apply to a wider class of dynamical systems generating time-periodic measures; however, establishing conditions for the existence and ergodicity of such measures in a more general setting (for SDE's or otherwise) is not trivial and is beyond the scope of this work.

[^0]Time-periodic probability measures associated with the time-asymptotic dynamics are arguably ubiquitous in many mathematical models. In particular, seasonal and diurnal cycles in climate models due to time-periodic forcing or retarded self-interactions in neural networks provide some of the obvious candidates and highlight the need for developing the linear response framework in the time-periodic setting. It is worth stressing that the need for rigorous formulation of the linear response and FDT for dissipative stochastic dynamical systems (in line with, for example, [59,60,64,77-79]) is justified by contemporary approaches to the simulation and reduced-order modelling of highdimensional, multi-scale dynamical phenomena. For example, comprehensive models for climate change prediction or molecular dynamics simulations involve stochastic components (e.g. [4,30,58,60,77,78,90]) to mimic the effects of unresolved dynamics, while reduced-order models typically involve stochastic noise terms (e.g. [15, 21, 49,57,72]. Here, similar to $[40,62]$, the presence of noise leads to improved regularity of the problem which simplifies key aspects of the analysis compared to deterministic, dissipative nonlinear systems (e.g. [9-11,35,80]). As a consequence, we are able to focus on systems that have other important features of realistic dynamics, namely a lack of ellipticity, non-compactness of state space, and a lack of global Lipschitz continuity of the coefficients in the underlying SDE. The results established below apply to a broad class of nonlinear functionals which include common quantities of interest, such as the mean and covariance of subsets of variables.

## 2 General set-up

Our framework relies on the theory of Markovian ${ }^{2}$ random dynamical systems (RDS), which provides a geometric link between stochastic analysis and dynamical systems. This relationship was established through the discovery (e.g. [8,54]) that for sufficiently regular coefficients $b, \sigma$ the stochastic differential equation (SDE)

$$
\begin{equation*}
\mathrm{d} X_{t}=b\left(t, X_{t}\right) \mathrm{d} t+\sigma\left(t, X_{t}\right) \mathrm{d} W_{t-s}, \quad X_{s} \in \mathbb{R}^{d} \tag{2.1}
\end{equation*}
$$

generates a stochastic flow $\{\phi(t, s, \cdot, \cdot): s, t \in \mathcal{I} \subseteq \mathbb{R}, s \leqslant t\}$ of homeomorphisms on $\mathbb{R}^{d}$ such that

$$
X_{t}^{s, x}(\omega)=\phi(t, s, \omega, x), \quad \mathbb{P}-\text { a.s. }
$$

for $x=X_{s}(\omega), \omega \in \Omega$ in the Wiener space $(\Omega, \mathcal{F}, \mathbb{P})$ with $W_{t}$ an $m$-dimensional Brownian motion. For $b=b(x), \sigma=\sigma(x)$, the SDE will be called autonomous, and non-autonomous otherwise. It turns out (e.g. [7,8]) that, for an autonomous SDE (in the above sense) with sufficiently regular coefficients there exists essentially a one-to-one correspondence between the SDE and an RDS; a rough but convenient interpretation (skipping the filtration) is that in the autonomous case there exists an RDS generating the SDE, which in turn generates the stochastic flow and vice versa.
One of the key concepts relevant for the analysis of the long-time behaviour of RDS is the extension of the notion of ergodicity to the random setting (e.g. [8,12,13,24,44,68-70]). These important results are established in the regime of (random) stationary measures and (random) stationary processes, in the case when the source of time dependence is only

[^1]due to the noise process (i.e. $b=b(x), \sigma=\sigma(x)$ in (2.1) and the SDE is autonomous in the jargon established above). Over the last decade significant progress has been made in the study of the long-time behaviour of SDE's generated by time-dependent vector fields (e.g. [20,31-33, 86, 87, 89]). Based on the insight from the latter results, we shall study the ergodicity of SDE's with time-periodic coefficients in order to establish fluctuationdissipation formulas through the linear response in the random periodic regime. Our strategy is to first prove the existence of a unique time-periodic measure under certain 'dissipative' assumptions on the SDE via a version of Lyapunov second method and coupling. The standard Lyapunov second method is a well-known and powerful technique for the investigation of stability of solutions of nonlinear dynamical systems in finite and infinite dimensions. An extension of this method to an RDS generated by an autonomous SDE is essentially due to Hasminskii (e.g. [44]); subsequent extensions include applications to SDE's with random switching (e.g. [67]) and to the case of non-trivial random stationary solutions and random attractors by Schmalfuss [84]. Importantly, this method involves the study of random invariant sets (under the considered dynamics) without the need for the explicit knowledge of solutions of the underlying SDE, and it is based solely on the vector fields encoded in the coefficients of the SDE even when the drift term, i.e. $b$ in (2.1), is only locally Lipschitz continuous. However, in the present non-autonomous, time-periodic set-up, the lack of stationarity and the unavoidable skew-product structure of the underlying dynamics pose additional challenges when dealing with ergodicity of time-asymptotic probability measures. The main issue which prevents one from using the 'classical' methods (e.g. [27,44]) for proving ergodicity the random periodic regime stems from the fact that these probability measures are defined on the skew-product fibre bundle (e.g. [23]) in the space of measures on the time-extended state space and that these measures are not mixing. Here, this complication is overcome by employing an extension of the classical Krylov-Bogolyubov procedure (see, for example, [8, §1.5]) which allows for dealing with ergodicity of probability measures on appropriate Poincaré sections in the narrow topology generated by the dual of an appropriate discrete-time transition semigroup, and then 'linking' the results via the continuous-time transition semigroup induced by the SDE dynamics on the space of skew-product probability measures. In the present case, the properties of the time-periodic measures established with the help of the Lyapunov's second method for dissipative SDE's allow us to dispense with explicit assumptions on the ergodicity in the Poincaré sections which would be otherwise required.
The rest of the article is organised as follows. In the remainder of this section, we fix the notation which is frequently used in the sequel. In Sect. 3, we recap some basic results and definitions, including the notion of a Random Dynamical System (RDS) generated by an SDE in finite dimensions, and we outline the notion of a random periodic process. In Sect. 4, we first prove the existence of stable random time-periodic solutions for a class of dissipative SDE's with time-periodic coefficients, and the existence of the associated timeperiodic measures (Sect. 4.2); sufficient conditions for ergodicity of such measures (in an appropriate sense) are established in Sect.4.3. Section 5 deals with the linear response theory in the above setting. The derivation of the linear response formula in the time-periodic setting is followed by the derivation of two classes of fluctuation-dissipation relationships: the first one applies to perturbations of dynamics with time-periodic ergodic probability measures required only to exist in the unperturbed dynamics, the second one involves sim-
pler formulas but it requires persistence of time periodicity in the perturbed probability measures.

### 2.1 Function spaces

Below, we outline function spaces which are used in the sequel.
Let $(\mathcal{X}, \mathrm{d})$ be a complete separable metric space. We consider either $\mathcal{X}=\mathbb{R}^{d}$ or the space $\mathcal{X}=\mathbb{R} \times \mathbb{R}^{d}$, or a flat cylinder $\mathcal{X}=[0, \tau) \times \mathbb{R}^{d}, 0<\tau<\infty,[0, \tau) \simeq \mathbb{R} \bmod \tau$, which arises when 'lifting' the dynamics generated by a non-autonomous SDE; in this section, we use $\mathcal{X}$ for all these spaces to unify the notation. Throughout the paper, we set $\mathbb{N}_{0}:=\{0,1,2, \ldots\}$ and $\mathbb{N}_{1}:=\{1,2, \ldots\}$.

- $(\Omega, \mathcal{F}, \mathbb{P})$ denotes the Wiener probability space where $\Omega:=\mathcal{C}_{0}\left(\mathbb{R} ; \mathbb{R}^{m}\right), m \in \mathbb{N}_{1}$; i.e. the abstract sample space $\Omega$ is identified with a linear subspace of continuous functions $\mathcal{C}\left(\mathbb{R} ; \mathbb{R}^{m}\right)$ which vanish at zero. $\mathcal{F}$ is the Borel $\mathfrak{S}$-algebra on $\Omega$ generated by open subsets in the compact-open topology defined via

$$
\varrho(\omega, \hat{\omega})=\sum_{n=0}^{\infty} \frac{1}{2^{n}} \frac{\|\omega-\hat{\omega}\|_{n}}{1+\|\omega-\hat{\omega}\|_{n}}, \quad\|\omega-\hat{\omega}\|_{n}:=\sup _{t \in[-n, n]}|\omega(t)-\hat{\omega}(t)|, \quad \omega, \hat{\omega} \in \Omega
$$

with $|\cdot|$ the Euclidean norm. Finally, $\mathbb{P}$ is the Wiener measure on $\mathcal{F}$. In such a set-up the canonical Wiener process (with two-sided time) on $(\Omega, \mathcal{F})$ with values in the Borelmeasurable space $\left(\mathbb{R}^{m}, \mathcal{B}\left(\mathbb{R}^{m}\right)\right)$ is defined as $W_{t}(\omega)=\omega(t), t \in \mathbb{R}$, via the identification of $\omega \in \Omega$ with functions $\omega(\cdot) \in \mathcal{C}_{0}\left(\mathbb{R} ; \mathbb{R}^{m}\right)$; see, for example, [8, Appendix A.2] and references therein for details.

- Given the probability space $(\Omega, \mathcal{F}, \mathbb{P})$ and $\mathcal{G} \subseteq \mathcal{F}, L^{p}(\Omega, \mathcal{G}, \mathbb{P}), 1 \leqslant p<\infty$ is the space of $\mathcal{G}$-measurable random variables $X: \Omega \rightarrow \mathbb{R}^{d}$ such that $\mathbb{E}|X|^{p}:=\int_{\Omega}|X(\omega)|^{p} \mathbb{P}(\mathrm{~d} \omega)<\infty$, and equipped with the norm $\|X\|_{p}:=\left(\mathbb{E}|X|^{p}\right)^{1 / p}$.
- Given the (Borel) measurable space $(\mathcal{X}, \mathcal{B}(\mathcal{X}))$, where $\mathcal{B}(\mathcal{X})$ denotes the Borel $\mathfrak{S}$-algebra over $\mathcal{X}$,
- $\mathbb{M}(\mathcal{X})$ is the space of measurable functions on $\mathcal{X}$,
- $\mathcal{C}(\mathcal{X})$ is the space of continuous functions on $\mathcal{X}$,
- $\mathbb{M}_{\infty}(\mathcal{X})$ is the space of bounded, measurable, real-valued, scalar functions on $\mathcal{X}$, i.e.

$$
\mathbb{M}_{\infty}(\mathcal{X}):=\left\{f: \mathcal{X} \rightarrow \mathbb{R}, f \in \mathbb{M}(\mathcal{X}):\|f\|_{\infty}<\infty\right\}, \quad\|f\|_{\infty}:=\sup _{x \in \mathcal{X}}|f(x)|
$$

- $\mathcal{C}_{\infty}(\mathcal{X})$ is the space of bounded, real-valued, scalar, continuous functions on $\mathcal{X}$, i.e.

$$
\mathcal{C}_{\infty}(\mathcal{X}):=\left\{f: \mathcal{X} \rightarrow \mathbb{R}, f \in \mathcal{C}(\mathcal{X}):\|f\|_{\infty}<\infty\right\}
$$

- The space $\mathcal{C}^{l}(\mathcal{X})$ contains $l$-times continuously differentiable real-valued functions.
- The space $\mathcal{C}_{\infty}^{l}(\mathcal{X})$ contains those functions in $\mathcal{C}^{l}(\mathcal{X})$ which are bounded.
- The space $\mathcal{C}^{l, k}\left(\mathcal{X}_{1} \times \mathcal{X}_{2}\right)$ denotes the space of functions which are $\mathcal{C}^{l}$ on $\mathcal{X}_{1}$, and $\mathcal{C}^{k}$ on $\mathcal{X}_{2}$. The space $\mathcal{C}_{\infty}^{l, k}$ contains bounded functions in $\mathcal{C}^{l, k}$.
- The space of bounded, real, Lipschitz continuous functions on $(\mathcal{X}, \mathrm{d})$ is denoted by

$$
\begin{aligned}
\operatorname{Lip}_{\infty}(\mathcal{X}) & :=\left\{f \in \mathcal{C}_{\infty}(\mathcal{X}):\|f\|_{\mathrm{BL}}<\infty\right\} \\
\|f\|_{\mathrm{BL}} & :=\max \left\{\|f\|_{\infty}, \operatorname{Lip}(f)\right\}, \quad \text { and } \quad \operatorname{Lip}(f):=\sup \left\{\frac{|f(y)-f(z)|}{\mathrm{d}(y, z)}: y \neq z, y, z \in \mathcal{X}\right\}
\end{aligned}
$$

- $\tilde{\mathcal{C}}^{l}, \delta(\mathcal{X}), l \in \mathbb{N}_{0}, 0<\delta \leqslant 1$, is the Fréchet space of functions $f: \mathcal{X} \rightarrow \mathcal{X}$, whose $l$-th derivatives are $\delta$-Hölder continuous, and which is furnished with the countable family of semi-norms

$$
\begin{aligned}
& \|f\|_{l, 0, N}:=\sup _{x \in \mathcal{X}} \frac{|\langle f(x), x\rangle|}{1+|x|^{2}}+\sum_{1 \leqslant|\beta| \leqslant l} \sup _{x \in B_{N}}\left|D^{\beta} f(x)\right| \\
& \|f\|_{l, \delta ; N}:=\|f\|_{l, 0 ; N}+\sum_{|\beta|=l} \sup _{x, y \in B_{N}, x \neq y} \frac{\left|D^{\beta} f(x)-D^{\beta} f(y)\right|}{|x-y|^{\delta}},
\end{aligned}
$$

where $|\cdot|$ is the Euclidean norm and $\langle\cdot, \cdot\rangle$ the dot product on $\mathcal{X}$, and $B_{N}=\{x \in \mathcal{X}:|x| \leqslant N\}$, $N \in \mathbb{N}_{1}$, is a closed ball in $\mathcal{X}$ with radius $N$, and

$$
D^{\beta} f(x):=\frac{\partial^{|\beta|} f}{\left(\partial x_{1}\right)^{\beta_{1}} \cdots\left(\partial x_{d}\right)^{\beta_{n}}}, \quad|\beta|:=\beta_{1}+\cdots+\beta_{n}, \quad \beta_{i} \in \mathbb{N}_{0}, \quad D^{0}=\mathbb{I}
$$

denotes the Fréchet derivative; $f \in \tilde{\mathcal{C}}^{l, \delta}(\mathcal{X})$ if $\|f\|_{l, \delta ; N}<\infty$ for all finite $N \in \mathbb{N}_{1}$.

- $\tilde{\mathcal{C}}_{b}^{l, \delta}(\mathcal{X}), l \in \mathbb{N}_{0}, 0<\delta \leqslant 1$, is the space of functions $f: \mathcal{X} \rightarrow \mathcal{X}$, whose $l$-th derivatives are $\delta$-Hölder continuous with the norm

$$
\begin{aligned}
& \|f\|_{l, 0}:=\sup _{x \in \mathcal{X}} \frac{|f(x)|}{1+|x|}+\sum_{1 \leqslant|\beta| \leqslant l} \sup _{x \in \mathcal{X}}\left|D^{\beta} f(x)\right| \\
& \|f\|_{l, \delta}:=\|f\|_{l, 0}+\sum_{|\beta|=l} \sup _{x, y \in \mathcal{X}, x \neq y} \frac{\left|D^{\beta} f(x)-D^{\beta} f(y)\right|}{|x-y|^{\delta}} .
\end{aligned}
$$

Functions $f \in \tilde{\mathcal{C}}_{b}^{l, \delta}(\mathcal{X})$ are such that $\|f\|_{l, \delta}<\infty$.

- $\tilde{\mathcal{C}}_{u b}^{l, \delta}(\mathcal{X}), l \in \mathbb{N}_{0}, 0<\delta \leqslant 1$, is the space of functions $f: \mathcal{X} \rightarrow \mathcal{X}$, whose $l$-th derivatives are $\delta$-Hölder continuous with the norm

$$
\begin{aligned}
& \|f\|_{l, 0}:=\sup _{x \in \mathcal{X}}|f(x)|+\sum_{1 \leqslant|\beta| \leqslant l} \sup _{x \in \mathcal{X}}\left|D^{\beta} f(x)\right| \\
& \|f\|_{l, \delta}:=\|f\|_{l, 0}+\sum_{|\beta|=l} \sup _{x, y \in \mathcal{X}, x \neq y} \frac{\left|D^{\beta} f(x)-D^{\beta} f(y)\right|}{|x-y|^{\delta}} .
\end{aligned}
$$

Functions $f \in \tilde{\mathcal{C}}_{u b}^{l, \delta}(\mathcal{X})$ are such that $\|f\|_{l, \delta}<\infty$.

## 3 Random periodic processes

In order to facilitate subsequent derivations, we recall definitions of random dynamical systems (RDS) generated by SDE's (see, for example, $[7,8,53,54]$ ), random periodic processes (see, for example, $[20,31-33,86,87,89]$ and transition evolutions generated by SDE's (see, for example, $[8,25,26,54]$ ). We also provide an intuitive example of a random periodic solution arising in the stochastic dynamics of periodically forced FitzHugh-Nagumo
model. All definitions below are restricted to $\mathbb{R}^{d}$ but a number of them are subsequently extended (explicitly or otherwise) to the skew-product representation in $[0, \tau) \times \mathbb{R}^{d}$ which is used to deal with the non-autonomous dynamics.

Definition 3.1 (Stochastic flow $[53,54])$ Let $\phi(t, s, \omega, x) \in \mathbb{R}^{d}$, $s, t \in \mathcal{I} \subseteq \mathbb{R}, x \in \mathbb{R}^{d}$, be a random field on a probability space $(\Omega, \mathcal{F}, \mathbb{P})$. The two-parameter family $\{\phi(t, s, \cdot, \cdot)$ : $s, t \in \mathcal{I} \subseteq \mathbb{R}\}$ is called a stochastic flow of homeomorphisms if there exists a null set $\mathcal{N} \subset \Omega$ such that for any $\omega \notin \mathcal{N}$, there exists a family of continuous maps $\{\phi(t, s, \omega, \cdot): s, t \in \mathcal{I}\}$ on $\mathbb{R}^{d}$ satisfying
(i) $\phi(t, s, \omega, \cdot)=\phi(t, u, \omega, \phi(u, s, \omega, \cdot))$ holds for any $s, t, u \in \mathcal{I}$,
(ii) $\phi(s, s, \omega, \cdot)=\operatorname{id}_{\mathcal{X}}$, for all $s \in \mathcal{I}$,
(iii) the $\operatorname{map} \phi(t, s, \omega, \cdot): \mathbb{R}^{d} \rightarrow \mathbb{R}^{d}$ is a homeomorphism for any $t, s \in \mathcal{I}$.

The map $\phi(t, s, \omega, \cdot)$ is a stochastic flow ofC ${ }^{l}$-diffeomorphisms, if it is a homeomorphism and $\phi(t, s, \omega, x)$ is $l$-times continuously differentiable with respect to $x \in \mathbb{R}^{d}$ for all $s, t \in \mathcal{I} \subseteq \mathbb{R}$ and the derivatives are continuous in $(s, t, x) \in \mathcal{I} \times \mathcal{I} \times \mathbb{R}^{d}$. The stochastic flow is referred to as 'forward' for $s \leqslant t$, and as 'backward' for $t \leqslant s$. In the sequel, we will confine the discussion to $(\Omega, \mathcal{F}, \mathbb{P})$ being the Wiener space defined in Sect. 2.1.

Definition 3.2 (Filtration generated by a stochastic flow) Given a probability space $(\Omega, \mathcal{F}, \mathbb{P})$, let $\mathcal{F}_{s}^{t} \subseteq \mathcal{F}$ be the smallest $\mathfrak{S}$-algebra on $\Omega$ generated by

$$
\cap_{\varepsilon>0} \mathfrak{S}(\phi(u, v, \cdot, \cdot): s-\varepsilon \leqslant u, v \leqslant t+\varepsilon)
$$

and containing all null sets of $\mathcal{F}$. The two-parameter filtration $\left\{\mathcal{F}_{s}^{t}: s \leqslant t\right\}$ is the filtration generated by the forward stochastic flow $\{\phi(t, s, \cdot, \cdot): s, t \in \mathcal{I} \subseteq \mathbb{R}, s \leqslant t\}$ and the filtered probability space is denoted by $\left(\Omega, \mathcal{F},\left(\mathcal{F}_{s}^{t}\right)_{s \leqslant t}, \mathbb{P}\right)$.

Definition 3.3 (Transition kernel) Consider the stochastic flow $\{\phi(t, s, \cdot, \cdot): s \leqslant t\}$ induced by the $\operatorname{SDE}$ (2.1). Given the Borel-measurable space $\left(\mathbb{R}^{d}, \mathcal{B}\left(\mathbb{R}^{d}\right)\right)$, the transition probability kernel $P(s, x ; t, \cdot)$ induced by solutions of $(2.1)$ is defined by

$$
\begin{equation*}
P(s, x ; t, A)=\mathbb{P}(\{\omega \in \Omega: \phi(t, s, \omega, x) \in A\}), \quad \forall s, t \in \mathcal{I}, s \leqslant t, A \in \mathcal{B}\left(\mathbb{R}^{d}\right) \tag{3.1}
\end{equation*}
$$

The transition kernel satisfies the Chapman-Kolmogorov equation

$$
\begin{equation*}
P(s, x ; t, A)=\int_{\mathbb{R}^{d}} P(u, y ; t, A) P(s, x ; u, \mathrm{~d} y) \tag{3.2}
\end{equation*}
$$

for any $s, t, u \in \mathcal{I}, s \leqslant u \leqslant t$, and for all $x \in \mathbb{R}^{d}, A \in \mathcal{B}\left(\mathbb{R}^{d}\right)$.
Definition 3.4 (Transition evolution and its dual) Given the forward stochastic flow $\{\phi(t, s, \cdot \cdot): s, t \in \mathcal{I} ; s \leqslant t\}$ and the transition kernel (3.1) induced by the solutions of (2.1), the operator $\mathcal{P}_{s, t}: \mathbb{M}_{\infty}\left(\mathbb{R}^{d}\right) \rightarrow \mathbb{M}_{\infty}\left(\mathbb{R}^{d}\right)$ called the transition evolution is defined by

$$
\begin{equation*}
\mathcal{P}_{s, t} \varphi(x)=\int_{\mathbb{R}^{d}} \varphi(y) P(s, x ; t, \mathrm{~d} y)=\mathbb{E}[\varphi(\phi(t, s, x))], \quad \forall s, t \in \mathcal{I}, s \leqslant t, x \in \mathbb{R}^{d} \tag{3.3}
\end{equation*}
$$

where we use the shorthand notation $\mathbb{E}[\varphi(\phi(t, s, x))]:=\int_{\Omega} \varphi(\phi(t, s, \omega, x)) \mathbb{P}(\mathrm{d} \omega)$. The action of transition evolutions to arbitrary measurable functions is extended in a standard way.

For any probability measure $\mu_{s} \in \mathcal{P}\left(\mathbb{R}^{d}\right)$, $s \in \mathcal{I}$, on $\left(\mathbb{R}^{d}, \mathcal{B}\left(\mathbb{R}^{d}\right)\right)$, the $L^{2}\left(\mu_{s}\right)$ dual $\mathcal{P}_{s, t}^{*}$ of the transition evolution $\mathcal{P}_{s, t}$ is defined by

$$
\begin{equation*}
\left(\mathcal{P}_{s, t}^{*} \mu_{s}\right)(A)=\int_{\mathbb{R}^{d}} P(s, x ; t, A) \mu_{s}(\mathrm{~d} x), \quad \forall s, t \in \mathcal{I}, s \leqslant t, A \in \mathcal{B}\left(\mathbb{R}^{d}\right) \tag{3.4}
\end{equation*}
$$

Consequently, with the help of (3.2), we have for any $s, u, t \in \mathcal{I}, s \leqslant u \leqslant t$,

$$
\begin{equation*}
\mu_{t}(A)=\left(\mathcal{P}_{s, t}^{*} \mu_{s}\right)(A)=\left(\mathcal{P}_{u, t}^{*} \mathcal{P}_{s, u}^{*} \mu_{s}\right)(A)=\left(\mathcal{P}_{u, t}^{*} \mu_{u}\right)(A), \quad \forall A \in \mathcal{B}\left(\mathbb{R}^{d}\right) \tag{3.5}
\end{equation*}
$$

Theorem 3.5 (Stochastic flows generated by solutions of SDE's) Suppose that the coefficients of the $\operatorname{SDE}(2.1)$ are such that $b(\cdot, x), \sigma(\cdot, x)$ are continuous for all $x \in \mathbb{R}^{d}$, and for all $t \in \mathbb{R}, b(t, \cdot)$, $\sigma_{k}(t, \cdot) \in \tilde{\mathcal{C}}^{l, \delta}\left(\mathbb{R}^{d}\right), l \in \mathbb{N}_{0}, 0<\delta \leqslant 1$, where $\left\{\sigma_{k}\right\}_{k=1}^{m}$, denote the columns of $\sigma$. If the initial condition $X_{s}$ in (2.1) is independent of the $\mathfrak{S}$-algebra generated by $W_{t-s}(\cdot), t \geqslant s$, and $\mathbb{E}\left[\left|X_{s}\right|^{2}\right]<\infty$, there exist unique global solutions of (2.1) which generate a forward stochastic flow of homeomorphisms $(l=0)$ or $\mathcal{C}^{l}$-diffeomorphisms $(l \geqslant 1)$ on $\mathbb{R}^{d},\{\phi(t, s, \cdot, \cdot): s, t \in \mathbb{R}, s \leqslant t\}$ such that

$$
\begin{equation*}
X_{t}^{s, x}(\omega)=\phi(t, s, \omega, x), \quad \forall s, t \in \mathbb{R}, s \leqslant t, \quad x \in \mathbb{R}^{d}, \quad \mathbb{P} \text { - a.s. } \tag{3.6}
\end{equation*}
$$

and which are adapted to the filtration $\left(\mathcal{F}_{s}^{t}\right)_{s \leqslant t}$ on $(\Omega, \mathcal{F}, \mathbb{P})$, see, for example, [54, Thm 3.4.6 and $\mathbb{4} 4.7$ ] with slight modifications. If, in addition, $\mathbb{E}\left|X_{s}\right|^{p}<\infty$, for some $2 \leqslant p<\infty$, then $\mathbb{E}\left[\left|X_{t}\right|^{p}\right]<\infty, s \leqslant t<\infty$. Stronger (e.g. dissipative) growth conditions may have to be imposed on the coefficients $(b, \sigma)$ in the SDE (2.1) to guarantee the existence of the absolute moments of the solutions for all time (see, for example, Remark 4.5 and "Appendix A").

Definition 3.6 (Infinitesimal generators) Let $f \in \mathcal{C}^{1,2}\left(\mathbb{R} \times \mathbb{R}^{d}, \mathbb{R}\right)$, and $t \mapsto \phi(t, s, \omega, x)$, $s \leqslant t, \omega \in \Omega$, be a solution of the SDE in (2.1). Considering the evolution of $f(t, \phi(t, s, \omega, x))$ allows one to represent the infinitesimal generator of solutions of (2.1) through the secondorder operator ${ }^{3}$ (e.g.[54]);

$$
\begin{equation*}
\mathcal{L}_{t} f(t, x)=\partial_{t} f(t, x)+\sum_{i=1}^{d} b_{i}(t, x) \partial_{x_{i}} f(t, x)+\frac{1}{2} \sum_{i, j=1}^{d} \sum_{k=1}^{m} \sigma_{i k}(t, x) \sigma_{j k}(t, x) \partial_{x_{i} x_{j}}^{2} f(t, x), \tag{3.7}
\end{equation*}
$$

where $(b, \sigma)$ are sufficiently regular drift and diffusion coefficients in the SDE (2.1). Analogously, for $g \in \mathcal{C}^{1,2}\left(\mathbb{R} \times \mathbb{R}^{d} \times \mathbb{R}^{d}, \mathbb{R}\right)$, the infinitesimal generator of the two-point motion [54, §4.2], $t \mapsto(\phi(t, s, \omega, x), \phi(t, s, \omega, y))$, of the flow $\{\phi(t, s, \cdot, \cdot): s \leqslant t\}$ can be represented through the second-order differential operator

[^2]\[

$$
\begin{align*}
\mathcal{L}_{t}^{(2)} g(t, x, y)= & \partial_{t} g(t, x, y)+\sum_{i=1}^{d}\left(b(t, x) \partial_{x_{i}} g(t, x, y)+b(t, y) \partial_{y_{i}} g(t, x, y)\right) \\
& +\frac{1}{2} \sum_{i, j=1}^{d} \sum_{k=1}^{m}\left(\sigma_{i k}(t, x) \sigma_{j k}(t, x) \partial_{x_{i} x_{j}}^{2} g(t, x, y)\right. \\
& +\sigma_{i k}(t, x) \sigma_{j k}(t, y) \partial_{x_{i} y_{j}}^{2} g(t, x, y) \\
& \left.+\sigma_{i k}(t, y) \sigma_{j k}(t, x) \partial_{y_{i} x_{j}}^{2} g(t, x, y)+\sigma_{i k}(t, y) \sigma_{j k}(t, y) \partial_{y_{i} y_{j}}^{2} g(t, x, y)\right) \tag{3.8}
\end{align*}
$$
\]

Definition 3.7 (Random dynamical system $[7,8]$ ) Given a probability space ( $\Omega, \mathcal{F}, \mathbb{P}$ ), a measurable random dynamical system $(\mathbf{R D S})$ on $\left(\mathbb{R}^{d}, \mathcal{B}\left(\mathbb{R}^{d}\right)\right)$ over a measurable dynamical system $(\mathbf{D S}), \Theta:=\left(\Omega, \mathcal{F}, \mathbb{P},\left(\theta_{t}\right)_{t \in \mathbb{R}}\right)$, satisfying ${ }^{4} \theta \mathbb{P}=\mathbb{P}$, is a map $\Phi: \mathcal{I} \times \Omega \times \mathbb{R}^{d} \rightarrow \mathbb{R}^{d}$ such that the following hold
(a) $(t, \omega, x) \mapsto \Phi(t, \omega, x)$ is measurable for all $t \in \mathcal{I} \subseteq \mathbb{R}$,
(b) $\Phi(0, \omega, \cdot)=\mathrm{id}_{\mathbb{R}^{d}}$ for all $\omega \in \Omega$,
(c) $\Phi(t+s, \omega, \cdot)=\Phi\left(t, \theta_{s} \omega, \Phi(s, \omega, \cdot)\right)$ for all $s, t \in \mathcal{I}, \omega \in \Omega$ (cocycle property),
(d) $\Phi$ is continuous if $(t, x) \mapsto \Phi(t, \omega, x)$ is continuous for all $t \in \mathcal{I}, x \in \mathbb{R}^{d}$,
(e) $\Phi$ is smooth of class $\mathcal{C}^{l}$, if $\Phi(t, \omega, x)$ is $l$-times differentiable w.r.t. $x \in \mathbb{R}^{d}$, and the derivatives are continuous w.r.t. $(t, x) \in \mathcal{I} \times \mathbb{R}^{d}$.

The canonical filtration on $(\Omega, \mathcal{F}, \mathbb{P})$ for the RDS is generated by $\left(\theta_{t}\right)_{t \in \mathbb{R}}$.
Definition 3.8 (Canonical DS for processes with stationary increments) Consider a probability space $\left(\Omega, \mathcal{F}, \mathbb{P}_{\xi}\right)$ with the measure $\mathbb{P}_{\xi}$ on $(\Omega, \mathcal{F})$ induced by the law of a stochastic process with continuous time $\xi=\left(\xi_{t}\right)_{t \in \mathbb{R}}, \xi_{t}: \Omega \rightarrow \mathbb{R}^{d}$. A process $\xi$ is said to have stationary increments if for any $t_{0} \leqslant \cdots \leqslant t_{n}, n \in \mathbb{N}_{1}$, the distribution of $\left(\xi_{t_{1}+t}-\xi_{t_{0}+t}, \ldots, \xi_{t_{n}+t}-\xi_{t_{n-1}+t}\right)$ is independent of $t \in \mathbb{R}$; i.e.

$$
\begin{equation*}
\theta(t) \mathbb{P}_{\xi}=\mathbb{P}_{\xi} \quad \text { for all } t \in \mathbb{R} \tag{3.9}
\end{equation*}
$$

where $\left(\theta_{t}\right)_{t \in \mathbb{R}}$ is a semigroup of time shifts. The corresponding measurable dynamical system $\Theta:=\left(\Omega, \mathcal{F}, \mathbb{P}_{\xi},\left(\theta_{t}\right)_{t \in \mathbb{R}}\right)$ is called the canonical dynamical system for the process with stationary increments; see, for example, [8, Appendix A.3] for details.

Proposition 3.9 (Canonical DS for Brownian motion/Wiener process) For the Wiener probability space $(\Omega, \mathcal{F}, \mathbb{P})$ defined in Sect. 2.1, the canonical dynamical system $\Theta=\left(\Omega, \mathcal{F}, \mathbb{P},\left(\theta_{t}\right)_{t \in \mathbb{R}}\right)$ for a stochastic process with stationary increments is given by

$$
\begin{equation*}
\theta_{t}: \Omega \rightarrow \Omega, \quad \theta_{s} \omega(t)=\omega(t+s)-\omega(s), \quad \forall s, t \in \mathbb{R} \tag{3.10}
\end{equation*}
$$

so that the set $\Omega=\mathcal{C}_{0}\left(\mathbb{R}, \mathbb{R}^{m}\right)$ is invariant w.r.t. the shifts $\left(\theta_{t}\right)_{t \in \mathbb{R}}$. The canonical stochastic process $W_{t}(\omega)=\omega(t), t \in \mathbb{R}$, with stationary independent increments is the Wiener process/Brownian motion (with two-sided time) which satisfies identically

$$
\begin{equation*}
W_{t}\left(\theta_{s} \omega\right)=W_{t+s}(\omega)-W_{s}(\omega), \quad \forall s, t \in \mathbb{R} \tag{3.11}
\end{equation*}
$$

$\overline{{ }^{4} \text { Here, the notation } \theta \mathbb{P}=\mathbb{P} \text { means that } \mathbb{P}\left(\left\{\omega \in \Omega: \theta_{t} \omega \in A\right\}\right)=\mathbb{P}(\{\omega \in \Omega: \omega \in A\}), \forall A \in \mathcal{F}, t \in \mathcal{I} \text {; i.e. the }}$ semigroup $\left(\theta_{t}\right)_{t \in \mathcal{I}}, \theta_{t}: \Omega \rightarrow \Omega$, preserves the measure $\mathbb{P}$; we restrict the definition of the RDS to $\mathcal{I}=\mathbb{R}$.

Proof See [8, Appendix A.3] for an outline or, for example, [73].
Remark 3.10 In the sequel, it will be more convenient to use (3.11) in the alternative form

$$
\begin{equation*}
W_{t}(\omega)=W_{t+s}\left(\theta_{-s} \omega\right)-W_{s}\left(\theta_{-s} \omega\right), \quad \forall s, t \in \mathbb{R} \tag{3.12}
\end{equation*}
$$

Assuming suitable regularity of the coefficients of autonomous SDE's, such as those in Theorem 3.5, together with appropriate adoption of two-sided stochastic calculus, the solutions of autonomous SDE's generate ${ }^{5}$ an RDS over $\Theta$ (e.g. [7,8,29, 46,54]). We will consider the non-autonomous dynamics of the $\operatorname{SDE}$ (2.1) with time-periodic coefficients as an RDS on a suitably extended space.

### 3.1 Time-periodic setting

In the sequel, we consider non-autonomous SDE's (2.1) on $\mathbb{R}^{d}$ with time-periodic coefficients; i.e. $b(t+\tau, \cdot)=b(t, \cdot), \sigma(t+\tau, \cdot)=\sigma(t, \cdot), 0<\tau<\infty$, satisfying the conditions in Theorem 3.5 so that (2.1) has global solutions generating the forward stochastic flow $\left\{\phi(t+s, s, \cdot \cdot): s \in \mathbb{R}, t \in \mathbb{R}^{+}\right\}$such that, for all $s \in \mathbb{R}, t \in \mathbb{R}^{+}$,

$$
\begin{equation*}
\phi(t+s+\tau, s+\tau, \omega, \cdot)=\phi\left(t+s, s, \theta_{\tau} \omega, \cdot\right) \quad \mathbb{P}-\text { a.s. } \tag{3.13}
\end{equation*}
$$

The above property follows from the time periodicity of the coefficients and the uniqueness of solutions of (2.1). The relationship in (3.13) is essential for constructing an RDS on $[0, \tau) \times \mathbb{R}^{d}$ from solutions of (2.1) with time-periodic coefficients, which is important for asserting the existence and ergodicity of time-periodic measures supported on random time-periodic paths defined below.

