

## A CLASS OF SELF-INTERACTING PROCESSES WITH APPLICATIONS TO GAMES AND REINFORCED RANDOM WALKS\*

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**Abstract.** This paper studies a class of non-Markovian and nonhomogeneous stochastic processes on a finite state space. Relying on a recent paper by Benaïm, Hofbauer, and Sorin [*SIAM J. Control Optim.*, 44 (2005), pp. 328–348] it is shown that, under certain assumptions, the asymptotic behavior of occupation measures can be described in terms of a certain set-valued deterministic dynamical system. This provides a unified approach to *simulated annealing* type processes and permits the study of new models of *vertex reinforced random walks* and new models of learning in games such as *Markovian fictitious play*.

**Key words.** stochastic approximation, processes with reinforcement, differential inclusions, learning in games, simulated annealing

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**1. Introduction.** Let  $E$  be a finite set called *the state space*,  $\mathbb{M} = \mathbb{M}(E)$  the set of Markov matrices over  $E$ , and  $\Sigma$  a compact convex subset of a Euclidean space called *the observation space*. The set  $\Sigma$  will be equipped with the distance induced by the Euclidean norm  $\|\cdot\|$  on the observation space. Let  $(\Omega, \mathcal{F}, \mathbb{P})$  be a probability space equipped with an increasing sequence of sub  $\sigma$ -fields  $\{\mathcal{F}_n, n \in \mathbb{N}\} : \mathcal{F}_n \subset \mathcal{F}_{n+1} \subset \mathcal{F}$ .

Our main object of interest is a discrete time random process  $(X, M, V) = ((X_n, M_n, V_n))$  defined on  $(\Omega, \mathcal{F}, \mathbb{P})$  taking values in  $E \times \mathbb{M}(E) \times \Sigma$  such that

- (i)  $(X, M, V)$  is *adapted* (to  $\{\mathcal{F}_n, n \in \mathbb{N}\}$ ), meaning that  $(X_n, M_n, V_n)$  is  $\mathcal{F}_n$ -measurable for each  $n$ ,
- (ii) for all  $y \in E$ ,

$$(1) \quad \mathbb{P}(X_{n+1} = y | \mathcal{F}_n) = M_n(X_n, y).$$

We refer to  $X_n$  (respectively,  $V_n$ ) as the *state* (respectively, the *observation*) variable at time  $n$  and to the sequence  $(M_n)$  as the *strategy*. We let

$$v_n = \frac{1}{n} \sum_{i=1}^n V_i$$

denote the empirical average up to time  $n$  of the sequence of observations.

A well-studied situation is when

$$(2) \quad M_n = K(v_n),$$

where  $K$  maps continuously probability vectors to irreducible Markov matrices and

$$V_{n+1} = H(X_{n+1}, v_n)$$

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for some map  $H : E \times \Sigma \rightarrow \Sigma$ . In such a case,  $(X_n)$  is called a “Markov chain controlled” by  $(v_n)$ , and the behavior of  $(v_n)$  can be analyzed through the ODE

$$(3) \quad \dot{v} = -v + \sum_x \pi(v)(x)H(x, v),$$

where  $\pi(v)$  is the invariant probability of  $K(v)$ . This approach to controlled Markov chains goes back to the work of Métivier and Priouret [20] (see also the books Benveniste, Métivier, and Priouret [10] and Duflo [12]), strongly influenced by the pioneered works of Ljung [19] and Kushner and Clarck [18] on the ODE’s method. It has been used in Benaïm [1] for analyzing certain vertex reinforced random walks on finite graphs.

The main purpose of this paper is to investigate the long-term behavior of  $(v_n)$  under less stringent assumptions than (2). In particular we are interested in situations where

- (a)  $M_n$  may depend on other (nonobservable or hidden) variables than  $v_n$  and
- (b) the closure of  $\{M_n : n \geq 0\}$  may contain degenerate (i.e., nonirreducible) Markov matrices.

Situation (a) typically occurs in game theory, where players may have only partial information on the actions played by their opponents, and (b) is motivated by stochastic optimization algorithms.

Relying on a recent paper by Benaïm, Hofbauer, and Sorin [5], it will be shown that under certain assumptions (involving estimates on the log-Sobolev and spectral gap constants of  $(M_n)$ ), the asymptotic behavior can be described in terms of a certain set-valued deterministic dynamical system that generalizes the ODE (3). Applications to nonhomogeneous Markov chains, vertex reinforced random walks, and learning processes in game theory will be given.

**1.1. Outline of contents.** The organization of the paper is as follows. Section 2 states the notation, the hypotheses, and the main result. Our main assumption (Hypothesis 2.1) is somewhat abstract, and more tractable conditions (expressed in terms of spectral gaps and log-Sobolev constants) are given in section 3. Section 4 is devoted to examples and applications. The proof of the main result is postponed to section 5.

**2. Notation, hypotheses, and main results.** A *probability vector* (or measure) over  $E$  is a map  $\mu : E \rightarrow \mathbb{R}^+$  such that  $\sum_x \mu(x) = 1$ . A *Markov matrix* is a map  $M : E \times E \rightarrow \mathbb{R}^+$  such that

$$\text{for all } x \in E, \sum_y M(x, y) = 1.$$

We let  $\Delta = \Delta(E)$  denote the space of probability vectors over  $E$  and  $\mathbf{M} = \mathbf{M}(E)$  denote the set of Markov matrices on  $E$ .

Given a function  $f : E \rightarrow \mathbb{R}$  and  $\mu \in \Delta$ , we use the notation

$$\mu f = \sum_x \mu(x)f(x).$$

A Markov matrix  $M$  on  $E$  acts on functions  $f$  and measures  $\mu$  according to the formulas

$$Mf(x) = \sum_y M(x, y)f(y),$$

$$\mu M(y) = \sum_x \mu(x)M(x, y).$$

We let  $M^n$  denote the Markov matrix obtained by matrix multiplication. Equivalently  $M^n f = M(M^{n-1} f)$  for  $n \geq 1$  with the convention that  $M^0 f = f$ .

Points  $x, y \in E$  are said to be *related* if there exist  $i, j \geq 0$  (depending on  $x$  and  $y$ ) such that  $M^i(x, y) > 0$  and  $M^j(y, x) > 0$ . An equivalence class for this relation is called a *recurrent class*. The Markov matrix  $M$  on  $E$  is said to be *indecomposable* if it has a unique recurrent class (possibly periodic), and it is said to be *irreducible* if this recurrent class is  $E$ .

By standard results, indecomposability of  $M$  implies that  $M$  possesses a unique *invariant probability measure*  $\pi$  characterized by the relation  $\pi M = \pi$ . Moreover, the *generator*  $L = -I + M$  has kernel  $\mathbb{R}1$ , and its restriction to  $\{f : \pi f = 0\}$  is an isomorphism. It then follows that  $-L$  admits a pseudo “inverse”  $Q$  characterized by

$$Q\mathbf{1} = 0$$

and

$$Q(I - M) = (I - M)Q = I - \Pi,$$

where  $\Pi \in \mathbb{M}$  denote the matrix defined by  $\Pi(x, y) = \pi(y)$ . To shorten notation, we also call  $Q$  the pseudoinverse of  $M$ .

Given a vector  $f$  and a matrix  $N$ , we set  $|f| = \max |f(x)|$  and  $|N| = \max_{x,y} |N(x, y)|$ .

Our main assumption is the following hypothesis.

*Hypothesis 2.1. The matrices  $(M_n)$  are indecomposable, and their pseudoinverses  $(Q_n)$  and invariant probabilities  $(\pi_n)$  satisfy almost surely*

(i)

$$\lim_{n \rightarrow \infty} \frac{|Q_n|^2 \log(n)}{n} = 0,$$

(ii)

$$\lim_{n \rightarrow \infty} |Q_{n+1} - Q_n| = 0,$$

(iii)

$$\lim_{n \rightarrow \infty} |\pi_{n+1} - \pi_n| = 0.$$

The verification of Hypothesis 2.1 is the subject of section 3, where sufficient and more tractable conditions will be detailed.

Let  $\hat{V}_n : E \rightarrow \Sigma$  be an  $\mathcal{F}_n$ -measurable map defined by

$$\hat{V}_n(x) = \frac{\mathbb{E}(V_{n+1} \mathbf{1}_{X_{n+1}=x} | \mathcal{F}_n)}{M_n(X_n, x)}$$

for  $M_n(X_n, x) \neq 0$ . In addition to Hypothesis 2.1, we assume the following hypothesis.

*Hypothesis 2.2.*

$$\lim_{n \rightarrow \infty} M_{n+1} Q_{n+1} (\hat{V}_{n+1} - \hat{V}_n) = 0$$

*almost surely.*

*Remark 2.3.* We give here some sufficient conditions ensuring Hypothesis 2.2. These conditions are usually much easier to verify.

- (i) Assume that  $x \mapsto \hat{V}_{n+1}(x) - \hat{V}_n(x)$  is a constant map (possibly depending on  $n$ ). Then Hypothesis 2.2 holds since  $Q_n 1 = 0$ . This sufficient condition will be used in section 4.
- (ii) More generally, let  $T\Sigma$  be the affine hull of  $\Sigma$  (the smallest affine space containing  $\Sigma$ ). Assume that for all  $n \in \mathbb{N}$ , there exists a vector  $A_n \in T\Sigma$  and a map  $B_n : E \rightarrow T\Sigma$  such that
  - (a) for all  $x \in E$ ,  $\hat{V}_{n+1}(x) - \hat{V}_n(x) = A_n + B_n(x)$ ,
  - (b)  $\limsup_{n \rightarrow \infty} |B_n| \sqrt{\frac{n}{\log(n)}} < \infty$  almost surely.
 Then  $|M_{n+1}Q_{n+1}((\hat{V}_{n+1} - \hat{V}_n))| = |M_{n+1}Q_{n+1}B_n| \leq |Q_{n+1}||B_n| \rightarrow 0$  almost surely by Hypothesis 2.1.
- (iii) Assume that  $M_n(x, y) = \pi_n(y)$ . Then  $M_{n+1}Q_{n+1} = 0$  so that Hypothesis 2.2 holds.

