# Efficient Markovian couplings: examples and counterexamples 

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#### Abstract

In this paper we study the notion of an efficient coupling of Markov processes. Informally, an efficient coupling is one which couples at the maximum possible exponential rate, as given by the spectral gap. This notion is of interest not only for its own sake, but also of growing importance arising from the recent advent of methods of "perfect simulation": it helps to establish the "price of perfection" for such methods. In general one can always achieve efficient coupling if the coupling is allowed to "cheat" (if each component's behaviour is affected by future behaviour of the other component), but the situation is more interesting if the coupling is required to be co-adapted. We present an informal heuristic for the existence of an efficient coupling, and justify the heuristic by proving rigorous results and examples in the contexts of finite reversible Markov chains and of reflecting Brownian motion in planar domains.


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## 1 Introduction

We will call a diffusion-coupling "efficient" if it can be used to obtain a sharp estimate for the spectral gap of the operator which is the generator of the
diffusion in question. The main results of this paper show that among well known couplings one can find both efficient and inefficient couplings. Moreover, we give examples of Markov processes for which there is no "efficient" Markovian coupling. We will present techniques which can be used to prove efficiency for many concrete examples of couplings.

Coupling techniques can be applied to obtain various estimates in probability and analysis, both in purely theoretical research and in situations directly related to applications ([32] provides a good introduction; see also [24, 25]). Their importance in applications has recently increased dramatically with the advent of coupling-based "perfect simulation" due to Propp and Wilson [40] and Fill [22]. It is now a matter of pressing importance better to understand the price which must be paid for using such coupling-based approaches - the "price of perfection". One measure of this price is the extent to which coupling occurs at a slower rate than the approach to equilibrium, and this therefore provides a strong motivation for the idea of efficiency and the explorations which we describe below. As pointed out to us by Terry Lyons, in high- or infinitedimensional settings it is of more interest to consider the relationship of perfect simulation to log-Sobolev inequalities (see [18] for a useful expository article on log-Sobolev in the context of finite Markov chains), and we hope to consider this in later work.

The reader is advised that here we are considering only couplings of Markov chains or diffusions which are co-adapted, which is to say that either one of the random processes behaves as a Markov process when we take into account the past of both the random processes in question. This is an important point: it is possible (by rather soft arguments) to produce efficient couplings in which a process is allowed to "cheat" by looking ahead into the future of the other process. See for example $[1,2]$ (which contain much else of relevance to the general concerns which prompted our paper). The imposition of the co-adapted property turns efficiency into a non-trivial notion: it also corresponds to reasonable (though not entirely inevitable) assumptions about how one might implement actual couplings for example in a perfect simulation context.

The concept of strong uniform times $[2,1]$ is also motivated by the desire to get a handle on rates of convergence to stationarity, but uses randomized stopping times rather than coupling ideas, and delivers total variation bounds rather than the $L^{2}$-inspired arguments discussed below. Both the kind of coupling considered here and also strong uniform times have led to practical simulation algorithms (respectively "Coupling From The Past" (CFTP) [40] and a sophisticated rejection sampler [17, 22]). Note however that Matthews [35] uses spectral decomposition to obtain a near-optimal strong stationary time.

The idea of an efficient coupling is illustrated in this paper by two kinds of examples. First we consider continuous time Markov chains with finite state space. These results apply to many "attractive systems," similar to the Ehrenfest model discussed in Example 2.11 below. We restrict ourselves to Markov chains reversible with respect to counting measure (hence with symmetric transition probability functions): the ideas of this paper extend to irreversible chains and Rajesh Nandy is investigating this. The second family of examples is con-
cerned with reflected Brownian motion in planar domains. This is related to work on applications of couplings to estimation of the spectral gap for diffusions on manifolds; recently $[8,10,12,48]$, though the basic idea dates back as far as [19]. An extensive bibliography of the notion of couplings as used in spectral gap theory is to be found in [11]: also see [43] for a useful introductory account of analytical approaches.

Wide applications of the coupling technique inevitably lead to the question, which couplings are "good" and which are not? Chen $[9,10]$ has contributed to this question, introducing a concept of "optimal" couplings. However the terminology is somewhat deceptive, as what is being optimized is a time-varying quantity rather auxiliary to any notion of rapid coupling. In general one expects there to be many different notions of good coupling, depending on whether one seeks high probability of early coupling, high probability of successful coupling, or low exponential moment $\mathbb{E}[\exp (\alpha \tau)]$ of coupling time $\tau$. The notion of "efficient" couplings introduced below isolates those co-adapted couplings which can be used to give a sharp estimate for the "spectral gap." We will show that some Chen-optimal couplings are not efficient because there may be no efficient couplings for some Markov chains. It is natural to expect, although we do not prove it, that some efficient couplings are not Chen-optimal.

We note here that there are of course many other ways of estimating rates of convergence other than coupling: see [42, 41] for examples closely tied to the demands of Markov chain Monte Carlo.

We now present a brief and informal review of the concepts of coupling and spectral gap, and their relationship. Consider a positive-recurrent Markov process $X$, symmetric with respect to some reference measure $m$. For many processes the following eigenfunction expansion holds for the density $p_{t}(x, y)$ relative to $m$ :

$$
\begin{equation*}
p(t, x, y)=c+g(x, y) e^{-\mu_{2} t}+R(t, x, y) . \tag{1.1}
\end{equation*}
$$

The first eigenvalue for the process generator is equal to 0 and the first eigenfunction is the constant function $c$ while $\mu_{2}$ stands for the second eigenvalue and $g$ is a combination of corresponding eigenfunctions; $g(x, y)=\sum_{\varphi} \varphi(x) \varphi(y)$ where the $\varphi$ are orthogonal eigenfunctions with eigenvalue $\mu_{2}$. The remainder $R(t, x, y)$ converges to 0 faster than $e^{-\mu_{2} t}$ as $t \rightarrow \infty$, uniformly in $x$ and $y$ for regular cases. Hence $\mu_{2}$, the "spectral gap" between the first and second eigenvalues, determines the speed of convergence of the transition distribution (density $p(t, x, \cdot)$ ) to the stationary distribution as $t \rightarrow \infty$.

Notice that we may replace $g(x, y)$ by $\varphi(x) \varphi(y)$ when the second eigenvalue is multiplicity-free. Notice also that reversibility considerably simplifies the above analysis, since otherwise the multipliers of the $g(x, y)$ term may include a factor which is polynomial in $t$. Fortunately reversibility holds in many of the most important applications.

A "coupling" is a pair of (typically dependent) copies of the Markov process $X$, the first one, $X^{1}$, starting from $x_{1}$ and the second one, $X^{2}$, starting from $x_{2}$. "Good" couplings are characterized by small coupling time $\tau$, the minimum
time $t$ for which $X_{t}^{1}$ and $X_{t}^{2}$ are equal. Applications of the coupling technique depend on the fact that we may, and we do, construct $X^{1}$ and $X^{2}$ in such a way that $X_{t}^{1}=X_{t}^{2}$ for all $t \geq \tau$.

The eigenfunction representation (1.1) now gives

$$
\begin{align*}
& p\left(t, x_{1}, y\right)-p\left(t, x_{2}, y\right)= \\
& \quad\left(g\left(x_{1}, y\right)-g\left(x_{2}, y\right)\right) e^{-\mu_{2} t} \quad+\quad R\left(t, x_{1}, y\right)-R\left(t, x_{2}, y\right), \tag{1.2}
\end{align*}
$$

while the coupling yields

$$
\begin{align*}
& \left|p\left(t, x_{1}, y\right) d y-p\left(t, x_{2}, y\right) d y\right| \\
& \quad=\left|\mathbb{P}\left(X_{t}^{1} \in d y \mid X_{0}^{1}=x_{1}\right)-\mathbb{P}\left(X_{t}^{2} \in d y \mid X_{0}^{2}=x_{2}\right)\right| \\
& \quad \leq \mathbb{P}\left(X_{t}^{1} \in d y, t<\tau \mid X_{0}^{1}=x_{1}\right)+\mathbb{P}\left(X_{t}^{2} \in d y, t<\tau \mid X_{0}^{2}=x_{2}\right) \tag{1.3}
\end{align*}
$$

Suppose that one can prove that $\mathbb{P}\left(\tau>t \mid X_{t}^{1}=x_{1}, X_{t}^{2}=x_{2}\right) \approx e^{-\mu^{*} t}$ for "generic" $x_{1} \neq x_{2}$. Given a suitable sense for "generic", we can combine (1.2) and (1.3) to show that $\mu^{*}$ is a lower bound for $\mu_{2}$ (and in most applications it is the lower bound which counts). We will call $\mu^{*}$ the coupling exponent.

The above argument has been used in various forms to estimate $\mu_{2}$, as for example in [12, Theorem 1.7], or [48].

The following informal definition and heuristic capture the spirit of the results and examples of this paper.

Informal Definition of Efficiency. We will call a coupling ( $X^{1}, X^{2}$ ) an efficient Markovian coupling if $\left(X^{1}, X^{2}\right)$ is a Markov process and $\mu^{*}=\mu_{2}$.

Informal Efficient Coupling Heuristic. A coupling ( $X^{1}, X^{2}$ ) is efficient if and only if, for all $t$, and given $\{t<\tau\}$, the conditional distributions of $\left(X_{t}^{1}, X_{t}^{2}\right)$ and $\left(X_{t}^{2}, X_{t}^{1}\right)$ are singular with respect to each other.

The above heuristic is not true in a rigorous sense, as will be shown in Section 2 , but it works in sufficiently many circumstances to make it "almost true."

The efficient coupling heuristic is closely related to "monotonicity" in the sense of [31, §II.2]. The connection will be made more precise in Theorem 2.6 below. The importance of monotonicity or ordering for effective estimation of the rate of convergence of a Markov chain to its stationary distribution is clear, for example, in [34]. The results in Section 2 below are closely related to those in that paper except that our focus is different - the couplings are the main object of study in this paper, rather than just an effective technical tool. The literature on estimating the rate of convergence for Markov processes is enormous. The forthcoming book [1] or the articles [26, 36] may serve as starting points. The importance of monotonicity is explained in [22, §4].

Also note that [1, Chapter 14 §7.1] describes another way of measuring efficiency for a Markov coupling, in the sense of mean coupling time of a graphbased Markov chain of size $n$ increasing fast with $n$.

We point out that our use of "monotonicity" is somewhat different from that in the sources quoted above. The difference is perhaps best explained by the example of obtuse and acute triangles, discussed in Section 3. The triangles in both families can be expressed as partially ordered sets in a similar way, but efficient couplings for reflected Brownian motion exist only in obtuse triangles, as far as we can tell.

On the practical side, being able to construct an efficient Markovian coupling does not guarantee having a good estimate for the rate of convergence of the process to the stationary distribution. Estimating the coupling exponent $\mu^{*}$ itself may be a hard task, especially when the state space and the transition probabilities do not have a simple structure. Note however that the ideas of perfect simulation [22, 40] finesse this problem away in suitably regular cases.

We will give several distinct formal definitions of efficiency, one for Markov processes with finite state space and continuous time, and two for reflected Brownian motion in planar domains. The goal of the paper is to introduce the idea of efficiency and some accompanying techniques, not to provide a rigid definition and a general theorem. We will adjust our definition to fit particular families of Markov processes and couplings.

The remainder of the paper consists of three sections. Section 2 is devoted to continuous time Markov chains with finite state space. Section 3 studies the mirror coupling of reflected Brownian motions in triangles. Section 4 presents a few informal examples involving mirror couplings for reflected Brownian motion in planar domains. This last section is specialized to a very narrow family of processes but it is at least partly justified by the fact that the technique developed for this case has been subsequently successfully applied in a different context; indeed the methods discussed below have already been used in [3] to prove a number of positive results on the "hot spots" conjecture of J. Rauch, and more recently to construct a counterexample to the conjecture $[5,6]$.

