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# Iterated averaging of three-scale oscillatory systems 

Gil Ariel, Bjorn Engquist, Seong Jun Kim and Richard Tsai


#### Abstract

A theory of iterated averaging is developed for a class of highly oscillatory ordinary differential equations (ODEs) with three well separated time scales. The solutions of these equations are assumed to be (almost) periodic in the fastest time scales. It is proved that the dynamics on the slowest time scale can be approximated by an effective ODE obtained by averaging out oscillations. In particular, the effective dynamics of the considered class of ODEs is always deterministic and does not show any stochastic effects. This is in contrast to systems in which the dynamics on the fastest time scale is mixing. The systems are studied from three perspectives: first, using the tools of averaging theory; second, by formal asymptotic expansions; and third, by averaging with respect to fast oscillations using nested convolutions with averaging kernels. The latter motivates a hierarchical numerical algorithm consisting of nested integrators.


## 1 Introduction

Averaging methods are some of the most widely used tools for the study and approximation of highly oscillatory ordinary differential equations (ODEs). In its most basic form, averaging means that following a change of variables that separates the system into slow coordinates and periodic fast ones, the effect of the fast oscillations can be integrated out, yielding an approximate effective equation. More precisely, consider the initial value problem

$$
\begin{equation*}
x^{\prime}=f\left(x, \varepsilon^{-1} t\right), \quad x(0)=x_{0} \tag{1}
\end{equation*}
$$

where $0<\varepsilon \leq \varepsilon_{0} \ll 1$ denotes the scale separation in the problem. It is assumed that a unique uniformly bounded solution $x(t) \in \mathbb{R}^{d}$ exists in a time segment $[0, T]$ which is independent of $\varepsilon$. Furthermore, $f(x, s)$ is sufficiently smooth and 1-periodic in $s$. See, for example $[11,40]$ for details. Then, $x(t)$ can be approximated by the effective averaged equation

$$
\begin{equation*}
\xi^{\prime}=F(\xi), \quad \xi(0)=x_{0} \tag{2}
\end{equation*}
$$

where $F(\xi)$ is the average of $f(x, s)$ over one period of $s$

$$
\begin{equation*}
F(\xi)=\langle f(\xi, \cdot)\rangle_{s} \equiv \int_{0}^{1} f(\xi, s) d s \tag{3}
\end{equation*}
$$

Furthermore, the approximation is of order $\varepsilon$ in the sup norm, i.e.

$$
\begin{equation*}
\sup _{0 \leq t \leq T}|x(t)-\xi(t)|=\mathscr{O}(\varepsilon) \tag{4}
\end{equation*}
$$

[^0]This elementary result has been generalized considerably to include, for example, chaotic and stochastic systems [7, $10,11,30,31,38,39,40]$. From a computational point of view, averaging methods inspire efficient numerical schemes for integrating the slow components of slow-fast systems without fully resolving all fast oscillations in $[0, T]$. See, for example $[1,2,8,14,15,16,18,19,23,25,41,42,44]$ and references therein.

The main purpose of this paper is twofold. First, we generalize the classical two-scale averaging theory (1)-(4) to a class of highly oscillatory systems whose solutions possess slow variables or observables as well as fast oscillations with frequencies that are separated by two large spectral gaps. We shall call such systems three-scale oscillatory systems. Our generalization is along the lines of the so called iterated homogenization problem [9, 10, 12, 29, 34]. Similarly, we develop a theory of iterated averaging which is approached from three different points of view: (1) application of long time asymptotic and second order averaging theorems [40, 43]; (2) formal perturbation expansions applied to an associated SDE; and (3) a direct proof by repeated convolution with averaging kernels across different time scales. The latter is the foundation for a hierarchical numerical algorithm consisting of nested integrators, which is the second main objective of the paper. We demonstrate how tiers of two-scale numerical methods can be "stacked-up" in order to construct a consistent approximation of the effective slow dynamics in three-scale systems. To our knowledge, very few algorithms have been developed considering directly three or more scales. We develop our algorithm using the framework of the heterogeneous multiscale method (HMM) for highly oscillatory ODEs [2, 16]. A proof of concept for the algorithm was previously suggested in [6] without the underlying mathematical theory of iterated averaging or proof of convergence.

To gain insight on the analysis of the numerical method described in [6], it is helpful to study the method when it is applied to an appropriately chosen and simpler problem which shows the similar features of the original model but is easier to analyze. We thus develop a simple model, motivated by Fermi, Pasta and Ulam (FPU) [24], which shows the multiscale behavior in time scales of integer powers of $\varepsilon$. The FPU model involves almost-periodic motion on the $\varepsilon$ time scale, motion of the soft nonlinear springs on the time scale independent of $\varepsilon$, and the slow energy exchange among stiff springs on the $\varepsilon^{-1}$ time scale; see e.g. [28]. We remark however that there of course exist dynamical systems whose multiscale features are not integer powers of $\varepsilon$, e.g. in weakly coupled, nearly resonant harmonic oscillators. A review of the near-resonance is beyond the scope of this paper, but a new approach for these general systems will be reported in a forthcoming paper.

On this account, we consider a system of ODEs written in the following form;

$$
\begin{cases}x^{\prime}=f_{1}(x, y, z)+\varepsilon f_{0}(x, y, z), & x(0)=x_{0}  \tag{5}\\ y^{\prime}=g_{1}(x, y, z)+\varepsilon g_{0}(x, y, z), & y(0)=y_{0} \\ z^{\prime}=\varepsilon^{-1} h_{2}(x, y, z), & z(0)=z_{0}\end{cases}
$$

where $f_{1}$ and $f_{0}$ are periodic in both $y$ and $z ; z(t) \in \mathbb{R}$ is the fast variable, and its time derivative in $\mathscr{O}\left(\varepsilon^{-1}\right)$ time interval remains $\mathscr{O}\left(\varepsilon^{-1}\right) ; y(t) \in \mathbb{R}$ is the intermediate variable with derivative bounded uniformly in any $\mathscr{O}\left(\varepsilon^{-1}\right)$ time intervals; $x \in \mathbb{R}^{d}$ is the slow variable, which effectively varies on the time scale of $\mathscr{O}\left(\varepsilon^{-1}\right)$ or longer. More precisely, there exists slowly changing quantity $\xi(t)$ over $t \in\left[0, \varepsilon^{-1} T\right]$ and $C>0$ such that

$$
\sup _{0 \leq t \leq \varepsilon^{-1} T}|x(t)-\xi(t)| \leq C \varepsilon
$$

where $T=\mathscr{O}(1)$. When only looking for the effective dynamics $\xi(t)$ of the slow variable $x(t)$ in the time interval $[0, T]$ which is independent of $\varepsilon$, one can simply truncate the $\varepsilon f_{0}$ term, and average over $y$ and $z$. On the other hand, in case of $\varepsilon^{-1} T$, we cannot neglect $\varepsilon f_{0}$ because its effect appears in this longer time interval. We will show that the condition for existence of $\xi(t)$ is that $f_{1}$ should have a zero average in $z$ for fixed $x$ and $y$.

For convenience we rescale (5) by replacing $t$ by $\varepsilon^{-1} t$ so that the slowest time scale of interest is independent of $\varepsilon$. When more than two time scales are considered as above, a new type of slow variables whose formal derivatives are not bounded along the flow of the dynamical system appears. Note that the time derivative of $x(t)$ now becomes formally unbounded for $0<\varepsilon \leq \varepsilon_{0}$. Such type of slow variables does not exist for problems with only two separated time scales. We refer to [6] for more detailed discussions. In addition, we emphasize that the effective behavior described by this type of slow variables should be taken into account on building multiscale algorithms. This will be discussed later in Section 5.2.

We further assume that the evolution of $x(t)$ is essentially constant in time intervals whose length is $o(1)$, i.e.,

$$
\left|g_{1}(x(t), y(t), z(t))-g_{1}\left(x\left(t_{0}\right), y(t), z(t)\right)\right| \leq C_{1} \varepsilon, \quad\left|h_{2}(x(t), y(t), z(t))-h_{2}\left(x\left(t_{0}\right), y(t), z(t)\right)\right| \leq C_{2} \varepsilon
$$

hold over $\mathscr{O}(\varepsilon), \mathscr{O}\left(\varepsilon^{2}\right)$ time intervals respectively. In this context, $y(t)$ and $z(t)$ is essentially independent of $x(t)$ and the following simplified model equation will be considered in the paper:

$$
\begin{equation*}
x^{\prime}=\varepsilon^{-1} f_{1}\left(x, \varepsilon^{-1} t, \varepsilon^{-2} t\right)+f_{0}\left(x, \varepsilon^{-1} t, \varepsilon^{-2} t\right), x(0)=x_{0} \tag{6}
\end{equation*}
$$

where $f_{1}\left(x, s_{1}, s_{2}\right)$ and $f_{2}\left(x, s_{1}, s_{2}\right)$ are sufficiently smooth (e.g. $C^{1}$ ) and 1-periodic in $s_{1}$ and $s_{2}$. It is assumed that a unique bounded solution exists in a time segment $[0, T]$. Both the bound and $T$ are independent of $\varepsilon$. The dynamics of (6) can be characterized across three well-separated time scales: a fast time scale involving time intervals with length of order $\varepsilon^{2}$, an intermediate scale of order $\varepsilon$ and a slow $\mathscr{O}(1)$ scale.

In the following, we shall use $s_{1}$ and $s_{2}$ exclusively for the variables which will be scaled respectively to $\varepsilon^{-1} t$ and $\varepsilon^{-2} t$. In the following, averages of a function $f\left(x, s_{1}, s_{2}\right)$ with respect to $s_{1}$ and $s_{2}$ with $x$ fixed are denoted respectively by

$$
\left\langle f\left(\xi, s_{1}, s_{2}\right)\right\rangle_{1}, \quad\left\langle f\left(\xi, s_{1}, s_{2}\right)\right\rangle_{2}
$$

The double average of $f\left(x, s_{1}, s_{2}\right)$ over both $s_{1}$ and $s_{2}$ is denoted by

$$
\left\langle f\left(x, s_{1}, s_{2}\right)\right\rangle_{12}
$$

We proceed formally. Assuming that for fixed $x$ and $s_{1}$ the average of $f_{1}$ with respect to $s_{2}$ vanishes, $x(t)$ can be written as a sum of a smooth function $w(t)$ whose first time derivative is bounded independent of $\varepsilon$, and fast oscillations around $w(t)$ whose amplitudes tend to zero as $\varepsilon \rightarrow 0[37,40]$. As a motivation for identifying an ODE for $w(t)$, consider the form

$$
x(t)=w(t)+\varepsilon \phi\left(w(t), \varepsilon^{-1} t, \varepsilon^{-2} t\right) .
$$

Substituting into (6) yields

$$
\begin{equation*}
w^{\prime}=\varepsilon^{-1}\left(f_{1}-\frac{\partial \phi}{\partial s_{2}}\right)+\left(\nabla_{x} f_{1}\right) \phi-\frac{\partial \phi}{\partial s_{1}}+f_{0}+\mathscr{O}(\varepsilon) \tag{7}
\end{equation*}
$$

In order to have that $w^{\prime}(t)$ is bounded independent of $\varepsilon$ we require that $\partial \phi / \partial s_{2}=f_{1}$. Consequently, $\phi$ should be taken to be the anti-derivative of $f_{1}$ with respect to $s_{2}$ and

$$
\begin{equation*}
w^{\prime}=\left(\nabla_{x} f_{1}\right) \phi-\frac{\partial \phi}{\partial s_{1}}+f_{0}+\mathscr{O}(\varepsilon) \tag{8}
\end{equation*}
$$

Since $\phi$ is periodic, $\partial \phi / \partial s_{1}$ averages to zero and therefore does not contribute to the effective equation. In order to eliminate small oscillations in $\phi$, the right hand side (RHS) of (8) should be averages with respect to $s_{1}$ and $s_{2}$. This leads to the averaged equation

$$
\begin{equation*}
\xi^{\prime}=F(\xi), \quad \xi(0)=x_{0} \tag{9}
\end{equation*}
$$

where

$$
\begin{align*}
& F=\left\langle\left(\nabla_{x} f_{1}\right) \phi+f_{0}\right\rangle_{12} \\
& \phi\left(\xi, s_{1}, s_{2}\right)=\int_{0}^{s_{2}} f_{1}\left(\xi, s_{1}, \tau\right) d \tau \tag{10}
\end{align*}
$$

Remark 1. The definition of $\phi$ is not unique. For example, it is straightforward to verify that taking $\phi\left(\xi, s_{1}, s_{2}\right)=$ $\int_{0}^{s_{2}} f_{1}\left(\xi, s_{1}, \tau\right) d \tau+a\left(\xi, s_{1}\right)$ where, for fixed $\xi, a\left(\xi, s_{1}\right)$ is periodic in $s_{1}$, one would arrive at the same averaged equation. Two canonical choices are either $a\left(\xi, s_{1}\right)=0$ which is used in Sections 2 and 4 , or $a\left(\xi, s_{1}\right)$ such that $\langle\phi\rangle_{2}=0$ which is used in Section 3.

Remark 2. Following the previous remark we note that in general, it is not true that writing $x=w+\varepsilon \phi+\mathscr{O}\left(\varepsilon^{2}\right)$.
The following theorem is proved in Section 2.
Theorem 1. Let $x(t)$ and $\xi(t)$ denote solutions of equations (6) and (9) respectively. If $f_{1}$ has a zero average with respect to $s_{2}$, then there exists a constant $C>0$ independent of $\varepsilon$ such that

$$
\sup _{0 \leq t \leq T}|x(t)-\xi(t)| \leq C \varepsilon .
$$

In fact, as will be seen later, the same theorem holds under weaker assumptions on the fast and intermediate time scales $s_{1}$ and $s_{2}$. The fast, $\varepsilon^{2}$ scale can be periodic, ergodic on a torus or almost periodic with frequencies that are bounded from below. On the intermediate $\mathscr{O}(\varepsilon)$ scale, the only requirement is that for fixed $\xi$, the effective dynamics obtained after integrating the fastest $\varepsilon^{2}$ time scale is ergodic and that a two-scale averaging principle [40] holds. This includes both oscillatory, stochastic and chaotic dynamics. Such generalizations are further discussed in Section 5.

