Ergodic Numerical Approximation to Periodic Measures of Stochastic Differential Equations

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Abstract

In this paper, we consider numerical approximation to periodic measure of a time periodic stochastic differential equations (SDEs) under weakly dissipative condition. For this we first study the existence of the periodic measure \( \rho_t \) and the large time behaviour of \( U(t+s,s,x) := \mathbb{E}\phi(X^{s,x}_t) - \int \phi d\rho_t \), where \( X^{s,x}_t \) is the solution of the SDEs and \( \phi \) is a test function being smooth and of polynomial growth at infinity. We prove \( U \) and all its spatial derivatives decay to 0 with exponential rate on time \( t \) in the sense of average on initial time \( s \). We also prove the existence and the geometric ergodicity of the periodic measure of the discretized semi-flow from the Euler-Maruyama scheme and moment estimate of any order when the time step is sufficiently small (uniform for all orders). We thereafter obtain that the weak error for the numerical scheme of infinite horizon is of the order 1 in terms of the time step. We prove that the choice of step size can be uniform for all test functions \( \phi \). Subsequently we are able to estimate the average periodic measure with ergodic numerical schemes.

Keywords: Periodic measure; Fokker-Planck equation; discretized semi-flows; geometrical ergodicity; weak approximation.

Mathematics Subject Classifications (2010): 37H99, 60H10, 60H35.

1 Introduction

Random periodicity is ubiquitous in the real world from daily temperature process to economic cycles. The concepts of random periodic paths and periodic measures were introduced and their ergodicity was obtained recently ([10],[11],[12],[15],[35]). They are two different indispensable ways in the pathwise sense and in distributions to describe random periodicity. The “equivalence” of the random periodic solutions and periodic measures and their characterisation in terms of purely imaginary eigenvalues of the infinitesimal generator of the Markovian semigroup were obtained in [12]. The presence of pure imaginary eigenvalues distinguishes the random periodic processes/periodic measures regime from that of the stationary processes/mixing invariant measures, in the latter case the Koopman-von Neumann Theorem says the infinitesimal generator has a unique eigenvalue 0 on the imaginary axis.

As in the case of deterministic dynamical systems where periodic motion has been in the central stage of its study, the relevance of random periodic paths and periodic measures to theoretical and applied problems arising in stochastic dynamical systems has begun to be realised. In particular, there has been progress in the study of some topics in stochastic dynamics e.g. bifurcations (Wang [32]), random attractors (Bates, Lu and Wang [2]), stochastic resonance
(Cherubini, Lamb, Rasmussen and Sato [6], Feng, Zhao and Zhong [13],[14], random horseshoes (Huang, Lian and Lu [13]), modelling the El Nño phenomenon (Chekroun, Simonnet and Ghil [5]), isochronicity of stochastic oscillations (Engel and Kuehn [7]), and invariant measures of quasi-periodic stochastic systems (Feng, Qu and Zhao [9]).

However, it is difficult to construct random periodic solutions explicitly for many problems. So numerical approximation is critical in the study of stochastic dynamics in addition to the study of random periodic dynamics theory. There are numerous works on numerical analysis of SDEs on a finite horizon ([21],[25],[20],[26]). A numerical analysis of approximation to invariant measures of SDEs through discretizing the pull-back, was given in [23],[28],[30],[31],[33],[34]. Numerical approximations to stable zero solutions of SDEs were given in [17],[21]. Despite the importance both on the theoretical and applied aspects of the random periodic regime, its numerical analysis has barely been developed. The only result we know is the pathwise approximations of the random periodic solutions of SDEs discussed in [8]. In this paper we study the weak approximation to periodic measures.

We consider the following non-autonomous stochastic differential equations on \( \mathbb{R}^d \)
\[
dX_t = b(t, X_t)dt + \sigma(X_t)dW_t, \quad t \geq s, \tag{1.1}
\]
with initial condition \( X_s = x \) where \( b : \mathbb{R} \times \mathbb{R}^d \rightarrow \mathbb{R}^d, \sigma : \mathbb{R}^d \rightarrow \mathbb{R}^{d \times d} \), \( W_t \) is a two-sided Wiener process in \( \mathbb{R}^d \) on the Wiener probability space \( (\Omega, \mathcal{F}, \mathbb{P}) \). We assume that \( b \) is \( \tau \)-periodic in the time variable and weakly dissipative in the space variable. For a technical reason, here we only consider the case when \( \sigma \) is time independent as we use the results in [13],[18]. Denote by \( X_i^{s,x} \) the solution of (1.1) throughout the paper.

The existence of the periodic measure was studied in [13]. Under the assumption that the drift term is weakly dissipative and the diffusion term is non-degenerate, it was proved that the periodic measure \( \rho_x \) exists and has a density function, denoted by \( q(s,x) \). To obtain \( q(s,x) \), one could solve the corresponding infinite horizon Fokker-Planck equation \( \frac{\partial}{\partial s} q(s,x) = L^*(s)q(s,x) \) with additional condition \( q(s+\tau,x) = q(s,x) \), where \( L^*(s) \) is the adjoint of the infinitesimal generator of the process \( X_t \). This partial differential equation is generally difficult to solve explicitly. But we will show that theoretically it plays an essential role in establishing the theory of numerical schemes of weak approximations.

We apply numerical schemes such as Euler-Maruyama method to estimate the periodic measure. For any fixed \( i \in \mathbb{Z} \), denoted by \( \{X_i^{s+\Delta t+n\Delta t}\}_{n=0,1,\ldots} \) the discrete approximation of the solution of (1.1) with step size \( \Delta t = \frac{s}{n} \) and \( X_i^{s,x} = x \). We prove that the discrete semi-flow is geometrically ergodic and has a periodic measure \( \tilde{\rho}_i^{\Delta t}(\cdot), \ i \in \mathbb{Z} \).

In this paper, we use the idea of lifting the flow and periodic measure to the cylinder \([0,\tau) \times \mathbb{R}^d\) proposed in [12]. With the help of this tool, our main result is to prove that the cumulative discretization errors is of the order of \( O(\Delta t) \) for the approximation of the average of periodic measure, i.e. for any \( \phi \in C^\infty_p (\mathbb{R}^d) \)
\[
\left| \int_{\mathbb{R}^d} \phi(x) \tilde{\rho}(dx) - \int_{\mathbb{R}^d} \phi(x) \tilde{\rho}^{\Delta t}(dx) \right| \leq C \Delta t, \tag{1.2}
\]
where \( \tilde{\rho} := \frac{1}{\tau} \int_0^\tau \rho_s ds \) ([12]), \( \tilde{\rho}^{\Delta t} := \frac{1}{\tau} \sum_{n=0}^{N-1} \rho^{\Delta t}_i, C^\infty_p (\mathbb{R}^d) \) is the space of smooth functions with the property that themselves and all their derivatives have at most polynomial growth at infinity. In fact, (1.2) only holds for \( \Delta t \) being small enough and the choice of step size can be uniform for all \( \phi \in C^\infty_p \). For this, the uniformity of the step size working for all moment estimates of \( \tilde{X}_i^{s,x} \) is derived in Proposition 4.3. The error estimate (1.2) can also be numerically verified.

The results in this paper are applicable for many physically relevant SDEs, for instance, Benzi-Parisi-Sutera-Vulpiani’s stochastic resonance model for the ice-age transition in climate
change dynamics is SDE (1.1), with $b(t, x) = x - x^3 + A \cos(Bt)$ and $\sigma(x) = \sigma$ being constant (\[1\]). It was proved that this model has a unique periodic measure (\[13\]). This result implies the transition between ice-age and interglacial climates. A partial differential equation for expected transition between ice-age and interglacial climates.

2 Preliminary results and notation

2.1 Lifts of semi-flows, random periodic paths and periodic measures

Denote by $(\Omega, \mathcal{F}, \mathbb{P}, (\theta_t)_{t \in \mathbb{R}})$ the metric dynamical system associated with the canonical probability space $(\Omega, \mathcal{F}, \mathbb{P})$ for Brownian motion $W$ in $\mathbb{R}^d$, where $\theta_t : \Omega \to \Omega$ defined by $(\theta_t \omega)(t) = W(t + s) - W(s)$, is measurably invertible for all $s \in \mathbb{R}$. Denote $\Delta := \{(t, s) \in \mathbb{R}^2, s \leq t\}$ and let $u : \Delta \times \Omega \times \mathbb{R}^d \to \mathbb{R}^d$ be a periodic stochastic semi-flow of period $\tau$ satisfying for all $(t, s) \in \Delta$ and $r \in [s, t]$

$$u(t, r, \omega) \circ u(r, s, \omega) = u(t, s, \omega), \quad \text{for almost all } \omega \in \Omega,$$

and

$$u(t + \tau, s + \tau, \omega) = u(t, s, \theta_\tau \omega), \quad \text{for almost all } \omega \in \Omega.$$

Here $\tau > 0$ is a deterministic real number. Solutions of stochastic differential equations (1.1) with coefficients being periodic in time with period $\tau$, when they exist and unique, generates a periodic semi-flow $u(t, s)x = X_{t}^{s,x}$, which satisfies the above two properties. As we consider periodic measures in this paper, so perfection is not needed here.

Consider the case when $u(t + s, s, \cdot)$ is a Markovian semi-flow on a filtered dynamical system $(\Omega, \mathcal{F}, \mathbb{P}, (\theta_t)_{t \in \mathbb{R}}, (\mathcal{F}_t^\omega)_{s \leq t})$, i.e., for any $s, t, r \in \mathbb{R}$, $s \leq t$, we have $\theta_r^{-1}\mathcal{F}_s^\omega = \mathcal{F}_{s+r}^\omega$ and $u(t + s, s, \cdot)$ is independent of $\mathcal{F}_\infty^\omega$, where $\mathcal{F}_\infty^\omega := \bigvee_{r \leq t} \mathcal{F}_r^\omega$. For any $\Gamma \in \mathcal{B}(\mathbb{R}^d)$, $t \in \mathbb{R}^+$, $s \in \mathbb{R}$, denote the transition probability of $u$ by $P(t + s, s, x, \Gamma) = \mathbb{P}(\{\omega : u(t + s, s, \omega)x \in \Gamma\})$. From the periodicity of semi-flow $u$ and the measure preserving property of $\theta_\tau$, the transition probability $P(t + s, s, x, \cdot)$ satisfies the periodic relation

$$P(t + s + \tau, s + \tau, x, \cdot) = P(t + s, s, x, \cdot). \quad (2.1)$$

Define for $\phi \in \mathcal{B}_b(\mathbb{R}^d)$, the space of bounded and Borel measurable function from $\mathbb{R}^d$ to $\mathbb{R}$,

$$\mathcal{T}(t + s, s)\phi(x) := \mathbb{E}\phi(X_{t+s}^{s,x}) = \int_{\mathbb{R}^d} \phi(y)P(t + s, s, x, dy), \quad t \geq 0.$$

Then it is well-known that $\mathcal{T}(t + s, s) : \mathcal{B}_b(\mathbb{R}^d) \to \mathcal{B}_b(\mathbb{R}^d)$ defines a semigroup and satisfies the $\tau$-periodic property:

$$\mathcal{T}(t + s + \tau, s + \tau) = \mathcal{T}(t + s, s),$$
This follows from (2.1) and the definition of $T(t + s, s)$ easily. Moreover for any probability measure $\rho \in \mathcal{P}(\mathbb{R}^d)$, the space of probability measures on $(\mathbb{R}^d, \mathcal{B}(\mathbb{R}^d))$, define
\[
(T^* (t + s, s) \rho)(\Gamma) = \int_{\mathbb{R}^d} P(t + s, s, x, \Gamma) \rho(dx).
\]
The definition of periodic measure of the periodic Markovian semi-group is given as follows. The existence of the periodic measure was proved in [13] for a wide class of SDEs.