**2.1. Adapted set-valued dynamical systems.** The purpose of this section is to introduce certain differential inclusions on  $\Sigma$  that will prove to be useful for analyzing the long-term behavior of  $(v_n)$ . Recall that we let  $\pi_n$  denote the invariant probability of  $M_n$ . Let

$$(4) \quad \theta_n = \pi_n \hat{V}_n = \sum_x \pi_n(x) \hat{V}_n(x).$$

We let  $C_n \subset \Sigma \times \Sigma$  denote the *topological support* of the law of  $(v_n, \theta_n)$ . That is the smallest closed set  $F \subset \Sigma \times \Sigma$  such that

$$\mathbb{P}((v_n, \theta_n) \in F) = 1.$$

Let  $\text{clos}\{C_n\}$  denote the set of all possible limit points  $z = \lim z_{n_k}$  with  $z_{n_k} \in C_{n_k}$  and  $n_k \rightarrow \infty$ . It is easily seen that  $\text{clos}\{C_n\}$  is a nonempty compact subset of  $\Sigma \times \Sigma$ .

A nonempty set  $G \subset \Sigma \times \Sigma$  is called a *graph* (or a bundle) over  $\Sigma$  if the projection

$$\begin{aligned} p : G &\rightarrow \Sigma, \\ (u, v) &\mapsto u \end{aligned}$$

is onto. A graph  $G$  over  $\Sigma$  defines a *set-valued function* mapping each point  $u \in \Sigma$  to a set  $G(u) = \{v \in \Sigma : (u, v) \in G\}$ .

DEFINITION 2.4. A set  $C \subset \Sigma \times \Sigma$  is said to be adapted to  $\{(v_n, \theta_n)\}$  (or simply adapted) if

- (i)  $C$  is a closed graph over  $\Sigma$ ,
- (ii) for all  $u \in \Sigma$ ,  $C(u)$  is a nonempty convex set,
- (iii)  $\text{clos}\{C_n\} \subset C$ .

To an adapted set  $C$ , we associate the differential inclusion

$$(5) \quad \dot{v} \in -v + C(v).$$

A solution to (5) is an absolutely continuous mapping  $v : \mathbb{R} \rightarrow \Sigma$  verifying  $\dot{v}(t) + v(t) \in C(v(t))$  for almost every  $t$ . A set  $A \subset \Sigma$  is said to be *invariant* if, for all  $x \in A$ , there exists a solution  $\mathbf{x}$  to (5) with  $\mathbf{x}(0) = x$  and such that  $\mathbf{x}(\mathbb{R}) \subset A$ .

Given a set  $A \subset \Sigma$  and  $(x, y) \in A^2$ , we write  $x \hookrightarrow_A y$  if, for every  $\varepsilon > 0$  and  $T > 0$ , there exists an integer  $n \in \mathbb{N}$ , solutions  $\mathbf{x}_1, \dots, \mathbf{x}_n$  to (5), and real numbers  $t_1, t_2, \dots, t_n$  greater than  $T$  such that

- (a)  $\mathbf{x}_i([0, t_i]) \subset A$ ,
- (b)  $\|\mathbf{x}_i(t_i) - \mathbf{x}_{i+1}(0)\| \leq \varepsilon$  for all  $i = 1, \dots, n-1$ ,
- (c)  $\|\mathbf{x}_1(0) - x\| \leq \varepsilon$  and  $\|\mathbf{x}_n(t_n) - y\| \leq \varepsilon$ .

DEFINITION 2.5. A set  $A \subset \Sigma$  is said to be internally chain transitive provided  $A$  is compact and  $x \rightsquigarrow_A y$  for all  $x, y \in A$ .

It is not hard to verify (see, e.g., Benaïm, Hofbauer, and Sorin [5, Lemma 3.5]) that an internally chain transitive set is invariant.

The limit set of  $(v_n)$  is the set  $L = L((v_n))$  consisting of all points  $p = \lim v_{n_k}$  for some sequence  $n_k \rightarrow \infty$ . The next Theorem 2.6 is the main result of the paper. Its proof heavily relies on Benaïm, Hofbauer, and Sorin [5] and is given in section 5.

THEOREM 2.6. Assume that Hypotheses 2.1 and 2.2 hold. Let  $C$  be an adapted graph. Then the limit set of  $(v_n)$  is an internally chain transitive set for the differential inclusion

$$\dot{v} \in -v + C(v).$$

**2.2. Background: How to use Theorem 2.6.** The notion of “internally chain transitive set” was introduced by Benaïm and Hirsch [3] in order to analyze the long-term behavior of certain perturbations of flows and has been recently extended to multivalued dynamical systems by Benaïm, Hofbauer, and Sorin [5]. We refer the reader to this paper for more details, examples, and properties. For convenience this section briefly reviews a few useful properties of internally chain transitive sets.

The differential inclusion (5) induces a set-valued dynamical system  $\{\Phi_t\}_{t \in \mathbb{R}}$  defined by

$$\Phi_t(x) = \{\mathbf{x}(t) : \mathbf{x} \text{ is a solution to (5) with } \mathbf{x}(0) = x \in \Sigma\}.$$

A nonempty compact set  $A$  is an attracting set if there exists a neighborhood  $U$  of  $A$  and a function  $\mathbf{t}$  from  $(0, \varepsilon_0)$  to  $\mathbb{R}^+$  with  $\varepsilon_0 > 0$  such that

$$\Phi_t(U) \subset A^\varepsilon$$

for all  $\varepsilon < \varepsilon_0$  and  $t \geq \mathbf{t}(\varepsilon)$ , where  $A^\varepsilon$  stands for the  $\varepsilon$ -neighborhood of  $A$ . If additionally  $A$  is invariant, then  $A$  is an attractor.

Given an attracting set (respectively, an attractor)  $A$ , its basin of attraction is the set

$$B(A) = \{x \in \Sigma : \exists t \geq 0, \Phi_t(x) \in U\}.$$

When  $B(A) = \Sigma$ ,  $A$  is a globally attracting set (respectively, a global attractor).

Given a closed invariant set  $S$ , the induced dynamical system  $\Phi^S$  on  $S$  is defined by

$$\Phi_t^S(x) = \{\mathbf{x}(t) : \mathbf{x} \text{ is a solution to (5) with } \mathbf{x}(0) = x \text{ and } \mathbf{x}(\mathbb{R}) \subset S\}.$$

An invariant set  $S$  is attractor free if there exists no proper subset  $A$  of  $S$  which is an attractor for  $\Phi^S$ .

Throughout the remainder of this section, we let  $L$  denote an internally chain transitive set (for instance, the limit set  $L = L(v_n)$ ). Properties of  $L$  will then be obtained through the next result (see Benaïm, Hofbauer, and Sorin [5, Lemma 3.5, Proposition 3.20 and Theorem 3.23]).

PROPOSITION 2.7.

- (i) The set  $L$  is nonempty, compact, invariant, and attractor free.
- (ii) If  $A$  is an attracting set with  $B(A) \cap L \neq \emptyset$ , then  $L \subset A$ .

Some useful properties of attracting sets or attractors are given in the two following propositions (see Benaïm, Hofbauer, and Sorin [5, Propositions 3.25 and 3.27]).

**PROPOSITION 2.8.** *Let  $\Lambda \subset \Sigma$  be compact with a bounded open neighborhood  $U$  and  $V : \bar{U} \rightarrow [0, \infty[$ . Assume the following conditions:*

- (i)  $\Phi_t(U) \subset U$  for all  $t \geq 0$ ,
- (ii)  $V^{-1}(0) = \Lambda$ ,
- (iii)  $V$  is continuous, and for all positive  $t$ ,  $x \in U \setminus \Lambda$  and  $y \in \Phi_t(x)$ , one has  $V(y) < V(x)$ .

*Then  $\Lambda$  contains an attractor whose basin contains  $U$ .*

The map  $V$  introduced in this proposition is called a *strong Lyapunov function* associated to  $\Lambda$ .

Let now  $\Lambda$  be a subset of  $\Sigma$  and  $U \subset \Sigma$  an open neighborhood of  $\Lambda$ . A continuous function  $V : U \rightarrow \mathbb{R}$  is called a *Lyapunov function* for  $\Lambda \subset \Sigma$  if  $V(y) < V(x)$  for all  $x \in U \setminus \Lambda$ ,  $y \in \Phi_t(x)$ ,  $t > 0$ , and  $V(y) \leq V(x)$  for all  $x \in \Lambda$ ,  $y \in \Phi_t(x)$ , and  $t \geq 0$ .

**PROPOSITION 2.9 (Lyapunov).** *Suppose  $V : U \rightarrow \mathbb{R}$  is a Lyapunov function for  $\Lambda$  and  $L \subset U$ . Assume that  $V(\Lambda)$  has an empty interior. Then  $L \subset \Lambda$ , and the restriction of  $V$  to  $L$  is constant.*

**3. Verification of Hypothesis 2.1.** This section is devoted to the verification of Hypothesis 2.1. The results given here will be used in section 4 to analyze specific situations.

**3.1. Estimates based on compactness.** Let  $M_{ind}(E)$  denote the open set of indecomposable Markov matrices.

**PROPOSITION 3.1.** *Suppose that the sequence  $(M_n)$  lies in a compact subset of  $M_{ind}(E)$  and verifies  $\lim_{n \rightarrow \infty} (M_{n+1} - M_n) = 0$ . Then Hypothesis 2.1 holds.*

This proposition is a direct consequence of the next lemma.

**LEMMA 3.2.** *Let  $TM(E)$  be the space of matrices  $K = K(x, y)$  such that  $\sum_y K(x, y) = 0$ . The map  $Q : M_{ind}(E) \rightarrow TM(E)$ , which associates to  $M$  its pseudoinverse, and the map  $\Pi : M_{ind}(E) \rightarrow \Delta$ , which associates to  $M$  its invariant measure, are smooth maps.*

*Proof.* Let  $\phi : M_{ind}(E) \times \Delta \rightarrow T\Delta$  be the smooth map defined by

$$\phi(M, \mu) = \mu(I - M)$$

with  $T\Delta = \{\mu : E \rightarrow \mathbb{R} : \sum_x \mu(x) = 0\}$ . Set  $M \in M_{ind}(E)$ . Then  $\Pi(M)$ , the invariant probability of  $M$ , is a solution to  $\phi(M, \mu) = 0$ . For all  $\nu \in T\Delta$ ,

$$\frac{\partial \phi}{\partial \mu}(M, \mu) \cdot \nu = \nu(I - M).$$

Hence, by uniqueness of the invariant probability measure,  $\frac{\partial \phi}{\partial \mu}(M, \mu)$  has kernel  $\{0\}$ . The fact that  $\Pi$  is smooth follows from the implicit function theorem.