### 1.1 A simple example

Consider the ODE system

$$
\begin{equation*}
x^{\prime}=2 \pi \varepsilon^{-1}\left[x \sin \left(2 \pi \varepsilon^{-2} t\right)+\cos \left(2 \pi \varepsilon^{-2} t\right)\right], x(0)=0 \tag{11}
\end{equation*}
$$

for $0 \leq t \leq 1$. By the method of variation of parameters, the exact solution of (11) is found to be

$$
\begin{equation*}
x(t)=2 \pi \varepsilon^{-1} e^{-\varepsilon \cos \left(2 \pi \varepsilon^{-2} t\right)} \int_{0}^{t} \cos \left(2 \pi \varepsilon^{-2} s\right) e^{\varepsilon \cos \left(2 \pi \varepsilon^{-2} s\right)} d s \tag{12}
\end{equation*}
$$

Expanding in $\varepsilon$,

$$
\begin{equation*}
x(t)=\pi t+\varepsilon\left[\sin \left(2 \pi \varepsilon^{-2} t\right)-\pi t \cos \left(2 \pi \varepsilon^{-2} t\right)\right]+\mathscr{O}\left(\varepsilon^{2}\right) \tag{13}
\end{equation*}
$$

Hence, $x(t)=\pi t+\mathscr{O}(\varepsilon)$ for $t \in[0,1]$.
A naive approach for obtaining an effective equation is to assume that in any sufficiently short time interval the solution of (11) can be approximated by integrating over the fast oscillations in the RHS of (11) while keeping $x$ fixed. This yields

$$
\begin{equation*}
\tilde{x}(t)=\left.2 \pi \varepsilon^{-1} \int_{0}^{t}\left[X \sin \left(2 \pi \varepsilon^{-2} s\right)+\cos \left(2 \pi \varepsilon^{-2} s\right)\right] d s\right|_{X=\tilde{x}(t)}=2 \pi \varepsilon^{-1}\left[X \mathscr{O}\left(\varepsilon^{2}\right)+\mathscr{O}\left(\varepsilon^{2}\right)\right] \tag{14}
\end{equation*}
$$

which implies that $\tilde{x}(t)=\mathscr{O}(\varepsilon)$ in $t \in[0,1]$.
In contrast, Theorem 1 yields the correct averaged equation for $x$ for $t \in[0,1]$ :

$$
X^{\prime}=\pi, \quad X(0)=0
$$

which clearly implies that

$$
\sup _{t \in[0,1]}|x(t)-X(t)| \leq(1+\pi) \varepsilon
$$

On the other hand, performing averages over fast oscillations can be approximated in a convenient and computationally efficient way by convolution of $x^{\prime}$ with appropriate compactly supported kernels. Inside the convolution, the value of $x$ is not exactly fixed but varies following the correct dynamics. This subtle change in the values of $x(t)$ allows for the kernel to capture the correct effective change of $x(t)$ in a longer time scale. This can be demonstrated in the case of example (11). Consider a cosine kernel, which is particularly convenient in this example

$$
K(t)=\frac{1}{2} \chi_{[-1,1]}(t)[1+\cos (\pi t)],
$$

where $\chi_{A}$ is the indicator function of a set $A$. Furthermore, for $\eta>0$, let $K_{\eta}$ denote a scaling of $K(t)$ to $[-\eta, \eta]$,

$$
K_{\eta}(t)=\frac{1}{\eta} K\left(\frac{t}{\eta}\right)
$$

Without loss of generality, we calculate the convolution of $K_{\eta}$ with $x^{\prime}$ in (11) at $t=0$

$$
\left(K_{\eta} * x^{\prime}\right)(0)=\frac{2 \pi}{\varepsilon} \int_{-\eta}^{\eta} K_{\eta}(-s)\left[x(s) \sin \left(2 \pi \varepsilon^{-2} s\right)+\cos \left(2 \pi \varepsilon^{-2} s\right)\right] d s
$$

Substituting in the expansion of $x(s)$ shown in Eq. (13), and using $\eta=n \varepsilon^{2}$ yields

$$
\begin{aligned}
\left(K_{\eta} * x^{\prime}\right)(0)= & \frac{\pi}{\eta \varepsilon} \int_{-n \varepsilon^{2}}^{n \varepsilon^{2}}\left[1+\cos \left(\frac{\pi s}{\eta}\right)\right]\left[\pi s \sin \left(\frac{2 \pi n s}{\eta}\right)+\cos \left(\frac{2 \pi n s}{\eta}\right)\right] d s \\
& +\frac{\pi}{\eta} \int_{-n \varepsilon^{2}}^{n \varepsilon^{2}}\left[1+\cos \left(\frac{\pi s}{\eta}\right)\right]\left[\sin ^{2}\left(\frac{2 \pi n s}{\eta}\right)-\frac{\pi s}{2} \sin \left(\frac{4 \pi n s}{\eta}\right)\right] d s+\mathscr{O}(\varepsilon) \\
= & \pi+\mathscr{O}(\varepsilon)
\end{aligned}
$$

which is, to leading order in $\varepsilon$, the correct derivative for the slow variable $x$.
In Section 4 these ideas are generalized to oscillatory three-scale systems in which kernels are applied iteratively to the different time scales. This can then be exploited for construction of efficient multiscale numerical schemes.

### 1.2 Formal asymptotic expansions

Formal asymptotic expansions of singularly perturbed operators have been successfully applied to a wide variety of problems [10, 39]. Of particular relevance are the applications to stochastic differential equations [15, 35, 38, 44]. Consider SDE systems of the form

$$
\left\{\begin{array}{l}
d x_{t}=\left[\varepsilon^{-1} f_{1}\left(x_{t}, y_{t}\right)+f_{0}\left(x_{t}, y_{t}\right)\right] d t  \tag{15}\\
d y_{t}=\varepsilon^{-2} a\left(x_{t}, y_{t}\right) d t+\varepsilon^{-1} \beta\left(x_{t}, y_{t}\right) d B_{t}
\end{array}\right.
$$

where $B_{t}$ is a standard Brownian motion in $\mathbb{R}^{d}$. The variable $x_{t} \in \mathbb{R}^{d}$ is a slow process that evolves according to an ODE with a fast random coefficient $y_{t} \in \mathbb{R}^{n}$. Under some ergodicity, smoothness and growth assumptions, $x_{t}$ can be approximated by an effective equation of the form

$$
\begin{equation*}
d X_{t}=F\left(X_{t}\right) d t+b\left(X_{t}\right) d B_{t} \tag{16}
\end{equation*}
$$

where $F$ and $b$ can be expressed as averages with respect to the fast process $y_{t}$ with $X$ fixed [38,39]. One of the interesting consequences of (16) is that if the dynamics of $y$ is mixing ${ }^{1}$, then the effective diffusion $b(x)$ may be non-zero even if $y_{t}$ is deterministic, i.e., $\beta=0$. For details and examples, see [17, 23, 26, 27, 35, 36, 39].

Many oscillatory dynamics are not mixing, even though for fixed $x_{t}$ the fast dynamics is ergodic on a lowdimensional invariant manifold. In Section 3 we show that the method of formal asymptotic expansions gives the correct vanishing effective diffusion coefficient, i.e., the effective slow dynamics is deterministic. Furthermore, the method reproduces the correct effective drift.

### 1.3 Numerical methods

The above discussion on applying averaging kernels across different time scales motivates a numerical method which applies our previous two-scale HMM algorithms [3, 4, 5] hierarchically to multiple ( $>2$ ) timescale systems. We consider the time scales $\mathscr{O}\left(\varepsilon^{2}\right), \mathscr{O}(\varepsilon)$ and $\mathscr{O}(1)$ and assume that both the $\varepsilon$ and $\varepsilon^{2}$ scales are oscillatory.

The HMM to be constructed should evaluate the effective rate of change of $x(t)$. For three-scale problems this requires averaging over the $\mathscr{O}\left(\varepsilon^{2}\right)$ as well as the $\mathscr{O}(\varepsilon)$ scale oscillations, thus obtaining a numerical approximation for the effective equation. See [6] for further details.

The hierarchical HMM structure is illustrated in Fig. 1. The downward pointing arrows depict the determination of an initial condition for a lower, fast scale from data in an upper tier working on a slower time scale. The upward pointing arrows from 2nd tier to 1 st tier and 1st tier to 0th tier relate the evaluation of averages with respect to $s_{2}$ and $s_{1}$, respectively. Below we detail the equations solved in each tier.

Let $\eta_{i}$ and $h_{i}$ denote the range of integration and step size used in the $i$-th tier, respectively. A chosen ODE solver in the 2nd tier numerically approximates the full ODE at the initial time $t_{n, m}=n h_{0}+m h_{1}$.

[^1]$$
x_{2}^{\prime}=\varepsilon^{-1} f_{1}\left(x_{2}, \varepsilon^{-1} t, \varepsilon^{-2} t\right)+f_{0}\left(x_{2}, \varepsilon^{-1} t, \varepsilon^{-2} t\right), \quad x_{2}\left(t_{n, m}\right)=X_{1}
$$
in a time segment $t \in\left[t_{n, m}-\eta_{2}, t_{n, m}+\eta_{2}\right]$ where $X_{1}\left(t_{n, m} ; X_{0}\left(t_{n}\right)\right)$ is an approximation of $x(t)$ at $t_{n, m}$ obtained from the 1 st tier. Denote the solution by $x_{2}\left(t ; X_{1}\right)$ and let
$$
\mathscr{F}_{1}\left(t_{n, m} ; X_{1}\right)=K_{\eta_{2}} *\left(\varepsilon^{-1} f_{1}\left(x_{2}\left(\cdot ; X_{1}\right), \varepsilon^{-1} \cdot, \varepsilon^{-2} \cdot\right)+f_{0}\left(x_{2}\left(\cdot ; X_{1}\right), \varepsilon^{-1} \cdot, \varepsilon^{-2} \cdot\right)\right)
$$

The 1 st tier numerically approximates the effective ODE for the $\mathscr{O}(\varepsilon)$ scale

$$
x_{1}^{\prime}=\mathscr{F}_{1}\left(t ; X_{1}\right), \quad x_{1}\left(t_{n}\right)=X_{0} .
$$

in the time interval $t \in\left[t_{n}-\eta_{1}, t_{n}+\eta_{1}\right]$ where $X_{0}\left(t_{n}\right)$ is an approximation of $x(t)$ at $t_{n}=n h_{0}$ obtained from the 0 th tier. Denote the solution $x_{1}\left(t ; X_{0}\right)$ and let

$$
\mathscr{F}_{0}\left(t_{n} ; X_{0}\right)=K_{\eta_{1}} * \mathscr{F}_{1}(\cdot) .
$$

Finally, the 0th tier numerically approximates the effective ODE for the $\mathscr{O}(1)$ scale

$$
x^{\prime}=\mathscr{F}_{0}(x), \quad x(0)=x_{0}
$$

in $t \in\left[0, \eta_{0}\right]=[0, T]$.


Fig. 1 An illustration of a three scale algorithm.

### 1.4 Layout

The layout of the paper is as follows. Section 2 details a proof of Theorem 1 using the tools of averaging theory. In Section 3, the same result is derived using formal asymptotic expansions to singular perturbations of SDEs in which the white noise is turned off. Section 4 proves that the effective dynamics can be approximated using convolution with respect to averaging kernels which are applied iteratively to the different time scales. The methods can be used in a numerical HMM scheme as described above. The HMM scheme suggests that the basic idea of iterated averaging can be applied in more general settings. Such generalizations are discussed in Section 5. A few examples of such generalizations are presented in Section 6. We conclude in Section 7.

## 2 A theory of iterated averaging

In this section we prove Theorem 1 which generalizes an averaging theorem for long time scales due to van der Burgh [40, 43]. Further generalizations are discussed in Section 5.

### 2.1 Basic estimate

As before, let $x(t)$ solve

Iterated averaging

$$
\begin{equation*}
x^{\prime}=\varepsilon^{-1} f_{1}\left(x, \varepsilon^{-1} t, \varepsilon^{-2} t\right)+f_{0}\left(x, \varepsilon^{-1} t, \varepsilon^{-2} t\right), \quad x(0)=x_{0} \tag{17}
\end{equation*}
$$

where $f_{1}\left(x, s_{1}, s_{2}\right)$ and $f_{0}\left(x, s_{1}, s_{2}\right)$ are sufficiently smooth and 1-periodic in $s_{1}$ and $s_{2}$. It is further assumed that the solution $x(t)$ exists, is unique and remains bounded independent of $\varepsilon$ for a time segment $[0, T]$ independent of $\varepsilon$. When $\left\langle f_{1}\right\rangle_{2}=0$, then $x(t)$ can be approximated by a slow trajectory and then our goal is to derive an approximate ODE for such a trajectory. As usual, by slow we mean that the first derivative of a time-dependent function is bounded independent of $\varepsilon$.

We consider two functions $w(t)$ and $y(t)$. Let $w(t)$ solve

$$
\begin{equation*}
w^{\prime}=h\left(w(t), \varepsilon^{-1} t, \varepsilon^{-2} t\right)+f_{0}\left(w(t), \varepsilon^{-1} t, \varepsilon^{-2} t\right), \quad w(0)=x_{0} \tag{18}
\end{equation*}
$$

where $h$ is defined by

$$
\begin{align*}
& h\left(x, s_{1}, s_{2}\right)=\left(\nabla_{x} f_{1}\right) \phi\left(x, s_{1}, s_{2}\right)-\frac{\partial \phi}{\partial s_{1}}\left(x, s_{1}, s_{2}\right),  \tag{19}\\
& \phi\left(x, s_{1}, s_{2}\right)=\int_{0}^{s_{2}} f_{1}\left(x, s_{1}, \tau\right) d \tau
\end{align*}
$$

Note that for fixed $x, h\left(x, s_{1}, s_{2}\right)$ and $\phi\left(x, s_{1}, s_{2}\right)$ are 1-periodic in $s_{1}$ and $s_{2}$.
Notation 1 We use the notation $\left(\nabla_{x} f_{1}\right) \phi\left(x, s_{1}, s_{2}\right)$ for the multiplication of $\phi$ by the derivative of $f_{1}$ with respect to $x$ and both are evaluated at $\left(x, s_{1}, s_{2}\right)$.

Let $y(t)$ solve

$$
\begin{equation*}
y^{\prime}=G\left(y, \varepsilon^{-1} t\right), \quad y(0)=x_{0}, \tag{20}
\end{equation*}
$$

where $G$ is given by

$$
\begin{equation*}
G\left(y, s_{1}\right)=\left\langle h\left(y, s_{1}, s_{2}\right)\right\rangle_{2}+\left\langle f_{0}\left(y, s_{1}, s_{2}\right)\right\rangle_{2} . \tag{21}
\end{equation*}
$$

We will show that for $t \in[0, T]$, there exist nonnegative constants $C_{0}$ and $C_{1}$ independent of $\varepsilon$ such that

$$
|x(t)-w(t)| \leq C_{0} \varepsilon \quad \text { and } \quad|w(t)-y(t)| \leq C_{1} \varepsilon
$$

Thus, we conclude by the triangle inequality for $t \in[0, T]$ that

$$
|x(t)-y(t)| \leq C \varepsilon
$$

We denote a generic positive constant by $C$ whose value may change between expressions but is independent of $\varepsilon$.
Lemma 1. The solutions $x(t)$ and $w(t)$ defined above satisfy

$$
|x(t)-w(t)| \leq \varepsilon C e^{\left(L_{f_{0}}+L_{h}\right) t}
$$

for $t \in[0, T]$ and $T>0$ which is independent of $\varepsilon . L_{f_{0}}$ and $L_{h}$ are Lipschitz constants for $f_{0}$ and $h$, respectively.