**Definition 2.1.** (13) The measure valued function $\rho : \mathbb{R} \to \mathcal{P}(\mathbb{R}^d)$ is called a $\tau$-periodic measure of the $\tau$-periodic Markovian semi-group $T$ if
\[
T^* (t + s, s) \rho_s = \rho_{t+s}, \quad \forall s \in \mathbb{R}, \ t \in \mathbb{R}^+.
\] (2.2)

The idea of lifting a stochastic periodic semi-flow to a cocycle on a cylinder in [12] plays an important role in this paper. As for this paper, the relevant part is briefly discussed below. Let $S = [0, \tau) \times \mathbb{R}^d$, the lifted cocycle arising from SDE (1.1) with coordinates $X_s = (s, X_s) \in S$ is given by
\[
d\tilde{X}_t = \tilde{b}(\tilde{X}_t)dt + \tilde{\sigma}(\tilde{X}_t)d\tilde{W}(t),
\]
where $\tilde{X}^{0,x}_0 = \tilde{x} = (s, x), \tilde{W} = (\tilde{W}_0, W), \tilde{W}_0$ is a one-dimensional Brownian motion which is independent of $W$, $\tilde{b}(\tilde{X}_t) = \left(\begin{array}{c} 1 \\ b(t, X_t) \end{array}\right)$, $\tilde{\sigma}(\tilde{X}_t) = \left(\begin{array}{cc} 0 & 0 \\ 0 & \sigma(X_t) \end{array}\right)$. One can enlarge the probability space $(\Omega, \mathcal{F}, \mathbb{P})$, still denoted by $(\Omega, \mathcal{F}, \mathbb{P})$, as the canonical probability space for $\mathbb{R}^{d+1}$ Brownian motion $\tilde{W}$.

It is easy to see that the infinitesimal generator of the process $\tilde{X}$ is given by
\[
\tilde{L} = \sum_{i=1}^d b_i(s, x) \frac{\partial}{\partial x_i} + \frac{1}{2} \sum_{i,j=1}^d a_{ij}(x) \frac{\partial^2}{\partial x_i \partial x_j} + \frac{\partial}{\partial s} =: \mathcal{L}(s) + \frac{\partial}{\partial s},
\] (2.3)
where $a = (a_{ij}) = \sigma \sigma^T$, and $U(t + s, s, x) := T(t + s, s) \phi(x)$ satisfies
\[
\frac{\partial}{\partial t} U(t + s, s, x) = \tilde{L} U(t + s, s, x), \quad U(s, s, x) = \phi(x),
\] (2.4)
provided $U$ is sufficiently smooth. Meanwhile, the transition probability and the periodic measure are lifted to
\[
\tilde{\rho}_s(C \times \Gamma) = \delta_{\{s \text{ mod } \tau\}}(C) \rho_s(\Gamma),
\]
\[
\tilde{P}(t, \tilde{x}, C \times \Gamma) = \delta_{\{t+s \text{ mod } \tau\}}(C) P(t + s, s, \Gamma) = \mathbb{P}(\omega : X^{t,z}_{t+s} \in \Gamma),
\]
where $C \in \mathcal{B}([0, \tau))$ and $\Gamma \in \mathcal{B}(\mathbb{R}^d)$. It was shown in [12] that $\tilde{P}$ generates a homogeneous semigroup $\tilde{T}$ defined by $(\tilde{T} \tilde{\phi})(\tilde{x}) = \int_{\mathbb{S}} \tilde{\phi}(\tilde{y}) \tilde{P}(t, \tilde{x}, d\tilde{y})$ and $\tilde{\rho}_s$ is a periodic measure of the lifted semigroup $\tilde{T}$. It was also noticed that $\tilde{\rho} = \frac{1}{\tau} \int_0^\tau \tilde{\rho}_s ds$ is an invariant measure of $\tilde{P}$ following a standard procedure of Fubini theorem. It is easy to see that for a measurable function $\phi : \mathbb{R}^d \to \mathbb{R}$,
\[
\int_0^\tau \int_{\mathbb{R}^d} \phi(x) \tilde{\rho}(dt, dx) = \frac{1}{\tau} \int_0^\tau \int_0^\tau \int_{\mathbb{R}^d} \phi(x) \delta_{\{s \text{ mod } \tau\}}(dt) \rho_s(dx) ds
\] (2.5)
\[= \frac{1}{\tau} \int_0^\tau \int_{\mathbb{R}^d} \phi(x) \rho_s(ds) dx.
\]
\[= \int_{\mathbb{R}^d} \phi(x) \tilde{\rho}(dx),
\]
where $\tilde{\rho} = \frac{1}{\tau} \int_0^\tau \rho_s ds$. This will be used in later part of this paper.
2.2 Assumptions and some preliminary estimates

Assume

**Condition (1)** The functions $b, \sigma$ are of class $C^\infty$ with $\sigma$ being bounded, $b$ and $\sigma$ having bounded derivatives of any order and $b$ being $\tau$-periodic with respect to time.

**Condition (2)** (Uniform ellipticity) There exists a positive constant $\alpha$ such that for any $x, y \in \mathbb{R}^d$, we have $\sum_{i,j} a_{ij}(y)x_i x_j \geq \alpha |x|^2$.

**Condition (3)** (Weak dissipativity) There exist constants $\beta > 0$ and $C > 0$ such that for any $t \in \mathbb{R}^+$ and any $x \in \mathbb{R}^d$, $x \cdot b(t, x) \leq -\beta |x|^2 + C$.

Under conditions (1)-(3), it was proved in [13] that the periodic measure $\rho : (-\infty, +\infty) \rightarrow \mathcal{P}(\mathbb{R}^d)$ exists and is geometrically ergodic:

$$\|P(n\tau + s, s, x) - \rho_s\|_{TV} \leq Ce^{-\delta n\tau}.$$ 

We now discuss the existence of the density function $q(s, x)$ of the periodic measure $\rho_s$. Set the Fokker-Planck operator as follows

$$\tilde{L}^*(s) \cdot q = -\sum_{i=1}^d \frac{\partial}{\partial x_i} (b_i(s, x) \cdot q) + \frac{1}{2} \sum_{i,j=1}^d \frac{\partial^2}{\partial x_i \partial x_j} (a_{ij}(x) \cdot q),$$

and

$$\tilde{L}^* = \tilde{L}^*(s) - \frac{\partial}{\partial s}.$$ 

**Proposition 2.2.** Assume Conditions (1), (2) and (3). Then the periodic measure $\rho_s$ has a density $q(s, x)$ with respect to the Lebesgue measure in $\mathbb{R}^d$, and the density is the unique bounded solution of the Fokker-Planck equation

$$\frac{\partial}{\partial s} q(s, x) = \tilde{L}^*(s) q(s, x), \quad (2.6)$$

satisfying that for any $s \in [0, \tau)$, $q(s + \tau, x) = q(s, x)$ and $q(s, x) \to 0$ as $|x| \to \infty$.

**Proof.** Under the assumption of this proposition, the $\tau$-periodic two-parameter Markov transition probability $P(t, s, x, \Gamma)$ has a density $p(t, s, x, y)$. Thus we have the representation of periodic measure $\rho_s$ as follows, for any $\Gamma \in \mathcal{B}(\mathbb{R}^d)$,

$$\rho_s(\Gamma) = \int_{\mathbb{R}^d} \int_{\Gamma} p(s + \tau, s, x, y) dy \rho_s(dx) = \int_{\Gamma} \int_{\mathbb{R}^d} p(s + \tau, s, x, y) \rho_s(dx) dy,$$

where we applied Fubini’s theorem. Hence we get the formula of the density of $\rho_s$ as

$$q(s, y) = \int_{\mathbb{R}^d} p(s + \tau, s, x, y) \rho_s(dx). \quad (2.7)$$

It is easy to prove the periodicity of the density $q(s, y)$ by the periodic property of both $p(t, s, x, y)$ and $\rho_s$. Moreover, we have that for any $\Gamma \in \mathcal{B}(\mathbb{R}^d)$,

$$\rho_{t+\tau}(\Gamma) = \int_{\Gamma} \int_{\mathbb{R}^d} p(t + \tau, s + \tau, z, y) \int_{\mathbb{R}^d} p(s + \tau, s, x, z) \rho_s(dx) dy dz dy$$

$$= \int_{\Gamma} \int_{\mathbb{R}^d} p(t, s, z, y) q(s, z) dz dy.$$
As the periodic measure satisfies $\rho_t(\Gamma) = \rho_{t+\gamma}(\Gamma)$, the above implies

$$q(t, y) = \int_{\mathbb{R}^d} p(t, s, z, y)q(s, z)dz.$$  \hspace{1cm} (2.8)

It is well known that $p(t, s, y, x)$ satisfies the Fokker-Planck equation $\partial_t p(t, s, y, x) = \tilde{L}^*(t) p(t, s, y, x)$. Therefore,

$$\partial_t q(t, x) = \int_{\mathbb{R}^d} \tilde{L}^*(t) p(t, s, y, x) q(s, y) dy$$

$$= -\sum_{i=1}^d \frac{\partial}{\partial x_i} \left( b_i(t, x) \int_{\mathbb{R}^d} p(t, s, y, x) q(s, y) dy \right)$$

$$+ \frac{1}{2} \sum_{i,j=1}^d \frac{\partial^2}{\partial x_i x_j} \left( a_{ij}(x) \int_{\mathbb{R}^d} p(t, s, y, x) q(s, y) dy \right)$$

$$= \tilde{L}^*(t) \int_{\mathbb{R}^d} p(t, s, y, x) q(s, y) dy = \tilde{L}^*(t) q(t, x),$$

which implies the density $q(s, y)$ satisfies the equation \((2.6)\). The claim that $q(t, y) \to 0$ as $y \to \infty$ follows from \((2.8)\) and the fact that when $|y| \to \infty$, we have $p(t, s, z, y) \to 0$. \hfill \Box

**Corollary 2.3.** If the density function $q(s, x)$ of periodic measure satisfies the equation \((2.6)\), then for any $\tau$-periodic function $f \in C_\infty^p$, we have

$$\int_0^\tau \int_{\mathbb{R}^d} \tilde{L} f(s, x) q(s, x) dx ds = 0.$$  \hspace{1cm}

**Proof.** The main ingredient of proof is to apply integration by parts. Note first

$$\int_0^\tau \int_{\mathbb{R}^d} b_i(s, x) \frac{\partial}{\partial x_i} f(s, x) q(s, x) dx ds = -\int_0^\tau \int_{\mathbb{R}^d} f(s, x) \frac{\partial}{\partial x_i} \left( b_i(s, x) q(s, x) \right) dx ds,$$

$$\int_0^\tau \int_{\mathbb{R}^d} a_{ij}(s, x) \frac{\partial^2}{\partial x_i x_j} f(s, x) q(s, x) dx ds = \int_0^\tau \int_{\mathbb{R}^d} f(s, x) \frac{\partial^2}{\partial x_i x_j} \left( a_{ij}(s, x) q(s, x) \right) dx ds.$$

Here we used the property that $q(s, x)$ vanishes as $x$ goes to $\infty$ when we performed the integration by parts. Applying the periodicity with respect to time $s$ of function $f$ and density function $q$ in the third part, we have

$$\int_0^\tau \int_{\mathbb{R}^d} \frac{\partial}{\partial s} f(s, x) q(s, x) dx ds = -\int_0^\tau \int_{\mathbb{R}^d} f(s, x) \frac{\partial}{\partial s} q(s, x) dx ds.$$

Therefore, by the Fokker-Planck equation on the density function $q(s, x)$, we have

$$\int_0^\tau \int_{\mathbb{R}^d} \tilde{L} f(s, x) q(s, x) dx ds = \int_0^\tau \int_{\mathbb{R}^d} f(s, x) \tilde{L}^* q(s, x) dx ds = 0.$$  \hfill \Box

**Proposition 2.4.** Assume Conditions (1) and (3). Then for any $p \in \mathbb{N}$, there exist strictly positive constants $C_p$ and $\gamma_p$, such that for any $t > 0$ and $x \in \mathbb{R}^d$,

$$\mathbb{E} \left| X_{t+\gamma}^{s,x} \right|^p \leq C_p (1 + |x|^p \exp(-\gamma_p)).$$  \hspace{1cm} (2.9)
Proof. Denote by $X_t := X_{t+s}$ for simplicity. Applying Itô’s formula and Conditions (1), (3), we have the estimate
\[
d (e^{\delta t} |X_t|^p) \leq (\delta - p\beta) e^{\delta t} |X_t|^p dt + p\sigma(X_t)e^{\delta t} |X_t|^{p-1} dW_t + \left( pC\sigma + \left( \frac{p}{2} \right) C_{\sigma}^2 \right) e^{\delta t} |X_t|^{p-2} dt,
\]
where $C_\sigma$ is the bound of function $\sigma$, $C$ and $\beta$ are the constants in the weakly dissipative condition. For convenience, here we denote $C_{\sigma}^2$.