We denote by  $\hat{\Pi}(M) \in M(E)$  the matrix defined by  $\hat{\Pi}(M)(x, y) = \Pi(M)(y)$ . The pseudoinverse of  $M$  is a solution to  $\psi(M, Q) = 0$ , where  $\psi : M_{ind}(E) \times TM(E) \rightarrow TM(E)$  is the smooth map defined by

$$\psi(M, Q) = Q(I - M) - (I - \hat{\Pi}(M)).$$

For all  $A \in TM(E)$ ,

$$\frac{\partial \psi}{\partial Q}(M, Q) \cdot A = A(I - M).$$

Hence, by uniqueness of the invariant probability measure,  $\frac{\partial \psi}{\partial Q}(M, Q)$  has kernel  $\{0\}$ . The fact that  $Q$  depends smoothly on  $M$  follows from the implicit function theorem.  $\square$

Let  $K$  be a continuous mapping from a compact set  $\Gamma$  into  $M(E)$  such that  $K(w)$  is indecomposable for all  $w \in \Gamma$ . Assume  $(w_n)$  is a sequence of  $\Gamma$ -valued random variables such that  $M_n = K(w_n)$ . If, in addition,  $\lim_{n \rightarrow \infty} (M_{n+1} - M_n) = 0$ , then Proposition 3.1 applies.

**3.2. Estimates based on log-Sobolev and spectral gap constants.** Propositions 3.3 and 3.4 below can be used to verify Hypothesis 2.1 when the sequence  $(M_n)$  is not bounded away from  $M_{ind}(E)$ . The strategy is then to verify assertions (i), (ii), and (iii) of Proposition 3.3. We will use the estimates given by Proposition 3.4 to verify assertion (i).

PROPOSITION 3.3. *Suppose that the matrices  $(M_n)$  are indecomposable and that their pseudoinverses  $(Q_n)$  and invariant probabilities  $(\pi_n)$  satisfy almost surely*

(i)

$$\lim_{n \rightarrow \infty} \frac{|Q_n|^2 \log(n)}{n} = 0,$$

(ii)

$$\limsup_{n \rightarrow \infty} |M_{n+1} - M_n| \frac{n}{\log(n)} < \infty,$$

(iii)

$$\limsup_{n \rightarrow \infty} |\pi_{n+1} - \pi_n| \sqrt{\frac{n}{\log(n)}} < \infty.$$

Then Hypothesis 2.1 holds.

*Proof.* The proof amounts to show that Hypothesis 2.1 (ii) holds. Set  $L_n = M_n - I$  and  $\Pi_n = \hat{\Pi}(M_n)$ . Using the characterization of  $Q_n$ , one has

$$Q_{n+1}(L_{n+1} - L_n) + (Q_{n+1} - Q_n)L_n = \Pi_{n+1} - \Pi_n.$$

Hence,

$$Q_{n+1}(L_{n+1} - L_n)Q_n + (Q_{n+1} - Q_n)L_nQ_n = (\Pi_{n+1} - \Pi_n)Q_n.$$

This implies (using  $Q_n\Pi_n = Q_n\Pi_{n+1} = 0$  and  $L_nQ_n = \Pi_n - I$ )

$$Q_{n+1}(M_{n+1} - M_n)Q_n + (Q_n - Q_{n+1}) = (\Pi_{n+1} - \Pi_n)Q_n.$$

Therefore,

$$|Q_n - Q_{n+1}| \leq c(|Q_{n+1}||Q_n||M_{n+1} - M_n| + |\pi_{n+1} - \pi_n||Q_n|)$$

for some constant  $c > 0$ , and conditions (i), (ii), and (iii) imply Hypothesis 2.1 (ii).  $\square$

Let  $M_{irr}(E)$  denote the open set of irreducible Markov matrices. Let  $M \in M_{irr}(E)$  with invariant probability  $\pi$ , and let  $f : E \rightarrow \mathbb{R}$ . The *variance*, *entropy*, and *energy* of  $f$  are defined, respectively, as

$$\begin{aligned} \text{var}(f) &= \pi(f^2) - (\pi f)^2, \\ \mathcal{L}(f) &= \sum_x f(x)^2 \log \left( \frac{f(x)^2}{\pi f^2} \right) \pi(x), \\ \mathcal{E}(f) &= \frac{1}{2} \sum_{x,y} (f(y) - f(x))^2 M(x,y) \pi(x). \end{aligned}$$

The *spectral gap* and *log-Sobolev* constants of  $M$  are then defined to be

$$\lambda = \min \left\{ \frac{\mathcal{E}(f)}{\text{var}(f)} : \text{var}(f) \neq 0 \right\},$$

$$\alpha = \min \left\{ \frac{\mathcal{E}(f)}{\mathcal{L}(f)} : \mathcal{L}(f) \neq 0 \right\}.$$

The following estimates follow from the quantitative results for finite Markov chains as given in Saloff-Coste [28].

PROPOSITION 3.4. *Let  $M \in \mathbf{M}_{\text{irr}}(E)$  with invariant probability  $\pi$ , log-Sobolev constant  $\alpha$ , and spectral gap  $\lambda$ . For all  $(x, y) \in E$ , the following estimates hold:*

(i)

$$|Q(x, y)| \leq \sqrt{\frac{\pi(y)}{\pi(x)} \frac{1}{\lambda}},$$

(ii)

$$|Q(x, y)| \leq \frac{1}{\alpha} \log_+ \left( \log \left( \frac{1}{\pi(x)} \right) \right) + \frac{e}{\lambda},$$

where  $\log_+(t) = \max(0, \log(t))$ .

In particular, denoting  $\pi_* = \min_x \pi(x)$ ,

$$|Q| \leq \frac{1}{\alpha} \left[ \log_+ \left( \log \left( \frac{1}{\pi_*} \right) \right) + \frac{e}{2} \right]$$

and

$$|Q| \leq \frac{1}{\lambda} \left[ \log_+ \left( \log \left( \frac{1}{\pi_*} \right) \right) \left( \frac{\log \left( \frac{1 - \pi_*}{\pi_*} \right)}{1 - 2\pi_*} \right) + e \right].$$

*Proof.* Let  $L = -I + M$ , and let  $\{P_t\}$  be the continuous time semigroup  $P_t = e^{tL}$ . Then  $Q$  can be written as

$$Q(x, y) = \int_0^\infty (P_t(x, y) - \pi(y)) dt.$$

The first assertion then easily follows from the estimate

$$|P_t(x, y) - \pi(y)| \leq \sqrt{\frac{\pi(y)}{\pi(x)}} e^{-\lambda t}$$

whose proof can be found in Saloff-Coste (see [28, Corollary 2.1.5]).

We now pass to the second assertion. If  $\pi(x) \geq e^{-2}$ , the inequality to be proved follows from inequality (i). Hence we assume that  $\pi(x) < e^{-2}$ , and we follow the line of the proof of Theorem 2.2.5 in Saloff-Coste [28]. For  $q \geq 1$ , we let  $\|\cdot\|_q$  denote the norm in  $l^q(\pi)$ . We let  $P_t^*$  denote the adjoint of  $P_t$  in  $l^2(\pi)$  and  $p_t(x, y) = p_t^*(y, x) = P_t(x, y)/\pi(y)$ . Let  $g_x$  denote the function given by  $g_x(y) = 0$  for  $x \neq y$  and  $g_x(x) = 1/\pi(x)$ . Then

$$|P_t(x, y) - \pi(y)| \leq \|p_t(x, \cdot) - 1\|_2 = \|(P_t^* - \pi)g_x\|_2.$$

Therefore,

$$\begin{aligned} |P_{t+s}(x, y) - \pi(y)| &\leq \|p_{t+s}(x, \cdot) - 1\|_2 \leq \|P_t^* - \pi\|_{2 \rightarrow 2} \|P_s^* g_x\|_2 \\ &\leq e^{-\lambda t} \|P_s^*\|_{k \rightarrow 2} \|g_x\|_k \end{aligned}$$

for any  $k \geq 1$ , where we have used the fact that  $\|P_t^* - \pi\|_{2 \rightarrow 2} \leq e^{-\lambda t}$ . Let  $q$  be the Hölder conjugate of  $k$ . Then  $\|P_s^*\|_{k \rightarrow 2} = \|P_s\|_{2 \rightarrow q}$ . Now choose  $q(s) = 1 + e^{2\alpha s}$ . By hypercontractivity (see Theorem 2.2.4 in Saloff-Coste [28]),  $\|P_s\|_{2 \rightarrow q(s)} \leq 1$  so that

$$|P_{t+s}(x, y) - \pi(y)| \leq e^{-\lambda t} \pi(x)^{-1/q(s)}.$$

Hence,

$$|Q(x, y)| \leq 2s + \frac{1}{\lambda} \pi(x)^{-1/q(s)}.$$

For  $s = \frac{1}{2\alpha} \log_+(\log(\frac{1}{\pi(x)}))$ , this gives the desired result.

The uniform bounds on  $|Q|$  follow from the rough estimates

$$\frac{1 - 2\pi_*}{\log\left(\frac{(1 - \pi_*)}{\pi_*}\right)} \lambda \leq \alpha \leq \frac{\lambda}{2}$$

given in Saloff-Coste (see [28, Lemma 2.2.2 and Corollary 2.2.10]). □

**4. Some applications.** In sections 4.1 and 4.2, we are interested in the long-term behavior of the *empirical occupation measure of the process*. We then let  $\Sigma = \Delta$ ,  $V_n = \delta_{X_n}$  and

$$v_n = \frac{1}{n} \sum_{i=1}^n \delta_{X_i}.$$

Hence,  $\hat{V}_n(x) = \delta_x$  and  $\theta_n = \pi_n$ . Note that, according to Remark 2.3 (i), Hypothesis 2.2 holds.