Definition 3.11 (Random periodic path of a stochastic flow [31,32,89]) A random periodic path of period $0<\tau<\infty$ generated by a stochastic flow $\{\phi(t+s, s, \cdot, \cdot): s \in \mathbb{R}$, $\left.t \in \mathbb{R}^{+}\right\}$is a measurable function $S: \mathbb{R} \times \Omega \rightarrow \mathbb{R}^{d}$ such that for any $s \in \mathbb{R}$ the following holds

$$
\begin{equation*}
S(\tau+s, \omega)=S\left(s, \theta_{\tau} \omega\right) \quad \text { and } \quad \phi(t+s, s, \omega, S(s, \omega))=S(t+s, \omega) \quad \mathbb{P} \text { - a.s. } \quad \forall t \in \mathbb{R}^{+} . \tag{3.14}
\end{equation*}
$$

Definition 3.12 (Random periodic path of $R D S[33,89]$ ) A random periodic path of period $0<\tau<\infty$ generated by an $\operatorname{RDS}, \Phi: \mathbb{R}^{+} \times \Omega \times \mathbb{R}^{d} \rightarrow \mathbb{R}^{d}$, is a measurable function $S: \mathbb{R} \times \Omega \rightarrow \mathbb{R}^{d}$ such that for any $s \in \mathbb{R}$ and almost all $\omega \in \Omega$ the following holds

$$
\begin{equation*}
S(\tau+s, \omega)=S\left(s, \theta_{\tau} \omega\right) \quad \text { and } \quad \Phi\left(t, \theta_{s} \omega, S(s, \omega)\right)=S(t+s, \omega) \quad \forall t \in \mathbb{R}^{+} \tag{3.15}
\end{equation*}
$$

Example 3.13 Let $b: \mathbb{R}^{d} \rightarrow \mathbb{R}^{d}, d \geqslant 2$, be a globally Lipschitz vector field, and consider the deterministic flow $\left\{\psi(t, \cdot): t \in \mathbb{R}^{+}\right\}$, defined via $\psi(t, \cdot) \equiv \phi(t, 0, \cdot)$ and generated by the autonomous ODE

$$
\begin{equation*}
\frac{\mathrm{d} Y_{t}}{\mathrm{~d} t}=b\left(Y_{t}\right), \quad Y_{0}=y \in \mathbb{R}^{d}, \quad t \in \mathbb{R}^{+} \tag{3.16}
\end{equation*}
$$

[^3]Assume that there exists a periodic solution $\mathfrak{Y}: \mathbb{R} \rightarrow \mathbb{R}^{d}$ of the ODE (3.16) of period $0<\tau<\infty$,

$$
\mathfrak{Y}(\tau+s)=\mathfrak{Y}(s) \quad \text { and } \quad \psi(t, \mathfrak{Y}(s))=\mathfrak{Y}(t+s), \quad s \in \mathbb{R}, \quad t \in \mathbb{R}^{+}
$$

Consider the stochastic process $X_{t}(\omega)=\mathfrak{Y}(t)+Z_{t}(\omega)$, where $Z_{t}$ solves the following SDE

$$
\begin{equation*}
d Z_{t}=\hat{b}\left(t, Z_{t}\right) \mathrm{d} t+\hat{\sigma}\left(t, Z_{t}\right) \mathrm{d} W_{t}, \quad Z_{0}=0, \quad t \in \mathbb{R}^{+} \tag{3.17}
\end{equation*}
$$

with time-periodic coefficients

$$
\hat{b}(t, z):=b(\mathfrak{Y}(t)+z)-b(\mathfrak{Y}(t)), \quad \hat{\sigma}(t, z):=\sigma(\mathfrak{Y}(t)+z)
$$

If $\mathfrak{Z}(t, \omega)$ is a random $\tau$-periodic solution of (3.17), then $S(t, \omega)=\mathfrak{Y}(t)+\mathfrak{Z}(t, \omega)$ is a random $\tau$-periodic solution of the autonomous SDE:

$$
\mathrm{d} X_{t}=b\left(X_{t}\right) \mathrm{d} t+\sigma\left(X_{t}\right) \mathrm{d} W_{t}, X_{0}=\mathfrak{Y}(0), \quad t \in \mathbb{R}^{+}
$$

Example 3.14 (Stochastic FitzHugh-Nagumo model with periodic current) Consider the following SDE with nontrivial random periodic solutions (see [32]) which has less restrictive conditions on the drift than those considered in the sequel:

$$
\begin{equation*}
\mathrm{d} X_{t}=A X_{t} \mathrm{~d} t+b\left(t, X_{t}\right) \mathrm{d} t+\sigma(t) \mathrm{d} W_{t-s}, \quad X_{s}=x \in \mathbb{R}^{2}, \quad s, t \in \mathbb{R}, s \leqslant t \tag{3.18}
\end{equation*}
$$

where

$$
A=\binom{1-1}{a-1}, \quad b(t, x, y)=\binom{-\frac{1}{3} x^{3}+B_{1} \sin (\tau t)}{c}, \quad \sigma(t)=\left(\begin{array}{cc}
\sqrt{2 \beta^{-1}}+B_{2} \cos (\tau t) & 0 \\
0 & 0
\end{array}\right)
$$

with $a<1, \beta>0, B_{1}, B_{2}, c \in \mathbb{R}, 0<\tau<\infty$, and $W_{t}=\left(W_{t}^{1}, 0\right)^{T}$, where $W_{t}^{1}$ is a two-sided Wiener process on $\mathbb{R}$. Let $X_{t}^{s, x}(\omega)=\phi(t, s, \omega, x), s \leqslant t$, be the solution of (3.18) represented via

$$
\phi(t, s, \omega, x)=e^{A(t-s)} x+\int_{s}^{t} e^{A(t-\zeta)} b(\zeta, \phi(\zeta, s, \omega, x)) \mathrm{d} \zeta+\int_{s}^{t} e^{A(t-\zeta)} \sigma(\zeta) \mathrm{d} W_{\zeta-s}(\omega)
$$

where $x \mapsto e^{A(t-s)} x$ is the solution of the linear ODE

$$
\frac{\mathrm{d} Y_{t}}{\mathrm{~d} t}=A Y_{t}, \quad Y_{s}=y \in \mathbb{R}^{2}, \quad s, t \in \mathbb{R}, s \leqslant t
$$

Consider the projections $P^{-}: \mathbb{R}^{2} \rightarrow E^{-}, P^{+}: \mathbb{R}^{2} \rightarrow E^{+}$, where the linear subspaces are

$$
E^{-}=\operatorname{span}\left\{y \in \mathbb{R}^{2}: A y=-\lambda y\right\}, \quad E^{+}=\operatorname{span}\left\{y \in \mathbb{R}^{2}: A y=\lambda y\right\}, \lambda:=\sqrt{1-a} .
$$

The process $S(t, \omega)$ defined by

$$
\begin{aligned}
S(t, \omega)= & \int_{-\infty}^{t} e^{A(t-\zeta)} P^{-} b(\zeta, S(\zeta, \omega)) \mathrm{d} \zeta-\int_{t}^{\infty} e^{A(t-\zeta)} P^{+} b(\zeta, S(\zeta, \omega)) \mathrm{d} \zeta \\
& +\int_{-\infty}^{t} e^{A(t-\zeta)} P^{-} \sigma(\zeta) \mathrm{d} W_{\zeta-s}(\omega)-\int_{t}^{\infty} e^{A(t-\zeta)} P^{+} \sigma(\zeta) \mathrm{d} W_{\zeta-s}(\omega),
\end{aligned}
$$

is a random $2 \pi / \tau$-periodic solution of the flow generated by the SDE (3.18); see, for example, [31,32], for further details or [20] for a simpler one-dimensional example.

## 4 Time-periodic ergodic measures for dissipative SDE's

In this section, we consider a class of non-autonomous SDE's (2.1) which generate stable random periodic paths. First, in Sect. 4.2 we prove the existence of a unique stable random periodic solution for a class of 'dissipative' ${ }^{6}$ SDE's with time-periodic coefficients, and we assert the existence of time-periodic measures induced by such dynamics (Theorem 4.7). Ergodicity (in an appropriate sense, and under typical regularity conditions) of these time-periodic measures are established in Theorem 4.11 of Sect. 4.3. We conclude with an example of a periodically forced stochastic Lorenz model, which is then used in Sect. 5 to illustrate the utility of fluctuation-dissipation formulas for time-periodic measures when considering the linear response of the dynamics to small perturbations.

### 4.1 Preliminaries, definitions, and assumptions

First, we recall the notion of a time-periodic probability measure which will be needed throughout the remainder of this paper.

Definition 4.1 (Time-periodic probability measure [33]) A measure-valued map given by $t \mapsto \mu_{s+t} \in \mathcal{P}\left(\mathbb{R}^{d}\right)$ and induced by the family $\left(\mathcal{P}_{s, s+t}^{*}\right)_{t \in \mathbb{R}^{+}}, s \in \mathbb{R}$, defined in (3.4) is referred to as a time-periodic probability measure of period $0<\tau<\infty$, if the following holds for any $s \in \mathbb{R}$

$$
\begin{equation*}
\mu_{s+t}=\mathcal{P}_{s, s+t}^{*} \mu_{s} \quad \text { and } \quad \mu_{s+\tau}=\mu_{s}, \quad \forall t \in \mathbb{R}^{+} \tag{4.1}
\end{equation*}
$$

Furthermore, $\mu_{s+t} \in \mathcal{P}\left(\mathbb{R}^{d}\right)$, $s \in \mathbb{R}, t \in \mathbb{R}^{+}$, is called a time-periodic measure with the minimal (or fundamental) period $\tau$, if $\tau$ is the smallest strictly positive number such that (4.1) holds. ${ }^{7}$

Proposition 4.2 Let $S: \mathbb{R} \times \Omega \rightarrow \mathbb{R}^{d}$ be a random periodic path (3.14) of a stochastic flow $\left\{\phi(t+s, s, \cdot, \cdot): s \in \mathbb{R}, t \in \mathbb{R}^{+}\right\}$on $\left(\mathbb{R}^{d}, \mathcal{B}\left(\mathbb{R}^{d}\right)\right)$ and consider a family of probability measures

$$
\mu_{s+t}(A):=\mathbb{P}(\{\omega: S(s+t, \omega) \in A\}), \quad \forall s \in \mathbb{R}, t \in \mathbb{R}^{+}, A \in \mathcal{B}\left(\mathbb{R}^{d}\right)
$$

Then, the family $\left(\mu_{s+t}\right)_{s \in \mathbb{R}, t \in \mathbb{R}^{+}}$consists of $\tau$-periodic probability measures on $\mathbb{R}^{d}$.
Proof This follows by a direct calculation combined with the properties of a random periodic path (3.14), since for all $s \in \mathbb{R}, t \in \mathbb{R}^{+}, A \in \mathcal{B}\left(\mathbb{R}^{d}\right)$, we have

$$
\begin{aligned}
\mu_{s+\tau}(A)=\mathbb{P}(\{\omega: S(s+\tau, \omega) \in A\}) & =\mathbb{P}\left(\left\{\omega: S\left(s, \theta_{\tau} \omega\right) \in A\right\}\right) \\
& =\mathbb{P}(\{\omega: S(s, \omega) \in A\})=\mu_{s}(A)
\end{aligned}
$$

The above results will be generalised to the dynamics in the extended state space in Sect. 4.1.1.

[^4]
### 4.1.1 Dynamics on the extended state space

A useful way of examining ergodicity of time-periodic measures induced by nonautonomous SDE's with time-periodic coefficients of period $\tau$ is to lift the original dynamics from $\mathbb{R}^{d}$ to the extended state space $[0, \tau) \times \mathbb{R}^{d},[0, \tau) \simeq \mathbb{R} \bmod \tau$, so that the resulting 'lifted' SDE is autonomous. Such a representation of the original dynamics does not necessarily simplify the formulation of the problem but the flows of the lifted solutions generate a cocycle ${ }^{8}$ in the skew-product variables on $[0, \tau) \times \mathbb{R}^{d}$; we refer to this extended state space as the 'flat cylinder'. Then, the lifted random periodic paths (3.14) of the stochastic flow induced by the non-autonomous SDE (2.1) can be associated with random periodic paths (satisfying (3.15)) of an RDS (see Definition 3.7) generated by the lifted flow in the skew-product representation on $[0, \tau) \times \mathbb{R}^{d}$; this fact allows to prove ergodicity (in an appropriate sense) of time-periodic measures supported on the random periodic paths on the fibre bundle ${ }^{9}$ on $\mathcal{P}\left([0, \tau) \times \mathbb{R}^{d}\right)$.

To this end, consider the solutions of the SDE (2.1) satisfying the conditions of Theorem 3.5 and assume that the coefficients of (2.1) are time-periodic with period $0<\tau<\infty$; we recast the solutions of (2.1) as an extended process $\tilde{X}_{t}(\omega)=\left(t, X_{t}^{s, x}(\omega)\right)^{T}$ in the skewproduct representation on $\mathbb{R} \times \mathbb{R}^{d}$ satisfying

$$
\begin{equation*}
\mathrm{d} \tilde{X}_{t}=\tilde{b}\left(\tilde{X}_{t}\right) \mathrm{d} t+\tilde{\sigma}\left(\tilde{X}_{t}\right) \mathrm{d} \tilde{W}_{t-s}, \quad \tilde{X}_{s}=(s, x) \in \mathbb{R} \times \mathbb{R}^{d}, \quad s \leqslant t \tag{4.2}
\end{equation*}
$$

where $\tilde{W}_{t-s}(\omega)=\left(0, W_{t-s}(\omega)\right), \omega \in \Omega$, and $W_{t-s}$ is the $m$-dimensional Brownian motion for the two-sided time (see Sect. 2.1 or [8]), and $\tilde{b}: \mathbb{R}^{d+1} \rightarrow \mathbb{R}^{d+1}$, $\tilde{\sigma}: \mathbb{R}^{d+1} \rightarrow \mathbb{R}^{(d+1) \times(m+1)}$, so that

$$
d\binom{\zeta_{t}}{X_{t}}=\binom{1}{b\left(\zeta_{t}, X_{t}^{s, x}\right)} \mathrm{d} t+\left(\begin{array}{cc}
0 & 0  \tag{4.3}\\
0 & \sigma\left(\zeta_{t}, X_{t}^{s, x}\right)
\end{array}\right) \mathrm{d} \tilde{W}_{t-s}, \quad s \leqslant t
$$

The dynamics in (4.2) or (4.3) can be represented in a more convenient form for the subsequent derivations by setting $t \rightarrow t+s$, so that

$$
\begin{equation*}
\mathrm{d} \tilde{X}_{t+s}=\tilde{b}\left(\tilde{X}_{t+s}\right) \mathrm{d} t+\tilde{\sigma}\left(\tilde{X}_{t+s}\right) \mathrm{d} \tilde{W}_{t+s}, \quad \tilde{X}_{s}=(s, x) \in \mathbb{R} \times \mathbb{R}^{d}, t \in \mathbb{R}^{+} \tag{4.4}
\end{equation*}
$$

where $\tilde{W}_{t+s}=\tilde{W}_{t+s}\left(\theta_{-s} \omega\right)$ is the Brownian motion satisfying (3.12). Finally, given the form of the coefficients $\tilde{b}, \tilde{\sigma}$, it is convenient to consider the dynamics induced by (4.4) on the flat cylinder $[0, \tau) \times \mathbb{R}^{d}$, where $[0, \tau) \simeq \mathbb{R} \bmod \tau$.

The RDS associated with the lifted dynamics (4.4) is generated in the skew-product representation (see, for example, $[8,23]$ ) on $[0, \tau) \times \mathbb{R}^{d}$ via

$$
\begin{equation*}
\tilde{\Phi}(t, \omega, \tilde{x}):=\left(t+s \bmod \tau, \quad \phi\left(t+s, s, \theta_{-s} \omega, x\right)\right) \quad \forall \tilde{x}:=(s, x) \in[0, \tau) \times \mathbb{R}^{d}, t \in \mathbb{R}^{+} \tag{4.5}
\end{equation*}
$$

The cocycle property ${ }^{10}$ of $\tilde{\Phi}$ in (4.5), i.e. $\tilde{\Phi}(t+r, \omega, \cdot)=\tilde{\Phi}\left(t, \theta_{r} \omega, \tilde{\Phi}(r, \omega, \cdot)\right)$ for all $r, t \in \mathbb{R}^{+}$, and a.a. $\omega \in \Omega$, can be verified by recalling that $t+r \bmod \tau=t+r-k \tau$, where

[^5]$k=\left\lfloor\frac{t+r}{\tau}\right\rfloor$, and utilising (3.13). Note that, unless (2.1) is autonomous, $\left\{\phi\left(t+s, s, \theta_{-s} \omega, \cdot\right)\right.$ : $\left.s \in \mathbb{R}, t \in \mathbb{R}^{+}\right\}$does not have the cocycle property, and hence, it does not generate an RDS on $\mathbb{R}^{d}$. The RDS representation of the non-autonomous dynamics of the $\operatorname{SDE}$ (2.1) will be useful in Sect. 4.3 when considering the ergodicity of measures supported on random periodic paths, and in the discussion of the linear response in Sect. 5.3.
The transition kernel and transition evolutions on $[0, \tau) \times \mathbb{R}^{d}$ are constructed as follows. For any $\tilde{x}:=(s, x) \in[0, \tau) \times \mathbb{R}^{d}, t \in \mathbb{R}^{+}$, and $\tilde{A} \in \mathcal{B}([0, \tau)) \otimes \mathcal{B}\left(\mathbb{R}^{d}\right)$, the transition kernel $\tilde{P}(\tilde{x} ; t, \cdot)$ associated with $\tilde{\Phi}$ is given by (see Definition 3.3)
\[

$$
\begin{align*}
\tilde{P}(\tilde{x} ; t, \tilde{A}) & :=\mathbb{P}(\{\omega: \tilde{\Phi}(t, \omega, \tilde{x}) \in \tilde{A}\}) \\
& =\mathbb{P}\left(\left\{\omega:\left(t+s \bmod \tau, \phi\left(t+s, s, \theta_{-s} \omega, x\right)\right) \in \mathcal{J} \times A\right\}\right) \\
& =\delta_{(t+s \bmod \tau)}(\mathcal{J}) \otimes P(s, x ; t+s, A) \tag{4.6}
\end{align*}
$$
\]

for all $\tilde{A} \equiv \mathcal{J} \times A \in \mathcal{B}([0, \tau)) \otimes \mathcal{B}\left(\mathbb{R}^{d}\right)$ and the transition kernel $P$ defined in (3.1).
The transition evolution $\left(\tilde{\mathcal{P}}_{t}\right)_{t \in \mathbb{R}^{+}}$induced by $\tilde{\Phi}$ and its dual $\left(\tilde{\mathcal{P}}_{t}^{*}\right)_{t \in \mathbb{R}^{+}}$are given by (cf. (3.3))

$$
\begin{align*}
\tilde{\mathcal{P}}_{t} \varphi(\tilde{x}) & :=\int_{[0, \tau) \times \mathbb{R}^{d}} \varphi(\tilde{y}) \tilde{P}(\tilde{x} ; t, \mathrm{~d} \tilde{y}), & \forall \varphi \in \mathbb{M}_{\infty}\left([0, \tau) \times \mathbb{R}^{d}\right),  \tag{4.7}\\
\tilde{\mu}_{t+r}(\tilde{A})=\left(\tilde{\mathcal{P}}_{t}^{*} \tilde{\mu}_{r}\right)(\tilde{A}) & :=\int_{[0, \tau) \times \mathbb{R}^{d}} \tilde{P}(\tilde{x} ; t, \tilde{A}) \tilde{\mu}_{r}(\mathrm{~d} \tilde{x}), & \forall \tilde{\mu}_{r} \in \mathcal{P}\left([0, \tau) \times \mathbb{R}^{d}\right), r \in \mathbb{R}^{+}, \tag{4.8}
\end{align*}
$$

with the short-hand notation $\tilde{\mu}_{r}(\mathrm{~d} \tilde{x})=\delta_{(r \bmod \tau)}(s) \mathrm{d} s \otimes \mu_{r}(\mathrm{~d} x)$ for probability measures in the skew-product fibre bundle on $\mathcal{P}\left([0, \tau) \times \mathbb{R}^{d}\right)$, where $\tilde{\mu}_{r} \in \mathcal{P}\left([0, \tau) \times \mathbb{R}^{d}\right)$ and $\mu_{r} \in \mathcal{P}\left(\mathbb{R}^{d}\right)$; see, for example, [23] for more details concerning the structure of skewproduct fibre bundles on spaces of probability measures. Extension of (4.7) to functions $\mathbb{M}\left([0, \tau) \times \mathbb{R}^{d}\right)$ can be carried out in a standard way.

Lemma 4.3 The families of transition evolutions $\left(\tilde{\mathcal{P}}_{t}\right)_{t \in \mathbb{R}^{+}}$and $\left(\tilde{\mathcal{P}}_{t}^{*}\right)_{t \in \mathbb{R}^{+}}$possess a semigroup structure. In particular, for $\tilde{\mu}_{t}=\delta_{(t \bmod \tau)} \otimes \mu_{t}$ in the skew-product fibre bundle on $\mathcal{P}\left([0, \tau) \times \mathbb{R}^{d}\right)$ the following holds

$$
\tilde{\mu}_{t+r+u}=\tilde{\mathcal{P}}_{t+r}^{*} \tilde{\mu}_{u}=\tilde{\mathcal{P}}_{t}^{*}\left(\tilde{\mathcal{P}}_{r}^{*} \tilde{\mu}_{u}\right)=\tilde{\mathcal{P}}_{t}^{*} \tilde{\mu}_{r+u} \quad \forall r, t, u \in \mathbb{R}^{+}
$$

If the $R D S\left\{\tilde{\Phi}(t, \cdot, \cdot): t \in \mathbb{R}^{+}\right\}$on $[0, \tau) \times \mathbb{R}^{d}$ has a random periodic path $t \rightarrow \tilde{S}(t, \omega)$ of period $0<\tau<\infty$, where $\tilde{S}(t, \omega)=(t \bmod \tau, S(t, \omega)), t \in \mathbb{R}, \omega \in \Omega$, and $t \rightarrow S(t, \omega)$ is a random periodic path of $\left\{\phi(t+s, s, \cdot, \cdot): t \in \mathbb{R}^{+}\right\}$on $\mathbb{R}^{d}$, then all probability measures in the family $\left(\tilde{\mu}_{t}\right)_{t \in \mathbb{R}^{+}}, \tilde{\mu}_{t} \in \mathcal{P}\left([0, \tau) \times \mathbb{R}^{d}\right)$, supported on such a path are $\tau$-periodic, i.e.

$$
\tilde{\mu}_{t+r}=\tilde{\mathcal{P}}_{r}^{*} \tilde{\mu}_{t}, \quad \tilde{\mu}_{t+\tau}=\tilde{\mu}_{t} \quad \forall t \in \mathbb{R}^{+},
$$

and

$$
\tilde{\mu}_{t}(\tilde{A})=\mu_{t}\left(A_{t}\right), \quad \forall \tilde{A} \in \mathcal{B}([0, \tau)) \otimes \mathcal{B}\left(\mathbb{R}^{d}\right), \quad A_{t}=\left\{x \in \mathbb{R}^{d}:(t \bmod \tau, x) \in \tilde{A}\right\} .
$$

Moreover, every such $\tau$-periodic measure is invariant under the action of the discrete dynamics induced by $\left(\tilde{\mathcal{P}}_{n \tau}^{*}\right)_{n \in \mathbb{N}_{0}}$, i.e.

$$
\tilde{\mathcal{P}}_{n \tau}^{*} \tilde{\mu}_{t}=\tilde{\mu}_{t} \quad \forall n \in \mathbb{N}_{0}, t \in \mathbb{R}^{+}
$$

Proof The first claim is a direct consequence of (4.8), and the proof follows either by using the cocycle property of $\tilde{\Phi}$ in the first line of (4.6) or by utilising the Chapman-Kolmogorov equation (3.2) for $P$ in the last line of (4.6).
Regarding the second claim, consider measures supported on the random periodic path $\tilde{S}$

$$
\begin{equation*}
\tilde{\mu}_{t+r}(\tilde{A}):=\mathbb{P}(\{\omega: \tilde{S}(t+r, \omega) \in \tilde{A}\}), \quad \tilde{A} \in \mathcal{B}([0, \tau)) \otimes \mathcal{B}\left(\mathbb{R}^{d}\right) \tag{4.9}
\end{equation*}
$$

Since $\tilde{S}$ is a random periodic path of the $\operatorname{RDS} \tilde{\Phi}$, we have for all $\tilde{A} \in \mathcal{B}([0, \tau)) \otimes \mathcal{B}\left(\mathbb{R}^{d}\right)$ that

$$
\tilde{\mu}_{t+r}(\tilde{A})=\mathbb{P}(\{\omega: \tilde{S}(t+r, \omega) \in \tilde{A}\})=\mathbb{P}\left(\left\{\omega: \tilde{\Phi}\left(r, \theta_{t} \omega, \tilde{S}(t, \omega)\right) \in \tilde{A}\right\}\right)=\tilde{\mathcal{P}}_{r}^{*} \tilde{\mu}_{t}(\tilde{A}),
$$

for all $r, t \in \mathbb{R}^{+}$by the general properties the of the random periodic path (3.15); this could also be obtained directly from (4.8) by using the invariance of $\tilde{S}$ under the action of $\tilde{\Phi}$. Moreover,

$$
\tilde{\mu}_{t+\tau}(\tilde{A})=\mathbb{P}(\{\omega: \tilde{S}(t+\tau, \omega) \in \tilde{A}\})=\mathbb{P}\left(\left\{\omega: \tilde{S}\left(t, \theta_{\tau} \omega\right) \in \tilde{A}\right\}\right)=\tilde{\mu}_{t}(\tilde{A})
$$

by the property (3.15). Thus, $\tilde{\mu}_{t}$ is a $\tau$-periodic measure for the $\operatorname{RDS}\left\{\tilde{\Phi}(t, \cdot, \cdot): t \in \mathbb{R}^{+}\right\}$ on $[0, \tau) \times \mathbb{R}^{d}$ which is supported on the random periodic path $\tilde{S}$.
The last two claims are simple consequences of the properties established above and the skew-product structure of probability measures supported on random periodic paths.

In the following sections, after outlining some general assumptions, we will investigate the existence and uniqueness of stable random periodic paths of the $\operatorname{RDS}\{\tilde{\Phi}(t, \cdot, \cdot)$ : $\left.t \in \mathbb{R}^{+}\right\}$, and we will prove the ergodicity of probability measures associated with the dynamics of the skew-product lift (4.4) of the dynamics in (2.1) under some standard regularity assumptions.

### 4.1.2 Assumptions

Throughout, we assume that the SDE (2.1) with time-periodic coefficients of period $0<\tau<\infty$ satisfies the conditions of Theorem 3.5, so that (2.1) has global solutions.

In order to establish the existence of stable random periodic paths in Sect. 4.2, we will require the following assumption:

Assumption 4.4 Let $V \in \mathcal{C}^{1,2}\left(\mathbb{R} \times \mathbb{R}^{d} ; \mathbb{R}^{+}\right)$s.t. $V(t, 0)=0$ for all $t \in \mathbb{R}$, satisfy the following:
(i) There exist $\lambda \in L^{1}(\mathbb{R} ; \mathrm{d} t)$, and a constant $\mathfrak{C} \geqslant 1$, such that for some $1<p<\infty$ and all $\xi, \eta \in L^{p}\left(\Omega, \mathcal{F}_{-\infty}^{t}, \mathbb{P}\right)$, we have

$$
\left\{\begin{array}{l}
\mathbb{E}|\xi|^{p} \leqslant \mathbb{E}[V(t, \xi)] \leqslant \mathfrak{C} \mathbb{E}|\xi|^{p}<\infty  \tag{4.10}\\
\mathbb{E}\left[\mathcal{L}^{(2)} V(t, \xi-\eta)\right] \leqslant \lambda(t) \mathbb{E}[V(t, \xi-\eta)]
\end{array}\right.
$$

where $\mathcal{L}^{(2)}$ is the two-point generator defined in (3.8) and associated with the SDE (2.1).
(ii) There exists $\bar{\lambda}>0$ such that

$$
\begin{equation*}
\limsup _{(t-s) \rightarrow \infty} \frac{1}{t-s} \int_{s}^{t} \lambda(u) \mathrm{d} u<-\bar{\lambda}<0 \tag{4.11}
\end{equation*}
$$

(iii) For the one-point motion $t \mapsto \phi(t, s, \omega, \xi)$ induced by (2.1) for $\omega \in \Omega, \xi \in \mathbb{R}^{d}$, and $s \leqslant t$, there exists $0<\mathfrak{D}<\infty$ independent of $s, t \in \mathbb{R}$ such that ${ }^{11}$ for all $\xi \in L^{p}\left(\Omega, \mathcal{F}_{-\infty}^{s}, \mathbb{P}\right)$

$$
\begin{equation*}
\limsup _{(t-s) \rightarrow \infty} \mathbb{E}[V(t, \phi(t, s, \xi)-\xi)] \leqslant \mathfrak{D} \tag{4.12}
\end{equation*}
$$

where $\mathbb{E}[V(\phi(t, s, \xi))]:=\int_{\Omega} V(\phi(t, s, \omega, \xi)) \mathbb{P}(\mathrm{d} \omega)$.
As pointed out later (Remark 4.12 in Sect. 4.3), this assumption is not strictly required for proving ergodicity of $\tau$-periodic probability measures. However, without showing the existence of random periodic paths (in this case, stable random periodic paths), the existence of $\tau$-periodic measures $\tilde{\mu}_{t} \in \mathcal{P}\left([0, \tau) \times \mathbb{R}^{d}\right)$ would have to be assumed a priori alongside the ergodicity of $\tilde{\mu}_{t}$ for all fixed $t \in[0, \tau)$ with respect to the discrete transition evolution $\left(\tilde{\mathcal{P}}_{n \tau}^{*}\right)_{n \in \mathbb{N}_{0}}$, as done in [33].

## Remark 4.5

(a) An important class of coefficients satisfying Assumption 4.4, which yield global solutions of (2.1) are specified in "Appendix A". In particular, we might take $b(t, \cdot) \in \tilde{\mathcal{C}}^{1, \delta}\left(\mathbb{R}^{d}\right)$ and $\sigma_{k}(t, \cdot) \in \tilde{\mathcal{C}}_{b}^{1, \delta}\left(\mathbb{R}^{d}\right), 0<\delta \leqslant 1, k=1, \ldots, m$, satisfying the following 'dissipative' condition

$$
\begin{equation*}
\langle b(t, x), x\rangle \leqslant L_{b_{1}}(t)-L_{b_{2}}(t)|x|^{2}, \quad\|\sigma(t, x)\|_{\mathrm{HS}}^{2} \leqslant L_{\sigma}(t)\left(1+|x|^{2}\right) \tag{4.13}
\end{equation*}
$$

where $L_{b_{1}}, L_{b_{2}}, L_{\sigma} \in \mathcal{C}_{\infty}\left(\mathbb{R}, \mathbb{R}^{+}\right)$. Here, $\langle\cdot, \cdot\rangle$ denotes the dot product on $\mathbb{R}^{d}$ and $\|\cdot\|_{\text {HS }}$ denotes the Hilbert-Schmidt norm (aka Frobenius norm) defined by $\|A\|_{\mathrm{HS}}^{2}=\operatorname{trace}\left(A A^{T}\right)$. Condition (4.12) is satisfied for (4.13) when (see Lemma A. 1 in "Appendix A")

$$
\begin{equation*}
\inf _{t \in \mathbb{R}}\left(L_{b_{2}}(t)-2^{\frac{p}{2}-1} L_{b_{1}}(t)-\frac{1}{2}\left(2^{\frac{p}{2}-1}+1\right) L_{\sigma}(t)(p-1)\right)>0 \tag{4.14}
\end{equation*}
$$

and it also leads to the global existence of the $p$-th absolute moment of the law of the associated SDE; tighter bounds can be obtained for $p=2,3$ as shown in Proposition A. 2 in "Appendix A". Condition (4.12) is reminiscent of the Haśminskiitype regularity condition [44] for the existence and uniqueness of global solutions of SDE's; sufficient conditions for verification of Haśminskii's conditions require the existence of real-valued functions $L_{b}(\cdot), L_{\sigma}(\cdot) \in \mathcal{C}_{\infty}\left(\mathbb{R} ; \mathbb{R}^{+}\right)$such that

$$
\begin{equation*}
\langle b(t, x), x\rangle \leqslant L_{b}(t)\left(1+|x|^{2}\right), \quad\|\sigma(t, x)\|_{\mathrm{HS}}^{2} \leqslant L_{\sigma}(t)\left(1+|x|^{2}\right) \tag{4.15}
\end{equation*}
$$

[^6]Coefficients satisfying (4.13) also satisfy (4.15), since for some $L_{b} \in \mathcal{C}_{\infty}\left(\mathbb{R} ; \mathbb{R}^{+}\right)$we have

$$
L_{b_{1}}(t)-L_{b_{2}}|x|^{2} \leqslant L_{b}\left(1+|x|^{2}\right)
$$

(b) Construction of the Lyapunov function $V$ satisfying Assumption 4.4 is often not straightforward. However, one can construct (e.g. [44,45,65]) a polynomial Lyapunov function growing at infinity as $|x|^{2 N}, N \in \mathbb{N}_{1}$, for a broad class of SDE's whose coefficients $b(\cdot, x), \sigma(\cdot, x)$ are continuous, and $b(t, \cdot) \in \tilde{\mathcal{C}}^{1, \delta}\left(\mathbb{R}^{d}\right)$, $\left\{\sigma_{k}(t, \cdot)\right\}_{1 \leqslant k \leqslant m} \in \tilde{\mathcal{C}}_{b}^{1, \delta}\left(\mathbb{R}^{d}\right)$ are such that

$$
\left\{\begin{array}{l}
\langle b(t, x)-b(t, y), x-y\rangle \leqslant-K_{t}|x-y|^{2}  \tag{4.16}\\
\|\sigma(t, x)-\sigma(t, y)\|_{\mathrm{HS}} \leqslant L_{t}|x-y| \\
\sup _{t \in \mathbb{R}}\left\{|b(t, 0)|+\|\sigma(t, 0)\|_{\mathrm{HS}}\right\}<\infty
\end{array}\right.
$$

where $0<L_{t}, K_{t}<\infty$, and

$$
\begin{equation*}
\limsup _{(t-s) \rightarrow \infty} \frac{1}{t-s} \int_{s}^{t} \lambda(u) \mathrm{d} u<0 \tag{4.17}
\end{equation*}
$$

with $\lambda(t)=-K_{t}+\frac{(p-1)}{2} p L_{t}^{2}$ for some $1<p<\infty$. The function $K_{t}$ is defined by

$$
K_{t}=\liminf _{R \rightarrow \infty} K_{t}(R),
$$

where $K_{t}: \mathbb{R} \rightarrow \mathbb{R}$ is a Borel function defined by

$$
K_{t}(R)=\inf \left\{-\frac{\langle b(t, x)-b(t, y), x-y\rangle}{|x-y|^{2}}: \quad|x-y|=R\right\}
$$

Many important classes of SDE's driven Levy processes (including the Brownian motion) satisfy the dissipative conditions (4.16)-(4.17); see $[44,45,65]$ for more details.

In order to study the ergodicity of $\tau$-periodic measures, we will require variants of the following standard conditions (e.g. [43]) to be satisfied:
(i) Relative compactness property of the transition kernel $P$ in (3.3).
(ii) Irreducibility of the transition kernel.
(iii) Strong Feller property ${ }^{12}$ of the transition evolution $\left(\mathcal{P}_{s, t}\right)_{t \geqslant s}$ (3.3).

Thus, we will require the following version of the Hörmander condition (e.g. [66,74]) in Sect. 4.3 in addition to Assumption 4.4:

Assumption 4.6 Denote by $\sigma_{k}, 1 \leqslant k \leqslant m$, the columns of $\sigma$ in (2.1), and assume that the following are satisfied for all $t \in \mathbb{R}$ :
(i) $b(t, \cdot) \in \tilde{\mathcal{C}}^{\infty}\left(\mathbb{R}^{d}\right)$ and $t \mapsto b(t, \cdot)$ is differentiable.

[^7](ii) $\sigma_{k}(t, \cdot) \in \tilde{\mathcal{C}}_{b}^{\infty}\left(\mathbb{R}^{d}\right), t \mapsto \sigma_{k}(t, x)$ is differentiable, and
\[

$$
\begin{equation*}
\left|\partial_{t} D_{x}^{\beta} \sigma_{k}(t, x)\right| \leqslant C<\infty, \quad(t, x) \in \mathbb{R} \times \mathbb{R}^{d} \tag{4.18}
\end{equation*}
$$

\]

for every multi-index $\beta$.
(iii) $\operatorname{Lie}\left(\sigma_{1}(t, \cdot), \cdots, \sigma_{m}(t, \cdot)\right)=\mathbb{R}^{d}$, for all $t \in \mathcal{I}$, where

$$
\operatorname{Lie}\left(\sigma_{1}(t, x), \cdots, \sigma_{m}(t, x)\right):=\operatorname{span}\left\{\sigma_{i},\left[\sigma_{i}, \sigma_{j}\right],\left[\sigma_{i},\left[\sigma_{j}, \sigma_{k}\right]\right], \cdots, 1 \leqslant i, j, k \leqslant m\right\}
$$

and $[F, G]$ is the Lie bracket between the vector fields $F$ and $G$ defined by

$$
[F, G](t, x):=D_{x} G(t, x) F(t, x)-D_{x} F(t, x) G(t, x) .
$$

### 4.2 Existence and uniqueness of time-periodic measures on stable random periodic paths

Given the preliminary results and assumptions outlined in Sect. 4.1, we have the following result on the existence of a $\tau$-periodic measure (Definition 4.1) for the lifted SDE in (4.2).

Theorem 4.7 Consider the forward stochastic flow $\{\phi(t, s, \cdot, \cdot): s, t \in \mathbb{R}, s \leqslant t\}$ generated by the SDE in (2.1) with time-periodic coefficients of period $0<\tau<\infty$, and satisfying the conditions of Theorem 3.5. If Assumption 4.4 holds, there exists a family $\left(\tilde{\mu}_{t}\right)_{t \in \mathbb{R}^{+}}$of $\tau$-periodic measures, $\tilde{\mu}_{t}=\delta_{(t \bmod \tau)} \otimes \mu_{t}, \mu_{t} \in \mathcal{P}\left(\mathbb{R}^{d}\right), \tilde{\mu}_{t} \in \mathcal{P}\left([0, \tau) \times \mathbb{R}^{d}\right)$, given by

$$
\begin{equation*}
\tilde{\mu}_{t}(\tilde{A}):=\mathbb{P}(\{\omega: \tilde{S}(t, \omega) \in \tilde{A}\}), \quad t \in \mathbb{R}^{+}, \tilde{A} \in \mathcal{B}([0, \tau]) \otimes \mathcal{B}\left(\mathbb{R}^{d}\right) \tag{4.19}
\end{equation*}
$$

which are supported on a unique random periodic path $\tilde{S}$ of the $R D S\left\{\tilde{\Phi}(t, \cdot, \cdot): t \in \mathbb{R}^{+}\right\}$ with $\tilde{\Phi}$ in (4.5) generated in the skew-product variables on $[0, \tau) \times \mathbb{R}^{d}$.