Proof. From (17) and (18), integrating with respect to time yields that

$$
\begin{aligned}
& x(t)-x_{0}=\int_{0}^{t} \varepsilon^{-1} f_{1}\left(x(\tau), \varepsilon^{-1} \tau, \varepsilon^{-2} \tau\right) d \tau+\int_{0}^{t} f_{0}\left(x(\tau), \varepsilon^{-1} \tau, \varepsilon^{-2} \tau\right) d \tau \\
& w(t)-x_{0}=\int_{0}^{t} h\left(w(\tau), \varepsilon^{-1} \tau, \varepsilon^{-2} \tau\right) d \tau+\int_{0}^{t} f_{0}\left(w(\tau), \varepsilon^{-1} \tau, \varepsilon^{-2} \tau\right) d \tau
\end{aligned}
$$

This leads to

$$
\begin{equation*}
|x(t)-w(t)| \leq\left|\int_{0}^{t} \varepsilon^{-1} f_{1}\left(x(\tau), \varepsilon^{-1} \tau, \varepsilon^{-2} \tau\right)-h\left(w(\tau), \varepsilon^{-1} \tau, \varepsilon^{-2} \tau\right) d \tau\right|+L_{f_{0}} \int_{0}^{t}|x(\tau)-w(\tau)| d \tau \tag{22}
\end{equation*}
$$

where $L_{f_{0}}$ is a Lipschitz constant for $f_{0}$. We will show that the first integral in (22) is bounded by

$$
L_{h} \int_{0}^{t}|x(\tau)-w(\tau)| d \tau+\mathscr{O}(\varepsilon)
$$

where $L_{h}$ is a Lipschitz constant for $h$. Then, we have

$$
|x(t)-w(t)| \leq\left(L_{f_{0}}+L_{h}\right) \int_{0}^{t}|x(\tau)-w(\tau)| d \tau+\mathscr{O}(\varepsilon)
$$

It follows from Gronwall's inequality that

$$
|x(t)-w(t)| \leq \varepsilon C e^{\left(L_{f_{0}}+L_{h}\right) t}
$$

To this end, let $n$ denote the largest integer such that $\varepsilon^{2} n \leqq t$. We have

$$
\varepsilon^{-1} \int_{0}^{t} f_{1}\left(x(\tau), \varepsilon^{-1} \tau, \varepsilon^{-2} \tau\right) d \tau=\sum_{j=0}^{n-1} \int_{\varepsilon^{2} j}^{\varepsilon^{2}(j+1)} \varepsilon^{-1} f_{1}\left(x(\tau), \varepsilon^{-1} \tau, \varepsilon^{-2} \tau\right) d \tau+\varepsilon^{-1} \int_{\varepsilon^{2} n}^{t} f_{1}\left(x(\tau), \varepsilon^{-1} \tau, \varepsilon^{-2} \tau\right) d \tau
$$

Since $f_{1}$ is bounded and $0 \leq t-n \varepsilon^{2}<\varepsilon^{2}$, the last term is $\mathscr{O}(\varepsilon)$. Denoting $t_{j}=\varepsilon^{2} j$ and using the periodicity of $f_{1}$ in $s_{2}$, each term in the sum can be written as

$$
\begin{equation*}
\varepsilon \int_{0}^{1} f_{1}\left(x\left(t_{j}+\varepsilon^{2} s\right), \varepsilon^{-1} t_{j}+\varepsilon s, s\right) d s \tag{23}
\end{equation*}
$$

By the fundamental theorem of calculus, we write

$$
\begin{align*}
x\left(t_{j}+\varepsilon^{2} s\right) & =x\left(t_{j}\right)+\int_{t_{j}}^{t_{j}+\varepsilon^{2} s} x^{\prime}(\tau) d \tau \\
& =x\left(t_{j}\right)+\int_{t_{j}}^{t_{j}+\varepsilon^{2} s}\left[\varepsilon^{-1} f_{1}\left(x(\tau), \varepsilon^{-1} \tau, \varepsilon^{-2} \tau\right)+f_{0}\left(x(\tau), \varepsilon^{-1} \tau, \varepsilon^{-2} \tau\right)\right] d \tau  \tag{24}\\
& =x\left(t_{j}\right)+\int_{t_{j}}^{t_{j}+\varepsilon^{2} s} \varepsilon^{-1} f_{1}\left(x(\tau), \varepsilon^{-1} \tau, \varepsilon^{-2} \tau\right) d \tau+\mathscr{O}\left(\varepsilon^{2}\right)
\end{align*}
$$

since $f_{0}$ is bounded independent of $\varepsilon$. After the substitution of (24) into (23) and using a Taylor expansion of $f_{1}\left(x\left(t_{j}+\right.\right.$ $\left.\varepsilon^{2} s\right), \varepsilon^{-1} t_{j}+\varepsilon s, s_{2}$ ) around $s=0$ while keeping $s_{2}$ fixed,

$$
\begin{aligned}
\varepsilon \int_{0}^{1} f_{1}\left(x\left(t_{j}+\varepsilon^{2} s\right), \varepsilon^{-1} t_{j}+\varepsilon s, s\right) d s & =\varepsilon \int_{0}^{1}\left\{f_{1}\left(x\left(t_{j}\right), \varepsilon^{-1} t_{j}, s\right)+\left(\int_{t_{j}}^{t_{j}+\varepsilon^{2} s} x^{\prime}(\tau) d \tau\right) \nabla_{x} f_{1}\left(x\left(t_{j}\right), \varepsilon^{-1} t_{j}, s\right)\right. \\
& \left.+\varepsilon s \frac{\partial f_{1}}{\partial s_{1}}\left(x\left(t_{j}\right), \varepsilon^{-1} t_{j}, s\right)\right\} d s+\mathscr{O}\left(\varepsilon^{3}\right) \\
& =\int_{0}^{1}\left(\int_{t_{j}}^{t_{j}+\varepsilon^{2} s} f_{1}\left(x(\tau), \varepsilon^{-1} \tau, \varepsilon^{-2} \tau\right) d \tau\right) \nabla_{x} f_{1}\left(x\left(t_{j}\right), \varepsilon^{-1} t_{j}, s\right) d s \\
& +\int_{0}^{1} \varepsilon^{2} s \frac{\partial f_{1}}{\partial s_{1}}\left(x\left(t_{j}\right), \varepsilon^{-1} t_{j}, s\right) d s+\mathscr{O}\left(\varepsilon^{3}\right)
\end{aligned}
$$

where we used the fact that $\left\langle f_{1}\right\rangle_{2}=0$. Hence, for $\varepsilon^{-1} \int_{0}^{t} f_{1} d \tau$,

$$
\begin{align*}
\varepsilon^{-1} \int_{0}^{t} f_{1}\left(x(\tau), \varepsilon^{-1} \tau, \varepsilon^{-2} \tau\right) d \tau & =\sum_{j=0}^{n-1}\left[\int_{0}^{1}\left(\int_{t_{j}}^{t_{j}+\varepsilon^{2} s} f_{1}\left(x(\tau), \varepsilon^{-1} \tau, \varepsilon^{-2} \tau\right) d \tau\right) \nabla_{x} f_{1}\left(x\left(t_{j}\right), \varepsilon^{-1} t_{j}, s\right) d s\right.  \tag{25}\\
& \left.+\int_{0}^{1} \varepsilon^{2} s \frac{\partial f_{1}}{\partial s_{1}}\left(x\left(t_{j}\right), \varepsilon^{-1} t_{j}, s\right) d s+\mathscr{O}\left(\varepsilon^{3}\right)\right]
\end{align*}
$$

Next, changing of variables to $s=\varepsilon^{-2} \tau$ in $\int_{0}^{t} h d \tau$ gives

$$
\begin{aligned}
\int_{0}^{t} h\left(w(\tau), \varepsilon^{-1} \tau, \varepsilon^{-2} \tau\right) d \tau & =\sum_{j=0}^{n-1}\left[\int_{t_{j}}^{t_{j+1}}\left(\int_{0}^{\varepsilon^{-2} \tau} f_{1}\left(w(\tau), \varepsilon^{-1} \tau, v\right) d v\right) \nabla_{x} f_{1}\left(w(\tau), \varepsilon^{-1} \tau, \varepsilon^{-2} \tau\right) d \tau\right. \\
& \left.-\int_{t_{j}}^{t_{j+1}} \int_{0}^{\varepsilon^{-2} \tau} \frac{\partial f_{1}}{\partial s_{1}}\left(w(\tau), \varepsilon^{-1} \tau, \tilde{v}\right) d \tilde{v} d \tau\right]+\mathscr{O}\left(\varepsilon^{2}\right) \\
& =\sum_{j=0}^{n-1}\left[\int_{j}^{j+1} \varepsilon^{2}\left(\int_{j}^{s} f_{1}\left(w\left(\varepsilon^{2} s\right), \varepsilon s, v\right) d v\right) \nabla_{x} f_{1}\left(w\left(\varepsilon^{2} s\right), \varepsilon s, s\right) d s\right. \\
& \left.-\int_{j}^{j+1} \varepsilon^{2}\left(\int_{0}^{s} \frac{\partial f_{1}}{\partial s_{1}}\left(w\left(\varepsilon^{2} s\right), \varepsilon s, \tilde{v}\right) d \tilde{v}\right) d s\right]+\mathscr{O}\left(\varepsilon^{2}\right)
\end{aligned}
$$

Changing variables back to $\tau=\varepsilon^{2} v$ in the first integral and $\tilde{v}=\tau$ in the second integral,

$$
\begin{align*}
\int_{0}^{t} h\left(w(\tau), \varepsilon^{-1} \tau, \varepsilon^{-2} \tau\right) d \tau & =\sum_{j=0}^{n-1}\left[\int_{j}^{j+1}\left(\int_{t_{j}}^{\varepsilon^{2} s} f_{1}\left(w\left(\varepsilon^{2} s\right), \varepsilon s, \varepsilon^{-2} \tau\right) d \tau\right) \nabla_{x} f_{1}\left(w\left(\varepsilon^{2} s\right), \varepsilon s, s\right) d s\right.  \tag{26}\\
& \left.-\varepsilon^{2} \int_{j}^{j+1} \int_{0}^{s} \frac{\partial f_{1}}{\partial s_{1}}\left(w\left(\varepsilon^{2} s\right), \varepsilon s, \tau\right) d \tau d s\right]+\mathscr{O}\left(\varepsilon^{2}\right)
\end{align*}
$$

We need to compare (25) and (26). Note that $w^{\prime}(t)$ is bounded independent of $\varepsilon$ and that $\tau \in[j, j+1]$.

$$
\left\{\begin{array}{l}
w\left(\varepsilon^{2} s\right)=w(\tau)+\mathscr{O}\left(\varepsilon^{2} s-\tau\right)=w(\tau)+\mathscr{O}\left(\varepsilon^{2}\right)  \tag{27}\\
w\left(\varepsilon^{2} s\right)=w\left(\varepsilon^{2} j\right)+\mathscr{O}\left(\varepsilon^{2} s-\varepsilon^{2} j\right)=w\left(t_{j}\right)+\mathscr{O}\left(\varepsilon^{2}\right)
\end{array}\right.
$$

and

$$
\left\{\begin{array}{l}
\varepsilon s-\varepsilon^{-1} \tau=\mathscr{O}(\varepsilon)  \tag{28}\\
\varepsilon s-\varepsilon^{-1} t_{j}=\mathscr{O}(\varepsilon)
\end{array}\right.
$$

Therefore, the first integration in (26) can be written as

$$
\begin{aligned}
& \int_{j}^{j+1}\left(\int_{t_{j}}^{\varepsilon^{2} s} f_{1}\left(w(\tau), \varepsilon^{-1} \tau, \varepsilon^{-2} \tau\right) d \tau\right) \nabla_{x} f_{1}\left(w\left(t_{j}\right), \varepsilon^{-1} t_{j}, s\right) d s+\mathscr{O}\left(\varepsilon^{3}\right) \\
= & \int_{0}^{1}\left(\int_{t_{j}}^{t_{j}+\varepsilon^{2} s} f_{1}\left(w(\tau), \varepsilon^{-1} \tau, \varepsilon^{-2} \tau\right) d \tau\right) \nabla_{x} f_{1}\left(w\left(t_{j}\right), \varepsilon^{-1} t_{j}, s\right) d s+\mathscr{O}\left(\varepsilon^{3}\right)
\end{aligned}
$$

which has the same form of the first term in (25) with an $\mathscr{O}\left(\varepsilon^{3}\right)$ error.
Now compare the second integrals. Applying integration by parts to the second term in (25) gives

$$
\begin{align*}
\int_{0}^{1} \varepsilon^{2} s \frac{\partial f_{1}}{\partial s_{1}}\left(x\left(t_{j}\right), \varepsilon^{-1} t_{j}, s\right) d s & =\left[\varepsilon^{2} s \int_{0}^{s} \frac{\partial f_{1}}{\partial s_{1}}\left(x\left(t_{j}\right), \varepsilon^{-1} t_{j}, \tau\right) d \tau\right]_{s=0}^{1}-\varepsilon^{2} \int_{0}^{1} \int_{0}^{s} \frac{\partial f_{1}}{\partial s_{1}}\left(x\left(t_{j}\right), \varepsilon^{-1} t_{j}, \tau\right) d \tau d s  \tag{29}\\
& =-\varepsilon^{2} \int_{0}^{1} \int_{0}^{s} \frac{\partial f_{1}}{\partial s_{1}}\left(x\left(t_{j}\right), \varepsilon^{-1} t_{j}, \tau\right) d \tau d s
\end{align*}
$$

On the other hand, in (26) using (27) and (28),

$$
\begin{align*}
-\varepsilon^{2} \int_{j}^{j+1} \int_{0}^{s} \frac{\partial f_{1}}{\partial s_{1}}\left(w\left(\varepsilon^{2} s\right), \varepsilon s, \tau\right) d \tau d s & =-\varepsilon^{2} \int_{j}^{j+1} \int_{0}^{s} \frac{\partial f_{1}}{\partial s_{1}}\left(w\left(t_{j}\right), \varepsilon^{-1} t_{j}, \tau\right) d \tau d s+\mathscr{O}\left(\varepsilon^{3}\right)  \tag{30}\\
& =-\varepsilon^{2} \int_{0}^{1} \int_{0}^{s} \frac{\partial f_{1}}{\partial s_{1}}\left(w\left(t_{j}\right), \varepsilon^{-1} t_{j}, \tau\right) d \tau d s+\mathscr{O}\left(\varepsilon^{3}\right)
\end{align*}
$$

Thus, it is shown that (29) and (30) are different up to an $\mathscr{O}\left(\varepsilon^{3}\right)$ error. Putting all estimates together and noting $n=\mathscr{O}\left(\varepsilon^{-2}\right)$, we conclude that

$$
\begin{equation*}
\left|\int_{0}^{t} \varepsilon^{-1} f_{1}\left(x(\tau), \varepsilon^{-1} \tau, \varepsilon^{-2} \tau\right)-h\left(w(\tau), \varepsilon^{-1} \tau, \varepsilon^{-2} \tau\right) d \tau\right| \leq L_{h} \int_{0}^{t}|x(\tau)-w(\tau)| d \tau+\mathscr{O}(\varepsilon) \tag{31}
\end{equation*}
$$

where $L_{h}$ is a Lipschitz constant for the function $h$. Substituting (31) into (22) yields the desired estimate.
The following lemma gives an estimate of how much $w(t)$ and $y(t)$ can be apart in $0 \leq t \leq T$.
Lemma 2. Let $w(t)$ and $y(t)$ be defined as above. Then, the following estimate holds

$$
\begin{equation*}
|w(t)-y(t)| \leq \varepsilon C e^{L_{G} t} \tag{32}
\end{equation*}
$$

for some constant $C>0$ which is independent of $\varepsilon$, where $L_{G}$ is a Lipschitz constant for $G$.
Proof. Consider $w^{\prime}-y^{\prime}$. Using (18) and (20) we have,

$$
\begin{aligned}
w^{\prime}-y^{\prime} & =f_{0}\left(w, \varepsilon^{-1} t, \varepsilon^{-2} t\right)-\left\langle f_{0}\right\rangle_{2}\left(y, \varepsilon^{-1} t\right)+h\left(w, \varepsilon^{-1} t, \varepsilon^{-2} t\right)-\langle h\rangle_{2}\left(y, \varepsilon^{-1} t\right) \\
& =\left[G\left(w, \varepsilon^{-1} t\right)-G\left(y, \varepsilon^{-1} t\right)\right]+z\left(w, \varepsilon^{-1} t, \varepsilon^{-2} t\right),
\end{aligned}
$$

where

$$
z\left(w, s_{1}, s_{2}\right)=\left[f_{0}\left(w, s_{1}, s_{2}\right)-\left\langle f_{0}\right\rangle_{2}\left(w, s_{1}\right)\right]+\left[h\left(w, s_{1}, s_{2}\right)-\langle h\rangle_{2}\left(w, s_{1}\right)\right] .
$$

Note that for fixed $w$ and $s_{1}, z\left(w, s_{1}, s_{2}\right)$ is 1-periodic in $s_{2}$ with a zero average, $\langle z\rangle_{2}=0$. Integrating to the time $t$

$$
w(t)-y(t)=\int_{0}^{t}\left[G\left(w(\tau), \varepsilon^{-1} \tau\right)-G\left(y(\tau), \varepsilon^{-1} \tau\right)\right] d \tau+r(t ; \varepsilon)
$$

where $r(t ; \varepsilon)=\int_{0}^{t} z\left(w(\tau), \varepsilon^{-1} \tau, \varepsilon^{-2} \tau\right) d \tau$. Taking absolute values

$$
|w(t)-y(t)| \leq \int_{0}^{t}\left|G\left(w(\tau), \varepsilon^{-1} \tau\right)-G\left(y(\tau), \varepsilon^{-1} \tau\right)\right| d \tau+|r(t ; \varepsilon)|
$$

