An Adaptive Block Bregman Proximal Gradient Method for Computing Stationary States of Multicomponent Phase-Field Crystal Model

Chenglong Bao^{1,3}, Chang Chen¹ and Kai Jiang^{2,*}

 Yau Mathematical Sciences Center, Tsinghua University, Beijing, 100084, China.
 School of Mathematics and Computational Science, Hunan Key Laboratory for Computation and Simulation in Science and Engineering, Xiangtan University, Xiangtan, Hunan, 411105, China.
 Yangi Leke Paijing Institute of Mathematical Science and Ambiastions, China

³ Yanqi Lake Beijing Institute of Mathematical Sciences and Applications, China.

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Abstract. In this paper, we compute the stationary states of the multicomponent phasefield crystal model by formulating it as a block constrained minimization problem. The original infinite-dimensional non-convex minimization problem is approximated by a finite-dimensional constrained non-convex minimization problem after an appropriate spatial discretization. To efficiently solve the above optimization problem, we propose a so-called adaptive block Bregman proximal gradient (AB-BPG) algorithm that fully exploits the problem's block structure. The proposed method updates each order parameter alternatively, and the update order of blocks can be chosen in a deterministic or random manner. Besides, we choose the step size by developing a practical linear search approach such that the generated sequence either keeps energy dissipation or has a controllable subsequence with energy dissipation. The convergence property of the proposed method is established without the requirement of global Lipschitz continuity of the derivative of the bulk energy part by using the Bregman divergence. The numerical results on computing stationary ordered structures in binary, ternary, and quinary component coupled-mode Swift-Hohenberg models have shown a significant acceleration over many existing methods.

AMS subject classifications: 54C40, 14E20

Key words: Multicomponent coupled-mode Swift-Hohenberg model, stationary states, adaptive block Bregman proximal gradient algorithm, convergence analysis, adaptive step size.

*Corresponding author. *Email addresses:* kaijiang@xtu.edu.cn (K. Jiang), clbao@mail.tsinghua.edu.cn (C. Bao), cc20@mails.tsinghua.edu.cn (C. Chen)

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1 Introduction

Multicomponent systems, such as alloys, soft matters, are an important class of materials, particularly for technical applications and processes. The microstructures of materials play a central role for a broad range of industrial application, such as the mechanical property of the quality and the durability, optical device, high-capacity data storage devices [23,37–39,48]. Advances in modeling and computation have significantly improved the understanding of the fundamental nature of microstructure and phase selection processes. Notable contributions have been made through using the phase-field methodology [20,21], which has been successful at examining mesoscale microstructure evolution over diffusive time scales. Recently, phase field crystal (PFC) models have been proposed to efficiently simulate eutectic solidification, elastic anisotropy, solute drag, quasicrystal formation, solute clustering and precipitation mechanisms [14, 38, 46]. Besides, binary and ternary component phase field models have attracted many research interests from the computation perspective [2, 3, 5, 11, 15, 43, 54].

The PFC model for a general class of multicomponent systems is formulated consisting *s* components in *d* dimensional space. The concentrations of the components are described by *s* vector-valued functions $\{\phi_i(\mathbf{r})\}_{i=1}^s = (\phi_1(\mathbf{r}), \dots, \phi_s(\mathbf{r}))$. The variable $\phi_\alpha(\mathbf{r})$, so-called order parameter, denotes the local fraction of phase α . The free energy functional of PFC model of a *s*-component system can be described by two contributions, a bulk free energy $F[\{\phi_i(\mathbf{r})\}_{i=1}^s]$ and an interaction potential $G[\{\phi_i(\mathbf{r})\}_{i=1}^s]$, which drive the density fields to become ordered by creating minimal in the free energy for these states. Formally, we can write the free energy functional of the multicomponent system as

$$E[\{\phi_i(\mathbf{r})\}_{i=1}^s;\Theta] = G[\{\phi_i(\mathbf{r})\}_{i=1}^s;\Theta] + F[\{\phi_i(\mathbf{r})\}_{i=1}^s;\Theta],$$
(1.1)

where Θ are relevant physical parameters. *F* has polynomial or logarithmic formulation [12, 16, 47] and *G* is the interaction potential, such as high-order differential terms or convolution terms [12, 42]. Usually, some constraints are imposed on the PFC model, such as the mass conservation or incompressibility which means the order parameter $\{\phi_i(\mathbf{r})\}_{i=1}^s$ belong to a feasible space.

To understand the fundamental nature of multicomponent systems, it often involves finding stationary states corresponding to ordered structures. Denote V_i ($i = 1, 2, \dots, s$) to be a feasible space of the *i*-th order parameter, the above problem is transformed into solving the minimization problem

min
$$E[\{\phi_i(\mathbf{r})\}_{i=1}^s;\Theta]$$
, s.t. $\phi_i(\mathbf{r}) \in V_i \ (i=1,2,\cdots,s)$, (1.2)

with different physical parameters Θ , which brings a tremendous computational burden.

Different methods have been proposed for computing the stationary states of multicomponent models and can be classified into two categories through different formulations and numerical techniques. One is to solve the steady nonlinear Euler-Lagrangian system of (1.2) through different spatial discretization approaches. The other class of approaches has been constructed via the formulation of the minimization problem (1.2). Among these methods, great efforts have been made for solving the nonlinear gradient flow equations. Numerically, a gradient flow equation is discretized in both space and time domains via different discretization techniques, and the stationary state is obtained with proper choices of initialization and step sizes. Typical energy stable schemes to gradient flows include convex splitting [49,52], stabilized factor methods [45], exponential time differencing scheme [17,18,32,56], and recently developed invariant energy quadrature [54], and scalar auxiliary variable approaches [43]. When designing the fully discretized scheme, choosing proper time steps greatly impacts the performance of the methods. Most existing methods fix the time step or obtain the time step using heuristic methods [41].

In this work, instead of designing a numerical scheme for the gradient flow, we formulate the infinite-dimensional problem (1.2) as a finite-dimensional block-wise nonconvex problem via appropriate spatial discretization schemes. Similar ideas have shown success in computing stationary states of many physical problems, such as the Bose-Einstein condensate [50], the calculation of density functional theory [34] and one component PFC models [27]. It is noted that extending single component systems to multicomponent models is not trivial due to the following two reasons. First, the direct extension may fail convergence, even for the simplest steep descent method [40]. Second, the update order's choice is not unique, leading to difficulty in analyzing the convergence of the numerical algorithm. Therefore, it deserves to develop specialized numerical algorithms for the multicomponent systems, which require carefully examine or extend the corresponding analysis in a single component system.

In this paper, based on the newly developed optimization techniques, we propose an adaptive block Bregman proximal gradient (AB-BPG) method for solving the discretized problem with multi-block structures. The proposed algorithm has the desired energy dissipation and mass conservation properties. Theoretically, we prove the convergence property of the proposed algorithm without the global Lipschitz constant requirement. It guarantees that the algorithm converges to a stationary state given any initial point. Our method updates each order parameter function with adaptive step sizes by the line search method in practice. Compared with our previous work [27] for single-component PFC models, the main contributions of this paper include:

- We propose an AB-BPG algorithm for arbitrary multiple components PFC models by involving in the block structures. The sequence of blocks can be updated either deterministically cyclic or randomly shuffled for each iteration. Moreover, the convergence property without the global Lipschitz constant assumption of the derivative of the bulk energy *F* is rigorously proved once each block is updated at least once in every *T* iterations;
- Together with a practical line search strategy, the sequence $\{\Phi^k = \{\phi_i^k\}_{i=1}^s\}$ generated by the proposed method has the generalized energy dissipation property, i.e., one of the following properties holds:

- 1. The sequence has the energy dissipation property;
- 2. There exist a subsequence $\{\Phi^{k_j}\} \subset \{\Phi^k\}$ and a constant $M \in \mathbb{N}$ such that $1 \le k_{j+1}-k_j \le M+1$ and $E(\Phi^{k_{j+1}})-E(\Phi^{k_j}) \ge 0, \forall j$;
- Extensive numerical experiments on computing the stationary states in the binary, ternary, and quinary coupled-mode Swift-Hohenberg model has shown the advantages of the proposed method in terms of computational efficiency.

The rest of this paper is organized as follows. Section 2 presents a concrete multicomponent PFC model, i.e., the coupled-mode Swift-Hohenberg (CMSH) model, and a spatial discretization formulation based on the projection method. In Section 3, we propose the adaptive block Bregman proximal gradient (AB-BPG) methods for solving the constrained non-convex multi-block problems with proved convergence. In Section 4, we apply the proposed approaches to the CMSH model with two choices of Bregman divergence. Numerical results are reported in Section 5 to illustrate the efficiency and accuracy of our algorithms.

2 Problem formulation

There are several multicomponent PFC models to describe the phase behaviors of alloys and soft-matters [2,3,19,24,26,28,36,38,38,39,48]. In this work, we consider the coupled-mode Swift-Hohenberg (CMSH) model of multicomponent systems, which extends the classical Swift-Hohenberg model from one length scale to multiple length scales [26, 28, 47]. The CMSH model allows the study of the formation and relative stability of periodic crystals and quasicrystals. Define the integral average

$$\int = \begin{cases} \frac{1}{|\Omega|} \int_{\Omega'} & \text{for periodic crystals,} \\ \lim_{R \to \infty} \frac{1}{|B_R|} \int_{B_R'} & \text{for quasicystals,} \end{cases}$$
(2.1)

where Ω is a bounded domain and B_R is a ball centered at origin with radii R. The free energy of the CMSH model for s component system is

$$E[\{\phi_j(\boldsymbol{r})\}_{j=1}^s] = \oint \left\{ \frac{1}{2} \sum_{j=1}^s [(\nabla^2 + q_j^2)\phi_j(\boldsymbol{r})]^2 + \sum_{\mathcal{I}_{s,n}} \tau_{i_1,i_2,\cdots,i_s} \prod_{j=1}^s \phi_j^{i_j}(\boldsymbol{r}) \right\} d\boldsymbol{r},$$
(2.2)

where ϕ_j is the *j*-th order parameter, $q_j > 0$ is the *j*-th characteristic length scale, $\tau_{i_1,i_2,...,i_s}$ is interaction intensity related to the physical conditions, and $\mathcal{I}_{s,n}$ is the index set defined as

$$\mathcal{I}_{s,n} := \left\{ (i_1, i_1, \cdots, i_s) : i_j \in \mathbb{N} \ (j = 1, 2, \cdots, s), \ 1 \le \sum_{j=1}^s i_j \le n \right\}.$$

C. Bao, C. Chen and K. Jiang / CSIAM Trans. Appl. Math., 3 (2022), pp. 133-171

Moreover, to conserve the average density, each order parameter ϕ_i satisfies

$$\oint \phi_j(\mathbf{r}) d\mathbf{r} = 0. \tag{2.3}$$

Theoretically, the ordered patterns including periodic and quasiperiodic structures correspond to local minimizers of the free energy functional (2.2) with respect to order parameters ϕ_i ($j = 1, 2, \dots, s$). Thus, denote

$$G_{j}[\phi_{j}] = \int \frac{1}{2} [(\nabla^{2} + q_{j}^{2})\phi_{j}(\mathbf{r})]^{2} d\mathbf{r}, \ F[\{\phi_{j}(\mathbf{r})\}_{j=1}^{s}] = \int \sum_{\mathcal{I}_{s,n}} \tau_{i_{1},i_{2},\cdots,i_{s}} \prod_{j=1}^{s} \phi_{j}^{i_{j}}(\mathbf{r}) d\mathbf{r}, \qquad (2.4)$$

we focus on solving the minimization:

min
$$E[\{\phi_j(\mathbf{r})\}_{j=1}^s] = \sum_{j=1}^s G_j[\phi_j(\mathbf{r})] + F[\{\phi_j(\mathbf{r})\}_{j=1}^s]$$

s.t. $\oint \phi_j(\mathbf{r}) d\mathbf{r} = 0, \quad j = 1, 2, \cdots, s.$ (2.5)

Throughout this paper, we assume that $n \le 4$ in (2.4) which means that bulk energy *F* is a 4^{th} order polynomial.

In this work, we pay attention to the periodic and quasiperiodic crystals and use the projection method [30] to discretize the CMSH free energy functional. After discretization, the infinite-dimensional problem (2.5) can be formulated to a finite-dimensional minimization problem in the form of

$$\min_{\hat{\Phi}} E(\hat{\Phi}) = \sum_{j=1}^{s} G_j(\hat{\phi}_j) + F(\hat{\Phi}), \quad \text{s.t.} \quad e_1^{\top} \hat{\phi}_j = 0, \quad j = 1, 2, \cdots, s,$$
(2.6)

where $\hat{\phi}_j \in \mathbb{C}^{N_j}$ is the truncated Fourier coefficients corresponding to ϕ_j and $\hat{\Phi} = \{\hat{\phi}_j\}_{j=1}^s \in \mathbb{C}^N$ with $N = \sum_{j=1}^s N_j$. $G_j(\hat{\phi}_j) = \frac{1}{2} \langle \hat{\phi}_j, \mathcal{D}_j \hat{\phi}_j \rangle$ and $\mathcal{D}_j \in \mathbb{C}^{N_j \times N_j}$ is a diagonal matrix. $F(\{\hat{\phi}_j\}_{j=1}^s)$ are *n*-dimensional convolutions in the reciprocal space. A direct evaluation of the nonlinear term $F(\{\hat{\phi}_j\}_{j=1}^s)$ is extremely expensive, however, $F(\{\hat{\phi}_j\}_{j=1}^s)$ is a simple multiplication in the *n*-dimensional physical space. Thus, the pseudospectral method takes the advantage of this observation by evaluating $G_j(\hat{\phi}_j)$ in the Fourier space and $F(\{\hat{\phi}_j\}_{j=1}^s)$ in the physical space via the Fast Fourier Transformation. For the self-containess, we leave the concrete discretization to Appendix A.