Proof First, for $\xi \in L^{p}\left(\Omega, \mathcal{F}_{-\infty}^{s}, \mathbb{P}\right), 1<p<\infty$, where $\mathcal{F}_{-\infty}^{s}:=\bigvee_{r \leqslant s} \mathcal{F}_{r}^{s}$, we show that $\{\phi(t, s, \omega, \xi): s, t \in \mathbb{R}, s \leqslant t\}$ converges to a random process $S(t, \omega) \in \mathbb{R}^{d}$ almost surely as $s \rightarrow-\infty$, and that $S(t, \omega)$ is bounded and independent of $\xi$. Next, we show that $t \mapsto S(t, \omega)$ is a unique stable random periodic path of period $0<\tau<\infty$ for $\{\phi(t, s, \omega, \cdot): s, t \in \mathbb{R}, s \leqslant t\}$. Finally, we conclude that the law of the random periodic path $\tilde{S}(t, \omega)=(t \bmod \tau, S(t, \omega))$ generates a $\tau$-periodic measure for the RDS generated by $\tilde{\Phi}$ on the flat cylinder $[0, \tau) \times \mathbb{R}^{d},[0, \tau) \simeq \mathbb{R} \bmod \tau$.
Existence of random periodic paths for the stochastic flow $\phi$. Set $\xi, \eta \in \mathbb{R}^{d}$ to be random variables on the filtered probability space $\left(\Omega, \mathcal{F}_{-\infty}^{s}, \mathbb{P}\right)$, s.t. $\xi, \eta \in L^{p}\left(\Omega, \mathcal{F}_{-\infty}^{s}, \mathbb{P}\right)$. Then, by Itô formula (e.g. Theorem 4.2.4 in [54] or Theorem 8.1 in [53]) and Assumption 4.4 we have for $s \leqslant t$

$$
\begin{aligned}
d \mathbb{E}[V(t, \phi(t, s, \xi)-\phi(t, s, \eta))] & =\mathbb{E}\left[\mathcal{L}^{(2)} V(t, \phi(u, s, \xi)-\phi(t, s, \eta))\right] \mathrm{d} t \\
& \leqslant \lambda(t) \mathbb{E}[V(t, \phi(t, s, \xi)-\phi(u, s, \eta))] \mathrm{d} t
\end{aligned}
$$

where $\mathbb{E}[V(t, \phi(t, s, \xi)-\phi(u, s, \eta))]:=\int_{\Omega} V(t, \phi(t, s, \omega, \xi)-\phi(u, s, \omega, \eta)) \mathbb{P}(\mathrm{d} \omega)$. Thus, by the first part of (4.10) and Gronwall's inequality, we arrive at

$$
\begin{align*}
\mathbb{E}|\phi(t, s, \xi)-\phi(t, s, \eta)|^{p} & \leqslant \mathbb{E}[V(t, \phi(t, s, \xi)-\phi(t, s, \eta))] \\
& \leqslant \mathbb{E}[V(s, \xi-\eta)] \exp \left(\int_{s}^{t} \lambda(u) \mathrm{d} u\right) \tag{4.20}
\end{align*}
$$

Finally, given the bound (4.20), for $r<s<t$, we have

$$
\begin{aligned}
\mathbb{E}|\phi(t, r, \xi)-\phi(t, s, \xi)|^{p} & =\mathbb{E}|\phi(t, s, \phi(s, r, \xi))-\phi(t, s, \xi)|^{p} \\
& \leqslant \mathbb{E}[V(s, \phi(s, r, \xi)-\xi)] \exp \left(\int_{s}^{t} \lambda(u) \mathrm{d} u\right)
\end{aligned}
$$

and, utilising the above with Assumption 4.4(iii), yields

$$
\begin{equation*}
\limsup _{r<s,(t-s) \rightarrow \infty} \mathbb{E}|\phi(t, r, \xi)-\phi(t, s, \xi)|^{p}=0 \tag{4.21}
\end{equation*}
$$

Thus, for $\xi \in L^{p}\left(\Omega, \mathcal{F}_{-\infty}^{s}, \mathbb{P}\right), 1<p<\infty$, the above bound implies that the $L^{p}$ limit of the flow $\{\phi(t, s, \cdot, \xi): s \leqslant t\}$ exists as $s \rightarrow-\infty$. Note that this limit is independent of the initial condition $\xi$ by (4.12). We denote this limit by the random process $S: \mathbb{R} \times \Omega \rightarrow \mathbb{R}^{d}$, so that

$$
\mathbb{E}|S(t)-\phi(t, s, \xi)|^{p} \rightarrow 0 \quad \text { as } s \rightarrow-\infty,
$$

for $\xi \in L^{p}\left(\Omega, \mathcal{F}_{-\infty}^{s}, \mathbb{P}\right)$, where $S(t):=S(t, \cdot)$. Then, by Chebyshev's first inequality (aka Markov's inequality; e.g. [6]), for any $\varepsilon>0$, we have

$$
\begin{equation*}
\mathbb{P}(\{\omega \in \Omega:|S(t, \omega)-\phi(t, s, \omega, \xi)| \geqslant \varepsilon\}) \leqslant \varepsilon^{-p} \mathbb{E}|S(t)-\phi(t, s, \xi)|^{p} \tag{4.22}
\end{equation*}
$$

which implies that the convergence is also in probability. Thus, there exists a subsequence $\left(s_{k}\right)_{k \in \mathbb{N}_{1}}$ in $\mathbb{R}$ with $s_{k} \rightarrow-\infty$ as $k \rightarrow \infty$ such that

$$
S(t, \omega)=\lim _{k \rightarrow \infty} \phi\left(t, s_{k}, \omega, \xi\right), \quad \mathbb{P}-\text { a.s. }
$$

To simplify notation, we write

$$
\begin{equation*}
S(t, \omega)=\lim _{s \rightarrow-\infty} \phi(t, s, \omega, \xi), \quad \mathbb{P}-\text { a.s. } \tag{4.23}
\end{equation*}
$$

Note that for $\xi \in L^{p}\left(\Omega, \mathcal{F}_{-\infty}^{s}, \mathbb{P}\right)$ with the norm $\|\cdot\|_{p}:=\left(\mathbb{E}|\cdot|^{p}\right)^{1 / p}$ we have

$$
\begin{aligned}
\|\phi(t, s, \xi)\|_{p} & \leqslant\|\phi(t, s, \xi)-\xi\|_{p}+\|\xi\|_{p} \\
& \leqslant(\mathbb{E}[V(t, \phi(t, s, \xi)-\xi)])^{\frac{1}{p}}+\|\xi\|_{p} \\
& \leqslant\left(\sup _{s \leqslant t} \mathbb{E}[V(t, \phi(t, s, \xi)-\xi)]\right)^{\frac{1}{p}}+\|\xi\|_{p}<\infty,
\end{aligned}
$$

by condition (4.12) of Assumption 4.4. Consequently, for any $t \in \mathbb{R}$, we have

$$
\begin{equation*}
\|S(t)\|_{p} \leqslant \limsup _{s \rightarrow-\infty}\|\phi(t, s, \xi)\|_{p}<\infty, \tag{4.24}
\end{equation*}
$$

implying that $S(t, \omega)$ is bounded in $L^{p}\left(\Omega, \mathcal{F}_{-\infty}^{t}, \mathbb{P}\right)$.
Next, we show that $t \rightarrow S(t, \omega)$ is a random periodic path of period $0<\tau<\infty$ for the stochastic flow $\{\phi(t, s, \cdot, \cdot): s \leqslant t\}$ using its $\tau$-periodic property (see equation (3.13) with appropriately changed variables); namely

$$
S(t+\tau, \omega)=\lim _{s \rightarrow-\infty} \phi(t+\tau, s, \omega, \xi)
$$

$$
\begin{align*}
& =\lim _{s \rightarrow-\infty} \phi(t+\tau, s-\tau+\tau, \omega, \xi) \\
& =\lim _{s \rightarrow-\infty} \phi\left(t, s-\tau, \theta_{\tau} \omega, \xi\right) \\
& =S\left(t, \theta_{\tau} \omega\right) \quad \mathbb{P}-\text { a.s. } \tag{4.25}
\end{align*}
$$

Then, by the continuity of $(t, s, x) \mapsto \phi(t, s, \cdot x)$ and the flow property, we have

$$
\begin{align*}
\phi(t+s, s, \omega, S(s, \omega)) & =\lim _{r \rightarrow-\infty} \phi(t+s, s, \omega, \phi(s, r, \omega, \xi)) \\
& =\lim _{r \rightarrow-\infty} \phi(t+s, r, \omega, \xi) \\
& =S(t+s, \omega) \quad \forall t \in \mathbb{R}^{+}, s \in \mathbb{R} \quad \mathbb{P} \text { - a.s. } \tag{4.26}
\end{align*}
$$

The equalities (4.25) and (4.26) imply that $S(t, \omega)$ is a random periodic path (3.14) of period $0<\tau<\infty$ of the stochastic flow $\left\{\phi(t+s, s, \cdot, \cdot): s \in \mathbb{R}, t \in \mathbb{R}^{+}\right\}$on $\mathbb{R}^{d}$.
Uniqueness: Let $S_{1}(t, \omega)$ and $S_{2}(t, \omega)$ be two random periodic paths of the forward stochas$\left.\overline{\text { tic flow }\{\phi(t}+s, s, \cdot, \cdot): s \in \mathbb{R}, t \in \mathbb{R}^{+}\right\}$on $\mathbb{R}^{d}$. We know from (4.26) that for $s \leqslant t$,

$$
\begin{array}{ll}
S_{1}(t, \omega)=\phi\left(t, s, \omega, S_{1}(s, \omega)\right) & \mathbb{P}-\text { a.s. } \\
S_{2}(t, \omega)=\phi\left(t, s, \omega, S_{2}(s, \omega)\right) & \mathbb{P} \text { - a.s. }
\end{array}
$$

Then, for $1<p<\infty$, we have

$$
\begin{aligned}
\left\|S_{1}(t)-S_{2}(t)\right\|_{p}^{p} & =\left\|\phi\left(t, s, S_{1}(s)\right)-\phi\left(t, s, S_{2}(s)\right)\right\|_{p}^{p} \\
& \leqslant \exp (-\bar{\lambda}(t-s)) \mathbb{E}\left[V\left(s, S_{1}(s)-S_{2}(s)\right)\right] \underset{s \rightarrow-\infty}{\longrightarrow}
\end{aligned}
$$

Thus, $S_{1}(t, \omega)=S_{2}(t, \omega)$ for all $t \in \mathbb{R} \quad \mathbb{P}$ - a.s.
Construction of $\tau$-periodic measure for the $\operatorname{RDS} \tilde{\Phi}:$ Let $\tilde{S}: \mathbb{R} \times \Omega \rightarrow[0, \tau) \times \mathbb{R}^{d}$, $[0, \tau) \simeq \mathbb{R} \bmod \tau$, be defined by

$$
\tilde{S}(r, \omega)=(r \bmod \tau, \quad S(r, \omega)), \quad \forall r \in \mathbb{R}^{+}
$$

or, alternatively $\tilde{S}(r, \omega)=(|r| \bmod \tau, \quad S(r, \omega)), \forall r \in \mathbb{R}$. Then,

$$
\begin{equation*}
\tilde{S}(r+\tau, \omega)=(r+\tau \bmod \tau, S(r+\tau, \omega))=\left(r \bmod \tau, S\left(r, \theta_{\tau} \omega\right)\right) \tag{4.27}
\end{equation*}
$$

and from (4.5) and (4.26) we have

$$
\begin{align*}
\tilde{\Phi}\left(t, \theta_{r} \omega, \tilde{S}(r, \omega)\right) & =\tilde{\Phi}\left(t, \theta_{r} \omega,(r \bmod \tau, \quad S(r, \omega))\right) \\
& =(t+r \bmod \tau, \quad \phi(t+r, r, \omega, S(r, \omega))) \\
& =(t+r \bmod \tau, \quad S(t+r, \omega)) \\
& =\tilde{S}(t+r, \omega), \quad \forall t, r \in \mathbb{R}^{+} \quad \mathbb{P}-\text { a.s. } \tag{4.28}
\end{align*}
$$

The equalities (4.27)-(4.28) and the lifted version of (3.15) imply that $\tilde{S}(t, \omega)$ is a random periodic path of period $\tau$ of the RDS generated by $\tilde{\Phi}$ (4.5) in the skew-product representation on the flat cylinder $[0, \tau) \times \mathbb{R}^{d}$.

Finally, let $\left(\tilde{\mu}_{t}\right)_{t \in \mathbb{R}^{+}}, \tilde{\mu}_{t} \in \mathcal{P}\left([0, \tau) \times \mathbb{R}^{d}\right)$ be defined by

$$
\tilde{\mu}_{t}(\tilde{A})=\mathbb{P}(\{\omega: \tilde{S}(t, \omega) \in \tilde{A}\}), \quad \forall t \in \mathbb{R}^{+}, \tilde{A} \in \mathcal{B}([0, \tau)) \otimes \mathcal{B}\left(\mathbb{R}^{d}\right)
$$

It follows from (4.27)-(4.28) and Lemma 4.3 that the probability measure $\tilde{\mu}_{t}$ is $\tau$-periodic under the action of the transition evolution $\left(\tilde{\mathcal{P}}_{t}^{*}\right)_{t \in \mathbb{R}^{+}}$which is induced by the $\operatorname{RDS}\left\{\tilde{\Phi}(t, \cdot, \cdot): t \in \mathbb{R}^{+}\right\}$on $[0, \tau) \times \mathbb{R}^{d}$. The skew-product structure of these measures in $\mathcal{P}\left([0, \tau) \times \mathbb{R}^{d}\right)$ arises from Lemma 4.3, or directly from (4.6), so that for any $\mathcal{J} \in \mathcal{B}([0, \tau)), A \in \mathcal{B}\left(\mathbb{R}^{d}\right)$

$$
\tilde{\mu}_{t}(\mathcal{J} \times A)=\delta_{(t \bmod \tau)}(\mathcal{J}) \otimes \mathbb{P}(\omega: S(t, \omega) \in A)=\delta_{(t \bmod \tau)}(\mathcal{J}) \otimes \mu_{t}(A)
$$

### 4.3 Ergodicity of time-periodic measures

In this section, we turn to establishing ergodicity of the $\tau$-periodic measures $\left(\tilde{\mu}_{t}\right)_{t \in \mathbb{R}^{+}}$, $\tilde{\mu}_{t} \in \mathcal{P}\left([0, \tau) \times \mathbb{R}^{d}\right)$, generated by the Markovian ${ }^{13} \operatorname{RDS}\left\{\tilde{\Phi}(t, \cdot, \cdot): t \in \mathbb{R}^{+}\right\}$which was constructed in (4.5) in the skew-product representation on the flat cylinder $[0, \tau) \times \mathbb{R}^{d}$ from the lifted flow of solutions of the SDE (2.1) with time-periodic coefficients. The existence of $\tau$-periodic measures supported on stable random periodic paths was established in Theorem 4.7. The lack of stationarity and the unavoidable skew-product structure of the underlying dynamics pose additional challenges when dealing with ergodicity of $\tilde{\mathcal{P}}_{t}^{*}$-invariant measures, as outlined below. The main theorem of this section (Theorem 4.11) is preceded by some preparatory results and definitions.

Definition 4.8 (Ergodic periodic measure [33]) A family of $\tau$-periodic measures $\left(\tilde{\mu}_{t}\right)_{t \in \mathbb{R}^{+}}$ on the extended state space $\left([0, \tau) \times \mathbb{R}^{d}, \mathcal{B}([0, \tau)) \otimes \mathcal{B}\left(\mathbb{R}^{d}\right)\right)$ is said to be ergodic if

$$
\begin{equation*}
\overline{\tilde{\mu}}=\frac{1}{\tau} \int_{0}^{\tau} \tilde{\mu}_{t} \mathrm{~d} t \tag{4.29}
\end{equation*}
$$

is ergodic with respect to the transition semigroup $\left(\tilde{\mathcal{P}}_{t}^{*}\right)_{t \in \mathbb{R}^{+}}$in (4.8).
One can check by the linearity of $\tilde{\mu}_{0} \mapsto \tilde{\mathcal{P}}_{t}^{*} \tilde{\mu}_{0}$ and Fubini's theorem that $\overline{\tilde{\mu}}$ is an invariant measure for the transition semigroup $\left(\tilde{\mathcal{P}}_{t}^{*}\right)_{t \in \mathbb{R}^{+}}$defined in (4.8); i.e. $\tilde{\mathcal{P}}_{t}^{*}$ - invariance of $\overline{\tilde{\mu}}$ implies $\tilde{\mathcal{P}}_{t}^{*} \tilde{\tilde{\mu}}=\tilde{\tilde{\mu}}$, for all $t \in \mathbb{R}^{+}$. Moreover, from the definition of a $\tau$-periodic measure $\tilde{\mu}_{t}$ in (4.19), induced by the $\operatorname{RDS}\left\{\tilde{\Phi}(t, \cdot, \cdot): t \in \mathbb{R}^{+}\right\}$on $[0, \tau) \times \mathbb{R}^{d}$, we have

$$
\begin{aligned}
\overline{\tilde{\mu}}(\tilde{A}) & =\frac{1}{\tau} \int_{0}^{\tau} \tilde{\mu}_{t}(\tilde{A}) \mathrm{d} t=\frac{1}{\tau} \int_{0}^{\tau} \mathbb{P}(\{\omega: \tilde{S}(t, \omega) \in \tilde{A}\}) \mathrm{d} t=\frac{1}{\tau} \mathbb{E}\left[\int_{0}^{\tau} \mathbb{I}_{\tilde{A}}(\tilde{S}(t, \cdot)) \mathrm{d} t\right] \\
& =\mathbb{E}\left[\frac{1}{\tau} m_{1}(\{t \in[0, \tau): \tilde{S}(t, \cdot) \in \tilde{A}\})\right], \quad \tilde{A} \in \mathcal{B}([0, \tau)) \otimes \mathcal{B}\left(\mathbb{R}^{d}\right),
\end{aligned}
$$

where $t \rightarrow \tilde{S}(t, \omega)=(t \bmod \tau, S(t, \omega)), t \in \mathbb{R}^{+}$, is a random periodic path (3.15) of an RDS generated by the lifted dynamics of the $\operatorname{SDE}(2.1)$ via $\tilde{\Phi}$ in (4.5), and $m_{1}$ is the

[^8]Lebesgue measure on $\mathbb{R}$. Thus, given the invariance of $\overline{\tilde{\mu}}$ under the action of the transition semigroup $\left(\tilde{\mathcal{P}}_{t}^{*}\right)_{t \in \mathbb{R}^{+}}$in (4.8), and the $\tau$-periodicity of $\tilde{\mu}_{t}$ (see Definition 4.1), one has

$$
\begin{aligned}
\mathbb{E}\left[\frac{1}{\tau} m_{1}(\{t \in[0, \tau): \tilde{S}(t, \cdot) \in \tilde{A}\})\right] & =\left(\tilde{\mathcal{P}}_{u}^{*} \tilde{\tilde{\mu}}\right)(\tilde{A})=\frac{1}{\tau} \int_{0}^{\tau}\left(\tilde{\mathcal{P}}_{u}^{*} \tilde{\mu}_{t}\right)(\tilde{A}) \mathrm{d} t \\
& =\frac{1}{\tau} \int_{0}^{\tau} \tilde{\mu}_{t+u}(\tilde{A}) \mathrm{d} t=\frac{1}{\tau} \int_{u}^{u+\tau} \tilde{\mu}_{t}(\tilde{A}) \mathrm{d} t \\
& =\mathbb{E}\left[\frac{1}{\tau} m_{1}(\{t \in[u, u+\tau): \tilde{S}(t, \cdot) \in \tilde{A}\})\right]
\end{aligned}
$$

for any $\tilde{A} \in \mathcal{B}([0, \tau)) \otimes \mathcal{B}\left(\mathbb{R}^{d}\right)$ and any $u \in \mathbb{R}^{+}$. This implies that the expected time spent by the random periodic path $t \mapsto \tilde{S}(t, \omega)$ in any set $\tilde{A} \in \mathcal{B}([0, \tau)) \otimes \mathcal{B}\left(\mathbb{R}^{d}\right)$ over a time interval of exactly one period is independent of the starting point.

Verification of ergodicity (in the sense of Definition 4.8) of $\tau$-periodic measures $\left(\tilde{\mu}_{t}\right)_{t \in \mathbb{R}^{+}}$ supported on the random periodic paths of $\tilde{\Phi}$ requires one to assert that the time-averaged measure $\tilde{\tilde{\mu}}$ in the skew-product fibre bundle on $\mathcal{P}\left([0, \tau) \times \mathbb{R}^{d}\right)$ is $\tilde{\mathcal{P}}_{t}^{*}$ - ergodic. This set-up arises from the need to deal with the random periodic nature of the underlying dynamics, and it prevents a direct application of the classical tools for asserting ergodicity in the (asymptotically) stationary case. In particular, it is well known (e.g. [25, Theorem 3.2.4]) that the following are equivalent: ${ }^{14}$
(i) A probability measure $\overline{\tilde{\mu}}$ is weakly mixing.
(ii) There exists $\mathcal{I} \subset[0, \infty)$ of relative measure 1 such that $\lim _{t \rightarrow \infty, t \in \mathcal{I}} \tilde{P}(t, \tilde{x}, \cdot) \rightarrow \overline{\tilde{\mu}}$ weakly.

Thus, given the form of the transition kernel $\tilde{P}$ in (4.6) and the underlying skew-product structure, it is clear that one cannot establish the mixing property in the random periodic regime. ${ }^{15}$ Thus, this key condition in Doob's Theorem [27] does not hold in the random periodic regime which, alongside the lack of irreducibility of the transition kernel, renders the Hasminskii's Theorem [44] for asserting regularity of the transition kernel (needed in Doob's Theorem) inapplicable.
Instead, the $\tilde{\mathcal{P}}_{t}^{*}$-ergodicity of $\tilde{\tilde{\mu}}$ can be verified by means of a proposition which was proved in [33, Lemma 2.18]; we repeat its statement below with a concise proof to make this section self-contained. The main benefit of utilising the proposition below when dealing with $\overline{\tilde{\mu}}$ is that it essentially relies on ergodicity of $\tau$-periodic measures $\tilde{\mu}_{t}$ for any fixed $t \in[0, \tau)$ with respect to the discrete dynamics induced by $\left(\tilde{\mathcal{P}}_{n \tau}^{*}\right)_{n \in \mathbb{N}_{0}}$; the subsequent use of the semigroup property of $\left(\tilde{\mathcal{P}}_{t}^{*}\right)_{t \in \mathbb{R}^{+}}$allows one to show the ergodicity of $\tilde{\tilde{\mu}}$. Importantly, the $\tilde{\mathcal{P}}_{n \tau}^{*}$-ergodicity of $\tilde{\mu}_{t}$ on the respective Poincaré sections with a fixed $t \in[0, \tau)$ turns the problem into a stationary one which can be dealt with using the standard methods. The result below provides an extension of the classical Krylov-Bogolyubov procedure (see, for example, [8, \$1.5]).

Proposition 4.9 Consider a family of $\tau$-periodic measures $\left(\tilde{\mu}_{t}\right)_{t \in \mathbb{R}^{+}}$on the extended state space $\left([0, \tau) \times \mathbb{R}^{d}, \mathcal{B}([0, \tau)) \otimes \mathcal{B}\left(\mathbb{R}^{d}\right)\right)$. The $\tilde{\mathcal{P}}_{t}^{*}$-invariant measure $\overline{\tilde{\mu}}$ in (4.29) is ergodic if and only if the following holds for any $\tilde{A} \in \mathcal{B}([0, \tau)) \otimes \mathcal{B}\left(\mathbb{R}^{d}\right)$

[^9]\[

$$
\begin{equation*}
\lim _{N \rightarrow \infty} \int_{[0, \tau) \times \mathbb{R}^{d}}\left|\int_{0}^{\tau}\left\{\frac{1}{N} \sum_{n=0}^{N-1} \tilde{P}(\tilde{x} ; t+n \tau, \tilde{A})-\tilde{\mu}_{t}(\tilde{A})\right\} \mathrm{d} t\right| \tilde{\tilde{\mu}}(\mathrm{d} \tilde{x})=0 . \tag{4.30}
\end{equation*}
$$

\]

Proof Recall from (e.g. [8]) that $\overline{\tilde{\mu}}$ is ergodic if $\tilde{\mathcal{P}}_{t} \mathbb{I}_{\tilde{A}}=\mathbb{I}_{\tilde{A}}$, $\overline{\tilde{\mu}}$ - a.e. $\tilde{A} \in \mathcal{B}([0, \tau)) \otimes \mathcal{B}\left(\mathbb{R}^{d}\right)$ implies that either $\overline{\tilde{\mu}}(\tilde{A})=0$ or $\tilde{\tilde{\mu}}(\tilde{A})=1$. First, we assume that (4.30) holds for any $\tilde{A} \in \mathcal{B}([0, \tau)) \otimes \mathcal{B}\left(\mathbb{R}^{d}\right)$ with $\tilde{P}(\tilde{x} ; t, \tilde{A})=\tilde{\mathcal{P}}_{t} \mathbb{I}_{\tilde{A}}(\tilde{x})=\mathbb{I}_{\tilde{A}}(\tilde{x})$. Then, it follows from (4.30) that

$$
\int_{[0, \tau) \times \mathbb{R}^{d}}\left|\mathbb{I}_{\tilde{A}}(\tilde{x})-\tilde{\tilde{\mu}}(\tilde{A})\right| \overline{\tilde{\mu}}(\mathrm{d} \tilde{x})=\int_{[0, \tau) \times \mathbb{R}^{d}}\left|\frac{1}{\tau} \int_{n \tau}^{(n+1) \tau} \tilde{P}(\tilde{x} ; t, \tilde{A}) \mathrm{d} t-\tilde{\tilde{\mu}}(\tilde{A})\right| \tilde{\tilde{\mu}}(\mathrm{d} \tilde{x})=0
$$

This implies that $\mathbb{I}_{\tilde{A}}(\tilde{x})$ is a constant for $\overline{\tilde{\mu}}$ - a.e. $\tilde{x} \in[0, \tau) \times \mathbb{R}^{d}$. Thus, either $\tilde{\tilde{\mu}}(\tilde{A})=0$ or $\overline{\tilde{\mu}}(\tilde{A})=1$. Conversely, assume that $\overline{\tilde{\mu}}$ is ergodic, then for any $\tilde{A} \in \mathcal{B}([0, \tau)) \otimes \mathcal{B}\left(\mathbb{R}^{d}\right)$

$$
\lim _{T \rightarrow \infty} \frac{1}{T} \int_{0}^{T} \tilde{P}(\tilde{x} ; t, \tilde{A}) \mathrm{d} t=\overline{\tilde{\mu}}(\tilde{A}) \quad \text { in } L^{2}(\overline{\tilde{\mu}})
$$

Therefore,

$$
\begin{equation*}
\lim _{N \rightarrow \infty} \frac{1}{N \tau} \sum_{n=0}^{N-1} \int_{0}^{\tau} \tilde{P}(\tilde{x} ; t+n \tau, \tilde{A}) \mathrm{d} t=\tilde{\tilde{\mu}}(\tilde{A}) \quad \text { in } L^{2}(\tilde{\tilde{\mu}}), \tag{4.31}
\end{equation*}
$$

and (4.30) follows from (4.31) and from the Cauchy-Schwartz inequality.
Consequently, the subsequent verification of the ergodicity of the $\tilde{\mathcal{P}}_{t}^{*}$ - invariant measure $\overline{\tilde{\mu}}$ on $[0, \tau) \times \mathbb{R}^{d}$ relies (explicitly or otherwise) on the semigroup property and periodicity of the transition semigroup $\left(\tilde{\mathcal{P}}_{t}^{*}\right)_{t \in \mathbb{R}^{+}}$, and on proving the strong Feller property of the transition evolution $\left(\mathcal{P}_{s, t}\right)_{t \geqslant s}$ in (3.3). Recall that the transition evolution $\left(\mathcal{P}_{s, t}\right)_{t \geqslant s}$ has the strong Feller property (i.e. $\mathcal{P}_{s, t} \varphi \in \mathcal{C}_{\infty}\left(\mathbb{R}^{d}\right)$ for any $\varphi \in \mathbb{M}_{\infty}\left(\mathbb{R}^{d}\right)$ ) if and only if
(i) $\left(\mathcal{P}_{s, t}\right)_{t \geqslant s}$ is Feller, i.e. $\mathcal{P}_{s, t}: \mathcal{C}_{\infty}\left(\mathbb{R}^{d}\right) \rightarrow \mathcal{C}_{\infty}\left(\mathbb{R}^{d}\right)$, and
(ii) For any $\varphi \in \mathcal{C}_{\infty}\left(\mathbb{R}^{d}\right)$ the family $\left(\mathcal{P}_{s, t} \varphi\right)_{t \geqslant s}$ is equicontinuous.

The first condition follows from the existence of the stochastic flow (see, for example, [44,54]); thus, we only derive the second item in Proposition 4.10 below.

Proposition 4.10 Suppose that Assumption 4.6 holds. Then, for any $t \in[s, s+T]$, there exist $0<C_{T}<\infty$ such that, for any $x, y \in \mathbb{R}^{d}$ and $\varphi \in \mathcal{C}_{\infty}\left(\mathbb{R}^{d}\right)$, we have

$$
\left|\mathcal{P}_{s, t} \varphi(x)-\mathcal{P}_{s, t} \varphi(y)\right| \leqslant C_{T}\|\varphi\|_{\infty}|x-y| .
$$

Proof The proof consists of a tedious but relatively straightforward extension of results which are well known in the autonomous case; for detailed derivations, involving some Malliavin calculus estimates; see Theorem B. 10 in "Appendix 5.3.2".

Given the above setting, we have the following main result of this section:
Theorem 4.11 Suppose that Proposition 4.10 and Assumption 4.4 hold. Then, the family of $\tau$-periodic measures $\left(\tilde{\mu}_{t}\right)_{t \in \mathbb{R}^{+}}, \tilde{\mu}_{t} \in \mathcal{P}\left([0, \tau) \times \mathbb{R}^{d}\right)$, in (4.19) is ergodic in the sense of Definition 4.8.

Remark 4.12 The requirement in the above theorem that Assumption 4.4 holds is inherited from the conditions required in Theorem 4.7 for the existence of stable random periodic paths on which the $\tau$-periodic measures $\left(\tilde{\mu}_{t}\right)_{t \in \mathbb{R}^{+}}$are supported; hence, the only additional condition in Theorem 4.11 is introduced by imposing Assumption 4.6 which is required in Proposition 4.10 to assert the strong Feller property of $\left(\mathcal{P}_{s, t}\right)_{t \geqslant s}$. If one dropped Assumption 4.4, the existence of $\tau$-periodic measures would have to be assumed a priori alongside the ergodicity of $\tilde{\mu}_{t}$ for all fixed $t \in[0, \tau)$ w.r.t. the discrete transition evolution $\left(\tilde{\mathcal{P}}_{n \tau}\right)_{n \in \mathbb{N}_{0}}$, as done in [33]. In the present case, the properties of the $\tau$-periodic measures derived explicitly in the previous section allow us to dispense with such assumptions.

Proof of Theorem 4.11. The proof is relatively long, and we divide it into four steps. Throughout, we skip the dependence on $\omega \in \Omega$ in all quantities involving expectations. Step I: First, we show that for a random periodic path $S: \mathbb{R} \times \Omega \rightarrow \mathbb{R}^{d}$ of the stochastic $\overline{\text { flow } \phi}$ on $\mathbb{R}^{d}$, and $\eta \in L^{p}\left(\Omega, \mathcal{F}_{-\infty}^{s}, \mathbb{P}\right), 1<p<\infty$, there exists $0<\tilde{\mathfrak{C}}<\infty$ such that

$$
\begin{equation*}
\|\phi(s+n \tau, s, \eta)-S(s+n \tau)\|_{p} \leqslant \tilde{\mathfrak{C}} \exp \left(\frac{1}{p} \int_{s}^{s+n \tau} \lambda(u) \mathrm{d} u\right), \quad n \in \mathbb{N}_{0} \tag{4.32}
\end{equation*}
$$

To see this, note that from the definition of the random periodic path of a stochastic flow (3.14) we have $S(s+n \tau, \omega)=\phi(s+n \tau, s, \omega, S(s, \omega)) \mathbb{P}$-a.s., so that

$$
\begin{align*}
\|\phi(s+n \tau, s, \eta)-S(s+n \tau)\|_{p} & =\|\phi(s+n \tau, s, \eta)-\phi(s+n \tau, s, S(s))\|_{p} \\
& \leqslant(\mathbb{E}[V(s, \eta-S(s))])^{\frac{1}{p}} \exp \left(\frac{1}{p} \int_{s}^{s+n \tau} \lambda(u) \mathrm{d} u\right) \\
& =\tilde{\mathbb{C}} \exp \left(\frac{1}{p} \int_{s}^{s+n \tau} \lambda(u) \mathrm{d} u\right), \quad n \in \mathbb{N}_{0}, \tag{4.33}
\end{align*}
$$

by Assumption 4.4(i) and the fact that $S(s) \in L^{p}\left(\Omega, \mathcal{F}_{-\infty}^{s}, \mathbb{P}\right), 1<p<\infty$, which was shown in the proof of Theorem 4.7.
Step II: We show that for $1<p<\infty$, there exists $0<\mathfrak{C}_{\tau}<\infty$, such that for $n \in \mathbb{N}_{0}$

$$
\begin{equation*}
\left|\mathcal{P}_{s, s+n \tau} \varphi(x)-\int_{\mathbb{R}^{d}} \varphi(y) \mu_{s}(\mathrm{~d} y)\right| \leqslant \mathfrak{C}_{\tau}\|\varphi\|_{\infty} \exp \left(\frac{1}{p} \int_{s}^{s+n \tau} \lambda(u) \mathrm{d} u\right), \quad \varphi \in \mathcal{C}_{\infty}\left(\mathbb{R}^{d}\right) \tag{4.34}
\end{equation*}
$$

where $\mu_{s}(A)=\mathbb{P}(\{\omega: S(s, \omega) \in A\}), A \in \mathcal{B}\left(\mathbb{R}^{d}\right)$.
To see this, we note that from the definition of the periodic measure $\mu_{s}$, we have that

$$
\int_{\mathbb{R}^{d}} \mathcal{P}_{s, s+n \tau} \varphi(y) \mu_{s}(\mathrm{~d} y)=\int_{\mathbb{R}^{d}} \varphi(y) \mu_{s}(\mathrm{~d} y), \quad \varphi \in \mathcal{C}_{\infty}\left(\mathbb{R}^{d}\right) ;
$$

i.e. $\mu_{s}$ is invariant under the action of the dual of the discrete transition evolution $\left(\mathcal{P}_{s, s+n \tau}^{*}\right)_{n \in \mathbb{N}_{0}}$. Thus, for $\psi \in \operatorname{Lip}_{\infty}\left(\mathbb{R}^{d}\right)$, we have for $1<p<\infty$,

$$
\begin{aligned}
\left|\mathcal{P}_{s, s+n \tau} \psi(x)-\int_{\mathbb{R}^{d}} \psi(y) \mu_{s}(\mathrm{~d} y)\right| & =\left|\int_{\mathbb{R}^{d}}\left(\mathcal{P}_{s, s+n \tau} \psi(x)-\mathcal{P}_{s, s+n \tau} \psi(y)\right) \mu_{s}(\mathrm{~d} y)\right| \\
& \leqslant\|\psi\|_{\mathrm{BL}} \int_{\mathbb{R}^{d}} \mathbb{E}|\phi(s+n \tau, s, x)-\phi(s+n \tau, s, y)| \mu_{s}(\mathrm{~d} y) \\
& =\|\psi\|_{\mathrm{BL}} \mathbb{E}|\phi(s+n \tau, s, x)-\phi(s+n \tau, s, S(s))|
\end{aligned}
$$

$$
\begin{align*}
& \leqslant\|\psi\|_{\mathrm{BL}}\left(\mathbb{E}|\phi(s+n \tau, s, x)-S(s+n \tau)|^{p}\right)^{\frac{1}{p}} \\
& \leqslant \tilde{\mathfrak{C}}\|\psi\|_{\mathrm{BL}} \exp \left(\frac{1}{p} \int_{s}^{s+n \tau} \lambda(u) \mathrm{d} u\right) \tag{4.35}
\end{align*}
$$

where we applied Hölder's inequality and estimate (4.32) in the last two lines, respectively.
Now, let $\varphi \in \mathcal{C}_{\infty}\left(\mathbb{R}^{d}\right)$ be given. Setting $\psi=\mathcal{P}_{s+n \tau, s+\tau+n \tau} \varphi=\mathcal{P}_{s, s+\tau} \varphi$ in (4.35), which holds due to (3.13), and using the invariance of $\mu_{s}$ under the transition evolution $\left(\mathcal{P}_{s, s+n \tau}^{*}\right)_{n \in \mathbb{N}_{0}}$, we obtain by Proposition 4.10 that

$$
\begin{align*}
\left|\mathcal{P}_{s, s+\tau+n \tau} \varphi(x)-\int_{\mathbb{R}^{d}} \mathcal{P}_{s, s+\tau} \varphi(y) \mu_{s}(\mathrm{~d} y)\right| & =\left|\int_{\mathbb{R}^{d}}\left(\mathcal{P}_{s, s+\tau+n \tau} \varphi(x)-\mathcal{P}_{s, s+\tau+n \tau} \varphi(y)\right) \mu_{s}(\mathrm{~d} y)\right| \\
& \leqslant \tilde{\mathfrak{C}}\left\|\mathcal{P}_{s+n \tau, s+\tau+n \tau} \varphi\right\|_{\mathrm{BL}} \exp \left(\frac{1}{p} \int_{s}^{s+n \tau} \lambda(u) \mathrm{d} u\right) \\
& =\tilde{\mathfrak{C}}\left\|\mathcal{P}_{s, s+\tau} \varphi\right\|_{\mathrm{BL}} \exp \left(\frac{1}{p} \int_{s}^{s+n \tau} \lambda(u) \mathrm{d} u\right) \\
& \leqslant \mathfrak{C}_{\tau}\|\varphi\|_{\infty} \exp \left(\frac{1}{p} \int_{s}^{s+n \tau} \lambda(u) \mathrm{d} u\right) \tag{4.36}
\end{align*}
$$

where $\mathfrak{C}_{\tau}=C_{\tau} \tilde{\mathfrak{C}}$, and $C_{\tau}$ is a constant appearing in Proposition 4.10.
Step III: Let $A \subset \mathbb{R}^{d}$ be a closed set, take $\varphi=\mathbb{I}_{A}$, and consider the sequence $\left(\varphi_{m}\right)_{m \in \mathbb{N}_{1}}$ of functions defined by

$$
\varphi_{m}(x)= \begin{cases}1, & \text { if } x \in A \\ 1-2^{m} \mathrm{~d}(x, A), & \text { if } \mathrm{d}(x, A) \leqslant 2^{-m} \\ 0, & \text { if } \mathrm{d}(x, A) \geqslant 2^{-m}\end{cases}
$$

where $\mathrm{d}(x, A)=\inf \{|x-y|: y \in A\}, x \in \mathbb{R}^{d}$. Then,

$$
\varphi_{m}(x) \rightarrow \varphi(x), \quad \text { as } m \rightarrow \infty, \text { for all } x \in \mathbb{R}^{d}
$$

Next, for $s \in[0, \tau)$, we have

$$
\mathcal{P}_{s, s+n \tau} \varphi_{m}(x) \rightarrow \mathcal{P}_{s, s+n \tau} \varphi(x)=\mathcal{P}_{s, s+n \tau} \mathbb{I}_{A}(x)
$$

which implies that $P(s, \cdot ; s+n \tau, A)=\mathcal{P}_{s, s+n \tau} \mathbb{I}_{A} \in \mathcal{C}_{\infty}\left(\mathbb{R}^{d}\right)$ and, since $\mu_{s}$ is invariant under $\left(\mathcal{P}_{s, s+n \tau}^{*}\right)_{n \in \mathbb{N}_{0}}$, (4.36) leads to

$$
\begin{equation*}
\left|P(s, x ; s+n \tau, A)-\mu_{s}(A)\right| \leqslant \mathfrak{C}_{\tau} \exp \left(\frac{1}{p} \int_{s}^{s+n \tau} \lambda(u) \mathrm{d} u\right) \tag{4.37}
\end{equation*}
$$