**4.1. Markov chains.** Let  $(M_n)$  be a deterministic (or  $\mathcal{F}_0$ -measurable) sequence of Markov matrices over  $E$ . A *nonhomogeneous* Markov chain with transition matrices  $(M_n)$  is an adapted process  $(X_n)$  on  $E$  verifying (1).

**PROPOSITION 4.1.** *Let  $L((\pi_n)) \subset \Delta$  denote the limit set of  $(\pi_n)$ , and let  $\text{conv}[L((\pi_n))]$  denote its convex hull. Suppose that Hypothesis 2.1 holds. Then  $L((v_n)) \subset \text{conv}[L((\pi_n))]$  with probability one.*

*Proof.* The set  $C = \Delta \times \text{conv}[L((\pi_n))]$  is adapted to  $(v_n, \pi_n)$ . The induced differential equation  $\dot{v} \in -v + \text{conv}[L((\pi_n))]$  has a unique global attractor  $\text{conv}[L((\pi_n))]$ . Hence, by Theorem 2.6 and Proposition 2.7 (ii),  $L((v_n)) \subset \text{conv}[L((\pi_n))]$ . □

**COROLLARY 4.2.** *Suppose that the sequence  $(M_n)$  lies in a compact subset of  $M_{\text{ind}}(E)$  and verifies  $M_{n+1} - M_n \rightarrow 0$ . Then the conclusion of Proposition 4.1 holds.*

*Proof.* The proof follows from Propositions 4.1 and 3.1. □

**COROLLARY 4.3.** *Assume that  $M_n \rightarrow M \in M_{\text{ind}}(E)$ . Then  $v_n \rightarrow \pi$  is the invariant probability of  $M$ .*

**4.1.1. Markov chains with rare transitions.** Among the well-studied chains that motivate our analysis are the *chains with rare transitions*.

Let  $M_0$  be an irreducible Markov matrix over  $E$ , reversible with respect to a reference probability  $\pi_0$ . That is,

$$\pi_0(x)M_0(x, y) = \pi_0(y)M_0(y, x).$$

We sometimes call such an  $M_0$  an *exploration matrix* since it provides a way to explore the state space.

Let  $W : E \times E \rightarrow \mathbb{R}$  be a map and  $(\beta_n)$  a sequence of positive numbers. Set

$$(6) \quad M_n(x, y) = M(\beta_n, x, y),$$

where

$$M(\beta, x, y) = \begin{cases} M_0(x, y)\psi[\exp(-\beta W(x, y))] & \text{if } x \neq y, \\ 1 - \sum_{y \neq x} M(\beta, x, y) & \text{if } x = y, \end{cases}$$

and

$$(7) \quad \psi(u) = \min(1, u)$$

or

$$\psi(u) = \frac{u}{1+u}.$$

In particular, let  $U : E \rightarrow \mathbb{R}$  be a map, and let

$$(8) \quad W(x, y) = U(y) - U(x);$$

then  $(M_n)$  are the transition matrices of the so-called *Metropolis–Hastings* ( $\beta_n = \beta$ ) or *simulated annealing* ( $\beta_n \rightarrow \infty$ ) algorithm (see Hajek [15], Holley and Stroock [17], and Miclo [21]).

Consider the Markov chain with rare transitions (6), where  $W$  is given by (8). For  $x, y \in E$ , a path  $\gamma$  from  $x$  to  $y$  is a sequence of points  $x_0 = x, x_1, \dots, x_n = y$  such that  $M_0(x_i, x_{i+1}) > 0$ . We let  $\Gamma_{x,y}$  denote the set of all paths from  $x$  to  $y$ . The *elevation* from  $x$  to  $y$  is defined as

$$\text{Elev}(x, y) = \min\{\max\{U(z) : z \in \gamma\} : \gamma \in \Gamma_{x,y}\},$$

and the *energy barrier* is defined as

$$(9) \quad U^\# = \max\{\text{Elev}(x, y) - U(x) - U(y) + \min U : x \in E, y \in E\}.$$

**PROPOSITION 4.4.** *Consider the Markov chain with rare transitions (6) with  $W$  given by (8). Assume that  $\beta_n = \beta(n)$ , where  $\beta : \mathbb{R}^+ \rightarrow \mathbb{R}^+$  is differentiable, and verify*

$$0 \leq \dot{\beta}(t) \leq \frac{A}{t}$$

for some  $A < 1/2U^\#$ . Then  $v_n \rightarrow \pi$  where

$$\pi(x) \propto \pi_0(x)\mathbf{1}_{\text{Argmin}U}(x).$$

*Proof.* Our first goal is to verify Hypothesis 2.1. Let  $\lambda(\beta)$  denote the spectral gap of  $M(\beta, \cdot, \cdot)$ . It follows from Theorem 2.1 in Holley and Stroock [17] that

$$(10) \quad \lim_{\beta \rightarrow \infty} \frac{\log(\lambda(\beta))}{\beta} = -U^\#.$$

The invariant probability measure of  $M(\beta, \cdot, \cdot)$  is the *Gibbs measure*

$$(11) \quad \pi_\beta(x) \propto \exp(-\beta U(x))\pi_0(x).$$

Since  $\beta_n \leq \beta_1 + A \log(n)$ , by application of the last inequality of Proposition 3.4, one gets that Hypothesis 2.1 (i) holds.

For  $x \neq y$ ,

$$\frac{\partial M(\beta, x, y)}{\partial \beta} = -M_0(x, y)W(x, y)\psi'(\exp(-\beta W(x, y)) \exp(-\beta W(x, y))).$$

Using the fact that  $|\psi'(t)t| \leq 1$ , one gets that

$$\left| \frac{\partial M(\beta, x, y)}{\partial \beta} \right| \leq c$$

for some  $c > 0$ . Hence, by the mean value theorem,

$$|M_{n+1} - M_n| \leq c|\beta_{n+1} - \beta_n| \leq (Ac)/n.$$

This proves assertion (ii) of Proposition 3.3. The proof of assertion (iii) is similar since

$$\left| \frac{\partial \pi_\beta(x)}{\partial \beta} \right| = |\pi_\beta(x)(U(x) - \sum_y \pi_\beta(y)U(y))| \leq 2\|U\|.$$

This concludes the verification of Hypothesis 2.1.

Here  $\pi_n(x) \propto \exp(-\beta_n U(x))\pi_0(x)$  so that  $\pi_n \rightarrow \pi$ . The result follows from Proposition 4.1.  $\square$

*Remark 4.5.* For general  $W$ , it is always possible to define a *quasipotential*  $U$  (defined in terms of  $W$  and  $M_0$ ) and an energy barrier  $U^\#$  (in general not given by (9)) such that both (10) and (11) hold. We refer the reader to Miclo [21] for more details and proofs. With this quasipotential and barrier, Proposition 4.4 holds.

**4.2. Vertex reinforced random walks.** Vertex reinforced random walks (VRRW) were first introduced by Pemantle (see [23, 24]).

Suppose  $\mathcal{F}_n = \sigma(X_1, \dots, X_n)$ . A general VRRW on  $E$  is defined by

$$M_n(x, y) = K_n(x, y, v_n),$$

where for each integer  $n$  and  $v \in \Delta$ ,  $K_n(\cdot, \cdot, v)$  is a deterministic Markov matrix over  $E$ , which specifies the rule of the reinforcement.

The following result was proved in Benaim [1].

**PROPOSITION 4.6.** *Assume that there exists a  $[0, 1]$ -valued sequence  $\epsilon_n$  converging to 0 at infinity such that  $K_n(x, y, v) = K(x, y, \epsilon_n, v)$ , that the map  $(\epsilon, v) \mapsto K(\cdot, \cdot, \epsilon, v)$  is continuous on  $[0, 1] \times \Delta$ , and that  $K(\cdot, \cdot, \epsilon, v)$  is indecomposable for each  $(\epsilon, v) \in [0, 1] \times \Delta$ . Let  $\pi(v)$  denote the invariant measure of  $K(x, y, 0, v)$ . Then the limit set of  $(v_n)$  is almost surely an internally chain transitive set of the differential equation*

$$(12) \quad \dot{v} = -v + \pi(v).$$

*Proof.* This follows from Proposition 3.1 and Theorem 2.6.  $\square$

**4.2.1. Linear reinforcement.** The original VRRW as defined by Pemantle [23, 24] corresponds to a *linear reinforcement*

$$M_n(x, y) \propto U(x, y) \left[ 1 + \sum_{i=1}^n \mathbf{1}_{X_i=y} \right],$$

where  $U$  is a matrix with nonnegative entries.

We will here assume that  $U$  has positive entries. Then, for each  $n$ ,  $M_n$  is irreducible. With the notation of the previous paragraph,

$$(13) \quad M_n(x, y) = K(x, y, 1/n, v_n),$$

where, for  $(\epsilon, v) \in [0, 1] \times \Delta$ ,

$$(14) \quad K(x, y, \epsilon, v) \propto U(x, y) [\epsilon + v(y)].$$

The mapping  $(\epsilon, v) \mapsto K(\cdot, \cdot, \epsilon, v)$  is continuous on  $[0, 1] \times \Delta$ .

On a finite graph, this process was first analyzed by Pemantle [24] for symmetric positive matrices ( $U(x, y) = U(y, x) > 0$ ) and later by Benaïm [1] for general positive matrices using Proposition 4.6. An example of what can be proved is the following result first due to Pemantle [24].

**PROPOSITION 4.7.** *Suppose  $U(x, y) = U(y, x) > 0$ . Then the limit set of  $(v_n)$  is a compact connected subset of the critical set of the map*

$$v \mapsto U(v, v) = \sum_{x, y} U(x, y) v(x) v(y).$$

*Proof.* This follows from the fact that  $v \mapsto U(v, v)$  is a strict Lyapunov function of (12) whose critical points are the zeroes of (12).  $\square$

When the matrix  $U$  has zero entries,  $K(x, y, 0, v)$  may no longer be indecomposable for some  $v \in \partial\Delta$ , and Proposition 4.6 cannot be applied. This makes the analysis of VRRW with linear reinforcement much more difficult. Beautiful results on  $\mathbb{Z}$  and  $\mathbb{Z}^d$  have been obtained by Pemantle and Volkov [26], Volkov [30], and Tarrès [29]. We refer the reader to Pemantle [25] for a survey and further references.