Apply Gronwall’s inequality on (2.10),
\[
\begin{align*}
\mathbb{E} e^{\delta (T \wedge \tau_N)} |X_{T \wedge \tau_N}|^p &\leq |x|^p + (\delta - p\beta) \mathbb{E} \int_0^{T \wedge \tau_N} e^{\delta t} |X_t|^p dt + C_{p,\sigma} \mathbb{E} \int_0^{T \wedge \tau_N} e^{\delta t} |X_t|^{p-2} dt \\
&\leq |x|^p + \frac{2C_{p,\sigma}}{p\delta \varepsilon^2} \mathbb{E}(e^{\delta (T \wedge \tau_N)} - 1) + K_1 \int_0^{T \wedge \tau_N} e^{\delta t} |X_t|^p dt,
\end{align*}
\]
where $K_1 = \delta - p\beta + \frac{(p-2)C_{p,\sigma}}{p\varepsilon^2}, \varepsilon < \left( \frac{p^2 \beta}{(p-2)C_{p,\sigma}} \right)^{\frac{p-2}{p}}$ is chosen such that $K_1 - \delta < 0$. The choice of the constant $\delta$ guarantees $K_1 > 0$.

Then we let $N$ go to $\infty$ to obtain
\[
\mathbb{E} e^{\delta T} |X_T|^p \leq |x|^p + \frac{2C_{p,\sigma}}{p\delta \varepsilon^2} (e^{\delta T} - 1) + K_1 \int_0^T e^{\delta t} \mathbb{E} |X_t|^p dt. 
\tag{2.10}
\]

Apply Gronwall’s inequality on (2.10),
\[
\mathbb{E} e^{\delta T} |X_T|^p \leq e^{K_1 T} \left( |x|^p - \frac{2C_{p,\sigma}}{p\delta \varepsilon^2} (e^{\delta T} - 1) \right) + \frac{2K_1 C_{p,\sigma}}{p(\delta - K_1)\varepsilon^2} (e^{\delta T} - e^{K_1 T}). 
\tag{2.11}
\]

Then (2.9) follows easily. \hfill \Box

Proposition 2.5. Assume Conditions (1) and (3). Then for any $p \in \mathbb{N}$, $\frac{1}{T} \int_0^T \int_{\mathbb{R}^d} |x|^p q(s, x) dx ds \leq C_p$, where $C_p$ is determined from Proposition 2.4.

Proof. For the density function of transition kernel $P(t+s, s, x, \cdot)$, there exists a constant $C$ such that $|p(t+s, s, x, y)| \leq C$, for any $t \geq 1$. Then by dominated convergent theorem and Theorem 3.7 in [13], for any compact set $K \subset \mathbb{R}^d$, we have
\[
\frac{1}{T} \int_0^T \lim_{n \to \infty} \mathbb{E} \left( |X_{n(t+s)}^x|^p \right) 1_K(X_{n(t+s)}^x) ds = \frac{1}{T} \int_0^T \lim_{n \to \infty} \int_K |y|^p p(n\tau + s, s, x, y) dy ds \\
= \frac{1}{T} \int_0^T \int_K |y|^p \lim_{n \to \infty} p(n\tau + s, s, x, y) dy ds = \frac{1}{T} \int_0^T \int_K |y|^p q(s, y) dy ds.
\]

Thus the average of periodic measure possesses finite moments of any order on any compact set $K$ from the estimates in Proposition 2.4. Note the bound can be independent of $K$. The result follows from taking limit $K \uparrow \mathbb{R}^d$ and Fatou’s Lemma. \hfill \Box
Consider the sequence \( \{X_{t_n}\}_{n \in \mathbb{N}} \) with \( t_n = n\tau \). We prove

**Proposition 2.6.** Assume Conditions (1) and (3), then there exists a constant \( r > 1 \) and a ball \( B(0,R) \), such that,

\[
\sup_{x \in B(0,R)} \mathbb{E} \left[ r \left| X_{t_{n+1}} \right|^2 - \left| X_{t_n} \right|^2 \left| X_{t_n} = x \right. \right] < 0.
\]

**Proof.** Apply the result (2.11) in Proposition 2.4 with \( p = 2 \), we have

\[
\mathbb{E} \left[ r \left| X_{t_{n+1}} \right|^2 - \left| X_{t_n} \right|^2 \left| X_{t_n} = x \right. \right] \leq \left( |x|^2 e^{(K_1-\delta)\tau} + C_{\beta,\sigma} (1 - e^{(K_1-\delta)\tau}) \right) r - |x|^2.
\]

In order to make the right hand side of the above negative, we need \( (1 - re^{(K_1-\delta)\tau}) |x|^2 > rC_{\beta,\sigma}(1 - e^{(K_1-\delta)\tau}) \). As \( K_1 - \delta < 0 \), there always exists a constant \( r \) to ensure \( 1 - re^{(K_1-\delta)\tau} > 0 \) for the given period \( \tau \). Then the ball \( B(0,R) \) is determined by taking \( R > \sqrt{\frac{rC_{\beta,\sigma}(1 - e^{(K_1-\delta)\tau})}{1 - re^{(K_1-\delta)\tau}}} \).

Let the function \( \phi \in C^\infty_p \) and \( U(t+s,s,x) = \mathbb{E}\phi(X_{t+s}^{t,s}) \). Then \( U \) satisfies the PDE (2.4). Considering the spatial differentiation of the solution with respect to \( x \), Kunita showed in [22] that the function \( U(t+s,s,x) \) satisfies that for any order \( n \in \mathbb{N} \), there exists an integer \( r_n \in \mathbb{N} \) such that for any \( T > 0 \), \( C_n(t) > 0 \),

\[
|D^n U(t+s,s,x)| \leq C_n(T)(1 + |x|^r_n), \forall t < T.
\]

From Proposition 2.3, the average of periodic measure possesses finite moments of any order. Together with (2.12), we have that the initial condition \( \phi \) and \( D^n U(t+s,s,x) \) belong to \( L^2(\mathbb{R}^{d+1}, \tilde{\rho}) \).

Note that the function \( \frac{1}{T} \int_0^T U(t+s,s,x)ds \) has the same spatial derivatives as \( \frac{1}{T} \int_0^T U(t+s,s,x)ds - \frac{1}{T} \int \phi(x)q(s,x)dxds \). Without loss of generality, in the following sections, we assume that

\[
\int_0^T \int \phi(x)q(s,x)dxds = 0.
\]

Note when \( \tilde{\phi}(x) = \phi(x) \), we have \( \int \tilde{\phi}(x) d\tilde{\rho}(x) = \frac{1}{T} \int_0^T \int \phi(x)q(s,x)dxds \), where \( \tilde{\rho} \) is the average of lifted periodic measure, which is the invariant measure of the lifted Markov semigroup. It is easy to know that \( \tilde{\rho}(dx) = q(s,x)dxds \).

For simplicity, in the following sections, we may often write \( \tilde{U}(t) \) or \( \tilde{U}(t+s,s) \) to represent the function \( U(t+s,s,x) \). We also often write \( b \) to represent \( b(s,x) \) as we have the uniform conditions for the function and any order of its derivatives in Condition (1). The operators \( \partial_i, \partial_{ij}, \nabla \) and \( D^k \) on function \( U(t+s,s,x) \) always refer to derivatives with respect to spatial coordinates. The derivatives with respect to initial time will stay as \( \partial_t \).

3 Exponential decay of initial distribution and spatial derivatives

3.1 Estimates on the average of \( U(t+s,s) \) on a ball

We always assume (2.13) in this section unless otherwise stated.
Lemma 3.1. Assume Conditions (1), (2) and (3). Then for any ball $B$, there exist strictly positive constants $C$ and $\lambda$ such that for any $t > 0$ and any $x \in B$, $U$ defined with $\phi$ satisfying (2.13) has the following estimate: 
\[
\frac{1}{\tau} \int_0^\tau |D^r U(t + s, s, x)|^2 q(s, x)dxds \leq C \exp(-\gamma t). 
\] (3.1)

Proof. First we apply mathematical induction to obtain that for any $p \in \mathbb{N}^+$, there exist constants $C_p > 0$, $\gamma_p > 0$ such that
\[
\frac{1}{\tau} \int_0^\tau \int_{\mathbb{R}^d} |D^p U(t + s, s, x)|^2 q(s, x)dxds \leq C_p \exp(-\gamma_p t). 
\] (3.1)

We start to prove the basis step, when $p = 1$. Consider the Markov chain $\{X_{t_n}\}_{n \in \mathbb{N}}$ with $t_n = n\tau$. In [13], it was proved that the transition kernel $P(s, s + k\tau, x, \Gamma)$ is irreducible. With the result of Proposition 2.6, one can find some compact set $K$ and a constant $\beta > 0$ such that for any $x \in K$, we have
\[
E \left( X_{t_n}^0, x \right)^2 - \frac{1}{r} \lambda x^2 < 0, 
\]
where $1/r < 1$. Now we take $V(x) = x^2$ as the norm-like function and from Proposition 2.4, we obtain that the norm-like function $V(x) = x^2$ is finite on the compact set $K$. Combining the above results, we have that
\[
(P(t_1, 0)V)(x) = E X_{t_1}^2 \leq \frac{1}{r} x^2 + \beta = \frac{1}{r} V(x) + \beta, 
\]
where $\beta$ is a positive number. Thus the condition of Theorem 3.7 in [13] is satisfied. So the Markov chain $\{X_{t_n}\}_{n \in \mathbb{N}}$ is geometrically ergodic. That is for those $\phi \in C^\infty_p$ with the assumption (2.13), there exist strictly positive constants $C$ and $\lambda$ such that for any $n$, 
\[
\frac{1}{\tau} \int_0^\tau \int_{\mathbb{R}^d} |E\phi(X_{t_n+s}^x)| q(s, x)dxds \leq Ce^{-\beta t_n}. 
\] (3.2)

As function $\phi$ has at most polynomial growth at infinity, we have $|\phi(x)| \leq C |x|^N$ for some integer $N \in \mathbb{N}$. By Proposition 2.4 there exist $C_0 > 0$, $\gamma > 0$ such that
\[
|U(t + s, s, x)| \leq C_0 (1 + |x|^N \exp(-\gamma t)). 
\] (3.3)

Applying (3.3) and (3.2), together with Proposition 2.5 we have that for any $n$,
\[
\frac{1}{\tau} \int_0^\tau \int_{\mathbb{R}^d} |U(t_n + s, s, x)|^2 q(s, x)dxds \leq C_0 \int_0^\tau \int_{\mathbb{R}^d} |U(t_n + s, s, x)|(1 + |x|^N \exp(-\gamma t_n))q(s, x)dxds 
\]
\[
\leq C_0 C \exp(-\lambda t_n) + \frac{C^2_0 \exp(-\gamma t_n)}{\tau} \int_0^\tau \int_{\mathbb{R}^d} |x|^N q(s, x)dxds 
\]
\[
\leq C_1 \exp(-\lambda_1 t_n). 
\] (3.4)