Compared to the single-component case, the problem (2.6) has the block structure in terms of the order parameters. Besides, the objective function is the summation of non-separable bulk energy and separable interaction energy that facilitates the numerical algorithm's design. In the next section, we aim at designing efficient optimization-based algorithms for solving the non-convex and multi-block problem (2.6). It is worth noting that the proposed AB-BPG method can also be applied to other spatial discretization methods.

3 The proposed method

In this section, we consider the minimization problem in the form of

$$\min_{X} E(X) = f(X) + \sum_{j=1}^{s} g_j(x_j)$$
s.t. $x_j \in S_j, \quad j = 1, 2, \cdots, s,$
(3.1)

where $\mathbf{x}_j \in \mathbb{C}^{N_j}$, S_j is the feasible space of variable \mathbf{x}_j , $X = {\mathbf{x}_j}_{j=1}^s \in \mathbb{C}^N$ with $N = \sum_{j=1}^s N_j$. It is easy to know that the problem (2.6) can be reduced to (3.1) by setting f = F, $g_j = G_j$ and $S_j = {\hat{\phi}_j : e_1^\top \hat{\phi}_j = 0}$. Throughout this paper, we make the following assumptions on the objective function (see the concrete definition of notations in the next subsection).

Assumption 3.1.

- 1. $f: \mathbb{C}^N \to (-\infty, \infty]$ is proper and continuously differential on \mathbb{C}^N but may not be convex.
- 2. $g_j : \mathbb{C}^{N_j} \to (-\infty, \infty]$ is proper, lower semicontinuous and convex.
- 3. $S_i \subseteq \text{dom}g_i$ is a nonempty, closed and convex set.
- 4. *E* is bounded below, and level bounded.
- 5. For all $X \in \text{dom}E$, $\pi_j(\mathcal{B}(X)) \subseteq \text{ridom}g_j$ where $\mathcal{B}(X)$ is the closed ball that contains the sub-level set $[E \leq E(X)] \cap \prod_{j=1}^s S_j$.

Before presenting our numerical algorithm, we first introduce some notations and useful definitions in the following analysis. Then, we give an abstract framework of the first order method with proved convergence. Two kinds of concrete numerical algorithms to the CMSH model based on the abstract formulation for solving (2.6) will be presented in Section 4.

3.1 Notations and definitions

We denote $\prod_{j=1}^{s} S_j := \{X = (x_1, \dots, x_s) : x_j \in S_j, \forall j = 1, \dots, s\}$ and the projection operator is defined as $\pi_j : \mathbb{C}^N \to \mathbb{C}^{N_j}, X \mapsto x_j$. For a subset $S \subseteq \mathbb{C}^N, \pi_j(S) := \{x_j : x_j = \pi_j(X), \forall X \in S\}$. Let $C^k(S)$ be the *k*-th continuously differential functions on *S*. The domain of a function $f : \mathbb{C}^N \to \mathbb{R}$ is defined as dom $f := \{x : f(x) < +\infty\}$ and the relative interior of dom *f* is defined as ridom $f := \{x \in \text{dom } f : \exists r > 0, B(x, r) \cap \text{affdom } f \subseteq \text{dom } f\}$, where affdom *f* is the smallest affine set that contains dom *f* and $B(x, r) := \{y : ||y - x|| \le r\}$. *f* is proper if $f > -\infty$ and dom $f \neq \emptyset$. For $\alpha \in \mathbb{R}$, $[f \le \alpha] := \{x : f(x) \le \alpha\}$ is the α -(sub)level set of *f*. We say that *f* is level bounded if $[f \le \alpha]$ is bounded for all $\alpha \in \mathbb{R}$. *f* is lower semicontinuous if all level set of *f* is closed. The subgradient of *f* at $x \in \text{dom } g$ is defined as $\partial f(x) = \{u : x \in \mathbb{R}\}$.

 $f(y)-f(x)-\langle u,y-x\rangle \ge 0, \forall y \in \text{dom} f$ }. For $a \in \mathbb{R}$, we denote $[a]^+ := \max\{0,a\}$. Moreover, the next table summarizes the notations used in this work.

Throughout this paper, we assume that $h: \mathbb{C}^N \to (-\infty, +\infty]$ is a strongly convex function and give the following useful definitions.

Definition 3.1 (Bregman divergence [33]). *The Bregman divergence with respect to* $h \in C^1(\text{intdom}h)$ *is defined as*

$$D_h(x,y) = h(x) - h(y) - \langle \nabla h(y), x - y \rangle, \quad \forall \ (x,y) \in \operatorname{dom} h \times \operatorname{intdom} h.$$
(3.2)

It is noted that $D_h(x,y) \ge 0$ and $D_h(x,y) = 0$ if and only if x = y due to the strongly convexity of h. Moreover, $D(\cdot,y)$ is also strongly convex for any fixed $y \in \text{intdom}h$. The Bregman divergence with $h = \|\cdot\|^2/2$ reduces to the Euclidean distance. Using the Bregman divergence, we can generalize the Lipschitz condition to the so-called relative smoothness as follows.

Definition 3.2 (Relative smoothness [8]). A function f is called R_f -smooth relative to h if there exists $R_f > 0$ such that $R_f h(x) - f(x)$ is convex for all $x \in \text{dom}h$.

When $h = \|\cdot\|^2/2$, the relative smoothness reduces to the Lipschitz smoothness which is an essential assumption in the analysis of many scheme in computing gradient flows (such as semi-implicit [22] or stabilized semi-implicit [45]). However, this assumption greatly limits its application range in practical computation. In this work, we overcome this difficulty from numerical optimization using this novel tool. To deal with the multiblock problem of form (3.1), we generalize the definition of relative smoothness to blockwise function as Definition 2.4 in [1].

Definition 3.3 (Block-wise relative smoothness). For a block-wise function f(X), we call f(X) is $(R_f^1, R_f^2, \dots, R_f^s)$ -smooth relative to (h_1, h_1, \dots, h_s) if for each *j*-th block and fixed $\mathbf{x}_{\neq j}$,

- $h_i: \mathbb{C}^{N_j} \to (-\infty, +\infty]$ is γ_i -strongly convex and dom $h_i \subseteq \{u: f(u, x_{\neq i}) < \infty\}$.
- $F_j(u) := f(u, x_{\neq j})$ is R_f^j -smooth relative to h_j with respect to u.

Now, we are ready to present the numerical algorithm in the next subsection.

3.2 AB-BPG algorithm

Our main idea is to develop a kind of block coordinate descent methods which minimize *E* cyclically over each of x_1, x_2, \dots, x_s while fixing the remaining blocks at their last updated values, i.e., the Gauss-Seidel fashion. Precisely, given feasible $X^k, X^{k-1} \in \prod_{j=1}^s S_j$, our AB-BPG method can pick $b_k \in \{1, 2, \dots, s\}$ deterministically or randomly, then $X^k = (x_1^k, x_2^k, \dots, x_s^k)$ is updated as follows

$$\begin{cases} \boldsymbol{x}_{i}^{k+1} = \boldsymbol{x}_{i}^{k}, & \text{if } i \neq b_{k}, \\ \boldsymbol{x}_{i}^{k+1} = \operatorname*{argmin}_{\boldsymbol{z} \in \mathcal{S}_{i}} \left\{ g_{i}(\boldsymbol{z}) + \langle \nabla_{i} f(\boldsymbol{y}^{k}, \boldsymbol{x}_{\neq i}^{k}), \boldsymbol{z} - \boldsymbol{y}^{k} \rangle + \frac{1}{\alpha_{k}} D_{h_{i}}(\boldsymbol{z}, \boldsymbol{y}^{k}) \right\}, & \text{if } i = b_{k}, \end{cases}$$
(3.3)

Notation	Definition
S	the total number of blocks
b_k	the update block selected at the <i>k</i> -th iteration
n_j^k	the number of updates to x_j within the first k iterations
x_j^k	the value of x_j after the <i>k</i> -th iteration
$ ilde{m{x}}_j^n$	the value of x_j after <i>n</i> -th update
y^k	the value of extrapolation point y at the k -th iteration
w_k	the extrapolation weight used at the k -th iteration
m_k	$m_k = \operatorname{argmin}_j \{ E(X^j) : [k - M]^+ \le j \le k \}$
X	$X = (x_1, x_2, \cdots, x_s)$
$x_{ eq j}$	the value of $(x_1, \cdots, x_{j-1}, x_{j+1}, \cdots, x_s)$
$\nabla_j f(X)$	the partial gradient of $f(X)$ with respect to x_j

Table 1: Summary of notations.

where $\alpha_k > 0$ is the step size and y^k is the extrapolation

$$\boldsymbol{y}^{k} = (1 + w_{k})\boldsymbol{x}_{i}^{k} - w_{k}\boldsymbol{x}_{i}^{\text{prev}} = (1 + w_{k})\tilde{\boldsymbol{x}}_{i}^{n_{i}^{k}} - w_{k}\tilde{\boldsymbol{x}}_{i}^{n_{i}^{k}-1}.$$
(3.4)

The extrapolation weight $w_k \in [0, \bar{w}]$ for some $\bar{w} > 0$ and x_i^{prev} is the value of x_i before it is updated to x_i^k . The definitions of \tilde{x}_i and n_i^k can be found in Table 1.

To ensure the convergence, we make mild assumptions for the distance generating function h_i , $i = 1, 2, \dots, s$ and the block update order [53].

Assumption 3.2. There exist strongly convex functions $h_j, j = 1, 2, \dots, s$, and positive constants $R_f^j, j = 1, 2, \dots, s$ such that $S_j \subseteq \text{intdom} h_j, j = 1, 2, \dots, s$, and f is $(R_f^1, R_f^2, \dots, R_f^s)$ -smooth relative to (h_1, h_1, \dots, h_s) .

Assumption 3.3. With any $T \ge s$ consecutive iterations, each block should be updated at least once, i.e., for any k, it has $\{1, 2, \dots, s\} \subseteq \{b_k, b_{k+1}, \dots, b_{k+T}\}$.

In the following context, we show some properties of the iterates (3.3). First, define the *j*-th block Bregman proximal gradient mapping $T_{\alpha}^{j} : \prod_{i=1}^{s} S_{i} \to \mathbb{C}^{N_{j}}$ as

$$T^{j}_{\alpha}(X) := \operatorname*{argmin}_{z \in \mathcal{S}_{j}} \Big\{ g_{j}(z) + \langle \nabla_{j} f(X), z - x_{j} \rangle + \frac{1}{\alpha} D_{h_{j}}(z, x_{j}) \Big\}.$$
(3.5)

The next lemma shows the well-posedness of T^{j}_{α} .

Lemma 3.1. Suppose Assumption 3.1 and Assumption 3.2 hold. The map T^{j}_{α} defined in (3.5) is nonempty and single-valued from $\prod_{i=1}^{s} S_i$ to S_j .

Proof. Since $f \in C^1(\mathbb{C}^n)$ and $S_j \subset intdomh_j$, $\nabla_j f(X)$ and $D_{h_i}(\cdot, x_j)$ is well-defined. Let

$$\Gamma_{j}(z) := \alpha g_{j}(z) + \alpha \langle \nabla_{j} f(X), z - x_{j} \rangle + D_{h_{j}}(z, x_{j}),$$

we know that Γ_j is strongly convex due to the convexity of $D(\cdot, x_j)$ and g_j . Thus, Γ_j is coercive [9, Corollary 11.17]. According to the Corollary 3.23 in [13], Γ_j achieves its minimum on S_j and the strongly convexity implies the uniqueness of minimum.

Remark 3.1. In Assumption 3.1, S_j is convex for all j. Thus, Lemma 3.1 implies that the iteration $\mathbf{x}_i^{k+1} = T_{\alpha_k}^i(\mathbf{y}^k, \mathbf{x}_{\neq i}^k)$ in (3.3) is well-defined as long as $X^k, X^{k-1} \in \prod_{j=1}^s S_j$ is set in the initialization.

The next lemma shows that the mapping T^{j}_{α} has the descent property.

Lemma 3.2 (Sufficient decrease property). Suppose Assumption 3.1 and Assumption 3.2 hold. Let $X = {x_i}_{i=1}^s \in \prod_{i=1}^s S_i$ and $x_i^+ = T_{\alpha}^j(X)$, then we have

$$E(X) - E(\mathbf{x}_{j}^{+}, \mathbf{x}_{\neq j}) \ge \left(\frac{1}{\alpha} - R_{f}^{j}\right) \frac{\gamma_{j}}{2} \|\mathbf{x}_{j}^{+} - \mathbf{x}_{j}\|^{2},$$
(3.6)

where γ_i is the strong convexity coefficient of h_i .