By the covering lemma (e.g. [6]), the inequality (4.37) holds for any $A \in \mathcal{B}\left(\mathbb{R}^{d}\right)$, and thus, for $\mathcal{J} \subseteq[0, \tau)$, we have

$$
\begin{aligned}
\int_{\mathcal{J}}\left|P(s, x ; s+n \tau, A)-\mu_{s}(A)\right| \mathrm{d} s & \leqslant \int_{0}^{\tau}\left|P(s, x ; s+n \tau, A)-\mu_{s}(A)\right| \mathrm{d} s \\
& \leqslant \mathfrak{C}_{\tau} \int_{0}^{\tau} \exp \left(\frac{1}{p} \int_{s}^{s+n \tau} \lambda(u) \mathrm{d} u\right) \mathrm{d} s
\end{aligned}
$$

$$
=\mathfrak{C}_{\tau} \int_{0}^{\tau} \exp \left(\frac{1}{p n \tau} \int_{s}^{s+n \tau} \lambda(u) \mathrm{d} u\right)^{n \tau} \mathrm{~d} s
$$

Now, we use the Chapman-Kolmogorov equation (3.2) for the transition probability to obtain

$$
\begin{aligned}
\left|\int_{\mathcal{J}}\left[P(s, x ; t+n \tau, A)-\mu_{t}(A)\right] \mathrm{d} t\right| & =\left|\left[\int_{\mathcal{J}} \int_{\mathbb{R}^{d}} P(t, y ; t+n \tau, A)-\mu_{t}(A)\right] P(s, x ; t, \mathrm{~d} y) \mathrm{d} t\right| \\
& \leqslant \int_{0}^{\tau} \int_{\mathbb{R}^{d}} \mathfrak{C}_{\tau} \exp \left(\frac{1}{p n \tau} \int_{t}^{t+n \tau} \lambda(u) \mathrm{d} u\right)^{n \tau} P(s, x ; t, \mathrm{~d} y) \mathrm{d} t \\
& =\mathfrak{C}_{\tau} \int_{0}^{\tau} \exp \left(\frac{1}{p n \tau} \int_{t}^{t+n \tau} \lambda(u) \mathrm{d} u\right)^{n \tau} \mathrm{~d} t .
\end{aligned}
$$

By condition (4.11) of Assumption 4.4, there exists $0<\beta<1,0<K<\infty$, such that

$$
\left|\int_{\mathcal{J}}\left(P(s, x ; t+n \tau, A)-\mu_{t}(A)\right) \mathrm{d} t\right| \leqslant \int_{\mathcal{J}}\left|P(s, x ; t+n \tau, A)-\mu_{t}(A)\right| \mathrm{d} t \leqslant K \beta^{n \tau}
$$

It then follows that

$$
\begin{equation*}
\frac{1}{\tau} \int_{0}^{\tau} \int_{\mathbb{R}^{d}}\left|\int_{\mathcal{J}}\left\{\frac{1}{N} \sum_{n=0}^{N-1} P(s, x ; t+n \tau, A)-\mu_{t}(A)\right\} \mathrm{d} t\right| \mu_{s}(\mathrm{~d} x) \mathrm{d} s \leqslant \frac{K}{N} \sum_{n=0}^{N-1} \beta^{n \tau} \underset{N \rightarrow \infty}{\longrightarrow} 0 \tag{4.38}
\end{equation*}
$$

Step IV: In this final step, with the help of Step III, we show the convergence of KrylovBogolyubov scheme for the $\tau$-periodic measures $\left(\tilde{\mu}_{t}\right)_{t \in \mathbb{R}^{+}}$on the cylinder $[0, \tau) \times \mathbb{R}^{d}$. For any $\mathcal{J} \times A \in \mathcal{B}([0, \tau)) \otimes \mathcal{B}\left(\mathbb{R}^{d}\right)$ we have

$$
\begin{aligned}
& \int_{[0, \tau) \times \mathbb{R}^{d}}\left|\int_{0}^{\tau}\left(\frac{1}{N} \sum_{n=0}^{N-1} \tilde{P}(\tilde{x} ; t+n \tau, \mathcal{J} \times A)-\tilde{\mu}_{t}(\mathcal{J} \times A)\right) \mathrm{d} t\right| \tilde{\tilde{\mu}}(\mathrm{d} \tilde{x}) \\
&= \frac{1}{\tau} \int_{0}^{\tau} \int_{\mathbb{R}^{d}}\left|\int_{0}^{\tau}\left(\frac{1}{N} \sum_{n=0}^{N-1} \tilde{P}((s, x) ; t+n \tau, \mathcal{J} \times A)-\tilde{\mu}_{t}(\mathcal{J} \times A)\right) \mathrm{d} t\right| \mu_{s}(\mathrm{~d} x) \mathrm{d} s \\
&= \frac{1}{\tau} \int_{0}^{\tau} \int_{\mathbb{R}^{d}}\left|\int_{0}^{\tau}\left(\frac{1}{N} \sum_{n=0}^{N-1} P(s, x ; t+s+n \tau, A)-\mu_{t}(A)\right) \delta_{(t+s \bmod \tau)}(\mathcal{J}) \mathrm{d} t\right| \mu_{s}(\mathrm{~d} x) \mathrm{d} s \\
&= \frac{1}{\tau} \int_{0}^{\tau} \int_{\mathbb{R}^{d}} \left\lvert\, \int_{0}^{\tau-s}\left(\frac{1}{N} \sum_{n=0}^{N-1} P(s, x ; t+s+n \tau, A)-\mu_{t}(A)\right) \delta_{(t+s)}(\mathcal{J}) \mathrm{d} t\right. \\
& \left.\quad+\int_{\tau-s}^{\tau}\left(\frac{1}{N} \sum_{n=0}^{N-1} P(s, x ; t+s+n \tau, A)-\mu_{t}(A)\right) \delta_{(t+s-\tau)}(\mathcal{J}) \mathrm{d} t \right\rvert\, \mu_{s}(\mathrm{~d} x) \mathrm{d} s \\
&= \frac{1}{\tau} \int_{0}^{\tau} \int_{\mathbb{R}^{d}} \left\lvert\, \int_{0}^{\tau-s}\left(\frac{1}{N} \sum_{n=0}^{N-1} P(s, x ; t+s+n \tau, A)-\mu_{t}(A)\right) \delta_{(t+s)}(\mathcal{J}) \mathrm{d} t\right. \\
& \left.\quad+\int_{-s}^{0}\left(\frac{1}{N} \sum_{n=1}^{N} P(s, x ; t+s+n \tau, A)-\mu_{t}(A)\right) \delta_{(t+s)}(\mathcal{J}) \mathrm{d} t \right\rvert\, \mu_{s}(\mathrm{~d} x) \mathrm{d} s \\
&= \frac{1}{\tau} \int_{0}^{\tau} \int_{\mathbb{R}^{d}} \left\lvert\, \int_{\mathcal{J}}\left(\frac{1}{N} \sum_{n=0}^{N-1} P(s, x ; t+n \tau, A)-\mu_{t}(A)\right) \mathrm{d} t\right. \\
& \left.-\frac{1}{N} \int_{-s}^{0}(P(s, x ; t+s, A)-P(s, x ; t+s+N \tau, A)) \delta_{(t+s)}(\mathcal{J}) \mathrm{d} t \right\rvert\, \mu_{s}(\mathrm{~d} x) \mathrm{d} s
\end{aligned}
$$

$$
\begin{aligned}
\leqslant & \frac{1}{\tau} \int_{0}^{\tau} \int_{\mathbb{R}^{d}}\left|\int_{\mathcal{J}}\left(\frac{1}{N} \sum_{n=0}^{N-1} P(s, x ; t+n \tau, A)-\mu_{t}(A)\right) \mathrm{d} t\right| \mu_{s}(\mathrm{~d} x) \mathrm{d} s \\
& \left.+\frac{1}{N \tau} \int_{0}^{\tau} \int_{\mathbb{R}^{d}} \right\rvert\, \int_{-s}^{0}(P(s, x ; t+s, A) \\
& -P(s, x ; t+s+N \tau, A)) \delta_{(t+s)}(\mathcal{J}) \mathrm{d} t \mid \mu_{s}(\mathrm{~d} x) \mathrm{d} s \underset{N \rightarrow \infty}{\longrightarrow} 0
\end{aligned}
$$

Remark 4.13 The invariance of the $\tau$-periodic probability measures under the discrete evolution $\left(\tilde{\mathcal{P}}_{n \tau}^{*}\right)_{n \in \mathbb{N}_{0}}$ on their respective Poincaré sections was pointed out in Lemma 4.3. It can be shown, as a consequence of [33, Theorem 4.11], that such $\tau$-periodic probability measures are ergodic w.r.t. the discrete evolution $\left(\tilde{\mathcal{P}}_{n \tau}^{*}\right)_{n \in \mathbb{N}_{0}}$ on their respective Poincaré sections; given that we require Assumption 4.6 to be satisfied, these measures are supported on all of $\mathbb{R}^{d}$. This fact will be useful in Sect. 5 concerned with ergodic averages in the context of the linear response.

Example 4.14 (Stochastic Lorenz model with periodicforcing) Consider a modified Lorenz system (e.g. [48]) given by

$$
\left\{\begin{array}{l}
\dot{x}=-\bar{\alpha} x+\bar{\alpha} y  \tag{4.39}\\
\dot{y}=-\bar{\alpha} x-\bar{\beta} y-x z \\
\dot{z}=-\bar{\gamma} z+x y-\bar{\gamma} \bar{\beta}^{-2}(\bar{\varrho}+\bar{\alpha})
\end{array}\right.
$$

with parameters $\bar{\alpha}, \bar{\beta}, \bar{\gamma}, \bar{\varrho}>0$. We set $\left(v_{1}, v_{2}, v_{3}\right):=(x, y, z) \in \mathbb{R}^{3}$ and consider the periodically-in-time and stochastically perturbed version of (4.39) for $t \in \mathbb{R}^{+}$in the form

$$
\begin{equation*}
\mathrm{d} v_{t}=b\left(t, v_{t}\right) \mathrm{d} t+\sigma\left(v_{t}\right) \mathrm{d} W_{t}=\left[-A v_{t}-G\left(v_{t}\right)+F(t)\right] \mathrm{d} t+\sigma\left(v_{t}\right) \mathrm{d} W_{t}, \quad v_{0} \in \mathbb{R}^{3} \tag{4.40}
\end{equation*}
$$

where

$$
\begin{aligned}
A & =\left[\begin{array}{ccc}
\bar{\alpha} & -\bar{\alpha} & 0 \\
\bar{\alpha} & \bar{\beta} & 0 \\
0 & 0 & \bar{\gamma}
\end{array}\right], G(v)=\left[\begin{array}{c}
0 \\
v_{1} v_{3} \\
-v_{1} v_{2}
\end{array}\right], \quad F(t)=\left[\begin{array}{c}
\bar{f}\left(1+\bar{\delta} \sin \left(\frac{2 \pi}{\tau} t\right)\right) \\
0 \\
-\bar{\gamma} \bar{\beta}^{-2}(\bar{\varrho}+\bar{\alpha})
\end{array}\right], \\
\sigma(v) & =\bar{\sigma}\left[\begin{array}{ccc}
v_{1} & 0 & 0 \\
0 & v_{2} & 0 \\
0 & 0 & v_{3}
\end{array}\right],
\end{aligned}
$$

with $|\bar{\delta}| \leqslant|\bar{f}|<\infty$ and $\bar{\sigma} \in \mathbb{R} \backslash\{0\}$ finite, $0<\tau<\infty$, and $W_{t}=\left(W_{t}^{1}, W_{t}^{2}, W_{t}^{3}\right)$ an independent Wiener process in $\mathbb{R}^{3}$. Although the above system is in the 'toy category', considering the effects of time-periodic forcing and stochastic perturbations is relevant in many atmosphere-ocean applications to model, for example, seasonal and diurnal cycles in climate models (e.g. [59,60,64,77-79]). It is well known that for $\bar{\sigma}=0$ the system (4.40) has an absorbing ball for all values of the parameters, since for $V(t, v)=|v|^{2}$ we have

$$
\begin{aligned}
\frac{1}{2} \frac{\mathrm{~d} V}{\mathrm{~d} t}=\langle b(v), v\rangle= & -\bar{\alpha}\left(v_{1}-\frac{F_{1}}{2 \bar{\alpha}}\right)^{2}-\bar{\beta} v_{2}^{2}-\bar{\gamma}\left(v_{3}+\frac{\bar{\varrho}+\bar{\alpha}}{2 \bar{\beta}^{2}}\right)^{2} \\
& +\frac{\bar{\alpha} \bar{\gamma} \bar{\beta}^{-2}(\bar{\varrho}+\bar{\alpha})^{2}+F_{1}^{2}}{4 \bar{\alpha}}
\end{aligned}
$$

where we skip the explicit time dependence and $F_{1}(t)=\bar{f}\left(1+\bar{\delta} \sin \left(\frac{2 \pi}{\tau} t\right)\right)$. Note that the drift and diffusion coefficients, $b, \sigma$, in (4.40) are smooth and satisfy the growth conditions (4.13) outlined in Remark 4.16(a); since for $0<\varkappa_{1}, \varkappa_{3}<\infty$, and $\bar{F}_{1}=\sup _{[0, \tau]}\left|f_{1}(t)\right|$ we have

$$
\langle b(v), v\rangle \leqslant-\bar{\alpha}\left(1-\frac{\bar{F}_{1}}{4 \bar{\alpha} \varkappa_{1}}\right) v_{1}^{2}-\bar{\beta} v_{2}^{2}-\bar{\gamma}\left(1-\frac{\bar{\varrho}+\bar{\alpha}}{4 \bar{\beta}^{2} \varkappa_{3}}\right) v_{3}^{2}+\bar{F}_{1} \varkappa_{1}+\bar{\gamma} \bar{\beta}^{-2}(\bar{\varrho}+\bar{\alpha}) \varkappa_{3}
$$

where we used the fact that $|x| \leqslant \varkappa+\frac{1}{4 \varkappa}|x|^{2}$ for $\varkappa>0$. Thus, we have

$$
\begin{equation*}
\langle b(v), v\rangle \leqslant L_{b_{1}}-L_{b_{2}}|v|^{2}, \quad\|\sigma(v)\|_{\mathrm{HS}}^{2} \leqslant L_{\sigma}\left(1+|v|^{2}\right) \tag{4.41}
\end{equation*}
$$

where $L_{\sigma}=\bar{\sigma}$ and

$$
\begin{equation*}
L_{b_{1}}=\varkappa_{1} \bar{F}_{1}+\varkappa_{3} \bar{\gamma} \bar{\beta}^{-2}(\bar{\varrho}+\bar{\alpha}), \quad L_{b_{2}}=\min \left(\bar{\beta}, \bar{\alpha}\left(1-\frac{\bar{F}_{1}}{4 \bar{\alpha} \varkappa_{1}}\right), \bar{\gamma}\left(1-\frac{\bar{\varrho}+\bar{\alpha}}{4 \bar{\beta}^{2} \varkappa_{3}}\right)\right) \tag{4.42}
\end{equation*}
$$

Thus, (4.40) has global solutions and it generates a stochastic flow of diffeomorphisms on $\mathbb{R}^{3}$.
Next, note that the linear part in (4.40) satisfies

$$
\langle A v, v\rangle_{\mathbb{R}^{3}} \geqslant \mathfrak{C}_{A}|v|^{2}, \quad \mathfrak{C}_{A}=\min \{\bar{\alpha}, \bar{\beta}, \bar{\gamma}\}
$$

and the nonlinear term $G(v)=B(v, v)$ is given by a bilinear map $B(v, w)=\left(0, v_{1} w_{3},-v_{1} w_{2}\right)$, $v, w \in \mathbb{R}^{3}$, which satisfies (see also [48])

$$
\left\{\begin{array}{l}
\langle B(v, w), w\rangle_{\mathbb{R}^{3}}=\left\langle\left(0, v_{1} w_{3},-v_{1} w_{2}\right),\left(w_{1}, w_{2}, w_{3}\right)\right\rangle_{\mathbb{R}^{3}}=0  \tag{4.43}\\
\langle B(v, w), u\rangle_{\mathbb{R}^{3}}=\left\langle\left(0, v_{1} w_{3},-v_{1} w_{2}\right),\left(u_{1}, u_{2}, u_{3}\right)\right\rangle_{\mathbb{R}^{3}}=-\langle B(v, u), w\rangle_{\mathbb{R}^{3}} \\
|B(v, w)| \leqslant|v||w|
\end{array}\right.
$$

Consider $V(t, v)=|v|^{p}$ for some $1<p<\infty$, so that

$$
\partial_{v_{i}} V(t, v)=p v_{i}|v|^{p-2}, \quad \partial_{v_{i} v_{j}}^{2} V(t, v)=p(p-2) v_{i} v_{j}|v|^{p-4}+\delta_{i j} p|v|^{p-2}
$$

Next, we have

$$
\begin{equation*}
\langle G(v)-G(w), v-w\rangle=\langle B(v-w, v), v-w\rangle \leqslant|v-w|^{2}|v| \tag{4.44}
\end{equation*}
$$

which follows from (4.43) after some simple manipulations, ${ }^{16}$ so that

$$
\begin{align*}
\mathcal{L}^{(2)} V(t, v-w)= & p\langle-A v+A w-G(v)+G(w), v-w\rangle_{\mathbb{R}^{3}}|v-w|^{p-2} \\
& +\frac{1}{2} \sum_{i j=1}^{3}\left[\left(v_{i}-w_{i}\right)\left(\sigma_{i}(v)-\sigma_{i}(w)\right)\left(\sigma_{j}(v)-\sigma_{j}(w)\right)\left(v_{i}-w_{j}\right)\right. \\
& \left.p(p-2)|v-w|^{p-4}+\delta_{i j} p\left(\sigma_{i}(v)-\sigma_{i}(w)\right)\left(\sigma_{j}(v)-\sigma_{j}(w)\right)|v-w|^{p-2}\right] \\
\leqslant & p|v||v-w|^{p}-p \beta|v-w|^{p}+\frac{1}{2} \bar{\sigma}^{2} p(p-1)|v-w|^{p}, \tag{4.45}
\end{align*}
$$

where $\mathcal{L}^{(2)}$ is the two-point generator associated with (4.40). Next, choose $p$ such that for $\nu, w \in L^{p+1}\left(\Omega, \mathcal{F}_{-\infty}^{t}, \mathbb{P}\right)$ and $0<\mathbb{E}|\nu-w|^{p}$. Then, from Jensen's inequality we have

$$
\begin{equation*}
0<\mathbb{E}|v-w|^{p} \leqslant\left(\mathbb{E}|v-w|^{p+1}\right)^{p /(p+1)}<\infty, \tag{4.46}
\end{equation*}
$$

while the Hölder's inequality leads to (with $\left.\|X\|_{p}:=\left(\mathbb{E}\left[|X|^{p}\right]\right)^{1 / p}\right)$

$$
\begin{equation*}
\mathbb{E}\left[|v||v-w|^{p}\right]=\left\||v||v-w|^{p}\right\|_{1} \leqslant\|v\|_{p+1}\left(\mathbb{E}|v-w|^{p+1}\right)^{p / p+1} \tag{4.47}
\end{equation*}
$$

The bounds (4.46) and (4.47) imply that there exists a constant $1 \leqslant \mathfrak{C}_{p}<\infty$ such that

$$
\begin{equation*}
\left(\mathbb{E}|v-w|^{p+1}\right)^{p /(p+1)}=\mathfrak{C}_{p} \mathbb{E}|v-w|^{p} . \tag{4.48}
\end{equation*}
$$

Combining (4.48), (4.47), and (4.45) leads to

$$
\mathbb{E}\left[\mathcal{L}^{(2)} V(t, v-w)\right] \leqslant-\lambda_{p} \mathbb{E}[V(t, v-w)],
$$

where $\lambda_{p}=p\left(\mathfrak{C}_{A}-\frac{1}{2} \bar{\sigma}(p-1)-\mathfrak{C}_{p}\|\nu\|_{p+1}\right)$. Now, for $v_{t}=\phi(t, s, \omega, \xi), w_{t}=\phi(t, s, \omega, \eta)$ solving (4.40), we have from the above

$$
\begin{equation*}
\mathbb{E}\left[\mathcal{L}^{(2)} V(t, \phi(t, s, \xi)-\phi(t, s, \eta))\right] \leqslant-\lambda_{p}(t, s) \mathbb{E}[V(t, \phi(t, s, \xi)-\phi(t, s, \eta))], \tag{4.49}
\end{equation*}
$$

so that combining Itô's lemma

$$
\mathrm{d} \mathbb{E}[V(t, \phi(t, s, \xi)-\phi(t, s, \eta))]=\mathbb{E}\left[\mathcal{L}^{(2)} V(t, \phi(t, s, \xi)-\phi(t, s, \eta))\right],
$$

with (4.49) we obtain

$$
\begin{equation*}
\mathbb{E}[V(t, \phi(t, s, \xi)-\phi(t, s, \eta))] \leqslant \mathbb{E}[V(s, \xi-\eta)] \exp \left(-\int_{s}^{t} \lambda_{p}(r, s) \mathrm{d} r\right) . \tag{4.50}
\end{equation*}
$$

[^10]



Fig. 1 Illustration of some aspects of the dynamics of the stochastic Lorenz model with time-periodic forcing (4.40) with the flow of solutions $\{\phi(t, 0, \omega, \cdot), t \geqslant 0\}$ in two different regimes. The top row corresponds to the regime in which the time-periodic measure exists and is supported on stable random periodic orbits of (4.40); the inset of the top-right figure shows a finite sample from this measure on a Poincare section (i.e. on the subspace $\mathbb{R}^{3}$ of $\mathcal{X}=[0, \tau) \times \mathbb{R}^{3}$ ). The top-left inset illustrates the relationship in (4.50) in the case when $\lim _{t \rightarrow \infty} \lambda_{p}(t)<0$ and random periodic orbits exist (see text and Theorem 4.7); colours denote path-wise evolution of $|\phi(t, 0, \omega, \xi)-\phi(t, 0, \omega, \eta)|^{2}$ for fixed $\xi, \eta$, and the dotted black line denotes $\mathbb{E}|\phi(t, 0, \cdot, \xi)-\phi(t, 0, \cdot, \eta)|^{2}$. The bottom row illustrates a regime where $\lim _{t \rightarrow \infty}|\phi(t, 0, \omega, \xi)-\phi(t, 0, \omega, \eta)|>0$ and existence of random periodic orbits and periodic measures cannot be guaranteed. Parameters in (4.40) are: (top row) $\bar{\alpha}=7.3, \bar{\beta}=26, \bar{\gamma}=7, \bar{\varrho}=10$,
$\bar{f}=100, \bar{\delta}=0.9, \tau=1, \bar{\sigma}=0.2$, and (bottom row) $\bar{\alpha}=10, \bar{\beta}=1, \bar{\gamma}=8 / 3, \bar{\varrho}=28$,
$\bar{f}=23, \bar{\delta}=0.9, \tau=1, \bar{\sigma}=0.2$

Thus, in order for Assumption 4.4(ii) to hold, it is sufficient to require that

$$
\mathfrak{C}_{A}-\frac{1}{2} \bar{\sigma}(p-1)-\limsup _{(t-s) \rightarrow \infty} \frac{1}{t-s} \int_{s}^{t} \mathfrak{C}_{p}(u, s)\|\phi(u, s, \xi)\|_{p+1} \mathrm{~d} u>0
$$

Finally, we choose $p=2$, so that $V(t, v)=|v|^{2}$ and note that (see "Appendix A")

$$
\lim _{(t-s) \rightarrow \infty}\|\phi(t, s, \xi)\|_{3}=(27)^{1 / 6}\left(\frac{L_{b_{1}}+L_{\sigma}}{L_{b_{2}}-L_{\sigma}}\right)^{1 / 2}
$$

so that simple but tedious algebraic manipulations lead to

$$
\begin{equation*}
\mathfrak{C}_{A}-\frac{1}{2} \bar{\sigma}-\overline{\mathfrak{C}}\left(\frac{L_{b_{1}}+L_{\sigma}}{L_{b_{2}}-L_{\sigma}}\right)^{1 / 2}>0, \quad \overline{\mathfrak{C}}=(27)_{(t-s) \rightarrow \infty}^{1 / 6} \lim _{(1)} \frac{1}{t-s} \int_{s}^{t} \mathfrak{C}_{2}(u, s) \mathrm{d} u \tag{4.51}
\end{equation*}
$$

where $L_{b_{1}}, L_{b_{2}}$, and $L_{\sigma}$ for the system (4.40) are given in (4.42).
Therefore, by Theorem 4.7, we conclude that the time-periodically forced stochastic Lorenz equation (4.40) admits a family of periodic measures $\left\{\mu_{t}: t \in[0, \tau)\right\} \subset \mathcal{P}\left(\mathbb{R}^{3}\right)$
supported on stable periodic solutions of (4.40) lifted to $[0, \tau) \times \mathbb{R}^{3}$. For $\bar{\sigma} \neq 0$, Assumption 4.6 holds in addition to Assumption 4.4, and Theorem 4.11 implies existence of ergodic $\tau$-periodic measures $\tilde{\mu}_{t}=\delta_{t \bmod \tau} \otimes \mu_{t} \in \mathcal{P}\left([0, \tau) \times \mathbb{R}^{3}\right)$ in the sense of Definition 4.8. Numerical illustration of the convergence in (4.50) is provided in Fig. 1 alongside a sample from the density of the ergodic measure supported on the attractor containing stable random periodic orbits.

## 5 Linear response in the random time-periodic regime

In this section, we derive a general formula for the linear response function which characterises the change of a statistical observable in response to small perturbations of SDE dynamics with a time-periodic ergodic probability measure. The results presented below build on and extend the derivations obtained for time-dependent stochastic systems in [62]. First, we derive the linear response formula associated with perturbations of dynamics with time-periodic measures, and we represent it via formulas exploiting the asymptotic statistical properties of the unperturbed dynamics; in line with terminology from statistical physics, these are termed fluctuation-dissipation formulas. In Theorem 5.14, we derive the fluctuation-dissipation formulas in the case when only the unperturbed dynamics has a time-periodic ergodic probability measure. In Theorem 5.16 we consider the linear response associated with perturbations of dynamics with a time-periodic ergodic probability measure under stronger conditions when perturbed dynamics also has a time-periodic ergodic measure. During the revision of the manuscript, we become aware of related results derived independently in [19] for non-autonomous SDE's; those results are complementary to ours since they are confined to a finite time interval in the non-autonomous case with a restricted class of perturbations, and they do not deal with perturbations of asymptotically time-periodic ergodic measures. We conclude with some examples of the linear response for the periodically forced stochastic Lorenz model used earlier in Example 4.14. In principle, the results discussed below apply to a wider class of SDE's generating time-periodic measures under less stringent conditions than those in Assumption 4.4; however, establishing the existence and ergodicity of such measures in a more general setting is not trivial and it is beyond the scope of this work.

### 5.1 Set-up and assumptions

Consider the following SDE on $\mathbb{R}^{d}$ for $t \geqslant s, s \in \mathbb{R}$,

$$
\begin{equation*}
\mathrm{d} X_{t}^{\alpha}=\hat{b}\left(\alpha(t), t, X_{t}^{\alpha}\right) \mathrm{d} t+\hat{\sigma}\left(\alpha(t), t, X_{t}^{\alpha}\right) \mathrm{d} W_{t-s}, \quad X_{s}^{\alpha}=x \tag{5.1}
\end{equation*}
$$

where the maps $t \mapsto \hat{b}(0, t, \cdot)=b(t, \cdot), t \mapsto \hat{\sigma}(0, t, \cdot)=\sigma(t, \cdot)$, are $\tau$-periodic and coincide with the coefficients in $(2.1) ; \alpha(\cdot) \in \mathcal{C}_{\infty}^{1}(\mathbb{R} ; \mathbb{R})$ will be assumed sufficiently small in the sequel. Allowing for the explicit time dependence in $\alpha(t)$ enables one to consider timedependent changes in the coefficients of (5.1) relative to those in the original dynamics for $\alpha=0$; for example, one can think of changes in the 'climatological' forcing (e.g. [1,3, $34,36-39,59,62,64]$ ) which is relevant for considerations in atmosphere-ocean science.
Similar to Sect. 4, we consider the Wiener probability space $(\Omega, \mathcal{F}, \mathbb{P}), \Omega:=\mathcal{C}_{0}\left(\mathbb{R}, \mathbb{R}^{m}\right)$, with $\mathcal{F}$ the Borel $\mathfrak{S}$-algebra on $\Omega$, and the probability measure $\mathbb{P}$ on $(\Omega, \mathcal{F})$ induced by the $m$-dimensional Wiener process $W_{t}$. Furthermore, we assume that there exists a proper interval $\mathcal{A} \subseteq \mathbb{R}$ containing $\alpha=0$ such that, for all $\alpha \in \mathcal{A}$, the coefficients $\hat{b}(\alpha, t, \cdot)$,
$\hat{\sigma}(\alpha, t, \cdot)$ are sufficiently regular for (5.1) to generate a forward stochastic flow on $\mathbb{R}^{d}$ (see Theorem 3.5); i.e.

$$
X_{t}^{\alpha}(\omega)=\phi^{\alpha}(t, s, \omega, x), \quad s \leqslant t \quad \mathbb{P} \text { - a.s. }
$$

The forward stochastic flow $\left\{\phi^{\alpha}(t, s, \cdot, \cdot): s \leqslant t\right\}$ induced by (5.1) has a one-point generator

$$
\begin{equation*}
\mathcal{L}_{\alpha}=\sum_{i=1} \hat{b}_{i}(\alpha(t), t, x) \partial_{x_{i}}+\frac{1}{2} \sum_{i, j=1} \hat{a}_{i j}(\alpha(t), t, x) \partial_{x_{i} x_{j}}^{2} \quad \hat{a}:=\hat{\sigma} \hat{\sigma}^{T} \tag{5.2}
\end{equation*}
$$

As in the previous sections (see (4.2)), we lift the $\operatorname{SDE}(5.1)$ to $\mathbb{R} \times \mathbb{R}^{d}$ to represent the dynamics as

$$
\begin{equation*}
\mathrm{d} \tilde{X}_{t}^{\alpha}=\tilde{b}\left(\alpha(t), \tilde{X}_{t}^{\alpha}\right) \mathrm{d} t+\tilde{\sigma}\left(\alpha(t), \tilde{X}_{t}^{\alpha}\right) \mathrm{d} \tilde{W}_{t-s}, \quad \tilde{X}_{s}^{\alpha}=\tilde{x}, \quad s \leqslant t \tag{5.3}
\end{equation*}
$$

where $\tilde{x}=(s, x) \in \mathbb{R} \times \mathbb{R}^{d}$, and

$$
\tilde{b}(\alpha, \tilde{x})=(1, \hat{b}(\alpha, s, x))^{T}, \quad \tilde{\sigma}(\alpha, \tilde{x})=\left(\begin{array}{lc}
0 & 0 \\
0 & \hat{\sigma}(\alpha, s, x)
\end{array}\right) .
$$

Finally, shifting the time $t \rightarrow t+s$ allows to represent the dynamics in (5.3) for $t \in \mathbb{R}^{+}$ as

$$
\begin{equation*}
\mathrm{d} \tilde{X}_{t+s}^{\alpha}=\tilde{b}\left(\alpha(t+s), \tilde{X}_{t+s}^{\alpha}\right) \mathrm{d} t+\tilde{\sigma}\left(\alpha(t+s), \tilde{X}_{t+s}^{\alpha}\right) \mathrm{d} \tilde{W}_{t+s}, \quad \tilde{X}_{s}^{\alpha}=\tilde{x} \tag{5.4}
\end{equation*}
$$

where $\tilde{W}_{t+s}=\tilde{W}_{t+s}\left(\theta_{-s} \omega\right)$ due to (3.12). In what follows, we will always assume that the SDE with $\alpha=0$ satisfies the conditions of Theorem 4.7.

The lifted process $\tilde{\Phi}^{\alpha}: \mathbb{R}^{+} \times \Omega \times \mathbb{R} \times \mathbb{R}^{d} \rightarrow \mathbb{R} \times \mathbb{R}^{d}$, is defined analogously to (4.5); namely

$$
\begin{equation*}
\tilde{\Phi}^{\alpha}(t, \omega, \tilde{x}):=\left(t+s, \quad \phi^{\alpha}\left(t+s, s, \theta_{-s} \omega, x\right)\right), \quad \tilde{x}:=(s, x) \in \mathbb{R} \times \mathbb{R}^{d}, t \in \mathbb{R}^{+} \tag{5.5}
\end{equation*}
$$

Note that if the flow $\phi^{\alpha}$ is induced by the solutions of (5.1) with the coefficients $\hat{b}(\alpha(t), t, x)$ and $\hat{\sigma}(\alpha(t), t, x)$ which are time periodic for all $\alpha \in \mathcal{A}, \tilde{\Phi}^{\alpha}$ can be represented on a flat cylinder $[0, \hat{\tau}) \times \mathbb{R}^{d},[0, \hat{\tau}) \approx \mathbb{R} \bmod \hat{\tau}, 0<\hat{\tau}<\infty$, and the results of Sect. 4 hold; one obvious case is for $\alpha=0$ when, by construction, the coefficients are time periodic with period $\tau$. If $\hat{\tau}=\tau$ for any $\alpha \in \mathcal{A}$, both the unperturbed $(\alpha=0)$ and perturbed ( $\alpha \neq 0$ ) dynamics can be considered on the same cylinder. We will consider such a case in the last theorem of this section (Theorem 5.16).
Similarly, the transition evolutions $\left(\tilde{\mathcal{P}}_{t}^{\alpha}\right)_{t \in \mathbb{R}^{+}}$and their duals $\left(\tilde{\mathcal{P}}_{t}^{\alpha *}\right)_{t \in \mathbb{R}^{+}}$can be defined through (4.7)-(4.8) with the transition kernel $\tilde{P}^{\alpha}(\tilde{x} ; t, \tilde{A}):=\mathbb{P}\left(\left\{\omega: \tilde{\Phi}^{\alpha}(t, \omega, \tilde{x}) \in \tilde{A}\right\}\right)$; namely

$$
\begin{align*}
\tilde{\mathcal{P}}_{t}^{\alpha} \varphi(\tilde{x}) & :=\int_{\mathbb{R} \times \mathbb{R}^{d}} \varphi(\tilde{y}) \tilde{P}^{\alpha}(\tilde{x} ; t, \mathrm{~d} \tilde{y}), & \forall \varphi \in \mathbb{M}_{\infty}\left(\mathbb{R} \times \mathbb{R}^{d}\right),  \tag{5.6}\\
\tilde{\mu}_{t+r}^{\alpha}(\tilde{A})=\left(\tilde{\mathcal{P}}_{t}^{\alpha *} \tilde{\mu}_{r}^{\alpha}\right)(\tilde{A}) & :=\int_{\mathbb{R} \times \mathbb{R}^{d}} \tilde{P}^{\alpha}(\tilde{x} ; t, \tilde{A}) \tilde{\mu}_{r}^{\alpha}(\mathrm{d} \tilde{x}), & \forall \tilde{\mu}_{r}^{\alpha} \in \mathcal{P}\left(\mathbb{R} \times \mathbb{R}^{d}\right), r \in \mathbb{R}^{+}, \tag{5.7}
\end{align*}
$$

with the short-hand notation $\tilde{\mu}_{r}^{\alpha}(\mathrm{d} \tilde{x})=\delta_{r}(s) \mathrm{d} s \otimes \mu_{r}^{\alpha}(\mathrm{d} x)$ for probability measures in the skew-product fibre bundle on $\mathcal{P}\left(\mathbb{R} \times \mathbb{R}^{d}\right)$, where $\tilde{\mu}_{r}^{\alpha} \in \mathcal{P}\left(\mathbb{R} \times \mathbb{R}^{d}\right)$ and $\mu_{r}^{\alpha} \in \mathcal{P}\left(\mathbb{R}^{d}\right)$; see, for example, [23] for more details concerning the structure of skew-product fibre bundles of probability measures. The definition in (5.6) can be extended to $\varphi \in \mathbb{M}\left(\mathbb{R} \times \mathbb{R}^{d}\right)$ in a standard fashion.
The generator of the lifted one-point motion is given by $\tilde{\mathcal{L}}_{\alpha}=\partial_{s}+\mathcal{L}_{\alpha}$ (see Definition 3.6). By construction (see, for example, [62]), one can check that if $\mu_{t}^{\alpha} \in \mathcal{P}\left(\mathbb{R}^{d}\right)$ is a solution of forward Kolmogorov equation with the operator $\mathcal{L}_{\alpha}^{*}$, then $\tilde{\mu}_{t}^{\alpha}=\delta_{t} \otimes \mu_{t}^{\alpha}$ solves the lifted forward Kolmogorov equation with $\tilde{\mathcal{L}}_{\alpha}^{*}$ in the skew-product fibre bundle embedded in $\mathcal{P}\left(\mathbb{R} \times \mathbb{R}^{d}\right)$.

In the sequel, we derive fluctuation-dissipation formulas associated with the linear response for time-asymptotic SDE's dynamics in the random time-periodic regime with the ergodic measure $\overline{\tilde{\mu}} \in \mathcal{P}\left([0, \tau) \times \mathbb{R}^{d}\right)$, as in Theorem 4.11. We start with the definition of a linear response function which approximates changes in the statistical observables due to sufficiently small perturbations of the unperturbed dynamics with a time-periodic ergodic probability measure.