In the following, we prove that the function $t \to \frac{1}{\tau} \int_{\mathbb{R}^d} |U(t + s, s, x)|^2 q(s, x)dxds$ is monotonic. For this, note
\[
\frac{d}{dt} |U(t + s, s, x)|^2 = \mathcal{L} |U(t + s, s, x)|^2 - a_{ij} \partial_j U(t + s, s, x) \partial_i U(t + s, s, x), 
\]
and
\[
\frac{1}{\tau} \int_0^\tau \int_{\mathbb{R}^d} \mathcal{L} |U(t + s, s, x)|^2 q(s, x)dxds = \frac{1}{\tau} \int_0^\tau \int_{\mathbb{R}^d} |U(t + s, s, x)|^2 \mathcal{L}^* q(s, x)dxds = 0. 
\]
It turns out from the elliptic condition (2) that
\[
\frac{d}{dt} \left( \frac{1}{\tau} \int_0^T \int_{\mathbb{R}^d} |\mathcal{U}(t + s, x)|^2 q(s, x) dx ds \right) \\
= \frac{1}{\tau} \int_0^T \int_{\mathbb{R}^d} \mathcal{L} |\mathcal{U}(t + s, x)|^2 q(s, x) dx ds \\
- \frac{1}{\tau} \int_0^T \int_{\mathbb{R}^d} a_{ij}(x) \partial_i \mathcal{U}(t + s, x) \partial_j \mathcal{U}(t + s, x) q(s, x) dx ds \\
\leq -\frac{\alpha}{\tau} \int_0^T \int_{\mathbb{R}^d} |\nabla \mathcal{U}(t + s, x)|^2 q(s, x) dx ds \leq 0.
\]
This implies that \(\frac{1}{\tau} \int_0^T \int_{\mathbb{R}^d} |\mathcal{U}(t + s, x)|^2 q(s, x) dx ds\) is decreasing in \(t\). Thus by (3.4), we have that for any \(t\),
\[
\frac{1}{\tau} \int_0^T \int_{\mathbb{R}^d} |\mathcal{U}(t + s, x)|^2 q(s, x) dx ds \leq C_2 \exp(-\lambda_1 t).
\] (3.5)
The above shows that the exponential contraction of \(\mathcal{U}(t + s, x)\) under the average of periodic measure holds for any \(t\).

On the other hand, by Condition (2) we have
\[
\frac{d}{dt} |\mathcal{U}(t + s, x)|^2 - \mathcal{L} |\mathcal{U}(t + s, x)|^2 \leq -\alpha |\nabla \mathcal{U}(t + s, x)|^2.
\] (3.6)
Multiplying the above inequality with \(e^{\delta t}\), and integrating both sides with respect to the average periodic measure \(\bar{\rho}\) and time \(t\), together with Corollary 2.3, we obtain for arbitrary \(T > 0\),
\[
\int_0^T e^{\delta t} \int_{\mathbb{S}} \frac{d}{dt} |\mathcal{U}(t + s, x)|^2 d\bar{\rho} dt + \int_0^T \alpha e^{\delta t} \int_{\mathbb{S}} |\nabla \mathcal{U}(t + s, x)|^2 d\bar{\rho} dt \leq 0.
\] (3.7)
Here \(\mathbb{S} = [0, \tau] \times \mathbb{R}^d\). Integration by parts on the first term of (3.7) gives us
\[
\int_0^T e^{\delta t} \int_{\mathbb{S}} \frac{d}{dt} |\mathcal{U}(t + s, x)|^2 d\bar{\rho} dt = e^{\delta T} \int_{\mathbb{S}} |\mathcal{U}(T + s, x)|^2 d\bar{\rho} - \int_{\mathbb{S}} |\mathcal{U}(s, x)|^2 d\bar{\rho} - \delta \int_0^T e^{\delta t} \int_{\mathbb{S}} |\mathcal{U}(t + s, x)|^2 d\bar{\rho} dt, \]
where we have the initial condition that \(\mathcal{U}(s, x) = \bar{\phi}(\bar{x})\). By Proposition 2.4 and \(\phi \in C^\infty_p\) of the function \(\phi\), we have a constant \(C_3 > 0\) such that \(\int_{\mathbb{S}} \left| \bar{\phi}(\bar{x}) \right|^2 d\bar{\rho} < C_3\). Consider (3.5) and take \(\delta < \lambda_1\),
\[
\delta \int_0^T e^{\delta t} \int_{\mathbb{S}} |\mathcal{U}(t + s, x)|^2 d\bar{\rho} dt \leq \delta \int_0^T e^{\delta t} C_2 e^{(-\lambda_1 t) dt} = \frac{C_2 \delta}{\lambda_1 - \delta} (1 - e^{(\delta - \lambda_1) T}) \leq C_4
\]
for a constant \(C_4 > 0\). Applying these results to (3.7), we obtain that for any \(T\) and any \(s \in [0, \tau]\),
\[
\int_0^T e^{\delta t} \int_{\mathbb{S}} |\nabla \mathcal{U}(t + s, x)|^2 d\bar{\rho} dt \leq \frac{1}{\alpha} \left( \int_{\mathbb{S}} |\mathcal{U}(s, x)|^2 d\bar{\rho} + \delta \int_0^T e^{\delta t} \int_{\mathbb{S}} |\mathcal{U}(t + s, x)|^2 d\bar{\rho} dt \right) \leq C_5, \] (3.8)
where $C_5 = C_3 + C_4$. Now let’s consider $|\nabla U(t + s, s, x)|^2$ and note that

$$\frac{d}{dt} |\nabla U(t + s, s, x)|^2 - \tilde{L} |\nabla U(t + s, s, x)|^2 = -a_{ij}(\partial_i U(t + s, s, x))(\partial_j U(t + s, s, x)) + 2(\partial_i b_i)(\partial_j U(t + s, s, x)))(\partial_j U(t + s, s, x)) + (\partial_i a_{ij})(\partial_j U(t + s, s, x))(\partial_j U(t + s, s, x)).$$

Applying Young’s inequality with $\varepsilon$, we have

$$(\partial_i a_{ij})(\partial_j U(t + s, s, x))(\partial_j U(t + s, s, x)) \leq \frac{\varepsilon}{2}(\partial_j U(t + s, s, x))^2 + \frac{(\partial_i a_{ij})(\partial_j U(t + s, s, x))^2}{2\varepsilon}.$$  \hfill (3.8)

From Conditions (1) and (2), we can choose $\varepsilon$ small enough such that $-\alpha + \frac{\varepsilon}{2} < 0$. It turns out that there exist strictly positive constants $C_6$ and $C_7$ such that

$$\frac{d}{dt} |\nabla U(t + s, s, x)|^2 - \tilde{L} |\nabla U(t + s, s, x)|^2 \leq -C_6 |D^2 U(t + s, s, x)|^2 + C_7 |\nabla U(t + s, s, x)|^2. \quad \hfill (3.9)$$

We choose $\gamma < \delta$ and multiply $e^{\gamma t}$ on both sides of the above inequality. It follows that

$$\int_0^T e^{\gamma t} \int_S \frac{d}{dt} |\nabla U(t + s, s, x)|^2 d\tilde{\rho} dt - \int_0^T e^{\gamma t} \int_S \tilde{L} |\nabla U(t + s, s, x)|^2 d\tilde{\rho} dt \leq -C_6 \int_0^T e^{\gamma t} \int_S |D^2 U(t + s, s, x)|^2 d\tilde{\rho} dt + C_7 \int_0^T e^{\gamma t} \int_S |\nabla U(t + s, s, x)|^2 d\tilde{\rho} dt.$$

Following Corollary 2.3, we see that $\int_0^T e^{\gamma t} \int_S \tilde{L} |\nabla U(t + s, s, x)|^2 d\tilde{\rho} dt = 0$. Thus

$$\int_0^T e^{\gamma t} \int_S \frac{d}{dt} |\nabla U(t + s, s, x)|^2 d\tilde{\rho} dt \leq C_7 \int_0^T e^{\gamma t} \int_S |\nabla U(t + s, s, x)|^2 d\tilde{\rho} dt.$$

Now by integration by parts, we note that

$$\int_0^T e^{\gamma t} \int_S \frac{d}{dt} |\nabla U(t + s, s, x)|^2 d\tilde{\rho} dt = e^{\gamma T} \int_S |\nabla U(T + s, s, x)|^2 d\tilde{\rho} - \int_0^T \int_S |\nabla U(s, s, x)|^2 d\tilde{\rho} - \gamma \int_0^T e^{\gamma t} \int_S |\nabla U(t + s, s, x)|^2 d\tilde{\rho} dt.$$

Then apply (3.8) with $\gamma$ small enough and the boundedness of $\int_S |\nabla \tilde{\phi}|^2 d\tilde{\rho}$ to have

$$e^{\gamma T} \int_S |\nabla U(T + s, s, x)|^2 d\tilde{\rho} \leq \int_S |\nabla \tilde{\phi}|^2 d\tilde{\rho} + (\gamma + C_7) \int_0^T e^{\gamma t} \int_S |\nabla U(t + s, s, x)|^2 d\tilde{\rho} dt \leq C_8.$$

Thus we obtained (3.1) for the case when $p = 1$. Now we continue to prove the induction step in the following content. Assume that for any $k \leq m$, there exist strictly positive constants $C_k$ and $\gamma_k$ such that for any $t > 0$,

$$\frac{1}{\tau} \int_0^T \int_{\mathbb{R}^d} |D^k U(t + s, s, x)|^2 q(s, x) dx ds \leq C_k \exp(-\gamma_k t).$$
Here we need to compare the expansion of the operators $\frac{d}{dt}$ and $\tilde{L}$ in the following:

$$|D^m U(t + s, s, x)|^2 = \sum_{|J| = m} (\partial_J U(t + s, s, x))^2,$$

where $J$ is the multi-index with length $|J| = m$. The multi-indices $J_a$ and $J_b$ are introduced for the following identity,

$$\frac{d}{dt} |D^m U(t + s, s, x)|^2 - \tilde{L} |D^m U(t + s, s, x)|^2 = -a_{ij} (\partial_j \partial_i U(t + s, s, x)) (\partial_j \partial_i U(t + s, s, x)) + \sum_{|J_a| + |J_b| \leq 2m + 1} \Phi_{J_a, J_b} \partial_J U(t + s, s, x) \partial_J U(t + s, s, x).$$

Here the notation $\Phi_{J_a, J_b}$ contains all the combinations of spatial derivatives on the functions $a$ and $b$ with respect multi-indices $J_a$ and $J_b$ under some specified $J$. It is obvious the length of $J_a$ and $J_b$ will not exceed $m + 1$. The boundedness of each elements in $\Phi_{J_a, J_b}$ comes from Condition (1). Therefore we will always have the following result by Young’s inequality,

$$\frac{d}{dt} |D^m U(t + s, s, x)|^2 - \tilde{L} |D^m U(t + s, s, x)|^2 \leq -C_1^m |D^{m+1} U(t + s, s, x)|^2 + C_2^m \sum_{k \leq m} |D^k U(t + s, s, x)|^2.$$

Then we choose a strictly positive constant $\delta_{m+1}$ small enough to proceed as in (3.8). Multiplying $e^{\delta_{m+1}t}$ on both sides and integrating with respect to $\tilde{\rho}$, we will have

$$\int_0^\infty e^{\delta_{m+1}t} \left( \int_{\mathbb{R}^d} |D^{m+1} U(t + s, s)|^2 \, d\tilde{\rho} \right) dt < \infty.$$

Consider a higher order

$$\frac{d}{dt} |D^{m+1} U(t + s, s, x)|^2 - \tilde{L} |D^{m+1} U(t + s, s, x)|^2 \leq -C_1^{m+1} |D^{m+2} U(t + s, s, x)|^2 + C_2^{m+1} \sum_{k \leq m+1} |D^k U(t + s, s, x)|^2.$$

By choosing $\gamma_{m+1} < \delta_{m+1}$ and following the same procedure as above, we have (3.1) for the case when $p = m + 1$. By induction principle, we proved the above result holds for any order of spatial derivatives of $U(t + s, s, x)$.