**4.2.2. Nonhomogeneous linear reinforcement.** Let  $(a_n)$  be a positive sequence, and denote  $r_n = \sum_{i=1}^n a_i$ . We will assume that  $\lim_{n \rightarrow \infty} \frac{r_{n+1}}{r_n} = 1$ . Consider the VRRW corresponding to

$$M_n(x, y) \propto U(x, y) \left[ 1 + \sum_{i=1}^n a_i \mathbf{1}_{X_i=y} \right],$$

where  $U$  is a matrix with positive entries. Equivalently,  $M_n(x, y) = K(x, y, \epsilon_n, w_n)$  with

$$(15) \quad K(x, y, \epsilon, w) \propto U(x, y) [\epsilon + w(y)],$$

$\epsilon_n = 1/r_n$ , and  $w_n = \frac{1}{r_n} \sum_{i=1}^n a_i \delta_{X_i}$ . Using Proposition 3.1, it is not hard to check that Hypotheses 2.1 and 2.2 (with  $V_i = \delta_{X_i}$ ) are satisfied so that Theorem 2.6 applies.

Since  $\delta_{X_i} = v_i + (i - 1)(v_i - v_{i-1})$ , using the convention  $r_0 = v_0 = 0$ ,

$$\begin{aligned} w_n &= \frac{1}{r_n} \sum_{i=1}^n (r_i - r_{i-1})v_i + \frac{1}{r_n} \sum_{i=1}^n (i - 1)(v_i - v_{i-1})a_i \\ &= v_n + \frac{1}{r_n} \sum_{i=1}^{n-1} r_i(v_i - v_{i+1}) + \frac{1}{r_n} \sum_{i=1}^{n-1} ia_{i+1}(v_{i+1} - v_i) \\ &= v_n - \frac{1}{r_n} \sum_{i=1}^n (r_i - ia_{i+1})(v_{i+1} - v_i). \end{aligned}$$

Since  $|v_{i+1} - v_i| \leq 2/i$ ,

$$|w_n - v_n| \leq \frac{2}{r_n} \sum_{i=1}^n \left| \frac{r_i}{i} - a_{i+1} \right|.$$

Consider now the two following classes of sequences  $(a_i)$ :

- (i)  $a_i = a(i)$ , where  $a$  is a nondecreasing continuous function such that for all positive  $s \in ]0, 1]$ ,  $\lim_{t \rightarrow \infty} \frac{a(ts)}{a(t)} = 1$ .
- (ii)  $a_i = a(i)$ , where  $a$  is a decreasing continuous function such that for all positive  $s \in ]0, 1]$ ,  $\lim_{t \rightarrow \infty} \frac{a(ts)}{a(t)} = 1$ , there exists  $b : [0, 1] \rightarrow \mathbb{R}^+$  measurable such that  $\int_0^1 b(s)ds < \infty$ , and for all  $(s, t) \in ]0, 1] \times \mathbb{R}^+$ ,

$$0 \leq \frac{a(ts)}{a(t)} - 1 \leq b(s).$$

For example,  $a_i = (\log(i + 1))^\alpha$  satisfies (i) for  $\alpha \geq 0$  and (ii) for  $\alpha < 0$ .

LEMMA 4.8. Assume (i) or (ii) holds. Then  $\lim_{n \rightarrow \infty} |w_n - v_n| = 0$ .

Proof. Note that it suffices to prove that  $\frac{r_i}{i} - a_{i+1} = o(a_i)$ . Assume first that (i) holds. Then

$$\begin{aligned} 0 &\leq a_{i+1} - \frac{r_i}{i} \\ &\leq a_i \left( \frac{a_{i+1}}{a_i} - 1 + \int_0^1 \left( 1 - \frac{a(is)}{a(i)} \right) ds \right) \\ &= o(a_i). \end{aligned}$$

Assume now that (ii) holds. Then

$$\begin{aligned} 0 &\leq \frac{r_i}{i} - a_{i+1} \\ &\leq a_i \left( 1 - \frac{a_{i+1}}{a_i} + \int_0^1 \left( \frac{a(is)}{a(i)} - 1 \right) ds \right) \\ &= o(a_i) \end{aligned}$$

by the dominated convergence theorem. □

Let  $\pi(\epsilon, v)$  denote the invariant probability of  $K(x, y, \epsilon, v)$  and  $\pi(v) = \pi(0, v)$ . The map  $(\epsilon, v) \mapsto \pi(\epsilon, v)$  is uniformly continuous. Then the previous lemma implies that when (i) or (ii) holds, since  $\pi_n = \pi(\epsilon_n, w_n)$ ,  $\lim_{n \rightarrow \infty} |\pi_n - \pi(v_n)| = 0$ . This last property with Theorem 2.6 implies the following theorem.

THEOREM 4.9. *Assume that (i) or (ii) holds. Then the limit set of  $(v_n)$  is almost surely an internally chain transitive set of the differential equation*

$$(16) \quad \dot{v} = -v + \pi(v).$$

Note that Proposition 4.7 also holds for sequences  $(a_i)$  satisfying (i) or (ii).

**4.2.3. Exponential reinforcement.** Let  $U : E \times E \rightarrow \mathbb{R}$  be a map. For  $x \in E$  and  $v \in \Delta$ , set

$$U(x, v) = \sum_{y \in E} U(x, y)v(y),$$

$$W(x, y, v) = U(y, v) - U(x, v),$$

$$K(\beta, x, y, v) = \begin{cases} M_0(x, y)\psi[\exp(-\beta W(x, y, v))] & \text{if } x \neq y, \\ 1 - \sum_{y \neq x} K(\beta, x, y, v) & \text{if } x = y, \end{cases}$$

and

$$(17) \quad K_n(x, y, v) = K(\beta_n, x, y, v).$$

Here  $M_0$  is an exploration matrix,  $(\beta_n)_n$  is a positive sequence, and  $\psi$  is given by (7). When  $\beta_n = \beta$ , such a VRRW can be seen as a discrete time version of the self-interacting diffusions on compact manifolds that have been thoroughly analyzed by Benaïm, Ledoux, and Raimond [7] and Benaïm and Raimond (see [8, 9]). When  $\beta_n = A \log(n)$ , the VRRW can be seen as a discrete time version of the self-interacting diffusions on compact manifolds studied by Raimond [27].

Let  $U^\#(\cdot, y)$  be the energy barrier as defined by (9) of the map  $x \mapsto U(x, y)$

THEOREM 4.10. *Consider the VRRW with exponential reinforcement defined by (17). Assume that  $\beta_n = \beta(n)$ , where  $\beta : \mathbb{R}^+ \rightarrow \mathbb{R}^+$  is differentiable, and verify*

$$0 \leq \dot{\beta}(t) \leq \frac{A}{t}$$

for some  $A < 1/2 \max\{U^\#(\cdot, y) : y \in E\}$ . Let

$$C(v) = \Delta(\text{Argmin}U(\cdot, v))$$

denote the set of probabilities supported by  $\text{Argmin}U(\cdot, v)$ . Then the limit set of  $(v_n)$  is an internally chain transitive set of

$$\dot{v} \in -v + C(v).$$

*Proof.* This is an application of Theorem 2.6. The verification of Hypothesis 2.1 is similar to the one given in Proposition 4.4. Details are left to the reader.

It is easily seen that  $C$  is a closed-valued set with convex values. For  $v \in \Delta$ , let

$$\pi_n[v](x) \propto \pi_0(x) \exp(-\beta_n U(x, v))$$

and

$$\pi[v](x) \propto \pi_0(x) \mathbf{1}_{\text{Argmin}U(\cdot, v)}(x).$$

The invariant probability of  $K_n$  is  $\pi_n[v_n]$  and

$$\lim_{n \rightarrow \infty} \pi_n[v](x) = \pi[v](x).$$

This proves that  $C$  is adapted to  $(v_n, \pi_n[v_n])$ , and the result follows from Theorem 2.6.  $\square$

**COROLLARY 4.11** (symmetric interaction). *Assume that the Hypotheses of Theorem 4.10 hold, and assume furthermore that  $U$  is symmetric (i.e.,  $U(x, y) = U(y, x)$ ). Then  $(v_n)$  converges almost surely to a connected component of the set*

$$\{v \in \Delta : v \in C(v)\}.$$

*Proof.* For  $u, v \in \Delta$ , set

$$U(u, v) = \sum_{x, y} U(x, y)u(x)v(y),$$

and let

$$H(v) = \frac{1}{2}U(v, v),$$

We claim that  $H$  is a Lyapunov function of the differential inclusion (5). Let  $t \mapsto v(t)$  be a solution to (5); then, for almost all  $t \geq 0$ ,

$$\begin{aligned} \frac{d}{dt}H(v(t)) &= \frac{1}{2}[U(\dot{v}(t), v(t)) + U(v(t), \dot{v}(t))] = U(\dot{v}(t), v(t)) \\ &= U(\dot{v}(t) + v(t), v(t)) - U(v(t), v(t)) \\ &= \min_x U(x, v(t)) - U(v(t), v(t)), \end{aligned}$$

where we have used the symmetry of  $U$ , the fact that  $\dot{v} + v \in C(v)$ , and the definition of  $C(v)$ . Since  $t \mapsto H(v(t))$  is locally Lipschitz, it is nondecreasing. If now  $t \mapsto H(v(t))$  is constant over a time interval, then  $v(t) \in C(v(t))$  over this time interval. This proves that  $H$  is a Lyapunov function for  $\Lambda = \{v \in \Delta : v \in C(v)\}$ . The result now follows from Proposition 2.9 (compare to Benaïm, Hofbauer, and Sorin [5, Theorem 5.5]), provided we show that  $H(\Lambda)$  has an empty interior.