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## 2 Couplings for symmetric Markov chains with finite state space

We devote this section to symmetric Markov processes with continuous time and a finite state space, where the reference measure $m$ is counting measure. This
simple case fully illustrates the main idea of our test for efficiency but avoids technical issues which arise when the state space is continuous.

Thus $X=\left\{X_{t}: t \geq 0\right\}$ is a continuous time symmetric Markov process with a finite state space $D$ and transition probabilities $p(t, y, x)=p(t, x, y)=$ $\mathbb{P}\left(X_{s+t}=y \mid X_{s}=x\right)$. The following eigenvalue expansion (1.1) holds for $X$ [14, p. 183]:

$$
\begin{equation*}
p(t, x, y)=c+g(x, y) e^{-\mu_{2} t}+R(t, x, y) \tag{2.1}
\end{equation*}
$$

Here $g(x, y)$ is a combination of eigenfunctions corresponding to the second eigenvalue $\mu_{2}$, as in (1.1) and the remainder $R(t, x, y)$ converges to 0 faster than $e^{-\mu_{2} t}$ when $t \rightarrow \infty$.

Suppose now that $\left(X^{1}, X^{2}\right)$ is a Markovian coupling for the process $X$. That is to say, each of the three processes $\left\{\left(X_{t}^{1}, X_{t}^{2}\right): t \geq 0\right\},\left\{X_{t}^{1}: t \geq 0\right\}$ and $\left\{X_{t}^{2}: t \geq 0\right\}$ is Markov with respect to the filtration generated jointly by $X^{1}$ and $X^{2}$, and the processes $X^{1}$ and $X^{2}$ have the same transition probabilities as $X$. We call the time $\tau=\inf \left\{t \geq 0: X_{t}^{1}=X_{t}^{2}\right\}$ the coupling time for $X^{1}$ and $X^{2}$. It is convenient to stipulate that $X_{t}^{1}=X_{t}^{2}$ for all $t \geq \tau$.

We also require that the coupling is invariant under the transposition of its components; this is to say that the transition probabilities for $\left(X^{1}, X^{2}\right)$ are the same as for $\left(X^{2}, X^{1}\right)$. In fact this entails no loss in generality, since standard stochastic control arguments (randomizing between an asymmetric coupling transition kernel and its transposition) show that coupling times are stochastically minorized by those obtained by transposition-invariant couplings.

Since the state space $D^{2}=D \times D$ for the Markov process $\left(X^{1}, X^{2}\right)$ is finite, we can apply the Perron-Frobenius theory for nonnegative matrices to the transition probability matrix of the coupling process $\left(X^{1}, X^{2}\right)$. From this we deduce that there exists a $\mu^{\prime}=\mu^{\prime}\left(x_{1}, x_{2}\right)$, the coupling exponent function, such that for all $t \geq 0$ and for all $\left(x_{1}, x_{2}\right) \in D^{2}$

$$
c_{1}\left(x_{1}, x_{2}\right) e^{-\mu^{\prime} t} \leq \mathbb{P}\left(\tau>t \mid\left(X_{0}^{1}, X_{0}^{2}\right)=\left(x_{1}, x_{2}\right)\right) \quad \leq \quad c_{2}\left(x_{1}, x_{2}\right) e^{-\mu^{\prime} t}
$$

(Of course for $x_{1}=x_{2}=x$ we take $c_{1}(x, x)=c_{2}(x, x)=0$ and then $\mu^{\prime}(x, x)$ is not well-defined.) We set

$$
\begin{equation*}
\mu^{*}=\min _{x_{1}, x_{2} \in D} \mu^{\prime}\left(x_{1}, x_{2}\right) \tag{2.2}
\end{equation*}
$$

the coupling exponent. The following simple inequality can be found for example in [12, Thm. 1.7], but we state and prove it here for the sake of completeness. Our main result about finite state space Markov processes uses a generalization of its proof.

Lemma 2.1 If $g\left(x_{1}, \cdot\right)$ and $g\left(x_{2}, \cdot\right)$ are not identical then $\mu^{\prime}\left(x_{1}, x_{2}\right) \leq \mu_{2}$. It follows that we always have $\mu^{*} \leq \mu_{2}$.

Proof : Choose $x_{1}, x_{2} \in D$ and $y \in D$ such that $g\left(x_{1}, y\right) \neq g\left(x_{2}, y\right)$. By equation (2.1),

$$
\begin{align*}
& p\left(t, x_{1}, y\right)-p\left(t, x_{2}, y\right)= \\
& \quad\left[g\left(x_{1}, y\right)-g\left(x_{2}, y\right)\right] e^{-\mu_{2} t}+R\left(t, x_{1}, y\right)+R\left(t, x_{2}, y\right) \tag{2.3}
\end{align*}
$$

Another representation for the same quantity comes from the coupling, namely,

$$
\begin{aligned}
& \left|p\left(t, x_{1}, y\right)-p\left(t, x_{2}, y\right)\right|=\left|\mathbb{P}\left(X_{t}^{1}=y \mid X_{0}^{1}=x_{1}\right)-\mathbb{P}\left(X_{t}^{2}=y \mid X_{0}^{2}=x_{2}\right)\right| \\
& \quad=\left|\mathbb{P}\left(X_{t}^{1}=y, t<\tau \mid X_{0}^{1}=x_{1}\right)-\mathbb{P}\left(X_{t}^{2}=y, t<\tau \mid X_{0}^{2}=x_{2}\right)\right| \\
& \quad \leq \quad \mathbb{P}\left(t<\tau \mid X_{0}^{1}=x_{1}, X_{0}^{2}=x_{2}\right) \\
& \quad \leq c e^{-\mu^{\prime}\left(x_{1}, x_{2}\right) t}
\end{aligned}
$$

Since this estimate and equation (2.3) both hold for arbitrarily large $t$, we can use the condition on our choices of $x_{1}, x_{2}, y$ to show that $\mu^{\prime}\left(x_{1}, x_{2}\right) \leq \mu_{2}$.

This shows that the worst-case coupling exponential decay is never faster than the exponential decay rate of convergence to equilibrium (to wit, the second eigenvalue), and thus motivates our definition of an efficient coupling.

Definition 2.2 Recall the coupling exponent $\mu^{*}=\min _{x_{1}, x_{2} \in D} \mu^{\prime}\left(x_{1}, x_{2}\right)$ defined by equation (2.2). The coupling $\left(X^{1}, X^{2}\right)$ will be called efficient if $\mu^{*}=\mu_{2}$.

Before we state and prove some tests for efficiency, we discuss its definition. The definition is intended to encapsulate a desirable property of couplings in the context of spectral gap estimation. Suppose that one can prove that $\mu^{\prime}\left(x_{1}, x_{2}\right) \geq$ $a$ for some $a$ and some pair of points $x_{1}, x_{2} \in D$. Does it necessarily follow that $\mu_{2} \geq a$ ? Example 2.3 below shows that the answer is "no." In order to prove this lower bound for the spectral gap it suffices, in view of Lemma 2.1, to show that $\mu^{\prime}\left(x_{1}, x_{2}\right) \geq a$ for some $x_{1}, x_{2} \in D$ with $g\left(x_{1}, \cdot\right)$ and $g\left(x_{2}, \cdot\right)$ not identical. A practical strategy might be to prove that $\mu^{\prime}\left(x_{1}, x_{2}\right) \geq a$ for all distinct $x_{1}, x_{2} \in D$. This will be illustrated in Example 2.11 below. We proceed with the aforementioned example showing that the bound $\mu^{\prime}\left(x_{1}, x_{2}\right) \geq a$ for a single pair $\left(x_{1}, x_{2}\right)$ does not necessarily imply the same bound for the spectral gap.

Example 2.3 Let $D=\{0,1, \ldots, 100\} \times\{0,1, \ldots, 10\}$. Suppose the process $X$ jumps only to its nearest neighbors in $D$, and let the jump rate be equal to 1 for every pair of neighbors. (So $X$ is a reflected simple symmetric random walk on a rectangular portion of the planar square lattice $\mathbb{Z}^{2}$.) Consider a coupling $\left(X^{1}, X^{2}\right)$ such that
(a) the first components of $X^{1}$ and $X^{2}$ are collectively independent of the second components, so that we can describe the joint evolution of $\left(X^{1}, X^{2}\right)$ by specifying how the first components behave and separately how the second components behave;
(b) the first components of $X^{1}$ and $X^{2}$ evolve independently till they first agree, after which they remain equal;
(c) the second components behave similarly (independent till they first agree, after which they stick together).

The components of $X$ are independent so it is not hard to check that the second eigenvalue $\mu_{2}$ for $X$ is the same as for its first component, However, if $X_{0}^{1}=$ $(a, b)$ and $X_{0}^{2}=(a, c)$ with $b \neq c$ then $\mu^{\prime}((a, b),(a, c))$ is the same as the second eigenvalue for the second component of $X$ and so it is strictly larger than $\mu_{2}$.

Of course it can be checked directly here that, in the notation of Lemma 2.1, $g((a, b), \cdot) \equiv g((a, c), \cdot)$.

Definition 2.4 We say that $\left(y_{1}, y_{2}\right)$ is accessible from $\left(x_{1}, x_{2}\right)$ if

$$
\mathbb{P}\left[\left(X_{t}^{1}, X_{t}^{2}\right)=\left(y_{1}, y_{2}\right) \mid\left(X_{0}^{1}, X_{0}^{2}\right)=\left(x_{1}, x_{2}\right)\right]>0
$$

for some (and, therefore, all) $t>0$.
Accessibility is clearly a transitive property: if $\left(y_{1}, y_{2}\right)$ is accessible from $\left(x_{1}, x_{2}\right)$ and $\left(z_{1}, z_{2}\right)$ is accessible from $\left(y_{1}, y_{2}\right)$ then $\left(z_{1}, z_{2}\right)$ is accessible from $\left(x_{1}, x_{2}\right)$. Note that accessibility of states in $D^{2}$ is a property of the coupling and not just $X$ : in fact it is the same as accessibility for the coupling Markov chain $\left(X^{1}, X^{2}\right)$ restricted to $D^{2} \backslash \Delta$, where $\Delta \subset D^{2}$ is the diagonal.

Definition 2.5 We will say that the coupling $\left(X^{1}, X^{2}\right)$ has the transposition property relative to $x_{1}$ if
(A) for all $x_{2}$ with $x_{1} \neq x_{2}$, and for all $y_{1}$ and $y_{2}$ with $y_{1} \neq y_{2}$, the accessibility of ( $y_{1}, y_{2}$ ) from ( $x_{1}, x_{2}$ ) implies accessibility of both $\left(x_{1}, x_{2}\right)$ and $\left(x_{2}, x_{1}\right)$ from $\left(y_{1}, y_{2}\right)$;
(B) for every $x_{2} \neq x_{1}$, there is at least one pair $\left(y_{1}, y_{2}\right)$, distinct from $\left(x_{1}, x_{2}\right)$, which is accessible from $\left(x_{1}, x_{2}\right)$ (so for $\left(X_{1}, X_{2}\right)$ restricted to $D^{2} \backslash \Delta$ there are no isolated states involving $x_{1}$ ).

We will say that $D^{2}$ is irreducible with respect to a given coupling if every state ( $y_{1}, y_{2}$ ) with $y_{1} \neq y_{2}$ is accessible from any other state ( $x_{1}, x_{2}$ ) with $x_{1} \neq$ $x_{2}$.