Let $L_{G}$ denote a Lipschitz constant for $G$. Then,

$$
\begin{equation*}
|w(t)-y(t)| \leq L_{G} \int_{0}^{t}|w(\tau)-y(\tau)| d \tau+|r(t ; \varepsilon)| \tag{33}
\end{equation*}
$$

In order to evaluate $r(t ; \varepsilon)$, let $n$ denote the largest integer such that $\varepsilon^{2} n \leq t$. We have

$$
\begin{equation*}
r(t ; \varepsilon)=\sum_{j=0}^{n-1} \int_{\varepsilon^{2} j}^{\varepsilon^{2}(j+1)} z\left(w(\tau), \varepsilon^{-1} \tau, \varepsilon^{-2} \tau\right) d \tau+\int_{\varepsilon^{2} n}^{t} z\left(w(\tau), \varepsilon^{-1} \tau, \varepsilon^{-2} \tau\right) d \tau \tag{34}
\end{equation*}
$$

Since $z$ is bounded and $0 \leqq t-n \varepsilon^{2}<\varepsilon^{2}$, the last term is $\mathscr{O}\left(\varepsilon^{2}\right)$. Denoting $t_{j}=\varepsilon^{2} j$ and using the periodicity of $z$ in $s$, each term in the sum can be written as

$$
\varepsilon^{2} \int_{0}^{1} z\left(w\left(t_{j}+\varepsilon^{2} s\right), \varepsilon^{-1} t_{j}+\varepsilon s, s\right) d s .
$$

Expanding in $\varepsilon$ yields,

$$
\varepsilon^{2} \int_{0}^{1} z\left(w\left(t_{j}\right), \varepsilon^{-1} t_{j}, s\right) d s=\varepsilon^{2}\left[\langle z\rangle_{2}\left(w\left(t_{j}\right), \varepsilon^{-1} t_{j}\right)+\mathscr{O}(\varepsilon)\right]=\mathscr{O}\left(\varepsilon^{3}\right)
$$

Therefore,

$$
\begin{equation*}
r(t ; \varepsilon)=\mathscr{O}\left(n \varepsilon^{3}\right)=\mathscr{O}(\varepsilon) \tag{35}
\end{equation*}
$$

Substituting into (33) yields

$$
|w(t)-y(t)| \leq L_{G} \int_{0}^{t}|w(\tau)-y(\tau)| d \tau+\mathscr{O}(\varepsilon)
$$

Using Gronwall's inequality we conclude Lemma 2.

### 2.2 Proof of Theorem 1

Lemma 1 and 2 imply that

$$
\begin{equation*}
\sup _{0 \leq t \leq T}|x(t)-y(t)| \leq \varepsilon C \tag{36}
\end{equation*}
$$

where $y(t)$ is the solution of

$$
y^{\prime}=G\left(y, \varepsilon^{-1} t\right), \quad y(0)=x_{0},
$$

with

$$
\begin{align*}
G\left(y, s_{1}\right) & =\left\langle h\left(y, s_{1}, s_{2}\right)\right\rangle_{2}+\left\langle f_{0}\left(y, s_{1}, s_{2}\right)\right\rangle_{2}, \\
h\left(y, s_{1}, s_{2}\right) & =\left(\nabla_{x} f_{1}\right) \phi-\frac{\partial \phi}{\partial s_{1}},  \tag{37}\\
\phi\left(y, s_{1}, s_{2}\right) & =\int_{0}^{s_{2}} f_{1}\left(y, s_{1}, \tau\right) d \tau .
\end{align*}
$$

Since the RHS of $G\left(y, s_{1}\right)$ is periodic in $s_{1}$, we are in a position to apply the two-scale averaging theorem [40], integrating out the intermediate $\mathscr{O}(\varepsilon)$ time scale. Noting that since $\phi$ is 1 -periodic in $s_{1}, \partial \phi / \partial s_{1}$ is also periodic and has zero average with respect to $s_{1}$. This leads to an averaged equation for $y(t)$ :

$$
\xi^{\prime}=F(\xi), \quad \xi(0)=x_{0},
$$

where

$$
\begin{equation*}
F(\xi)=\left\langle\left(\nabla_{x} f_{1}\right) \phi\left(\xi, s_{1}, s_{2}\right)\right\rangle_{12}+\left\langle f_{0}\left(\xi, s_{1}, s_{2}\right)\right\rangle_{12} \tag{38}
\end{equation*}
$$

and we have

$$
\begin{equation*}
\sup _{0 \leq t \leq T}|y(t)-\xi(t)| \leq \varepsilon C . \tag{39}
\end{equation*}
$$

Combining (36) with (39) completes the proof of Theorem 1.
Remark 3. It is not difficult to generalize this result to systems with non-commensurate and widely separated frequencies.

## 3 Formal asymptotic expansions

In this section we analyze the multiscale structure of a system using the operator formalism firstly developed by Papanicolaou et al. as a formal asymptotic expansion for singular perturbations of SDEs [10, 38]. This approach has been further generalized and applied to many different problems, for example, in [35, 39, 44]. For the case of Hamiltonian dynamics, including integrable periodic systems, a rigorous version of formalism is presented in [21, 22] and references therein. The derivation in this section is formal; nonetheless, it is instructive and provides intuitive explanation for Theorem 1.

### 3.1 Stochastic differential equations

For completeness, we begin by reviewing singular perturbation expansions of SDEs. The resulting effective equations are then considered in the case in which all diffusion coefficients are formally set to zero.

Consider SDE systems of the form

$$
\left\{\begin{array}{l}
d x_{t}=\left[\varepsilon^{-1} f_{1}\left(x_{t}, y_{t}\right)+f_{0}\left(x_{t}, y_{t}\right)\right] d t  \tag{40}\\
d y_{t}=\varepsilon^{-2} a\left(x_{t}, y_{t}\right) d t+\varepsilon^{-1} \beta\left(x_{t}, y_{t}\right) d B_{t},
\end{array}\right.
$$

where $x_{t} \in \mathbb{R}^{d}$ is a slow process that evolves according to an ODE with a fast random coefficient $y_{t} \in \mathbb{R}^{n}$. We assumed that $\beta(x, y) \beta^{T}(x, y)$ is uniformly positive definite in $\mathbb{R}^{d} \times \mathbb{R}^{n}$. Furthermore, we assume that for fixed $x_{t}$ the dynamics of
$y_{t}$ is ergodic on an invariant set $\Sigma_{x}$ with a unique invariant measure $d \mu_{x}$. The expectations with respect to the invariant measures are denoted $\langle\cdot\rangle_{y}$, in which the $x$ dependence is suppressed. A necessary condition for $x$ to be slow is that $\left\langle f_{1}\right\rangle_{y}=0$. Otherwise, $x$ exhibits non-trivial dynamics on the $\mathscr{O}(\varepsilon)$ time scale. It is well known that under some suitable conditions on (40), $x(t)$ satisfies an effective SDE that is independent of $\varepsilon$,

$$
\begin{equation*}
d X_{t}=F\left(X_{t}\right) d t+b\left(X_{t}\right) d B_{t} \tag{41}
\end{equation*}
$$

where $F$ and $b$ can be expressed as averages with respect to the fast process $y(t)$ with $X$ fixed $[10,38,39,44]$. We begin with a brief overview of the relevant results of [38]. For details the reader is referred to [39] and references therein.

The backwards equation that governs the evolution of a probability density, $\phi$, of the initial conditions

$$
\left\{\begin{array}{l}
\partial_{t} u=\mathscr{L} u \\
u(0, x, y)=\varphi(x, y)
\end{array}\right.
$$

where $\mathscr{L}$, the generator of (40) can be written as

$$
\begin{align*}
& \mathscr{L}=\varepsilon^{-2} \mathscr{L}_{2}+\varepsilon^{-1} \mathscr{L}_{1}+\mathscr{L}_{0} \\
& \mathscr{L}_{0}=f_{0} \cdot \nabla_{x} \\
& \mathscr{L}_{1}=f_{1} \cdot \nabla_{x}  \tag{42}\\
& \mathscr{L}_{2}=a \cdot \nabla_{y}+\frac{1}{2} \beta \beta^{T}: \nabla_{y} \nabla_{y}^{T}
\end{align*}
$$

where $A: B$ denotes formally the trace of the matrix $A B^{T}$. Next, consider a formal asymptotic expansion of $u$ in $\varepsilon$

$$
u(t, x, y)=u_{0}(x)+\varepsilon u_{1}+\varepsilon^{2} u_{2}+\ldots,
$$

where we assumed that the leading order term $u_{0}$ only depends on the slow process. Substituting into the backwards equation yields

$$
\begin{align*}
& \mathscr{L}_{2} u_{0}=0  \tag{43}\\
& \mathscr{L}_{2} u_{1}=-\mathscr{L}_{1} u_{0}  \tag{44}\\
& \mathscr{L}_{2} u_{2}=\partial_{t} u_{0}-\mathscr{L}_{1} u_{1}-\mathscr{L}_{0} u_{0} . \tag{45}
\end{align*}
$$

The leading order equation (43) is automatically satisfied since $u_{0}$ only depends on $x$. Let $\mathscr{L}_{2}^{*}$ be the $L^{2}$ adjoint of $\mathscr{L}_{2}$, and assume that the Null space of $\mathscr{L}_{2}^{*}$ is a one dimensional subspace, spanned $\rho_{x}$, the density of the invariant measure $d \mu_{x} .{ }^{2}$ Applying the Fredholm alternative, equation (44) has a solution if $f_{1} \nabla_{x} u_{0}$ is perpendicular to the Null space of $\mathscr{L}_{2}^{*}$. In other words, the projection of $f_{1} \nabla_{x} u_{0}$ on Null $\mathscr{L}_{2}^{*}$ should vanish. This projection amounts to averaging with respect to the invariant measure of $y$ (at fixed $x$ ), which is also equivalent to taking the standard $L_{2}$ inner product in $\Sigma_{x}$ with $\rho$. This yields the condition

$$
\left\langle f_{1}(x, y)\right\rangle_{y}=0
$$

which implies the reasonable requirement that the average of $f_{1}$ vanishes. Otherwise, $x$ oscillates with large amplitudes on the $\varepsilon$ time scale and thus cannot be approximated by a slow variable. Then, (44) has a unique solution such that $\left\langle u_{1}\right\rangle_{y}=0$. We formally write

$$
u_{1}=-\mathscr{L}_{2}^{-1}\left[\mathscr{L}_{1} u_{0}\right]
$$

Applying again the Fredholm alternative, equation (45) also has a solution if the RHS is perpendicular to the Null space of $L_{2}^{*}$, i.e.

$$
\left\langle\partial_{t} u_{0}-\mathscr{L}_{1} u_{1}-\mathscr{L}_{0} u_{0}\right\rangle_{y}=0
$$

Substituting in $\mathscr{L}_{0}, \mathscr{L}_{1}$ and $u_{1}$ we obtain

$$
\begin{equation*}
\partial_{t} u_{0}=F(x) \nabla_{x} u_{0}+B(x): \nabla_{x} \nabla_{x} u_{0}, \tag{46}
\end{equation*}
$$

[^2]where
\[

$$
\begin{align*}
F(x) & =\left\langle-f_{1} \cdot \nabla_{x} \mathscr{L}_{2}^{-1}\left[f_{1}\right]+f_{0}\right\rangle_{y}  \tag{47}\\
B(x) & =\left\langle-f_{1} \mathscr{L}_{2}^{-1}\left[f_{1}\right]\right\rangle_{y}
\end{align*}
$$
\]

We identify (46) as the backwards equation associated with the effective $\operatorname{SDE}(41)$ and $2 B(x)=b(x) b(x)^{T}$. If $\beta \beta^{T}$ is uniformly positive, then inverting $\mathscr{L}_{2}$ amounts to solving a uniformly elliptic cell problem

$$
\begin{equation*}
\left[a \cdot \nabla_{y}+\frac{1}{2} \beta \beta^{T}: \nabla_{y} \nabla_{y}\right] r(x, y)=f_{1}(x, y) \tag{48}
\end{equation*}
$$

with appropriate boundary conditions, e.g., periodic on a torus. The equation indeed has a unique solution. See [39] for further details.

### 3.2 Periodic ODEs

We now formally set $\beta=0$ and consider a case in which for fixed $x$ the dynamics of $y$ is periodic. Hence, we can think of $y$ as an fast oscillator. Assume that $a(x, y)$ is uniformly positive in $\mathbb{R}^{d} \times \mathbb{R}^{n}, \inf _{x \in \mathbb{R}^{d}, y \in \mathbb{R}^{n}}|a(x, y)|>0$. Loosely speaking, this means that the system remains highly oscillatory at all times. Therefore, for fixed $x$ trajectories are closed loops and the invariant set $\Sigma_{x}$ is a one-dimensional manifold in $\mathbb{R}^{d}$ that depends on the initial $y$ and on $x$. The invariant measure $d \mu_{x}$, which is supported on $\Sigma_{x}$ is absolutely continuous with respect to the Lebesgue measure on $\Sigma_{x}$. The ergodic assumption holds with respect to the same manifold and measure. A rigorous treatment of this problem is beyond the scope of this manuscripts. For two-scale systems the reader is referred to [21,22] and references therein.

Instead, we continue formally trying to identify the effective equation for $x_{t}$. The main difference between the periodic and the stochastic cases can be seen from two complementary points of view. First, as a dynamical system, the fast process in the $\operatorname{SDE}(40)$ is mixing (for fixed $x$ ). This is no longer the case with periodic systems which are ergodic but not mixing. Second, as a homogenization problem, since $\beta \beta^{T}=0$ the cell problem (48) is no longer elliptic.