By (2.7), we can conclude that $q(s, x) > 0$ as $p(s+\tau, s, y, x) > 0$ for any $s \in \mathbb{R}$ and $x, y \in \mathbb{R}^d$. We can also prove the continuity of $q(t, x)$ from the continuity of $p(t + \tau, t, y, x)$ in $x$. Thus the density function $q(s, x)$ is strictly positive continuous function on any ball $B = B(0, R)$. It turns out that there exists $C > 0$ such that

$$\frac{1}{\tau} \int_0^\tau \|\partial_J U(t + s, s, x)\|^2_{L^2(B)} \, ds \leq C \int_0^\tau \int_{\mathbb{R}^d} |\partial_J U(t + s, s, x)|^2 q(s, x) dx \, ds.$$

By the Sobolev embedding $W^{k, 2}(B) \hookrightarrow C(B)$ for $k > \frac{d}{2}$ (II), we have that

$$\frac{1}{\tau} \int_0^\tau |U(t + s, s, x)| \, ds \leq \frac{1}{\tau} \int_0^\tau \int_B |D^k U(t + s, s, x)|^2 \, dx \, ds \leq C_k \exp(-\lambda_k t),$$

for any $x \in B$, where $k > \frac{d}{2}$. The proof is completed. □
3.2 Estimates on the average of $\mathcal{U}(t+s,s)$ in $L^2(\pi_r)$

In Section 3.1, we obtained the exponential contraction of $\frac{1}{\tau} \int_0^\tau |\mathcal{U}(t+s,s,x)| ds$ in any ball $B$ when we assumed $\int_0^\tau \int_{\mathbb{R}^d} \phi(x)q(s,x)dxds = 0$. To consider the behaviour outside of the ball $B$, we need to introduce the weight $\pi_r(s,x)$ with some integer $r$ determined later,

$$\pi_r(s,x) = 1/(2 + |x|^2 + \cos(\frac{2\pi s}{\tau}))^r.$$ 

We consider its gradient and partial derivatives with respect to time $s$,

$$\nabla \pi_r(s,x) = -\frac{2rx}{2 + |x|^2 + \cos(\frac{2\pi s}{\tau})}\pi_r(s,x), \quad \partial_s \pi_r(s,x) = \frac{2\pi \sin(\frac{2\pi s}{\tau})}{2 + |x|^2 + \cos(\frac{2\pi s}{\tau})}\pi_r(s,x).$$

In general, it is easy to see that for any multi-index $J$ and any integer $r$, there exist smooth functions $\psi_{J,r}(s,x)$ and $\psi_{s,r}(s,x)$ such that,

$$\partial_J \pi_r(s,x) = \psi_{J,r}(s,x) \pi_r(s,x), \quad \partial_s \pi_r(s,x) = \psi_{s,r}(s,x) \pi_r(s,x),$$

where $\psi_{J,r}(s,x) \to 0$ and $\psi_{s,r}(s,x) \to 0$ when $|x| \to +\infty$.

**Lemma 3.2.** Assume Conditions (1), (2) and (3), there exist strictly positive constants $C$ and $\lambda$ such that for any $t > 0$, we have

$$\frac{1}{\tau} \int_0^\tau \int_{\mathbb{R}^d} |\mathcal{U}(t+s,s,x)|^2 \pi_r(s,x)dxds \leq C \exp(-\lambda t).$$

**Proof.** Recall (2.12) to lead that for any integer $n \geq 0$, it is possible to choose an integer $r_n$ such that, for any $0 \leq m \leq n$, $t \geq 0$, we have $|D^m \mathcal{U}(t+s,s,x)| \pi_{r_n}(s,x) \in L^2(\mathbb{R}^d)$. We denote the multi-index $I$ for the derivative $\partial$ with the length $|I|$. Consider the integer $M_I$ defined by $|I| = [M_I - d/2]$, and the property of the weight $\pi_r$, we have that there exists an integer $r_0$ such that for any $t > 0$, any $r \geq r_0$ and any $m \leq M_I$, $D^m (\mathcal{U}(t+s,s,x)\pi_r(s,x)) \in L^2(\mathbb{R}^d)$. It is easy to see the periodicity of the function $\mathcal{U}(t+s,s,x)\pi_r(s,x)$ with respect to the initial time $s$. Note any order of its spatial derivatives are also $\tau$-periodic in $s$, so by integration by parts formula and periodicity,

$$\int_0^\tau \int_{\mathbb{R}^d} \frac{d}{dt} |\mathcal{U}(t+s,s)|^2 \pi_r dxds$$

$$= -\int_0^\tau \int_{\mathbb{R}^d} (\partial_s b_i) |\mathcal{U}(t+s,s)|^2 \pi_r dxds - \int_0^\tau \int_{\mathbb{R}^d} b_i |\mathcal{U}(t+s,s)|^2 (\partial_i \pi_r)dxds - \int_0^\tau \int_{\mathbb{R}^d} |\mathcal{U}(t+s,s)|^2 (\partial_s \pi_r) dxds$$

$$- \int_0^\tau \int_{\mathbb{R}^d} (\partial_i a_{ij}) \mathcal{U}(t+s,s)(\partial_j \mathcal{U}(t+s,s))\pi_r dxds - \int_0^\tau \int_{\mathbb{R}^d} a_{ij}(\partial_i \mathcal{U}(t+s,s))(\partial_j \mathcal{U}(t+s,s))\pi_r dxds$$

$$- \int_0^\tau \int_{\mathbb{R}^d} a_{ij}\mathcal{U}(t+s,s)(\partial_j \mathcal{U}(t+s,s))(\partial_i \pi_r) dxds.$$

By Condition (2) and the property of the weight $\pi(s,x)$, we have that

$$\int_0^\tau \int_{\mathbb{R}^d} \frac{d}{dt} |\mathcal{U}(t+s,s)|^2 \pi_r dxds$$




The result follows then from the Gronwall’s inequality.

\[ \text{where } C \]

and

\[ \text{there exists an integer } k \]

Now choosing the ball \( B \) the following result from (3.10),

\[ \text{that } \Phi \]

\( a,b \)

where \( \phi \) is a bounded function depending on functions \( a, b \) and their derivatives, \( \Phi \) is a function which could depend on functions \( \psi_{I,r} \). It is easy to prove that \( \Phi_{a,b} \) is independent of \( r \). We also know that \( \Phi \) tends to 0 when \( |x| \) goes to \( \infty \). Therefore, we choose \( r \geq r_0 \) large enough to obtain,

\[ \lim_{|x| \to \infty} \left( \Phi_{a,b}(s,x) + \Phi_\psi(s,x) + \frac{2r \cdot x \cdot b(s,x)}{2 + |x|^2 + \cos\left(\frac{2\pi s}{r}\right)} \right) < 0. \quad (3.10) \]

Now choosing the ball \( B = B(0, R) \) with \( R \) being large enough, which depends on the integer \( r \), we have the following result from (3.10),

\[ \int_0^T \int_{\mathbb{R}^d} \left( \Phi_{a,b}(s,x) + \Phi_\psi(s,x) + \frac{2r \cdot x \cdot b(s,x)}{2 + |x|^2 + \cos\left(\frac{2\pi s}{r}\right)} \right) |U(t+s,s)|^2 \pi_r dxds \]

\[ \leq C_1 \int_0^T \int_B |U(t+s,s)|^2 \pi_r dxds - C_2 \int_0^T \int_{B^c} |U(t+s,s)|^2 \pi_r dxds \]

\[ \leq (C_1 + C_2) \int_0^T \int_B |U(t+s,s)|^2 \pi_r dxds - C_2 \int_0^T \int_{\mathbb{R}^d} |U(t+s,s)|^2 \pi_r dxds, \]

where \( C_1, C_2 > 0 \) are constants. Therefore, by Lemma 3.1

\[ \frac{d}{dt} \int_0^T \int_{\mathbb{R}^d} |U(t+s,s)|^2 \pi_r dxds \leq -C_2 \int_0^T \int_{\mathbb{R}^d} |U(t+s,s)|^2 \pi_r dxds + C_3 \exp(-\lambda t). \]

The result follows then from the Gronwall’s inequality. \( \square \)

### 3.3 Exponential decay of the spatial derivatives of the solution

**Theorem 3.3.** Assume Conditions (1), (2) and (3), and \( \phi \in C^{\infty}_{\mathbf{c}} \). Then for any multi-index \( I \), there exists an integer \( k_1 \), strictly positive constants \( \Gamma_I \) and \( \gamma_I \) such that \( \frac{1}{\Gamma_I} \int_0^\tau |\partial_I U(t+s,s,x)| ds \leq \Gamma_I (1 + |x|^{k_1}) \exp(-\gamma_I t) \).

**Proof.** The process of the proof is similar to Lemma 3.1. We first apply induction method on each order of spatial derivatives of \( U(t+s,s,x) \). It guaranteed first the exponential contraction in any ball \( B(0,R) \). Now we consider the behaviour outside of the ball to have,

\[ \int_{B^c} \hat{L} |U(t+s,s)|^2 \pi_r \hat{d} \hat{x} \]

\[ = \int_{B^c} \left( \Phi_{a,b} + \Phi_\psi + \frac{2r \cdot x \cdot b}{2 + |x|^2 + \cos\left(\frac{2\pi s}{r}\right)} \right) |U(t+s,s)|^2 \pi_r d\hat{x} < 0, \]
if we choose the ball large enough. Thus we have some positive $C_0$ and $\lambda_0$ such that
\[ \int_S \hat{\mathcal{L}} |\mathcal{U}(t + s, s)|^2 \pi_r d\hat{x} < \int_B \hat{\mathcal{L}} |\mathcal{U}(t + s, s)|^2 \pi_r d\hat{x} \leq C_0 \exp(-\lambda_0 t). \] (3.11)

On the other hand, we integrate with respect to $d\hat{x}$ with weight $\pi_r$, multiply $e^{\delta t}$ and integrate with respect to $t$ from 0 to $T$ to have
\[ e^{\delta T} \int_S |\hat{\mathcal{L}} |\mathcal{U}(t + s, s)|^2 \pi_r d\hat{x} + C \int_0^T e^{\delta t} \int_S |\nabla \mathcal{U}(t + s, s)|^2 \pi_r d\hat{x} dt \]
\[ \leq \int_S |\hat{\phi}(\hat{x})|^2 \pi_r d\hat{x} + \delta \int_0^T e^{\delta t} \int_S |\mathcal{U}(t + s, s)|^2 \pi_r d\hat{x} dt + \int_0^T e^{\delta t} \int_S \hat{\mathcal{L}} |\mathcal{U}(t + s, s)|^2 \pi_r d\hat{x} dt. \]

By the estimates (3.5) and (3.11), we can choose constant $\delta$ small enough to obtain
\[ \int_0^T e^{\delta t} \int_S |\nabla \mathcal{U}(t + s, s)|^2 \pi_r d\hat{x} dt \leq C. \]

Similarly we consider (3.9) to have
\[ e^{\gamma T} \int_S |\nabla \mathcal{U}(T + s, s)|^2 \pi_r d\hat{x} + C \int_0^T e^{\gamma t} \int_S |D^2 \mathcal{U}(t + s, s)|^2 \pi_r d\hat{x} dt \]
\[ \leq \int_S |\nabla \hat{\phi}(\hat{x})|^2 \pi_r d\hat{x} + (\gamma + C_2) \int_0^T e^{\gamma t} \int_S |\nabla \mathcal{U}(t + s, s)|^2 \pi_r d\hat{x} dt + \int_0^T e^{\gamma t} \int_S \hat{\mathcal{L}} |\nabla \mathcal{U}(t + s, s)|^2 \pi_r d\hat{x} dt. \]

which gives us the conclusion that $\int_S |\nabla \mathcal{U}(t + s, s)|^2 \pi_r d\hat{x} \leq Ce^{-\gamma t}$. It is easy to repeat the process for any $m \in \mathbb{N}$ with positive constants $C_m$ and $\gamma_m$ to obtain $\int_S |D^m \mathcal{U}(t + s, s)|^2 \pi_r d\hat{x} \leq C_m e^{-\gamma_m t}$. Then we proved the conclusion of the theorem by the weighted Sobolev embedding Theorem with $\pi_r(s, x) d\hat{x}$ instead of the the density function of average periodic measure $q(s, x) d\hat{x}$.