Note that the transposition property relative to $x_{1}$ can also be reduced to considerations about state classification for the coupling chain $\left(X^{1}, X^{2}\right)$, bearing in mind that we are only considering coupling chains which are symmetric under the permutation $\left(x_{1}, x_{2}\right) \leftrightarrow\left(x_{2}, x_{1}\right)$. It can be shown that the transposition property is equivalent to the requirement that, for any $x_{2}$ with $x_{2} \neq x_{1}$, the communicating class of ( $x_{1}, x_{2}$ ) under the chain ( $X^{1}, X^{2}$ ) restricted to $D^{2} \backslash \Delta$ is essential and is saturated under the symmetry $\left(x_{1}, x_{2}\right) \leftrightarrow\left(x_{2}, x_{1}\right)$.

The following fundamental result uses these properties to identify many chains which are efficient and many chains which are not.

## Theorem 2.6

(i) If the coupling $\left(X^{1}, X^{2}\right)$ has the transposition property relative to a point $x_{1} \in D$ then $\mu^{\prime}\left(x_{1}, x_{2}\right)<\mu_{2}$ for some $x_{2} \in D^{2}$ distinct from $x_{1}$, and so the coupling is not efficient.
(ii) Suppose that for some $x_{1}, x_{2} \in D$ there exists a function $f: D \rightarrow \mathbb{R}$ with the property that $f\left(X_{t}^{1}\right)-f\left(X_{t}^{2}\right)$ almost surely remains strictly positive for $t<\tau$, given $X_{0}^{1}=x_{1}$ and $X_{0}^{2}=x_{2}$. Then $\mu^{\prime}\left(x_{1}, x_{2}\right) \geq \mu_{2}$.

The following statement follows immediately from Theorem 2.6.

## Corollary 2.7

(i) If $D^{2}$ is irreducible for the coupling $\left(X^{1}, X^{2}\right)$ then this coupling is not efficient.
(ii) Suppose that for every pair of distinct points $x_{1}, x_{2} \in D$ there exists a function $f: D \rightarrow \mathbb{R}$ with the property that $f\left(X_{t}^{1}\right)-f\left(X_{t}^{2}\right)$ almost surely remains strictly positive for $t<\tau$., given $X_{0}^{1}=x_{1}$ and $X_{0}^{2}=x_{2}$. Then the coupling is efficient.

Remark : An interesting example to which Corollary 2.7 applies is the independence sampler discussed in [46]. Here one can compute the eigenvalues (and indeed the transition matrix) explicitly [33, 45] and verify (at least for finite state space) that the Markov chain is efficient. As pointed out by Cai [7], this chain possesses a monotonicity structure. Cai uses this monotonicity to build a CFTP algorithm, but it also guarantees efficiency as above.

We will show in Examples 2.9 and 2.10 that neither of the conditions in parts (i) and (ii) of Theorem 2.6 is necessary. Theorem 2.6 should be compared with [40, §5], which implies that bounded monotone Markov chains are efficient.

## Proof of Theorem 2.6:

(i) Suppose that the coupling $\left(X^{1}, X^{2}\right)$ has the transposition property relative to some $x_{1}$. Fix any point $x_{2} \in D$ distinct from $x_{1}$ and assume that $\left(X_{0}^{1}, X_{0}^{2}\right)=\left(x_{1}, x_{2}\right)$. Consider a coupling $\left(\tilde{X}^{1}, \tilde{X}^{2}\right)$ which is independent of $\left(X^{1}, X^{2}\right)$, having the same transition probabilities as $\left(X^{1}, X^{2}\right)$ but starting from $\left(\tilde{x}_{2}, x_{1}\right)$ rather than from $\left(x_{1}, x_{2}\right)$. Let $\tilde{\tau}$ denote the coupling time for $\left(\tilde{X}^{1}, \tilde{X}^{2}\right)$.
We will estimate the chance that $\left(X^{1}, X^{2}\right)$ and ( $\tilde{X}^{1}, \tilde{X}^{2}$ ) have not met before time $s$, given $\{\tau>s, \tilde{\tau}>s\}$. The invariance of the coupling under the transposition of its components, the transitivity of the accessibility property and the transposition property relative to $x_{1}$ can be used to show that the accessibility of any $\left(y_{1}, y_{2}\right)$ from $\left(x_{1}, x_{2}\right)$ implies accessibility of $\left(y_{2}, y_{1}\right)$ from $\left(x_{1}, x_{2}\right)$. Consider any integer $k>0$. Since we have assumed that $\left(X_{0}^{1}, X_{0}^{2}\right)=\left(x_{1}, x_{2}\right)$ and $\left(\tilde{X}_{0}^{1}, \tilde{X}_{0}^{2}\right)=\left(x_{2}, x_{1}\right)$, at time $t$ both processes
can only take those pairs of values from which they can reach ( $x_{1}, x_{2}$ ) (we are using the transposition property here). The state space is finite, so the processes can reach $\left(x_{1}, x_{2}\right)$ within an arbitrarily small time, less than $1 / 4$ say, with some strictly positive probability not depending on their values at time $t=k$. We will need a stronger version of this statement. Condition the processes $\left(X_{t}^{1}, X_{t}^{2}\right)$ and $\left(\tilde{X}_{t}^{1}, \tilde{X}_{t}^{2}\right)$ on their values at times $t=k$ and $t=k+1$ (we consider only the values that can be taken with strictly positive probabilities). The transposition property can be used again to show that there are possible trajectories which take the processes to the intermediate state $\left(x_{1}, x_{2}\right)$ for all $t \in[k+1 / 4, k+3 / 4]$. The finiteness of the state space can be now invoked to see that such an event has a probability bounded below. More precisely, there exists a probability $p>0$ such that for all integers $k>0$ and all $s \geq k+1$,

$$
\begin{aligned}
& \mathbb{P}\left[X_{t}^{1}=\tilde{X}_{t}^{1}=x_{1}, X_{t}^{2}=\tilde{X}_{t}^{2}=x_{2} \text { for } k+1 / 4 \leq t \leq k+3 / 4\right. \\
& \left.\quad \mid X_{k}^{1}, X_{k+1}^{1}, X_{k}^{2}, X_{k+1}^{2}, \tilde{X}_{k}^{1}, \tilde{X}_{k+1}^{1}, \tilde{X}_{k}^{2}, \tilde{X}_{k+1}^{2}, \tau>s, \tilde{\tau}>s\right]>p .
\end{aligned}
$$

We are interested only in large $s$ so we will assume that $s>1$. Let $j$ be the largest integer with $j \leq s$. Condition the processes $\left(X^{1}, X^{2}\right)$ and $\left(\tilde{X}^{1}, \tilde{X}^{2}\right)$ on their values at times $0,1,2, \ldots, j$ and on the event $\{\tau>s, \tilde{\tau}>s\}$. It follows that

$$
\begin{align*}
& \mathbb{P}\left[\left(X_{t}^{1}, X_{t}^{2}\right) \neq\left(\tilde{X}_{t}^{1}, \tilde{X}_{t}^{2}\right) \text { for } 0 \leq t<s \mid \tau>s, \tilde{\tau}>s\right] \\
& \quad \leq(1-p)^{j} \leq e^{-\mu s}, \tag{2.4}
\end{align*}
$$

for appropriate $\mu>0$.
Let $\sigma$ be the smallest $t$ such that $\left(X_{t}^{1}, X_{t}^{2}\right)=\left(\tilde{X}_{t}^{1}, \tilde{X}_{t}^{2}\right)$. Thus inequality (2.4) gives

$$
\begin{equation*}
\mathbb{P}\left[\sigma>s \mid \tau>s, \tilde{\tau}>s, X_{0}^{1}=x_{1}, X_{0}^{2}=x_{2}\right] \leq e^{-\mu s} \tag{2.5}
\end{equation*}
$$

We now make a simple observation about the relative behaviour of the processes $\left(X^{1}, X^{2}\right)$ and $\left(\tilde{X}^{1}, \tilde{X}^{2}\right)$ after the meeting time $\sigma$. We use the strong Markov property applied to $\left(X^{1}, X^{2}, \tilde{X}^{1}, \tilde{X}^{2}\right)$ at time $\sigma$, to deduce the following:

$$
\begin{aligned}
& \mathbb{P}\left[X_{s}^{1}=y, s \geq \sigma \mid X_{\sigma}^{1}, X_{\sigma}^{2}, \tilde{X}_{\sigma}^{1}, \tilde{X}_{\sigma}^{2}, \sigma, s<\tau, s<\tilde{\tau}, X_{0}^{1}=x_{1}, \tilde{X}_{0}^{1}=x_{2}\right] \\
= & \mathbb{P}\left[\tilde{X}_{s}^{1}=y, s \geq \sigma \mid X_{\sigma}^{1}, X_{\sigma}^{2}, \tilde{X}_{\sigma}^{1}, \tilde{X}_{\sigma}^{2}, \sigma, s<\tau, s<\tilde{\tau}, X_{0}^{1}=x_{1}, \tilde{X}_{0}^{1}=x_{2}\right]
\end{aligned}
$$

(we suppress the conditioning on $X_{0}^{2}=x_{2}$ and $\tilde{X}_{0}^{2}=x_{1}$, since by definition we have $X_{0}^{2}=\tilde{X}_{0}^{1}$ and $\tilde{X}_{0}^{2}=X_{0}^{1}$ ). Hence we can use integration, and rewrite the conditioning in terms of $\left(X_{0}^{1}, X_{0}^{2}\right)$, respectively ( $\tilde{X}_{0}^{1}, \tilde{X}_{0}^{2}$ ), to show that

$$
\begin{aligned}
& \mathbb{P}\left[X_{s}^{1}=y, s \geq \sigma \mid s<\tau, s<\tilde{\tau}, X_{0}^{1}=x_{1}, X_{0}^{2}=x_{2}\right] \\
& \quad=\quad \mathbb{P}\left[\tilde{X}_{s}^{1}=y, s \geq \sigma \mid s<\tau, s<\tilde{\tau}, \tilde{X}_{0}^{1}=x_{2}, \tilde{X}_{0}^{2}=x_{1}\right]
\end{aligned}
$$

We now generalize the proof of Lemma 2.1. First we find points $x_{2}, y \in D$ such that $g\left(x_{1}, y\right) \neq g\left(x_{2}, y\right)$. Note that such points exist using the orthogonality of the eigenfunctions $\varphi$ in $g(x, y)=\sum_{\varphi} \varphi(x) \varphi(y)$, and the fact that eigenfunctions corresponding to the second eigenvalue must change sign.
We recall (2.3), namely,

$$
\begin{align*}
& p\left(s, x_{1}, y\right)-p\left(s, x_{2}, y\right)= \\
& \quad\left[g\left(x_{1}, y\right)-g\left(x_{2}, y\right)\right] e^{-\mu_{2} s}+R\left(s, x_{1}, y\right)+R\left(s, x_{2}, y\right) . \tag{2.6}
\end{align*}
$$