For fixed $x$, consider the periodic solution of

$$
\frac{d \zeta}{d t}=a(x, \zeta(t))
$$

with a suitable initial condition on $\Sigma_{x}$. The period is denoted $\tau_{x} . \zeta(t)$ transverses the exact periodic trajectory of $y(t)$ and with the correct invariant measure. Therefore, averages with respect to $d \mu_{x}$ can be written as the time average over a single period of $\zeta$

$$
\langle h(x, y)\rangle_{y}=\frac{1}{\tau_{x}} \int_{0}^{\tau_{x}} h(x, \zeta(t)) d t=\frac{1}{\tau_{x}} \int_{\Sigma_{x}} h(x, \zeta(t)) \frac{1}{a(x, \zeta(t))} d \zeta
$$

Furthermore, recall the hierarchy of operators (42). Substituting $\beta=0, \mathscr{L}_{2}$ takes the form

$$
\mathscr{L}_{2}=a(x, y) \cdot \nabla_{y} .
$$

We note that, for any $y \in \Sigma_{x}, \mathscr{L}_{2}$ is the directional derivative of $\zeta$ along the tangent direction to $\Sigma_{x}$ and that $\left|\mathscr{L}_{2} \zeta\right|$ is inversely proportional to the density of $d \mu_{x}$

$$
\mathscr{L}_{2} h(x, y)=a(x, y) \cdot \nabla_{y} h(x, y)=\left.\frac{d}{d t} h(x, \zeta(t))\right|_{\zeta(t)=y}
$$

This implies that the inverse of $\mathscr{L}_{2}$ can be described in terms of integration with respect to time along the trajectory of $\zeta(t)$

$$
H(x, y)=\mathscr{L}_{2}^{-1} h(x, y)=\int_{0}^{t(y)} h(x, \zeta(s)) d s+C(x)
$$

where $t(y)$ is the unique time in which $\zeta(t)=y$ within one period of $\zeta$. Following the Fredholm alternative, we pick the unique inverse that is perpendicular to Null $\mathscr{L}_{2}^{*}$, i.e., we require

$$
\langle H(x, y)\rangle_{y}=\frac{1}{\tau_{x}} \int_{0}^{\tau_{x}} H(x, \zeta(t)) d t=0
$$

This fixes the constant

$$
C(x)=-\frac{1}{\tau_{x}} \int_{0}^{\tau_{x}} \int_{0}^{t} H(x, \zeta(s)) d s d t
$$

In particular, we recognize that $\mathscr{L}_{2}^{-1}\left[f_{1}\right]=\phi$ as given by

$$
\begin{equation*}
\phi(x, y)=\int_{0}^{y} f_{1}(x, \tau) d \tau-\left\langle\int_{0}^{y} f_{1}(x, \tau) d \tau\right\rangle_{y} \tag{49}
\end{equation*}
$$

so that $\langle\phi\rangle_{y}=0$. Substituting into (47) yields

$$
\begin{align*}
& F(x)=\left\langle-f_{1} \cdot \nabla_{x} \phi+f_{0}\right\rangle_{y}  \tag{50}\\
& B(x)=\left\langle-f_{1} \phi\right\rangle_{y} .
\end{align*}
$$

In order to identify (50) with the averaged equation (37), we need a simple Lemma.
Lemma 3. Let $h(s)$ denote an S-periodic function with zero average and let $H(s)$ be an anti-derivative of $h, H^{\prime}=h$. Then,

$$
\begin{equation*}
\int_{0}^{S} h(s) H(s) d s=0 \tag{51}
\end{equation*}
$$

Proof. First, we note that since $h(s)$ has zero average, $\int_{0}^{S} h(\tau) d \tau=0$, its anti-derivative is also $S$-periodic

$$
H(S+s)=H(0)+\int_{0}^{S+s} h(\tau) d \tau=H(s)
$$

Then, using integration by parts

$$
\int_{0}^{S} h(s) H(s) d s=\left[H^{2}(s)\right]_{0}^{S}-\int_{0}^{S} H(s) h(s) d s
$$

The first term on the right vanishes, which proves (51).
Since $\phi$ is the anti-derivative of $f_{1}$, an immediate consequence is that $B(x)=-\left\langle f_{1}(x, y) \phi(x, y)\right\rangle_{y}=0$. Hence, the effective dynamics of $x_{t}$ is deterministic. More precisely, the variance of the stochastic perturbation is of order $\varepsilon$. Furthermore,

$$
0=\nabla_{x}\left\langle f_{1}(x, y) \phi(x, y)\right\rangle_{y}=\left\langle\nabla_{x} f_{1}(x, y) \phi(x, y)\right\rangle_{y}+\left\langle f_{1}(x, y) \nabla_{x} \phi(x, y)\right\rangle_{y}
$$

We conclude that the effective drift and diffusion coefficient can be written as

$$
\begin{align*}
& F(x)=\left\langle\left(\nabla_{x} f_{1}\right) \cdot \phi+f_{0}\right\rangle_{y}  \tag{52}\\
& B(x)=0 .
\end{align*}
$$

Thus, we obtain the consistent form of the effective drift as proven by Theorem 1. Note that since (40) contains only two scales $\mathscr{O}(1)$ and $\mathscr{O}\left(\varepsilon^{2}\right), \frac{\partial \phi}{\partial s_{1}}=0$ in (37).

## 4 Iterated averaging with multiple kernels

The goal of this section is to generalize the framework of averaging kernels studied in [2, 3, 4, 5, 20] to include three or more time scales. In particular, we prove that averaging of three scale oscillatory problems can be approximated via convolution with respect to kernels with known support, differentiability properties and moments. Let $\mathbb{K}^{p, q}$ denote the space of normalized $C^{q}$ functions, supported on $[-1,1]$ that have $p$ vanishing moments, i.e.,

$$
\int_{[-1,1]} K(t) t^{r} d t= \begin{cases}1, & r=0  \tag{53}\\ 0, & 1 \leq r \leq p\end{cases}
$$

Recall that for $\eta>0, K_{\eta}(t)$ denotes a rescaling of $K$ as $K_{\eta}(t)=\eta^{-1} K\left(\eta^{-1} t\right)$.
We will have an error estimate for approximating the double average $\langle f\rangle_{12}$ with two convolutions $K_{\eta_{1}} *\left[K_{\eta_{2}} * f\right]$. Applied to the RHS of (6), the convolution approximates $F(\xi)$ of Theorem 1. The proposition below establishes the accuracy of iterated averaging with multiple kernels under the scaling condition $\varepsilon^{2} \ll \eta_{2} \ll \varepsilon \ll \eta_{1} \ll 1$.
Proposition 1 Let $K \in \mathbb{K}^{p, q}$ and $f_{0}, f_{1} \in C^{\max \{q, p+1\}}$. Then, for $K \in \mathbb{K}^{p, q}$, there exists a constant $C>0$ such that

$$
\begin{equation*}
\left|K_{\eta_{1}} *\left\{K_{\eta_{2}} *\left(\varepsilon^{-1} f_{1}+f_{0}\right)(\cdot)\right\}(t)-F(\xi(t))\right| \leq C\left(\frac{\varepsilon^{2 q-1}}{\eta_{2}^{q}}+\frac{\varepsilon^{q}}{\eta_{1}^{q}}+\frac{\eta_{2}^{p+1}}{\varepsilon^{p+2}}+\eta_{1}^{p+1}\right) \max _{j=0 \ldots q}\left\|K^{(j)}\right\|_{1} \tag{54}
\end{equation*}
$$

for some constant $C>0$, where $F$ is given by (37) as in Theorem $1, K^{(j)}(\cdot)$ denotes the $j$-th derivative of $K$ and $\|\cdot\|_{1}$ is the $L_{1}$ norm.

Throughout, $C$ denotes a generic positive constant whose value may change between expressions.

### 4.1 Estimation of the effective force

Let $f: \mathbb{R}^{n+1} \rightarrow \mathbb{R}$ denote a scalar function of the vector argument $\left(x, s_{1}, s_{2}, \cdots, s_{n}\right)$ and 1-periodic in $s_{1}, \cdots, s_{n}$. To be consistent with previous notation, averaging with respect to $s_{k}$ is respectively denoted by

$$
\begin{equation*}
\langle f\rangle_{k}=\int_{0}^{1} f\left(x, \cdots, s_{k}, \cdots\right) d s_{k}, \quad k=1,2, \cdots, n \tag{55}
\end{equation*}
$$

Motivated by the averaging techniques in [2], we approximate $\langle f\rangle_{n}$ using a kernel. First let us prove the following lemma.

Lemma 4. If $f\left(t, s_{1}, s_{2}, \cdots, s_{n}\right)=a\left(t, s_{1}, s_{2}, \cdots, s_{n-1}\right) b\left(s_{n}\right)$, where $n \in \mathbb{N}, b\left(s_{n}+1\right)=b\left(s_{n}\right), \int_{0}^{1} b\left(s_{n}\right) d s_{n}=0$, $a \in$ $C^{q}\left(\mathbb{R}^{n}\right)$ and $\max _{0 \leq j \leq q}\left\|a^{(j)}\right\|_{\infty} \leq M$, then for any $K \in \mathbb{K}^{p, q}$ and $\eta=\mathscr{O}\left(\varepsilon^{k}\right), k>n-1$,

$$
\left|K_{\eta} * f\left(\cdot, \varepsilon^{-1} \cdot, \varepsilon^{-2} \cdot, \cdots, \varepsilon^{-n} \cdot\right)\right| \leq C M\left(\frac{\varepsilon^{n}}{\eta}\right)^{q} \max _{j=0 \ldots q}\left\|K^{(j)}\right\|_{1}
$$

Proof. Let $\tilde{K}_{\eta}(x, y)=K_{\eta}(x-y) a\left(y, \varepsilon^{-1} y, \varepsilon^{-2} y, \cdots, \varepsilon^{-n+1} y\right)$.

$$
\begin{aligned}
\left|K_{\eta} * f\right| & =\left|\int \tilde{K}_{\eta}(t, s) b\left(\varepsilon^{-n} s\right) d s\right| \\
& \leq \varepsilon^{n q} \int\left|\partial_{y}^{q} \tilde{K}_{\eta}(t, s) b^{[q]}\left(\varepsilon^{-n} s\right)\right| d s
\end{aligned}
$$

where $b^{[j]}(t)=\int_{0}^{t} b^{[j-1]}(s) d s-\int_{0}^{1} \int_{0}^{t} b^{[j-1]}(s) d s d t$ and $\left\|b^{[j]}(t)\right\|_{\infty} \leq\|b\|_{\infty}$. For the kernel part,

$$
\begin{aligned}
& \int\left|\partial_{y}^{q} \tilde{K}_{\eta}(t, s)\right| d s=\int \left\lvert\, \sum_{k_{1}+k_{2}+\cdots+k_{n+1}=q}\binom{q}{k_{1}, k_{2}, \cdots, k_{n+1}}\right. \\
\cdot & \left.\left(-\eta^{-1}\right)^{k_{1}} \cdot K_{\eta}^{\left(k_{1}\right)}(t-s) \cdot \prod_{i=1}^{n} \frac{1}{\varepsilon^{k_{i+1}(i-1)}} \partial_{i}^{k_{i+1}} a\left(s, \varepsilon^{-1} s, \varepsilon^{-2} s, \cdots, \varepsilon^{-n+1} s\right) \right\rvert\, d s \\
\leq & \frac{C}{\eta^{q}} \max _{j=0 \ldots q}\left\|a^{(j)}\right\|_{\infty} \max _{j=0 \ldots q}\left\|K^{(j)}\right\|_{1}
\end{aligned}
$$

where $\binom{q}{k_{1}, k_{2}, \cdots, k_{n+1}}$ is the multinomial coefficient and defined by $\frac{q!}{k_{1}!k_{2}!\cdots k_{n+1}!}$.

Lemma 5. For $r=\left(r_{1}, r_{2}, \cdots, r_{n}\right) \in \mathbb{Z}_{+}^{n}$ an ordered $n$-tuple of nonnegative integers, assume $\partial^{r} f\left(t, s_{1}, \cdots, s_{n}\right)$ is continuous and bounded by $C_{f}$ for $r=0, \cdots, \sigma$, and $\sigma \geq 1 .{ }^{3}$ Then, for any $K \in \mathbb{K}^{p, q}$ and $\eta=\mathscr{O}\left(\varepsilon^{k}\right)$ with $k>n-1$, there exists $C>0$ such that

$$
\begin{equation*}
\left|K_{\eta} *\left(f\left(\cdot, \varepsilon^{-1} \cdot, \varepsilon^{-2} \cdot, \cdots, \varepsilon^{-n} \cdot\right)-\langle f\rangle_{n}\right)\right| \leq C\left(\frac{\varepsilon^{n q}}{\eta^{q}}+\frac{\eta^{\sigma}}{\varepsilon^{(n-1) \sigma}}\right) \max _{j=0 \ldots q}\left\|K^{(j)}\right\|_{1} . \tag{56}
\end{equation*}
$$

Proof. Let

$$
g\left(t, s_{1}, s_{2}, \cdots, s_{n}\right)=f\left(t, s_{1}, s_{2}, \cdots, s_{n}\right)-\langle f\rangle_{n}
$$

$g$ is 1-periodic with respect to $s_{i}$ and $\partial^{r} g\left(t, s_{1}, s_{2}, \cdots, s_{n}\right)$ are continuous and bounded for $r \in \mathbb{Z}_{+}^{n},|r|=0, \cdots, \sigma$. In considering

$$
K_{\eta} * g=\int_{t-\eta}^{t+\eta} K_{\eta}\left(t-s_{n}\right) \cdot g\left(s_{n}, \varepsilon^{-1} s_{n}, \varepsilon^{-2} s_{n}, \cdots, \varepsilon^{-n} s_{n}\right) d s_{n}
$$

we expand $g\left(s_{n}, \varepsilon^{-1} s_{n}, \varepsilon^{-2} s_{n}, \cdots, \varepsilon^{-(n-1)} s_{n}, t\right)$ around $s_{n}=0$ while keeping $t$ fixed. We denote $\partial_{x_{i}} g$ by the partial derivative with respect to the $i$-th component of $g$ respectively.

$$
\begin{aligned}
g\left(s_{n}, \varepsilon^{-1} s_{n}, \varepsilon^{-2} s_{n}, \cdots, \varepsilon^{-n} s_{n}\right) & =\sum_{j=0}^{\sigma-1} \frac{1}{j!}\left[\sum_{i=1}^{n} \frac{s_{n}}{\varepsilon^{i-1}} \partial_{x_{i}}\right]^{j} g\left(0, \cdots, 0, \varepsilon^{-n} s_{n}\right) \\
& +\frac{1}{\sigma!}\left[\sum_{i=1}^{n} \frac{s_{n}}{\varepsilon^{i-1}} \partial_{x_{i}}\right]^{\sigma} g\left(\mu_{1}, \varepsilon^{-1} \mu_{2}, \cdots, \varepsilon^{-n+1} \mu_{n}, \varepsilon^{-n} s_{n}\right)
\end{aligned}
$$

where $\left(\mu_{1}, \mu_{2}, \cdots, \mu_{n}\right)$ is in the open line segment joining $\mathbf{0}$ and $\left(s_{n}, \varepsilon^{-1} s_{n}, \cdots, \varepsilon^{-n+1} s_{n}\right)$ in $\mathbb{R}^{n}$. The key idea consists of in writing the expansion as a sum of two parts. Without loss of generality, we set $t=0$ and write $K_{\eta} * g$ as $I_{1}+I_{2}$ where

$$
\begin{aligned}
& I_{1}=\sum_{j=0}^{\sigma-1} \int_{-\eta}^{\eta} K_{\eta}\left(-s_{n}\right) \cdot \frac{1}{j!}\left[\sum_{i=1}^{n} \frac{s_{n}}{\varepsilon^{i-1}} \partial_{x_{i}}\right]^{j} g\left(0, \cdots, 0, \varepsilon^{-n} s_{n}\right) d s_{n} \\
& I_{2}=\int_{-\eta}^{\eta} K_{\eta}\left(-s_{n}\right) \cdot \frac{1}{\sigma!}\left[\sum_{i=1}^{n} \frac{s_{n}}{\varepsilon^{i-1}} \partial_{x_{i}}\right]^{\sigma} g\left(\mu_{1}, \varepsilon^{-1} \mu_{2}, \cdots, \varepsilon^{-n+1} \mu_{n}, \varepsilon^{-n} s_{n}\right) d s_{n}
\end{aligned}
$$

By using Lemma $4, I_{1}$ is estimated by

$$
\left|I_{1}\right| \leq C\left(\frac{\varepsilon^{n}}{\eta}\right)^{q} \cdot \max _{j=0 \ldots q}\left\|K^{(j)}\right\|_{1} \cdot \sum_{j=0}^{\sigma-1} \frac{1}{j!}\left[\sum_{i=1}^{n} \frac{\eta}{\varepsilon^{i-1}}\right]^{j} .
$$

Finding the leading order term in the summation, $I_{2}$ is estimated by

$$
\left|I_{2}\right| \leq C\left(\frac{\eta}{\varepsilon^{n-1}}\right)^{\sigma} \sup _{t, s_{1}, s_{2}, \cdots, s_{n}} \sup _{|r|=\sigma}\left|\partial^{r} f\left(t, s_{1}, \cdots, s_{n}\right)\right| \cdot\|K\|_{1}
$$

Putting these estimates together, we find that there exists a positive constant $C$ such that

$$
\left|K_{\eta} * g\right| \leq C\left(\frac{\varepsilon^{n q}}{\eta^{q}}+\frac{\eta^{\sigma}}{\varepsilon^{(n-1) \sigma}}\right) \max _{j=0 \ldots q}\left\|K^{(j)}\right\|_{1}
$$

We now compare iterated averaging of (6) with the averaged equation (9). Before we move on to the next step, we simplify our notation of the forces by writing $f\left(t, \varepsilon^{-1} t, \varepsilon^{-2} t\right)$ instead of $f\left(x(t), \varepsilon^{-1} t, \varepsilon^{-2} t\right)$. This is possible since, after solving (6) at the 1 st tier, $x(t)$ is known up to a prescribed accuracy $\Delta$. Theorem 2 shows that the error between $x$ and $x_{1}$ (solution of the 1 st tier) is $\mathscr{O}(\Delta)$.