Definition 5.1 (Linear response function) Assuming that $\mathbb{E}\left[\varphi\left(\tilde{\Phi}^{\alpha}\right)\right] \in L^{1}\left(\tilde{\mu}_{0}\right)$, consider a family of statistical observables

$$
\begin{equation*}
\mathbb{F}_{\varphi}^{\tilde{\mu}_{0}}(t, \alpha)=\int_{[0, \tau) \times \mathbb{R}^{d}} \mathbb{E}\left[\varphi\left(\tilde{\Phi}^{\alpha}(t, \cdot, \tilde{x})\right)\right] \tilde{\mu}_{0}(\mathrm{~d} \tilde{x})=\left\langle\tilde{\mathcal{P}}_{t}^{\alpha} \varphi, \tilde{\mu}_{0}\right\rangle, \quad \varphi \in \mathcal{C}_{\tau}\left(\mathbb{R} \times \mathbb{R}^{d}\right) \tag{5.8}
\end{equation*}
$$

where

$$
\begin{equation*}
\mathcal{C}_{\tau}^{2}\left(\mathbb{R} \times \mathbb{R}^{d}\right):=\left\{\varphi \in \mathcal{C}^{2}\left(\mathbb{R} \times \mathbb{R}^{d}\right): \varphi(t+\tau, \cdot)=\varphi(t, \cdot), t \in \mathbb{R}\right\} \tag{5.9}
\end{equation*}
$$

the transition evolution $\left\{\tilde{\mathcal{P}}_{t}^{\alpha}\right\}_{t \in \mathbb{R}^{+}}$is induced by $\tilde{\Phi}^{\alpha}$ in (5.5), and $\tilde{\mu}_{0} \in \mathcal{P}\left([0, \tau) \times \mathbb{R}^{d}\right)$ is the probability measure on the initial condition in (5.4), which is assumed throughout to be given by the $\tau$-periodic ergodic measure associated with the 'unperturbed' dynamics with $\alpha=0$. If there exists a locally integrable function $\mathcal{R}_{\varphi}^{\tilde{\mu}_{0}}$ such that the Gateaux derivative of $\tilde{\mathbb{F}}_{\varphi}^{\tilde{\mu}_{0}}(\cdot, \alpha)$ at $\alpha=0$ satisfies

$$
\begin{equation*}
\Delta \mathbb{F}_{\varphi, \vartheta}^{\tilde{\mu}_{0}}(t):=\left.\frac{d}{d \varepsilon} \mathbb{F}_{\varphi}^{\tilde{\mu}_{0}}(t, \varepsilon \vartheta)\right|_{\varepsilon=0}=\int_{0}^{t} \mathcal{R}_{\varphi}^{\tilde{\mu}_{0}}(t-r, r) \vartheta(r) \mathrm{d} r, \tag{5.10}
\end{equation*}
$$

for $\vartheta \in \mathcal{C}_{\infty}^{1}\left(\mathbb{R}^{+}, \mathbb{R}\right), \vartheta(0)=0$, we say that $\mathcal{R}_{\varphi}^{\tilde{\mu}_{0}}$ is a linear response function due to perturbations of the statistical observable $\mathbb{F}_{\varphi}^{\tilde{\mu}_{0}}$.

In other words, $\mathcal{R}_{\varphi}^{\tilde{\mu}_{0}}$ can be defined if the functional $\mathbb{F}_{\varphi}^{\tilde{\mu}_{0}}(\cdot, \alpha)$ is Gateaux differentiable at $\alpha=0$ in the direction of $\vartheta$, and the Gateaux derivative is linear and continuous in the neighbourhood of $\alpha=0$. The formula (5.10) can be interpreted as an $\mathcal{O}(\varepsilon)$ approximation of the change of the statistical observable $\mathbb{F}_{\varphi}^{\tilde{\mu}_{0}}$ in response to a sufficiently small perturbation $\varepsilon \vartheta(t)$ around $\alpha=0$. The explicit time dependence in the perturbation $\vartheta(t)$ enables one to consider the linear response to small time-dependent changes in the coefficients of (5.4) relative to those in the original dynamics for $\alpha=0$; for example, one can consider changes in the climatological forcing (e.g. [1,3,34,36-39,59,62,64]).

Throughout the remainder of this section, we impose the following regularity conditions which reduce to Assumption 4.4 when $\alpha=0$, and which imply the smoothing property (e.g. [26]) of the transition evolutions $\left(\tilde{\mathcal{P}}_{t}^{\alpha}\right)_{t \in \mathbb{R}^{+}}$(in the $x$-component of the extended state space $\mathbb{R} \times \mathbb{R}^{d}$ ):

Assumption 5.2 Assume that there exists a proper interval $\mathcal{A} \subseteq \mathbb{R}$ containing $\alpha=0$, and that the following conditions are satisfied for all $s \leqslant t, s \in \mathbb{R}$, and for all $\alpha \in \mathcal{A}$ :
(i) $D_{\alpha}^{n} \hat{b}(\alpha, t, \cdot) \in \tilde{\mathcal{C}}^{\infty}\left(\mathbb{R}^{d}\right),\left|D_{\alpha}^{n} \hat{b}(\alpha, t, x)\right| \leqslant \mathfrak{C}_{\ell}\left(1+|x|^{\ell}\right), 0 \leqslant \ell<\infty, 0 \leqslant \mathfrak{C}_{\ell}<\infty$, $n \in \mathbb{N}_{0}$, and $t \mapsto \hat{b}(\cdot, t, \cdot)$ is differentiable on $\mathcal{A} \times \mathbb{R} \times \mathbb{R}^{d}$.
(ii) $D_{\alpha}^{n} \hat{\sigma}_{k}(\alpha, t, \cdot) \in \tilde{\mathcal{C}}_{b}^{\infty}\left(\mathbb{R}^{d}\right), n \in \mathbb{N}_{0}$, and $t \mapsto \hat{\sigma}_{k}(t, \cdot, \cdot)$ is differentiable on $\mathcal{A} \times \mathbb{R} \times \mathbb{R}^{d}$, and $\left|\partial_{t} \partial_{\alpha}^{n} D_{x}^{\beta} \hat{\sigma}_{k}(\alpha, t, x)\right|<\mathfrak{C}<\infty, \quad(\alpha, t, x) \in \mathcal{A} \times \mathbb{R} \times \mathbb{R}^{d}, \quad 1 \leqslant k \leqslant m$, for any multi-index $\beta$, and $\hat{\sigma}_{k}, 1 \leqslant k \leqslant m$ the columns of $\hat{\sigma}$.
(iii) $\operatorname{Lie}\left(\hat{\sigma}_{1}(\alpha, t, x), \cdots, \hat{\sigma}_{m}(\alpha, t, x)\right)=\mathbb{R}^{d}$, for all $s \leqslant t$, where

$$
\operatorname{Lie}\left(\hat{\sigma}_{1}, \cdots, \hat{\sigma}_{m}\right):=\operatorname{span}\left\{\hat{\sigma}_{i},\left[\hat{\sigma}_{i}, \hat{\sigma}_{j}\right],\left[\hat{\sigma}_{i},\left[\hat{\sigma}_{j}, \hat{\sigma}_{k}\right]\right], \cdots, 1 \leqslant i, j, k \leqslant m\right\}
$$

and $[F, G](\alpha, t, x)$ is the Lie bracket between the vector fields $F$ and $G$ defined by

$$
[F, G](\alpha, t, x):=D_{x} G(\alpha, t, x) F(\alpha, t, x)-D_{x} F(\alpha, t, x) G(\alpha, t, x) .
$$

Remark 5.3 Assumption 5.2, which is a version of the Hörmander condition, implies the existence of a smooth density of the time-marginal probability measure on $\left(\mathbb{R}^{d}, \mathcal{B}\left(\mathbb{R}^{d}\right)\right)$ induced by the law of the solutions of (5.1); i.e. $\mu_{t}^{\alpha}(\mathrm{d} x)=\rho_{t}^{\alpha}(x) \mathrm{d} x, \rho_{t}^{\alpha} \in \mathcal{C}^{\infty}\left(\mathbb{R}^{d}\right) \cap L_{+}^{1}\left(\mathbb{R}^{d}\right)$ for all $s \leqslant t$. In order to simplify the subsequent derivations, we will abuse notation and set the following

$$
\begin{equation*}
\tilde{\mu}_{t}^{\alpha}(\mathrm{d} \tilde{x}) \equiv \tilde{\rho}_{t}^{\alpha}(\tilde{x}) \mathrm{d} \tilde{x} \equiv \delta_{t}(s) \mathrm{d} s \otimes \rho_{t}^{\alpha}(x) \mathrm{d} x \tag{5.11}
\end{equation*}
$$

when dealing with probability measures $\tilde{\mu}_{t}^{\alpha} \in \mathcal{P}\left(\mathbb{R} \times \mathbb{R}^{d}\right)$, $\tilde{\mu}_{t}^{\alpha}=\delta_{t} \otimes \mu_{t}^{\alpha}$, of the lifted process in the skew-product fibre bundle associated with the SDE (5.4). This intuitive convention is consistent with the convention introduced in (4.6) for transition evolutions.

In addition to Assumption 5.2, the last theorem of this section, which is concerned with the linear response when both the unperturbed and perturbed measures are $\tau$-periodic and ergodic, will require the following assumption (which coincides with Assumption 4.4 for $\alpha=0$ ).

Assumption 5.4 Let $V \in \mathcal{C}^{1,2}\left(\mathbb{R} \times \mathbb{R}^{d}, \mathbb{R}^{+}\right)$such that $V(t, 0)=0$ for all $t \in \mathbb{R}$, and the coefficients $\hat{b}(\alpha, t, x)$, $\hat{\sigma}(\alpha, t, x)$, in (5.1) be such that for $\alpha \in \mathcal{A}$, where $\mathcal{A} \subseteq \mathbb{R}$ is a proper interval containing $\alpha=0$, the following hold:
(i) There exist functions $\lambda_{\alpha} \in L^{1}(\mathbb{R})$ with $\mathcal{A} \ni \alpha \mapsto \lambda_{\alpha}(t)$ bounded for $t \in \mathbb{R}$, and a constant $1 \leqslant \mathfrak{C}<\infty$ such that for all $\xi, \eta \in L^{p}\left(\Omega, \mathcal{F}_{-\infty}^{s}, \mathbb{P}\right)$ and some $1<p<\infty$ we
have

$$
\left\{\begin{array}{l}
\mathbb{E}|\xi|^{p} \leqslant \mathbb{E}[V(t, \xi)] \leqslant \mathfrak{C}|\xi|^{p},  \tag{5.12}\\
\mathbb{E}\left[\mathcal{L}_{\alpha}^{(2)} V(t, \xi-\eta)\right] \leqslant \lambda_{\alpha}(t) \mathbb{E}[V(t, \xi-\eta)]
\end{array}\right.
$$

where $\mathcal{L}_{\alpha}^{(2)}$ is defined analogously to the two-point generator $\mathcal{L}^{(2)}$ in (3.8) but based on the coefficients $\hat{b}(\alpha, t, \cdot), \hat{\sigma}(\alpha, t, \cdot)$ of (5.1); moreover, $\mathcal{L}_{\alpha}^{(2)} \equiv \mathcal{L}^{(2)}$ for $\alpha=0$.
(ii) There exists $\overline{\bar{\lambda}}>0$ such that

$$
\begin{equation*}
\sup _{\alpha \in \mathcal{A}}\left\{\limsup _{(t-s) \rightarrow \infty} \frac{1}{t-s} \int_{s}^{t} \lambda_{\alpha}(u) \mathrm{d} u\right\}<-\overline{\bar{\lambda}}<0 \tag{5.13}
\end{equation*}
$$

(iii) Given the one-point motion $t \mapsto \phi^{\alpha}(t, s, \omega, \xi)$ induced by (5.1) for $\omega \in \Omega, \xi \in \mathbb{R}^{d}$, and $s \leqslant t$, there exists $0<\mathfrak{D}_{\alpha}<\infty$ independent of $s, t \in \mathbb{R}$ such that ${ }^{17}$ for all $\xi \in L^{p}\left(\Omega, \mathcal{F}_{-\infty}^{s}, \mathbb{P}\right)$

$$
\begin{equation*}
\limsup _{(t-s) \rightarrow \infty} \mathbb{E}\left[V\left(t, \phi^{\alpha}(t, s, \xi)-\xi\right)\right]<\mathfrak{D}_{\alpha} \tag{5.14}
\end{equation*}
$$

where $\mathbb{E}\left[V\left(\phi^{\alpha}(t, s, \xi)\right)\right]:=\int_{\Omega} V\left(\phi^{\alpha}(t, s, \omega, \xi)\right) \mathbb{P}(\mathrm{d} \omega)$.

### 5.2 Preparatory lemmas

We start with the following standard and preparatory results which utilise relatively wellknown results from $[26,85]$, and are aimed at representing $\tilde{\mathcal{P}}_{t}^{\alpha} \varphi(\tilde{x})-\tilde{\mathcal{P}}_{t} \varphi(\tilde{x})$ in the form amenable to further analysis in the context of the linear response. The main results are derived in Sect. 5.3.

Lemma 5.5 Suppose that the coefficients $\hat{b}(\alpha, t, x), \hat{\sigma}(\alpha, t, x)$ in the $\operatorname{SDE}$ (5.1) satisfy Assumption 5.2 so that global solutions of (5.1) exist for all time, and $\varphi \in \mathcal{C}^{2}\left(\mathbb{R} \times \mathbb{R}^{d}\right)$ is such that for any fixed $\tilde{x}=(s, x) \in \mathbb{R} \times \mathbb{R}^{d}$ and for all $\alpha \in \mathcal{A}, \mathbb{E}\left[D_{x}^{\beta} \varphi\left(\tilde{\Phi}^{\alpha}(t, \tilde{x})\right)\right]<\infty$, $|\beta| \leqslant 2$, where $\left\{\tilde{\Phi}^{\alpha}(t, \cdot, \cdot): t \in \mathbb{R}^{+}\right\}$in (5.5) is generated by the lifted SDE (5.4) on $\mathbb{R} \times \mathbb{R}^{d}$.
Then, the function $v(r, \tilde{x}):=\tilde{\mathcal{P}}_{t-r}^{\alpha} \varphi(\tilde{x})$ with $\tilde{x}:=(s, x) \in \mathbb{R} \times \mathbb{R}^{d}$, and $\tilde{\mathcal{P}}_{t}^{\alpha}$ defined in (5.6), is the unique solution of the backward Kolmogorov equation with the terminal condition

$$
\left\{\begin{array}{l}
\partial_{r} v(r, \tilde{x})=-\tilde{\mathcal{L}}_{\alpha} v(r, \tilde{x}), \quad 0 \leqslant r \leqslant t  \tag{5.15}\\
v(t, \tilde{x})=\varphi(\tilde{x}) .
\end{array}\right.
$$

Moreover, for any $\varphi \in \mathcal{C}_{\infty}^{2}\left(\mathbb{R} \times \mathbb{R}^{d}\right)$, there exists a constant $\mathfrak{C}>0$ such that

$$
\begin{equation*}
\left\|\tilde{\mathcal{P}}_{t-r}^{\alpha} \varphi\right\|_{2, \infty} \leqslant \mathfrak{C}\|\varphi\|_{2, \infty}, \quad 0 \leqslant r \leqslant t \tag{5.16}
\end{equation*}
$$

where $\|\varphi\|_{2, \infty}=\|\varphi\|_{\infty}+\sum_{1 \leqslant|\beta| \leqslant 2}\left\|D_{x}^{\beta} \varphi\right\|_{\infty}$.

[^11]Proof See, for example, [54, Thm 4.8.11] or [17,26,85]. Sufficient conditions for $\mathbb{E}\left[D_{x}^{\beta} \varphi\left(\tilde{\Phi}^{\alpha}\right)\right]<\infty,|\beta| \leqslant 2, \varphi \in \mathcal{C}^{2}\left(\mathbb{R} \times \mathbb{R}^{d}\right), \alpha \in \mathcal{A}$, are given in Proposition 5.8. It can also be shown that, under Assumption 5.2, the problem (5.15) has unique classical solutions for $\varphi \in \mathcal{C}\left(\mathbb{R} \times \mathbb{R}^{d}\right)$ due to the smoothing property of $\left(\tilde{\mathcal{P}}_{t}^{\alpha}\right)_{t \in \mathbb{R}^{+}}$(e.g. [26]; the general case can be obtained by approximating $\varphi \in \mathcal{C}\left(\mathbb{R} \times \mathbb{R}^{d}\right)$ by $\varphi_{n} \in \mathcal{C}_{\infty}\left(\mathbb{R} \times \mathbb{R}^{d}\right)$ converging uniformly to $\varphi$ on compact subsets of $\mathbb{R} \times \mathbb{R}^{d}$.

Lemma 5.6 Suppose that the conditions of Lemma 5.5 hold. Then,

$$
\begin{equation*}
\tilde{\mathcal{P}}_{t}^{\alpha} \varphi(\tilde{x})-\tilde{\mathcal{P}}_{t} \varphi(\tilde{x})=\int_{0}^{t} \tilde{\mathcal{P}}_{r}\left(\tilde{\mathcal{L}}_{\alpha}-\tilde{\mathcal{L}}\right) \tilde{\mathcal{P}}_{t-r}^{\alpha} \varphi(\tilde{x}) \mathrm{d} r, \quad 0 \leqslant r \leqslant t, \quad \tilde{x} \in \mathbb{R} \times \mathbb{R}^{d} \tag{5.17}
\end{equation*}
$$

Proof It can be obtained from Lemma 5.5 that the function $u(r, \tilde{x})=\tilde{\mathcal{P}}_{t-r}^{\alpha} \varphi(\tilde{x})-\tilde{\mathcal{P}}_{t-r} \varphi(\tilde{x})$, $0 \leqslant r \leqslant t, \tilde{x} \in \mathbb{R} \times \mathbb{R}^{d}$, uniquely solves the inhomogeneous Cauchy problem (see, for example, $[17,26,85]$ )

$$
\left\{\begin{array}{l}
\partial_{r} u(r, \tilde{x})=-\tilde{\mathcal{L}} u(r, \tilde{x})-\alpha(r) \hat{\mathfrak{f}}_{t, \alpha}(r, \tilde{x})  \tag{5.18}\\
u(t, \tilde{x})=0
\end{array}\right.
$$

where $\alpha \in \mathcal{C}_{\infty}^{1}(\mathbb{R}, \mathbb{R}), \alpha(r) \in \mathcal{A}$, and

$$
\begin{equation*}
\alpha(r) \hat{\mathfrak{f}}_{t, \alpha}(r, \tilde{x}):=\left(\tilde{\mathcal{L}}_{\alpha}-\tilde{\mathcal{L}}\right) \tilde{\mathcal{P}}_{t-r}^{\alpha} \varphi(\tilde{x}) . \tag{5.19}
\end{equation*}
$$

Next, consider the solutions of (5.4) with $\alpha=0$ represented through (5.5) as $\tilde{\Phi}^{0}=\tilde{\Phi}(r, \omega, \tilde{x})$ where $\tilde{\Phi}$ is defined in (4.5) and solves the SDE (4.4). Then, by Itô's formula

$$
\begin{align*}
\mathrm{d} u\left(r, \tilde{\Phi}_{r}\right) & =\left[\partial_{r} u\left(r, \tilde{\Phi}_{r}\right)+\tilde{\mathcal{L}} u\left(r, \tilde{\Phi}_{r}\right)\right] \mathrm{d} r+D_{x} u\left(r, \tilde{\Phi}_{r}\right)^{T} \tilde{\sigma}\left(r, \tilde{\Phi}_{r}\right) \mathrm{d} \tilde{W}_{r} \\
& =-\alpha(r) \hat{f}_{t, \alpha}\left(r, \tilde{\Phi}_{r}\right) \mathrm{d} r+D_{x} u\left(r, \tilde{\Phi}_{r}\right)^{T} \tilde{\sigma}\left(r, \tilde{\Phi}_{r}\right) \mathrm{d} \tilde{W}_{r} \tag{5.20}
\end{align*}
$$

where $\tilde{\Phi}_{r} \equiv \tilde{\Phi}(r, \omega, \tilde{x})$ to simplify notation. Combining (5.20) with (5.18), and using the explicit form of $u(t, \tilde{x})$ leads to

$$
\tilde{\mathcal{P}}_{t}^{\alpha} \varphi(\tilde{x})-\tilde{\mathcal{P}}_{t} \varphi(\tilde{x})=\int_{0}^{t} \alpha(r) \mathbb{E}\left[\hat{\mathfrak{f}}_{t, \alpha}(r, \tilde{\Phi}(r, \tilde{x})] \mathrm{d} r=\int_{0}^{t} \tilde{\mathcal{P}}_{r}\left(\tilde{\mathcal{L}}_{\alpha}-\tilde{\mathcal{L}}\right) \tilde{\mathcal{P}}_{t-r}^{\alpha} \varphi(\tilde{x}) \mathrm{d} r\right.
$$

The above identity is well defined due to the underlying assumptions, and it is discussed further in Proposition 5.8.

Definition 5.7 ( $\alpha$-linearised generator) Given the infinitesimal generator $\tilde{\mathcal{L}}_{\alpha}=\partial_{s}+\mathcal{L}_{\alpha}$ with $\mathcal{L}_{\alpha}$ defined in (5.2) and ( $\hat{b}, \hat{\sigma}$ ) satisfying Assumption 5.2, the $\alpha$-linearised generator is defined by

$$
\begin{equation*}
\tilde{\mathcal{V}} \varphi(\tilde{x})=\mathfrak{b}(\tilde{x}) D_{x} \varphi(\tilde{x})+\frac{1}{2} \operatorname{Tr}\left(\mathfrak{a}(\tilde{x}) D_{x}^{2} \varphi(\tilde{x})\right), \quad \tilde{x}=(s, x), \quad \varphi \in \mathcal{C}^{2}\left(\mathbb{R} \times \mathbb{R}^{d}\right) \tag{5.21}
\end{equation*}
$$

where

$$
\mathfrak{b}(\tilde{x})=\left.\partial_{\alpha} \hat{b}(\alpha, s, x)\right|_{\alpha=0^{\prime}} \quad \mathfrak{a}(\tilde{x}):=\sigma(s, x) H^{T}(s, x)+\sigma^{T}(s, x) H(s, x),
$$

with $H_{i k}(s, x)=\left.\partial_{\alpha} \hat{\sigma}_{i k}(\alpha, s, x)\right|_{\alpha=0}, 1 \leqslant k \leqslant m, 1 \leqslant i \leqslant d$, and $\hat{\sigma}(0, t, \cdot)=\sigma(t, \cdot)$ as in (5.1). The $L^{2}\left(\tilde{\mu}_{t}\right)$ dual of $\tilde{\mathcal{V}}$ is given by

$$
\begin{equation*}
\tilde{\mathcal{V}}^{*} \tilde{\rho}_{t}(\tilde{x})=-D_{x}\left(\mathfrak{b}(\tilde{x}) \tilde{\rho}_{t}(\tilde{x})\right)+\frac{1}{2} \operatorname{Tr}\left(D_{x}^{2}\left(\mathfrak{a}(\tilde{x}) \tilde{\rho}_{t}(\tilde{x})\right)\right) \tag{5.22}
\end{equation*}
$$

where $\tilde{\rho}_{t}(\tilde{x})$ is understood in the sense of (5.11).
Note that the properties of $\mathfrak{a}(\tilde{x})$ and $\mathfrak{b}(\tilde{x})$ are fully controlled through Assumption 5.2.
Proposition 5.8 Suppose that the conditions of Lemma 5.5 are satisfied, and consider a function $\mathfrak{f}_{t, \alpha}:[0, T) \times \mathbb{R} \times \mathbb{R}^{d} \rightarrow \mathbb{R}, \alpha \in \mathcal{A}$, defined by

$$
\begin{equation*}
\mathfrak{f}_{t, \alpha}(r, \tilde{x})=\tilde{\mathcal{V}} \tilde{\mathcal{P}}_{t-r}^{\alpha} \varphi(\tilde{x}), \quad 0 \leqslant r \leqslant t \leqslant T \tag{5.23}
\end{equation*}
$$

Then, $\mathfrak{f}_{t, \alpha} \in \mathcal{C}\left(\mathbb{R} \times \mathbb{R}^{d}\right)$ and $\mathfrak{f}_{t, \alpha}<\infty$ for any fixed $\tilde{x} \in \mathbb{R} \times \mathbb{R}^{d}$.
If Assumption 5.2 holds for $0 \leqslant \ell<\infty$, and the initial condition in the lifted SDE (5.4) has $p \geqslant \max (2, \ell)$ finite moments, then there exists a constant $\mathfrak{C}=\mathfrak{C}(T, k, \varphi)>0$ such that for any fixed $\tilde{x} \in \mathbb{R} \times \mathbb{R}^{d}$ one has

$$
\begin{equation*}
\sup _{0 \leqslant r \leqslant t \leqslant T} \mathbb{E}\left|\mathfrak{f}_{t, \alpha}(r, \tilde{\Phi}(r, \tilde{x}))\right| \leqslant \mathfrak{C}<\infty \tag{5.24}
\end{equation*}
$$

where $\left\{\tilde{\Phi}(r, \cdot, \cdot): r \in \mathbb{R}^{+}\right\}$is the RDS (4.5) generated by the $\operatorname{SDE}(5.4)$ with $\alpha=0$.
The sufficient condition for (5.24) to hold for $T \rightarrow \infty$ is that $(\hat{b}, \hat{\sigma})$ satisfy the dissipative conditions (4.13)-(4.14) with $p=\max (2, \ell)$. If $\varphi(\cdot, x) \leqslant \mathfrak{C}_{l}\left(1+|x|^{l}\right), 0 \leqslant \ell<\infty, 0 \leqslant \mathfrak{C}_{\ell}<\infty$, then (5.24) holds for $T \rightarrow \infty$ when ( $\hat{b}, \hat{\sigma})$ satisfy (4.13)-(4.14) with $p=\max (2, \ell, l)$; i.e. the conditions $\mathbb{E}\left[D_{x}^{\beta} \varphi\left(\tilde{\Phi}^{\alpha}\right)\right]<\infty,|\beta| \leqslant 2$, in Lemma 5.5 can be replaced by assuming a polynomial growth of $\varphi$.

Remark 5.9 Note that for a dissipative dynamics satisfying (4.13) the dissipation coefficient $L_{b_{2}}$ might not be large enough to satisfy (4.14) with a given $p \geqslant 2$. Thus, not all dissipative dynamics automatically satisfy Proposition 5.8 for all time.

Proof For $(s, x) \in \mathbb{R} \times \mathbb{R}^{d}$, one can obtain directly from (5.21) that

$$
\begin{equation*}
\mathfrak{f}_{t, \alpha}(r, \tilde{x})=\tilde{\mathcal{V}} \tilde{\mathcal{P}}_{t-r}^{\alpha} \varphi(\tilde{x})=\mathfrak{b}(\tilde{x}) D_{x} \tilde{\mathcal{P}}_{t-r}^{\alpha} \varphi(\tilde{x})+\frac{1}{2} \operatorname{Tr}\left(\mathfrak{a}(\tilde{x}) D_{x}^{2} \tilde{\mathcal{P}}_{t-r}^{\alpha} \varphi(\tilde{x})\right) \tag{5.25}
\end{equation*}
$$

The regularity and growth conditions of the coefficients $(\hat{b}, \hat{\sigma})$ imposed in Assumption 5.2 ensure the existence of global solutions to (5.4) which are represented via $\tilde{\Phi}^{\alpha}(t-r, \cdot, \tilde{x})$, $t-r \in \mathbb{R}^{+}$, and generate $\tilde{\mathcal{P}}_{t-r}^{\alpha}$ in (5.6). If $p \geqslant 2$ moments of the initial condition of (5.4) are finite and $\mathbb{E}\left[D_{x}^{\beta} \varphi\left(\tilde{\Phi}^{\alpha}(t, \tilde{x})\right)\right]<\infty,|\beta| \leqslant 2, t \leqslant T$, then by the assumption on $(\hat{b}, \hat{\sigma})$ and Lemma 5.5 , we have $\mathfrak{f}_{t, \alpha} \in \mathcal{C}\left(\mathbb{R} \times \mathbb{R}^{d}\right)$ and $\mathfrak{f}_{t, \alpha}<\infty$ for any fixed $\tilde{x} \in \mathbb{R} \times \mathbb{R}^{d}$.
Regarding (5.24), the polynomial growth of $\hat{b}$ combined with the standard calculation utilising Itô's formula guarantees the existence of $\max (2, \ell)$ finite moments of the solution for $T<\infty$ (Theorem 3.5). Thus, (5.24) follows by the Cauchy-Schwarz inequality applied to $\mathbb{E}\left|\mathfrak{f}_{t, \alpha}\right|$ and the finiteness of the moments of $\mathfrak{a}, \mathfrak{b}$ for $T<\infty$.

Considering (5.24) for $T \rightarrow \infty$, may require additional dissipative constraints on the drift $\hat{b}$, as outlined below. Moreover, we specify two explicit classes of $\varphi \in \mathcal{C}_{\infty}^{2}\left(\mathbb{R} \times \mathbb{R}^{d}\right)$ for
which $\mathbb{E}\left[D_{x}^{\beta} \varphi\left(\tilde{\Phi}^{\alpha}(t, \tilde{x})\right)\right]<\infty,|\beta| \leqslant 2$, (5.24) holds for all time. First, if $\varphi \in \mathcal{C}_{\infty}^{2}\left(\mathbb{R} \times \mathbb{R}^{d}\right)$, then

$$
\begin{equation*}
\sup _{0 \leqslant r \leqslant t \leqslant T} \mathbb{E} \mid \mathfrak{f}_{t, \alpha}\left(r, \tilde{\Phi}(r, \tilde{x}) \mid \leqslant C_{T} \sup _{0 \leqslant r \leqslant t \leqslant T}\left\|\tilde{\mathcal{P}}_{t-r}^{\alpha} \varphi\right\|_{2, \infty}\right. \tag{5.26}
\end{equation*}
$$

By the second part of Lemma 5.5, the term on the right of (5.26) is bounded by $C\|\varphi\|_{2, \infty}$ so that one can set explicitly $\mathfrak{C}=C_{T} C\|\varphi\|_{2, \infty}$ in (5.24). Given the existence of global solutions of (5.4) for all time (Assumption 5.2 and Theorem 3.5), the bound (5.24) can be extended to $T \rightarrow \infty$ provided that the dissipative conditions (4.13)-(4.14) hold for ( $\hat{b}, \hat{\sigma}$ ) with $p=\max (2, \ell)$; so that $\mathbb{E}|\mathfrak{b}(\tilde{\Phi}(r, \tilde{x}))|<\infty, \mathbb{E}|\mathfrak{a}(\tilde{\Phi}(r, \tilde{x}))|<\infty, r \in \mathbb{R}^{+}$, in (5.25). More generally, if $\varphi \in \mathcal{C}^{2}\left(\mathbb{R} \times \mathbb{R}^{d}\right)$ and $\varphi(\cdot, x) \leqslant \mathfrak{C}_{l}\left(1+|x|^{l}\right)$ with $0 \leqslant \ell<\infty, 0 \leqslant \mathfrak{C}_{\ell}<\infty$, the bound (5.24) can be extended to $T \rightarrow \infty$ provided that the dissipative conditions (4.13)-(4.14) hold for $(\hat{b}, \hat{\sigma})$ with $p=\max (2, \ell, l)$. Both assertions can be obtained through derivations analogous to Lemma A. 1 in "Appendix A".

Corollary 5.10 Let $\hat{\mathrm{f}}_{t, \alpha}:[0, T) \times \mathbb{R} \times \mathbb{R}^{d} \rightarrow \mathbb{R}, \alpha \in \mathcal{A}$, be defined by

$$
\begin{equation*}
\alpha(r) \hat{f}_{t, \alpha}(r, \tilde{x}):=\left(\tilde{\mathcal{L}}_{\alpha}-\tilde{\mathcal{L}}\right) \tilde{\mathcal{P}}_{t-r}^{\alpha} \varphi(\tilde{x}), \quad \tilde{x} \in \mathbb{R} \times \mathbb{R}^{d}, \quad 0 \leqslant r \leqslant t \leqslant T \tag{5.27}
\end{equation*}
$$

as in (5.19) of Lemma 5.6; so that $\hat{\mathfrak{f}}_{t, \alpha}(r, \tilde{x})=\mathfrak{f}_{t, \alpha}(r, \tilde{x})+\mathcal{O}(\alpha)$ with $\mathfrak{f}_{t, \alpha}$ defined in Proposition 5.8. Then, under the same conditions as those in Proposition 5.8, $\hat{\mathfrak{f}}_{t, \alpha} \in \mathcal{C}\left(\mathbb{R} \times \mathbb{R}^{d}\right)$ and $\hat{\mathfrak{f}}_{t, \alpha}<\infty$ for any fixed $\tilde{x} \in \mathbb{R} \times \mathbb{R}^{d}$. Furthermore, for $p \geqslant 2$ chosen as in Proposition 5.8 one has

$$
\begin{equation*}
\sup _{0 \leqslant r \leqslant t \leqslant T} \mathbb{E}\left|\hat{f}_{t, \alpha}(r, \tilde{\Phi}(r, \tilde{x}))\right|<\infty, \quad \tilde{x} \in \mathbb{R} \times \mathbb{R}^{d} \tag{5.28}
\end{equation*}
$$

which can be extended to $T \rightarrow \infty$ in a way analogous to that in Proposition 5.8.
Proof This is a direct consequence of Proposition 5.8 and the fact that the $\mathcal{O}(\alpha)$ terms involve $D_{x} \tilde{\mathcal{P}}_{t-r}^{\alpha} \varphi(\tilde{x}), D_{x}^{2} \tilde{\mathcal{P}}_{t-r}^{\alpha} \varphi(\tilde{x})$, with coefficients given by $D_{\alpha}^{n} \hat{b}$ and $D_{\alpha}^{n} \hat{\sigma}$ which are controlled through Assumption 5.2 and the polynomial bound on the growth of the $\alpha$-derivatives of $\hat{b}$.

### 5.3 Linear response and fluctuation-dissipation formulas for time-periodic measures

Here, in Sect. 5.3.1 we derive a general expression for the linear response function characterising the change in the statistical observable (5.8) to small perturbations of dynamics of an SDE whose time-asymptotic dynamics is characterised by time-periodic ergodic probability measures (see Sect. 4). This is followed in Sect. 5.3 .2 by deriving a more tractable representation of the response function in terms of fluctuation-dissipation-type formulas which allow one to express the change in the statistical observables through statistical characteristics of the unperturbed dynamics.

### 5.3.1 The linear response

First, we derive a general formula for the linear response function associated with perturbations of the time-asymptotic dynamics of the SDE (4.4). The derived formula is equivalent to the one obtained formally in [62].

Theorem 5.11 (Linear response) Suppose that Assumption 5.2 holds, Assumption 5.4 holds for $\alpha=0$, and that Proposition 5.8 is satisfied. Consider the family of transition semigroups $\left\{\tilde{\mathcal{P}}_{t}^{\alpha}:(t, \alpha) \in \mathbb{R}^{+} \times \mathcal{A}\right\}$ induced by the $\operatorname{SDE}$ (5.4) which for $\alpha=0$ admit $\tau$-periodic ergodic measures $\left(\tilde{\mu}_{t}\right)_{t \geqslant 0}, \tilde{\mu}_{t} \in \mathcal{P}\left([0, \tau) \times \mathbb{R}^{d}\right)$.

Then, given the observable $\mathbb{F}_{\varphi}^{\tilde{\mu}_{0}}(t, \alpha)=\left\langle\tilde{\mathcal{P}}_{t}^{\alpha} \varphi, \tilde{\mu}_{0}\right\rangle$ in (5.8), for any $\varphi \in \mathcal{C}_{\tau}^{2}\left(\mathbb{R} \times \mathbb{R}^{d}\right)$ such that $\mathbb{E}\left[D_{x}^{\beta} \varphi\left(\tilde{\Phi}^{\alpha}(t, \cdot)\right)\right] \in L^{1}\left(\tilde{\mu}_{0}\right),|\beta| \leqslant 2$, and the perturbation $\alpha(\cdot)=\varepsilon \vartheta(\cdot) \in \mathcal{C}_{\infty}^{1}\left(\mathbb{R}^{+}, \mathbb{R}\right)$, $\vartheta(0)=0$, such that $\varepsilon \vartheta \in \mathcal{A}$, the following holds

$$
\begin{equation*}
\Delta \mathbb{F}_{\varphi, \vartheta}^{\tilde{\mu}_{0}}(t)=\int_{0}^{t} \mathcal{R}_{\varphi}^{\tilde{\mu}_{0}}(t-r, r) \vartheta(r) \mathrm{d} r, \tag{5.29}
\end{equation*}
$$

with the linear response function given by

$$
\begin{equation*}
\mathcal{R}_{\varphi}^{\tilde{\mu}_{0}}(t-r, r)=\int_{[0, \tau) \times \mathbb{R}^{d}} \tilde{\mathcal{P}}_{t-r} \varphi(\tilde{x})\left(\tilde{\mathcal{V}}^{*} \tilde{\rho}_{r}\right)(\tilde{x}) d \tilde{x}, \tag{5.30}
\end{equation*}
$$

where $\tilde{\mathcal{V}}^{*}$ is defined in (5.22), and $\tilde{\rho}_{r}(\tilde{x}) \mathrm{d} \tilde{x}$ is understood in the sense of (5.11).
Remark $5.12 \Delta \mathbb{F}_{\varphi, \vartheta}^{\tilde{\mu}_{0}}(t)$ can be interpreted as the approximate $\mathcal{O}(\varepsilon)$ response of the observable $\mathbb{F}_{\varphi}^{\tilde{\mu}_{0}}(t, 0)$ in (5.4) to a sufficiently small perturbation $\alpha(t)=\varepsilon \vartheta(t), \vartheta(0)=0$. Note that the linear response formula is solely based on quantities defined for the 'unperturbed' dynamics $(\alpha=0)$. Moreover, the perturbation $\alpha(t)=\varepsilon \vartheta(t)$ does not necessarily have to factorise the coefficients of (5.1) as, for example, $\hat{b}(\alpha(t), t, x)=\hat{b}(0, t, x)+\varepsilon \mathfrak{F}(x) \vartheta(t)$. For example, consider the following $\hat{b}(\alpha(t), t, x)=-\left(2-\sin ^{2}(x(1+\alpha(t))) x \cos ^{2}(t)\right.$ in (5.1) and the derivations below.