The estimate based on coupling is more complicated in the present case. We use symmetry, inequality (2.5), and follow the method described in the proof of Lemma 2.1. We find

$$
\begin{aligned}
& \left|p\left(s, x_{1}, y\right)-p\left(s, x_{2}, y\right)\right| \\
= & \left|\mathbb{P}\left(X_{s}^{1}=y \mid X_{0}^{1}=x_{1}\right)-\mathbb{P}\left(X_{s}^{2}=y \mid X_{0}^{2}=x_{2}\right)\right| \\
= & \mid \mathbb{P}\left(X_{s}^{1}=y, s<\tau \mid X_{0}^{1}=x_{1}, X_{0}^{2}=x_{2}\right) \\
& \quad-\mathbb{P}\left(X_{s}^{2}=y, s<\tau \mid X_{0}^{1}=x_{1}, X_{0}^{2}=x_{2}\right) \mid \\
= & \mid \mathbb{P}\left(X_{s}^{1}=y \mid s<\tau, X_{0}^{1}=x_{1}, X_{0}^{2}=x_{2}\right) \mathbb{P}\left(s<\tau \mid X_{0}^{1}=x_{1}, X_{0}^{2}=x_{2}\right) \\
& -\mathbb{P}\left(X_{s}^{2}=y \mid s<\tau, X_{0}^{1}=x_{1}, X_{0}^{2}=x_{2}\right) \mathbb{P}\left(s<\tau \mid X_{0}^{1}=x_{1}, X_{0}^{2}=x_{2}\right) \mid \\
\leq & \mid \mathbb{P}\left(X_{s}^{1}=y \mid s<\tau, X_{0}^{1}=x_{1}, X_{0}^{2}=x_{2}\right) \\
& \quad-\mathbb{P}\left(X_{s}^{2}=y \mid s<\tau, X_{0}^{1}=x_{1}, X_{0}^{2}=x_{2}\right) \mid \times c\left(x_{1}, x_{2}\right) e^{-\mu^{\prime}\left(x_{1}, x_{2}\right) s} \\
= & \mid \mathbb{P}\left(X_{s}^{1}=y \mid s<\tau, s<\tilde{\tau}, X_{0}^{1}=x_{1}, X_{0}^{2}=x_{2}\right) \\
& \quad-\mathbb{P}\left(\tilde{X}_{s}^{1}=y \mid s<\tau, s<\tilde{\tau}, \tilde{X}_{0}^{1}=x_{2}, \tilde{X}_{0}^{1}=x_{1}\right) \mid \times c e^{-\mu^{\prime} s}
\end{aligned}
$$

This last step uses the fact that $\left(X^{1}, X^{2}\right),\left(\tilde{X}^{1}, \tilde{X}^{2}\right)$ are independent to justify the insertion of conditioning on both $s<\tau$ and $s<\tilde{\tau}$ for both probabilities. We now use the conditional probability identity noted above to cancel between the two conditional probabilities to yield:

$$
\begin{aligned}
& \mid \mathbb{P}\left(X_{s}^{1}=y \mid s<\tau, s<\tilde{\tau}, X_{0}^{1}=x_{1}, X_{0}^{2}=x_{2}\right) \\
& \quad-\mathbb{P}\left(\tilde{X}_{s}^{1}=y \mid s<\tau, s<\tilde{\tau}, \tilde{X}_{0}^{1}=x_{2}, \tilde{X}_{0}^{2}=x_{1}\right) \mid \times c e^{-\mu^{\prime} s} \\
= & \mid \mathbb{P}\left(X_{s}^{1}=y, s<\sigma \mid s<\tau, s<\tilde{\tau}, X_{0}^{1}=x_{1}, X_{0}^{2}=x_{2}\right) \\
& \quad-\mathbb{P}\left(\tilde{X}_{s}^{1}=y, s<\sigma \mid s<\tau, s<\tilde{\tau}, \tilde{X}_{0}^{1}=x_{2}, \tilde{X}_{0}^{2}=x_{1}\right) \mid \times c e^{-\mu^{\prime} s} \\
\leq & \mathbb{P}\left(s<\sigma \mid s<\tau, s<\tilde{\tau}, X_{0}^{1}=x_{1}, \tilde{X}_{0}^{1}=x_{2}\right) c e^{-\mu^{\prime} s} \\
\leq & e^{-\mu s} c e^{-\mu^{\prime} s} .
\end{aligned}
$$

The last estimate and (2.6) hold for arbitrarily large $s$, so $\mu_{2} \geq \mu^{\prime}\left(x_{1}, x_{2}\right)+$ $\mu>\mu^{\prime}\left(x_{1}, x_{2}\right)$ and we see that the coupling is not efficient.
(ii) Find a function $f: D \rightarrow \mathbb{R}$ corresponding to $x_{1}, x_{2}$, as in the statement of the theorem. Let $\rho=\inf \{f(x)-f(y): f(x)>f(y)\}$. Note that $\rho>0$
because $D$ is finite. If $X_{0}^{1}=x_{1}$ and $X_{0}^{2}=x_{2}$ then $f\left(X_{t}^{1}\right)-f\left(X_{t}^{2}\right) \geq \rho$ for $t<\tau$.
Let $n$ be the smallest index for an eigenvalue $\mu_{n}$ such that $\varphi_{n}\left(x_{1}\right) \neq$ $\varphi_{n}\left(x_{2}\right)$. Such an index must exist because otherwise the eigenfunction expansions would be identical for $p\left(t, x_{1}, y\right)$ and $p\left(t, x_{2}, y\right)$ and, consequently, these functions would be identical; this is not the case since $x_{1} \neq x_{2}$.
By replacing the function $f(x)$ with $f(x) \exp (\alpha f(x))$, for an appropriate $\alpha$, if necessary, we may assume that $S=\sum_{y \in D} f(y) \varphi_{n}(y) \neq 0$. From an eigenfunction expansion similar to (2.1) but listing higher order terms we obtain,

$$
\begin{align*}
& \mathbb{E}\left[f\left(X_{t}^{1}\right) \mid X_{0}^{1}=x_{1}\right]-\mathbb{E}\left[f\left(X_{t}^{2}\right) \mid X_{0}^{2}=x_{2}\right] \\
= & \sum_{y \in D} f(y) p\left(t, x_{1}, y\right)-\sum_{y \in D} f(y) p\left(t, x_{2}, y\right) \\
= & \sum_{y \in D} f(y)\left\{\left[\varphi_{n}\left(x_{1}\right)-\varphi_{n}\left(x_{2}\right)\right] \varphi_{n}(y) e^{-\mu_{n} t}+\tilde{R}\left(t, x_{1}, y\right)+\tilde{R}\left(t, x_{2}, y\right)\right\} \\
= & S\left[\varphi_{n}\left(x_{1}\right)-\varphi_{n}\left(x_{2}\right)\right] e^{-\mu_{n} t}+\hat{R}\left(t, x_{1}, x_{2}, y\right), \tag{2.7}
\end{align*}
$$

where $\hat{R}\left(t, x_{1}, x_{2}, y\right)$ goes to 0 faster than $e^{-\mu_{n} t}$ as $t \rightarrow \infty$. Recalling the definition of $\mu^{\prime}$,

$$
\begin{aligned}
& \mathbb{E}\left[f\left(X_{t}^{1}\right) \mid X_{0}^{1}=x_{1}\right]-\mathbb{E}\left[f\left(X_{t}^{2}\right) \mid X_{0}^{2}=x_{2}\right] \\
& =\mathbb{E}\left[f\left(X_{t}^{1}\right) \mathrm{I}_{\{t<\tau\}} \mid X_{0}^{1}=x_{1}\right]-\mathbb{E}\left[f\left(X_{t}^{2}\right) \mathrm{I}_{\{t<\tau\}} \mid X_{0}^{2}=x_{2}\right] \\
& =\mathbb{E}\left[\left(f\left(X_{t}^{1}\right)-f\left(X_{t}^{2}\right)\right) \mathrm{I}_{\{t<\tau\}} \mid X_{0}^{1}=x_{1}, X_{0}^{2}=x_{2}\right] \\
& \geq \rho \mathbb{E}\left[\mathrm{I}_{\{t<\tau\}} \mid X_{0}^{1}=x_{1}, X_{0}^{2}=x_{2}\right] \\
& \geq \rho c\left(x_{1}, x_{2}\right) e^{-\mu^{\prime}\left(x_{1}, x_{2}\right) t} .
\end{aligned}
$$

Comparing this with (2.7) for large $t$ shows that $\mu^{\prime}\left(x_{1}, x_{2}\right) \geq \mu_{n} \geq \mu_{2}$.
This concludes the proof of the theorem.
Note that Corollary 2.7(ii) follows because $\mu^{*}=\min \mu^{\prime}\left(x_{1}, x_{2}\right) \geq \mu_{2}$. Since we always have $\mu^{*} \leq \mu_{2}$, we see $\mu^{*}=\mu_{2}$ as required.

The following notation will be used for the rest of the section. For distinct $d_{1}, d_{2} \in D$, we will denote the jump rate from $d_{1}$ to $d_{2}$ by $q\left(d_{1}, d_{2}\right)$, i.e.,

$$
q\left(d_{1}, d_{2}\right)=\lim _{s \rightarrow 0} \frac{1}{s} \mathbb{P}\left(X_{t+s}=d_{2} \mid X_{t}=d_{1}\right)
$$

Note that by symmetry we have $q\left(d_{1}, d_{2}\right)=q\left(d_{2}, d_{1}\right)$. Consider any coupling $\left(X^{1}, X^{2}\right)$ for $X$. In a slight abuse of notation we will also use $q$ for the transition rates for the coupling process $\left(X^{1}, X^{2}\right)$ : for distinct pairs $\left(d_{1}, d_{2}\right)$ and $\left(d_{3}, d_{4}\right)$
we set

$$
\begin{aligned}
& q\left(\left(d_{1}, d_{2}\right),\left(d_{3}, d_{4}\right)\right)= \\
& \lim _{s \rightarrow 0} \frac{1}{s} \mathbb{P}\left[\left(X_{t+s}^{1}, X_{t+s}^{2}\right)=\left(d_{3}, d_{4}\right) \mid\left(X_{t}^{1}, X_{t}^{2}\right)=\left(d_{1}, d_{2}\right)\right] .
\end{aligned}
$$

Since the processes $X^{1}$ and $X^{2}$ are Markov and have the same transition probabilities as $X$, for all $d_{1}, d_{2}, d_{3} \in D$ with $d_{1} \neq d_{3}$ we must have

$$
\begin{equation*}
\sum_{d_{4} \in D} q\left(\left(d_{1}, d_{2}\right),\left(d_{3}, d_{4}\right)\right)=q\left(d_{1}, d_{3}\right) . \tag{2.8}
\end{equation*}
$$

For the same reason, if $d_{2} \neq d_{4}$ then

$$
\begin{equation*}
\sum_{d_{3} \in D} q\left(\left(d_{1}, d_{2}\right),\left(d_{3}, d_{4}\right)\right)=q\left(d_{2}, d_{4}\right) . \tag{2.9}
\end{equation*}
$$

We will say that $X^{1}$ and $X^{2}$ make independent jumps from $\left(d_{1}, d_{2}\right)$ if

$$
q\left(\left(d_{1}, d_{2}\right),\left(d_{3}, d_{2}\right)\right)=q\left(d_{1}, d_{3}\right)
$$

and

$$
q\left(\left(d_{1}, d_{2}\right),\left(d_{1}, d_{4}\right)\right)=q\left(d_{2}, d_{4}\right)
$$

for all $d_{3} \neq d_{1}, d_{4} \neq d_{2}$.
Consider a simple example with the state space $D=\{0,1, \ldots, 100\}^{2}$. Suppose that $X$ is a continuous time Markov process on $D$ such that its jumps form the simple random walk reflected on the "boundary" of $D$. Let $X^{1}$ and $X^{2}$ be run as independent copies of $X$ until their coupling time $\tau$. It is no surprise to note that ( $X^{1}, X^{2}$ ) is not efficient, by Corollary 2.7 (i). However, this "independent" coupling is efficient when the state space is ordered. Moreover, a very weak condition ensures efficiency in this case, so that the family of efficient couplings is rather large: it is required only that the coupling maintains the ordering. We state the following result for skip-free chains:

Corollary 2.8 Suppose the state space $D$ is a finite subinterval of the integers $\mathbb{Z}$ and $X$ can jump from $x$ only to $x-1$ or $x+1$, for every $x$. (So $X$ is a finite-state-space generalized birth-death process.) Assume that $\left(X_{t}^{1}, X_{t}^{2}\right)$ almost surely never jumps to $\left(X_{t-}^{2}, X_{t-}^{1}\right)$. Then $\left(X^{1}, X^{2}\right)$ is an efficient coupling. In particular, the coupling is efficient if $X^{1}$ and $X^{2}$ have independent jumps until the coupling time $\tau$.