[^3]Theorem 2. Let $x_{1}(t)$ denote an approximation of $x(t)$ in the 1 st tier using a two-scale HMM in the time interval $t \in\left[t_{n}-\eta_{1}, t_{n}+\eta_{1}\right]$. Given $0<\varepsilon<\varepsilon_{0}$ and a prescribed accuracy $\Delta_{1}$, there exists $C>0$ such that

$$
\begin{equation*}
\sup _{t \in\left[t_{n}-\eta_{1}, t_{n}+\eta_{1}\right]}\left|x(t)-x_{1}(t)\right| \leq C \Delta_{1} . \tag{57}
\end{equation*}
$$

Proof. By considering 1st and 2nd tier as the two-scale HMM solvers, we generalize the error analysis discussed in $[3,6,20]$. We denote the order of accuracy, step size and length of integration in $i$-th tier by $m_{i}, h_{i}$ and $\eta_{i}$ respectively.

At the 1 st tier, the global accuracy of integrating the original ODE (6) to time $\eta_{1}(\ll \varepsilon)$ is given by

$$
C \max \left\{\frac{\eta_{1} h_{1}^{m_{1}}}{\varepsilon^{m_{1}+1}}, \frac{\eta_{1} \eta_{2} h_{2}^{m_{2}}}{h_{1} \varepsilon^{2 m_{2}+1}}, \frac{\eta_{1} \varepsilon^{2 q}}{h_{1} \eta_{2}^{q+1}}\right\}
$$

for some $C>0$. The errors from each evaluation at the 2 nd tier accumulate by taking $h_{1}^{-1} \eta_{1}$ steps. Thus we can balance the required accuracy $\Delta_{1}$ with different sources of errors. Note that the maximal possible accuracy is $\Delta_{1}=\varepsilon^{2}$ since this is the error introduced by the averaged equation on the $\mathscr{O}(\varepsilon)$ time.

Next, by iterated use of Lemma 5, we show that $\left\langle f\left(t, s_{1}, s_{2}\right)\right\rangle_{12}$ is well approximated by $K_{\eta_{1}} *\left[K_{\eta_{2}} * f\right](t)$.
Lemma 6. Let $f_{0}\left(t, s_{1}, s_{2}\right)$ be 1-periodic in $s_{1}$ and $s_{2}$. For $r=\left(r_{1}, r_{2}\right) \in \mathbb{Z}_{+}^{2}$, assume that $\partial^{r} f_{0}\left(t, s_{1}, s_{2}\right)$ is continuous and bounded for $|r|=0, \cdots, \sigma$, and $\sigma \geq 1$. Then, for any $K \in \mathbb{K}^{p, q}, \varepsilon^{2} \ll \eta_{2} \ll \varepsilon \ll \eta_{1} \ll 1$, there exists a constant $C>0$ such that

$$
\left|K_{\eta_{1}} *\left[K_{\eta_{2}} *\left(f_{0}\left(\cdot, \varepsilon^{-1} \cdot, \varepsilon^{-2} \cdot\right)-\left\langle f_{0}\left(t, s_{1}, s_{2}\right)\right\rangle_{12}\right)(\cdot)\right](t)\right| \leq C\left(\frac{\varepsilon^{2 q}}{\eta_{2}^{q}}+\frac{\varepsilon^{q}}{\eta_{1}^{q}}+\frac{\eta_{2}^{\sigma}}{\varepsilon^{\sigma}}+\eta_{1}^{\sigma}\right) \max _{j=0 \ldots q}\left\|K^{(j)}\right\|_{1} .
$$

Proof. Let

$$
\begin{equation*}
g\left(t, s_{1}, s_{2}\right)=f_{0}\left(t, s_{1}, s_{2}\right)-\left\langle f_{0}\right\rangle_{12}(t) . \tag{58}
\end{equation*}
$$

Note that $\langle g\rangle_{12}=\int_{0}^{1} \int_{0}^{1}\left(f_{0}\left(t, s_{1}, s_{2}\right)-\left\langle f_{0}\right\rangle_{12}\right) d s_{2} d s_{1}=0$ and that $g\left(x, s_{1}, s_{2}\right)$ is 1-periodic with respect to $s_{1}, s_{2}$ and $\partial^{r} g\left(x, s_{1}, s_{2}\right)$ is continuous and bounded for $r \in \mathbb{Z}_{+}^{2},|r|=0, \cdots, \sigma$. Iterated convolution with two kernels yields

$$
K_{\eta_{1}} *\left[\left(K_{\eta_{2}} * g\left(\cdot, \varepsilon^{-1} \cdot, \varepsilon^{-2} \cdot\right)\right)(\cdot)\right](t)=\int_{t-\eta_{1}}^{t+\eta_{1}} K_{\eta_{1}}\left(t-s_{1}\right) \underbrace{\left[\int_{s_{1}-\eta_{2}}^{s_{1}+\eta_{2}} K_{\eta_{2}}\left(s_{1}-s_{2}\right) \cdot g\left(s_{2}, \varepsilon^{-1} s_{2}, \varepsilon^{-2} s_{2}\right) d s_{2}\right]}_{=I_{1}\left(s_{1}\right)} d s_{1}
$$

We iterate the argument in Lemma 5. First, there exists $C>0$ such that

$$
\begin{equation*}
\left|I_{1}\left(s_{1}\right)-\langle g\rangle_{2}\left(s_{1}, \varepsilon^{-1} s_{1}\right)\right| \leq C\left(\frac{\varepsilon^{2 q}}{\eta_{2}^{q}}+\frac{\eta_{2}^{\sigma}}{\varepsilon^{\sigma}}\right) \max _{j=0 \ldots q}\left\|K^{(j)}\right\|_{1} \tag{59}
\end{equation*}
$$

Recalling that $\langle g\rangle_{12}=0$,

$$
\begin{equation*}
I_{1}=I_{1}-\langle g\rangle_{2}\left(s_{1}, \varepsilon^{-1} s_{1}\right)+\langle g\rangle_{2}\left(s_{1}, \varepsilon^{-1} s_{1}\right)-\langle g\rangle_{12} \tag{60}
\end{equation*}
$$

Thus

$$
\begin{align*}
& \left|\int_{t-\eta_{1}}^{t+\eta_{1}} K_{\eta_{1}}\left(t-s_{1}\right) \cdot I_{1} d s_{1}\right| \\
= & \left|\int_{t-\eta_{1}}^{t+\eta_{1}} K_{\eta_{1}}\left(t-s_{1}\right)\left(I_{1}-\langle g\rangle_{2}\left(s_{1}, \varepsilon^{-1} s_{1}\right)+\langle g\rangle_{2}\left(s_{1}, \varepsilon^{-1} s_{1}\right)-\langle g\rangle_{12}\right) d s_{1}\right|  \tag{61}\\
\leq & C\left(\frac{\varepsilon^{2 q}}{\eta_{2}^{q}}+\frac{\eta_{2}^{\sigma}}{\varepsilon^{\sigma}}\right) \max _{j=0 \ldots q}\left\|K^{(j)}\right\|_{1}+\left|\int_{t-\eta_{1}}^{t+\eta_{1}} K_{\eta_{1}}\left(t-s_{1}\right)\left(\langle g\rangle_{2}\left(s_{1}, \varepsilon^{-1} s_{1}\right)-\langle g\rangle_{12}\right) d s_{1}\right| .
\end{align*}
$$

Second, define $\hat{g}\left(\hat{s}_{1}, s_{1}\right)=\langle g\rangle_{2}\left(\hat{s}_{1}, \varepsilon^{-1} s_{1}\right)-\langle g\rangle_{12}$. Hence, $\hat{g}\left(\hat{s}_{1}, s_{1}\right)$ is 1-periodic in the second variable and the average over $s_{1}$ is zero. A second application of Lemma 5 yields existence of $C_{2}>0$ such that

$$
\begin{equation*}
\left|\int_{t-\eta_{1}}^{t+\eta_{1}} K_{\eta_{1}}\left(t-s_{1}\right) \hat{g}\left(\hat{s}_{1}, s_{1}\right) d s_{1}\right| \leq C_{2}\left(\frac{\varepsilon^{q}}{\eta_{1}^{q}}+\eta_{1}^{\sigma}\right) \max _{j=0 \ldots q}\left\|K^{(j)}\right\|_{1} . \tag{62}
\end{equation*}
$$

Hence, we can find a positive constant $C$ such that

$$
\left|K_{\eta_{1}} *\left[K_{\eta_{2}} *\left(f_{0}\left(\cdot, \varepsilon^{-1} \cdot, \varepsilon^{-2} \cdot\right)-\left\langle f_{0}\left(t, s_{1}, s_{2}\right)\right\rangle_{12}\right)(\cdot)\right](t)\right| \leq C\left(\frac{\varepsilon^{2 q}}{\eta_{2}^{q}}+\frac{\varepsilon^{q}}{\eta_{1}^{q}}+\frac{\eta_{2}^{\sigma}}{\varepsilon^{\sigma}}+\eta_{1}^{\sigma}\right) \max _{j=0 \ldots q}\left\|K^{(j)}\right\|_{1} .
$$

This concludes the proof of the theorem.
Recall the three scale problem (6) and its averaged equation (9):

$$
\begin{aligned}
& x^{\prime}=\varepsilon^{-1} f_{1}\left(x, \varepsilon^{-1} t, \varepsilon^{-2} t\right)+f_{0}\left(x, \varepsilon^{-1} t, \varepsilon^{-2} t\right), \quad x(0)=x_{0} \\
& \xi^{\prime}=\left\langle\left(\nabla_{x} f_{1}\right) \phi\left(\xi, s_{1}, s_{2}\right)\right\rangle_{12}+\left\langle f_{0}\left(\xi, s_{1}, s_{2}\right)\right\rangle_{12}, \quad \xi(0)=x_{0}
\end{aligned}
$$

Key to the following theorem is the vanishing of $\left\langle f_{1}\left(x(t), s_{1}, s_{2}\right)\right\rangle_{2}$.
Theorem 3. Let $f_{1}\left(x(t), s_{1}, s_{2}\right)$ be 1-periodic in $s_{1}$ and $s_{2}$, and have a zero average with respect to $s_{2}$. For $r=\left(r_{1}, r_{2}\right) \in$ $\mathbb{Z}_{+}^{2}$, assume $\partial^{r} f_{1}\left(x(t), s_{1}, s_{2}\right)$ are continuous and bounded for $|r|=0, \cdots, \sigma+1$, and $\sigma \geq 1$. Then, for any $K \in \mathbb{K}^{p, q}$ and suitable choice of $\eta_{1}, \eta_{2}$ such that $\varepsilon^{2} \ll \eta_{2} \ll \varepsilon \ll \eta_{1} \ll 1$, the following estimate holds.

$$
\begin{equation*}
\left|K_{\eta_{1}} *\left\{K_{\eta_{2}} *\left(\varepsilon^{-1} f_{1}\left(x(\cdot), \varepsilon^{-1} \cdot, \varepsilon^{-2} \cdot\right)-\langle h\rangle_{12}(\xi)\right)(\cdot)\right\}(t)\right| \leq C\left(\frac{\varepsilon^{2 q-1}}{\eta_{2 q}}+\frac{\varepsilon^{q}}{\eta_{1}^{q}}+\frac{\eta_{2}^{\sigma}}{\varepsilon^{\sigma+1}}+\eta_{1}^{\sigma-1}\right) \max _{j=0 \ldots q}\left\|K^{(j)}\right\|_{1} \tag{63}
\end{equation*}
$$

where

$$
\begin{aligned}
h\left(\xi, s_{1}, s_{2}\right) & =\left(\nabla_{x} f_{1}\right) \phi\left(\xi, s_{1}, s_{2}\right) \\
\phi\left(\xi, s_{1}, s_{2}\right) & =\int_{0}^{s_{2}} f_{1}\left(\xi, s_{1}, \tau\right) d \tau
\end{aligned}
$$

Proof. We begin with the first convolution $K_{\eta_{2}} *\left(\varepsilon^{-1} f_{1}-\langle h\rangle_{12}\right)$. Lemma 1 allows ones to write $x(t)$ as

$$
\begin{equation*}
x(t)=w(t)+\varepsilon \psi(t)+\varepsilon \phi(t) \tag{64}
\end{equation*}
$$

where $\psi(t)=\psi\left(w(t), \varepsilon^{-1} t, \varepsilon^{-2} t\right), \phi(t)=\phi\left(w(t), \varepsilon^{-1} t, \varepsilon^{-2} t\right)$, and $\psi(t)$ is bounded independent of $\varepsilon$.

$$
\begin{aligned}
& K_{\eta_{2}} *\left(\varepsilon^{-1} f_{1}-\langle h\rangle_{12}\right)\left(s_{1}\right) \\
= & \int_{s_{1}-\eta_{2}}^{s_{1}+\eta_{2}} K_{\eta_{2}}\left(s_{1}-s_{2}\right) \cdot\left(\varepsilon^{-1} f_{1}\left(w\left(s_{2}\right)+\varepsilon \psi\left(s_{2}\right)+\varepsilon \phi\left(s_{2}\right), \varepsilon^{-1} s_{2}, \varepsilon^{-2} s_{2}\right)-\langle h\rangle_{12}\right) d s_{2} \\
= & I_{1}+I_{2}+\mathscr{O}(\varepsilon)
\end{aligned}
$$

where we set

$$
\begin{align*}
& I_{1}=\int_{s_{1}-\eta_{2}}^{s_{1}+\eta_{2}} K_{\eta_{2}}\left(s_{1}-s_{2}\right) \cdot\left(\varepsilon^{-1} f_{1}\left(w\left(s_{2}\right)+\varepsilon \psi\left(s_{2}\right), \varepsilon^{-1} s_{2}, \varepsilon^{-2} s_{2}\right)\right) d s_{2} \\
& I_{2}=\int_{s_{1}-\eta_{2}}^{s_{1}+\eta_{2}} K_{\eta_{2}}\left(s_{1}-s_{2}\right) \cdot\left(\nabla_{x} f_{1}\left(w\left(s_{2}\right)+\varepsilon \psi\left(s_{2}\right), \varepsilon^{-1} s_{2}, \varepsilon^{-2} s_{2}\right) \phi\left(s_{2}\right)-\langle h\rangle_{12}\right) d s_{2} . \tag{65}
\end{align*}
$$

Each term in the integrations is estimated in a similar way. Before we move on, recall from Theorem 2 that we identify $f_{1}\left(x\left(s_{2}\right), \varepsilon^{-1} s_{2}, \varepsilon^{-2} s_{2}\right)$ with $f_{1}\left(s_{2}, \varepsilon^{-1} s_{2}, \varepsilon^{-2} s_{2}\right)$. In (65), this simplification is also allowed since by solving the 1 st and 2nd tier, we know $x(t)$ over $\left[t-\eta_{1}, t+\eta_{1}\right]$, and so does $w(t)+\varepsilon \psi(t)$. First, applying Lemma 5 to $I_{1}$ with $\left\langle f_{1}\right\rangle_{2}=0$ yields

$$
\begin{equation*}
\left|I_{1}\right| \leq\left(\frac{\varepsilon^{2 q-1}}{\eta_{2}^{q}}+\frac{\eta_{2}^{\sigma}}{\varepsilon^{\sigma+1}}\right) \max _{j=0 \ldots q}\left\|K^{(j)}\right\|_{1} \tag{66}
\end{equation*}
$$

Second, for $I_{2}$, we now return to consider convolving with $K_{\eta_{1}}$ and $K_{\eta_{2}}$. Note that $\langle h\rangle_{12}$ is a function of $\xi(\cdot)$. Theorem 1 yields that

$$
\sup _{0 \leq t \leq T}|w(t)-\xi(t)| \leq \sup _{0 \leq t \leq T}(|w(t)-x(t)|+|x(t)-\xi(t)|) \leq C \varepsilon .
$$

Then we have $\left\langle\nabla_{x} f_{1} \cdot \phi-\langle h\rangle_{12}\right\rangle_{12}=0+\mathscr{O}(\varepsilon)$, and this allows us to use Lemma 6 and thus to get an estimate for $\left|K_{\eta_{1}} * I_{2}(\cdot)\right|$.

$$
\begin{equation*}
\left|K_{\eta_{1}} * I_{2}(\cdot)\right| \leq C\left(\frac{\varepsilon^{2 q}}{\eta_{2}^{q}}+\frac{\varepsilon^{q}}{\eta_{1}^{q}}+\frac{\eta_{2}^{\sigma}}{\varepsilon^{\sigma}}+\eta_{1}^{\sigma}\right) \max _{j=0 \ldots q}\left\|K^{(j)}\right\|_{1} . \tag{67}
\end{equation*}
$$

Since we differentiate $f_{1}$ with respect to $x, \partial^{r} \nabla_{x} f_{1} \cdot \phi$ are continuous and bounded for $|r|=0, \cdots, \sigma$. Putting estimates (66) and (67) together, Theorem 3 follows.