The following remark applies to general case without assumption (2.13).

**Remark 3.4.** The proof of the previous theorem also gives us the result that there exist some integer $l \in \mathbb{N}$ and constants $\Gamma > 0$, $\gamma > 0$, such that for any $t$ and $x$,
\[ \frac{1}{\tau} \int_0^\tau \mathcal{U}(t + s, s, x) ds - \int_S \hat{\phi}(\hat{x}) d\hat{p}(\hat{x}) \leq \Gamma(1 + |x|^l) \exp(-\gamma t). \] (3.12)

### 4 Ergodicity for discretized semi-flows of Euler-Maruyama scheme

We consider Euler-Maruyama numerical scheme with step size $\Delta t = \frac{\tau}{N} > 0$ for SDE (1.1):
\[ \hat{X}_{-kT+i+1}\Delta t} = \hat{X}_{-kT+i}\Delta t + b(i\Delta t, \hat{X}_{-kT+i+1}\Delta t) \Delta t + \sigma(\hat{X}_{-kT+i+1}\Delta t) \Delta W_i, \] (4.1)
with $\hat{X}_{-kT+i+1}\Delta t} = x$, where $i = 0, 1, 2, \ldots$, $\Delta W_i = W_{-kT+i+1}\Delta t} - W_{-kT+i}\Delta t}$. There are several methods to generate the stochastic increment $\Delta W_i$. But in order to obtain the ergodicity of numerical schemes, in this paper we apply Gaussian distribution in the approximation i.e. $\Delta W = \sqrt{\Delta t} \mathcal{N}(0, 1)$. Denote transition probability
\[ \hat{P}( -kT+i\Delta t}, -kT+i\Delta t}, x, \Gamma) = \mathbb{P} \left( \hat{X}_{-kT+i}\Delta t} \in \Gamma \right). \]
It is easy to see that \( \hat{P}(-k\tau + i\Delta t, -k\tau, x, \Gamma) = \hat{P}(i\Delta t, 0, x, \Gamma) \). One can easily extend the numerical scheme (4.1) to \( \hat{X}_t^{j\Delta t} \), \( i \geq j \), \( j \in \mathbb{Z} \) with \( \hat{X}_t^{j\Delta t} = x \) and its transition probability to \( \hat{P}(i\Delta t, j\Delta t, x, \cdot) \), \( i \geq j \), \( j \in \mathbb{Z} \). The corresponding semigroup \( \hat{T}(i\Delta t, j\Delta t) \), \( i \geq j \), \( j \in \mathbb{Z} \) can be generated from the transition probability in a standard way. 

**Proof.** We first consider the one-dimensional case. Condition (3'), which is the same as Condition (3) for \( p \rightarrow 0 \) and \( d = 1 \). By Condition (1), if the function \( b(t, \cdot) \) is of linear growth \( |b(t, x)| \leq L|x| + C \), where \( L, C > 0 \).

**Remark 4.1.** By Condition (1), if the function \( b(t, 0) \) is bounded for any \( t > 0 \), then \( b \) is of linear growth \( |b(t, X_t)| \leq L|X_t| + C \), where \( L, C > 0 \).

**Remark 4.2.** Under Condition (3), the conclusion in the following proposition still holds for sufficient small step size \( \Delta t < \Delta t_c \) with some \( \Delta t_c > 0 \) that may depend on the growth order of the test function. In order to obtain a uniform \( \Delta t_c \), we consider the following slightly stronger condition. But in the case of one-dimension, Condition (3') is the same as Condition (3). This means that in the case of one-dimension and Condition (3), a uniform \( \Delta t < \Delta t_c \) is obtained with respect to all polynomial growth test functions.

**Condition (3')** For all \( i = 1, 2, \ldots, d \), there exist constants \( \beta_i > 0 \) and \( C_{\beta_i} > 0 \) such that for any \( t \in \mathbb{R}^+ \) and any \( x_i \in \mathbb{R} \) such that \( b_i(t, x_i) \leq -\beta_i |x_i|^2 + C_{\beta_i} \).

**Proposition 4.3.** Assume Conditions (1), (3') and the boundedness of \( b(t, 0) \) for any \( t > 0 \), then for any integer \( p \) and any \( \delta > 0 \), there exist constants \( C_p, \tilde{C}_p, \gamma, \tilde{\gamma} > 0 \) such that for any \( 0 < \Delta t \leq \frac{2}{L^2} \), \( x \in \mathbb{R}^d \) and \( n \in \mathbb{N} \),

\[
\mathbb{E} \left| \hat{X}_{-k\tau + i\Delta t} \right|_p^p \leq C_p \left( 1 + |x|^p \exp(-\gamma\Delta t) \right),
\]

and

\[
\mathbb{E} \left[ \left| \hat{X}_{-k\tau + (i+1)\Delta t} \right|_p^p \right| \hat{F}_i \leq (1 - \tilde{\gamma}\Delta t) \left| \hat{X}_{-k\tau + i\Delta t} \right|_p^p + \tilde{C}_p,
\]

where \( \hat{F}_i = \mathcal{F}_{-k\tau + i\Delta t} \).

**Proof.** We first consider the one-dimensional case. Condition (3'), which is the same as Condition (3) in this case, implies that for any \( x > \sqrt{\frac{C_{\beta}}{\beta}} \), \( x \cdot b(t, x) \leq -\beta |x|^2 + C_{\beta} < 0 \). It then follows that when \( x > \sqrt{\frac{C_{\beta}}{\beta}} \),

\[
(1 - L\Delta t)x - C\Delta t \leq x + b(t, x)\Delta t \leq (1 - \beta\Delta t)x + \frac{C_{\beta}\Delta t}{x} \leq (1 - \beta\Delta t)x + \sqrt{\beta C_{\beta}\Delta t}.
\]

Thus,

\[
|x + b(t, x)\Delta t| \leq \max\{|1 - \beta\Delta t|, |1 - L\Delta t|\} |x| + C_1\Delta t,
\]

where \( C_1 \) is independent of \( \Delta t \). One can obtain the same result for \( x < -\sqrt{\frac{C_{\beta}}{\beta}} \). It is not hard to verify that \( L \geq \beta \). Then, for \( \tilde{\gamma}_1 = \min\{\beta, \delta\} > 0 \), we can see that for \( 0 < \Delta t \leq \frac{2}{L\tilde{\gamma}_1^2} \),

\[
\max\{|1 - \beta\Delta t|, |1 - L\Delta t|\} < 1 - \tilde{\gamma}_1\Delta t.
\]
It then follows from (4.3) and (4.4) that

$$|x + b(t, x)\Delta t| \leq |1 - \hat{\gamma}_1 \Delta t| |x| + C_1 \Delta t. \tag{4.5}$$

Fix any sufficiently small $\hat{\varepsilon} > 0$, choose $\varepsilon_k$ such that $\frac{p - k}{p} \varepsilon_k = \hat{\varepsilon}^k$. Now for any given integer $p$, by Young’s inequality with positive $\hat{\varepsilon}$, by Young’s inequality with positive

$$1 + \frac{1}{p} + \frac{1}{q} = 1,$$

it then follows from (4.5) and (4.6) that

$$|1 - \hat{\gamma}_1 \Delta t|^{p-k} |x|^{p-k} C_1^k \leq |1 - \hat{\gamma}_1 \Delta t|^p |x|^p \left( \frac{p - k}{p} \varepsilon_k^{\frac{p}{k}} \right) + (kC_1)/(|p\varepsilon_k^2|), \tag{4.6}$$

we only consider the cases where $|1 - \hat{\gamma}_2 \Delta t| > 0$, such that

$$1 + 1 = 2.$$
\[ \leq \left( 1 - \hat{\gamma}_3 \Delta t \right)^{\frac{p}{2}} \left| \hat{X}_{-k\tau + i\Delta t} - \hat{X}_{-k\tau + (i+1)\Delta t} \right|^p + C_3 \Delta t \right) I_A \left( \hat{X}_{-k\tau + i\Delta t} - \hat{X}_{-k\tau + (i+1)\Delta t} \right), \]

with some constant \( C_3 \) independent of \( \Delta t \), where \( \hat{\varepsilon} \) is chosen small enough such that \( (1 + \hat{\varepsilon} \Delta t)^{\frac{p}{2}} |1 - \hat{\gamma}_3 \Delta t|^p \leq 1 - \hat{\gamma}_3 \Delta t |< 1 \) for some \( \hat{\gamma}_3 > 0 \). Moreover by linear growth condition of \( b \) and \( |1 + L \Delta t| \leq 1 + L \cdot \frac{3}{2} = 3, \) we can obtain

\[ \mathbb{E} \left[ \left| \hat{X}_{-k\tau + (i+1)\Delta t} \right|^p \right] \leq \left( 3^p \left( \frac{C_3}{\beta} \right)^{\frac{p}{2}} + C \right) I_A \left( \hat{X}_{-k\tau + i\Delta t} \right), \]

where \( L \) is the bound of function \( b \)'s first derivative (or coefficient of global Lipschitz). We combine the above two estimates to obtain

\[ \mathbb{E} \left[ \left| \hat{X}_{-k\tau + (i+1)\Delta t} \right|^p \right] \leq |1 - \hat{\gamma}_3 \Delta t|^{\frac{p}{2}} \left| \hat{X}_{-k\tau + i\Delta t} \right|^p + \hat{C}_p, \]

with \( \hat{C}_p = \max \left\{ C_3 \Delta t, 3^p \left( \frac{C_3}{\beta} \right)^{\frac{p}{2}} + C \right\} \) and \( \hat{\gamma}_3 = \hat{\gamma}_3 \). Therefore,

\[ |1 - \hat{\gamma}_3 \Delta t|^{\frac{p}{2}} \mathbb{E} \left[ \left| \hat{X}_{-k\tau + n\Delta t} \right|^p \right] \]

\[ = |x|^p + \sum_{i=1}^n \left| 1 - \hat{\gamma}_3 \Delta t \right|^{\frac{p}{2}} \mathbb{E} \left[ \left| \hat{X}_{-k\tau + i\Delta t} \right|^p - \left| 1 - \hat{\gamma}_3 \Delta t \right|^{\frac{p}{2}} \mathbb{E} \left[ \left| \hat{X}_{-k\tau + (i-1)\Delta t} \right|^p \right] \right] \]

\[ \leq |x|^p + \hat{C}_p \sum_{i=1}^n \left| 1 - \hat{\gamma}_3 \Delta t \right|^{\frac{p}{2}} \leq |x|^p + C_p, \]

where \( C_p = \frac{\hat{C}_p}{1 - |1 - \hat{\gamma}_3 \Delta t|^{\frac{p}{2}}} \). Finally from \( |1 - \hat{\gamma}_3 \Delta t| < \exp(-\hat{\gamma}_3 \Delta t) \), we have that \( \mathbb{E} \left[ \left| \hat{X}_{-k\tau + n\Delta t} \right|^p \right] \leq |x|^p \exp(-\frac{\hat{\gamma}_3 \Delta t}{2}) + C_p. \)

For the multi-dimensional case, we apply Condition (3') to have the estimates with coefficients \( C_{i,p} \) and \( \gamma_i \) for each \( i = 1, 2, \ldots, d \). Then the final conclusion follows. \( \square \)

**Proposition 4.4.** Assume the conditions in Proposition 3.3 and Condition (2), then the Euler-Maruyama scheme (4.1) is geometrically ergodic for all step-size \( 0 < \Delta t < \frac{2}{k} \), i.e. there exists a periodic measure \( \mu^* : \mathbb{R}^d \rightarrow \mathcal{P}(\mathbb{R}^d) \) such that

\[ \left\| \hat{P}(i\Delta t, -k\tau + i\Delta t, \cdot) - \hat{\mu}(\cdot) \right\|_{TV} \leq C e^{-\delta k\tau}, \quad k \in \mathbb{N}, \]

for some constants \( C, \delta > 0 \).