Proof : The corollary follows from Corollary 2.7 (ii). It suffices to use either $f(x) \equiv x$ or $f(x) \equiv-x$.

The next two examples show that neither of the conditions in parts (i) and (ii) of Theorem 2.6 is necessary.

Example 2.9 We construct a Markov process, and a coupling with $\mu^{\prime}\left(x_{1}, x_{2}\right) \geq$ $\mu_{2}$ for a specific pair of points $x_{1}$ and $x_{2}$, although there is no function $f$ satisfying condition (ii) of Theorem 2.6 for this pair of points. We also show that $\varphi_{2}\left(x_{1}\right) \neq \varphi_{2}\left(x_{2}\right)$, since otherwise this example would not be an improvement on Example 2.3 (where we have $\mu^{\prime}\left(x_{1}, x_{2}\right)>\mu_{2}$ for some points but only for those with $\left.\varphi_{2}\left(x_{1}\right)=\varphi_{2}\left(x_{2}\right)\right)$.

Fix some large $n$ and let the state space of the process be

$$
D=\left\{a_{0}^{1}, a_{1}^{1}, a_{2}^{1}, a_{3}^{1}, a_{0}^{2}, a_{1}^{2}, a_{2}^{2}, a_{3}^{2}, \ldots, a_{0}^{n}, a_{1}^{n}, a_{2}^{n}, a_{3}^{n}\right\} .
$$

Fig. 1 illustrates all possible jumps for the process. We take $q\left(a_{0}^{j}, a_{0}^{j+1}\right)=1$ for $j=1,2, \ldots, n-1$, and $q\left(d_{1}, d_{2}\right)=\tilde{q}$ for all other vertices $d_{1}$ and $d_{2}$ connected by a wedge in the graph. We will choose a value for $\tilde{q}$ later in the example.
[Figure 1 about here.]
Consider a coupling with the jump rates

$$
q\left(\left(a_{m}^{j}, a_{m}^{k}\right),\left(a_{l}^{j}, a_{l}^{k}\right)\right)=q\left(a_{m}^{j}, a_{l}^{j}\right)
$$

for all $j, k, l$ and $m$, with $j \neq k$. Suppose that

$$
\begin{array}{rll}
q\left(\left(a_{0}^{j}, a_{0}^{k}\right),\left(a_{0}^{j+1}, a_{0}^{k}\right)\right) & =1, & \\
q \neq k, j=1, \ldots, n-1, \\
q\left(\left(a_{0}^{j}, a_{0}^{k}\right),\left(a_{0}^{j-1}, a_{0}^{k}\right)\right) & =1, & \\
q \neq k, j=2, \ldots, n, \\
q\left(\left(a_{0}^{j}, a_{0}^{k}\right),\left(a_{0}^{j}, a_{0}^{k+1}\right)\right) & =1, & \\
q\left(\left(a_{0}^{j}, a_{0}^{k}\right),\left(a_{0}^{j}, a_{0}^{k-1}\right)\right) & =1, & j \neq k, k=1, \ldots, n-1, \\
\end{array}
$$

We require that $X^{1}$ and $X^{2}$ make independent jumps from all other points $\left(d_{1}, d_{2}\right)$.

If $\left(X_{0}^{1}, X_{0}^{2}\right)=\left(a_{l}^{j}, a_{m}^{j}\right)$ for some $l \neq m$ then $\mathbb{P}\left(\left(X_{s}^{1}, X_{s}^{2}\right)=\left(a_{m}^{j}, a_{l}^{j}\right)\right)>0$ for every $s>0$, so there does not exist a function $f: D \rightarrow \mathbb{R}$ such that $f\left(X_{t}^{1}\right)-f\left(X_{t}^{2}\right)$ is strictly positive until the coupling time. However, we will argue that $\mu^{\prime}\left(a_{l}^{j}, a_{m}^{j}\right)=\mu_{2}$ for some $j$ and $l \neq m$, if $\tilde{q}$ is large enough.

If both $X_{t}^{1}$ and $X_{t}^{2}$ lie on the "spine" $\left\{a_{0}^{1}, a_{0}^{2}, \ldots, a_{0}^{n}\right\}$ at some time $s$, then from this time on, they will make excursions into side alleys of the form $\left\{a_{0}^{j}, a_{1}^{j}, a_{2}^{j}, a_{3}^{j}\right\}$ at the same time and they will return from those excursions to the spine at the same time. It follows that if $\left(X_{s}^{1}, X_{s}^{2}\right)=\left(a_{0}^{j}, a_{0}^{k}\right)$ with $j<k$ then $X_{t}^{1}$ will lie to the left of $X_{t}^{2}$ for all $t \in[s, \tau)$. We let $f\left(a_{m}^{j}\right)=j$ and use Theorem 2.6 (ii) to see that $\mu^{\prime}\left(a_{0}^{j}, a_{0}^{k}\right) \geq \mu_{2}$ for $j \neq k$.

We will estimate $\mu^{\prime}\left(a_{0}^{j}, a_{0}^{k}\right)$. First suppose that the processes $X^{1}$ and $X^{2}$ start from distinct points on the spine. Then their evolution may be described as that of two independent copies of $X$ along the spine except that they may make simultaneous excursions into the side alleys of the form $\left\{a_{0}^{j}, a_{1}^{j}, a_{2}^{j}, a_{3}^{j}\right\}$. Those side excursions can only delay the coupling time $\tau$ for $X^{1}$ and $X^{2}$ so $\tau$ is
stochastically minorized by the coupling time for a pair of independent random walks on the spine, reflected at the endpoints of the spine. Hence, by Corollary 2.8,

$$
\mathbb{P}\left[\tau>t \mid\left(X_{0}^{1}, X_{0}^{2}\right)=\left(a_{0}^{j}, a_{0}^{k}\right)\right] \geq c(j, k) e^{-\hat{\mu} t}
$$

for all $j \neq k$, where $\hat{\mu}$ is the second eigenvalue for the process restricted to the spine. Note that $\hat{\mu}$ does not depend on $\tilde{q}$.

Next suppose that $\left(X^{1}, X^{2}\right)$ starts from $\left(a_{l}^{j}, a_{m}^{j}\right)$ for some $j$ and $l \neq m$. If we choose sufficiently large $\tilde{q}$ then the processes $X^{1}$ and $X^{2}$ will rapidly and independently jump along the edges connecting the elements of the family $\left\{a_{0}^{j}, a_{1}^{j}, a_{2}^{j}, a_{3}^{j}\right\}$. It is clear that they will rapidly couple, before leaving this set. As a consequence, for sufficiently large $\tilde{q}$ one can choose $c_{1}$ such that for all $j, l$ and $m$,

$$
\mathbb{P}\left[\tau>t \mid\left(X_{0}^{1}, X_{0}^{2}\right)=\left(a_{l}^{j}, a_{m}^{j}\right)\right] \leq c_{1} e^{-2 \hat{\mu} t}
$$

It follows that for large $t$,

$$
\begin{aligned}
& \mathbb{P}\left[\tau>t \mid\left(X_{0}^{1}, X_{0}^{2}\right)=\left(a_{l}^{j}, a_{m}^{j}\right)\right] \leq c_{1} e^{-2 \hat{\mu} t} \leq \min _{r \neq k} c(r, k) e^{-\hat{\mu} t} \leq \\
& \quad \leq \min _{r \neq k} \mathbb{P}\left[\tau>t \mid\left(X_{0}^{1}, X_{0}^{2}\right)=\left(a_{0}^{r}, a_{0}^{k}\right)\right]
\end{aligned}
$$

Hence, $\min _{j} \min _{l \neq m} \mu^{\prime}\left(a_{l}^{j}, a_{m}^{j}\right) \geq \max _{r \neq k} \mu^{\prime}\left(a_{0}^{r}, a_{0}^{k}\right) \geq \mu_{2}$.
Finally we will show that $g\left(a_{l}^{j}, \cdot\right)$ and $g\left(a_{m}^{j}, \cdot\right)$ are not identical for some $j$ and $l \neq m$. Suppose that the converse holds. Fix some $j$ and note that if $g\left(a_{l}^{j}, \cdot\right)=g\left(a_{m}^{j}, \cdot\right)$ for all $l$ and $m$ then $g\left(a_{1}^{j}, \cdot\right)$ is an average in the first argument of the other $g\left(a_{m}^{j}, \cdot\right)$, and indeed $g$ is harmonic in its first argument at $a_{1}^{j}$. As a function of its first argument, $g$ is an eigenfunction corresponding to $\mu_{2}$; therefore we must have $g\left(a_{l}^{j}, \cdot\right)=0$ for all $l$. If this is true for all $j$ then $g$ is identically equal to 0 , which is a contradiction.

We conclude that for some $j$ and $l \neq m$ we have $\mu^{\prime}\left(a_{l}^{j}, a_{m}^{j}\right) \geq \mu_{2}$ and $g\left(a_{l}^{j}, \cdot\right), g\left(a_{m}^{j}, \cdot\right)$ not identical. However there is no function $f$ which would satisfy condition (ii) of Theorem 2.6 for $a_{l}^{j}$ and $a_{m}^{j}$.

Example 2.10 We present a Markov process, a coupling and a pair of points $x_{1}, x_{2}$ with $\mu^{\prime}\left(x_{1}, x_{2}\right)<\mu_{2}$ although neither $x_{1}$ nor $x_{2}$ satisfies the transposition property (i) of Theorem 2.6 (i). Let $D=\{-100,-99, \ldots, 100\} \times\{0,1, \ldots, 50\}$ be a bounded portion of $\mathbb{Z}^{2}$. Let the process $X$ be able to jump only to its nearest neighbors in $D$, and let the jump rate be equal to 1 for every pair of neighbors. We consider a coupling $\left(X^{1}, X^{2}\right)$ with the jump rates

$$
q(((-j, k),(j, k)),((-l, m),(l, m))=1
$$

for $j \geq 1,|j-l|+|k-m|=1$. The processes $X^{1}$ and $X^{2}$ have independent jumps from all points which are not symmetric with respect to the vertical axis $K$.

Suppose that $X^{1}$ and $X^{2}$ start from distinct points $x_{1}$ and $x_{2}$ which are not symmetric with respect to $K$. Then these processes are independent copies of $X$ until the time $T$ when $X^{1}$ hits $X^{2}$ or the symmetric image of $X^{2}$ with respect to $K$. The second eigenvalue $\mu_{2}$ for $X$ is the same as the second eigenvalue for its first component. The following few assertions are quite clear but we will not go into a detailed proof of them as it would take too much space. We have $\mathbb{P}(T>t) \geq c e^{-\mu t}$ for some $c$ and $\mu<\mu_{2}$. This is justified by comparing $T$ to $U$, where $U$ is the time when $X$ hits $K$, and by the fact that $\mathbb{P}(U>t) \leq c e^{-\mu_{2} t}$. We obtain $\mu^{\prime}\left(x_{1}, x_{2}\right) \leq \mu<\mu_{2}$.

Let $x_{3}$ be the point symmetric to $x_{1}$ with respect to $K$. Note that $\left(x_{1}, x_{3}\right)$ is accessible from ( $x_{1}, x_{2}$ ) for this particular coupling but neither ( $x_{1}, x_{2}$ ) nor $\left(x_{2}, x_{1}\right)$ is accessible from $\left(x_{1}, x_{3}\right)$. Hence neither $x_{1}$ nor $x_{2}$ has the transposition property.