We conclude this section by proving Proposition 1 which is the cornerstone of our numerical method.
Proof. (Proposition 1) The RHS of the estimate in Theorem 3 dominates that of Lemma 6. Therefore, having $p(<\sigma)$ vanishing moments yields (54).

Remark 4. Theorem 1 is only valid up to times $T$ independent of $\varepsilon$. However, in special cases in which additional cancellation or self averaging occurs, iterated averaging with kernels may give an consistent approximation for the effective behavior of ODEs for longer time intervals.

## 5 Generalizations

This paper is focused on three scale problems modeled by

$$
x^{\prime}=\varepsilon^{-1} f_{1}\left(x, \varepsilon^{-1} t, \varepsilon^{-2} t\right)+f_{0}\left(x, \varepsilon^{-1} t, \varepsilon^{-2} t\right), \quad x(0)=x_{0},
$$

restricting $\left\langle f_{1}\right\rangle_{2}=0$. However, the discussions at the preceding sections suggest several possible generalizations.

### 5.1 Almost-periodic dynamics

The three-scale averaging theorem can be generalized to include dynamics in which the fast $\mathscr{O}\left(\varepsilon^{2}\right)$ or $\mathscr{O}(\varepsilon)$ time scales are not necessarily periodic but rather ergodic on a torus. The periodicity of $s_{2}$ is only taken into account when evaluating the remainder term $r(t ; \varepsilon)$ in (34). In the case of a torus, $r(t ; \varepsilon)$ can be written as a finite sum of periodic functions whose periods are incommensurable. Thus, estimate (35) still holds. A similar generalization can be obtained for almost-periodic functions whose spectrum is bounded away from zero. See, for example [13, 40].

### 5.2 The 3-tier HMM using slow charts

In the proof of Theorem 1, as well as while applying the expansion formalism, it was necessary to assume that the average of the singular $\mathscr{O}\left(\varepsilon^{-1}\right)$ part in $x^{\prime}$ vanishes, $\left\langle f_{1}\right\rangle_{2}=0$. Nonetheless, it can be shown that the estimate of section 2.1 does hold, but only on a short time segment of $\mathscr{O}(\varepsilon)$ length, i.e.

$$
\sup _{0 \leq \leq \leq \varepsilon T}|x(t)-\xi(t)|=\mathscr{O}(\varepsilon) .
$$

In addition, the HMM procedure which utilizes the iterated averaging estimate indicates that in fact, in order to generate a consistent approximation of a slow variable, one only needs to evaluate its derivative on a short time segment of order $\varepsilon$. Then, if the dynamics on the intermediate $\mathscr{O}(\varepsilon)$ is again oscillatory (i.e. periodic or as above), additional averaging on the $\varepsilon$ time scale may average this divergence out. This self averaging property can be captured by the iterated averaging procedure.

To this end, we first need to define slow variables. Formally, slow variables of a dynamical system involving three or more time scales are defined as below [6].

Definition 1. A smooth time dependent function $\alpha:[0, T] \mapsto \mathbb{R}$ is said to evolve on the $\varepsilon^{k}$ time scale in $[0, T]$ for some integer $k$ and for $0<\varepsilon \leq \varepsilon_{0}$, if there exists a smooth function $\beta:[0, T] \mapsto \mathbb{R}$ and constants $C_{0}$ and $C_{1}$ such that

$$
\sup _{t \in[0, T]}\left|\frac{d}{d t} \beta(t)\right| \leq C_{0} \varepsilon^{-k} \quad \text { and } \sup _{t \in[0, T]}|\alpha(t)-\beta(t)| \leq C_{1} \varepsilon
$$

Definition 2. A function $\xi(x)$ is said to evolve on the $\varepsilon^{k}$ time scale along the trajectories of (6) in $[0, T]$ if the time dependent function $\xi\left(x\left(t ; \varepsilon, x_{0}\right)\right)$ evolves on the $\varepsilon^{k}$ time scale in $[0, T]$. For brevity, we will refer to variables that evolve on the $\varepsilon^{0}$ time scale as slow.

The 3-tier HMM using slow variables shares a similar strategy described in Section 1.3, which implements a recursive two-level solver, but the fast oscillations we need to average over are not explicitly given. As [3], with the need for identifying hidden slow variables, we approximate an averaged equation for slow variables by time averaging the microscopic evolution using a suitable kernel. Recall that since the time scales $\mathscr{O}\left(\varepsilon^{2}\right), \mathscr{O}\left(\varepsilon^{1}\right)$, and $\mathscr{O}(1)$ are considered, we need to identify three sets of variables which evolve on each time scale respectively.

Suppose we obtain such a coordinate system using, e.g., the method described in [6]. We denote this system of coordinates $\xi=\left(\xi_{0}, \xi_{1}, \xi_{2}\right)$, where $\xi_{i}=\left(\xi_{i}{ }^{1}, \ldots, \xi_{i}^{d_{i}}\right)$ are the variables evolving on the $\varepsilon^{i}$ time scale, and $\sum_{i=0}^{2} d_{i}=d$. One should take the coordinates of $\xi$ as slow as possible, i.e., if $\phi$ evolves on both $\varepsilon^{0}$ and $\varepsilon^{1}$ time scales, then $\phi \in \xi_{0}$. In terms of the new coordinates the ODE system takes the form

$$
\begin{align*}
& \xi_{0}^{\prime}=\varepsilon^{-1} f_{1}\left(\xi_{0}, \xi_{1}, \xi_{2}\right)+f_{0}\left(\xi_{0}, \xi_{1}, \xi_{2}\right),\left\langle f_{1}\right\rangle_{2}=0 \\
& \xi_{1}^{\prime}=\varepsilon^{-1} g_{1}\left(\xi_{0}, \xi_{1}, \xi_{2}\right)+g_{0}\left(\xi_{0}, \xi_{1}, \xi_{2}\right),\left\langle g_{1}\right\rangle_{2}>C_{1}>0  \tag{68}\\
& \xi_{2}^{\prime}=\varepsilon^{-2} h_{2}\left(\xi_{0}, \xi_{1}, \xi_{2}\right)
\end{align*}
$$

where $\langle\cdot\rangle_{2}$ denotes averaging with respect to the invariant measure for $\xi_{2}$ on fixed $\xi_{0}$ and $\xi_{1}$ and $C_{1}$ is independent of $\varepsilon$, i.e., $\left\langle g_{1}\right\rangle_{2}$ is bounded away from 0 independent of $\varepsilon$. Note that both (68) and (??) are identical. We assume that no resonances, passage through resonances or turning points exists as these may cause hidden slow variables and the decomposition of states into time scales may not be trivial, as discussed in [6, 32].

Then outline of the 3-tier HMM with slow variables is as follows. For simplicity of notation, we suppress the superscript in $\xi_{i}$. As before, we denote the descretized time $t_{n, m}=n h_{0}+m h_{1}$ and $t_{n}=n h_{0}$. We concentrate on the forward Euler and a symmetric kernel.

1. Determination of slow variables:

Find a coordinate system $\xi(x)=\left(\xi_{0}(x), \xi_{1}(x), \xi_{2}(x)\right)$ where $\xi_{i}$ are the variables evolving on the $\varepsilon^{i}$ time scale. See [6] for details. Set $n=0$.
2. Multilevel evolution:

- (0th tier) At $t=t_{n}$, set $\tilde{X}_{0}=x_{1,0}=x_{n}$.
a. (1st tier) For $m=0$ to $k\left(=\eta_{1} / h_{1}\right)$, set $t=t_{n, m}=n h_{0}+m h_{1}$,
i. (2nd tier) Solve the full ODE (6) in $t \in\left[t_{n, m}-\eta_{2}, t_{n, m}+\eta_{2}\right]$ with initial conditions $\tilde{X}_{0}$.
ii. (2nd tier) Force estimation in $\mathscr{O}\left(\varepsilon^{1}\right)$ scale: approximate $\xi_{i}^{\prime}\left(t_{n, m}\right), i=0,1$ by

$$
\begin{equation*}
\left\langle\xi_{i}^{\prime}\right\rangle_{\eta_{2}}\left(t_{n, m}\right)=\left(K_{\eta_{2}} * \xi_{i}^{\prime}\right)\left(t_{n, m}\right) \tag{69}
\end{equation*}
$$

b. (1st tier) $x_{1, m+1}=x_{1, m}+h_{1} \delta x_{1}$, where $\delta x_{1}$ is the least squares solution to the linear system

$$
\begin{equation*}
\delta x_{1} \cdot \nabla \xi_{i}=\left\langle\xi_{i}^{\prime}\right\rangle_{\eta_{2}}, \tag{70}
\end{equation*}
$$

for all $i=0,1$. Redefine $\tilde{X}_{0}=x_{1, m+1}$.
c. End FOR
d. (1st tier) Force estimation in $\mathscr{O}(1)$ scale: approximate $\xi_{0}^{\prime}\left(t_{n}\right)$, by

$$
\begin{equation*}
\left\langle\xi_{0}^{\prime}\right\rangle_{\eta_{1}}\left(t_{n}\right)=K_{\eta_{1}} *\left(\left\langle\xi_{0}^{\prime}\right\rangle_{\eta_{2}}(\cdot)\right)\left(t_{n}\right) . \tag{71}
\end{equation*}
$$

- (0th tier) $x_{n+1}=x_{n}+h_{0} \delta x_{0}$, where $\delta x_{0}$ is the least squares solution to the linear system

$$
\begin{equation*}
\delta x_{0} \cdot \nabla \xi_{0}=\left\langle\xi_{0}^{\prime}\right\rangle \eta_{1} \tag{72}
\end{equation*}
$$

$$
\text { Set } \tilde{X}_{0}=x_{n+1}
$$

3. $n=n+1$ and repeat 2 .

A rigorous proof of Section 5.2 is beyond the scope of the current paper. However, we refer the reader to [2,3] for designing multiscale algorithms that compute the effective behavior of two-scale highly oscillatory dynamical systems by using slow variables. A related example is presented in Section 6.2.

### 5.3 Stochastic effects

The theory of asymptotic expansions and the nested-HMM integrators approach can be extended to a setting in which the intermediate scale is stochastic. In fact, the only requirements for applying the numerical method 1.3 is that the fastest $\mathscr{O}\left(\varepsilon^{2}\right)$ scale is oscillatory - thus ensuring that the effective slow scale is deterministic, and that the intermediate $\mathscr{O}(\varepsilon)$ time scale is ergodic. Consider (68). If the dynamics of $\xi_{1}$ is stochastic, then, on the $\mathscr{O}(1)$ time scale averaging with kernels needs to be replaced by an alternative method such as stochastic HMM [17, 44]. Such an example is presented in Section 6.3.

## 6 Examples

In this section, we numerically apply the iterated HMM approach described in Section 1.3 to deterministic and stochastic systems with three scales. A basic example which has the form of (6) is studied in Section 6.1. Following that, we concentrate on examples which demonstrate the applicability of our method to the generalizations discussed in the previous section. The classical example of two coupled harmonic oscillators in resonance is illustrated in Section 6.2. In this example, one of the slow variables has formally unbounded derivatives as $\varepsilon \rightarrow 0$, but it evolves on the $\varepsilon^{0}$ time scale due to a zero-average of $\varepsilon^{-1}$ term. Section 6.3 is a deterministic example for the generalization discussed in Section 5. The period of the fast oscillator on the $\mathscr{O}\left(\varepsilon^{2}\right)$ time scale changes according to the $\varepsilon$ scale variable. Lastly, an interesting stochastic system whose period of the fastest oscillator changes randomly on the $\varepsilon$ scale is given in Section 6.4.

Our multiscale algorithm is constructed as a family of multilevel ( $>2$ ) solvers which resolve the different time scales and use kernels to estimate the effective force of the slower scales. In Section 6.1 through 6.3 , we use a symmetric $C^{\infty}$ kernel. We see that for a smooth kernel the computational cost is independent of $\varepsilon$ - see [6] for discussions about accuracy and efficiency. In the stochastic example 6.4 , a $\mathbb{K}^{2,7}$ kernel is used.

### 6.1 Example 1

We begin with a simple example of a three-scale system, which is similar to (12)

$$
\begin{equation*}
x^{\prime}=\varepsilon^{-1}\left[x \sin \left(\varepsilon^{-2} t\right)+\cos \left(\varepsilon^{-2} t\right) \sin ^{2}\left(\varepsilon^{-1} t\right)\right]+\cos ^{2}(x)+\cos \left(\varepsilon^{-1} t\right), \quad x(0)=1 \tag{73}
\end{equation*}
$$

Applying Theorem 1 to (73), we have an averaged equation for $x(t)$,

$$
\begin{equation*}
\bar{x}^{\prime}=1 / 4+\cos ^{2}(\bar{x}), \quad \bar{x}(0)=1 \tag{74}
\end{equation*}
$$

The different three-time scales of (73) are illustrated in Figures 2 and 3. The solution $x(t)$ undergoes small-amplitude fast oscillations around the slow trajectory over the interval [0,10]. As proved in Section 2 and 4, the 3-tier HMM approximates $\bar{x}(t)$ which remains close (of order $\varepsilon$ ) to the slow trajectory of $x(t)$. We apply an exponential kernel $K^{\exp } \in \mathbb{K}^{1, \infty}([-1,1])$. See $[3,4,20]$ for details. In Figure 3, we compare the results of 3-tier HMM with both $\bar{x}(t)$ and $x(t)$ obtained by the explicit Runge-Kutta 4th order method. HMM is about 12 times faster than RK4 applied to (73) directly with the step size $h=\varepsilon^{2} / 5$. The computational effort of HMM is independent of $\varepsilon$ once $\eta_{i}$ and $h_{i}$ are fixed. However, for classical numerical methods moving from $\varepsilon=10^{-3}$ to $\varepsilon=10^{-4}$ multiplies the computational effort by 100.