Proof First, we set $\alpha=\varepsilon \vartheta \in \mathcal{A}$ and show that under the assumptions of the proposition the following holds for any $\varphi \in \mathcal{C}_{\tau}^{2}\left(\mathbb{R} \times \mathbb{R}^{d}\right)$ such that $\mathbb{E}\left[D_{x}^{\beta} \varphi\left(\tilde{\Phi}^{\alpha}(t, \cdot)\right)\right] \in L^{1}\left(\tilde{\mu}_{0}\right),|\beta| \leqslant 2$ (see Proposition 5.8 for a sufficient condition for this to hold for all time),

$$
\lim _{\varepsilon \rightarrow 0} \frac{1}{\varepsilon}\left(\tilde{\mathcal{P}}_{t}^{\varepsilon \vartheta(t)} \varphi(\tilde{x})-\tilde{\mathcal{P}}_{t} \varphi(\tilde{x})\right)=\int_{0}^{t} \vartheta(r) \tilde{\mathcal{P}}_{r}\left(\tilde{\mathcal{V}} \tilde{\mathcal{P}}_{t-r} \varphi\right)(\tilde{x}) \mathrm{d} r, \quad 0 \leqslant r \leqslant t, \vartheta \in \mathcal{C}_{\infty}^{1}\left(\mathbb{R}^{+}, \mathbb{R}\right)
$$

with $\tilde{\mathcal{V}}$ defined in (5.21). To this end, it follows from Lemma 5.6 that for $\varepsilon>0$ sufficiently small

$$
\frac{1}{\varepsilon}\left(\tilde{\mathcal{P}}_{t}^{\varepsilon \vartheta(t)} \varphi(\tilde{x})-\tilde{\mathcal{P}}_{t} \varphi(\tilde{x})\right)=\frac{1}{\varepsilon} \int_{0}^{t} \tilde{\mathcal{P}}_{r}\left(\tilde{\mathcal{L}}_{\varepsilon \vartheta(r)}-\tilde{\mathcal{L}}\right) \tilde{\mathcal{P}}_{t-r}^{\varepsilon \vartheta(r)} \varphi(\tilde{x}) \mathrm{d} r=\int_{0}^{t} \vartheta(r) \mathbb{E}\left[\hat{\mathfrak{f}}_{t, \varepsilon}(r, \tilde{\Phi}(r, \tilde{x}))\right] \mathrm{d} r
$$

where $\tilde{\Phi}(r, \omega, \tilde{x})$ represents the solution of (5.4) with $\alpha=0$ (or the solution of (4.4)).
By Proposition 5.8, Corollary 5.10, and the dominated convergence theorem, we have

$$
\begin{aligned}
\lim _{\varepsilon \rightarrow 0} \frac{1}{\varepsilon}\left(\tilde{\mathcal{P}}_{t}^{\varepsilon \vartheta(t)} \varphi(\tilde{x})-\tilde{\mathcal{P}}_{t} \varphi(\tilde{x})\right) & =\int_{0}^{t} \vartheta(r) \mathbb{E}\left[\lim _{\varepsilon \rightarrow 0} \hat{\mathrm{f}}_{t, \varepsilon}(r, \tilde{\Phi}(r, \tilde{x})] \mathrm{d} r\right. \\
& =\int_{0}^{t} \vartheta(r) \mathbb{E}\left[\tilde{\mathcal{V}} \tilde{\mathcal{P}}_{t-r} \varphi(\tilde{\Phi}(r, \tilde{x})]\right. \\
& =\int_{0}^{t} \vartheta(r) \tilde{\mathcal{P}}_{r}\left(\tilde{\mathcal{V}} \tilde{\mathcal{P}}_{t-r} \varphi(\tilde{x})\right) \mathrm{d} r .
\end{aligned}
$$

Using Fubini's theorem and Proposition 5.8, we have

$$
\begin{aligned}
\lim _{\varepsilon \rightarrow 0} \frac{1}{\varepsilon}\left\langle\tilde{\mathcal{P}}_{t}^{\varepsilon \vartheta(t)} \varphi-\tilde{\mathcal{P}}_{t} \varphi, \tilde{\mu}_{0}\right\rangle & =\int_{0}^{t} \vartheta(r) \int_{[0, \tau) \times \mathbb{R}^{d}} \tilde{\mathcal{P}}_{r}\left(\tilde{\mathcal{V}} \tilde{\mathcal{P}}_{t-r} \varphi(\tilde{x})\right) \tilde{\mu}_{0}(\mathrm{~d} \tilde{x}) \mathrm{d} r \\
& =\int_{0}^{t} \vartheta(r) \int_{[0, \tau) \times \mathbb{R}^{d}} \tilde{\mathcal{V}} \tilde{\mathcal{P}}_{t-r} \varphi(\tilde{x})\left(\tilde{\mathcal{P}}_{r}^{*} \tilde{\mu}_{0}\right)(\mathrm{d} \tilde{x}) \mathrm{d} r \\
& =\int_{0}^{t} \vartheta(r) \int_{[0, \tau) \times \mathbb{R}^{d}} \tilde{\mathcal{V}} \tilde{\mathcal{P}}_{t-r} \varphi(\tilde{x}) \tilde{\mu}_{r}(\mathrm{~d} \tilde{x}) \mathrm{d} r .
\end{aligned}
$$

By the Hörmander Lie bracket condition in Assumption 5.2, and ergodicity of the timeperiodic measures $\tilde{\mu}_{r}$, there exists $0<\rho_{r} \in \mathcal{C}^{\infty}\left(\mathbb{R}^{d}\right) \cap L_{+}^{1}\left(\mathbb{R}^{d}\right)$ such that $\left(\tilde{\mathcal{P}}_{r}^{*} \tilde{\mu}_{0}\right)(\mathrm{d} \tilde{x})=$ $\tilde{\mu}_{r}(\mathrm{~d} \tilde{x})=\tilde{\rho}_{r}(\tilde{x}) \mathrm{d} \tilde{x}$, where $\tilde{\rho}_{r}(\tilde{x}) \mathrm{d} \tilde{x}$ is understood in the sense of (5.11). Thus, for any $\varphi \in \mathcal{C}_{\tau}^{2}\left(\mathbb{R} \times \mathbb{R}^{d}\right)$, such that $\mathbb{E}\left[D_{x}^{\beta} \varphi\left(\tilde{\Phi}^{\alpha}(t, \cdot)\right)\right] \in L^{1}\left(\tilde{\mu}_{0}\right),|\beta| \leqslant 2$, we have

$$
\begin{aligned}
\lim _{\varepsilon \rightarrow 0} \frac{1}{\varepsilon}\left\langle\tilde{\mathcal{P}}_{t}^{\varepsilon \vartheta(t)} \varphi-\tilde{\mathcal{P}}_{t} \varphi, \tilde{\mu}_{0}\right\rangle & =\int_{0}^{t} \vartheta(r) \int_{[0, \tau) \times \mathbb{R}^{d}} \tilde{\mathcal{V}} \tilde{\mathcal{P}}_{t-r} \varphi(\tilde{x}) \tilde{\rho}_{r}(\tilde{x}) \mathrm{d} \tilde{x} \mathrm{~d} r \\
& =\int_{0}^{t}\left(\int_{[0, \tau) \times \mathbb{R}^{d}} \tilde{\mathcal{P}}_{t-r} \varphi(\tilde{x})\left(\tilde{\mathcal{V}}^{*} \tilde{\rho}_{r}\right)(\tilde{x}) \mathrm{d} \tilde{x}\right) \vartheta(r) \mathrm{d} r .
\end{aligned}
$$

Finally, by the definition of the response functional $\mathcal{R}_{\varphi}^{\tilde{\mu}_{0}}$ we have for $\vartheta \in \mathcal{C}_{\infty}^{1}\left(\mathbb{R}^{+}, \mathbb{R}\right)$, that

$$
\begin{aligned}
\int_{0}^{t} \mathcal{R}_{\varphi}^{\tilde{\mu}_{0}}(t-r, r) \vartheta(r) \mathrm{d} r & =\lim _{\varepsilon \rightarrow 0} \frac{1}{\varepsilon}\left[\mathbb{F}_{\varphi}^{\tilde{\mu}_{0}}(t, \varepsilon \vartheta(t))-\mathbb{F}_{\varphi}^{\tilde{\mu}_{0}}(t, 0)\right] \\
& =\lim _{\varepsilon \rightarrow 0} \frac{1}{\varepsilon}\left\langle\tilde{\mathcal{P}}_{t}^{\varepsilon \vartheta(t)} \varphi-\tilde{\mathcal{P}}_{t} \varphi, \tilde{\mu}_{0}\right\rangle \\
& =\int_{0}^{t}\left(\int_{[0, \tau) \times \mathbb{R}^{d}} \tilde{\mathcal{P}}_{t-r} \varphi(\tilde{x})\left(\tilde{\mathcal{V}}^{*} \tilde{\rho}_{r}\right)(\tilde{x}) \mathrm{d} \tilde{x}\right) \vartheta(r) \mathrm{d} r
\end{aligned}
$$

where $\tilde{\rho}_{r}(\tilde{x}) \mathrm{d} \tilde{x}$ is understood in the sense of (5.11).

### 5.3.2 Fluctuation-dissipation formulas

Given the general framework for the linear response in the time-periodic regime, we now derive a set of more tractable representations of the response function (5.30) via formulas exploiting the time-asymptotic statistical properties of the unperturbed dynamics (4.4), or (5.1) with $\alpha=0$; in line with the terminology inherited from statistical physics, these are termed 'fluctuation-dissipation' formulas. The first set of results in Theorem 5.14 shadows and formalises formulas derived in [62], while the results in a more restrictive Theorem 5.16 concern the linear response in situations when the ' $\alpha$-perturbations' do not destroy the time periodicity and ergodicity of the dynamics in the sense that the coefficients in (5.1) remain $\tau$-periodic for all $\alpha \in \mathcal{A}$. It turns out the two results are related in a specific way.

Definition 5.13 (Correlation function) Given the $\operatorname{RDS}\left\{\tilde{\Phi}(t, \cdot, \cdot): t \in \mathbb{R}^{+}\right\}$on $[0, \tau) \times \mathbb{R}^{d}$, and $\varphi, \psi \in \mathcal{C}_{\tau}^{2}\left(\mathbb{R} \times \mathbb{R}^{d}\right)$, such that for any fixed $\tilde{x} \in[0, \tau) \times \mathbb{R}^{d}$, $\mathbb{E}[\varphi(\tilde{\Phi}(t, \tilde{x}))], \mathbb{E}[\psi(\tilde{\Phi}(t, \tilde{x}))]<\infty$, the correlation of the random variables $\varphi(\tilde{\Phi}(t, \cdot, \tilde{x}))$ and $\psi(\tilde{\Phi}(r, \cdot, \tilde{x}))$ is given by

$$
\begin{equation*}
\mathbb{E}[\varphi(\tilde{\Phi}(t, \tilde{x})) \psi(\tilde{\Phi}(r, \tilde{x}))]=\tilde{\mathcal{P}}_{r}\left(\psi(\tilde{x}) \tilde{\mathcal{P}}_{t-r} \varphi(\tilde{x})\right), \quad 0 \leqslant r \leqslant t \tag{5.31}
\end{equation*}
$$

where $\left(\tilde{\mathcal{P}}_{t}\right)_{t \in \mathbb{R}^{+}}$is defined in (4.7), and (5.31) follows from the Markov property of the $\operatorname{RDS} \tilde{\Phi}$. The correlation function based on $\varphi, \psi \in \mathcal{C}_{\tau}^{2}\left(\mathbb{R} \times \mathbb{R}^{d}\right)$ and $\tilde{\mu} \in \mathcal{P}\left([0, \tau) \times \mathbb{R}^{d}\right)$ is defined as

$$
\begin{equation*}
\mathcal{K}_{\varphi, \psi}^{\tilde{\mu}}(t-r, r):=\int_{[0, \tau) \times \mathbb{R}^{d}} \tilde{\mathcal{P}}_{r}\left(\psi \tilde{\mathcal{P}}_{t-r} \varphi\right) d \tilde{\mu}=\int_{[0, \tau) \times \mathbb{R}^{d}} \psi \tilde{\mathcal{P}}_{t-r} \varphi d\left(\tilde{\mathcal{P}}_{r}^{*} \tilde{\mu}\right), \quad 0 \leqslant r \leqslant t \tag{5.32}
\end{equation*}
$$

Theorem 5.14 (FDT I) Suppose that Assumption 5.2 holds, Assumption 5.4 holds for $\alpha=0$, and Proposition 5.8 is satisfied. Then, for any $\varphi \in \mathcal{C}_{\tau}^{2}\left(\mathbb{R} \times \mathbb{R}^{d}\right)$ such that $\mathbb{E}\left[D_{x}^{\beta} \varphi\left(\tilde{\Phi}^{\alpha}(t, \cdot)\right)\right] \in L^{1}\left(\tilde{\mu}_{0}\right),|\beta| \leqslant 2$, the following holds:
(i) There exists a family $\left(\tilde{\mu}_{t}\right)_{t \geqslant 0}, \tilde{\mu}_{t} \in \mathcal{P}\left([0, \tau) \times \mathbb{R}^{d}\right)$ of $\tau$-periodic probability measures, and a uniquely $\tilde{\mathcal{P}}_{t}^{*}$-ergodic probability measure, $\overline{\tilde{\mu}} \in \mathcal{P}\left([0, \tau) \times \mathbb{R}^{d}\right)$, which is associated with the $R D S\left\{\tilde{\Phi}(t, \cdot, \cdot): t \in \mathbb{R}^{+}\right\}$on $[0, \tau) \times \mathbb{R}^{d}$ and generated by the $\operatorname{SDE}$ (5.4) with $\alpha=0$.
(ii) The linear response function in (5.30) is given by

$$
\begin{equation*}
\mathcal{R}_{\varphi}^{\tilde{\mu}_{0}}(t-r, r)=\mathcal{K}_{\varphi, \mathbb{B}_{r}}^{\tilde{\mu}_{0}}(t-r, r)=\int_{[0, \tau) \times \mathbb{R}^{d}} \underset{\mathbb{B}_{r}}{ }(\tilde{x}) \tilde{\mathcal{P}}_{t-r} \varphi(\tilde{x}) \tilde{\rho}_{r}(\tilde{x}) \mathrm{d} \tilde{x}, \quad \mathbb{B}_{r}(\tilde{x})=\frac{\tilde{\mathcal{V}}^{*} \tilde{\rho}_{r}(\tilde{x})}{\tilde{\rho}_{r}(\tilde{x})} \tag{5.33}
\end{equation*}
$$

where $0 \leqslant r \leqslant t, \tilde{\mu}_{r}(\mathrm{~d} \tilde{x})=\tilde{\rho}_{r}(\tilde{x}) \mathrm{d} \tilde{x}$ is understood in the sense of $(5.11)$, the correlation function $\mathcal{K}_{\varphi,(\cdot)}^{\tilde{\mu}_{0}}$ is defined in (5.32), and the operator $\tilde{\mathcal{V}}^{*}$ is defined in (5.22).
(iii) The linear response for perturbations of observables based on the ergodic measure $\overline{\tilde{\mu}}$ is

$$
\begin{equation*}
\overline{\mathcal{R}}_{\varphi}(t-r)=\overline{\mathcal{K}}_{\varphi, \mathbb{B}}(t-r)=\int_{[0, \tau) \times \mathbb{R}^{d}} \mathbb{B}(\tilde{x}) \tilde{\mathcal{P}}_{t-r} \varphi(\tilde{x}) \overline{\tilde{\rho}}(\tilde{x}) \mathrm{d} \tilde{x}, \quad \mathbb{B}(\tilde{x})=\frac{\tilde{\mathcal{V}}^{*} \tilde{\tilde{\rho}}(\tilde{x})}{\tilde{\tilde{\rho}}(\tilde{x})} \tag{5.34}
\end{equation*}
$$

where $0 \leqslant r \leqslant t, \overline{\mathcal{R}}_{\varphi}(t-r):=\mathcal{R}_{\varphi}^{\overline{\tilde{\mu}}}(t-r, 0)$, and $\overline{\mathcal{K}}_{\varphi, \mathbb{B}}(t-r):=\mathcal{K}_{\varphi, \mathbb{B}}^{\overline{\tilde{\mu}}}(t-r, 0)$.
The above hold for all time if Proposition 5.8 is satisfied for all time.
Proof Part (i) is a direct consequence of Theorem 4.7 and Theorem 4.11. For Part (ii), we have from the representation of the response functional $\mathcal{R}_{\varphi}^{\tilde{\mu}_{0}}$ in (5.30) together with the operator $\tilde{\mathcal{V}}^{*}$ in (5.22) that the following holds for $0 \leqslant r \leqslant t$

$$
\begin{aligned}
& \mathcal{R}_{\varphi}^{\tilde{\mu}_{0}}(t-r, r) \\
& \quad=\int_{[0, \tau) \times \mathbb{R}^{d}} \tilde{\mathcal{P}}_{r} \tilde{\mathcal{V}}\left(\tilde{\mathcal{P}}_{t-r} \varphi(\tilde{x})\right) \tilde{\mu}_{0}(\mathrm{~d} \tilde{x})=\int_{[0, \tau) \times \mathbb{R}^{d}} \tilde{\mathcal{V}}\left(\tilde{\mathcal{P}}_{t-r} \varphi(\tilde{x})\right) \tilde{\rho}_{r}(\tilde{x}) \mathrm{d} \tilde{x} \\
& \quad=\int_{[0, \tau) \times \mathbb{R}^{d}}\left\{G_{i}(\tilde{x}) \partial_{x_{i}} \tilde{\mathcal{P}}_{t-r} \varphi(\tilde{x})+\frac{1}{2}\left[\sigma_{i k}(\tilde{x}) H_{j k}(\tilde{x})+\sigma_{j k}(\tilde{x}) H_{i k}(\tilde{x})\right] \partial_{x_{i} x_{j}}^{2} \tilde{\mathcal{P}}_{t-r} \varphi(\tilde{x})\right\} \tilde{\rho}_{r}(\tilde{x}) \mathrm{d} \tilde{x} \\
& \quad=\int_{[0, \tau) \times \mathbb{R}^{d}}\left\{-\partial_{x_{i}}\left[G_{i}(\tilde{x}) \tilde{\rho}_{r}(\tilde{x})\right]+\frac{1}{2} \partial_{x_{i} x_{j}}^{2}\left(\left[\sigma_{i k}(\tilde{x}) H_{j k}(\tilde{x})+\sigma_{j k}(\tilde{x}) H_{i k}(\tilde{x}] \tilde{\rho}_{r}(\tilde{x})\right)\right\} \tilde{\mathcal{P}}_{t-r} \varphi(\tilde{x}) \mathrm{d} \tilde{x}\right. \\
& \quad=\int_{[0, \tau) \times \mathbb{R}^{d}} \frac{\tilde{\mathcal{V}}^{*} \tilde{\rho}_{r}(\tilde{x})}{\tilde{\rho}_{r}(\tilde{x})} \tilde{\mathcal{P}}_{t-r} \varphi(\tilde{x}) \tilde{\rho}_{r}(\tilde{x}) \mathrm{d} \tilde{x}=\int_{[0, \tau) \times \mathbb{R}^{d}} \mathbb{B}_{r}(\tilde{x}) \tilde{\mathcal{P}}_{t-r} \varphi(\tilde{x}) \tilde{\rho}_{r}(\tilde{x}) \mathrm{d} \tilde{x}=\mathcal{K}_{\varphi, \mathbb{B}_{r}}^{\tilde{\mu}_{0}}(t-r, r),
\end{aligned}
$$

where $\varphi \in \mathcal{C}_{\tau}^{2}\left(\mathbb{R} \times \mathbb{R}^{d}\right), \mathbb{E}\left[D_{x}^{\beta} \varphi\left(\tilde{\Phi}^{\alpha}(t, \cdot)\right)\right] \in L^{1}\left(\tilde{\mu}_{0}\right),|\beta| \leqslant 2$, as in Theorem 5.11 (see also Proposition 5.8 for a sufficient condition), and $\tilde{\rho}_{r}(\tilde{x}) \mathrm{d} \tilde{x}$ is understood in the sense of (5.11); i.e.

$$
\begin{align*}
\mathcal{R}_{\varphi}^{\tilde{\mu}_{0}}(t-r, r) & =\int_{[0, \tau) \times \mathbb{R}^{d}}\left(\tilde{\mathcal{V}}^{*} \rho_{r}(x)\right)\left(\mathcal{P}_{r, t} \varphi(t, x)\right) d x \\
& =\int_{[0, \tau) \times \mathbb{R}^{d}} \mathbb{B}_{r}(x)\left(\mathcal{P}_{r, t} \varphi(t, x)\right) \rho_{r}(x) d x, \quad 0 \leqslant r \leqslant t \tag{5.35}
\end{align*}
$$

As regards Part (iii), notice that due to the fact that $\tilde{\mathcal{P}}_{r}^{*} \overline{\tilde{\mu}}=\overline{\tilde{\mu}}$ for any $r \in[0, \tau)$, we have

$$
\begin{aligned}
\mathcal{R}_{\varphi}^{\bar{\mu}}(t-r, r) & =\int_{[0, \tau) \times \mathbb{R}^{d}} \tilde{\mathcal{P}}_{r} \tilde{\mathcal{V}}\left(\tilde{\mathcal{P}}_{t-r} \varphi\right) \tilde{\tilde{\mu}}(\mathrm{d} \tilde{x})=\int_{[0, \tau) \times \mathbb{R}^{d}} \tilde{\mathcal{V}}\left(\tilde{\mathcal{P}}_{t-r} \varphi\right) \tilde{\mathcal{P}}_{r}^{*} \overline{\tilde{\mu}}(\mathrm{~d} \tilde{x}) \\
& =\int_{[0, \tau) \times \mathbb{R}^{d}} \tilde{\mathcal{V}}\left(\tilde{\mathcal{P}}_{t-r} \varphi\right) \tilde{\tilde{\mu}}(\mathrm{d} \tilde{x})=\mathcal{R}_{\varphi}^{\bar{\mu}}(t-r, 0)
\end{aligned}
$$

and the desired result can be derived by following analogous derivations to those above.
Remark 5.15
(i) Theorem 5.14 implies that for the $\operatorname{RDS}\left\{\tilde{\Phi}(t, \cdot, \cdot): t \in \mathbb{R}^{+}\right\}$in (4.5) induced by the lifted $\operatorname{SDE}(4.4)$ on $[0, \tau) \times \mathbb{R}^{d}$, the change in the value of an observable $\mathbb{F}_{\varphi}^{\tilde{\mu}_{0}}(t, \alpha)=\left\langle\tilde{\mathcal{P}}_{t} \varphi, \tilde{\mu}_{0}\right\rangle$ in (5.8) in response to a sufficiently small and regular perturbation can be represented by the correlation function utilising the unperturbed dynamics/fluctuations. The operator $\tilde{\mathcal{V}}$ defined in (5.21) does not depend on time due to the $\tau$-periodicity of the coefficients of (5.1) at $\alpha=0$ and the skew-product formulation on $[0, \tau) \times \mathbb{R}^{d}$.
(ii) The response functions (5.33) and (5.34) evaluated on the unperturbed dynamics are amenable to practical approximations via the appropriate long-time averages.
By $\tilde{\mathcal{P}}_{n \tau}^{*}$ - ergodicity of $\tilde{\mu}_{r}$ for any fixed $r \in[0, \tau)$, and $r \leqslant t$ we have (see Remark 4.13)

$$
\begin{aligned}
\mathcal{R}_{\varphi}^{\tilde{\mu}_{0}}(t-r, r) & =\int_{[0, \tau) \times \mathbb{R}^{d}} \mathbb{B}_{r}(\tilde{x}) \tilde{\mathcal{P}}_{t-r} \varphi(\tilde{x}) \tilde{\rho}_{r}(\tilde{x}) \mathrm{d} \tilde{x} \\
& =\lim _{n \rightarrow \infty} \frac{1}{N} \sum_{n=0}^{N} \mathbb{B}_{r}\left(\tilde{\mathcal{P}}_{n \tau}(\tilde{x})\right) \tilde{\mathcal{P}}_{t-r} \varphi\left(\tilde{\mathcal{P}}_{n \tau}(\tilde{x})\right) \\
& =\lim _{n \rightarrow \infty} \frac{1}{N} \sum_{n=0}^{N} \mathbb{B}_{r}\left(\mathcal{P}_{t, t+n \tau}(x)\right)\left(\mathcal{P}_{r, t} \varphi\left(t, \mathcal{P}_{t, t+n \tau}(x)\right)\right)
\end{aligned}
$$

where $\varphi \in \mathcal{C}_{\tau}^{2}\left(\mathbb{R} \times \mathbb{R}^{d}\right), \mathbb{E}\left[D_{x}^{\beta} \varphi\left(\tilde{\Phi}^{\alpha}(t, \cdot)\right)\right] \in L^{1}\left(\tilde{\mu}_{0}\right),|\beta| \leqslant 2$, and in the last line above, we used explicitly the skew-product formulation (5.11) and the representation (5.35); $\left(\mathcal{P}_{r, t}\right)_{t \geqslant r}$ is the family of transition evolutions defined in (3.3). Similarly, by $\tilde{\mathcal{P}}_{t}^{*}$ - ergodicity of $\overline{\tilde{\mu}}$ (see Theorem 4.11) we have for any $u \in \mathbb{R}^{+}$

$$
\begin{aligned}
\overline{\mathcal{R}}_{\varphi}(u) & =\int_{[0, \tau) \times \mathbb{R}^{d}} \mathbb{B}(\tilde{x}) \tilde{\mathcal{P}}_{t-r} \varphi(\tilde{x}) \overline{\tilde{\rho}}(\tilde{x}) \mathrm{d} \tilde{x} \\
& =\lim _{n \rightarrow \infty} \frac{1}{N} \int_{0}^{N} \mathbb{B}\left(\mathcal{P}_{0, \zeta}(x)\right) \mathcal{P}_{0, u} \varphi\left(\zeta, \mathcal{P}_{0, \zeta}(x)\right) \mathrm{d} \zeta
\end{aligned}
$$

(iii) Note that the function $\mathbb{B}_{r}$ in (5.30) of Theorem 5.14 is unique almost everywhere. To see this, suppose that there exists $\tilde{\mathbb{B}}_{r} \in L^{1}\left(\tilde{\mu}_{r}\right)$ such that

$$
\mathcal{K}_{\varphi, \mathbb{B}_{r}}^{\tilde{\mu}}(t-r, r)=\mathcal{K}_{\varphi, \tilde{\mathbb{B}}_{r}}^{\tilde{\mu}}(t-r, r), \quad \varphi \in \mathcal{C}_{c}^{2}\left([0, \tau) \times \mathbb{R}^{d}\right), \quad 0 \leqslant r \leqslant t
$$

This implies that

$$
\left\langle\tilde{\mathcal{P}}_{t-r} \varphi, \mathbb{B}_{r}-\tilde{\mathbb{B}}_{r}\right\rangle_{\tilde{\mu}_{r}}=0, \quad \varphi \in \mathcal{C}_{c}^{2}\left([0, \tau) \times \mathbb{R}^{d}\right), \quad 0 \leqslant r \leqslant t
$$

Given that $\varphi \in \mathcal{C}_{c}^{2}\left([0, \tau) \times \mathbb{R}^{d}\right)$ is bounded, taking limit as $r \rightarrow t$ in the above and applying the dominated convergence theorem, we obtain $\mathbb{B}_{t}=\tilde{\mathbb{B}}_{t} \tilde{\mu}_{t}$-a.e. by arbitrariness of $\varphi$.

Throughout the remainder of this section, we shall assume that $t \mapsto \hat{b}(\alpha, t, x), \hat{\sigma}(\alpha, t, x)$, are $\tau$-periodic for all $(\alpha, x) \in \mathcal{A} \times \mathbb{R}^{d}$. Thus, under Assumptions 5.4 and 5.2 , the family of measures $\left\{\tilde{\mu}_{t}^{\alpha}: t \in \mathbb{R}^{+}, \alpha \in \mathcal{A}\right\}$ is $\tau$-periodic and ergodic (as in the case of Sect. 4.2) and, as a consequence, the time-averaged ( $\tilde{\mathcal{P}}_{t}^{\alpha *}$ - ergodic) measures $\overline{\tilde{\mu}}^{\alpha}$ (see (4.29)) satisfy the stationary PDE

$$
\begin{equation*}
\int_{[0, \tau) \times \mathbb{R}^{d}} \tilde{\mathcal{L}}_{\alpha} \varphi(\tilde{x}) \overline{\tilde{\mu}}^{\alpha}(\mathrm{d} \tilde{x})=0, \quad \alpha \in \mathcal{A}, \quad \varphi \in \mathcal{D}\left(\tilde{\mathcal{L}}_{\alpha}\right) \cap \mathcal{E}_{\alpha} \tag{5.36}
\end{equation*}
$$

where $\tilde{\mathcal{L}}_{\alpha}$ is the generator of the $\operatorname{RDS}\left\{\tilde{\Phi}^{\alpha}(t, \cdot, \cdot): t \in \mathbb{R}^{+}\right\}$on $[0, \tau) \times \mathbb{R}^{d}$, for $\alpha \in \mathcal{A}$ (in the same form as $\tilde{\mathcal{L}}$ in Definition 3.6), and the domain $\mathcal{D}\left(\tilde{\mathcal{L}}_{\alpha}\right)$ is defined by

$$
\begin{align*}
& \mathcal{D}\left(\tilde{\mathcal{L}}_{\alpha}\right):= \mathbb{H}^{1}\left([0, \tau) \times \mathbb{R}^{d} ; \tilde{\tilde{\mu}}^{\alpha}\right) \\
&=\left\{\varphi:[0, \tau) \times \mathbb{R}^{d} \rightarrow \mathbb{R}: \varphi(0, \cdot)=\varphi(\tau, \cdot)\right. \text { and } \\
& \int_{[0, \tau) \times \mathbb{R}^{d}}|\varphi(s, x)|^{2} \overline{\tilde{\mu}}^{\alpha}(\mathrm{d} s \mathrm{~d} x) \\
&\left.\quad+\int_{[0, \tau) \times \mathbb{R}^{d}}|D \varphi(s, x)|^{2} \overline{\tilde{\mu}}^{\alpha}(\mathrm{d} s \mathrm{~d} x)<\infty, \alpha \in \mathcal{A}\right\} .  \tag{5.37}\\
& \mathcal{E}_{\alpha}:=\left\{\varphi \in \mathcal{C}^{2}\left([0, \tau) \times \mathbb{R}^{d}\right): \varphi(0, \cdot)=\varphi(\tau, \cdot),\right. \text { and } \\
&\left.\mathbb{E}\left[D_{x}^{\beta} \varphi\left(\tilde{\Phi}^{\alpha}(t, \cdot)\right)\right] \in L^{1}\left(\overline{\tilde{\mu}}^{\alpha}\right),|\beta| \leqslant 2, \alpha \in \mathcal{A}\right\} . \tag{5.38}
\end{align*}
$$

Sufficient conditions for $\varphi \in \mathcal{E}_{\alpha}$ were given in Proposition 5.8.
Given the ergodic measure $\overline{\tilde{\mu}}$ (4.29) of the unperturbed dynamics (i.e. (5.4) with $\alpha=0$ ) and the above set-up, the linear response $\Delta \mathbb{F}_{\varphi, \vartheta}^{\tilde{\tilde{\mu}}}(t)$ of the statistical observable $\mathbb{F}_{\varphi}^{\tilde{\tilde{\mu}}}(t, \alpha)$ in (5.8) due to a sufficiently small perturbation $\varepsilon \vartheta(t)$ around $\alpha=0$ such that it preserves the $\tau$-periodicity of the unperturbed dynamics is summarised as follows.

Theorem 5.16 (FDT II) Let Assumptions 5.2-5.4 be satisfied for $\alpha \in \mathcal{A}$, and suppose that Proposition 5.8 holds. Assume further that $t \mapsto \hat{b}(t, \alpha, x), t \mapsto \hat{\sigma}(t, \alpha, x)$ are $\tau$-periodic for all $(\alpha, x) \in \mathcal{A} \times \mathbb{R}^{d}$. Then, the following hold:
(i) For every $\alpha \in \mathcal{A}$, $\left\{\tilde{\mu}_{t}^{\alpha}: t \in \mathbb{R}^{+}\right\}$is a family of $\tau$-periodic measures induced by the one-point motion $\tilde{x} \mapsto \tilde{\Phi}^{\alpha}(t, \omega, \tilde{x})$, $\omega \in \Omega$, with the uniquely $\tilde{\mathcal{P}}_{t}^{\alpha *}$-ergodic measure $\tilde{\tilde{\mu}}^{\alpha}$ satisfying (5.36).
(ii) The map $\alpha \mapsto \overline{\tilde{\mu}}^{\alpha}(\mathrm{d} \tilde{x})=\overline{\tilde{\rho}}^{\alpha}(\tilde{x}) \mathrm{d} \tilde{x}$ (with $\overline{\tilde{\rho}}^{\alpha}$ understood in the sense of $(5.11)$ ) is weakly differentiable at $\alpha=0$ for all $\tilde{x} \in[0, \tau) \times \mathbb{R}^{d}$, and the linear response (5.30) associated with perturbations of observables based on the $\tilde{\mathcal{P}}_{t}^{*}$-ergodic measure $\tilde{\tilde{\mu}}$ is given by

$$
\begin{equation*}
\overline{\mathcal{R}}_{\varphi}(t-r):=\mathcal{R}_{\varphi}^{\tilde{\tilde{\mu}}}(t-r, 0)=\partial_{r} \mathcal{K}_{\varphi, \mathbb{W}}^{\tilde{\tilde{\mu}}}(t-r, r), \quad 0 \leqslant r \leqslant t \tag{5.39}
\end{equation*}
$$

for $\varphi \in \mathcal{D}\left(\tilde{\mathcal{L}}_{\alpha}\right) \cap \mathcal{E}_{\alpha}$, with the correlation function $\mathcal{K}_{\varphi, \mathbb{W}}^{\overline{\tilde{\mu}}}$ in $(5.32) \mathbb{W} \in \mathcal{C}_{\infty}\left([0, \tau) \times \mathbb{R}^{d}\right)$ given by

$$
\begin{equation*}
\mathbb{W}(\tilde{x})=\frac{\eta(\tilde{x})}{\tilde{\tilde{\rho}}(\tilde{x})}, \quad \text { s.t. } \quad\langle\eta, \varphi\rangle:=\left.\left\langle\partial_{\alpha} \overline{\tilde{\rho}}^{\alpha}, \varphi\right\rangle\right|_{\alpha=0^{\prime}} \tag{5.40}
\end{equation*}
$$

with $\partial_{\alpha} \overline{\tilde{\rho}}^{\alpha}$ understood in the weak sense.
Proof Part (i) is a direct consequence of Theorem 4.7 and Theorem 4.11 given the fact that Assumptions 5.2-5.4 hold for $\alpha$ in a proper interval $\mathcal{A}$ containing $\alpha=0$. For Part (ii), we proceed as at the beginning of the proof of Theorem 5.11, except that due to Part (i), both the unperturbed and the perturbed measures are $\tau$-periodic. Thus, for $\vartheta \in \mathcal{C}_{\infty}^{1}\left(\mathbb{R}^{+} ; \mathbb{R}\right)$ and $\varepsilon>0$ sufficiently small so that $\varepsilon \vartheta \in \mathcal{A}$, and for $\varphi \in \mathcal{E}_{\alpha}$ there exists a constant $0<C_{\varepsilon, \varphi}<\infty$ such that

$$
\begin{align*}
\left\langle\varphi, \tilde{\mu}_{t}^{\varepsilon \vartheta}\right\rangle-\left\langle\varphi, \tilde{\mu}_{t}\right\rangle & =\int_{[0, \tau) \times \mathbb{R}^{d}}\left(\tilde{\mathcal{P}}_{t}^{\varepsilon \vartheta} \varphi(\tilde{x})-\tilde{\mathcal{P}}_{t} \varphi(\tilde{x})\right) \tilde{\mu}_{0}(\mathrm{~d} \tilde{x}) \\
& =\int_{[0, \tau) \times \mathbb{R}^{d}}\left(\int_{0}^{t} \varepsilon \vartheta(r) \mathbb{E}\left[\hat{f}_{t, \varepsilon}(r, \tilde{\Phi}(r, \cdot, \tilde{x}))\right] \mathrm{d} r\right) \tilde{\mu}_{0}(\mathrm{~d} \tilde{x}) \\
& \leqslant \varepsilon t\|\vartheta\|_{\infty} C_{\varepsilon, \varphi}, \tag{5.41}
\end{align*}
$$

where the bound is due to Proposition 5.8. Averaging both sides over $t \in[0, \tau)$, we have

$$
\begin{equation*}
\lim _{\varepsilon \downarrow 0} \frac{1}{\varepsilon}\left\langle\varphi, \overline{\tilde{\mu}}^{\varepsilon \vartheta}-\overline{\tilde{\mu}}\right\rangle<\infty, \quad \varphi \in \mathcal{E}_{\alpha} . \tag{5.42}
\end{equation*}
$$

Thus, $\overline{\tilde{\mu}}^{\alpha}$ is weakly differentiable. Next, by the Hörmander condition in Assumption 5.2, we have $\overline{\tilde{\mu}}^{\alpha}(\mathrm{d} \tilde{x})=\overline{\tilde{\rho}}^{\alpha}(\tilde{x}) \mathrm{d} \tilde{x}$ so that for any $\vartheta \in \mathcal{C}_{\infty}^{1}\left(\mathbb{R}^{+}, \mathbb{R}\right)$

$$
\begin{equation*}
\lim _{\varepsilon \downarrow 0} \frac{1}{\varepsilon}\left\langle\varphi, \overline{\tilde{\rho}}^{\varepsilon \vartheta}-\tilde{\tilde{\rho}}\right\rangle<\infty, \quad \varphi \in \mathcal{E}_{\alpha} \tag{5.43}
\end{equation*}
$$

which implies that $\overline{\tilde{\rho}}^{\alpha}$ is weakly differentiable at $\alpha=0$ (with $\overline{\tilde{\rho}}^{\alpha}$ understood in the sense of (5.11) to simplify notation). Furthermore, (5.36) yields

$$
\begin{equation*}
\left\langle\overline{\tilde{}}^{\alpha}, \tilde{\mathcal{L}}_{\alpha} \varphi\right\rangle=0, \quad \alpha \in \mathcal{A}, \quad \varphi \in \mathcal{D}\left(\tilde{\mathcal{L}}_{\alpha}\right) \cap \mathcal{E}_{\alpha} \tag{5.44}
\end{equation*}
$$

Differentiating (5.44) with respect to the parameter $\alpha$ (in the weak sense), we obtain

$$
\begin{equation*}
\left.\left\langle\partial_{\alpha} \overline{\tilde{\rho}}^{\alpha},-\tilde{\mathcal{L}} \varphi\right\rangle\right|_{\alpha=0}=\langle\overline{\tilde{\rho}}, \tilde{\mathcal{V}} \varphi\rangle=\left\langle\tilde{\mathcal{V}}^{*} \overline{\tilde{\rho}}, \varphi\right\rangle . \tag{5.45}
\end{equation*}
$$

Next, we set $\langle\eta,-\tilde{\mathcal{L}} \varphi\rangle:=\left.\left\langle\partial_{\alpha} \overline{\tilde{\rho}}^{\alpha},-\tilde{\mathcal{L}} \varphi\right\rangle\right|_{\alpha=0}$, and note that $\tilde{\mathcal{P}}_{t} \varphi \in \mathcal{D}\left(\tilde{\mathcal{L}}_{\alpha}\right) \cap \mathcal{E}_{\alpha}, t \geqslant 0$, for any $\varphi \in D\left(\tilde{\mathcal{L}}_{\alpha}\right) \cap \mathcal{E}_{\alpha}$ (due to Assumption 5.2 and the associated smoothing property of $\left(\mathcal{P}_{s, t}\right)_{s \leqslant t}$ generating $\left(\tilde{\mathcal{P}}_{t}\right)_{t \in \mathbb{R}^{+}}$in (4.7); e.g. Proposition 4.10 and [26]). Thus, we have

$$
\begin{equation*}
\mathcal{K}_{\varphi, \mathbb{W}}^{\tilde{\tilde{\mu}}}(t-r, r)=\left\langle\tilde{\mathcal{P}}_{r}\left(\mathbb{W} \tilde{\mathcal{P}}_{t-r} \varphi\right), \overline{\tilde{\rho}}\right\rangle=\left\langle\tilde{\mathcal{P}}_{t-r} \varphi, \mathbb{W} \overline{\tilde{\rho}}\right\rangle=\left\langle\tilde{\mathcal{P}}_{t-r} \varphi, \eta\right\rangle \tag{5.46}
\end{equation*}
$$

where $\mathbb{W}$ is given in (5.40), and subsequently

$$
\partial_{r} \mathcal{K}_{\varphi, \mathbb{W}}(t-r, r)=\partial_{r}\left\langle\tilde{\mathcal{P}}_{t-r} \varphi, \eta\right\rangle=\left\langle-\tilde{\mathcal{L}} \tilde{\mathcal{P}}_{t-r} \varphi, \eta\right\rangle .
$$