Our next example is a continuous version of [21, Exercise 5.9].
Example 2.11 (Ehrenfest model or random walk on hypercube) Consider two urns with $n$ marked balls distributed among them. At every arrival time for a Poisson process, a ball is randomly chosen from among all balls in both urns and moved to the other urn.

A formal description of the model is the following. The state space $D$ for our process is the set of all binary sequences $\left(i_{1}, i_{2}, \ldots, i_{n}\right)$ of length $n$, so each $i_{j}$ is equal to 0 or to 1 . Let $U_{k}, k \geq 1$, be independent exponential (mean 1) random variables and set $T_{k}=U_{1}+\cdots+U_{k}$. Consider random variables $N_{k}$ which are independent of each other and of the $T_{k}$ 's, and which are uniformly distributed over the fixed range $\{1,2, \ldots, n\}$. Finally, let $\left\{J_{k}\right\}$ be a sequence of random variables, independent of each other, of the $T_{k}$ 's, and of the $N_{k}$ 's, and such that $\mathbb{P}\left(J_{k}=0\right)=\mathbb{P}\left(J_{k}=1\right)=1 / 2$. We define the process $X$ on $D$ by prescribing its initial value $X_{0}=\left(i_{1}, i_{2}, \ldots, i_{n}\right)$ and by specifying its jumps; $X$ jumps at times $T_{k}$ (and only at these times), the jump at time $T_{k}$ taking the process from $X_{T_{k}-}=\left(j_{1}, j_{2}, \ldots, j_{N_{k}}, \ldots, j_{n}\right)$ to $X_{T_{k}}=\left(j_{1}, j_{2}, \ldots, J_{k}, \ldots, j_{n}\right)$. If the process jumped instead to $\left(j_{1}, j_{2}, \ldots, 1-j_{N_{k}}, \ldots, j_{n}\right)$ at time $T_{k}$ then we would have obtained a model directly corresponding to the informal "urn" representation. However, the two processes $X$, corresponding to two kinds of jumps, can be transformed into each other by speeding up or slowing down the clock for $X$.

We construct a coupling ( $X^{1}, X^{2}$ ) by using just one family of random variables $\left\{T_{k}, N_{k}, J_{k}\right\}_{k \geq 1}$ for both processes $X^{1}$ and $X^{2}$. Specifically, the transition probabilities for $\left(X^{1}, X^{2}\right)$ are specified by the requirement that the process ( $X^{1}, X^{2}$ ) jumps at times $T_{k}$ (and only at these times) from

$$
\left(X_{T_{k}-}^{1}, X_{T_{k}-}^{2}\right)=\left(\left(j_{1}^{1}, j_{2}^{1}, \ldots, j_{N_{k}}^{1}, \ldots, j_{n}^{1}\right),\left(j_{1}^{2}, j_{2}^{2}, \ldots, j_{N_{k}}^{2}, \ldots, j_{n}^{2}\right)\right)
$$

to

$$
\left(X_{T_{k}}^{1}, X_{T_{k}}^{2}\right)=\left(\left(j_{1}^{1}, j_{2}^{1}, \ldots, J_{k}, \ldots, j_{n}^{1}\right),\left(j_{1}^{2}, j_{2}^{2}, \ldots, J_{k}, \ldots, j_{n}^{2}\right)\right)
$$

It is immediate that $\left(X^{1}, X^{2}\right)$ and also $X^{1}$ and $X^{2}$ are all Markov processes; moreover the latter two have the same transition probabilities as $X$.

Suppose that $X_{0}^{1}=\left(j_{1}^{1}, j_{2}^{1}, \ldots, j_{n}^{1}\right)$. Let $f\left(i_{1}, i_{2}, \ldots, i_{n}\right)$ be the number of $k$ such that $i_{k}=j_{k}^{1}$. It is elementary to check that if $X_{0}^{1} \neq X_{0}^{2}$ then $f\left(X_{t}^{1}\right)-f\left(X_{t}^{2}\right)$ stays strictly positive for all $t<\tau$. By Corollary 2.7 (ii), our coupling is efficient.

Recall $\mu^{*}$ defined before Lemma 2.1 in equation (2.2). In order to estimate $\mu^{*}$ we consider the worst case scenario, i.e., that initially no components of $X_{0}^{1}$ and $X_{0}^{2}$ are equal. For a fixed $k$, the waiting time for the $k$-th components of $X^{1}$ and $X^{2}$ to meet is exponential with mean $n$. Once these components meet, they will be equal to each other forever, although they will not be constant. It is an easy consequence of the theory of Poisson point processes with independent marks that the waiting times for different components are independent. The probability that a specified pair of components have not merged by time $t$ is equal to $e^{-t / n}$ so the probability that there exists at least one such pair is equal to $1-\left(1-e^{-t / n}\right)^{n}$ which is between (1/2)ne $e^{-t / n}$ and $n e^{-t / n}$ for large $t$. We conclude that $\mu^{*}=1 / n$, and since our coupling is efficient, we see that the spectral gap is also equal to $1 / n$.

Consider now the asymptotics when $n \rightarrow \infty$. In this example, the mean time to coupling is not of order $1 / \mu^{*}=n$. The expected waiting time between the $k$-th coupling of a pair of components of $X^{1}$ and $X^{2}$ and the ( $k+1$ )-st coupling is $n /(n-k)$, so that the expected time until all components are coupled is of order $\sum_{k=0}^{n-1} n /(n-k) \approx n \log n$, which is larger than $1 / \mu^{*}$ by a factor of $\log n$.

Example 2.12 We construct a symmetric Markov process $X_{t}$ for which there are no efficient Markovian couplings. The state space consists of 11 points,

$$
D=\left\{a_{1}, a_{2}, a_{3}, b_{1}, b_{2}, b_{3}, c_{1}, c_{2}, \ldots, c_{5}\right\} .
$$

The edges in Fig. 2 show all possible jumps for the process; in other words, points $d_{1}, d_{2} \in D$ are connected by an edge in Fig. 2 if and only if $q\left(d_{1}, d_{2}\right)>0$. The jump rates are $q\left(a_{1}, a_{2}\right)=1, q\left(a_{2}, a_{3}\right)=2$, and $q\left(d_{1}, d_{2}\right)=4$ for all other edges in Fig. 2.
[Figure 2 about here.]
Consider any coupling $\left(X^{1}, X^{2}\right)$ for this process. Suppose that $d_{1}, d_{2} \in$ $D \backslash\left\{a_{2}\right\}$ are such that $q\left(d_{1}, d_{2}\right)>0$, and, therefore, $q\left(d_{1}, d_{2}\right)=4$. Since $q\left(a_{2}, a_{1}\right)=1$ and $q\left(a_{2}, a_{3}\right)=2$, the identities (2.8) and (2.9) imply that $q\left(\left(d_{1}, a_{2}\right),\left(d_{2}, a_{1}\right)\right) \leq 1$ and $q\left(\left(d_{1}, a_{2}\right),\left(d_{2}, a_{3}\right)\right) \leq 2$. By another application of equations (2.8)-(2.9) we deduce from these inequalities that $q\left(\left(d_{1}, a_{2}\right),\left(d_{2}, a_{2}\right)\right)>$ 0 , since $q\left(a_{2}, a_{1}\right)=1, q\left(a_{2}, a_{3}\right)=2$, and the sum of all $q\left(\left(d_{1}, a_{2}\right),\left(d_{2}, a_{i}\right)\right)$ equals 4.

We will describe an evolution of the process $\left(X^{1}, X^{2}\right)$ before the coupling time which may happen with positive probability, no matter what coupling is used.

First we consider the case when one of the processes $X^{1}$ or $X^{2}$ starts from $a_{2}$. Without loss of generality suppose that $X_{0}^{1}=a_{2}$ and $X_{0}^{2}=d_{0} \neq a_{2}$. It follows from what we have just proved that $q\left(\left(d_{0}, a_{2}\right),\left(d_{2}, a_{2}\right)\right)>0$, where $d_{2}$ is any "neighbor" of $d_{0}$ in the "rectangle" $D \backslash\left\{a_{2}\right\}$. Hence, there is a positive
probability that $X^{2}$ will move through a sequence of points in $D \backslash\left\{a_{2}\right\}$ and reach $b_{2}$ at time $t_{1}$, before the time when $X^{1}$ leaves $a_{2}$.

We note that for any coupling $\left(X^{1}, X^{2}\right)$ and any position $\left(d_{1}, d_{2}\right)$ of this process at time $s$, if $X$ can jump with a positive probability from $d_{1}$ to $d_{3}$, then there is a positive probability that the first jump of $X^{1}$ after time $s$ will take it to $d_{3}$, and, moreover, the jump of $X^{1}$ will occur before or at the same time when $X^{2}$ makes its first jump after time $s$. Thus, there is a positive probability that $X^{1}$ will jump to $a_{1}$ at some time $t_{2}>t_{1}$, but $X^{2}$ will not jump within interval $\left(t_{1}, t_{2}\right]$ at all or it will have only one jump, at time $t_{2}$. Hence, there is positive probability that for some $t_{2}$ we have $X_{t_{2}}^{1}=a_{1}$ and $X_{t_{2}}^{2}=b_{k}$ for some $k$.

The same argument shows that, with a positive probability, $X^{1}$ will jump after time $t_{2}$ from $a_{1}$ to $c_{1}, c_{2}$ and $c_{3}$, at times $t_{3}, t_{4}$ and $t_{5}$, while $X_{t}^{2}$ will have at most 3 jumps for $t \in\left(t_{2}, t_{5}\right.$ ], and, moreover, all of the jumps of $X^{2}$ will occur at times $t_{3}, t_{4}$ and/or $t_{5}$. It may happen that $X^{2}$ hits $a_{1}$ or $a_{3}$ at some time $t_{6} \in\left(t_{2}, t_{5}\right]$. If this is the case then $X_{t_{6}}^{2} \in\left\{a_{1}, a_{3}\right\}$ while $X_{t_{6}}^{1} \in\left\{c_{1}, c_{2}, c_{3}\right\}$.

If $X^{2}$ does not hit $a_{1}$ or $a_{3}$ before or at the time $t_{5}$ then $X_{t_{5}}^{2} \in\left\{b_{1}, b_{2}, b_{3}\right\}$ and $X_{t_{5}}^{1}=c_{3}$. In this case, a possible evolution of the process after time $t_{5}$ is that $X^{2}$ will make one or two jumps that will take this process to either $a_{1}$ or $a_{3}$, whichever is closer to $X_{t_{5}}^{2}$. Let $t_{7}$ be the time when $X^{2}$ hits $a_{1}$ or $a_{3}$. With positive probability, $X^{1}$ will make at most two jumps during the time interval $\left(t_{5}, t_{7}\right]$ and so we will have $X_{t_{7}}^{1} \in\left\{c_{1}, c_{2}, c_{3}, c_{4}, c_{5}\right\}$. We see that with positive probability, there is a finite time $t_{8}$, equal either to $t_{6}$ or $t_{7}$, such that $X_{t_{8}}^{2} \in\left\{a_{1}, a_{3}\right\}$ and $X_{t_{8}}^{1} \in\left\{c_{1}, c_{2}, c_{3}, c_{4}, c_{5}\right\}$.

With positive probability, the process $X^{2}$ can jump to $a_{2}$ at time $t_{9}>t_{8}$, while $X^{1}$ will not jump during $\left(t_{8}, t_{9}\right)$. We will have $X_{t_{9}}^{2}=a_{2}$ and $X_{t_{9}}^{1} \in$ $D \backslash\left\{a_{2}\right\}$. Then $X_{t}^{1}$ may reach $d_{0}$ (the initial position of $X^{2}$ ) at time $t_{10}>t_{9}$, before $X^{2}$ leaves $a_{2}$, by the argument presented earlier in the proof. We have shown that with a positive probability, the coupling ( $X^{1}, X^{2}$ ) may go from $\left(a_{2}, d_{0}\right)$ to ( $d_{0}, a_{2}$ ), before the coupling time.