Proof. Due to the block-wise relative smoothness of f(X), the function $F_j(u) := R_f^j h_j(u) - f(u, x_{\neq j})$ is convex for any fixed $x_{\neq j}$. Thus, for any $z \in \text{intdom} h_j$, we have

$$f(\boldsymbol{u},\boldsymbol{x}_{\neq j}) \leq f(\boldsymbol{z},\boldsymbol{x}_{\neq j}) + \langle \nabla_j f(\boldsymbol{z},\boldsymbol{x}_{\neq j}), \boldsymbol{u} - \boldsymbol{z} \rangle + R_f^J D_{h_j}(\boldsymbol{u},\boldsymbol{z}).$$
(3.7)

Together with the definition of mapping T^{j}_{α} , we know that

$$E(X) = f(X) + g_{j}(\mathbf{x}_{j}) + \sum_{i \neq j} g_{i}(\mathbf{x}_{i})$$

$$= f(X) + \left[\langle \nabla_{j}f(X), \mathbf{z} - \mathbf{x}_{j} \rangle + \frac{1}{\alpha} D_{h_{j}}(\mathbf{z}, \mathbf{x}_{j}) + g_{j}(\mathbf{z}) \right]_{\mathbf{z} = \mathbf{x}_{j}} + \sum_{i \neq j} g_{i}(\mathbf{x}_{i})$$

$$\geq f(X) + \langle \nabla_{j}f(X), \mathbf{x}_{j}^{+} - \mathbf{x}_{j} \rangle + \frac{1}{\alpha} D_{h_{j}}(\mathbf{x}_{j}^{+}, \mathbf{x}_{j}) + g_{j}(\mathbf{x}_{j}^{+}) + \sum_{i \neq j} g_{i}(\mathbf{x}_{i})$$

$$\geq f(\mathbf{x}_{j}^{+}, \mathbf{x}_{\neq j}) - R_{f}^{j} D_{h_{j}}(\mathbf{x}_{j}^{+}, \mathbf{x}_{j}) + \frac{1}{\alpha} D_{h_{j}}(\mathbf{x}_{j}^{+}, \mathbf{x}_{j}) + g_{j}(\mathbf{x}_{j}^{+}) + \sum_{i \neq j} g_{i}(\mathbf{x}_{i})$$

$$= E(\mathbf{x}_{j}^{+}, \mathbf{x}_{\neq j}) + \left(\frac{1}{\alpha} - R_{f}^{j}\right) D_{h_{j}}(\mathbf{x}_{j}^{+}, \mathbf{x}_{j}) \geq E(\mathbf{x}_{j}^{+}, \mathbf{x}_{\neq j}) + \left(\frac{1}{\alpha} - R_{f}^{j}\right) \frac{\gamma_{j}}{2} ||\mathbf{x}_{j}^{+} - \mathbf{x}||^{2}.$$

The second inequality holds by setting $u = x_j^+$ and $z = x_j$ in (3.7), and the last inequality holds from the γ_i -convexity of h_i .

Restart technique. Let $i = b_k$, the iteration (3.3) can be written as $X^{k+1} = T^i_{\alpha_k}(\boldsymbol{y}^k, \boldsymbol{x}^k_{\neq i})$, then Lemma 3.2 implies that $E(\boldsymbol{y}^k, \boldsymbol{x}^k_{\neq i}) \ge E(X^{k+1})$ as long as $\alpha_k \in (0, 1/R^i_f]$. However, this descent property is not enough for investigating the iterates $\{E(X^k)\}$ as the relationship between $E(\boldsymbol{y}^k, \boldsymbol{x}^k_{\neq i})$ and $E(X^k)$ is not clear. Thus, it may lead to the energy oscillation in the sequence $\{E(X^k)\}$. To overcome this weakness, we propose a restart technique by setting $(\boldsymbol{y}^k, \boldsymbol{x}^k_{\neq i}) = X^k$ when the sharp oscillation is detected. In concrete, given $\alpha_k > 0$ and $i = b_k$, define

$$\boldsymbol{z}^{k} = \Gamma_{\alpha_{k}}^{i}(\boldsymbol{y}^{k}, \boldsymbol{x}_{\neq i}^{k}) = \operatorname*{argmin}_{\boldsymbol{z} \in \mathcal{S}_{i}} \left\{ g_{i}(\boldsymbol{z}) + \langle \nabla_{i} f(\boldsymbol{y}^{k}, \boldsymbol{x}_{\neq i}^{k}), \boldsymbol{z} - \boldsymbol{y}^{k} \rangle + \frac{1}{\alpha_{k}} D_{h_{i}}(\boldsymbol{z}, \boldsymbol{y}^{k}) \right\}.$$
(3.8)

Given some non-negative integer constant M, we define

$$m_k = \underset{[k-M]^+ \le j \le k}{\operatorname{argmax}} E(X^j), \tag{3.9}$$

and set $X^{k+1} = X^k$, $w_{k+1} = 0$ if the following inequality

$$E(X^{m_k}) - E(z^k, x^k_{\neq i}) \ge \sigma \|x^k_i - z^k\|^2$$
(3.10)

does not hold where $\sigma > 0$ is a small constant. Otherwise, we obtain X^{k+1} via $x_i^{k+1} = z^k$, $x_i^{k+1} = x_j^k$ $(j \neq i)$ and update $w_{k+1} \in [0, \bar{w}]$.

Remark 3.2. When M = 0, it guarantees that $\{E(X^k)\}$ is monotone decreasing. When M > 0, the scheme has the generalized descent property, i.e., the subsequence $\{E(X^{m_k})\}$ is decreasing (see Lemma 3.3).

Step size estimation. Let $i = b_k$, Lemma 3.2 shows that $E(\mathbf{y}^k, \mathbf{x}_{\neq i}^k) \ge E(X^{k+1})$ is ensured by step size $\alpha_k \in (0, 1/R_f^i]$ which may be too conservative. Thus, we propose a non-monotone backtracking line search method [35] for finding the appropriate step α_k which is initialized by the similar idea of BB method [7], i.e.,

$$\alpha_{k} = \begin{cases} \alpha_{0}, & w_{k} = 0, \\ \frac{\langle u_{i}^{k}, u_{i}^{k} \rangle}{\langle u_{i}^{k}, v_{i}^{k} \rangle} \text{ or } \frac{\langle v_{i}^{k}, u_{i}^{k} \rangle}{\langle v_{i}^{k}, v_{i}^{k} \rangle}, & w_{k} \neq 0, \end{cases}$$
(3.11)

where $u_i^k = y^k - x_i^k$ and $v_i^k = \nabla_i f(y^k, x_{\neq i}^k) - \nabla_i f(X^k)$. Let $\eta \in (0, \sigma]$ be a constant and z^k is obtained from (3.8), we adopt the step size $\alpha_k \in [\alpha_{\min}, \alpha_{\max}]$ whenever the following inequality holds

$$\max(E(\boldsymbol{y}^{k}, \boldsymbol{x}^{k}_{\neq i}), E(X^{m_{k}})) - E(\boldsymbol{z}^{k}, \boldsymbol{x}^{k}_{\neq i}) \ge \eta \|\boldsymbol{y}^{k} - \boldsymbol{z}^{k}\|^{2}.$$
(3.12)

Let

$$R := \max_{j=1,2,\cdots,s} \{R_j\} \quad \text{and} \quad \gamma = \max_{j=1,2,\cdots,s} \{\gamma_j\}.$$

Lemma 3.2, the inequality (3.12) holds whenever $0 < \alpha_{\min} < \gamma/(2\eta + \gamma R)$. Thus, the line search scheme will terminate in finite iterations. In summary, we present the detailed algorithm for estimating step sizes in Algorithm 1 and the proposed AB-BPG method in Algorithm 2.

Algorithm 1 Estimation of α_k at y^k with respect to block *i*

- 1: **Inputs:** X^k , y^k , $\zeta \in (0,1)$ and $\eta > 0$, $\alpha_0, \alpha_{\min}, \alpha_{\max} > 0$.
- 2: Initialize α_k by (3.11);
- 3: Calculate the smallest index $\ell \ge 0$ such that (3.12) holds at z^k defined in (3.8) with step size $\zeta^{\ell} \alpha_k \geq \alpha_{\min}$.
- 4: **Output:** step size $\alpha_k = \max(\min(\zeta^{\ell} \alpha_k, \alpha_{\max}), \alpha_{\min})$.

Algorithm 2 AB-BPG method

1: Initialize $X^{-1} = X^0 \in \prod_{j=1}^s S_j$, $\sigma \ge \eta > 0$ and $w_0 = 0$, $\bar{w}, \alpha_0, M \ge 0$, k = 0.

- 2: while stopping criterion is not satisfied do
- Pick $i = b_k \in \{1, 2, \dots, s\}$ in a deterministic or random manner; 3:
- Update $y^k = (1+w_k)\tilde{x}_i^{n_i^k} w_k\tilde{x}_i^{n_i^k-1}$ Obtain α_k via Algorithm 1. 4:
- 5:
- Calculate z^k via (3.8). 6:
- if (3.10) holds then 7:
- Set $x_i^{k+1} = z^k, x_i^{k+1} = x_i^k (j \neq i)$ and choose $w_{k+1} \in [0, \bar{w}]$. 8:
- else 9:
- Restart by setting $X^{k+1} = X^k$ and $w_{k+1} = 0$. 10:
- end if 11.
- k = k + 1. 12:
- 13: end while

The next lemma establishes the generalized descent property of the sequence $\{X^k\}$ generated by Algorithm 2.

Lemma 3.3. Suppose Assumption 3.1 and Assumption 3.2 hold. Let $\{X^k\}$ be the sequence generated by Algorithm 2. Then, we have $m_{k+1} \ge m_k$ and $\{E(X^{m_k})\}$ is non-increasing.

Proof. By the definition of m_k , it is easy to know $m_{k+1} \ge m_k$. If the non-restart condition (3.10) does not hold, we know $X^{k+1} = X^k$ which implies

$$E(X^{m_k}) - E(X^{k+1}) \ge 0. \tag{3.13}$$

If the non-restart condition (3.10) holds, we have

$$E(X^{m_k}) - E(X^{k+1}) \ge \sigma \|X^k - X^{k+1}\|^2 \ge 0.$$
(3.14)

Combing (3.13) with (3.14), we obtain

$$E(X^{m_{k+1}}) = \max\{E(X^{j}) | [k+1-M]^* \le j \le k+1)\}$$

$$\le \max\{E(X^{j}) | [k-M]^* \le j \le k\} = E(X^{m_k}).$$

This completes the proof.

Remark 3.3. Lemma 3.3 implies that the line search approach ensures a general energy dissipation property associating with the constant $M \ge 0$. For M = 0, the energy sequence $\{E(X^k)\}$ is monotone decreasing. For M > 0, the Algorithm 2 can find a controllable subsequence with energy dissipation.

3.3 Convergence analysis

In this subsection, we give a rigorous proof of energy and sequence convergence for Algorithm 2. Since $m_0 = 0$, the sequence $\{X^k\}$ generated by Algorithm 2 is contained in the sub-level set $[E \le E(X^0)]$. From the Assumption 3.1, we know $[E \le E(X^0)]$ is compact. Together with Lemma 3.1, we have $\{X^k\} \subseteq [E \le E(X^0)] \cap \prod_{j=1}^s S_j \subseteq \mathcal{B}(X^0)$ where $\mathcal{B}(X^0)$ is the closed ball that contains $[E \le E(X^0)] \cap \prod_{j=1}^s S_j$. The next lemma establishes that E(X) is Lipschitz continuous on $\mathcal{B}(X^0)$.

Lemma 3.4. Suppose Assumption 3.1 and Assumption 3.2 hold. Then, there exists $L_E > 0$ such that E(X) is L_E -Lipschitz continuous on $\mathcal{B}(X^0)$ for all X^0 .

Proof. Since $\mathcal{B}(X^0)$ is a compact subset of \mathbb{C}^N , then $\pi_j(\mathcal{B}(X^0)) \subseteq \mathbb{C}^{N_j}$ is closed. It's easy to know that $\pi_j(\mathcal{B}(X^0))$ is bounded by the fact that $||\mathbf{x}_j|| \leq ||X||$. Thus, g_j is L_g^j -Lipschitz continuous on $\pi_j(\mathcal{B}(X^0)) \subseteq$ ridom g_j [9, Corollary 8.41]. As a result, we have

$$\left| \sum_{j=1}^{s} g_{j}(\mathbf{x}_{j}) - \sum_{j=1}^{s} g_{j}(\mathbf{y}_{j}) \right| \leq \sum_{j=1}^{s} |g_{j}(\mathbf{x}_{j}) - g_{j}(\mathbf{y}_{j})| \leq \sum_{j=1}^{s} L_{g}^{j} ||\mathbf{x}_{j} - \mathbf{y}_{j}||$$
$$\leq s \left(\max_{j=1,2,\cdots,s} L_{g}^{j} \right) ||X - Y||, \quad \forall X, Y \in \mathcal{B}(X^{0}).$$

Together with the fact that $F \in C^1(\mathbb{C}^N)$, we conclude that there exists $L_E > 0$ such that $E = F + \sum_{i=1}^{s} g_i$ is L_E -Lipschitz continuous on the compact set $\mathcal{B}(X^0)$.