**Proof.** To check the local Doeblin condition ([24], [27]), we need to prove the transition kernel \( \hat{P}(s + \tau, s, A) = \hat{P}(s + N \Delta t, s, A) = \hat{P}(s + N \Delta t, s, A) = \hat{P}(X_{s+N\Delta t} = A, \hat{X}_{s+N\Delta t} = x) \) possesses a density function \( \hat{p}(t, s, x) \) satisfying \( \inf_{x, y \in K} \hat{p}(s + \tau, s, x, y) > 0 \) for some non-empty set \( K \in \mathcal{B} \) with Lebesgue measure \( \Lambda(K) > 0 \). To apply the result of Theorem 3.5 in [13] to prove the local Doeblin condition of the transition kernel, we only need to prove there is a non-empty compact set \( K \), such that for any \( i = 0, \ldots, N - 1 \) and any non-empty open set \( \Gamma \), we have

\[ \inf_{x \in K} \hat{P}(s + (i+1)\Delta t, s + i\Delta t, x, \Gamma) > 0. \]
Consider the numerical approximation in the time interval \([s + i\Delta t, s + (i + 1)\Delta t]\). For simplicity, we denote by \(s_i = s + i\Delta t, i = 0, \ldots, N - 1\) and \(\tilde{X}_i^x = \tilde{X}_{s_i}^x\). So

\[
\tilde{X}_{i\Delta t}^x - \tilde{X}_0^x = b(s_i, \tilde{X}_i^x)\Delta t + \sigma(\tilde{X}_i^x)\Delta W_t,
\]

where \(\Delta W_t = W_{t\Delta t} - W_0\). Let \(\tilde{X}_{i\Delta t}^x\) be defined by (4.8) conditioned on \(\tilde{X}_0^x = x\). Then the law of \(\tilde{X}_{i\Delta t}^x\) is \(\tilde{P}(s + (i + 1)\Delta t, s + i\Delta t, x, \cdot) = \tilde{P}(s_i + \Delta t, s_i, x, \cdot)\). Note \(b(s_i, x)\) and \(\sigma(x)\) are non-random and given, thus \(\tilde{P}(s_i + \Delta t, s_i, x, \cdot)\) is simply the Gaussian distribution with mean \(x + b(s_i, x)\Delta t \in \mathbb{R}^d\) and covariance matrix \(\sigma \sigma^T(x)\Delta t \in \mathbb{R}^{d \times d}\). The covariance matrix is uniformly non-degenerate, thus for any non-empty open set \(\Gamma \in \mathbb{R}^d\), we have \(\tilde{P}(s_i + \Delta t, s_i, x, \Gamma) > 0\) and the function \(x \mapsto \tilde{P}(s_i + \Delta t, s_i, x, \Gamma)\) is continuous. Thus for any compact set \(K \subset \mathbb{R}^d\), we have (4.7).

By Theorem 3.5 in [13], we obtain the local Doeblin condition of \(\tilde{P}(s + N\Delta t, s, x, \cdot)\). Note condition \(0 < \Delta t < \frac{2}{N}\) implies there exists \(\delta > 0\) such that \(0 < \Delta t < \frac{2}{N + 3}\). So Proposition 4.3 holds and estimate \([4.2]\) implies Lyapunov condition with Lyapunov function \(V(x) = x^2\). Then by Theorem 3.3 in [13], we deduce the ergodicity of the numerical scheme and the convergence to the periodic measure \(\tilde{\rho}^{\Delta t}\).

Similar in the continuous time case, we can lift the discrete semi-flow and periodic measure to \(\{0, 1, 2, \ldots, N\} \times \mathbb{R}^d\) as follows

\[
\tilde{X}_{i\Delta t}^{j\Delta t} = (i, \tilde{X}_{i\Delta t}^x), \quad i \geq j,
\]

\[
\tilde{P}(i, (j, x), \{k\} \times \Gamma) = \delta_{i+j \text{ mod } N}(k)\tilde{P}((i+j)\Delta t, j\Delta t, x, \Gamma), \quad i \geq j,
\]

\[
\tilde{\rho}^{\Delta t}(\{k\} \times \Gamma) = \delta_{i \text{ mod } N}(k)\tilde{\rho}_i^\Delta(\Gamma), \quad i \in \mathbb{Z}.
\]

Then \(\tilde{X}\) is a cocycle and \(\tilde{P}\) is the transition probability of \(\tilde{X}\). Moreover, \(\tilde{\rho}^{\Delta t}\) is the periodic measure of \(\tilde{P}\), i.e.

\[
\sum_{l=0}^{N-1} \int_{\mathbb{R}^d} \tilde{P}(i, (l, x), \{k\} \times \Gamma)\tilde{\rho}^{\Delta t}(\{l\} \times dx) = \tilde{\rho}^{\Delta t}(\{k\} \times \Gamma).
\]

Define \(\tilde{\rho}^{\Delta t} = \frac{1}{N} \sum_{j=0}^{N-1} \tilde{\rho}_j^{\Delta t}\). Then it is easy to see that

\[
\sum_{l=0}^{N-1} \int_{\mathbb{R}^d} \tilde{P}(i, (l, x), \{k\} \times \Gamma)\tilde{\rho}^{\Delta t}(\{l\} \times dx) = \tilde{\rho}^{\Delta t}(\{k\} \times \Gamma),
\]

i.e. \(\tilde{\rho}\) is the invariant measure of \(\tilde{P}(i), i \in \mathbb{N}\). Moreover, for any measurable function \(\phi : \mathbb{R}^d \to \mathbb{R}\),

\[
\sum_{k=0}^{N-1} \int_{\mathbb{R}^d} \phi(x)\tilde{\rho}^{\Delta t}(\{k\} \times dx) = \sum_{k=0}^{N-1} \int_{\mathbb{R}^d} \phi(x)\frac{1}{N} \sum_{j=0}^{N-1} \delta_j(\tilde{\rho}_j^{\Delta t})(dx) = \frac{1}{N} \sum_{k=0}^{N-1} \int_{\mathbb{R}^d} \phi(x)\tilde{\rho}_k^{\Delta t}(dx) = \sum_{k=0}^{N-1} \int_{\mathbb{R}^d} \phi(x)\tilde{\rho}_k^{\Delta t}(dx),
\]

where \(\tilde{\rho}^{\Delta t} = \frac{1}{N} \sum_{j=0}^{N-1} \tilde{\rho}_j^{\Delta t}\).
5 Error estimate for the approximation to periodic measures

In autonomous systems, there are some established results in the ergodicity of numerical schemes (Mattingly, Stuart and Higham [23]; Grorud and Talay [16]; Talay [28, 29]). But in our non-autonomous model, due to the lacking of weakly mixing property, those approaches may not give immediately the error between \( \int_0^\tau \phi(x) \tilde{\rho}(dx) \) and \( \int_0^\tau \phi(x) \tilde{\rho}^\Delta(dx) \) over \([0, \tau]\). We develop the following approach using integration with respect to initial time to obtain the error estimate of invariant measure. To approximate the average of periodic measure, we need to consider the long time behaviour of the SDE (1.1) by pullback of the initial time to \( s = 0 \).

Assume conditions in Proposition 4.4. Then for any step size \( \Delta t < \frac{\tau}{N} \), we have

\[
\left| \int_{\mathbb{R}^d} \phi(x) \tilde{\rho}(dx) - \int_{\mathbb{R}^d} \phi(x) \tilde{\rho}^N(dx) \right| = O(\Delta t).
\]

Theorem 5.1. Assume conditions in Proposition 4.4. Then for any step size \( \Delta t = \tau/N, N \in \mathbb{N} \), satisfying \( \Delta t < \frac{\tau}{2} \) and any function \( \phi \in C^p \), we have:

\[
\left| \int_{\mathbb{R}^d} \phi(x) \tilde{\rho}(dx) - \int_{\mathbb{R}^d} \phi(x) \tilde{\rho}^\Delta(dx) \right| \leq \varepsilon, \quad a.s.,
\]

(5.1)

\[
\left| \int_{\mathbb{R}^d} \phi(x) \tilde{\rho}(dx) - \int_{\mathbb{R}^d} \phi(x) \tilde{\rho}^\Delta(dx) \right| \leq \varepsilon, \quad a.s.,
\]

(5.2)

Proof. Define \( U(s + i \Delta t, s, x) = \mathbb{E} \phi(X_{s+i\Delta t}^{s,x}) \). Then

\[
U(s, s - k \tau, x) = \mathbb{E} \phi(X_{s-k\tau,x}^{s,x}), \quad U(s, s, \hat{X}_{s-k\tau,x}^{s,x}) = \phi(\hat{X}_{s-k\tau,x}^{s,x}), \quad a.s.
\]

(5.4)

By the periodicity of \( U(t, s, x) \) with respect to initial time \( s \), i.e. \( U(t, s - k \tau, x) = U(t + k \tau, s, x) \), it is always possible to move the initial time into \([0, \tau]\). Now we consider the following Itô-Taylor expansion:

\[
\mathbb{E} U(s + k \tau, s + (i + 1) \Delta t, \hat{X}_{s+(i+1)\Delta t}^{s,x}) = \mathbb{E} U(s + k \tau, s + (i + 1) \Delta t, \hat{X}_{s}^{s,x})
\]

\[
+ \mathcal{L}(s + (i + 1) \Delta t) \mathbb{E} U(s + k \tau, s + (i + 1) \Delta t, \hat{X}_{s+(i+1)\Delta t}^{s,x}) \Delta t + R_1^{s,x}(\Delta t)^2.
\]

(5.5)

Denote \( s' = s + (i + 1) \Delta t \) and \( t' = k \tau - (i + 1) \Delta t \). Then it is obvious that

\[
\mathbb{E} U(s + k \tau, s + i \Delta t, x) = \mathbb{E} U((s' - \Delta t) + (t' + \Delta t), s' - \Delta t, x).
\]

Therefore, we have the following Itô-Taylor expansion:

\[
\mathbb{E} U(s + k \tau, s + i \Delta t, \hat{X}_{s+(i+1)\Delta t}^{s,x}) = \mathbb{E} U(s + k \tau, s + (i + 1) \Delta t, \hat{X}_{s+(1+1)\Delta t}^{s,x})
\]

\[
+ \left( \frac{\partial}{\partial t} - \frac{\partial}{\partial s} \right) \mathbb{E} U(s + t, s, \hat{X}_{s+t+\Delta t}^{s,x}) \Delta t \bigg|_{s=s', t=t'} + R_2^{s,x}(\Delta t)^2.
\]

(5.6)

The coefficients \( R_1^{s,x} \) and \( R_2^{s,x} \) have the following form:

\[
\mathbb{E} \left[ \psi(\hat{X}_{s+(i+1)\Delta t}^{s,x}) - \psi(\hat{X}_{s+i\Delta t}^{s,x}) \right] = \mathbb{E} \left[ \psi(\hat{X}_{s+(i+1)\Delta t}^{s,x}) - \psi(\hat{X}_{s+i\Delta t}^{s,x}) \right]
\]

(5.7)
where $0 < \vartheta < 1$ and the function $\psi(x)$ is a product of functions $b$, $\sigma$ and their derivatives. One can obtain the boundedness of $\psi(x)$ from Condition (1). Combining (5.3) and (5.6), we have

$$
\psi \leq \frac{1}{\tau} \int_{0}^{T} R_i^s \, ds
$$

and their derivatives. One

where $R_i^s = R_i^s - R_i^t$. Combining this with Proposition 4.3 and Theorem 3.3, there exists a constant $\lambda > 0$ and an integer $l \in \mathbb{N}$, such that

$$
\sum_{i=0}^{kN-1} \frac{1}{\tau} \int_{0}^{T} R_i^s \, ds \leq \frac{C}{\tau} \sup_{i \geq 0} \mathbb{E}\left(1 + \left| \hat{X}_{-kN} \right|^l + \left| \hat{X}_{-kN+1} \right|^l \right)^{kN-1} \sum_{i=0}^{kN-1} \exp \left(-\lambda(kN-i)\Delta t\right).
$$