Since $\tilde{\mathcal{P}}_{t-r} \varphi \in \mathcal{D}\left(\tilde{\mathcal{L}}_{\alpha}\right) \cap \mathcal{E}_{\alpha}$ for $0 \leqslant r \leqslant t$, we have by (5.45) that for $\varphi \in \mathcal{D}\left(\tilde{\mathcal{L}}_{\alpha}\right) \cap \mathcal{E}_{\alpha}$

$$
\begin{aligned}
\partial_{r} \mathcal{K}_{\varphi, \mathbb{W}}(t-r, r)=\left\langle-\tilde{\mathcal{L}} \tilde{\mathcal{P}}_{t-r} \varphi, \eta\right\rangle & =\left\langle\tilde{\mathcal{P}} \tilde{\mathcal{P}}_{t-r} \varphi, \tilde{\tilde{\rho}}\right\rangle \\
& =\int_{[0, \tau) \times \mathbb{R}^{d}}\left(\tilde{\mathcal{V}} \tilde{\mathcal{P}}_{t-r} \varphi\right)(\tilde{x}) \tilde{\tilde{\mu}}(\mathrm{d} \tilde{x})=\mathcal{R}_{\varphi}^{\tilde{\mu}}(t-r, 0) .
\end{aligned}
$$

Remark 5.17 Note that Theorem 5.14 is more general than Theorem 5.16 in the sense that it only requires time periodicity of the coefficients of the $\operatorname{SDE}$ (5.1) and the existence of time-periodic probability measure for the unperturbed dynamics (i.e. for $\alpha=0$ in (5.1)) but it does not preclude the perturbed dynamics to have time-periodic measures. Thus, a natural question arises as to the connection between the linear response functions $\overline{\mathcal{R}}_{\varphi}$ in (5.34) of Theorem 5.14 and (5.39) of Theorem 5.16, respectively, in the case when both the unperturbed and the perturbed dynamics (i.e. for $\alpha \in \mathcal{A}$ in (5.1)) have time-periodic ergodic measures of period $\tau$. The desired connection stems from the fact that under Assumption 5.2 the identity (5.45) leads to

$$
\langle\mathbb{W},-\tilde{\mathcal{L}} \varphi\rangle_{\overline{\tilde{\mu}}}=\langle\mathbb{W} \overline{\tilde{\rho}},-\tilde{\mathcal{L}} \varphi\rangle=\langle\eta,-\tilde{\mathcal{L}} \varphi\rangle=\left\langle\tilde{\mathcal{V}}^{*} \overline{\tilde{\rho}}, \varphi\right\rangle=\left\langle\frac{\tilde{\mathcal{V}}^{*} \tilde{\tilde{\rho}}}{\tilde{\tilde{\rho}}}, \overline{\tilde{\rho}} \varphi\right\rangle=\langle\mathbb{B}, \varphi\rangle_{\overline{\tilde{\mu}}}
$$

for any $\varphi \in \mathcal{D}\left(\tilde{\mathcal{L}}_{\alpha}\right) \cap \mathcal{E}_{\alpha}$, where $\tilde{\mathcal{L}}=\partial_{s}+\mathcal{L}$ is the generator of the one-point motion $\tilde{x} \mapsto \tilde{\Phi}(t, \omega, \tilde{x}), t \geqslant 0$, on the extended state space $[0, \tau) \times \mathbb{R}^{d}$. Thus, $\mathbb{B}=-\tilde{\mathcal{L}}^{\tilde{\tilde{\mu}} * \mathbb{W}}$, where $\tilde{\mathcal{L}}^{\tilde{\tilde{\mu}} *}$ is the $L^{2}(\tilde{\tilde{\mu}})$ dual of $\tilde{\mathcal{L}}$. In fact, the above result also implies that $\mathbb{W}$ in (5.40) of Theorem 5.16 is not unique in contrast to $\mathbb{B}$ in (5.34) of Theorem 5.14. To see this, assume that there exists $\tilde{\mathbb{W}} \in L^{1}(\tilde{\tilde{\mu}})$ on $[0, \tau) \times \mathbb{R}^{d}$ such that

$$
\partial_{r} \mathcal{K}_{\varphi, \mathbb{W}}^{\overline{\tilde{\mu}}}(t-r, r)=\partial_{r} \mathcal{K}_{\varphi, \tilde{\mathbb{W}}}^{\overline{\tilde{\mu}}}(t-r, r), \quad \varphi \in \mathcal{C}_{c}^{\infty}\left([0, \tau) \times \mathbb{R}^{d}\right), \quad 0 \leqslant r \leqslant t
$$

which implies that (see (5.46))

$$
\frac{d}{\mathrm{~d} r}\left\langle\tilde{\mathcal{P}}_{t-r} \varphi, \mathbb{W}-\left.\tilde{\mathbb{W}}\right|_{\tilde{\tilde{\mu}}}=0, \quad \varphi \in \mathcal{C}_{c}^{\infty}\left([0, \tau) \times \mathbb{R}^{d}\right), \quad 0 \leqslant r \leqslant t\right.
$$

Since $\mathcal{C}_{c}^{\infty}\left([0, \tau) \times \mathbb{R}^{d}\right) \subset \mathcal{D}(\tilde{\mathcal{L}})$, it follows by the smoothing property of $\tilde{\mathcal{P}}_{t-r}$ (under Assumption 5.2) and the dominated convergence theorem that

$$
\langle\tilde{\mathcal{L}} \varphi, \mathbb{W}-\tilde{\mathbb{W}}\rangle_{\tilde{\tilde{\mu}}}=0
$$

Since $\mathbb{W}-\tilde{\mathbb{W}} \in L^{1}(\overline{\tilde{\mu}}), \mathbb{W}-\tilde{\mathbb{W}}$ is a.e. constant; hence, $\mathbb{W}$ satisfying Theorem 5.16 is not unique.


Fig. 2 The responses $\Delta \mathbb{E}_{x}(t), \Delta \mathbb{E}_{y}(t), \Delta \mathbb{E}_{z}(t)$ of the expectations $\mathbb{E}\left[x_{t}\right], \mathbb{E}\left[y_{t}\right], \mathbb{E}\left[z_{t}\right]$, associated with the dynamics (4.40) perturbed to (5.47) via the spatially uniform, time-aperiodic perturbation (5.51) with $\varepsilon=0.05$, and $t_{0}=80 \tau, \Delta T=12$. The left column shows the direct simulation of ( 5.47 ), and the right column shows the comparison between the exact response and the linear response given by (5.48) and the response function given by (5.50). The parameters of the unperturbed system (4.40) are $\bar{\alpha}=7.3, \bar{\beta}=26, \bar{\gamma}=7, \bar{\varrho}=10$,
$\bar{f}=100, \bar{\delta}=0.9, \tau=1, \bar{\sigma}=0.2$; see Fig. 1. Further details are discussed in the main text of Example 5.18

Example 5.18 (Stochastic Lorenz model with periodic forcing) We return to the stochastic Lorenz model with time-periodic forcing used in Example 4.14, and we consider a simple case of perturbed dynamics in (4.40) in the form

$$
\begin{equation*}
\mathrm{d} \nu_{t}^{\alpha}=\left[-A v_{t}^{\alpha}-G\left(v_{t}^{\alpha}\right)+F(t)+\varepsilon \mathfrak{F}\left(v_{t}^{\alpha}\right) \vartheta(t)\right] \mathrm{d} t+\sigma\left(v_{t}^{\alpha}\right) \mathrm{d} W_{t} \tag{5.47}
\end{equation*}
$$

where $\vartheta \in \mathcal{C}_{\infty}^{1}\left(\mathbb{R}^{+}, \mathbb{R}\right), \vartheta(0)=0, v \mapsto \mathfrak{F}(v) \in \mathcal{C}_{\infty}^{\infty}\left(\mathbb{R}^{d}\right)$, and $A, G(v), F(t)$ are defined as in (4.40) for $v_{t}^{\alpha}=\left(x_{t}^{\alpha}, y_{t}^{\alpha}, z_{t}^{\alpha}\right)$. Here, the last term in the drift represents a perturbation of the dynamics (4.40) due to $\alpha(t)=\varepsilon \vartheta(t)$, which is chosen for simplicity to be in the space-time factorised form (recall, however, that the perturbation in this framework can take a more general form; see Remark 5.12).
Note that, analogously to (4.40), the periodically forced dynamics in (5.47) satisfies Assumption 5.2, and recall that in appropriate parameter regimes of (4.40) (or in (5.47) with $\alpha=0$ ) there exists a time-periodic ergodic measure $\mu_{t} \in \mathcal{P}\left(\mathbb{R}^{d}\right)$ with a smooth density $\rho_{t}$ with respect to the Lebesgue measure on $\mathbb{R}^{3}$ for all $t \in[0, \tau), 0<\tau<\infty$.
Assuming that the conditions of Theorem 5.11 hold, the linear response $\Delta \mathbb{F}_{\varphi, \vartheta}^{\tilde{\mu}_{0}}(t)$ in (5.29)

$$
\begin{equation*}
\Delta \mathbb{F}_{\varphi, \vartheta}^{\tilde{\mu}_{0}}(t)=\int_{0}^{t} \mathcal{R}_{\varphi}^{\tilde{\mu}_{0}}(t-r, r) \vartheta(r) \mathrm{d} r, \quad \vartheta \in \mathcal{C}_{\infty}^{1}\left(\mathbb{R}^{+}, \mathbb{R}\right), \quad \vartheta(0)=0 \tag{5.48}
\end{equation*}
$$

of the observable $\varphi$ to the perturbation $\varepsilon \mathfrak{F}(v) \vartheta(t)$ in (5.47) in the random time-periodic ergodic regime is determined by convolving the response function $\mathcal{R}_{\varphi}^{\tilde{\mu}_{0}}$ with $\vartheta$, where (see Theorem 5.14)

$$
\begin{equation*}
\mathcal{R}_{\varphi}^{\tilde{\mu}_{0}}(t-r, r)=\int_{[0, \tau) \times \mathbb{R}^{3}} \partial_{v}\left(\mathfrak{F}(v) \tilde{\rho}_{r}(\tilde{v})\right) \tilde{\mathcal{P}}_{t-r} \varphi(\tilde{v}) \mathrm{d} \tilde{v}=\mathcal{K}_{\varphi, \mathbb{B}_{r}}^{\tilde{\mu}_{0}}(t-r, r), \quad 0 \leqslant r \leqslant t, \tag{5.49}
\end{equation*}
$$

which is solely based on the unperturbed dynamics. In the above expression $\tilde{\mu}_{t}(\tilde{v})=\tilde{\rho}_{t}(\tilde{v}) \mathrm{d} \tilde{v}, \tilde{v}=(s, v) \in[0, \tau) \times \mathbb{R}^{3}$, and $\mathcal{K}_{\varphi, \mathbb{B}_{r}}^{\tilde{\mu}_{0}}(t-r, r)$ is the correlation function (5.32) of the random variables $\varphi\left(\tilde{v}_{t}\right)$, and $\mathbb{B}_{r}(\tilde{v})=-\partial_{v}\left(\mathfrak{F}(v) \tilde{\rho}_{r}(\tilde{v})\right) / \tilde{\rho}_{r}(\tilde{v})$, with $\tilde{\mu}_{t}(\tilde{v})=\tilde{\rho}_{t}(\tilde{v}) \mathrm{d} \tilde{v}$ understood in the sense of (5.11). Thus, the response function can be written in a form amenable to computations as (see Remark 5.15(ii))

$$
\begin{align*}
\mathcal{R}_{\varphi}^{\tilde{\mu}_{0}}(t-r, r) & =\int_{[0, \tau) \times \mathbb{R}^{d}}\left(\tilde{\mathcal{V}}^{*} \rho_{r}(v)\right)\left(\mathcal{P}_{r, t} \varphi(t, v)\right) d v \\
& =\int_{[0, \tau) \times \mathbb{R}^{d}} \mathbb{B}_{r}(v)\left(\mathcal{P}_{r, t} \varphi(t, v)\right) \rho_{r}(v) d v, \quad 0 \leqslant r \leqslant t \tag{5.50}
\end{align*}
$$

where $\tilde{\mathcal{V}}^{*}$ is defined in (5.22), $\mathbb{B}_{r}$ is given in (5.33), and $\left(\mathcal{P}_{r, t}\right)_{t \geqslant r}$ is the family of transition evolutions defined in (3.3). Furthermore, (5.50) can be evaluated in a more practical fashion via appropriate ergodic averages, as discussed in Remark 5.15(ii).

For simplicity of the numerical illustration, we consider the linear response of the expectation of the solutions to (4.40) to the perturbation $\varepsilon \mathfrak{F}\left(v^{\alpha}\right) \vartheta(t)$ introduced in the drift coefficient of (5.47). Given the dynamics (5.47) and $\varphi(v)=v$, setting $p=2$ is sufficient for Assumption 4.4 and Proposition 5.8 to hold (i.e. Theorems 4.7, 4.11, 5.11, 5.14 will hold for all time if the perturbation maintains dissipativity, which is the case here). In the examples shown in Figs. 2 and 3, we denote the expectation of the solutions to (5.47) by $\mathbb{E}\left[x_{t}\right], \mathbb{E}\left[y_{t}\right], \mathbb{E}\left[z_{t}\right]$, and we consider the response of the expectation to a spatially uniform perturbation $\varepsilon \mathfrak{F}\left(v^{\alpha}\right) \vartheta(t)$ with $\mathfrak{F}\left(v^{\alpha}\right)=(\bar{f}, 0,0)^{T}$ and

$$
\begin{equation*}
\vartheta(t)=\Theta(t) \cos (2 \pi / 3.3 t) \tag{5.51}
\end{equation*}
$$

where

$$
\Theta(t)=\left\{\begin{array}{cl}
1 & t>t_{0}+\Delta T  \tag{5.52}\\
\frac{2}{\Delta T}\left(t-t_{0}\right)-\frac{1}{\Delta T^{2}}\left(t-t_{0}\right)^{2} & t_{0} \leqslant t \leqslant t_{0}+\Delta T \\
\frac{2}{\Delta T}\left(t-t_{0}\right)+\frac{1}{\Delta T^{2}}\left(t-t_{0}\right)^{2} & t_{0}-\Delta T \leqslant t \leqslant t_{0} \\
-1 & t<-t_{0}-\Delta T
\end{array}\right.
$$

The simulations were performed for (5.47) with the same parameter values as those in Example 4.14 in the random time-periodic regime, and the perturbation with $t_{0}=80 \tau, \Delta T=1$ with the amplitudes set to $\varepsilon=0.05$ in Fig. 2 , and $\varepsilon=0.25$ in Fig. 3. The unperturbed initial time-periodic measure at $t=0$, i.e. $\tilde{\mu}_{0}=\delta_{0} \otimes \mu_{0}$, was approximated from long-time simulations of an ensemble of solutions to (5.47) with $\alpha=0$ at $t=\tau n, n \in \mathbb{N}_{0}$. To simplify the notation, the linear response $\Delta \mathbb{F}_{\varphi, \vartheta}^{\tilde{\mu}_{0}}(t)$ in (5.29) of the expectations is denoted by, respectively, $\Delta \mathbb{E}_{x}(t), \Delta \mathbb{E}_{y}(t), \Delta \mathbb{E}_{z}(t)$. The linear response was estimated with the help of the fluctuation-dissipation formula (5.50), where $\mathcal{K}_{\varphi, \mathbb{B}_{r}}(t-r, r)$ in (5.32) exploits the statistical correlations in the time-asymptotic dynamics of the unperturbed system (4.40) via (5.31). As expected from the theory, the linear response provides a good approximation for a sufficiently small perturbation (Fig. 2), and it deteriorates with the increasing amplitude of the perturbation (Fig. 3); the accuracy of the approximation improves still with the decreasing amplitude of the perturbation but we do not show these unsurprising results.


Fig. 3 The responses $\Delta \mathbb{E}_{x}(t), \Delta \mathbb{E}_{y}(t), \Delta \mathbb{E}_{z}(t)$ of $\mathbb{E}\left[x_{t}\right], \mathbb{E}\left[y_{t}\right], \mathbb{E}\left[z_{t}\right]$ associated with the dynamics (4.40) to a spatially uniform, time-aperiodic perturbation (5.51) in (5.47) with $\varepsilon=0.25$; the remaining parameters are as in Fig. 2. See the main text in Example 5.18 for more details

In the simple example illustrated in Fig. 4, we consider the response of the expectation of the solution to (5.47) in the stable random time-periodic regime to a spatially uniform time-periodic perturbation $\varepsilon \mathfrak{F}\left(v^{\alpha}\right) \vartheta(t)$ with $\varepsilon=0.1, \mathfrak{F}\left(v^{\alpha}\right)=(\bar{f}, 0,0)^{T}$ and

$$
\begin{equation*}
\vartheta(t)=H(t-80.25 \tau) \cos ^{2}(2 \pi / \tau t), \tag{5.53}
\end{equation*}
$$

where $H(t)$ is the Heaviside step function. As in the previous examples, the accuracy of the linear response via the FDT formulas (5.49) or (5.50) combined with (5.48) improves for decreasing magnitude of the perturbation but we do not show these unsurprising results. Note that in this case the perturbed measures are also $\tau$-periodic and one could consider the FDT for the ergodic measure $\overline{\tilde{\mu}}$ as in Theorem 5.16; such considerations of the linear response are more relevant in the abstract analysis and we do not pursue such a scenario here (see, however, Remark 5.17).
Finally, it needs to be stressed that the direct numerical evaluation of the correlation function $\mathcal{K}_{\varphi, \mathbb{B}_{r}}(t-r, r)$ in (5.49) is, in general, very computationally intensive due to the need for estimating the time-dependent density $\rho_{t}(v), v \in \mathbb{R}^{3}, t \in[0, \tau)$; more practical implementations rely on various approximations (e.g. a Gaussian approximation of the underlying density), and they were discussed in [62] in the time-periodic setting, and in [ $1-3,34,36-39,59,62,64]$ in the stationary setting. These references also consider much more elaborate examples than what we could consider in this work.

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Fig. 4 The responses $\Delta \mathbb{E}_{x}(t), \Delta \mathbb{E}_{y}(t), \Delta \mathbb{E}_{z}(t)$ of $\mathbb{E}\left[x_{t}\right], \mathbb{E}\left[y_{t}\right]$, and $\mathbb{E}\left[z_{t}\right]$ associated with the dynamics (4.40) to a spatially uniform, time-periodic perturbation (5.53) in (5.47) with $\varepsilon=0.1$; the remaining parameters are as in Fig. 2. See the main text in Example 5.18 for more details

## Appendix A: Growth conditions and existence of absolute moments of solutions of non-autonomous SDE's

Here, we provide explicit examples of two classes of time-periodic coefficients in the SDE (2.1) which satisfy Assumption 4.4.

Lemma A. $1 \operatorname{Let}\{\phi(t, s, \cdot, \cdot): t \geqslant s\}$ be a stochastic flow generated by the $\operatorname{SDE}(2.1)$, and let $V \in \mathcal{C}^{1,2}\left(\mathbb{R} \times \mathbb{R}^{d} ; \mathbb{R}^{+}\right)$be a Lyapunov function satisfying the first part of (4.10) with some $1<p<\infty$. Assume that the following 'dissipative' growth conditions hold on the coefficients of (2.1)

$$
\begin{equation*}
\langle b(t, x), x\rangle \leqslant L_{b_{1}}(t)-L_{b_{2}}(t)|x|^{2}, \quad\|\sigma(t, x)\|_{\mathrm{HS}}^{2} \leqslant L_{\sigma}(t)\left(1+|x|^{2}\right) \tag{A.1}
\end{equation*}
$$

Suppose further that there exist bounded functions $L_{b_{1}}(\cdot), L_{b_{2}}(\cdot), L_{\sigma}(\cdot) \in \mathcal{C}_{\infty}\left(\mathbb{R} ; \mathbb{R}^{+}\right)$such that

$$
\inf _{t \in \mathbb{R}}\left(L_{b_{2}}(t)-2^{\frac{p}{2}-1} L_{b_{1}}-\frac{1}{2}\left(2^{\frac{p}{2}-1}+1\right) L_{\sigma}(t)(p-1)\right)>0
$$

Then, there exists a stochastic flow $\{\phi(t, s, \cdot, \cdot): s \leqslant t\}$ on $\mathbb{R}^{d}$ induced by the solutions of (2.1), which has a finite p-th absolute moment for all time. Furthermore, the following holds

$$
\limsup _{(t-s) \rightarrow \infty} \mathbb{E}[V(t, \phi(t, s, x)-x)]<\infty
$$

in Assumption 4.4(iii).
Proof It can be established (in a similar way to that in [54, Theorem 3.4.6]) that the growth conditions (A.1) lead to existence and uniqueness of global solutions of

$$
\begin{equation*}
\mathrm{d} X_{t}^{s, x}=b\left(t, X_{t}^{s, x}\right) \mathrm{d} t+\sigma\left(t, X_{t}^{s, x}\right) \mathrm{d} W_{t-s}, \quad X_{s}^{s, x}=x, \tag{A.2}
\end{equation*}
$$

such that $X_{t}^{s, x}(\omega)=\phi(t, s, \omega, x) \mathbb{P}$-a.s., $t \geqslant s$. In order to prove the main part of the Lemma, consider first $g(x)=|x|^{p}, p \geqslant 2$; then,

$$
\begin{align*}
\mathcal{L}_{t}|x|^{p} & =p|x|^{p-2} \sum_{i=1}^{d} b_{i}(t, x) x_{i}+\frac{1}{2} p|x|^{p-4} \sum_{i, j=1}^{d}\left\{|x|^{2} \delta_{i j}+(p-2) x_{i} x_{j}\right\}\left(\sigma \sigma^{T}\right)_{i j}(t, x, x) \\
& \leqslant p|x|^{p-2}\left(L_{b_{1}}(t)-L_{b_{2}}(t)|x|^{2}\right)+\frac{1}{2} L_{\sigma}(t) p(p-1)\left(1+|x|^{2}\right)|x|^{p-2} \\
& =p\left(L_{b_{1}}(t)+\frac{1}{2} L_{\sigma}(t)(p-1)\right)|x|^{p-2}-p\left(L_{b_{2}}(t)-\frac{1}{2} L_{\sigma}(t)(p-1)\right)|x|^{p} \tag{A.3}
\end{align*}
$$

Next, since $|x|^{p-2} \leqslant\left(1+|x|^{2}\right)^{\frac{p}{2}-1} \leqslant\left(1+|x|^{2}\right)^{\frac{p}{2}} \leqslant 2^{\frac{p}{2}-1}\left(1+|x|^{p}\right)$, we get

$$
\begin{aligned}
\mathcal{L}_{t}|x|^{p} \leqslant & 2^{\frac{p}{2}-1} p\left(L_{b_{1}}(t)+\frac{1}{2} L_{\sigma}(t)(p-1)\right) \\
& -p\left(L_{b_{2}}(t)-2^{\frac{p}{2}-1} L_{b_{1}}-\frac{1}{2}\left(2^{\frac{p}{2}-1}+1\right) L_{\sigma}(t)(p-1)\right)|x|^{p},
\end{aligned}
$$

which can be written as

$$
\begin{equation*}
\mathcal{L}_{t}|x|^{p} \leqslant \mathfrak{a}_{p}-\mathfrak{b}_{p}|x|^{p} \tag{A.4}
\end{equation*}
$$

with coefficients

$$
\begin{align*}
& \mathfrak{a}_{p}=p 2^{\frac{p}{2}-1} \sup _{t \in \mathbb{R}}\left(L_{b_{1}}(t)+\frac{1}{2} L_{\sigma}(t)(p-1)\right)  \tag{A.5}\\
& \mathfrak{b}_{p}=p \inf _{t \in \mathbb{R}}\left(L_{b_{2}}(t)-2^{\frac{p}{2}-1} L_{b_{1}}-\frac{1}{2}\left(2^{\frac{p}{2}-1}+1\right) L_{\sigma}(t)(p-1)\right) . \tag{A.6}
\end{align*}
$$

It turns out that sharper bounds can be obtained for $p=2,3$; these are derived in Proposition A.2.
Next, consider $X_{t}^{s, x}(\omega)=\phi(t, s, \omega, x)$ solving (A.2) and $g(x)=|x|^{p}, p \geqslant 2$. Then, Itô's Lemma and the bound (A.4) lead to

$$
\begin{equation*}
d \mathbb{E}\left[\left|X_{t}^{s, x}\right|^{p}\right]=\mathbb{E}\left[\mathcal{L}_{t}\left|X_{t}^{s, x}\right|^{p}\right] \mathrm{d} t \leqslant\left(\mathfrak{a}_{p}-\mathfrak{b}_{p} \mathbb{E}\left[\left|X_{t}^{s, x}\right|^{p}\right]\right) \mathrm{d} t . \tag{A.7}
\end{equation*}
$$

Therefore, based on the differential form of Gronwall's inequality, (A.7) yields

$$
\begin{equation*}
\mathbb{E}\left|X_{t}^{s, x}\right|^{p} \leqslant e^{-\mathfrak{b}_{p}(t-s)} \mathbb{E}|x|^{p}+\frac{\mathfrak{a}_{p}}{\mathfrak{b}_{p}}\left(1-e^{-\mathfrak{b}_{p}(t-s)}\right) \tag{A.8}
\end{equation*}
$$

Consequently, for $p \geqslant 2$, and $L_{b_{1}}, L_{b_{2}}, L_{\sigma}$ such that $\mathfrak{b}_{p}>0$ in (A.6), we have

$$
0 \leqslant \lim _{(t-s) \rightarrow \infty} \mathbb{E}|\phi(t, s, x)|^{p} \leqslant \frac{\mathfrak{a}_{p}}{\mathfrak{b}_{p}}<\infty
$$

For $1<p<2$ we use Hölder's inequality and obtain

$$
\mathbb{E}\left[\left|X_{t}^{s, x}\right|^{p}\right] \leqslant\left(\mathbb{E}\left[\left|X_{t}^{s, x}\right|^{2 p}\right]\right)^{\frac{1}{2}}
$$

Thus, analogous derivations to those in (A.3) onwards can be carried out for $p^{\prime}=2 p>2$, with $1<p<2$.

Finally, note that for $Y_{t}^{s, x}(\omega)=\phi(t, s, \omega, x)-x$ we obtain

$$
\mathrm{d} Y_{t}^{s, x}=b\left(t, Y_{t}^{s, x}+x\right) \mathrm{d} t+\sigma\left(t, Y_{t}^{s, x}+x\right) \mathrm{d} W_{t}, \quad Y_{s}^{s, x}=0
$$

so that analogous calculations in conjunction with Assumption (4.4) lead to

$$
\mathbb{E}[V(t, \phi(t, s, x)-x)]^{p} \leqslant \mathfrak{C} \mathbb{E}|\phi(t, s, x)-x|^{p} \leqslant \mathfrak{C} \frac{\mathfrak{a}_{p}}{\mathfrak{b}_{p}}\left(1-e^{\mathfrak{b}_{p}(t-s)}\right)
$$

Consequently, for $p>1$ and $0<L_{b_{1}}, L_{b_{2}}, L_{\sigma}<\infty$ such that $\mathfrak{b}_{p}>0$

$$
\begin{aligned}
& 0 \leqslant \lim _{s \rightarrow-\infty} \sup _{s \leqslant t} \mathbb{E}[V(t, \phi(t, s, x)-x)]<\infty, \\
& 0 \leqslant \lim _{t \rightarrow \infty} \sup _{s \leqslant t} \mathbb{E}[V(\phi(t, s, x)-x)]<\infty
\end{aligned}
$$

Proposition A. 2 (Sharper conditions for existence of absolute moments $p=2,3$ for SDE's with dissipative growth conditions) Consider the same set-up as in Lemma A.1. Suppose further that there exist bounded functions $L_{b_{1}}(\cdot), L_{b_{2}}(\cdot), L_{\sigma}(\cdot) \in \mathcal{C}_{\infty}\left(\mathbb{R} ; \mathbb{R}^{+}\right)$such that

$$
\tilde{\mathfrak{b}}_{2}=2 \inf _{t \in \mathbb{R}}\left(L_{b_{2}}(t)-\frac{1}{2} L_{\sigma}(t)\right)>0
$$

Then, there exists a stochastic flow $\{\phi(t, s, \cdot, \cdot): s \leqslant t\}$ on $\mathbb{R}^{d}$ induced by the global solutions of (2.1), which has a finite second absolute moment for all time. Moreover, if for some $x>0$

$$
\begin{equation*}
\tilde{\mathfrak{b}}_{3}=3 \inf _{t \in \mathbb{R}}\left(L_{b_{2}}(t)-L_{\sigma}(t)-\frac{4}{27 \varkappa^{2}}\left(L_{b_{1}}(t)+L_{\sigma}(t)\right)\right)>0, \tag{A.9}
\end{equation*}
$$

then the stochastic flow has a finite third moment for all time.
Proof For $p=2$, we proceed in a way similar to (A.3), and we have

$$
\begin{align*}
\mathcal{L}_{t}|x|^{2} & \leqslant 2\left(L_{b_{1}}(t)-L_{b_{2}}(t)|x|^{2}\right)+L_{\sigma}(t)\left(1+|x|^{2}\right) \\
& =2\left(L_{b_{1}}(t)+\frac{1}{2} L_{\sigma}(t)\right)-2\left(L_{b_{2}}(t)-\frac{1}{2} L_{\sigma}(t)\right)|x|^{2} \leqslant \tilde{\mathfrak{a}}_{2}-\tilde{\mathfrak{b}}_{2}|x|^{3} \tag{A.10}
\end{align*}
$$

where

$$
\begin{align*}
& \tilde{\mathfrak{a}}_{2}=2 \sup _{t \in \mathbb{R}}\left(L_{b_{1}}(t)+\frac{1}{2} L_{\sigma}(t)\right),  \tag{A.11}\\
& \tilde{\mathfrak{b}}_{2}=2 \inf _{t \in \mathbb{R}}\left(L_{b_{2}}(t)-\frac{1}{2} L_{\sigma}(t)\right) . \tag{A.12}
\end{align*}
$$

Thus, based on (A.8) the second absolute moment exists for all time if $\tilde{\mathfrak{b}}_{2}>0$.
For $p=3$ we have

$$
\begin{aligned}
\mathcal{L}_{t}|x|^{3} & \leqslant 3|x|\left(L_{b_{1}}(t)-L_{b_{2}}(t)|x|^{2}\right)+3|x| L_{\sigma}(t)\left(1+|x|^{2}\right) \\
& =3\left(L_{b_{1}}(t)+L_{\sigma}(t)\right)|x|-3\left(L_{b_{2}}(t)-L_{\sigma}(t)\right)|x|^{3}
\end{aligned}
$$

$$
\begin{align*}
& \leqslant 3 \varkappa\left(L_{b_{1}}(t)+L_{\sigma}(t)\right)-3\left(L_{b_{2}}(t)-L_{\sigma}(t)-\frac{4}{27 \varkappa^{2}}\left(L_{b_{1}}(t)+L_{\sigma}(t)\right)\right)|x|^{3} \\
& \leqslant \tilde{\mathfrak{a}}_{3}-\tilde{\mathfrak{b}}_{3}|x|^{3}, \tag{A.13}
\end{align*}
$$

with

$$
\begin{align*}
& \tilde{\mathfrak{a}}_{3}=3 \varkappa \sup _{t \in \mathbb{R}}\left(L_{b_{1}}(t)+L_{\sigma}(t)\right),  \tag{A.14}\\
& \tilde{\mathfrak{b}}_{3}=3 \inf _{t \in \mathbb{R}}\left(L_{b_{2}}(t)-L_{\sigma}(t)-\frac{4}{27 \varkappa^{2}}\left(L_{b_{1}}(t)+L_{\sigma}(t)\right)\right), \tag{A.15}
\end{align*}
$$

where we used the fact that $|x| \leqslant \varkappa+\frac{4}{27 \varkappa^{2}}|x|^{3}, \varkappa>0$. Thus, based on (A.8) the third absolute moment exists for all time if $\tilde{\mathfrak{b}}_{3}>0$.

Remark A. 3 It is worth noting that, for $L_{b_{1}}, L_{b_{2}}, L_{\sigma}$ constant, and such that $\tilde{\mathfrak{b}}_{3}>0$, the upper bound on the asymptotic moment $\mathbb{E}|\phi(t, s, x)|^{p}$ for $p=3$ is optimised for $\varkappa^{2}=12\left(\frac{L_{b_{1}}+L_{\sigma}}{L_{b_{2}}-L_{\sigma}}\right)$ so that

$$
\begin{equation*}
\min _{\varkappa>0} \frac{\tilde{\mathfrak{a}}_{3}}{\tilde{\mathfrak{b}}_{3}}=(27)^{1 / 2}\left(\frac{L_{b_{1}}+L_{\sigma}}{L_{b_{2}}-L_{\sigma}}\right)^{3 / 2} \tag{A.16}
\end{equation*}
$$

Moreover,

$$
\begin{equation*}
\frac{\tilde{\mathfrak{a}}_{2}}{\tilde{\mathfrak{b}}_{2}} \leqslant\left(\frac{L_{b_{1}}+L_{\sigma}}{L_{b_{2}}-L_{\sigma}}\right) \leqslant\left(\min _{\varkappa>0} \frac{\tilde{\mathfrak{a}}_{3}}{\tilde{\mathfrak{b}}_{3}}\right)^{2 / 3}=(27)^{1 / 3}\left(\frac{L_{b_{1}}+L_{\sigma}}{L_{b_{2}}-L_{\sigma}}\right) ; \tag{A.17}
\end{equation*}
$$

this fact merely reflects the Jensen's inequality for the second and third absolute moments, i.e. $\mathbb{E}\left|X_{t}^{s, x}\right|^{2} \leqslant\left(\mathbb{E}\left|X_{t}^{s, x}\right|^{3}\right)^{2 / 3}$, but it is useful in Example 4.14.

Lemma A. $4 \operatorname{Let}\{\phi(t, s, \cdot, \cdot): t \geqslant s\}$ be a stochastic flow generated by the $\operatorname{SDE}$ (2.1) and let $V \in \mathcal{C}^{1,2}\left(\mathbb{R} \times \mathbb{R}^{d} ; \mathbb{R}^{+}\right)$be a Lyapunov function satisfying the first part of condition (4.10). Suppose further that there exist $L_{b}(\cdot), L_{\sigma}(\cdot), C(\cdot) \in \mathcal{C}_{\infty}\left(\mathbb{R} ; \mathbb{R}^{+}\right)$such that

$$
\begin{equation*}
\langle b(t, x), x\rangle \leqslant L_{b}(t)\left(1+|x|^{2}\right), \quad\|\sigma(t, x)\|_{\mathrm{HS}}^{2} \leqslant L_{\sigma}(t)\left(1+|x|^{2}\right) \tag{A.18}
\end{equation*}
$$

and

$$
\begin{equation*}
0<\limsup _{(t-s) \rightarrow \infty} \exp \left(\int_{s}^{t} C(u, p) \mathrm{d} u\right)<\infty \tag{A.19}
\end{equation*}
$$

where $C(t, p)=L_{b}(t)+\frac{1}{2}(p-1) L_{\sigma}(t)$, for some $1<p<\infty$.
Then, for $x \in \mathbb{R}^{d}$, there exist global solutions of (2.1) such that the stochastic flow $\phi$ induced by the solutions of (2.1) has a finite $p$-th absolute moment for all time. Furthermore, the following holds

$$
\limsup _{(t-s) \rightarrow \infty} \mathbb{E}[V(t, \phi(t, s, x)-x)]<\infty
$$

in Assumption 4.4(iii).

Proof First, we suppose that $p \geqslant 2$ and set $g(x)=1+|x|^{2}$ and $\varphi(x)=g(x)^{\frac{p}{2}}$; then,

$$
\begin{aligned}
\mathcal{L}_{t} \varphi(x)= & p g(x)^{\frac{p}{2}-1} \sum_{i=1}^{d} b_{i}(t, x) x_{i}+\frac{1}{2} p g(x)^{\frac{p}{2}-2} \sum_{i, j=1}^{d} \\
& \times\left\{g(x) \delta_{i j}+(p-2) x_{i} x_{j}\right\}\left(\sigma \sigma^{T}\right)_{i j}(t, x, x) .
\end{aligned}
$$

By the growth conditions (A.18) on the coefficients $b, \sigma$, we obtain

$$
\mathcal{L}_{t} \varphi(x) \leqslant p C(t, p) \varphi(x)
$$

where $C(t, p)=L_{b}(t)+c(t)+\frac{1}{2}(p-1) L_{\sigma}(t)$. Next, let $Y_{t, s}^{x}(\omega)=\phi(t, s, \omega, x)-x$, it follows that $Y_{t, s}^{x}(x)$ solves the following SDE

$$
\mathrm{d} Y_{t, s}^{x}=b\left(t, Y_{t, s}^{x}+x\right) \mathrm{d} t+\sigma\left(t, Y_{t, s}^{x}+x\right) \mathrm{d} W_{t}, \quad Y_{s, s}^{x}=0
$$

By Itô's formula, we have

$$
\begin{aligned}
\mathbb{E}\left[\varphi\left(Y_{t, s}^{x}\right)\right] & =\varphi(0)+\mathbb{E}\left[\int_{s}^{t} \mathcal{L}_{t} \varphi\left(Y_{u, s}^{x}\right) \mathrm{d} u\right] \\
& \leqslant \varphi(0)+p \int_{s}^{t} C(u, p) \mathbb{E}\left[\varphi\left(Y_{u, s}^{x}\right)\right] \mathrm{d} u
\end{aligned}
$$

By Gronwall's inequality, we have

$$
\mathbb{E}\left[\varphi\left(Y_{t, s}^{x}\right)\right] \leqslant \varphi(0) \exp \left(p \int_{s}^{t} C(u, p) \mathrm{d} u\right)
$$

But $Y_{t, s}^{x}(\omega)=\phi(t, s, \omega, x)-x \mathbb{P}$-a.s., $\varphi(x)=g(x)^{\frac{p}{2}}=\left(1+|x|^{2}\right)^{\frac{p}{2}}$ and $\varphi(0)=1$, and thus,

$$
\mathbb{E}\left[\left(1+|\phi(t, s, \omega, x)-x|^{2}\right)^{\frac{p}{2}}\right] \leqslant \exp \left(p \int_{s}^{t} C(u, p) \mathrm{d} u\right)
$$

Next, note for $p \geqslant 2,|x|^{p} \leqslant\left(1+|x|^{2}\right)^{\frac{p}{2}}$ and by the assumption that $V(t, x) \leqslant \mathfrak{C}|x|^{p}$, we obtain

$$
\begin{aligned}
\mathbb{E}[V(t, \phi(t, s, x)-x)] & \leqslant \mathfrak{C} \mathbb{E}\left[\left(1+|\phi(t, s, x)-x|^{2}\right)^{\frac{p}{2}}\right] \\
& \leqslant \mathfrak{C} \exp \left(p \int_{s}^{t} C(u, p) \mathrm{d} u\right)
\end{aligned}
$$

Since $C(\cdot, p) \in \mathcal{C}(\mathbb{R} ; \mathbb{R})$ such that $\limsup _{s \rightarrow-\infty} \exp \left(\int_{s}^{t} C(u, p) \mathrm{d} u\right)<\infty$, then for $p \geqslant 2$, we have

$$
\begin{equation*}
\mathbb{E}[V(t, \phi(t, s, x)-x)] \leqslant \mathfrak{C} \limsup _{s \rightarrow-\infty} \exp \left(p \int_{s}^{t} C(u, p) \mathrm{d} u\right)<\infty . \tag{A.20}
\end{equation*}
$$

The case where $1 \leqslant p<2$, we use Hölder's inequality, namely

$$
\mathbb{E}[V(t, \phi(t, s, x)-x)] \leqslant \mathfrak{C} \mathbb{E}\left[|\phi(t, s, x)-x|^{p}\right] \leqslant \mathfrak{C}\left(\mathbb{E}\left[|\phi(t, s, x)-x|^{2 p}\right]\right)^{\frac{1}{2}}
$$

and the rest follows, since in this case $2 p \geqslant 2$.