Lemma 3.5. Suppose Assumption 3.1 and Assumption 3.2 hold. Let $\{X^k\}$ be the sequence generated by Algorithm 2, there exists $E^* > -\infty$ such that

$$\lim_{k \to \infty} \|X^{k+1} - X^k\| = 0, \quad \lim_{k \to \infty} E(X^k) = E^*.$$
(3.15)

Proof. We show the proof similar to the framework in [35]. Since $\{E(X^{m_k})\}$ is non-increasing from Lemma 3.3 and E(X) is bounded below, there exists $E^* > -\infty$ such that $\lim_{k\to\infty} E(X^{m_k}) = E^*$. Define $d^k = X^{k+1} - X^k$ and combine (3.13) with (3.14), we have

$$E(X^{m_k}) - E(X^{k+1}) \ge \sigma \|d^k\|.$$
(3.16)

Assume $k \ge M$, we prove that the following relations hold for any finite $j \ge 1$ by induction

$$\lim_{k \to \infty} \|d^{m_k - j}\| = 0, \quad \lim_{k \to \infty} E(X^{m_k - j}) = E^*.$$
(3.17)

Substituting *k* by $m_k - 1$ in (3.16), we obtain

$$\sigma \|d^{m_k-1}\|^2 \leq E(X^{m_{m_k-1}}) - E(X^{m_k}) \leq E(X^{m_{k-M-1}}) - E(X^{m_k}),$$

where the last inequality holds since $m_k \ge k - M$. Let $k \rightarrow \infty$, we get

$$\lim_{k\to\infty}\|d^{m_k-1}\|=0$$

Since *E* is *L*_{*E*}-Lipschitz continuous on $\mathcal{B}(X^0)$ and $\{X^k\} \subseteq \mathcal{B}(X^0)$, we have

$$|E(X^{m_k-1}) - E(X^{m_k})| \le L_E ||X^{m_k-1} - X^{m_k}|| = L_E ||d^{m_k-1}|| \to 0 \quad (k \to \infty).$$

Thus, one has

$$\lim_{k\to\infty} E(X^{m_k-1}) = \lim_{k\to\infty} E(X^{m_k}) = E^*,$$

which implies that (3.17) holds for j=1. Suppose now that (3.17) holds for some j>1, we show that it also holds for j+1. Substituting k by m_k-j-1 in (3.16), it gives

$$\sigma \|d^{m_k-j-1}\|^2 \leq E(X^{m_{m_k-j-1}}) - E(X^{m_k-j}) \leq E(X^{m_{k-M-j-1}}) - E(X^{m_k-j}).$$

Together with (3.17), it means that

$$\lim_{k\to\infty}\|d^{m_k-j-1}\|=0.$$

Similarly, we have

$$\lim_{k \to \infty} E(X^{m_k - j - 1}) = \lim_{k \to \infty} E(X^{m_k - j} - d^{m_k - j}) = \lim_{k \to \infty} E(X^{m_k - j}) = E^*.$$

Then (3.17) is also holds for j+1. By induction, we prove that relations (3.17) hold for any finite $j \ge 1$.

From the definition of m_k , we have the fact that

$$m_{k+1} - m_k \le M + 1, \quad \forall k \ge 0,$$
 (3.18)

Thus, we obtain $\{X^k\}_{k=M+1}^{\infty} \subset \bigcup_{j=0}^{M+1} \{X^{m_k-j}\}_{k=M+1}^{\infty}$. Using (3.17), it follows that

$$\lim_{k\to\infty} E(X^{m_k-j}) = E^*, \quad \forall j \in [0,M].$$

Then, we have

$$\lim_{k\to\infty} E(X^k) = E^*.$$

Moreover, we know

$$\lim_{k\to\infty} \|d^k\|^2 \leq \frac{1}{\sigma} \lim_{k\to\infty} \left(E(X^{m_k}) - E(X^k) \right) \to 0 \quad (k\to\infty).$$

Then, we obtain $\lim_{k\to\infty} ||d^k|| = 0$.

Remark 3.4. If M=0, Lemma 3.5 can be easily obtained since $m_k \equiv k$ and we do not require to assume $\pi_i(\mathcal{B}(X^0)) \subseteq \text{ridom} g_i$ in Assumption 3.1.

Lemma 3.6. Suppose Assumption 3.1, Assumption 3.2 and Assumption 3.3 hold. Let $\{X^k\}$ be the sequence generated by Algorithm 2. Then, there exists a positive constant C such that

dist
$$(0,\partial E(X^k)) \le C \sum_{l=k-3T+1}^k ||X^k - X^{k-1}||, \quad \forall k > 3T,$$
 (3.19)

where dist $(\mathbf{0},\partial E(X^k))$:= inf $\{\|y\|: y \in \partial E(X^k)\}$.

Proof. If $X^k = X^{k-1}$, we only need to consider (3.19) holds at X^{k-1} and it is easy to know $X^{k-2} \neq X^{k-1}$ from the monotonicity when $w_{k-1}=0$ as shown in Lemma 3.2. Thus, we only consider the case $X^k \neq X^{k-1}$.

It is noted that $\partial E(X) = \nabla f(X) + U$ where $U = (u_1, \dots, u_s)$, $u_i \in \partial g_i(X)$, $1 \le i \le s$. We first assume that the non-restart condition (3.10) is satisfied at the *k*-th iteration. For each $i \in \{1, 2, \dots, s\}$, we denote l_i^k as the last iteration at which the update of *i*-th block is achieved within the first *k*-th iteration, i.e., $l_i^k = \operatorname{argmax}\{\ell | b_\ell = i, \ell \le k\}$. Note that $l_{b_k}^k = k$ and $x_i^{l_k^k} = \tilde{x}_i^{n_k^k}$. By the optimal condition of the proximal subproblem (3.3), we have

$$\mathbf{0} \in \partial g_{i}(\mathbf{x}^{l_{i}^{k}}) + \nabla_{i} f(\mathbf{y}^{l_{i}^{k}-1}, \mathbf{x}^{l_{i}^{k}-1}_{\neq i}) + \frac{1}{\alpha_{l_{i}^{k}-1}} \left(\nabla h_{i}(\mathbf{x}^{l_{i}^{k}}) - \nabla h_{i}(\mathbf{y}^{l_{i}^{k}-1}) \right).$$
(3.20)

Since that $\mathcal{B}(X^0)$ is compact, we let

$$\rho_h := \max_{j=1,2,\cdots,s} \left(\max_{X \in \mathcal{B}(X^0)} \| \nabla^2 h_j(X) \| \right)$$

146

and ρ_f be the local Lipschitz constant of ∇f on $\mathcal{B}(X^0)$. Due to $\{X^k\} \subseteq \mathcal{B}(X^0)$, we know

$$\boldsymbol{y}^{l_i^k-1} = \tilde{\boldsymbol{x}}_i^{n_i^k-1} + w_{l_i^k-1}(\tilde{\boldsymbol{x}}_i^{n_i^k-1} - \tilde{\boldsymbol{x}}_i^{n_i^k-2}) \in \mathcal{B}_0(X^0).$$
(3.21)

Together with (3.20) and (3.21), we get

$$\inf_{\substack{\boldsymbol{u}_{i} \in \partial g_{i}(\boldsymbol{x}_{i}^{k}) \\ = \left\| \nabla_{i}f(X^{k}) - \nabla_{i}f(\boldsymbol{y}^{l_{i}^{k}-1}, \boldsymbol{x}^{l_{i}^{k}-1}_{\neq i}) - \frac{1}{\alpha_{l_{i}^{k}-1}} \left(\nabla h_{i}(\boldsymbol{x}^{l_{i}^{k}}_{i}) - \nabla h_{i}(\boldsymbol{y}^{l_{i}^{k}-1}) \right) \right\| \\ \leq \left\| \nabla_{i}f(X^{k}) - \nabla_{i}f(\boldsymbol{y}^{l_{i}^{k}-1}, \boldsymbol{x}^{l_{i}^{k}-1}_{\neq i}) \right\| + \frac{1}{\alpha_{l_{i}^{k}-1}} \left\| \nabla h_{i}(\boldsymbol{x}^{l_{i}^{k}}_{i}) - \nabla h_{i}(\boldsymbol{y}^{l_{i}^{k}-1}) \right\| \\ \leq \rho_{f} \left\| X^{k} - (\boldsymbol{y}^{l_{i}^{k}-1}, \boldsymbol{x}^{l_{i}^{k}-1}_{\neq i}) \right\| + \frac{\rho_{h}}{\alpha_{\min}} \left\| \boldsymbol{x}^{l_{i}^{k}}_{i} - \boldsymbol{y}^{l_{i}^{k}-1}_{i} \right\|. \tag{3.22}$$

If $i = b_k$, it follows that $l_i^k = k$. Then, we have

$$\|X^{k} - (\boldsymbol{y}^{l_{i}^{k}-1}, \boldsymbol{x}^{l_{i}^{k}-1}_{\neq i})\| = \|X^{l_{i}^{k}} - (\boldsymbol{y}^{l_{i}^{k}-1}, \boldsymbol{x}^{l_{i}^{k}-1}_{\neq i})\| = \|\boldsymbol{x}^{l_{i}^{k}}_{i} - \boldsymbol{y}^{l_{i}^{k}-1}_{i}\|, \quad i = b_{k}.$$
(3.23)

If $i \neq b_k$, we know

$$\begin{aligned} \|X^{k} - (\boldsymbol{y}^{l_{i}^{k}-1}, \boldsymbol{x}^{l_{i}^{k}-1}_{\neq i})\| &\leq \sum_{l=l_{i}^{k}+1}^{k} \|X^{l} - X^{l-1}\| + \|X^{l_{i}^{k}} - (\boldsymbol{y}^{l_{i}^{k}-1}, \boldsymbol{x}^{l_{i}^{k}-1}_{\neq i})\| \\ &= \sum_{l=l_{i}^{k}+1}^{k} \|X^{l} - X^{l-1}\| + \|\boldsymbol{x}^{l_{i}^{k}}_{i} - \boldsymbol{y}^{l_{i}^{k}-1}_{i}\|. \end{aligned}$$
(3.24)

Combining (3.22), (3.23) and (3.24), we have

$$\inf_{\boldsymbol{u}_i \in \partial g_i(\boldsymbol{x}_i^k)} \| \nabla_i f(X^k) + \boldsymbol{u}_i \| \leq \begin{cases} \rho_1 \| \boldsymbol{x}_i^{l_i^k} - \boldsymbol{y}_i^{l_i^k - 1} \|, & i = b_k, \\ \rho_f \sum_{l=l_i^k + 1}^k \| X^l - X^{l-1} \| + \rho_1 \| \boldsymbol{x}_i^{l_i^k} - \boldsymbol{y}_i^{l_i^k - 1} \|, & i \neq b_k, \end{cases}$$

where $\rho_1 = \rho_f + \rho_h / \alpha_{\min}$. Moreover, for each *i*, we know that $k - l_i^k \leq T$ by the Assumption 3.3 and there exist $k_i^1, k_i^2 \in [k-3T,k]$ such that $\mathbf{x}_i^{k_i^1} = \tilde{\mathbf{x}}_i^{n_i^k-1}, \mathbf{x}_i^{k_i^2} = \tilde{\mathbf{x}}_i^{n_i^k-2}$. Thus, we have $\|\mathbf{x}_i^{l_i^k} - \mathbf{y}_i^{l_i^k-1}\| \leq \|\tilde{\mathbf{x}}_i^{n_i^k} - \tilde{\mathbf{x}}_i^{n_i^k-1} - \tilde{\mathbf{x}}_i^{n_i^k-2}\|$ and

$$dist(\mathbf{0},\partial E(X^{k})) = \inf_{U \in \partial g_{i}} \left(\sum_{i=1}^{s} \|\nabla_{i}F(X^{k}) + \boldsymbol{u}_{i}\|^{2} \right)^{1/2} \leq \sum_{i=1}^{s} \inf_{\boldsymbol{u}_{i} \in \partial g_{i}} \|\nabla_{i}F(X^{k}) + \boldsymbol{u}_{i}\|$$
$$\leq \rho_{f} \sum_{i \neq b_{k}} \sum_{l=l_{i}^{k}+1}^{k} \|X^{l} - X^{l-1}\| + \rho_{1} \sum_{i=1}^{s} \|\boldsymbol{x}_{i}^{l_{i}^{k}} - \boldsymbol{y}_{i}^{l_{i}^{k}-1}\|$$

$$\leq \rho_f(s-1) \sum_{l=k-3T+1}^k \|X^l - X^{l-1}\| + \rho_1 \sum_{i=1}^s \left(\|\tilde{\mathbf{x}}_i^{n_i^k} - \tilde{\mathbf{x}}_i^{n_i^k-1}\| + \bar{w} \|\tilde{\mathbf{x}}_i^{n_i^k-1} - \tilde{\mathbf{x}}_i^{n_i^k-2}\| \right) \\ \leq \rho_f(s-1) \sum_{l=k-3T+1}^k \|X^l - X^{l-1}\| + (1+\bar{w})\rho_1 \sum_{l=k-3T+1}^k \|X^l - X^{l-1}\|.$$

Define $C := (s-1)\rho_f + (1+\bar{w})\rho_1$, we obtain (3.19).

Together with Lemma 3.5 and Lemma 3.6, we immediately have the following subsequence convergent property.

Theorem 3.1. Suppose Assumption 3.1, Assumption 3.2 and Assumption 3.3 hold. Let $\{X^k\}$ be the sequence generated by Algorithm 2, then any limit point X^* of $\{X^k\}$ is a critical point of E, *i.e.*, $\mathbf{0} \in \partial E(X^*)$.

Proof. From Remark 3.2, we know $\{X^k\} \subset [E \leq E(X^0)]$ and thus is bounded. Then, the set of limit points of $\{X^k\}$ is nonempty. For any limit point $X^* = (x_1^*, x_2^*, \dots, x_s^*)$, there exists a subsequence $\{X^{k_j}\}$ such that $\lim X^{k_j} = X^*$. By Lemma 3.5 and Lemma 3.6, we immediately obtain

$$\lim_{j \to \infty} \operatorname{dist}(\mathbf{0}, \partial E(X^{k_j})) = 0.$$
(3.25)

Moreover, from Lemma 3.4, we know E(X) is L_E -Lipschitz smooth on $\mathcal{B}(X^0)$. Using the fact that $\{X^k\} \subseteq \mathcal{B}(X^0)$ and $\mathcal{B}(X^0)$ is compact, we know $X^* \in \mathcal{B}(X^0)$. Thus, we get $E(X^{k_j}) \rightarrow E(X^*)$ as $j \rightarrow \infty$. For any $u_j \in \partial E(X^{k_j})$, we know

$$E(X) \ge E(X^{k_j}) + \langle u_j, X - X^{k_j} \rangle, \quad \forall X \in \operatorname{dom} E,$$

which implies

$$\left\{u: u = \lim_{j \to \infty} u_j, u_j \in \partial E(X^{k_j})\right\} \subset \partial E(X^*)$$

by setting $j \rightarrow \infty$. Together with (3.25), we conclude that $\mathbf{0} \in \partial E(X^*)$.