Let $k$ go to infinity and $\Delta t$ be small enough, we have

$$
\sum_{i=0}^{kN-1} \frac{1}{\tau} \int_{0}^{T} |R_i^s| (\Delta t)^2 \, ds \leq \tilde{C}(1 + |x|^l)(\Delta t).
$$

(5.9)

This is then followed by

$$
\lim_{n \to \infty} \frac{1}{n} \sum_{k=1}^{n} \left| \frac{1}{\tau} \int_{0}^{T} \mathbb{E}(\hat{X}_{s-kN}^s(x)) \, ds - \frac{1}{\tau} \int_{0}^{T} \mathcal{U}(s, s - kN, x) \, ds \right| \leq \tilde{C}(1 + |x|^l)(\Delta t).
$$

(5.11)

Thus (5.3) follows as the left hand side of (5.11) is independent of $x$ and $\varepsilon$. \hfill \Box

**Remark 5.2.** Consider the modification of Benzi-Parisi-Sutera-Vulpiani’s stochastic resonance model mentioned in the introduction, the coefficients of Linear growth is estimated as $L \leq \frac{\partial b(s, x)}{\partial x} \bigg|_{x=4} \leq 48$. Hence the step size $\Delta t < \frac{1}{24}$ will satisfy our theorem.
6 Numerical examples

In this section, we carry out some numerical experiments to support the theoretical results obtained in the last section. We give the error analysis for the numerical scheme of the average of periodic measures of two specific models arising in modelling daily temperature and climate dynamics respectively. For each example, we firstly generate discrete random periodic paths $\hat{X}_{s+k\Delta t}^{s,x}$, $k = 0, \ldots, N - 1$ and test the convergence with different initial values. The numerical error we calculate in this section is

$$
\left\| \frac{1}{N} \sum_{k=0}^{N-1} \phi \left( \hat{X}_{s+k\Delta t}^{s,x} \right) - \int_{S} \phi(x) d\hat{\rho}(x) \right\|,
$$

(6.1)

which consists of three parts of errors: those influenced by the finiteness of $N$, the discretization error of time integral on $s$ in (5.10) and the error given in (5.10). Here $\phi \in C^{\infty}_{p}$. The main task is to estimate the error in (6.1) in terms of rate with respect to $\Delta t$. We choose large enough $N$ to reduce its impact on the error. We also compare numerically the errors with different initial time $s$ and find that the convergence of solutions to the random periodic paths in both models are very fast, as seen in Figure 2 and Figure 5, so the effect of the initial time and position to reduce its impact on the error. We also compare numerically the errors with different initial values. The numerical example 6.1.

To present the error of our approximation scheme, we study the following temperature model considered by F. Benth and J. Benth [3],

$$
dX_t = (a_0 + a_1 \cos(2\pi t - a_2)/365) - \pi X_t dt + \sigma dW_t,
$$

with $a_0 = 6.4$, $a_1 = 10.4$, $a_2 = -166$ and $\sigma = 0.3$. From the discussion in [13], it is known that the periodic measure of this model exists and is a Gaussian distribution with mean

$$
\frac{a_0}{\pi} + a_1 \frac{k \sin(k(t - a_2)) + \pi \cos(k(t - a_2))}{k^2 + \pi^2},
$$

and variance $\frac{\sigma^2}{\pi^2}$, where $k = \frac{2\pi}{365}$. This is the case where the periodic measure is known explicitly. But a numerical experiment of calculating the numerical error is carried out here in order to verify the accuracy of our scheme. To simplify the calculation, we take the test function $\phi(x) = x^2$ as $E[X^2] = E[X]^2 + \text{Var}[X]$. Under the average of periodic measure, one can derive the exact value

$$
\int_{\mathbb{R}^d} x^2 d\hat{\rho}(x) = \frac{a_0^2}{\pi^2} + \frac{a_1^2}{2(\pi^2 + k^2)} + \frac{\sigma^2}{2\pi}.
$$

On the other hand, conducting numerical approximation with Euler-Maruyama scheme, we obtain $\hat{X}_{s+k\Delta t}^{s,x}$ for a range of different $\Delta t < 2/\pi = 0.63662$ varying from 0.2 to 0.004. We apply the Euler-Maruyama scheme with same length of time of 10000 periods for different step size $\Delta t$. The error is presented in Table 4 and as a log-log graph in Figure 4. Our numerical results show very good order 1 line fitting. Note the exact value $\int_{\mathbb{R}^d} x^2 d\hat{\rho}(x) = 10.266021$ (rounded off in 6 decimal places).

In Figure 4 we present two numerical approximations, $\hat{X}^{0.10}_s$ and $\hat{X}^{0.10}_s$, to random periodic path with different initial values. These two trajectories are generated with the same realisation of noise. They merge before $t = 2$ and show inconspicuous difference after that. We then
Table 1: Numerical results of Example 6.1 where approximations and numerical errors are rounded off in 6 decimal places

| Step sizes | $\Delta t = 0.2$ | $\Delta t = 0.1$ | $\Delta t = 0.08$ | $\Delta t = 0.05$ | $\Delta t = 0.04$ |
|------------|------------------|------------------|-------------------|------------------|------------------|
| Approximation | 10.560044 | 10.383899 | 10.357855 | 10.318126 | 10.307453 |
| Numerical error | 0.294024 | 0.117879 | 0.091834 | 0.052105 | 0.041433 |
| CPU(seconds) | 175.41 | 371.39 | 465.71 | 756.28 | 951.37 |

| Step sizes | $\Delta t = 0.02$ | $\Delta t = 0.01$ | $\Delta t = 0.008$ | $\Delta t = 0.005$ | $\Delta t = 0.004$ |
|------------|------------------|------------------|-------------------|------------------|------------------|
| Approximation | 10.283706 | 10.277765 | 10.275399 | 10.271213 | 10.270654 |
| Numerical error | 0.017686 | 0.011745 | 0.009379 | 0.005192 | 0.004634 |
| CPU(seconds) | 1600.47 | 3711.75 | 4407.75 | 6551.42 | 7600.27 |

In practice, we run the computation with 10 different step sizes simultaneously under ”multiprocessing” package of Python with 7 cores of CPU. The above results of error analysis took 7600.286 seconds of computing time where the CPU time of each step size is given in Table 1. One can see the majority of time was consumed under the small step sizes such as $\Delta t = 0.004$ and $\Delta t = 0.005$.

If necessary, one can split the approximation of random periodic path with small step sizes into several jobs. This works well due to ergodicity and fast convergence to random periodic path under our scheme. We do not need it in Example 6.1 as the computing time is reasonably short. But the possibility to split the computation into several independent jobs plays a crucial role in the case when a model has a large period. For such a problem, we need to consider the long time behaviour of $N\tau$ where both $N$ and $\tau$ are large. We will see that in the following example.

**Example 6.2.** We consider Benzi-Parisi-Sutera-Vulpiani’s climate dynamics model given by
Figure 2: Paths of the temperature model (Example 6.1)

Figure 3: Comparisons between approximation of periodic measure with $\Delta t = 0.01$ and 10000 periods and its theoretical result (Example 6.1)
SDE \((1.1)\) with \(b(t, x) = x - x^3 + A \cos(Bt)\) and \(\sigma(x) = \sigma\). The coefficients are chosen as \(A = 0.12, B = 0.001\) and \(\sigma = 0.285\) as discussed in [10] and [14]. We make a time scaling by taking \(b(t, x) = 0.4 \pi (x - x^3 + 0.12 \cos(0.0004 \pi t))\) and \(\sigma(x) = 0.285 \times \sqrt{0.4 \pi} \approx 0.3195\), so the period \(\tau = 5000\) in the system. Thus when we apply numerical approximation to the model, we can ensure our time step size dividing the period \(\tau = 5000\). It is also mollified to satisfy the global Lipschitz assumption in this paper. For this, what we could do is to modify the function \(x - x^3\) by a linear function \(128 - 47x\) when \(x \geq 4\) and \(-128 - 47x\) when \(x \leq -4\) and smooth this function by mollifier \(\eta(x) = \eta(x)\), where \(\eta(x) = \begin{cases} C\exp\left(\frac{1}{|x|^2}-1\right) & |x| \leq 1 \\
0 & \text{otherwise} \end{cases}\) such that \(\int \eta(x)dx = 1\). But this adds a lot of computing time as integration of convolution is needed in every step of the computation.

In our approximation, the drift term is \(\hat{b}(t, x) = 0.4 \pi ((1-\exp(-\frac{-5000}{285}))(x-x^3)+0.12 \cos(0.0004 \pi t))\). This function makes a very good approximation to function \(b(t, x)\) when \(|x| \leq 4\) though it is not the case globally. Note this function is Lipschitz and smooth. As we mentioned in the introduction, our modified model provides the same climate dynamics as the original one of Benzi-Parisi-Sutera-Vulpiani. Our numerical simulations presented in Figure [3] for the modified equation provide strong evidence that is the case as seen in the Figure 5 that the trajectory rarely goes outside of \([-4,4]\). In fact, during the approximation with 10000 periods, we did not find any point in the whole data set running out of this interval.

Note this model does not have an explicit solution, so we cannot carry out the error analysis as we did in Example [6.7]. To overcome this difficulty we use the approximation of solution with \(\Delta t = 0.001\) to replace our exact solution in the error analysis. To make the computation more efficient, we split the approximation involving \(\frac{5000 \times 10000}{0.001} = 5 \times 10^{10}\) iterations to 8 individual jobs of \(6.25 \times 10^9\) iterations with independent Brownian motions. The results are shown in the left hand side table of Table [3]. We then conduct the numerical experiment for step size varying from 1/125 to 1/50. The error is in the right hand side table of Table [3] and the log-log graph is presented in Figure [4]. We carry out numerical simulation with 10000 periods for each step size and our total cost is \(1555\times 7.5\) seconds with 7 cores.

We consider numerical simulation with \(\Delta t = 0.01\) to show the stochastic resonance phenomenon in Figure [3]. The simulations start from different initial condition but the same realisation of noise in each sub-graph, where one can see the convergence to random periodic path is also very fast. Together with ergodicity, it provides the possibility of splitting 10000 periods into several independent approximations for the step size \(\Delta t = 0.001\). Without the split, our total computing time would be about \(5 \times 10^5\) seconds as 6 cores of CPU are idle for long time.

We also generate the periodic measure approximations from two paths each with 10000 periods and the step size \(\Delta t = 0.01\) under two different realisations of noise. The distributions of periodic measure are presented in Figure [4]. One can see the distributions produced by two different Brownian motions are very similar. There are some minor differences due to insufficient amount of data. If we utilise sufficiently large amount of computations, the differences will eventually disappear.

We can apply time scaling on the model to rescale its period to a much smaller number. By doing this, the Lipschitz coefficient of the drift term will become very large. According to the upper bound of step size in Theorem 5.1, the total cost of approximation will remain the same as the step size has to be very small.
| Step sizes | ∆t = 1/50 | ∆t = 1/64 | ∆t = 1/80 |
|------------|-----------|-----------|-----------|
| Approximation | 1.0104580 | 1.0106331 | 1.0109203 |
| Numerical error | 0.0014689 | 0.0012938 | 0.0010066 |
| CPU(seconds) | 28630.88 | 36177.61 | 45509.97 |

Table 2: Numerical results of Example 6.2 where numerical results and errors are rounded off in 7 decimal places.

Figure 4: Error of approximation to average of periodic measure versus step size in log-log graph (Example 6.2)

Figure 5: Paths of stochastic resonance model (Example 6.2)
Figure 6: Approximations of periodic measure with $\Delta t = 0.01$ and 10000 periods (Example 6.2)

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