Table 1 Section 6.1. Parameters for the 3-tier HMM of example 1.

| $\varepsilon=10^{-3}$ | $\eta_{i}$ | $h_{i}$ | Method | Kernel |
| :---: | :---: | :---: | :---: | :---: |
| 2nd tier | $18 \varepsilon^{2}$ | $\varepsilon^{2} / 5$ | RK4 | $K^{\exp } \in \mathbb{K}^{1, \infty}([-1,1])$ |
| 1st tier | $18 \varepsilon$ | $\varepsilon^{1} / 5$ | RK2 | $K^{\exp } \in \mathbb{K}^{1, \infty}([-1,1])$ |
| 0th tier | 10 | $1 / 3$ | RK2 | - |




Fig. 2 Section 6.1. The dynamics of (73) on the (Left) $\varepsilon^{2}$ time scale and (Right) $\varepsilon^{1}$ time scale $\left(\varepsilon=10^{-3}\right)$. Plots are obtained by RK4 with $h=\varepsilon^{2} / 100$.


Fig. 3 Section 6.1. The dynamics of (73) on the $\varepsilon^{0}$ time scale. $x(t)$ and $\bar{x}(t)$ are represented by a full line where both are almost indistinguishable. The results of 3-tier HMM are indicated by circles.

### 6.2 Example 2

Consider the following system describing two coupled harmonic oscillators in resonance [6].

As depicted in Figure 4, all four state variables oscillate with a period which is of the order of $\varepsilon^{2}$. Hence, $x_{1}, y_{1}, x_{2}$ and $y_{2}$ evolve on the $\varepsilon^{2}$ time scale.

In order to find a slow coordinate system, we change to polar coordinates $\left(x_{i}, y_{t}\right) \mapsto\left(I_{i}, \varphi_{i}\right), i=1,2$ and introduce a polynomial variable $\theta$ that describes the $1: 1$ resonance between the oscillators

$$
\begin{aligned}
I_{1} & =x_{1}^{2}+y_{1}^{2} \\
I_{2} & =x_{2}^{2}+y_{2}^{2} \\
\theta & =x_{1} x_{2}+y_{1} y_{2} \\
\cos \varphi_{1} & =x_{1} / \sqrt{I_{1}}
\end{aligned}
$$

The corresponding time derivatives are

$$
\begin{aligned}
I_{1}^{\prime} & =2 \varepsilon^{-1} x_{1} y_{2}^{2}-6 x_{1}^{2} x_{2}^{2}+y_{1}^{2} \\
I_{2}^{\prime} & =-2 I_{2}+4 x_{1}^{2} y_{2}^{2} \\
\theta^{\prime} & =\varepsilon^{-1}\left(x_{2} y_{2}^{2}+y_{1} x_{2}-x_{1} y_{2}\right)+\left(-0.5 y_{1} y_{2}-x_{1} x_{2}-3 x_{1} x_{2}^{3}+2 x_{1}^{2} y_{1} y_{2}\right) \\
\varphi_{1}^{\prime} & =\varepsilon^{-2}
\end{aligned}
$$

It appears as if $\left(I_{1}, I_{2}, \theta, \varphi_{1}\right)$ is a chart in which $\varphi_{1}$ evolves of the $\varepsilon^{2}$ time scale, $I_{1}$ and $\theta$ evolve on the $\varepsilon$ time scale while $I_{2}$, which has a bounded derivative, evolves on the $\mathscr{O}(1)$ scale. The dynamics of the three slow variables $I_{1}, I_{2}$ and $\theta$ on the $\mathscr{O}(\varepsilon)$ scale is depicted on the right in Figure 4. The figure suggests that both $I_{1}$ and $I_{2}$ are practically constant on the $\varepsilon$ scale. Indeed, it can be shown that the average of $x_{1} y_{2}^{2}$ on any segment of length $\mathscr{O}(\varepsilon)$ and larger is of order $\varepsilon^{2}$. Therefore, the $\varepsilon^{-1}$ term in $I_{1}^{\prime}$ has a zero average. As a result, the averaged $I_{1}^{\prime}$ is bounded independent of $\varepsilon$ and $I_{1}$ evolves on the $\mathscr{O}(1)$ time scale, rather than the expected $\mathscr{O}(\varepsilon)$.

The time evolution of $I_{1}$ and $I_{2}$ on the slowest $\mathscr{O}(1)$ time scale is depicted in Figure 5. In addition, the figure shows the results of the 3-tier HMM integrator described in Section 1.3. We refer to the solver integrating the $\varepsilon^{i}$ scale as the $i$-th tier. The step-size and length of integration of the $i$-th tier are denoted $h_{i}, \eta_{i}$, respectively. The HMM algorithm approximates the slow $\mathscr{O}(1)$ dynamics using macroscopic steps which are independent of $\varepsilon$. The integration is done using a fourth order method (in the macroscopic step size) and its efficiency is essentially independent of $\varepsilon$. Simulation parameters are detailed in Table 2.

Table 2 Section 6.2. Parameters for the 3-tier HMM of example 2.

| $\varepsilon=10^{-3}$ | $\eta_{i}$ | $h_{i}$ | Method | Kernel |
| :---: | :---: | :---: | :---: | :---: |
| 2nd tier | $70.1 \varepsilon^{2}$ | $\varepsilon^{2} / 10$ | RK4 | $K^{\text {exp }} \in \mathbb{K}^{1, \infty}([-1,1])$ |
| 1st tier | $70.1 \varepsilon^{1}$ | $\varepsilon^{1} / 10$ | RK4 | $K^{\exp } \in \mathbb{K}^{1, \infty}([-1,1])$ |
| 0th tier | 10 | $1 / 3$ | RK4 | - |



Fig. 4 Section 6.2. The dynamics of (75) on the (Left) $\varepsilon^{2}$ time scale and (Right) $\varepsilon^{1}$ time scale. $\varepsilon=10^{-3}$.


Fig. 5 Section 6.2. The dynamics of (75) on the $\varepsilon^{0}$ time scale. Plus signs are results of the 3-tier HMM.

### 6.3 Example 3

Consider the following deterministic system describing two coupled fast harmonic oscillators and a slow dependent mode.

$$
\left\{\begin{array}{l}
x_{1}^{\prime}=-\varepsilon^{-2}\left(1+0.5 \sin y_{2}\right) x_{2}+(1-z)\left(x_{1}^{2}+x_{2}^{2}\right)^{-1} x_{1}  \tag{76}\\
x_{2}^{\prime}=\varepsilon^{-2}\left(1+0.5 \sin y_{2}\right) x_{1}, \\
y_{1}^{\prime}=-\varepsilon^{-1} y_{2}-0.5\left(1+x_{1}^{2}-z\right) y_{1}, \\
y_{2}^{\prime}=\varepsilon^{-1} y_{1} \\
z^{\prime}=-\left(1+0.5 x_{1}^{2}\right) z+y_{2}^{2}
\end{array}\right.
$$

with initial conditions $\left(x_{1}(0), x_{2}(0), y_{1}(0), y_{2}(0), z(0)\right)=(1,0,1,1.5,0.5)$. The system describes two coupled harmonic oscillators ( $x_{1}, x_{2}$ ) and ( $y_{1}, y_{2}$ ) with $\mathscr{O}\left(\varepsilon^{2}\right)$ and $\mathscr{O}(\varepsilon)$ periods, respectively. However, the period of the fastest $\mathscr{O}\left(\varepsilon^{2}\right)$ oscillator depends on $y_{2}$ and is therefore changing on the slower $\varepsilon$ scale. Figure 6 (Left) demonstrates the different period of $x_{1}$ and $x_{2}$ over $2.5 \times 10^{-7}$ duration with $\varepsilon=10^{-4}$. This is an example for the first generalization suggested in Section 5 in which the fastest oscillation exhibits non-trivial dynamics of the intermediate $\varepsilon$ scale.

The system admits three slow variables that evolve on the $\varepsilon^{0}$ scale: $z$, and the square amplitudes of the harmonic oscillators, $I_{1}=x_{1}^{2}+x_{2}^{2}$ and $I_{2}=y_{1}^{2}+y_{2}^{2}$. A numerical algorithm for identifying polynomial slow variables is described in [3]. Hence, we have a coordinate system $\left(\xi_{0}, \xi_{1}, \xi_{2}\right)$ in which $\xi_{i}$ evolves on the $\varepsilon^{i}$ time scale:

$$
\begin{aligned}
& \xi_{0}=\left\{I_{1}, I_{2}, z\right\}, \\
& \xi_{1}=\left\{y_{2}\right\}, \\
& \xi_{2}=\varphi \in S^{1} .
\end{aligned}
$$

As before, we refer to the solver integrating the $\varepsilon^{i}$ scale as the $i$-th tier. The step-size and length of integration of the $i$-th tier are denoted $h_{i}, \eta_{i}$, respectively. The dynamics of the slow variables $I_{1}, I_{2}$ and $z$, as well as the 3 -tier HMM approximation is depicted in Fig. 6 (Right). See Table 2 for simulation parameters.

Table 3 Section 6.3. Parameters for the 3-tier HMM of example 3.

| $\varepsilon=10^{-4}$ | $\eta_{i}$ | $h_{i}$ | Method | Kernel |
| :---: | :---: | :---: | :---: | :---: |
| 2nd tier | $75.1 \varepsilon^{2}$ | $\varepsilon^{2} / 10$ | RK4 | $K^{\exp } \in \mathbb{K}^{1, \infty}([-1,1])$ |
| 1st tier | $75.1 \varepsilon^{1}$ | $\varepsilon^{1} / 10$ | RK2 | $K^{\exp } \in \mathbb{K}^{1, \infty}([-1,1])$ |
| Oth tier | 10 | $1 / 2$ | RK2 | - |



Fig. 6 Section 6.3. (Left) The period of the fastest $\mathscr{O}\left(\varepsilon^{2}\right)$ oscillator is changing on the slower $\varepsilon$ scale. (Right) The dynamics of (76) on the $\varepsilon^{0}$ time scale. Plus signs: 3-tier HMM. Solid line: a reference solution using the RK4 method with step size $h=\varepsilon^{2} / 50$. HMM runs about 1150 times faster. $\varepsilon=10^{-4}$.

### 6.4 Example 4

Consider the following system in which a fast harmonic oscillator has a randomly changing period.

$$
\left\{\begin{array}{l}
d x_{1}=-\left(\varepsilon^{-2}(1+0.5 \sin y) x_{2}+x_{1}(1-z)\right) d t  \tag{77}\\
d x_{2}=\varepsilon^{-2}(1+0.5 \sin y) x_{1} d t \\
d y=-\varepsilon^{-1} y d t+\varepsilon^{-1 / 2} z d \mathrm{~B}_{t} \\
d z=-\left(\left(1+x_{1}^{2}\right) z-y\right) d t
\end{array}\right.
$$

with initial conditions: $\left(x_{1}(0), x_{2}(0), y(0), z(0)\right)=(2,0,1,1)$. In this example, $\left(x_{1}, x_{2}\right)$ is a harmonic oscillator with an $\mathscr{O}\left(\varepsilon^{2}\right)$ period. However, the period changes randomly through a random variable $y$ which is an Ornstein-Uhlenbeck process evolving on the $\mathscr{O}\left(\varepsilon^{1}\right)$ time scale. The system has two slow variables that evolve on the $\mathscr{O}\left(\varepsilon^{0}\right)$ scale: $z$ and $I_{1}=x_{1}^{2}+x_{2}^{2}$. Thus, we find a coordinate system; $\xi_{0}=\left\{I_{1}, z\right\}, \xi_{1}=\{y\}$ and $\xi_{2}=\varphi \in S^{1}$ in which $\xi_{i}$ evolve on the $\mathscr{O}\left(\varepsilon^{i}\right)$ scale.

In order to demonstrate that the effective dynamics of $z$ and $I_{1}$ is indeed deterministic, Figure 7 (Left) shows the standard deviations of $I_{1}$ and $z$ as a function of $\varepsilon$. As expected, it is of order $\sqrt{\varepsilon}$.

Fig. 7 (Right) compares the results computed by the proposed HMM with those by the semi-implicit Euler method [33]. The sample averages of $I_{1}$ and $z$ against $t$ are plotted with a solid line (Euler) and plus signs (HMM). We estimate the errors of the method by comparing the standard deviation of sample paths. Taking $\varepsilon=10^{-4}$, for the semi-implicit Euler, we take 1,000 paths over [0, 4] and decrease step size until the desired accuracy is reached, $\left(\max \left\{\sigma\left(I_{1}\right)\right\}+\max \{\sigma(z)\}\right) / 2=0.1$. This requires $h=\varepsilon^{2} / 100$. For the 3 -tier HMM, we compute 100 independent paths with $h_{0}=0.5$ and calculate the standard deviation for each. The kernel was constructed from Chebyshev polynomials to have exactly seven continuous derivatives and a single vanishing moment:

$$
K^{2,7}(t)=4157010 \chi_{[0,1]}(t)\left(42 t^{2}-44 t+11\right) t^{8}(t-1)^{8}
$$

where $\chi_{[0,1]}$ is the characteristic function of the interval $[0,1]$. HMM parameters are shown in Table 3. With these parameters, HMM achieves $\left(\max \left\{\sigma\left(I_{1}\right)\right\}+\max \{\sigma(z)\}\right) / 2=0.025$. Even if HMM has four times less standard deviation, it runs about 1,000 times faster than the semi-implicit Euler. In addition, we note that the dominant error of 3-tier HMM comes from $h_{0}$ and decreases with smaller $h_{0}$.

Table 4 Section 6.4. Parameters for 3-tier HMM

| $\varepsilon=10^{-4}$ | $\eta_{i}$ | $h_{i}$ | Method | Kernel |
| :---: | :---: | :---: | :---: | :---: |
| 2nd tier | $50 \varepsilon^{2}$ | $\varepsilon^{2} / 10$ | semi-implicit Euler | $K^{2,7} \in \mathbb{K}^{2,7}([0,1])$ |
| 1st tier | $50 \varepsilon^{1}$ | $\varepsilon^{1} / 10$ | Euler | $K^{2,7} \in \mathbb{K}^{2,7}([0,1])$ |
| 0th tier | 10 | $1 / 10$ | RK2 | - |



Fig. 7 Section 6.4. (Left) o-markers are $\max _{t \in[0,4]}\left\{\sigma\left(I_{1}\right)\right\}$. x-markers are $\max _{t \in[0,4]}\{\sigma(z)\}$. The dashed line is a guide for the eye with slope $1 / 2$. (Right) The evolution of the slow variables in example 3. Plus signs: HMM. The solid line is a reference solution calculated as explained in the text.

## 7 Summary

We developed an iterated averaging theory for oscillatory dynamical systems involving three widely separated time scales and the relevant multiscale method for computing the effective behavior. In such multiple time scale problems, we identified a new type of slow variables which do not have formally bounded derivatives. The effective behavior for such variables are studied intensively in two ways: one is a formal approach via the tools of averaging theory, and the other involves homogenization techniques based on singular perturbation expansions and consequent matching of variables. We showed that the results of the developed averaging theory can be efficiently approximated computationally via convolutions of the dynamical system's solutions with a smooth compactly supported kernel. With the developed averaging strategies, we proposed an HMM which is built hierarchically from our previously developed HMMs for two-scale problems. Several numerical examples were presented that demonstrate the efficacy of the proposed algorithms.

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[^1]:    ${ }^{1}$ Loosely speaking, mixing means that for any two possible states of $y_{t}$, the occurrence of the states is independent if a sufficient amount of time $t$ is given. For the precise definition and properties see [45]

[^2]:    ${ }^{2}$ The assumptions holds for the case of ergodic rotations on a torus in which we are interested

[^3]:    ${ }^{3}|r|=r_{1}+r_{2}+\cdots+r_{n}$ and $\partial^{r}=\left(\partial_{x_{1}}\right)^{r_{1}}\left(\partial_{x_{2}}\right)^{r_{2}} \cdots\left(\partial_{x_{n}}\right)^{r_{n}}$.