Let  $v \in \Lambda \cap \text{int}(\Delta)$ . Since the mapping  $x \mapsto U(x, v)$  is constant, for all  $w \in \Delta$ ,  $U(w, v) = U(v, v)$ . Therefore,  $H(v) = U(w, v)$  for all  $w \in \Delta$ . It follows that  $H$  restricted to  $\Lambda \cap \text{int}(\Delta)$  is a constant map. The same reasoning applies to prove that  $H$  restricted to each face of  $\Delta$  is a constant map. We have thus proved that  $H(\Lambda)$  takes finitely many values.  $\square$

*Remark 4.12.* Corollary 4.11 still holds true under the weaker assumption that the map  $v \mapsto U(x, v)$  is smooth and convex in  $v$ .

**COROLLARY 4.13.** *Assume that  $U$  is symmetric and nonnegative and that*

$$\text{Ker}(U) \cap T\Delta = \{0\}.$$

*Then  $\{v \in \Delta : v \in C(v)\}$  reduces to a singleton  $v^*$ , and  $(v_n)$  converges almost surely to  $v^*$ .*

*Proof.* Let  $v \in C(v), w \in \Delta$ , and  $h = w - v$ . Since  $v \in C(v), U(v, h) \geq 0$ . Thus  $U(w, w) - U(v, v) = 2U(v, h) + U(h, h) \geq 0$ , proving that  $v$  is a global minimum of  $v \mapsto U(v, v)$ . Since  $U(h, h) > 0$  for  $h = w - v \neq 0$ , such a global minimum is unique.  $\square$

**4.3. Games.** Consider a two-players game. We let  $E_1$  (respectively,  $E_2$ ) denote the finite *set of actions* available to player 1 (respectively, player 2) and

$$U = (U^1, U^2) : E_1 \times E_2 \rightarrow \mathbb{R} \times \mathbb{R}$$

denote the payoff function of the game. If player 1 and player 2 choose, respectively, the actions  $x \in E_1$  and  $y \in E_2$ , then player 1 gets  $U^1(x, y)$  and player 2 gets  $U^2(x, y)$ .

Let  $((X_n, Y_n))$  denote the sequence of plays. In noncooperative game theory, we assume that  $((X_n, Y_n))$  is adapted to some filtration  $(\mathcal{F}_n)$  and that at the beginning of round  $n+1$ , players have no information on the action to be played by their opponents: for all  $(x, y) \in E_1 \times E_2$  and  $n \in \mathbb{N}$

$$\mathbb{P}(X_{n+1} = x, Y_{n+1} = y | \mathcal{F}_n) = \mathbb{P}(X_{n+1} = x | \mathcal{F}_n) \mathbb{P}(Y_{n+1} = y | \mathcal{F}_n).$$

**4.3.1. Markovian fictitious play.** For  $x \in E_1$  and  $v^2 \in \Delta(E_2)$ , set

$$U^1(x, v^2) = \sum_{z \in F} U^1(x, z) v^2(z).$$

Let

$$v_n^2 = \frac{1}{n} \sum_{i=1}^n \delta_{Y_i}.$$

A well-studied strategy known as “fictitious play” consists for player 1 to play at time  $n+1$  an action maximizing  $U^1(\cdot, v_n^2)$ , that is,

$$(18) \quad X_{n+1} \in \text{Argmax} U^1(\cdot, v_n^2).$$

This strategy relies on the idea that in the absence of information on the next move of his opponent, player 1 assumes that he (the opponent) will play accordingly to the past empirical distribution of his moves. While fictitious play was originally proposed in 1951 by Brown [11] as an algorithm to compute Nash equilibria, it has been recently rediscovered as a “learning model” (see Fudenberg and Kreps [13] and Fudenberg and Levine [14]) and has been extensively studied (see Monderer and Shapley [22]; Benaïm and Hirsch [4]; Hofbauer and Sandholm [16]; and Benaïm, Hofbauer, and Sorin [5, 6]; see also Pemantle [25] for an overview and further references).

Fictitious plays requires solving the maximization problem (18) at each stage of the game. However,

- (a) if the cardinal of  $E_1$  is too large,
- (b) or if players have computational limitations,
- (c) or if players are not allowed to play every action at each time,

then such a computation may be problematic or impossible. An alternative strategy proposed first in Benaïm, Hofbauer, and Sorin [6], based on pairwise comparison of payoffs, is as follows: The strategy of player 1 is such that  $\mathbb{P}(X_{n+1} = y | \mathcal{F}_n) = M_n(X_n, y)$  with  $M_n$  the Markov matrix defined by

$$(19) \quad M_n(x, y) = \begin{cases} M_0(x, y) \psi[\exp(-\beta_n W_n(x, y))] & \text{if } x \neq y, \\ 1 - \sum_{y \neq x} M_n(x, y) & \text{if } x = y, \end{cases}$$

where

$$W_n(x, y) = U^1(x, v_n^2) - U^1(y, v_n^2),$$

$M_0$  is an exploration matrix,  $\psi$  is given by (7), and  $\beta_n$  is an increasing positive sequence. Such a strategy will be called a *Markovian fictitious play strategy*.

Adopting the viewpoint of player 1, we choose as an observation space

$$\Sigma = \Delta(E_1) \times \Delta(E_2)$$

and as an observation variable

$$V_n = (\delta_{X_n}, \delta_{Y_n}).$$

Hence  $(v_n)$  is the *empirical frequency of the actions* played up to time  $n$ , and

$$\hat{V}_n(x) = (\delta_x, \nu_n),$$

where  $\nu_n = E(\delta_{Y_{n+1}} | \mathcal{F}_n)$ .

We let  $U^{1,\#}(y)$  denote the energy barrier, as defined by (9), of the map  $x \mapsto U^1(x, y)$ .

**THEOREM 4.14.** *Assume that player 1 plays a Markovian fictitious play strategy as given by (19). Assume that  $\beta_n = \beta(n)$ , where  $\beta$  is differentiable, and  $\lim_{t \rightarrow \infty} \beta(t) = \infty$ , and verify*

$$0 \leq \dot{\beta}(t) \leq \frac{A}{t}$$

for some  $A < 1/2 \max\{U^{1,\#}(y) : y \in E_2\}$ .

For  $v = (v^1, v^2) \in \Delta(E_1) \times \Delta(E_2)$ , let

$$C_1(v^2) = \Delta(\text{Argmax} U^1(\cdot, v^2))$$

and

$$C(v) = C_1(v^2) \times \Delta(E_2).$$

Then the limit set of  $(v_n)$  is an internally chain transitive set of

$$\dot{v} \in -v + C(v).$$

*Proof.* This is still an application of Theorem 2.6. The verification of Hypothesis 2.1 is similar to the one given in Proposition 4.4. Let

$$\pi_n[v^2](x) \propto \pi_0(x) \exp(\beta_n U^1(x, v^2))$$

and

$$\pi[v^2](x) \propto \pi_0(x) \mathbf{1}_{\text{Argmax}(U^1(\cdot, v^2))}(x).$$

Then the invariant probability of  $M_n$  is  $\pi_n = \pi_n[v_n^2]$  and  $\theta_n = \pi_n \hat{V}_n = (\pi_n, \nu_n)$  with  $\nu_n = E(\delta_{Y_{n+1}} | \mathcal{F}_n)$ . Since  $\pi_n[v^2] \rightarrow \pi[v^2] \in C^1(v^2)$ , it follows that  $C$  is an adapted graph.  $\square$

Much more can be said under the assumption that **both** players adopt a Markovian fictitious play strategy:  $P(X_{n+1} = y | \mathcal{F}_n) = M_n^1(X_n, y)$  and  $P(Y_{n+1} = y | \mathcal{F}_n) = M_n^2(Y_n, y)$  with  $M_n^1$  and  $M_n^2$  the Markov matrices defined by (with  $i \in \{1, 2\}$ )

$$(20) \quad M_n^i(x, y) = \begin{cases} M_0(x, y) \psi[\exp(-\beta_n^i W_n^i(x, y))] & \text{if } x \neq y, \\ 1 - \sum_{y \neq x} M_n^i(x, y) & \text{if } x = y, \end{cases}$$

where

$$\begin{aligned} W_n^1(x, y) &= U^1(x, v_n^2) - U^1(y, v_n^2), \\ W_n^2(x, y) &= U^2(v_n^1, x) - U^2(v_n^1, y), \end{aligned}$$

$M_0^i$  is an exploration matrix,  $\psi$  is given by (7), and  $\beta_n^i$  is an increasing positive sequence.

Let  $\text{Conv}(U)$  denote the convex hull in  $\mathbb{R}^2$  of the set  $\{U(x, y) : x \in E_1, y \in E_2\}$  of all possible payoffs. We now choose

$$\Sigma = \Delta(E_1) \times \Delta(E_2) \times \text{Conv}(U)$$

as an observation space and

$$V_n = (\delta_{X_n}, \delta_{Y_n}, U(X_n, Y_n))$$

as the observation variable. Hence

$$\hat{V}_n(x, y) = (\delta_x, \delta_y, U(x, y)).$$

**THEOREM 4.15.** *Assume that both players adopt a Markovian fictitious play strategy. Assume that for  $i \in \{1, 2\}$ ,  $\beta_n^i = \beta^i(n)$ , where  $\beta^i$  is differentiable, and verify*

$$0 \leq \dot{\beta}^i(t) \leq \frac{A^i}{t}$$

for some  $A^i < 1/2 \max\{U^{i,\#}(y) : y \in E_{3-i}\}$ .