When M = 0 and E is a KL function [10], the sub-sequence convergence can be strengthen to the whole sequence convergence.

Definition 3.4 (KL function [10]). E(x) is the KL function if for all $\bar{x} \in \text{dom}\partial E := \{x: \partial E(x) \neq z\}$ \emptyset , there exists $\eta > 0$, a neighborhood U of \bar{x} and $\psi \in \Psi_{\eta} := \{\psi \in C[0,\eta) \cap C^{1}(0,\eta), where \psi \text{ is } u \in V_{\eta} := \{\psi \in C[0,\eta) \cap C^{1}(0,\eta), where \psi \}$ concave, $\psi(0) = 0$, $\psi' > 0$ on $(0,\eta)$ such that for all $x \in U \cap \{x : E(\bar{x}) < E(x) < E(\bar{x}) + \eta\}$, the following inequality holds,

$$\psi'(E(x) - E(\bar{x}))\operatorname{dist}(\mathbf{0}, \partial E(x)) \ge 1.$$
(3.26)

148

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Theorem 3.2 (Sequence convergence). Suppose Assumption 3.1, Assumption 3.2, Assumption 3.3 hold. Let $\{X^k\}$ be the sequence generated by Algorithm 2 with M = 0. If E is a KL function, then there exists some X^* such that

$$\lim_{k \to \infty} X^k = X^*, \quad \mathbf{0} \in \partial E(X^*). \tag{3.27}$$

Proof. The proof is in Appendix B.

In the following context, we introduce two kinds of h and present corresponding numerical algorithms for solving (2.6).

4 Application to the CMSH model

As discussed above, let f = F, $g_i = G_i$ and $S_i = \{\hat{\boldsymbol{\varphi}}: e_1^\top \hat{\boldsymbol{\varphi}} = 0\}$, then the problem (2.6) reduces to (3.1). In this section, we apply the Algorithm 2 to solve the finite dimensional CMSH model (2.6). Let $i = b_k$, a key component of efficiently implementing AB-BPG method is fast solving the following constrained subproblem

$$\min_{\hat{\boldsymbol{\phi}}\in\mathcal{S}_{i}}G_{i}(\hat{\boldsymbol{\phi}}) + \langle \nabla_{i}F(\hat{\boldsymbol{\psi}}^{k},\hat{\boldsymbol{\phi}}_{\neq i}^{k}),\hat{\boldsymbol{\phi}}-\hat{\boldsymbol{\phi}}_{i}^{k}\rangle + \frac{1}{\alpha_{k}}D_{h_{i}}(\hat{\boldsymbol{\phi}},\hat{\boldsymbol{\phi}}_{i}^{k}),$$
(4.1)

where $\hat{\psi}^k = (1+w_k)\hat{\phi}_i^k - w_k\hat{\phi}_i^{\text{prev}}$ is the extrapolation and $\hat{\phi}_i^{\text{prev}}$ is the value of $\hat{\phi}_i$ before it is updated to $\hat{\phi}_i^k$. Solving (4.1) depends on the form of h_i . In the following context, we let

$$h_j(\mathbf{x}) := h(\mathbf{x}) = \frac{a}{4} \|\mathbf{x}\|^4 + \frac{1}{2} \|\mathbf{x}\|^2, \quad a \ge 0$$
(4.2)

for all j and propose two classes of numerical algorithms for solving (2.5) based on different choices of a. For each algorithm, we will show that the subproblem (4.1) is well defined and can be solved efficiently.

Case I: a = 0. In this case, $h_j(x) = ||x||^2/2$ and the Bregman divergence of D_h becomes the Euclidean distance, i.e.,

$$D_h(x,y) = \frac{1}{2} ||x - y||^2.$$
(4.3)

Thus, the iteration scheme (4.1) is reduced to the accelerated block proximal gradient method [6, 53]. We can find the closed-form solution of the constrained minimization problem (4.1) by applying Lemma 4.1 in [27].

Lemma 4.1. Given $\alpha_k > 0$, Φ^k and $\hat{\psi}^k \in S_i$, the subproblem (4.1) with $h(\mathbf{x}) = ||\mathbf{x}||^2/2$ is welldefined and has analytical solution

$$\hat{\boldsymbol{\phi}}_{i}^{k+1} = (\alpha_{k}\mathcal{D}_{i}+I)^{-1} \left(\hat{\boldsymbol{\psi}}^{k} - \alpha_{k}\mathcal{P}_{1}\nabla_{i}F(\hat{\boldsymbol{\psi}}^{k}, \hat{\boldsymbol{\phi}}_{\neq i}^{k}) \right),$$
(4.4)

where \mathcal{D}_i is a N_i -order diagonal matrix as defined in (A.7) and $\mathcal{P}_1 = I - e_1 e_1^{\top}$ is the projection onto \mathcal{S}_i .

From the feasibility assumption, it is noted that $\hat{\psi}^k \in S_i$ holds as long as the initial point $\Phi^0 \in \prod_{i=1}^s S_i$. The concrete algorithm is given in Algorithm 3 with K = 2.

Case II: a > 0. Since $F(\hat{\Phi})$ can be represented as a 4^{th} -degree polynomial function, it is known that *F* is not block wise relatively smooth with respect to $h(x) = ||x||^2/2$. In this case, we choose a, b > 0. The next lemma shows the optimal condition of minimizing (4.1), which can be obtained from Lemma 4.2 in [27].

Lemma 4.2. Given $\alpha_k > 0$, Φ^k and $\hat{\psi}^k \in S_i$, the subproblem (4.1) with $h(\mathbf{x}) = \frac{a}{4} ||\mathbf{x}||^4 + \frac{b}{2} ||\mathbf{x}||^2 + 1$ (*a*,*b* > 0) is well-defined and has the closed form as follows

$$\hat{\boldsymbol{\psi}}_{i} = [\alpha_{k}\mathcal{D}_{i} + (ap^{*} + b)I]^{-1} \left(\nabla h(\hat{\boldsymbol{\psi}}^{k}) - \alpha_{k}\mathcal{P}_{1}\nabla F(\hat{\boldsymbol{\psi}}^{k}, \hat{\boldsymbol{\phi}}_{\neq i}^{k}) \right),$$
(4.5)

where \mathcal{D}_i is given in (A.7) and p^* is the fixed point of $p = \|\hat{\phi}_i^{k+1}\|^2 := r(p)$.

It is noted that the iterate (4.5) requires solving a nonlinear scalar equation, which can efficiently be solved by many existing solvers. In our implementation, the Newton method is used. The concrete algorithm is given in Algorithm 3 with K = 4.

Algorithm 3 AB-BPG-K method for PFC model

Require: $\hat{\Phi}^0 = \hat{\Phi}^{-1} \in \prod_{i=1}^s S_i$, $\rho \in (0,1)$, $\sigma \ge \eta > 0$ and $w_0 = 0$, $\bar{w}, \alpha_0, M > 0$, k = 0. 1: while stopping criterion is not satisfied do Pick $i = b_k \in \{1, 2, \dots, s\}$ in a deterministic or random manner 2: Update $\hat{\psi}^k = (1+w_k)\hat{\phi}_i^k - w_k\hat{\phi}_i^{pre}$ 3: Estimate α_k by Algorithm 1 4: if K = 2 then 5: Calculate $\boldsymbol{z}^{k} = (\alpha_{k}\mathcal{D}_{i}+I)^{-1}\left(\boldsymbol{\hat{\psi}}^{k}-\alpha_{k}\mathcal{P}_{1}\nabla_{i}F(\boldsymbol{\hat{\psi}}^{k},\boldsymbol{\hat{\phi}}_{\neq i}^{k})\right)$ 6: else if K = 4 then 7: Calculate the fixed point of (4.5). 8: Calculate $z^k = [\alpha_k \mathcal{D}_i + (ap^* + b)I]^{-1} \left(\nabla h(\hat{\psi}^k) - \alpha_k \mathcal{P}_1 \nabla_i F(\hat{\psi}^k, \hat{\phi}^k_{\neq i}) \right)$ 9: end if 10: if $E(\hat{\Phi}^{m_k}) - E(z^k, \hat{\phi}^k_{\neq i}) \ge \sigma \| \hat{\phi}^k_i - z^k \|^2$ then 11: $\hat{\boldsymbol{\phi}}_{i}^{k+1} = \boldsymbol{z}^{k}, \, \hat{\boldsymbol{\phi}}_{i}^{k+1} = \hat{\boldsymbol{\phi}}_{i}^{k} \, (j \neq i) \text{ and choose } w_{k+1} \in [0, \bar{w}].$ 12: 13: else Restart by setting $\hat{\Phi}^{k+1} = \hat{\Phi}^k$ and $w_{k+1} = 0$. 14: 15: end if k = k + 1.16: 17: end while

4.1 Convergence analysis of Algorithm 3

The convergence analysis of Algorithm 2 can be directly applied for Algorithm 3 if all required assumptions in Theorem 3.1 are satisfied. We first show that the energy function

E of CMSH model satisfies Assumption 3.1, Assumption 3.3. Then, Assumption 3.2 is analyzed for Case (P2) and Case (P4) independently.

Lemma 4.3. Let $E(\hat{\Phi}) = F(\hat{\Phi}) + \sum_{i=1}^{s} G_i(\hat{\phi}_i)$ be the energy function of (2.6). Then, it satisfies

- 1. E is bounded below and level bounded,
- 2. ridom $G_i = \mathbb{C}^{N_j}$, thus $\pi_i(\mathcal{B}(\hat{\Phi}^0)) \subseteq \operatorname{ridom} G_i$ for all j.

Proof. From the continuity and the coercive property of *F*, i.e., $F(\hat{\Phi}) \rightarrow +\infty$ as $\hat{\Phi} \rightarrow \infty$, the sub-level set $[E \leq \alpha]$ is compact for any $\alpha \in \mathbb{R}$. Since dom $G_i = \mathbb{C}^{N_j}$, we directly get that ridom $G_j = \mathbb{C}^{N_j}$.

Lemma 4.4. Let $F(\hat{\Phi})$ be defined in (2.4). Then, we have

- 1. If h is chosen as case (P2), then F is block-wise relative smooth with respect to $h_i \equiv h$ in any compact set $[E \leq E(X^0)]$.
- 2. If h is chosen as the case (P4), then F is block-wise relatively smooth to $h_i \equiv h$.

Proof. Denote $\hat{\Phi}^{\otimes k} := \hat{\Phi} \otimes \hat{\Phi} \otimes \cdots \otimes \hat{\Phi}$ where \otimes is the tensor product. Then, $F(\hat{\Phi})$ is the 4^{th} -degree polynomial, i.e., $F(\hat{\Phi}) = \sum_{k=2}^{4} \langle \mathcal{A}_k, \hat{\Phi}^{\otimes k} \rangle$ where the k^{th} -degree monomials are arranged as a k-order tensor \mathcal{A}_k . For any compact set $[E \leq E(X^0)], \nabla F$ is bounded and thus F is block-wise relative smooth with respect to any polynomial function in $[E \leq E(X^0)]$ which includes case (P2). Moreover, for any fixed $\hat{\phi}_{\neq j}$, $F_j(\hat{\psi}) = F(\hat{\psi}, \hat{\phi}_{\neq j})$ is still a 4^{th} -degree polynomial. When h is chosen as (P4), according to Lemma 2.1 in [33], there exists $R_F^j > 0$ such that $F_j(\hat{\psi})$ is R_F^j -smooth relative to h.

Combining Lemma 4.3, Lemma 4.4 with Theorem 3.1, we can directly give the convergence analysis of Algorithm 3.

Theorem 4.1. Let $E(\hat{\Phi}) = F(\hat{\Phi}) + \sum_{j=1}^{s} G_j(\hat{\phi}_j)$ be the energy function which is defined in (2.6). Then any limit point $\hat{\Phi}^*$ of $\{\hat{\Phi}^k\}$ is a critical point of E, i.e., $\nabla E(\hat{\Phi}^*) = 0$, if the sequence $\{\hat{\Phi}^k\}$ generated by Algorithm 3 satisfies one of the following conditions:

- 1. K = 2 in Algorithm 3 and $\{\hat{\Phi}^k\}$ is bounded.
- 2. K = 4 in Algorithm 3.

It is noted that when *h* is chosen as (P2), we cannot bounded the growth of *F* as *F* is a fourth order polynomial. Thus, the boundedness assumption of $\{\hat{\Phi}^k\}$ is imposed which is similar to the requirement as the semi-implicit scheme [45].