Next we consider an arbitrary initial position $\left(d_{1}, d_{2}\right)$ for $\left(X^{1}, X^{2}\right)$ with $d_{1} \neq d_{2}$. We assume that $d_{1}, d_{2} \in D \backslash\left\{a_{2}\right\}$ because the other case has been discussed in the first part of the proof. With positive probability, the process $X^{1}$ may keep jumping in the clockwise direction around $D \backslash\left\{a_{2}\right\}$, while $X^{2}$ will jump only at the same times when $X^{1}$ jumps. The same remark applies to jumps of $X^{1}$ in the counter-clockwise direction; it also applies when we reverse the roles of $X^{1}$ and $X^{2}$. An elementary argument now shows that for any starting position, one or both processes may hit $\left\{a_{1}, a_{3}\right\}$ before the coupling time. First suppose that only one of them hits the set $\left\{a_{1}, a_{3}\right\}$. Then this process may jump to $a_{2}$ before or at the same time when the other process jumps. At this instant, one of the processes will be at $a_{2}$ while the other will be in $D \backslash\left\{a_{2}\right\}$. The other possibility is that both processes hit $\left\{a_{1}, a_{3}\right\}$ at the same instant, before the coupling time. Then, since $q\left(a_{3}, a_{2}\right)>q\left(a_{1}, a_{2}\right)$, and using equations (2.8)-(2.9), there is a positive chance that the process at the point $a_{3}$ will jump to $a_{2}$ before the other process jumps to $a_{2}$. Hence, just as in the first part of the proof, we will have one of the processes at $a_{2}$ and the other
at a point of $D \backslash\left\{a_{2}\right\}$. The process which happens to be in $D \backslash\left\{a_{2}\right\}$ may go to any other point of $D \backslash\left\{a_{2}\right\}$ strictly before the other process leaves $a_{2}$.

We have proved that for every $\left(d_{1}, d_{2}\right)$, and every $d_{3} \in D \backslash\left\{a_{2}\right\}$, either $\left(a_{2}, d_{3}\right)$ is accessible from $\left(d_{1}, d_{2}\right)$ or $\left(d_{3}, a_{2}\right)$ is accessible from the same point. By the first part of the proof, $\left(a_{2}, d_{3}\right)$ is accessible from $\left(d_{3}, a_{2}\right)$ and vice versa, so by transitivity, both $\left(a_{2}, d_{3}\right)$ and $\left(d_{3}, a_{2}\right)$ are accessible from $\left(d_{1}, d_{2}\right)$. In particular, if $\left(d_{1}, d_{2}\right)$ is accessible from $\left(a_{2}, d_{3}\right)$ then both $\left(a_{2}, d_{3}\right)$ and $\left(d_{3}, a_{2}\right)$ are accessible from $\left(d_{1}, d_{2}\right)$. In other words, every coupling $\left(X_{t}^{1}, X_{t}^{2}\right)$ has the transposition property relative to $a_{2}$. We conclude that no coupling is efficient for this Markov chain, by Theorem 2.6 (i). Since for every $\left(d_{1}, d_{2}\right)$ with $d_{1} \neq d_{2}$, $\left(a_{2}, d_{3}\right)$ is accessible from $\left(d_{1}, d_{2}\right)$, it easily follows that $\mu^{\prime}\left(d_{1}, d_{2}\right)<\mu_{2}$ for every pair of distinct points $d_{1}, d_{2} \in D$.

Remark : We list three open problems inspired by Example 2.12.

1. Give necessary and sufficient conditions in terms of $q(x, y)$ for the existence of an efficient coupling for a continuous-time symmetric Markov process with finite state space.
2. A quantitative version of Problem 1 is the following. Let $\bar{\mu}=\sup \mu^{*}$, where the supremum is taken over all couplings for a given Markov process. Can we have $\bar{\mu}<\mu_{2}$ ? If so, how can one calculate $\mu_{2}-\bar{\mu}$ starting from $q(x, y)$ ?
3. If the state space $D$ has only 3 elements then there exists an efficient coupling (we omit an easy proof). Does an efficient coupling exist for every Markov process whose state space is a loop: namely, the state space is $\{0,1, \ldots, n\}$ for some $n$ and $q(j, k)>0$ if and only if $|j-k|=1$ or $n$ ? We conjecture that the answer is no if $n \geq 3$ (Mountford and Cranston [37] have now produced a counterexample to the question, as well as discussing many other interesting related questions, so our conjecture is correct.).

As with many coupling problems, it may help to think of this problem in terms of a game. Suppose that player $\mathcal{K}$ begins at state $x_{1}$ and player $\mathcal{W}$ begins at a different state $x_{2}$. $\mathcal{W}$ can make various moves according to possibilities admitted by the $Q$-matrix of the Markov process in question, and weighted by the off-diagonal terms of the $Q$-matrix. $\mathcal{K}$ must predeclare his moves in terms of either what $\mathcal{W}$ might do or electing to move independently, dividing weights between various possibilities so as $(a)$ to add up to the weights prescribed by the off-diagonal terms of the $Q$-matrix, and (b) such that a move predeclared in terms of a $\mathcal{W}$-move must have weight no greater than that of the $\mathcal{W}$-move. (Thus a viable strategy for $\mathcal{K}$ corresponds to an admissible coupling.) Player $\mathcal{W}$ wins the game if he can choose a sequence of $\mathcal{W}$ - and $\mathcal{K}$-moves compatible with $\mathcal{K}$ 's predeclared options, and such that at some stage the positions of $\mathcal{W}$ and $\mathcal{K}$ interchange from what they were at a previous occasion. By Theorem 2.6(i), there is no efficient coupling exactly when for some initial point $x_{1}$ for $\mathcal{K}$ it is the case that player $\mathcal{W}$ can win the game whatever viable strategy is chosen by $\mathcal{K}$.

## 3 Reflected Brownian motion in a triangle

We will illustrate the concept of efficiency for Markovian couplings for continuous processes by a detailed analysis of two couplings for reflected Brownian motion in a triangle. We have chosen this example as the role of "partial ordering" is clear in this case. Moreover, our methods developed for this example have laid a foundation for some results on the "hot spots" conjecture of J. Rauch $[3,6,5]$. We remark in passing that not much can be said about the exact values of eigenvalues for reflected Brownian motion in a triangle: see [38, 39] for the equilateral case.

We will discuss "synchronous" and "mirror" couplings. The synchronous couplings have been studied, for example, by Cranston and Le Jan [15, 16]. The mirror couplings seem to be a more effective tool than the synchronous couplings for estimating the spectral gap [48]. We will start with synchronous couplings. Our results are not as complete in this case as in the case of mirror couplings and for this reason this part of the presentation is less technical.

First we will give a construction of the synchronous coupling. Note that the notation is changed from the last section in the following respect: previously $X^{1}$ and $X^{2}$ denoted copies of $X$, while in this section they will stand for the components of the two-dimensional process $X$. The coupling process will consist of $X$ and an identically distributed (but not independent!) copy $Y$.

Let $\tilde{X}=\left(\tilde{X}^{1}, \tilde{X}^{2}\right)$ be a 2-dimensional Brownian motion with $\tilde{X}_{0}=\left(x^{1}, x^{2}\right)$ where $x^{2}>0$. Let $L_{t}^{X}=-\left(0 \wedge \min _{s \leq t} \tilde{X}_{s}^{2}\right)$. Then the Skorokhod Lemma [28, Lemma 3.6.14] implies that the process $X$, defined by $X_{t}=\left(\tilde{X}_{t}^{1}, \tilde{X}_{t}^{2}+L_{t}^{X}\right)$, is a reflected Brownian motion in the upper half-plane. Let $\tilde{Y}_{t}=\left(\tilde{Y}_{t}^{1}, \tilde{Y}_{t}^{2}\right)=\left(\tilde{X}_{t}^{1}+\right.$ $\left.\left(y^{1}-x^{1}\right), \tilde{X}_{t}^{2}+\left(y^{2}-x^{2}\right)\right)$, where $y_{2}>0$. Then $\tilde{Y}$ is a Brownian motion starting from $\left(y^{1}, y^{2}\right)$. Arguing as before, $L_{t}^{Y}=-\left(0 \wedge \min _{s \leq t} \tilde{Y}_{s}^{2}\right)$ can be used to define a reflected Brownian motion $Y$ by means of the formula $Y_{t}=\left(\tilde{Y}_{t}^{1}, \tilde{Y}_{t}^{2}+L_{t}^{Y}\right)$. Let $K_{t}$ be the straight line passing through the two planar points $X_{t}$ and $Y_{t}$ and let $\angle K_{t} \in(-\pi / 2, \pi / 2]$ be the angle between $K_{t}$ and the horizontal axis. (If $X_{t}=Y_{t}$ then let $K_{t}$ be the horizontal line passing through $X_{t}$.) If $\angle K_{0}=\pi / 2$ then $K_{t}$ will stay perpendicular to the horizontal axis until $X_{t}$ and $Y_{t}$ meet at some time $u$ and then $X_{t}=Y_{t}$ for $t>u$. If $\angle K_{0} \neq \pi / 2$ then $\angle K_{t}$ will converge in a monotone way to 0 as $t \rightarrow \infty$ and, moreover, $\angle K_{t}$ will be constantly equal to 0 after both $X_{t}$ and $Y_{t}$ hit the horizontal axis. The pair $(X, Y)$ of reflected Brownian motions will be called a synchronous coupling.

The above construction generalizes in a straightforward way to polygonal domains $D$ : given a pair of starting points $x, y \in \bar{D}$ we can define a pair of reflected Brownian motions $(X, Y)$ in $D$ with $\left(X_{0}, Y_{0}\right)=(x, y)$, and such that $X-Y$ remains constant on every interval during which both processes stay in the interior of the domain. It should be noted that neither of $X$ and $Y$ can hit any vertices of $\partial D$; this follows either by the results of Varadhan and Williams [47] or indeed by viewing $X, Y$ in polar coordinates centered at each of the finitely many polygonal vertices. Finally it is not hard to prove (for example using Brownian excursion theory based on excursions from the side visited immediately prior to coupling) that with positive probability there will
be $u$ such that $X_{u}=Y_{u}$ if and only if $K_{0}$ is perpendicular to one of the sides of the polygon $\partial D$ or if $\partial D$ contains a perpendicular pair of line segments. If such a $u$ exists then $X_{t}=Y_{t}$ for all $t \geq u$.

Chen explains how the spectral gap can be estimated using couplings [10, Theorem 6.2]. We will discuss two concrete implementations of Chen's theorem in the case of reflecting Brownian motion in a triangle.