For  $v = (v^1, v^2, u) \in \Delta(E_1) \times \Delta(E_2) \times \text{Conv}(U)$ , let

$$C(v) = \{(\alpha, \beta, \gamma) \in \Sigma : \alpha \in C_1(v^2), \beta \in C_2(v^1), \gamma = U(\alpha, \beta)\},$$

where  $C_1(v^2)$  is like in Theorem 4.14 and  $C_2(v^1)$  is analogously defined for player 2. Then the limit set of  $(v_n)$  is an internally chain transitive set of

$$\dot{v} \in -v + C(v).$$

*Proof.* Let  $(M_n^i)$  denote the strategy of player  $i$ . Let  $\pi_n^i, \lambda_n^i$  be the invariant measure and spectral gap of  $M_n^i$ . On the state space  $E_1 \times E_2$ , the strategy of the pair of players is  $M_n = M_n^1 \otimes M_n^2$  for which the invariant measure is  $\pi_n = \pi_n^1 \otimes \pi_n^2$  and the spectral gap is  $\lambda_n = \min(\lambda_n^1, \lambda_n^2)$ . Thus Hypothesis 2.1 holds for  $(M_n)$ . The rest of the proof is similar to the proof of Theorem 4.14 and is left to the reader.  $\square$

**COROLLARY 4.16** (zero sum games). *Suppose that  $U^2 = -U^1$ . Then under the assumption of Theorem 4.15,  $(v_n^1, v_n^2)$  converges almost surely to the set of Nash equilibria*

$$\{(v_1, v_2) : v_1 \in C_1(v^2), v_2 \in C_2(v^1)\},$$

and  $(U^1(X_n, Y_n))$  converges almost surely to the value of the game

$$u^* = \max_{v^1 \in \Delta(E_1)} \min_{v^2 \in \Delta(E_2)} U^1(v^1, v^2) = \min_{v^1 \in \Delta(E_1)} \max_{v^2 \in \Delta(E_2)} U^1(v^1, v^2).$$

*Proof.* This follows from Theorem 2.6, Proposition 2.7 (ii) and the fact that the set  $\{(v_1, v_2, u) : v_1 \in C_1(v^2), v_2 \in C_2(v^1), u = u^*\}$  is a global attractor of the

differential inclusion, as proved in full generality by Benaïm, Hofbauer, and Sorin [5].  $\square$

**COROLLARY 4.17** (potential games). *Suppose that  $U^2 = U^1$ . Then under the assumption of Theorem 4.15,  $(v_n^1, v_n^2)$  converges almost surely to a connected subset of the set of Nash equilibria*

$$\{(v_1, v_2) : v_1 \in C_1(v^2), v_2 \in C_2(v^1)\}$$

on which  $U^1$  is constant, and  $(U^1(X_n, Y_n))$  converges almost surely toward this constant.

*Proof.* The proof follows from Theorem 2.6, Proposition 2.9, and the fact that  $U^1 = U^2$  is a Lyapunov function of the differential inclusion. The proof of this later point is given by Benaïm, Hofbauer, and Sorin [5, Theorem 5.5]. It is similar to the proof of Corollary 4.11.  $\square$

**4.3.2. A remark on Hypothesis 2.2.** We give here a simple example showing the necessity of Hypothesis 2.2.

Consider the zero sum game, where  $E_1 = E_2 = \{0, 1\}$ ,  $U^1 = -U^2$ , and

$$U^1 = \begin{bmatrix} 0 & -1 \\ -1 & 0 \end{bmatrix}.$$

Let  $V_n = U^1(X_n, Y_n)$  be the payoff to player 1 at time  $n$ . One has

$$\hat{V}_n(x) = U^1(x, 1)\nu_n + U^1(x, 0)(1 - \nu_n)$$

with  $\nu_n = E(Y_{n+1} | \mathcal{F}_n)$ .

Suppose player 1 adopts the strategy given by

$$M_n = M = \begin{bmatrix} \epsilon & 1 - \epsilon \\ 1 - \epsilon & \epsilon \end{bmatrix}$$

for some  $0 < \epsilon < 1$ . Then  $\pi_n = \pi$  with  $\pi(0) = \pi(1) = 1/2$ , and

$$\theta_n = \pi_n \hat{V}_n = -1/2,$$

regardless of the strategy played by player 2.

Suppose now that player 2 plays  $Y_{n+1} = X_n$  for all  $n \geq 1$ . For  $\epsilon \neq 1/2$ , Hypothesis 2.2 is not verified, and the prediction given by (a wrong application of) Theorem 2.6 fails since

$$v_n \rightarrow \sum_{x,y} \pi(x)M(x, y)U^1(x, y) = -(1 - \epsilon).$$

**5. Proof of Theorem 2.6.** Let  $F$  denote a set-valued function mapping each point  $x \in \mathbb{R}^m$  to a set  $F(x) \subset \mathbb{R}^m$ . We call  $F$  a *standard set, valued map*, provided it verifies the three following conditions:

- (i)  $F$  is a closed set-valued map. That is,

$$\text{Graph}(F) = \{(x, y) : y \in F(x)\}$$

is a closed subset of  $\mathbb{R}^m \times \mathbb{R}^m$ .

- (ii)  $F$  has nonempty compact convex set values, meaning that  $F(x)$  is a nonempty compact convex subset of  $\mathbb{R}^m$  for all  $x \in \mathbb{R}^m$ .

(iii) There exists  $c > 0$  such that for all  $x \in \mathbb{R}^m$ ,

$$\sup_{z \in F(x)} \|z\| \leq c(1 + \|x\|),$$

where  $\|\cdot\|$  denotes any norm on  $\mathbb{R}^m$ .

Given a standard set-valued map  $F$ , set

$$F^\delta(u) = \{w \in \mathbb{R}^m : \exists v \in \mathbb{R}^m : d(u, v) \leq \delta, d(w, F(v)) \leq \delta\}.$$

The following proposition follows from the results of Benaïm, Hofbauer, and Sorin [5].

**PROPOSITION 5.1.** *Let  $(x_n)$  and  $(U_n)$  be discrete time processes living in  $\mathbb{R}^m$  and  $(\gamma_n)$  a sequence of nonnegative numbers. Let  $(F_n)$  be a sequence of set-valued maps, and let  $F$  be a standard set-valued map. Assume that*

(i)

$$x_{n+1} - x_n - \gamma_{n+1}U_{n+1} \in \gamma_{n+1}F_n(x_n),$$

(ii)

$$\sum_n \gamma_n = \infty, \quad \lim_{n \rightarrow \infty} \gamma_n = 0,$$

(iii) for all  $T > 0$ ,

$$\lim_{n \rightarrow \infty} \sup \left\{ \left\| \sum_{i=n}^{k-1} \gamma_{i+1}U_{i+1} \right\| : k = n+1, \dots, m(\tau_n + T) \right\} = 0,$$

where

$$\tau_n = \sum_{i=1}^n \gamma_i$$

and

$$(21) \quad m(t) = \sup\{k \geq 0 : t \geq \tau_k\},$$

(iv)  $\sup_n \|x_n\| = M < \infty$ ,

(v) for all  $\delta > 0$ , there exists  $n_0$  such that

$$F_n(x_n) \subset F^\delta(x_n)$$

for all  $n \geq n_0$ .

Then the limit set of  $(x_n)$  is an attractor free set of the dynamics induced by  $F$ .

*Remark 5.2.* This proposition is purely deterministic. If the  $(x_n)$ ,  $(U_n)$  are random processes, the assumptions have to be understood almost surely.

*Remark 5.3.* If condition (v) is strengthened to  $F_n = F$ , Proposition 5.1 follows from Proposition 1.3 and Theorem 4.3 of Benaïm, Hofbauer, and Sorin [5]. Under the weaker hypothesis (v), it suffices to verify that the arguments given in the proof of Proposition 1.3 adapt verbatim.

With the notation of the preceding sections, write

$$v_{n+1} - v_n = \frac{1}{n+1}[-v_n + V_{n+1}] = \frac{1}{n+1}[-v_n + \theta_n + U_{n+1}],$$

where

$$(22) \quad U_{n+1} = V_{n+1} - \theta_n.$$

Hence, conditions (i), (ii), and (iv) of the previous proposition are satisfied with  $F_n(u) = -u + C_n(u)$  and  $\gamma_n = \frac{1}{n}$ . Condition (v) follows from the next lemma.

LEMMA 5.4. *Let  $C$  be adapted. For  $u \in \Sigma$  and  $\delta > 0$ , set*

$$C^\delta(u) = \{w \in \Sigma : \exists v \in \Sigma : d(u, v) \leq \delta, d(w, C(v)) \leq \delta\}.$$

Then for all  $\delta > 0$ , there exists  $n_0$  such that

$$C_n(u) \subset C^\delta(u)$$

for all  $n \geq n_0$  and  $u \in p(C_n)$ .

*Proof.* Let  $\Gamma_n = p(C_n)$ . Assume to the contrary that there exist sequences  $u_{n_k} \in \Gamma_{n_k}$  and  $v_{n_k} \in C_{n_k}(u_{n_k})$  such that  $n_k \rightarrow \infty$  and  $v_{n_k} \notin C^\delta(u_{n_k})$ . By compactness we may assume that  $u_{n_k} \rightarrow u, v_{n_k} \rightarrow v \in C(u)$ . Hence, for  $n_k$  large enough,  $d(u_{n_k}, u) < \delta$  and  $d(v_{n_k}, v) < \delta$ , proving that  $v_{n_k} \in C^\delta(u_{n_k})$ .  $\square$

To conclude the proof of Theorem 2.6, it remains to verify condition (iii) of Proposition 5.1.

LEMMA 5.5. *Under Hypotheses 2.1 and 2.2, the sequence  $(U_n)$  defined by (22) verifies (iii) of Proposition 5.1.*

*Proof.* Set  $\frac{1}{n+1}U_{n+1} = \epsilon_{n+1}^0 + \epsilon_{n+1}$  with

$$\epsilon_{n+1}^0 = \frac{1}{n+1}(V_{n+1} - \hat{V}_n(X_{n+1}))$$

and

$$\begin{aligned} \epsilon_{n+1} &= \frac{1}{n+1}(\hat{V}_n(X_{n+1}) - \pi_n \hat{V}_n) \\ &= \frac{1}{n+1}(Q_n - M_n Q_n) \hat{V}_n(X_{n+1}), \end{aligned}$$

where the last equality follows from the definition of  $Q_n$ . Now write  $\epsilon_{n+1} = \sum_{i=1}^4 \epsilon_{n+1}^i$ , where

$$\begin{aligned} \epsilon_{n+1}^1 &= \frac{1}{n+1}[Q_n \hat{V}_n(X_{n+1}) - M_n Q_n \hat{V}_n(X_n)], \\ \epsilon_{n+1}^2 &= \frac{1}{n+1}M_n Q_n \hat{V}_n(X_n) - \frac{1}{n}M_n Q_n \hat{V}_n(X_n), \\ \epsilon_{n+1}^3 &= \frac{1}{n}M_n Q_n \hat{V}_n(X_n) - \frac{1}{n+1}M_{n+1} Q_{n+1} \hat{V}_{n+1}(X_{n+1}), \\ \epsilon_{n+1}^4 &= \frac{1}{n+1}M_{n+1} Q_{n+1}(\hat{V}_{n+1} - \hat{V}_n)(X_{n+1}), \\ \epsilon_{n+1}^5 &= \frac{1}{n+1}[M_{n+1} Q_{n+1} \hat{V}_n(X_{n+1}) - M_n Q_n \hat{V}_n(X_{n+1})]. \end{aligned}$$

For  $i = 0, \dots, 5$ , let

$$\epsilon_n^i(T) = \sup \left\{ \left\| \sum_{j=n}^{k-1} \epsilon_{j+1}^i \right\| : k = n+1, \dots, m(\tau_n + T) \right\}.$$

We are now going to prove that these six terms converge almost surely toward 0. This will imply that (iii) of Proposition 5.1 is verified.