If M=0, Theorem 3.2 implies that the sub-sequence convergence can be strengthened by the requirement of the KL property of *E*. According to Example 2 in [10], it is easy to know that $E(\hat{\Phi})$ in our model is a semi-algebraic function, then it is a KL function by Theorem 2 in [10]. Thus, the sequence convergence of Algorithm 3 with M=0 is obtained as follows. **Theorem 4.2.** Let $E(\hat{\Phi}) = F(\hat{\Phi}) + \sum_{j=1}^{s} G_j(\hat{\phi}_j)$ be the energy function defined in (2.6) and M=0in Algorithm 3. If the sequence $\{\hat{\Phi}^k\}$ generated by Algorithm 3 satisfies one of the conditions in Theorem 4.1, then there exists some $\hat{\Phi}^*$ such that $\lim_{k \to \infty} \hat{\Phi}^k = \hat{\Phi}^*$ and $\nabla E(\hat{\Phi}^*) = 0$.

From the above analysis, we know that the proposed AB-BPG method has the proven convergence to some stationary point compared to the current gradient flow based methods. Moreover, it has shown that the generated sequence has the (generalized) energy dissipation and mass conservation properties.

5 Numerical results

In this section, we apply the AB-BPG-K approaches (Algorithm 3) to binary, ternary, quinary component systems based on the CMSH model. The efficiency and accuracy of our methods are demonstrated through comparing with existing methods, including the first-order semi-implicit scheme (SIS), the second-order Adam-Bashforth with Lagrange extrapolation approach (BDF2) [26], the scalar auxiliary variable (SAV) method [43], the stabilized scalar auxiliary variable (S-SAV). Note that these employed methods all guarantee the equality constraint, i.e., mass conservation.

The step sizes α_k in the AB-BPG approaches are adaptively obtained via the linear search technique. To be fair, adaptive time stepping are applied to gradient flow methods. The SIS and BDF2 scheme use the adaptive time stepping [55]

$$\alpha_k = \max\left\{\alpha_{\min}, \frac{\alpha_{\max}}{\sqrt{1+\rho|E'(\Phi^k)|^2}}\right\},\tag{5.1}$$

where *E* is the energy functional defined in the model. The constant $\alpha_{\min}, \alpha_{\max}$, and ρ are set to $\alpha_{\min} = 0.001, \alpha_{\max} = 0.1$ and $\rho = 50$ as [55] suggested. The adaptive SAV scheme is implemented as Algorithm 1 in [44], where a semi-implicit first-order SAV scheme and a semi-implicit second-order SAV scheme based on Crank-Nicolson are computed at one step iteration. The adaptive time stepping is in the form of

$$\alpha_{k+1} = \max\left\{\alpha_{\min}, \min\left\{\rho\left(\frac{\operatorname{tol}}{e_{k+1}}\right)^{1/2}\alpha_k, \alpha_{\max}\right\}\right\},\tag{5.2}$$

where ρ is a default safety coefficient, tol is a reference tolerance, e_{k+1} is the relative error between a first-order SAV scheme and a second-order SAV scheme at each time level. e_{n+1} is the relative error between first-order SAV scheme and second-order scheme at each time level. The original parameters (ρ =0.9, tol=10⁻³, α_{min} =10⁻⁵, α_{max} =10⁻²) taken in work [44] are inefficient to compute stationary states. In our implementation, we carefully choose appropriate parameters case by case to ensure better numerical performance. Compared to the adaptive SAV scheme, the S-SAV scheme adds a first-order stabilization term $S_1(\Phi^{k+1} - \Phi^k)$ and a second-order stabilization term $S_2(\Phi^{k+1} - 2\Phi^k + \Phi^{k-1})$ [45] to the SAV-SI scheme and the SAV-CN scheme, respectively. In our computation, we let $S_1 = S_2 = 10$ and use the adaptive times stepping (5.2) with $\rho = 0.9$, tol = 10^{-3} , $\alpha_{\min} = 10^{-5}$, $\alpha_{\max} = 1$.

Our methods, the adaptive SIS and the adaptive BDF2 use a Gauss-Seidel manner to update order parameters in a fixed cyclic order. While the adaptive SAV and adaptive S-SAV method keep the same as Algorithm 1 in [44], which uses a Jacobian manner to update all order parameters simultaneously in one step iteration. In our implementation, all approaches are stopped when $\|\nabla E(\hat{\Phi})\|_{\infty} < 10^{-7}$ or the energy difference between two iterations is less than 10^{-14} . All experiments were performed on a workstation with a 3.20 GHz CPU (i7-8700, 12 processors). All codes were written by MATLAB without parallel implementation.

5.1 Binary component systems

We first choose s = 2 in (2.2) and take the two-dimensional decagonal quasicrystal as an example to examine our approaches' performance. The decagonal quasicrystal can be embedded into a four-dimensional periodic structure. Therefore, we carry out the projection method in four-dimensional space. The 4-order invertible matrix B associated with the four-dimensional periodic structure is chosen as I_4 . The corresponding computational domain in physical space is $[0,2\pi)^4$. The projection matrix \mathcal{P} in (A.5) of the decagonal quasicrystals is

$$\mathcal{P} = \left(\begin{array}{ccc} 1 & \cos(\pi/5) & \cos(2\pi/5) & \cos(3\pi/5) \\ 0 & \sin(\pi/5) & \sin(2\pi/5) & \cos(3\pi/5) \end{array}\right).$$
(5.3)

We use 38⁴ plane wave functions to discretize the binary CMSH energy functional. The parameters in (2.2) are given in Table 2. The initial configuration of order parameters is chosen as references [27, 28] suggest. The stationary quasicrystals, including physical space morphology and Fourier spectra, are given in Fig. 1.

Table 2: The non-zero model parameters used in computing binary decagonal quasicrystal.

$$c = 20, q_1 = 1, q_2 = 2\cos(\pi/5), \tau_{0,2} = \tau_{2,0} = -0.1, \tau_{0,3} = \tau_{3,0} = -0.3, \tau_{1,2} = \tau_{2,1} = -2.2, \tau_{0,4} = \tau_{4,0} = \tau_{1,1} = \tau_{2,2} = \tau_{1,3} = \tau_{3,1} = 1.$$

5.1.1 Algorithm study

In this subsection, we take the binary CMSH model as an example to show our proposed AB-BPG method's performance by choosing different hyperparameters, including the choice of the step sizes, the descent subsequence, and the block update manner. Similar results exist in the following ternary and quinary cases. For simplicity, we only present the results in the computing binary CMSH model.



Figure 1: The stationary decagonal quasicrystal in the binary CMSH model. (a): The spatial distribution of density, where the red and blue colors correspond to the rich concentration of ϕ_1 and ϕ_2 , respectively; (b): The Fourier spectra of ϕ_1 ; (c): The Fourier spectra of ϕ_2 . Only Fourier spectral points whose coefficient intensities are larger than 10^{-4} are presented.

The step size of the AB-BPG methods is adaptively obtained by Algorithm 1. In our implementation, we set $\alpha_0 = 0.1$, $\varsigma = (\sqrt{5}-1)/2$, $\eta = 10^{-12}$, $\alpha_{\min} = 10^{-6}$ and $\alpha_{\max} = 10$ in Algorithm 1. Fig. 2(a) illustrates the adaptive step sizes versus the iteration when M = 0 and a = 1.

In Fig. 2(b), it shows the total energy of the sequence $\{\hat{\Phi}^k\}$ and the subsequence $\{\hat{\Phi}^{m_k}\}$ defined in (3.9) versus the iterations with M = 5 and a = 0. It is observed that the subsequence $\{E(\hat{\Phi}^{m_k}\})$ is monotone decreasing which validates the generalized energy dissipation property proved in Lemma 3.3.

In the next, we consider the different choices of the update blocks. One is the cyclic rule that updates $\hat{\Phi}_1$ and $\hat{\Phi}_2$ alternatively in the Gauss-Seidel fashion. The other is to randomly choose the update block, and $\hat{\Phi}_1$ and $\hat{\Phi}_2$ should be updated at least once in 10 consecutive iterations. In this case, we set a = M = 0, and we independently run the random update rule 50 times and report its terminated iteration numbers for each trial. The result is shown in Fig. 2(c). Compared to the random update manner, it is observed that the cyclic rule needs fewer iterations for the desired accuracy. Therefore, in the following simulations, we only consider the cyclic update order when carrying out the AB-BPG algorithms.

5.1.2 Comparison with other methods

We compare AA-BPG methods with alternative methods, including adaptive SIS, adaptive BDF2, adaptive SAV and adaptive S-SAV. Theoretically, the SAV method always has a modified energy dissipation through adding a sufficiently large positive scalar auxiliary variable *C* which guarantees the boundedness of the bulk energy term. In practice, the original energy dissipation property might depend on the selection of *C*. When computing the decagonal quasicrystal in the binary CMSH model, the adaptive SAV scheme keep the original energy dissipate when $C = 10^8$. The times stepping for adaptive SAV scheme is obtained by formula (5.2) with $\alpha_{\min} = 10^{-5}$, $\alpha_{\max} = 0.5$, $\rho = 0.9$, tol = 10^{-3} . For



Figure 2: (a): The adaptive step size obtained by the AB-BPG-4 (M=0, a=1). The mean step sizes of ϕ_1 and ϕ_2 are 0.6033 and 0.7062, respectively. (b): The tendency of energy $\{E^k = E(\hat{\Phi}^k)\}$ and energy $\{E^{m_k} = E(\hat{\Phi}^{m_k})\}$ in AB-BPG-2 (M=5, a=0). The green o's mark the position of (m_k, E^{m_k}). (c): Numerical behavior of the AB-BPG-2 (M=0, a=0) via updating order parameters in a random manner. The blue dotted line and the pink points denote the convergent iterations of the AB-BPG methods via cyclic and random order update, respectively.

Μ	а	Iterations	CPU Time (s)	Gradient error (10^{-8})	$ E-E_s $ (10 ⁻¹⁴)
	0	54	72.56	9.74	4.07
0	0.1	46	66.04	9.42	51.54
	1	45	69.17	7.95	13.78
	0	65	86.62	4.91	0.76
5	0.1	59	85.10	9.25	23.08
	1	57	82.79	7.99	18.62
	0	65	85.88	4.01	0.06
10	0.1	59	85.31	9.60	95.61
	1	56	79.38	7.13	2.98
Adaptive SIS		293	163.10	9.98	152.28
Adaptive BDF2		305	276.98	9.77	142.73
Adaptive SAV		94	283.14	8.86	117.56
Adaptive S-SAV		236	742.46	9.81	141.65

Table 3: Numerical results of computing binary decagonal quasicrystal.

the AB-BPG approaches, we choose different values of M in (3.9) and a in Bregman divergence (4.2) for comparison. The linear search technique can obtain the step size of the AB-BPG methods.

Table 3 shows the corresponding numerical results of the AB-BPG, adaptive SIS, adaptive BDF2, adaptive SAV and adaptive S-SAV schemes. We choose a reference energy value $E_s = -1.54929536255898 \times 10^{-2}$, which is computed via the semi-implicit scheme



Figure 3: Numerical behaviors of our algorithms, adaptive SIS, adaptive BDF2, adaptive SAV, adaptive S-SAV for computing binary decagonal quasicrystal. First row: M=0; Second row: M=5; Left: Relative energy over iterations; Middle: Relative energy over CPU time; Right: Gradient error over iterations; The green \times s mark where restarts occurred.

using 56⁴ plane wave functions. The minimal iteration steps and the least CPU time are emphasized via the bold font. From Table 3, one can find that the proposed approaches are superior to other methods. In particular, when a = 1, the AB-BPG method with M = 0takes 45 iterations (69.17 seconds) to reduce the gradient error of 10^{-7} , which is 2.35 times faster than the adaptive SIS, 4 times than the adaptive BDF2 and adaptive SAV schemes, 10.73 times than the adaptive S-SAV scheme. Moreover, a great deal of the adaptive time stepping for gradient flow methods is that empirical parameters are involved in the formula (5.1) and (5.2), where the best-performing parameters usually cannot be guided in advance. Unsuitable choice of parameters may lead to too small step size (inefficient performance) or too large step size (divergence). While in our AB-BPG methods, efficient performances are shown in a wild range of choice of parameters. More importantly, the original energy dissipation and convergence are both guaranteed within any choice of parameters in our methods. Fig. 3 presents the iteration process of relative energy difference, CPU time, and gradient error of different approaches.

5.2 Ternary component systems

We consider the ternary component system when s = 3 in the CMSH model. A periodic structure of the sigma phase, which is a complicated spherical packed structure discovered in multicomponent material systems [31], is used to examine the performance of

Table 4: The non-zero model parameters used in computing ternary sigma phase. $c=1, q_1=q_2=q_3=1, \tau_{2,0,0}=\tau_{0,2,0}=\tau_{0,0,2}=-0.2, \tau_{3,0,0}=\tau_{0,3,0}=\tau_{0,0,3}=-0.3,$





Figure 4: The stationary sigma phase in ternary CMSH model from two perspectives.

our algorithms. For such a pattern, we implement our algorithm on a bounded computational domain $[0,27.7884) \times [0,27.7884) \times [0,14.1514)$ and $200 \times 200 \times 100$ plane wave functions are used to for discretization. The initial value of each component of the sigma phase is input, as suggested in [4,51]. The parameters in the ternary CMSH model are given in Table 4. The stationary sigma phase is present in Fig. 4.