## Appendix B: Strong Feller property for flows induced by non-autonomous SDE's

In order to prove Theorem 4.10, which is a generalisation of standard results to nonautonomous SDE's, we first outline some basic notions from Malliavin calculus; the actual proof is given in and discussed in Appendix 5.3.2.

## Malliavin calculus estimates

Establishing the strong Feller property of Markov evolutions $\left(\mathcal{P}_{s, t}\right)_{t \geqslant s}$ in our setting requires some estimates rooted in Malliavin calculus. We recall the main concepts and results on the Wiener space $(\Omega, \mathcal{F}, \mathbb{P})$; see (e.g. [40,43,66,71,74, 88]) for a comprehensive treatment. To this end, consider the Hilbert space $\mathcal{H}=L^{2}\left([s, \infty) ; \mathbb{R}^{m}\right)$ equipped with the inner product

$$
\left\langle\eta_{1}, \eta_{2}\right\rangle_{\mathcal{H}}=\int_{s}^{\infty} \eta_{1}(t) \cdot \eta_{2}(t) \mathrm{d} t
$$

For a Hilbert space $E$ and a real number $p \geqslant 1, L^{p}(\Omega ; E)$ the space of $E$-valued random variable $\xi$ such that $\mathbb{E}\left(\|\xi\|_{E}^{p}\right):=\int_{\Omega}\|\xi\|_{E}^{p} d \mathbb{P}<\infty$. Also, we set $L^{\infty-}(\Omega ; E):=\bigcap_{1 \leqslant p<\infty} L^{p}(\Omega ; E)$.

Following the approach due to Malliavin (e.g. [66,74]), we introduce a derivative operator $\mathcal{D}$ for a random variable $G$ on the space $L^{\infty-}(\Omega ; E)$. We say that $G \in \mathbb{D}^{1, \infty}(E)$ if there exists $\mathcal{D} G \in L^{\infty-}(\Omega ; \mathcal{H} \otimes E)$ such that for any $\eta \in \mathcal{H}$,

$$
\lim _{\varepsilon \rightarrow 0} \mathbb{E}\left\|\frac{G\left(\omega+\varepsilon \int_{s}^{c} \eta(\ell) \mathrm{d} \ell\right)-G(\omega)}{\varepsilon}-\langle\mathcal{D} G, \eta\rangle_{\mathcal{H}}\right\|_{E}^{p}=0
$$

holds for every $p \geqslant 1$. In this case, one defines the Malliavin derivative of $G$ in the direction of $\eta \in \mathcal{H}$ by $\mathcal{D}^{\eta} G:=\langle\mathcal{D} G, \eta\rangle_{\mathcal{H}}$. For any $p \geqslant 1$, we define the Sobolev space $\mathbb{D}^{1, p}(E)$ as the completion of $\mathbb{D}^{1, \infty}(E)$ under the norm

$$
\|G\|_{1, p, E}=\left(\mathbb{E}\|G\|_{\mathcal{H}}^{p}\right)^{1 / p}+\left(\mathbb{E}\|\mathcal{D} G\|_{\mathcal{H} \otimes E}^{p}\right)^{1 / p} .
$$

We define the $k$-th Malliavin derivative by $\mathcal{D}^{k} G=\mathcal{D}\left(\mathcal{D}^{k-1} G\right)$, which is a random variable with values in $\mathcal{H}^{\otimes k} \otimes E$. For any integer $k \geqslant 1$, the Sobolev space $\mathbb{D}^{k, p}(E)$ is the completion of $\mathbb{D}^{k, \infty}(E)$ under the norm

$$
\|G\|_{k, p, E}=\|G\|_{k-1, p, E}+\left\|\mathcal{D}^{k} G\right\|_{1, p, \mathcal{H}^{\otimes k} \otimes E}
$$

It turns out that $\mathcal{D}$ is a closed operator from $L^{p}(\Omega ; E)$ to $L^{p}(\Omega ; \mathcal{H} \otimes E)$. The adjoint $\delta$ of the operator $\mathcal{D}$ called the divergence operator is continuous from $\mathbb{D}^{1, p}(\mathcal{H} \otimes E)$ to $L^{p}(\Omega ; E)$ for any $p>1$, with the duality relationship given as

$$
\begin{equation*}
\mathbb{E}\left[\langle\mathcal{D G}, u\rangle_{\mathcal{H} \otimes E}\right]=\mathbb{E}\left[\langle G, \delta(u)\rangle_{E}\right], \tag{B.1}
\end{equation*}
$$

for any $G \in \mathbb{D}^{1, p}(\mathcal{H} \otimes E)$ and $u \in \mathbb{D}^{1, q}(\mathcal{H} \otimes E)$, with $\frac{1}{p}+\frac{1}{q}=1$.
Throughout the remaining part of this section, we assume the following notation:

- $\quad C$ is a generic constant which may depend on $T$, the exponent $p>1$, the initial point $x$ and fixed element $\eta$ of the Hilbert space $\mathcal{H}=L^{2}\left([s, \infty) ; \mathbb{R}^{m}\right)$.
- $\left(H_{n}\right)$ denotes a class of coefficients $b, \sigma_{k}, 1 \leqslant k \leqslant m$, where $\sigma_{k}$ are columns of $\sigma$, such that $b(t, \cdot) \in \tilde{\mathcal{C}^{n}}, \sigma_{k}(t, \cdot) \in \tilde{\mathcal{C}_{b}^{n}}\left(\mathbb{R}^{d}\right)$.

Proposition B. 1 Suppose the coefficients b, $\sigma$ of the $\operatorname{SDE}(2.1)$ are in the class $\left(H_{2}\right)$. Then, for any $t \geqslant s$, we have $\phi(t, s, \cdot, \cdot) \in \mathbb{D}^{1, \infty}\left(\mathbb{R}^{d}\right)$ and the Malliavin derivative $\mathcal{D}^{\eta} \phi(t, s)$ of $\phi(t, s)$ in the direction of $\eta=\left(\eta^{1}, \eta^{2}, \cdots, \eta^{m}\right) \in \mathcal{H}$ is the unique solution of the following affine $S D E$

$$
\left\{\begin{aligned}
\begin{array}{rl}
\mathrm{d} \mathcal{D}^{\eta} \phi(t, s)=D_{x} b(t, \phi(t, s)) \mathcal{D}^{\eta} \phi(t, s) \mathrm{d} t+ & \sum_{k=1}^{m} D_{x} \sigma_{k}(t, \phi(t, s)) \mathcal{D}^{\eta} \phi(t, s) \mathrm{d} W_{t}^{k} \\
& +\sum_{k=1}^{m} \sigma_{k}(t, \phi(t, s)) \eta^{k}(t) \mathrm{d} t, \quad t>s, \\
\mathcal{D}^{\eta} \phi(s, s)=0 .
\end{array}
\end{aligned}\right.
$$

Corollary B. 2 (Chain rule, cf. [74]) Suppose that condition $\left(H_{2}\right)$ holds true. Then, for any $\eta \in \mathcal{H}, p \geqslant 2$ and for any $\varphi \in \mathcal{C}_{\infty}^{2}\left(\mathbb{R}^{d}\right)$, we have

$$
\lim _{\varepsilon \rightarrow 0} \mathbb{E}\left|\frac{\varphi\left(\phi^{\varepsilon \eta}(t, s)\right)-\varphi(\phi(t, s))}{\varepsilon}-D_{x} \varphi(\phi(t, s)) \mathcal{D}^{\eta} \phi(t, s)\right|^{p}=0
$$

where $\phi^{\varepsilon \eta}(t, s), t \geqslant s, \varepsilon \in(0,1)$ is the solution of the following perturbed $S D E$

$$
\left\{\begin{array}{l}
\mathrm{d} \phi^{\varepsilon \eta}(t, s)=b\left(t, \phi^{\varepsilon \eta}(t, s)\right) \mathrm{d} t+\sum_{k=1}^{m} \sigma_{k}\left(t, \phi^{\varepsilon \eta}(t, s)\right) \mathrm{d} W_{t}^{k}+\varepsilon \sum_{k=1}^{m} \sigma_{k}\left(t, \phi^{\varepsilon \eta}(t, s)\right) \eta^{k}(t) \mathrm{d} t \\
\phi^{\varepsilon \eta}(s, s),=x \in \mathbb{R}^{d}
\end{array}\right.
$$

Moreover, $\varphi(\phi(t, s)) \in \mathbb{D}^{1, \infty}(\mathbb{R})$ and $\mathcal{D} \varphi(\phi(t, s))=D_{x} \varphi(\phi(t, s)) \mathcal{D} \phi(t, s)$.
Definition B. 3 (Mean square gradient) Let $G(x): \Omega \rightarrow \mathbb{R}^{d}$ be a measurable function for all $x \in \mathbb{R}^{d}$ and $i \in F$. We say that the mean square gradient of $G(x)$ with respect to $x$ exists if there is a linear map $A(x): \Omega \rightarrow \mathbb{R}^{d \times d}$ such that for any $v \in \mathbb{R}^{d}$,

$$
\lim _{\varepsilon \rightarrow 0} \mathbb{E}\left|\frac{G(x+\varepsilon v)-G(x)}{\varepsilon}-A(x) v\right|^{2}=0
$$

We denote the mean square gradient matrix $A(x)$ by $D_{x} G(x)$.
Theorem B. 4 (e.g. $[54,66,74])$ Assume the condition $\left(H_{2}\right)$ holds. Let $\phi(t, s, \omega, x), t \geqslant s$, be the solution of the $S D E(2.1)$. Then, the mean square gradient of $\phi(t, s, \cdot, x)$ with respect to $x$ exists. If we define $J_{t, s}=D_{x} \phi(t, s, \cdot, x)$, then

$$
\left\{\begin{array}{l}
\mathrm{d} J_{t, s}=D_{x} b(t, \phi(t, s, \cdot, x)) J_{t, s}+\sum_{k=1}^{m} D_{x} \sigma_{k}(t, \phi(t, s, \cdot, x)) J_{t, s} \mathrm{~d} W_{t}^{k}, \quad t \geqslant s  \tag{B.2}\\
J_{s, s}=I
\end{array}\right.
$$

where I is a $d \times d$ identity matrix. Moreover, the inverse $j_{t, s}^{-1}$ of $J_{t, s}$ exists and satisfy

$$
\left\{\begin{align*}
& \mathrm{d} J_{t, s}^{-1}=-J_{t, s}^{-1}\left(D_{x} b(t, \phi(t, s, \cdot, x))-\sum_{k=1}^{m} D_{x} \sigma_{k}(t, \phi(t, s, \cdot x)) D_{x} \sigma_{k}(t, \phi(t, s, \cdot, x))\right) \mathrm{d} t  \tag{B.3}\\
&-\sum_{k=1}^{m} J_{t, s}^{-1} D_{x} \sigma_{k}(t, \phi(t, s, \cdot, x)) \mathrm{d} W_{t}^{k} \\
& J_{s, s}^{-1}=I .
\end{align*}\right.
$$

We shall refer to the mean square gradient $\left\{J_{t, s}=D_{x} \phi(t, s, \cdot, \cdot): s \leqslant t\right\}$ as derivative flow of $\{\phi(t, s, \cdot, \cdot): s \leqslant t\}$. Next, we provide a crucial $L^{p}$ bound for the derivative flow $J_{t, s}$ and that of its inverse $J_{t, s}^{-1}$.

Lemma B. 5 Suppose the condition $\left(H_{2}\right)$ holds. Then, for any $p \geqslant 2$, there exists a positive constant $C=C(T, p)$ such that

$$
\begin{equation*}
\mathbb{E}\left(\sup _{s \leqslant t, u \leqslant s+T}\left|J_{t, u}\right|^{p}\right) \leqslant C \text { and } \mathbb{E}\left(\sup _{s \leqslant t u \leqslant s+T}\left|J_{t u}^{-1}\right|^{p}\right) \leqslant C \text {. } \tag{B.4}
\end{equation*}
$$

Now, let $\mathcal{D}_{r} \phi(t, s)$ be the solution of the following SDE:

$$
\left\{\begin{align*}
& \mathcal{D}_{u} \phi(t, s)= \sigma(u, \phi(u, s))+\int_{u}^{t} \mathcal{D}_{x} b(\ell, \phi(\ell, u)) \mathcal{D}_{u} \phi(\ell, u) \mathrm{d} \ell  \tag{B.5}\\
&+\sum_{k=1}^{m} \int_{u}^{t} D_{x} \sigma_{k}(\ell, \phi(\ell, u)) \mathcal{D}_{u} \phi(\ell, u) \mathrm{d} W_{\ell}^{k}, \quad \text { for } t \geqslant u \\
& \mathcal{D}_{u} \phi_{t, s}=0, \quad \text { for } s \leqslant t<u
\end{align*}\right.
$$

Comparing the SDE's (B.2) and (B.5), we obtain the following by the variation of parameters formula

$$
\left\{\begin{array}{l}
\mathcal{D}_{u} \phi(t, s)=J_{t, u} \sigma(\phi(t, s)), \quad s \leqslant u \leqslant t \leqslant s+T \\
\mathcal{D}_{u} \phi(t, s)=0, \quad u>t
\end{array}\right.
$$

Next, we recall a result on the Malliavin differentiability of the derivative flow $J_{t, s}, t \geqslant s$. To this end, let's denote by $\mathcal{D}_{u}^{\ell}$ the Malliavin derivative with respect to the $\ell$-th component of the Brownian motion $W$ at time $u$.

Lemma B. 6 Suppose that the condition $\left(H_{3}\right)$ holds. Then, for all $s \leqslant t \leqslant s+T$, $J_{t, s} \in \mathbb{D}^{1, \infty}\left(\mathbb{R}^{d} \otimes \mathbb{R}^{d}\right)$ and for any $p \geqslant 2$, there exists a positive constant $C=C(T, p, x)$, such that for all $j=1, \cdots, m$ and $u \in[s, s+T]$,

$$
\mathbb{E}\left[\sup _{s \leqslant t \leqslant s+T}\left|\mathcal{D}_{u}^{j} J_{t, s}\right|^{p}\right] \leqslant C
$$

Moreover, for any $t \leqslant s+T, X(t, s) \in \mathbb{D}^{2, \infty}\left(\mathbb{R}^{d}\right)$ and for any $p \geqslant 2$, there exists a positive constant $C=C(T, p, x)$ such that for all $j, l=1, \cdots, m$ and $\varsigma, u \leqslant t$

$$
\mathbb{E}\left|\mathcal{D}_{u}^{j}\left(\mathcal{D}_{\varsigma}^{l} \phi(t, s)\right)\right|^{p} \leqslant C
$$

Remark B. 7 If $\left(H_{\infty}\right)$ holds true, then $\phi(t, s) \in \mathbb{D}^{\infty}\left(\mathbb{R}^{d}\right)$ and $J_{t, s} \in \mathbb{D}^{\infty}\left(\mathbb{R}^{d} \otimes \mathbb{R}^{d}\right)$.
Denote by $(\mathcal{D} \phi(t, s))^{T}$ the transpose of the Malliavin derivative $\mathcal{D} \phi(t, s)$. From the relationship $\mathcal{D}_{u} \phi(t, s)=J_{t, u} \sigma(u, \phi(u, s))$, we have $\left(\mathcal{D}_{u} \phi(t, s)\right)^{T}=\sigma(u, \phi(u, s))^{T} J_{t, u}^{T}$.

Definition B. 8 (Malliavin covariance; e.g. $[66,74]$ ) The Malliavin covariance $M_{t, s}$ of the random vector $\phi(t, s)$ is defined by

$$
\begin{aligned}
M_{t, s}=\left\langle\mathcal{D} \phi(t, s),(\mathcal{D} \phi(t, s))^{T}\right\rangle_{\mathcal{H}} & =\int_{s}^{t} J_{t, \ell} \sigma(\ell, \phi(\ell, s)) \sigma(\ell, \phi(\ell, s))^{T} J_{t, \ell}^{T} \mathrm{~d} \ell \\
& =J_{t, s}\left[\int_{s}^{t} J_{\ell, s}^{-1} \sigma(\ell, \phi(\ell, s)) \sigma(\ell, \phi(\ell, s))^{T}\left(J_{\ell, s}^{-1}\right)^{T} \mathrm{~d} \ell\right] J_{t, s}^{T} \\
& =J_{t, s} C_{t, s} T_{t, s}^{T}
\end{aligned}
$$

where $C_{t, s}$ is defined by

$$
C_{t, s}=\int_{s}^{t} J_{\ell, s}^{-1} \sigma(\ell, \phi(\ell, s)) \sigma(\ell, \phi(\ell, s))^{T}\left(J_{\ell, s}^{-1}\right)^{T} \mathrm{~d} \ell
$$

is the so-called reduced Malliavin covariance of $\phi(t, s)$.
We conclude this section by elucidating the invertibility of the Malliavin covariance almost surely and its integrability of all negative orders.

Proposition B. 9 (e.g. [40,74,88]) Suppose Assumption 4.6 holds. Then, for every $t \geqslant s$, the Malliavin covariance matrix $M_{t, s}$ of the random vector $\phi(t, s)$ is invertible $\mathbb{P}$-a.s., and $\mathbb{E}\left[\operatorname{det}\left(M_{t, u}^{-p}\right)\right]<\infty$, for every $t, u \in[s, s+T], T>0$, and $p>1$. Moreover, for any $x \in \mathbb{R}^{d}, s \leqslant t$, the law of $\phi(t, s, \cdot, x)$ is absolutely continuous with respect to the Lebesgue measure on $\mathbb{R}^{d}$ and the probability density is smooth.

## Strong Feller property for non-autonomous dynamics

A transition evolution denoted by $\left(\mathcal{P}_{s, t}\right)_{t \geqslant s}$ (3.3) and induced by a stochastic flow $\{\phi(t, s, \cdot \cdot \cdot): s \leqslant t\}$ has the strong Feller property (i.e. $\mathcal{P}_{s, t} \varphi \in \mathcal{C}_{\infty}\left(\mathbb{R}^{d}\right)$ for any $\left.\varphi \in \mathbb{M}_{\infty}\left(\mathbb{R}^{d}\right)\right)$ if and only if
(a) $\quad\left(\mathcal{P}_{s, t}\right)_{t \geqslant t}$ is a Feller semigroup; i.e. $\mathcal{P}_{s, t}: \mathcal{C}_{\infty}\left(\mathbb{R}^{d}\right) \rightarrow \mathcal{C}_{\infty}\left(\mathbb{R}^{d}\right)$, and
(b) For any $\varphi \in \mathcal{C}_{\infty}\left(\mathbb{R}^{d}\right)$ the family $\left(\mathcal{P}_{s, t} \varphi\right)_{t \geqslant s}$ is equicontinuous.

The first condition follows from the existence of the stochastic flow (see, for example, [44, 54]); here, we are concerned with flows associated with solutions of the non-autonomous SDE (2.1). Thus, we shall only derive the second item.
Intuitively, the strong Feller property states that for sufficiently close initial data $x, y$ and any realisation $\omega$ of the past driving noise, one can construct a coupling between two solutions $\phi(t, s, \omega, x)$ and $\phi(t, s, \omega, y)$ such that with probability close to 1 as $x \rightarrow y$, one has $\phi(t, s, \omega, x)=\phi(t, s, \omega, y)$, for $t \geqslant s$ (e.g. [42,43]). One way of achieving such a coupling (e.g. [41,43]) is via a change of measure on the driving process for one of the two solutions such that the noises $W_{t}^{x}$ and $W_{t}^{y}$ driving the solutions $\phi(t, s, \omega, x)$ and $\phi(t, s, \omega, y)$, are related by $\mathrm{d} W_{t}^{x}=\mathrm{d} W_{t}^{y}+\eta_{t}^{x, y} \mathrm{~d} t$, where $\eta_{t}^{x, y}$ is a control process that steers the solution $\phi(t, s, \omega, x)$ towards the solution $\phi(t, s, \omega, y)$. If one sets $y=x+\varepsilon \eta$ and looks for a control of the
form $\eta_{t}^{x, y}=\varepsilon \eta$, then in the limit as $\varepsilon \rightarrow 0$, the scheme will induce a deformation onto the solution $\phi(t, s, \omega, x)$ after time $t$ in the form of Malliavin derivative of $\phi(t, s, \omega, x)$ at $\omega \in \Omega$ in the direction of $\eta \in \mathcal{H}, \mathcal{H}=L^{2}\left([s, \infty) ; \mathbb{R}^{d}\right)$, i.e. $\langle\mathcal{D} \phi(t, s, \omega, x), \eta\rangle_{\mathcal{H}}=\mathcal{D}^{\eta} \phi(t, s, \omega, x)$. On the other hand, the effect of the perturbation of initial condition by $v$ is given by the directional derivative of the solution $\phi(t, s, \omega, x)$ at $x$ along $v$, i.e. $D_{x} \phi(t, s, \omega, x) v=J_{s, t}(\omega, x) v$. In order to assert the strong Feller property, one has to find a control $\eta^{\nu}$ (e.g. [42,43]) such that

$$
\begin{equation*}
\left\langle\mathcal{D} \phi(t, s, \omega, x), \eta^{\nu}\right\rangle_{\mathcal{H}}=J_{t, s}(\omega, x) v, \tag{B.6}
\end{equation*}
$$

where for brevity of notation we skip the explicit dependence on $\omega$ and $x$.
Theorem B.10 Suppose that Assumption 4.6 hold true. Then, for any $t \in[s, s+T]$, there exist $C_{T}>0$ such that, for any $x, y \in \mathbb{R}^{d}$ and $\varphi \in \mathcal{C}_{\infty}\left(\mathbb{R}^{d}\right)$, we have

$$
\left|\mathcal{P}_{s, t} \varphi(x)-\mathcal{P}_{s, t} \varphi(y)\right| \leqslant C_{T}\|\varphi\|_{\infty}|x-y| .
$$

Proof First, we find a control satisfying (B.6). To this end, for any $v \in \mathbb{R}^{d}$ with $|v|=1$, let $\eta^{\nu}=(\mathcal{D} \phi(t, s))^{T} M_{t, s}^{-1} J_{t, s} v$. Then,

$$
\begin{aligned}
\left\langle\mathcal{D} \phi(t, s), \eta^{v}\right\rangle & =\left\langle\mathcal{D} \phi(t, s),(\mathcal{D} \phi(t, s))^{T} M_{t, s}^{-1} J_{t, s} v\right\rangle_{\mathcal{H}} \\
& =\left\langle\mathcal{D} \phi(t, s),(\mathcal{D} \phi(t, s))^{T}\right\rangle_{\mathcal{H}} M_{t, s}^{-1} J_{t, s} v=J_{t, s} v .
\end{aligned}
$$

Next, we show that $\eta^{\nu} \in \mathbb{D}^{1, p}(\mathcal{H})$ for any $p \geqslant 2$. In fact, by chain rule of differentiation,

$$
\begin{aligned}
\mathcal{D}_{\varsigma}^{k} \eta^{v}= & \left(\mathcal{D}_{\zeta}^{k}(\mathcal{D} \phi(t, s))^{T}\right) M_{t, s}^{-1} J_{t, s} v+(\mathcal{D} \phi(t, s))^{T} M_{t, s}^{-1}\left(\mathcal{D}_{\zeta}^{k} J_{t, s}\right) v \\
& +(\mathcal{D} \phi(t, s))^{T}\left(\mathcal{D}_{\varsigma}^{k} M_{t, s}^{-1}\right) J_{t, s} v \\
= & \left(\mathcal{D}_{\zeta}^{k}(\mathcal{D} \phi(t, s))^{T}\right) M_{t, s}^{-1} J_{t, s} v+(\mathcal{D} \phi(t, s))^{T} M_{t, s}^{-1}\left(\mathcal{D}_{\varsigma}^{\ell} J_{t, s}\right) v \\
& -(\mathcal{D} \phi(t, s))^{T} M_{t, s}^{-1}\left[\left\langle\mathcal{D}_{\zeta}^{k}(\mathcal{D} \phi(t, s)),(\mathcal{D} \phi(t, s))^{T}\right\rangle_{\mathcal{H}}\right. \\
& \left.+\left\langle\mathcal{D} \phi(t, s), \mathcal{D}_{\varsigma}^{k}(\mathcal{D} \phi(t, s))^{T}\right\rangle_{\mathcal{H}}\right] M_{t, s}^{-1} J_{t, s} v .
\end{aligned}
$$

By Lemmas B.5, B.6, and Proposition B. 9 , we arrive at

$$
\begin{equation*}
\mathbb{E}\left\|\eta^{\nu}\right\|_{\mathcal{H}}^{p}+\mathbb{E}\left\|\mathcal{D} \eta^{\nu}\right\|_{\mathcal{H} \otimes \mathcal{H}}^{p} \leqslant \mathbb{E}\left\|\eta^{\nu}\right\|_{\mathcal{H}}^{p}+\sum_{k=1}^{m} \mathbb{E}\left[\int_{s}^{t}\left\|\mathcal{D}_{\varsigma}^{k} \eta^{v}\right\|_{\mathcal{H}}^{p} d \varsigma\right]<\infty \tag{B.7}
\end{equation*}
$$

Recalling that $\eta_{\varsigma, s}^{\nu}=\sigma(\varsigma, X(\varsigma, s))^{T} J_{t, 5}^{T} M_{t, s}^{-1} J_{t, s} v$. Then, for $\varphi \in \mathcal{C}_{\infty}^{1}\left(\mathbb{R}^{d}\right)$, we have

$$
\begin{align*}
D_{x}\left(\mathcal{P}_{t, s} \varphi\right)(x) v & =\mathbb{E}\left[D_{x}[\varphi(\phi(t, s, x))] \nu\right]=\mathbb{E}\left[D_{x} \varphi(\phi(t, s, x)) J_{t, s}(\cdot, x) v\right] \\
& =\mathbb{E}\left[\left\langle D_{x} \varphi(\phi(t, s, x)) \mathcal{D} X_{t}^{s, x}(\omega), \eta^{\nu}\right\rangle_{\mathcal{H}}\right]=\mathbb{E}\left[\left\langle\mathcal{D} \varphi(\phi(t, s, x)), \eta^{\nu}\right\rangle_{\mathcal{H}}\right] \\
& =\mathbb{E}\left[\left\langle\mathcal{D} \varphi(\phi(t, s, x)), \eta^{v}\right\rangle_{\mathcal{H}}\right]=\mathbb{E}\left[\varphi(\phi(t, s, x)) \delta\left(\eta^{\nu}\right)\right] \\
& =\mathbb{E}\left[\varphi(\phi(t, s, x)) \int_{s}^{t} \sigma(\varsigma, \phi(\varsigma, s))^{T} J_{t, \zeta}^{T} M_{t, s}^{-1} J_{t, s} \star \mathrm{~d} W_{\varsigma}\right] \\
& =\mathbb{E}\left[\varphi(\phi(t, s, x)) \int_{s}^{t} \eta_{\zeta, s}^{v} \star \mathrm{~d} W_{\varsigma}\right] \tag{B.8}
\end{align*}
$$

where in the first and second lines, we applied chain rule for mean square gradient and Malliavin derivative, respectively, and third line is Malliavin integration by parts formula
(e.g. [74]) and, the stochastic integral in the fourth or fifth line is interpreted in the sense of Skorokhod, i.e.
$\int_{s}^{t} \eta_{\zeta, s}^{v} \star \mathrm{~d} W_{\zeta}$ is the divergence of the process $\left\{\eta_{\zeta, s}^{v} \mathbb{I}_{[s, t]}(\varsigma): \varsigma \geqslant s\right\}$ (see equation (B.1)).
Next, since $\mathcal{C}_{\infty}^{1}\left(\mathbb{R}^{d}\right)$ is dense in $\mathcal{C}_{\infty}\left(\mathbb{R}^{d}\right)$, we have

$$
\left(\varphi_{n}\right)_{n \in \mathbb{N}} \subset \mathcal{C}_{\infty}^{1}\left(\mathbb{R}^{d}\right), \varphi_{n} \underset{n \rightarrow \infty}{\longrightarrow} \varphi \in \mathcal{C}_{\infty}\left(\mathbb{R}^{d}\right)
$$

so that

$$
\left\{\begin{array}{l}
\lim _{n \rightarrow \infty} \mathcal{P}_{s, t} \varphi_{n}(x)=\mathcal{P}_{s, t} \varphi(x),  \tag{B.9}\\
\lim _{n \rightarrow \infty} D_{x}\left(\mathcal{P}_{s, t} \varphi_{n}\right)(x) v=\mathbb{E}\left[\varphi(\phi(t, s, x)) \int_{s}^{t} \eta_{\zeta, s}^{\nu} \star \mathrm{d} W_{\zeta}\right] .
\end{array}\right.
$$

On the other hand, by Proposition B.9, there exists a function $0<p_{s, t} \in \mathcal{C}_{\infty}^{\infty}\left(\mathbb{R}^{d}\right) \times \mathcal{C}_{\infty}^{\infty}\left(\mathbb{R}^{d}\right)$ such that $\mathbb{P}(\{\omega: \phi(t, s, \omega, x) \in \mathrm{d} y\})=P(s, x ; t, \mathrm{~d} y)=p_{s, t}(x, y) \mathrm{d} y$. This implies that

$$
\begin{align*}
\lim _{n \rightarrow \infty} D_{x}\left(\mathcal{P}_{s, t} \varphi_{n}\right)(x) v & =\lim _{n \rightarrow \infty} \int_{\mathbb{R}^{d}} \varphi_{n}(y) D_{x} p_{s, t}(x, y) v \mathrm{~d} y \\
& =\int_{\mathbb{R}^{d}} \varphi(y) D_{x} p_{s, t}(x, y) v \mathrm{~d} y=D_{x}\left(\mathcal{P}_{s, t} \varphi\right)(x) v \tag{B.10}
\end{align*}
$$

Comparing (B.9) and (B.10), we have that (B.8) holds for all $\varphi \in \mathcal{C}_{\infty}\left(\mathbb{R}^{d}\right)$.
Next, the Cauchy-Schwartz inequality yields

$$
\begin{equation*}
\left|D_{x}\left(\mathcal{P}_{s, t} \varphi\right)(x) v\right| \leqslant \sqrt{\left(\mathcal{P}_{s, t} \varphi^{2}\right)(x)}\left(\mathbb{E}\left|\int_{s}^{t} \eta_{\zeta, s}^{\nu} \star \mathrm{d} W_{\zeta}\right|^{2}\right)^{1 / 2}, \quad \varphi \in \mathcal{C}_{\infty}\left(\mathbb{R}^{d}\right) \tag{B.11}
\end{equation*}
$$

By generalised Itô isometry (cf. [74]), we have

$$
\begin{aligned}
\mathbb{E}\left|\int_{s}^{t} \eta_{\zeta, s}^{v} \star \mathrm{~d} W_{\varsigma}\right|^{2} & =\mathbb{E}\left(\int_{s}^{t}\left|\eta_{\zeta, s}^{v}\right|^{2} d \varsigma\right)+\mathbb{E}\left(\int_{s}^{t} \int_{s}^{t}\left\langle\mathcal{D}_{\xi} \eta_{\xi, s}^{v}, \mathcal{D}_{\varsigma} \eta_{\zeta, s}^{v}\right\rangle_{\mathbb{R}^{d} \otimes \mathbb{R}^{d}} d \xi d \varsigma\right) \\
& \leqslant \mathbb{E}\left(\int_{s}^{t}\left|\eta_{\zeta, s}^{v}\right|^{2} d \varsigma\right)+\mathbb{E}\left(\int_{s}^{t} \int_{s}^{t}\left\|\mathcal{D}_{\xi} \eta_{\zeta, s}^{v}\right\|_{\mathbb{R}^{d} \otimes \mathbb{R}^{d}}^{2} d \xi d \varsigma\right) \\
& =\mathbb{E}\left\|\eta^{v}\right\|_{\mathcal{H}}^{2}+\sum_{k=1}^{m} \mathbb{E}\left(\int_{s}^{t}\left\|\mathcal{D}_{\xi}^{k} \eta^{v}\right\|_{\mathcal{H}}^{2} d \xi\right) .
\end{aligned}
$$

Then, by the inequality (B.7) and (B.11), there exists $C_{T}>0$ such that

$$
\begin{equation*}
\left|D_{x}\left(\mathcal{P}_{s, t} \varphi\right)(x) v\right| \leqslant C_{T}\|\varphi\|_{\infty}|v|, \quad x, v \in \mathbb{R}^{d}, \varphi \in \mathcal{C}_{\infty}\left(\mathbb{R}^{d}\right) \tag{B.12}
\end{equation*}
$$

Finally, let $z_{\ell}=\ell x+(1-\ell) y, \quad \ell \in[0,1]$ and set $v=x-y$. Then, by the mean value theorem and inequality (B.12), we have

$$
\left|\mathcal{P}_{s, t} \varphi(x)-\mathcal{P}_{s, t} \varphi(y)\right| \leqslant \int_{0}^{1}\left|D_{x}\left(\mathcal{P}_{s, t} \varphi\right)\left(z_{\ell}\right) v\right| \mathrm{d} \ell \leqslant C_{T}\|\varphi\|_{\infty}|x-y|
$$

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[^0]:    ${ }^{1}$ The meaning of the term 'climate' used in most theoretical work in atmosphere-ocean science is loosely related to properties of the probability measure induced by the time-asymptotic dynamics.

[^1]:    ${ }^{2}$ Here, the notion of a 'Markovian RDS' means that there exists a version of the RDS which has the Markov property w.r.t. the canonical filtration generated on the Wiener space by the Brownian motion discussed later.

[^2]:    ${ }^{3}$ Strictly, $\mathcal{L}_{t}$ in (3.7) coincides with the generator of (2.1) on $f \in \mathcal{C}_{c}^{1,2}\left(\mathbb{R} \times \mathbb{R}^{d}\right.$, $\left.\mathbb{R}^{+}\right)$, but it is well defined for $f \in \mathcal{C}^{1,2}\left(\mathbb{R} \times \mathbb{R}^{d}, \mathbb{R}^{+}\right)$and we refer to $\mathcal{L}_{t}$ as the generator throughout; the same applies to $\mathcal{L}_{t}^{(2)}$ in (3.8).

[^3]:    ${ }^{5}$ The generation of an RDS from an SDE requires a 'perfection of the crude cocycle' associated with the SDE (see, for example, [8, Theorem 2.3.26]); here, this important technical nuance does not require an explicit discussion.

[^4]:    ${ }^{6}$ See (4.13) for one such class which we focus on in this work.
    ${ }^{7}$ Sufficient conditions for the existence of the minimal period, which are satisfied here, are established in [33, $\left.\$ 5\right]$.

[^5]:    ${ }^{8}$ See Definition 3.7.
    ${ }^{9}$ See, for example, [23] for a detailed description of such structures on spaces of probability measures.
    ${ }^{10}$ To be more accurate, the so-called crude cocycle property can be easily verified from the flow induced by the SDE, and the crude cocycle needs to be 'perfected' in order to generate an RDS over the DS for the Brownian motion (see, for example, [8, Theorem 2.3.26]); here, this important technical nuance does not require an explicit discussion.

[^6]:    ${ }^{11}$ This condition can be replaced by a stronger but a more concrete constraint on the global existence of the $p$-th absolute moment of $\phi$; see Lemma A. 1 in "Appendix A".

[^7]:    ${ }^{12}$ The transition evolution $\left(\mathcal{P}_{s, t}\right)_{t \geqslant s}$ on a complete separable metric space $\mathcal{X}$ has strong Feller property if for $\varphi \in \mathbb{M}_{\infty}(\mathcal{X})$, one has $\mathcal{P}_{s, t} \varphi \in \mathcal{C}_{\infty}(\mathcal{X}), \forall s \leqslant t$, i.e. $\mathcal{P}_{s, t}: \mathbb{M}_{\infty}(\mathcal{X}) \rightarrow \mathcal{C}_{\infty}(\mathcal{X}), \forall s \leqslant t$.

[^8]:    ${ }^{13}$ Here, the notion of a 'Markovian RDS' means that there exists a version of the RDS which has the Markov property w.r.t the filtration generated on the Wiener space by the canonical DS for the Wiener process with $\tilde{W}_{t+s}\left(\theta_{-s} \omega\right)$ for all $s \in[0, \tau), t \in \mathbb{R}^{+}$; see Proposition 3.9.

[^9]:    ${ }^{14}$ These statements are not restricted to the skew-product representation of time-periodic measures.
    ${ }^{15}$ As before, we exclude the stationary regime from the random periodic regime by requiring that fundamental period $0<\tau<\infty$; see Definition 4.1.

[^10]:    ${ }^{16}$ This identity is obtained with the help of (4.43) by noticing that one has

    $$
    \begin{aligned}
    \langle B(v-w, v), v-w\rangle & =\langle B(v, v), v-w\rangle-\langle B(w, v), v-w\rangle \\
    & =\langle B(v, v), v-w\rangle-\langle B(w, w), v-w\rangle-\langle B(w, v-w), v-w\rangle \\
    & =\langle G(v), v-w\rangle-\langle G(w), v-w\rangle
    \end{aligned}
    $$

    where the last term in the second line vanishes due to the fact that $\langle B(u, w), w\rangle_{\mathbb{R}^{3}}=0 \forall u, w \in \mathbb{R}^{3}$.

[^11]:    ${ }^{17}$ This condition can be replaced by a stronger but more concrete constraint on the global existence of the $p$-th absolute moment of $\phi^{\alpha}$; see Lemma A. 1 in "Appendix A".