Note that since  $\Sigma$  is compact, there exists a finite constant  $R$  such that for all  $n$ ,

$$\|V_n\| + \sum_x \|\hat{V}_n(x)\| \leq R.$$

The sequence  $(\epsilon_n^0)$  is a martingale difference with

$$\mathbb{E}(\|\epsilon_{n+1}^0\|^2 | \mathcal{F}_n) \leq R^2 / (n+1)^2.$$

Therefore, by Doob's convergence theorem for  $L^2$  martingales,  $\lim_{n \rightarrow \infty} \epsilon_n^0(T) = 0$  almost surely.

The sequence  $(\epsilon_n^1)$  is a martingale difference with  $\|\epsilon_{n+1}^1\| \leq R|Q_n|/(n+1)$ . Thus by a classical application of exponential martingale inequality (inequality (18) in Benaïm [2]), we have for all positive  $\alpha$

$$\mathbb{P}(\epsilon_n^1(T) \geq \alpha) \leq c \exp\left(\frac{-\alpha^2}{c \sum_{i=n}^{m(\tau_n+T)} (R^2|Q_i|^2/i^2)}\right)$$

for some positive constant  $c$ . By Hypothesis 2.1, for any  $\epsilon > 0$  and  $n \geq 3$  (using that  $ne^T - 1 \leq m(\tau_n + T) \leq (n+1)e^T - 1$ ),

$$\sum_{i=n}^{m(\tau_n+T)} (R^2|Q_i|^2/i^2) \leq \sum_{i=n}^{m(\tau_n+T)} \frac{1}{i} \frac{\epsilon}{\log(i)} \leq \frac{\epsilon(T+1)}{\log(n)}.$$

Thus, choosing  $\epsilon < \frac{\alpha^2}{c(T+1)}$ , we prove

$$\sum_n \mathbb{P}(\epsilon_n^1(T) \geq \alpha) < \infty$$

and  $\lim_{n \rightarrow \infty} \epsilon_{n+1}^1 = 0$  almost surely by the Borel–Cantelli lemma.

For  $n+1 \leq k \leq m(\tau_n + T)$ ,

$$\sum_{j=n}^{k-1} \epsilon_{j+1}^2 = \sum_{j=n}^{k-1} \frac{M_j Q_j \hat{V}_j(X_j)}{(j+1)j}.$$

Thus

$$\epsilon_n^2(T) \leq R \sum_{j=n}^{m(\tau_n+T)} \frac{|Q_j|}{j(j+1)} \leq R(T+1) \sup_{j \geq n} \frac{|Q_j|}{j}.$$

By Hypothesis 2.1, this goes to zero almost surely when  $n \rightarrow \infty$ .

For  $n+1 \leq k \leq m(\tau_n + T)$ ,

$$\sum_{j=n}^{k-1} \epsilon_j^3 = \frac{1}{n} M_n Q_n \hat{V}_n(X_n) - \frac{1}{k} M_k Q_k \hat{V}_k(X_k)$$

so that

$$\epsilon_n^3(T) \leq 2R \sup_{i \geq n} \frac{1}{i} |Q_i|$$

and  $\epsilon_n^3(T) \rightarrow 0$  almost surely as  $n \rightarrow \infty$  by Hypothesis 2.1.

The term  $\epsilon_n^4(T)$  is dominated by

$$(T + 1) \sup_{i \geq n} \sup_x |M_{i+1} Q_{i+1} (\hat{V}_{i+1} - \hat{V}_i)(x)|$$

which converges almost surely toward 0 as  $n \rightarrow \infty$  by Hypothesis 2.2.

Finally, since  $M_n Q_n = Q_n + \Pi_n - I$ ,

$$\epsilon_n^5 = \frac{1}{n+1} \left[ (Q_{n+1} - Q_n) \hat{V}_n(X_{n+1}) + (\Pi_{n+1} - \Pi_n) \hat{V}_n \right].$$

Hence

$$\epsilon_n^5(T) \leq R(T+1) \sup_{i \geq n} (|Q_{i+1} - Q_i| + |\pi_{i+1} - \pi_i|) \rightarrow 0$$

almost surely by Hypothesis 2.1. This completes the proof of (iii).  $\square$

#### REFERENCES

- [1] M. BENAÏM, *Vertex reinforced random walks and a conjecture of pemantle*, Ann. Probab., 25 (1997), pp. 361–392.
- [2] M. BENAÏM, *Dynamics of stochastic approximation algorithms*, in Séminaire de Probabilités XXXIII, Lecture Notes in Math. 1709, Springer, 1999, pp. 1–68.
- [3] M. BENAÏM AND M.W. HIRSCH, *Asymptotic pseudotrajectories and chain recurrent flows, with applications*, J. Dynam. Differential Equations, 8 (1996), pp. 141–176.
- [4] M. BENAÏM AND M.W. HIRSCH, *Mixed equilibria and dynamical systems arising from fictitious play in perturbed games*, Games Econom. Behav., 29 (1999), pp. 36–72.
- [5] M. BENAÏM, J. HOFBAUER, AND S. SORIN, *Stochastic approximations and differential inclusions*, SIAM J. Control Optim., 44 (2005), pp. 328–348.
- [6] M. BENAÏM, J. HOFBAUER, AND S. SORIN, *Stochastic approximations and differential inclusions. Part II: Applications*, Math. Oper. Res., 31 (2006), pp. 673–695.
- [7] M. BENAÏM, M. LEDOUX, AND O. RAIMOND, *Self-interacting diffusions*, Probab. Theory Related Fields, 122 (2002), pp. 1–41.
- [8] M. BENAÏM AND O. RAIMOND, *Self-interacting diffusions II: Convergence in Law*, Ann. Inst. H. Poincaré Probab. Statist., 6 (2003), pp. 1043–1055.
- [9] M. BENAÏM AND O. RAIMOND, *Self-interacting diffusions III: Symmetric interactions*, Ann. Probab., 33 (2005), pp. 1717–1759.
- [10] A. BENVENISTE, M. MÉTIVIER, AND P. PRIOURET, *Stochastic Approximation and Adaptive Algorithms*, Springer, Berlin, 1990.
- [11] G. BROWN, *Iterative solution of games by fictitious play*, in Activity Analysis of Production and Allocation, T.C. Koopmans, ed., Wiley, 1951, pp. 374–376.
- [12] M. DUFLO, *Algorithmes stochastiques*, Math. Appl. (Berlin) 23, Springer, Berlin, 1966.
- [13] D. FUDENBERG AND K. KREPS, *Learning mixed equilibria*, Games Econom. Behav., 5 (1993), pp. 320–367.
- [14] D. FUDENBERG AND D.K. LEVINE, *The Theory of Learning in Games*, MIT Press, Cambridge, MA, 1998.
- [15] B. HAJEK, *Cooling schedules for optimal annealing*, Math. Oper. Res., 13 (1982), pp. 311–329.
- [16] J. HOFBAUER AND W. SANDHOLM, *On the global convergence of stochastic fictitious play*, Econometrica, 70 (2002), pp. 2265–2294.
- [17] R. HOLLEY AND D. STROOCK, *Simulated annealing via Sobolev inequalities*, Comm. Math. Phys., 115 (1988), pp. 553–568.
- [18] H.J. KUSHNER AND C.C. CLARCK, *Stochastic Approximation for Constrained and Unconstrained Systems*, Springer, Berlin, 1978.
- [19] L. LJUNG, *Analysis of recursive stochastic algorithms*, IEEE Trans. Automat. Control, AC-22 (1977), pp. 551–575.
- [20] M. MÉTIVIER AND P. PRIOURET, *Théorèmes de convergence presque sûre pour une classe d’algorithmes stochastiques à pas décroissant*, Probab. Theor. Relat. Fields, 74 (1987), pp. 403–428.
- [21] L. MICLO, *Recuit simulé sans potentiel sur un ensemble fini*, in Séminaire de probabilités XXVI, Springer, Berlin, 1992, pp. 47–60.

- [22] D. MONDERER AND L.S. SHAPLEY, *Fictitious play property for games with identical interests*, J. Econom. Theory, 68 (1996), pp. 258–265.
- [23] R. PEMANTLE, *Random Processes with Reinforcement*, Ph. Dissertation, Massachusetts Institute of Technology, Cambridge, MA, 1988.
- [24] R. PEMANTLE, *Vertex reinforced random walk*, Probab. Theor. Relat. Fields, 92 (1992), pp. 117–136.
- [25] R. PEMANTLE, *A survey of random processes with reinforcement*, Probability Surv., 4 (2007), pp. 1–79.
- [26] R. PEMANTLE AND S. VOLKOV, *Vertex-reinforced random walk on  $\mathbb{Z}$  has finite range*, Ann. Probab., 27 (1999), pp. 1368–1388.
- [27] O. RAIMOND, *Self-interacting diffusions: A simulated annealing version*, Probab. Theor. Relat. Fields., 144 (2009), pp. 247–279.
- [28] L. SALOFF-COSTE, *Lectures on finite Markov chains*, in Lectures on Probability Theory and Statistics 1996, Lecture Notes in Math. 1665, Springer, Berlin, 1997, pp. 201–413.
- [29] P. TARRÈS, *VRRW on  $\mathbb{Z}$  eventually get stuck at a set of five points*, Ann. Probab., 32 (2004), pp. 1455–1478.
- [30] S. VOLKOV, *Vertex-reinforced random walks on arbitrary graphs*, Ann. Probab., 29 (2001), pp. 66–91.