Similar to the binary case, we report the AB-BPG algorithms' performance with different choices of *a* and *M*. In the adaptive SAV method, the auxiliary parameter *C* is set to 10^8 and the parameters in formula (5.2) are chosen as $\alpha_{\min} = 10^{-5}$, $\alpha_{\max} = 0.2$, $\rho = 0.8$, tol = 10^{-3} . Table 5 presents the corresponding numerical results and Fig. 5 gives the iteration process of different methods. The reference energy value $E_s = -1.16910245253091$ is obtained numerically via semi-implicit scheme by using $256 \times 256 \times 128$ plane wave functions. As is evident from these results, the AB-BPG methods demonstrate a great advantage in computing such a complicated periodic structure over other schemes. More precisely, when M = 5 and a = 1, the AB-BPG method spends 416 iterations to achieve the prescribed error, which converges 16.5 times faster than the adaptive SIS and adaptive SDF2 methods, 3 times than the adaptive SAV method and 6.5 times than the adaptive S-SAV method. From the cost of CPU time, the AB-BPG algorithm with M = 5, a = 1 takes 1358.93 seconds to achieve an accuracy of 10^{-7} in error, almost 10.7%, 10.8%, 11.9% and 6.5% the time employed by the adaptive SIS, adaptive BDF2, adaptive SAV and adaptive S-SAV schemes.

5.3 Quinary component systems

The last multicomponent system considered in this paper is the five component CMSH model. We take a two-dimensional chessboard-shaped tiling phase and three-

М	а	Iterations	CPU Time (s)	Gradient error (10^{-8})	$ E-E_s $ (10 ⁻¹²)
	0	758	2185.19	6.47	3.22
0	0.1	770	2549.49	9.15	3.57
	1	593	1952.63	9.71	4.79
	0	739	2068.65	9.50	3.54
5	0.1	611	1991.94	8.23	2.01
	1	416	1358.93	8.51	0.68
	0	677	1896.60	7.13	2.40
10	0.1	595	1959.80	9.22	0.58
	1	560	1865.38	7.32	1.31
Adaptive SIS		6905	12656.52	9.96	6.28
Adaptive BDF2		6896	12549.01	9.99	5.96
Adaptive SAV		1276	11406.12	9.94	6.11
Adaptive S-SAV		2737	20773.66	9.99	5.94

Table 5: Numerical results of computing ternary sigma phase.



Figure 5: Numerical behaviors of the AB-BPG, adaptive SIS, adaptive BDF2, adaptive SAV and adaptive S-SAV algorithms for computing ternary sigma phase. First row: M=0; Second row: M=20; Left: Relative energy over iterations; Middle: Relative energy over CPU time; Right: Gradient error over iterations; The green \times s mark where restarts occurred.

dimensional body-centered cubic (BCC) spherical structure to examine AB-BPG approaches' performance.

5.3.1 Chessboard-shaped tiling

When computing the chessboard-shaped tiling, the parameters in five-component CMSH model are given in Table 6. The corresponding computational domain in physical space is $\Omega = [0,2\pi)^2$ and 1024×1024 plane wave functions are used to discretize the computational domain. The initial solution of *j*-th component is

$$\phi_j(\mathbf{r}) = \sum_{\mathbf{h} \in \Lambda_0^j} \hat{\phi}(\mathbf{h}) e^{i\mathbf{h}^\top \mathbf{r}}, \quad \mathbf{r} \in \Omega,$$
(5.4)

where initial lattice points set $\Lambda_0^j \subset \mathbb{Z}^2$ can be found in the Table 7 only on which the Fourier coefficients $\hat{\phi}(h)$ located are nonzero. The convergent stationary morphology is given in Fig. 6.

Table 8 presents the numerical results of the AB-BPG, adaptive SIS, adaptive BDF2, adaptive SAV and adaptive S-SAV methods. The scalar auxiliary parameter C of adaptive

Table 6: The non-zero model parameters when computing quinary chessboard-shaped tiling.

$$c = 10, q_1 = q_2 = \dots = q_5 = 1,$$

$$\tau_{3,0,0,0,0} = \tau_{0,3,0,0,0} = \tau_{0,0,3,0,0} = \tau_{0,0,0,3,0} = \tau_{0,0,0,0,3} = -0.10,$$

$$\tau_{4,0,0,0,0} = \tau_{0,4,0,0,0} = \tau_{0,0,4,0,0} = \tau_{0,0,0,4,0} = \tau_{0,0,0,0,4} = 0.10,$$

$$\tau_{1,0,1,0,0} = -0.70, \tau_{0,1,0,1,1} = 0.05, \tau_{1,1,0,0,1} = -0.12, \tau_{0,1,0,1,0} = -0.44.$$

Table 7: The initial lattice points of each component when computing quinary chessboard-shaped tiling.

	ϕ_1	ϕ_2	ϕ_3	ϕ_4	ϕ_5
Λ_0^j	(±1,0)	$(0,\pm 1)$	(±2,0)	(0,±2)	(0,0)



Figure 6: The stationary chessboard-shaped phase in quinary CMSH model. The morphologies of ϕ_1, ϕ_3 (left), ϕ_2, ϕ_4 (middle) and ϕ_5 (right).

М	а	Iterations	CPU Time (s)	Gradient error (10^{-8})	$ E-E_s $ (10 ⁻¹⁴)
	0	111	75.56	6.97	2.90
0	0.1	142	101.52	5.44	1.48
	1	204	145.30	6.60	1.19
	0	193	149.67	6.28	0.46
5	0.1	174	131.11	9.65	1.38
	1	239	171.46	9.91	1.72
	0	186	120.03	8.34	2.62
10	0.1	188	135.34	6.56	0.84
	1	253	186.38	5.37	0.52
Adaptive SIS		18208	8957.85	282.01	4617.51
Adaptive BDF2		17723	8277.81	346.25	7073.17
Adaptive SAV		927	2184.01	9.92	94.23
Adaptive S-SAV		6655	16131.94	9.99	94.51

Table 8: Numerical results of computing the quinary chessboard-shape tiling.

SAV scheme is set to 10^{10} and the parameters in formula (5.2) are taken as $\alpha_{\min} = 10^{-5}$, $\alpha_{\max} = 0.7$, $\rho = 0.9$, tol = 10^{-3} . The reference energy value $E_s = -0.57163687783216$ is obtained via semi-implicit scheme by using 2048×2048 plane wave functions. Correspondingly, Fig. 7 presents the iteration process including the relative energy difference, the CPU times and the gradient error against iterations, respectively. These results demonstrate the superiority of the AB-BPG algorithms over the adaptive SIS, adaptive BDF2 and adaptive SAV methods. As Table 8 shows, the AB-BPG method with M = 0, a = 0 has the best performance. Even though the AB-BPG method costs much time as the adaptive SIS and adaptive BDF2 approaches per iteration due to the linear search technique, its adaptive step size compensates the extra work converging to 10^{-7} , almost 0.8%, 0.9%, 3.5% and 0.4% of CPU time spent by the adaptive SIS, adaptive BDF2, adaptive SAV methods, and adaptive S-SAV, respectively.

5.3.2 BCC

As last, we consider the quinary BCC structure. The parameters in the five component CMSH model are given in Table 9. A bounded domain $[0,2\sqrt{2}\pi)^3$ is used as the computational box and 128³ plane wave functions are employed to compute the BCC spherical phase. The initial values can be found in [29]. Fig. 8 shows the convergent stationary solution of different order parameter, which are all BCC phases but with different periodicity.

Table 10 and Fig. 9 compare the numerical behaviors of AB-BPG, adaptive SIS, adaptive BDF2, adaptive SAV and adaptive S-SAV methods. The scalar auxiliary parameter



Figure 7: Numerical behaviors of the AB-BPG and adaptive SIS, adaptive BDF2, adaptive SAV and adaptive S-SAV for computing the quinary chessboard-shape tiling. First row: M=0; Second row: M=5; Left: Relative energy over iterations; Middle: Relative energy over CPU times; Right: Gradient error over iterations; The green \times s mark where restarts occurred.

Table 9: The non-zero model parameters used in computing quinary BCC spherical structure.

$$c = 1, q_1 = 1, q_2 = 1.5, q_3 = 2, q_4 = 2.5, q_5 = 3,$$

$$\tau_{3,0,0,0,0} = -0.1, \tau_{0,3,0,0,0} = -0.6, \tau_{0,0,3,0,0} = -0.4, \tau_{0,0,0,3,0} = -0.2, \tau_{0,0,0,0,3} = -0.1,$$

$$\tau_{4,0,0,0,0} = \tau_{0,0,4,0,0} = \tau_{0,4,0,0,0} = \tau_{0,0,0,4,0} = \tau_{0,0,0,0,4} = 0.1,$$

$$\tau_{1,0,1,0,0} = 0.4, \tau_{0,1,0,1,0} = 0.3, \tau_{0,1,1,1,0} = -0.2, \tau_{1,1,0,0,1} = 0.8.$$

C of two SAV scheme is set to 10^{10} . For adaptive SAV scheme, the parameters in formula (5.2) are taken as $\alpha_{\min} = 10^{-5}$, $\alpha_{\max} = 0.2$, $\rho = 0.9$, tol = 10^{-3} . The reference energy $E_s = -1.22314417498279$ is obtained via semi-implicit scheme by using 256³ plane wave functions. Again, for this case, the proposed AB-BPG methods are still superior to the compared algorithms. More precisely, the best performance of AB-BPG (M = 0, a = 0) spends 182 iterations to achieve the prescribed error which converges 40 times faster than the adaptive SIS, 30 times than the adaptive BDF2 method, 5.1 times than the adaptive SAV scheme and 10.1 times than the adaptive S-SAV scheme. From the cost of CPU times, the AB-BPG algorithm with M = 0 takes 247.22 seconds to achieve an accuracy of 10^{-7} in the gradient error, almost 3.8%, 5.1%, 7.0% and 2.8% of the CPU time used by the adaptive SIS, adaptive BDF2, adaptive SAV schemes.



Figure 8: The stationary BCC phase in quinary CMSH model. The subfigures from left to right are morphology of ϕ_1 to ϕ_5 .

М	а	Iterations	CPU Times	Gradient error (10^{-8})	$ E-E_s $ (10 ⁻¹³)
	0	182	247.22	8.17	1.99
0	0.1	212	309.12	9.58	1.71
	1	232	344.20	6.66	1.38
	0	259	343.63	3.96	1.37
5	0.1	262	381.79	5.70	2.14
	1	347	515.89	5.84	1.28
10	0	277	371.94	5.56	1.33
	0.1	304	446.41	8.04	1.40
	1	353	505.98	7.38	1.33
Adaptive SIS		7443	6446.39	52.18	439.44
Adaptive BDF2		5628	4802.81	71.90	832.20
Adaptive SAV		933	3498.08	9.94	17.20
Adaptive S-SAV		2204	8766.82	9.95	17.19

Table 10: Numerical results of computing the quinary BCC structure.

6 Conclusion

In this paper, an AB-BPG algorithm is proposed to compute the stationary states of multicomponent phase-field crystal model with mass conservation. Compared to most existing methods, the new approaches consider the block structure of multicomponent models, and the update manner can be chosen deterministically or randomly. Using modern optimization methods, including the inertia acceleration approach, the restart technique, and the line search method, the proposed AB-BPG method has the general dissipation property with efficient implementation. Moreover, with the help of the Bregman divergence, it is proved that the generated sequence converges to a stationary point without the requirement of the global Lipschitz assumption. Extensive numerical experiments on computing stationary periodic crystals and quasicrystals in the binary, ternary, and



Figure 9: Numerical behaviors of the AB-BPG, adaptive SIS and adaptive BDF2, adaptive SAV and adaptive S-SAV methods for computing the quinary BCC phase. First row: M=0; Second row: M=5; Left: Relative energy over iterations; Middle: Relative energy over the CPU times; Right: Gradient error over iterations; The green \times s mark where restarts occurred.

quinary coupled-mode Swift-Hohenberg model have shown a significant acceleration over many existing methods.

In this paper, we have implemented our algorithms for computing multicomponent model with polynomial type potentials $F[\{\phi_j\}_{j=1}^s]$. In fact, the proposed methods can be applied to deal with non-polynomial type potentials by choosing suitable $\{h_j\}$ to make our Assumption 3.2 satisfied. This will soon be found in our future work. Besides, although our methods have shown efficient performance in our numerical experiments, we want to theoretically prove the convergence rates in the future. It is worth mentioning that in this work, we only consider efficient methods for computing stationary states. Besides, the physical evolution simulation is an important topic in PFC models. It is known that the classical proximal gradient method is exactly the semi-implicit scheme for gradient flows. From this perspective, we will develop our method to simulate the whole evolution of process of crystal growth in the further.

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Appendices

A Projection method

The projection method is a general framework to study the periodic and quasiperiodic crystals [25,30]. Each of the *d*-dimensional periodic system can be described by a Bravais lattice

$$\mathcal{R}_d = \{A_d \boldsymbol{n}_d, \, \boldsymbol{n}_d \in \mathbb{Z}^d\},$$

where $A_d \in \mathbb{R}^{d \times d}$ is invertible. The fundamental domain or so-called unit cell of the periodic system is

$$\Omega = \{ A_d \zeta, \zeta \in \mathbb{R}^d, \zeta_j \in [0,1), j = 1, \cdots, d \}.$$

The associated primitive reciprocal vectors, $B_d = (b_1, \dots, b_d)$, $b_j \in \mathbb{R}^d$, $j = 1, \dots, d$, satisfy the dual relationship $A_d B_d^T = I_d$ where I_d is the *d*-order identify matrix. The reciprocal lattice is

$$\mathcal{R}_d^* = \{ \boldsymbol{B}_d \boldsymbol{h}_d, \, \boldsymbol{h}_d \in \mathbb{Z}^d \}.$$

The corresponding periodic function $\psi(\mathbf{r})$, which has translational invariance with respect to the Bravais lattice \mathcal{R} , i.e., $\psi(\mathbf{r}) = \psi(\mathbf{r} + \mathcal{R})$, has the following Fourier expansion

$$\psi(\mathbf{r}) = \sum_{\mathbf{h} \in \mathbb{Z}^d} \hat{\psi}(\mathbf{h}) e^{i(\mathbf{B}\mathbf{h})^T \mathbf{r}},\tag{A.1}$$

where

$$\hat{\psi}(\boldsymbol{h}) = \int \psi(\boldsymbol{r}) e^{-(\boldsymbol{B}\boldsymbol{h})^T \boldsymbol{r}} d\boldsymbol{r}, \quad \boldsymbol{r} \in \Omega.$$
(A.2)

Quasiperiodic functions, or more general almost periodic functions, are an extension of periodic functions. The definition of quasiperiodic functions is given as follows [25]:

Definition A.1 (Quasiperiodic function). A *d*-dimensional continuous complex-valued function $f(\mathbf{r})$ is quasiperiodic, if there exists a *n*-dimensional periodic function $F(\mathbf{r}_s)$, $\mathbf{r}_s \in \mathbb{R}^n$, $n \ge d$, such that

$$f(\mathbf{r}) = F(\mathcal{P}^T \mathbf{r}),$$

where $\mathcal{P} \in \mathbb{R}^{d \times n}$ is the projection matrix. The column vectors of \mathcal{P} are rationally independent.

From the definition, one can find that the *d*-dimensional quasiperiodic system is a *d*-dimensional subspace of a *n*-dimensional periodic structure. As describes above, the Fourier series of $F(r_s)$ is

$$F(\mathbf{r}_s) = \sum_{\mathbf{h} \in \mathbb{Z}^n} \hat{F}(\mathbf{h}) e^{i(\mathbf{B}\mathbf{h})^T \mathbf{r}_s},$$
(A.3)

where $B \in \mathbb{R}^{n \times n}$ is associated to the periodicity of the *n*-dimensional periodic system. $\hat{F}(h) = \int F(\mathbf{r}_s) e^{-(Bh)^T \mathbf{r}_s} d\mathbf{r}_s, \mathbf{r}_s \in \Omega_s, \Omega_s = \{B^{-T}\alpha, \alpha \in \mathbb{R}^n, \alpha_j \in [0,1), j = 1, \dots, n\}$. Let $\hat{F}(h)$ be $\hat{f}(h)$, the projection method for a *d*-dimensional quasiperiodic function over \mathbb{R}^d can be written as [30]

$$f(\mathbf{r}) = \sum_{\mathbf{h} \in \mathbb{Z}^n} \hat{f}(\mathbf{h}) e^{i(\mathcal{P}B\mathbf{h})^\top \mathbf{r}}, \quad \mathbf{r} \in \mathbb{R}^d.$$
(A.4)

If consider periodic crystals, the projection matrix becomes the *d*-order identity matrix, then the projection reduces to the common Fourier spectral method.

Similarly, the order parameter in multicomponent systems can be expanded as follows

$$\phi_j(\mathbf{r}) = \sum_{\mathbf{h} \in \mathbb{Z}^n} \hat{\phi}_j(\mathbf{h}) e^{i(\mathcal{P}B\mathbf{h})^\top \mathbf{r}}, \quad j = 1, 2, \cdots, s.$$
(A.5)

In numerical implementation, we truncate the Fourier coefficients to satisfy

$$X_{j} = \left\{ \{ \hat{\phi}_{j}(\boldsymbol{h}) \}_{\boldsymbol{h} \in \mathbb{Z}^{n}} : \hat{\phi}_{j}(\boldsymbol{h}) = 0, \forall |h_{l}| > \frac{N_{l,j}}{2}, l = 1, 2, \cdots, n \right\},$$

where $N_{l,j}$ is chosen to be even for convenience. Let $\hat{\phi}_j = (\hat{\phi}_{1,j}, \hat{\phi}_{2,j}, \cdots, \hat{\phi}_{N_j,j})^\top \in \mathbb{C}^{N_j}$ with $N_j = \prod_{l=1}^n (N_{l,j}+1)$. Let $\hat{\Phi} = \{\hat{\phi}_j\}_{j=1}^s \in \mathbb{C}^N$ with $N = \sum_{j=1}^s N_j$. Using the projection method discretization, the energy functional (2.4) reduce to

$$G_{h,j}(\hat{\phi}_{j}) = \frac{1}{2} \sum_{h_{j,1}+h_{j,2}=0} \left[q_{j}^{2} - (\mathcal{P}\mathbf{B}h)^{\top} (\mathcal{P}\mathbf{B}h) \right]^{2} \hat{\phi}_{j}(h_{j,1}) \hat{\phi}_{j}(h_{j,2}), \quad j = 1, 2, \cdots, s,$$

$$F_{h}(\hat{\Phi}) = \sum_{\mathcal{I}_{s,n}} \tau_{i_{1},i_{2},\cdots,i_{s}} \sum_{\sum_{j,k}h_{j,k}=0} \prod_{j=1}^{s} \left(\prod_{k=1}^{i_{j}} \hat{\phi}_{j}(h_{j,k}) \right),$$

where $h_{j,k} \in \mathbb{Z}^n$, $\hat{\phi}_j \in X_j$, $j = 1, 2, \dots, s$. For simplicity, we omit the subscription h in $G_{h,j}$ and F_h in the following context. Then the discretized energy functional can be stated as

$$E(\{\hat{\boldsymbol{\phi}}_{j}\}_{j=1}^{s}) = \sum_{j=1}^{s} G_{j}(\hat{\boldsymbol{\phi}}_{j}) + F(\{\hat{\boldsymbol{\phi}}_{j}\}_{j=1}^{s}),$$
(A.6)

where $G_i(\hat{\phi}_i) = \frac{1}{2} \langle \hat{\phi}_i, \mathcal{D}_i \hat{\phi}_i \rangle$ and $\mathcal{D}_i \in \mathbb{C}^{N_i \times N_j}$ is a diagonal matrix with nonnegative entries

$$(\mathcal{D}_j)_h = [q_j^2 - (\mathcal{P}Bh)^\top (\mathcal{P}Bh)]^2 \text{ with } h \in \mathbb{Z}^n \text{ and } \hat{\phi}_j(h) \in X_j.$$
 (A.7)

 $F(\{\hat{\phi}_j\}_{j=1}^s)$ are *n*-dimensional convolutions in the reciprocal space. In summary, the discretized version of (2.5) has the form

$$\min_{\hat{\Phi}} E(\hat{\Phi}) = \sum_{j=1}^{s} G_j(\hat{\phi}_j) + F(\hat{\Phi}), \quad \text{s.t.} \quad e_1^{\top} \hat{\phi}_j = 0, \quad j = 1, 2, \cdots, s.$$
(A.8)

B Proof of Theorem 3.2

Before we prove the convergent property, we first present a useful lemma for our analysis.

Lemma B.1 (Uniformized Kurdyka-Lojasiewicz property [10]). Let *C* be a compact set and *E* defined in (2.6) be bounded below. Assume that *E* is constant on *C*. Then, there exist $\epsilon > 0$, $\eta > 0$, and $\psi \in \Psi_{\eta}$ such that for all $\bar{u} \in C$ and all $u \in \Gamma_{\eta}(\bar{u}, \epsilon)$, one has,

$$\psi'(E(u) - E(\bar{u})) \operatorname{dist}(\mathbf{0}, \partial E(u)) \ge 1, \tag{B.1}$$

where $\Psi_{\eta} = \{\psi \in C[0,\eta) \cap C^{1}(0,\eta) \text{ and } \psi \text{ is concave, } \psi(0) = 0, \psi' > 0 \text{ on } (0,\eta) \}$ and $\Gamma_{\eta}(x,\epsilon) = \{y | \|x-y\| \le \epsilon, E(x) < E(y) < E(x) + \eta \}.$

Proof. The proof is based on the fact that *E* satisfies the so-called Kurdyka-Lojasiewicz property on C [10].

Proof of Theorem 3.2. Define two sets $\Omega_2 = \{k | w_k = 0\}$ and $\Omega_1 = \mathbb{N} \setminus \Omega_2$. Since M = 0, we know $m_k \equiv k$. From the restart technique (3.10), the following sufficient decrease property holds

$$E(X^{k}) - E(X^{k+1}) \ge \sigma \|X^{k} - X^{k+1}\|^{2}, \forall k.$$
(B.2)

Let $S(X^0)$ be the set of limiting points of the sequence $\{X^k\}$ starting from X^0 . By the boundedness of $\{X^k\}$ and $S(X^k) = \bigcap_{n \in \mathbb{N}} \bigcup_{k \ge n} \{X^k\}$, it follows that $\{X^k\}$ is a non-empty and compact set. From Lemma 3.5, we know

$$E(X) = E^*, \quad \forall X \in S(X^0). \tag{B.3}$$

Then, for all ε_1 , $\eta > 0$, there exists $k_0 > 0$ such that

dist
$$(X^k, S(X^0)) \le \varepsilon_1, \quad E^* < E(X^k) < E^* + \eta, \quad \forall k > k_0.$$
 (B.4)

Applying Lemma B.1, for $\forall k > k_0$ we have

$$\psi'(E(X^k) - E^*) \operatorname{dist}(\mathbf{0}, \partial E(X^k)) \ge 1.$$
(B.5)

By the convexity of ψ , we have

$$\psi(E(X^k) - E^*) - \psi(E(X^{k+1}) - E^*) \ge \psi'(E(X^k) - E^*)(E(X^k) - E(X^{k+1})).$$
(B.6)

Denote $\Delta_{k,k+1} = \psi(E(X^k) - E^*) - \psi(E(X^{k+1}) - E^*)$. From (B.2), we have

$$\Delta_{k,k+1} \ge \frac{\sigma \|X^{k+1} - X^k\|^2}{\operatorname{dist}(\mathbf{0}, \partial E(X^k))}, \quad \forall k > k_0.$$
(B.7)

Next, we prove $\sum_{k=0}^{\infty} ||X^{k+1} - X^k|| < \infty$. For any $\varepsilon > 0$, we take $N > \max\{k_0, 3T\}$ large enough such that

$$\bar{C}\psi(E(X^{N+1}) - E^*) < \varepsilon/2, \tag{B.8}$$

$$\sum_{k=N-3T+2}^{N} \|X^{k+1} - X^{k}\| < \varepsilon/2, \tag{B.9}$$

where $\bar{C} = 3TC/\sigma$. The existence of *N* is ensured by the fact that $\lim_{k\to\infty} E(X^k) = E^*$, $\psi(0) = 0$ and $\lim_{k\to\infty} ||X^{k+1} - X^k|| = 0$. We now show that for all $\forall K > N$, the following inequality holds

$$\sum_{k=N+1}^{K} \|X^{k+1} - X^{k}\| < \varepsilon, \tag{B.10}$$

which implies $\sum_{k=0}^{\infty} ||X^{k+1} - X^k|| < \infty$ by Cauchy principle of convergence. In fact, together with Lemma 3.6 and (B.7), we obtain

$$\Delta_{k,k+1} \ge \frac{\|X^{k+1} - X^k\|^2}{C_1 \sum_{l=k-3T+1}^k \|X^l - X^{l-1}\|}, \quad \forall k > N,$$
(B.11)

where $C_1 = C/\sigma$. By using the geometric inequality, we get

$$2\|X^{k+1} - X^k\| \le \frac{1}{3T} \sum_{l=k-3T+1}^k \|X^l - X^{l-1}\| + \bar{C}\Delta_{k,k+1}.$$
 (B.12)

Summing up (B.12) for $k = N + 1, \dots, K$, it follows that

$$2\sum_{k=N+1}^{K} \|X^{k+1} - X^{k}\| \leq \sum_{k=N+1}^{K} \left(\frac{1}{3T} \sum_{l=k-3T+1}^{k} \|X^{l} - X^{l-1}\| + \bar{C}\Delta_{k,k+1}\right)$$
$$= \frac{1}{3T} \sum_{k=N+1}^{K} \sum_{l=-3T+1}^{0} \|X^{l+k} - X^{l+k-1}\| + \bar{C}\Delta_{N+1,K+1}$$
$$= \frac{1}{3T} \sum_{l=-3T+1}^{0} \sum_{k=N+l+1}^{K+l} \|X^{k} - X^{k-1}\| + \bar{C}\Delta_{N+1,K+1}$$

$$\leq \frac{1}{3T} \sum_{l=-3T+1}^{0} \sum_{k=N-3T+2}^{K} \|X^{k} - X^{k-1}\| + \bar{C}\Delta_{N+1,K+1}$$

$$= \sum_{k=N-3T+2}^{K} \|X^{k} - X^{k-1}\| + \bar{C}\Delta_{N+1,K+1}$$

$$\leq \sum_{k=N-3T+1}^{K} \|X^{k+1} - X^{k}\| + \bar{C}\Delta_{N+1,K+1},$$

where the first equality is from the fact that $\Delta_{p,q} + \Delta_{q,r} = \Delta_{p,r}$. Therefore, we have

$$\sum_{k=N+1}^{K} \|X^{k+1} - X^{k}\| \le \sum_{k=N-3T+1}^{N} \|X^{k+1} - X^{k}\| + \bar{C}\psi(E(X^{N+1}) - E^{*}) < \varepsilon.$$
(B.13)

As a result, we obtain $\sum_{k=0}^{\infty} ||X^{k+1} - X^k|| < \infty$ and $\lim_{k\to\infty} X^k = X^*$. From Lemma 3.5 and the continuity of E(X) on $\mathcal{B}(X^0)$, we get

$$\lim_{k \to \infty} E(X^k) = E(X^*) = E^*.$$
(B.14)

This completes the proof.

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168

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