Let $p(t, x, y)$ be the transition densities for reflecting Brownian motion in $D$. We have

$$
\begin{equation*}
p(t, x, y)=\sum_{n=1}^{\infty} \varphi_{n}(x) \varphi_{n}(y) e^{-\mu_{n} t} \tag{3.1}
\end{equation*}
$$

where $\mu_{n}$ is the $n$-th eigenvalue for the Laplacian in $D$ with Neumann boundary conditions and $\varphi_{n}$ is the corresponding eigenfunction. Recall that $\mu_{1}=0$ and $\varphi_{1}$ is a constant function. It is a classical result ("Mercer's Theorem") that

$$
\begin{equation*}
p(t, x, y)=c_{1}+\varphi_{2}(x) \varphi_{2}(y) e^{-\mu_{2} t}+R(t, x, y) \tag{3.2}
\end{equation*}
$$

where $R(t, x, y)$ converges to 0 faster than $e^{-\mu_{2} t}$ as $t \rightarrow \infty$ uniformly in $x$ and $y$ (see [3, Proposition 2.1] for a recent proof). To be more precise, the function $\varphi_{2}$ in (3.2) is an eigenfunction corresponding to $\mu_{2}$ but in the case of eigenvalue multiplicity it may be a linear combination $g(x, y)=\sum_{\varphi} \varphi(x) \varphi(y)$ of several eigenfunctions in (3.1) corresponding to $\mu_{2}$, analogous to $g$ in equation (1.1). However for the sake of simplicity, and because the generic case will not exhibit such multiplicity, we consider only the case $g(x, y)=\varphi_{2}(x) \varphi_{2}(y)$ in the following. Suppose that one can prove that for some $\mu \geq 0$ and $x, y \in D$,

$$
\begin{equation*}
\mathbb{E}\left(\left|X_{t}-Y_{t}\right| \mid X_{0}=x, Y_{0}=y\right) \leq c(x, y) e^{-\mu t}, \quad t \geq 0 \tag{3.3}
\end{equation*}
$$

(We choose to consider $\mathbb{E}\left[\left|X_{t}-Y_{t}\right| \mid X_{0}=x, Y_{0}=y\right]$ rather than $\mathbb{P}\left[\tau>t \mid X_{0}=\right.$ $\left.x, Y_{0}=y\right]$ because typically $\mathbb{P}\left[\tau>t \mid X_{0}=x, Y_{0}=y\right]=0$ for synchronous couplings: see [16].) One would expect that $\mu$ is then a lower bound for $\mu_{2}$. We examine this assertion in the next example and lemma.

Example 3.1 Consider the long rectangle $D=[0,1] \times[0,100]$ and let $x=$ $(1 / 4,1), y=(3 / 4,1)$. For these $x$ and $y$, and the synchronous coupling, (3.3) will hold with $\mu=\pi$. This follows from the fact that the line $K_{t}$ will always stay parallel to the horizontal axis and so we are effectively dealing with a 1dimensional Neumann problem on the interval $[0,1]$ for which $\pi$ is the second eigenvalue. The second eigenvalue for the Laplacian in $D$ is the same as for the interval $[0,100]$, i.e., $\mu_{2}=\pi / 100$. Hence, (3.3) may hold for some $\mu>\mu_{2}$ and some $x, y \in D$.

Lemma 3.2 If $\varphi_{2}(x) \neq \varphi_{2}(y)$ and (3.3) holds, then $\mu \leq \mu_{2}$.

Proof : Consider $x, y \in D$ such that $\varphi_{2}(x)>\varphi_{2}(y)$. Suppose that $X_{0}=x$ and $Y_{0}=y$. The function $\varphi_{2}$ is not identically equal to 0 so,

$$
\int_{D} \exp \left(c_{1} z_{1}+c_{2} z_{2}\right) \varphi_{2}\left(z_{1}, z_{2}\right) d z_{1} d z_{2}>0
$$

for some $c_{1}, c_{2} \in \mathbb{R} \backslash\{0\}$. Since $D$ is bounded there exists $c_{3}>0$ such that $c_{3}^{-1}$ is a Lipschitz constant for $\exp \left(c_{1} x^{1}+c_{2} x^{2}\right)$, and so

$$
\mathbb{E}\left[\left|X_{t}-Y_{t}\right|\right] \geq c_{3}\left[\mathbb{E}\left[\exp \left(c_{1} X_{t}^{1}+c_{2} X_{t}^{2}\right)\right]-\mathbb{E}\left[\exp \left(c_{1} Y_{t}^{1}+c_{2} Y_{t}^{2}\right)\right]\right]
$$

Then, by equation (3.2),

$$
\begin{aligned}
& \mathbb{E}\left(\left|X_{t}-Y_{t}\right| \mid X_{0}=x, Y_{0}=y\right) \\
\geq & c_{3}\left[\mathbb{E} \exp \left(c_{1} X_{t}^{1}+c_{2} X_{t}^{2}\right)-\mathbb{E} \exp \left(c_{1} Y_{t}^{1}+c_{2} Y_{t}^{2}\right)\right] \\
= & c_{3} \int_{D} \exp \left(c_{1} z_{1}+c_{2} z_{2}\right) p(t, x, z) d z-c_{3} \int_{D} \exp \left(c_{1} z_{1}+c_{2} z_{2}\right) p(t, y, z) d z \\
= & c_{3} \int_{D}\left[\varphi_{2}(x)-\varphi_{2}(y)\right] \exp \left(c_{1} z_{1}+c_{2} z_{2}\right) \varphi_{2}(z) e^{-\mu_{2} t} d z+R(t, x, y) \\
= & c_{4}(x, y) e^{-\mu_{2} t}+R(t, x, y),
\end{aligned}
$$

where $c_{4}(x, y)>0$. (Note that here $R(t, x, y)$ is actually the integrated sum of two terms of the form of $R(t, x, y)$ in equation (3.2).) Since this estimate and inequality (3.3) hold for arbitrarily large $t$, we see that $\mu \leq \mu_{2}$.

In view of Example 3.1 and Lemma 3.2, we propose the following definition.
Definition 3.3 We will call a synchronous coupling $(X, Y)$ of reflected Brownian motions in $D$ efficient if for some $x$ and $y$ with $\varphi_{2}(x) \neq \varphi_{2}(y)$, the estimate (3.3) holds with $\mu=\mu_{2}$.

Note that the above definition of efficiency is different from that given in Section 2 for Markov processes with finite state space. We require in Definition 3.3 that (3.3) holds with $\mu=\mu_{2}$ only for some $x$ and $y$ with $\varphi_{2}(x) \neq \varphi_{2}(y)$, not for all. This is as opposed to the obvious extension of Definition 2.2, which would require $\mu=\mu_{2}$ for all $x, y$ with $\varphi_{2}(x) \neq \varphi_{2}(y)$. The change of the definition is dictated by technical considerations: the stronger condition seems to be very hard to verify with the exception of some trivial examples.

Theorem 3.4 If a triangle $D$ has an obtuse angle (i.e., strictly greater than $\pi / 2)$ then the synchronous coupling for the reflected Brownian motion in $D$ is efficient.

Proof : This follows the idea of Corollary 2.7 (ii). Suppose that $D$ is an obtuse triangle. We will suppose without loss of generality that the longest side of $\partial D$ lies on the horizontal axis. By convention, the angle $\angle L$ between a straight line $L$ and the horizontal axis will lie in $(-\pi / 2, \pi / 2]$. Let the angles formed by the sides $I_{1}, I_{2}$ and $I_{3}$ of the triangle with the horizontal axis be $0, \alpha \in(-\pi / 2,0)$ and $\beta \in(0, \pi / 2)$.

A remark following [23, Corollary (6.31)] may be applied to the operator $\Delta+\mu_{2}$. Since this operator has analytic coefficients, the eigenfunctions must be analytic. In particular, an eigenfunction cannot be constant on a non-empty open set unless it is identically equal to 0 . Since $\varphi_{2}$ is not constant, we can find $x, y \in D$ (with $\left.\varphi_{2}(x) \neq \varphi_{2}(y)\right)$ such that the angle $\angle K$ between the line $K$ passing through these points and the horizontal axis belongs to $(\alpha, \beta)$. Moreover, we choose the points $x=\left(x_{1}, x_{2}\right)$ and $y=\left(y_{1}, y_{2}\right)$ so that $x_{1}<y_{1}$.

Let $(X, Y)$ be a synchronous coupling of reflecting Brownian motions in $D$ which starts from $(x, y)$. Let $K_{t}$ be the line passing through $X_{t}$ and $Y_{t}$. It follows easily from the construction of the synchronous coupling at the beginning of the section that $\angle K_{t}$ will monotonically move towards $\alpha$, as long as one of the processes is reflecting on the side $I_{2}$. Likewise, $\angle K_{t}$ can monotonically approach 0 or $\beta$, depending on the side where the reflection is taking place. We conclude that $\angle K_{t}$ will stay within $[\alpha, \beta]$ for all $t$. This part of the proof uses the obtuse property of the triangle in a crucial way.

Since $\angle K_{t} \in[\alpha, \beta] \subset(-\pi / 2, \pi / 2)$, we have $\left|X_{t}-Y_{t}\right| \leq k\left(Y_{t}^{1}-X_{t}^{1}\right)$ for some $k<\infty$. By equation (3.2),

$$
\begin{aligned}
& \mathbb{E}\left(\left|X_{t}-Y_{t}\right| \mid X_{0}=x, Y_{0}=y\right) \\
\leq & k \mathbb{E}\left(Y_{t}^{1} \mid X_{0}=x, Y_{0}=y\right)-k \mathbb{E}\left(X_{t}^{1} \mid X_{0}=x, Y_{0}=y\right) \\
= & k\left(\int_{D} p(t, y, z) z_{1} d z-\int_{D} p(t, x, z) z_{1} d z\right) \\
= & k\left(\int_{D}\left[\varphi_{2}(x)-\varphi_{2}(y)\right] z_{1} \varphi_{2}(z) e^{-\mu_{2} t} d z\right)+R(t, x, y) \\
\leq & c_{1}(x, y) e^{-\mu_{2} t} .
\end{aligned}
$$

It follows that inequality (3.3) holds with $\mu=\mu_{2}$, and so the coupling is efficient.

Remark : We will prove in Theorem 3.7 below that the mirror coupling is not efficient for reflected Brownian motion in a triangle if all its angles are acute (smaller than $\pi / 2$ ) and distinct. One may ask if the same is true for the synchronous coupling. We presently do not know the answer although we guess that synchronous couplings are also inefficient for triangles with acute angles. We will outline below an argument which shows that the synchronous coupling has a property similar to the "transposition property" discussed in Section 2, which is the basis for the proofs of inefficiency in Theorems 2.6 (i) and 3.7 (ii). Then we will show why this property alone is not sufficient to complete the proof of inefficiency.

Assume that all angles of the triangle $D$ are less than $\pi / 2$. Let us consider two possible scenarios for motions of the line $K_{t}$ passing through $X_{t}$ and $Y_{t}$. Suppose that $\left|X_{0}-Y_{0}\right|$ is small and both starting points are close to the center of $D$. In the first scenario the particles $X_{t}$ and $Y_{t}$ move around each other by the angle $\pi$ during a time interval $\left(t_{1}, t_{4}\right)$. This happens thus: one of the processes first reflects on the side $I_{1}$ so that $K_{t}$ becomes parallel to this side at time $t_{1}$, then one reflects on $I_{2}$ until $K_{t}$ becomes parallel to $I_{2}$ at time $t_{2}$, then one reflects on $I_{3}$ until $K_{t}$ becomes parallel to $I_{3}$ at time $t_{3}$, and finally one reflects on $I_{1}$, until $K_{t}$ again becomes parallel to $I_{1}$ at time $t_{4}$. From elementary geometry and the acuteness of all angles of $D$ it follows that $X_{t_{1}}^{1}-Y_{t_{1}}^{1}$ and $X_{t_{4}}^{1}-Y_{t_{4}}^{1}$ have opposite signs. Let $\alpha_{j}$ be the angle of $D$ opposite to $I_{j}$. Then trigonometry and the synchronized property of the coupling combine to show that

$$
\left|X_{t_{2}}-Y_{t_{2}}\right|=\cos \alpha_{3}\left|X_{t_{1}}^{1}-Y_{t_{1}}^{1}\right|
$$

By repeating this remark we see that

$$
\left|X_{t_{4}}-Y_{t_{4}}\right|=\cos \alpha_{2} \cos \alpha_{1} \cos \alpha_{3}\left|X_{t_{1}}-Y_{t_{1}}\right|
$$

In the second scenario the particles do not revolve around each other. The processes instead reflect on $I_{1}$ until $K_{t}$ is parallel to $I_{1}$ at time $t_{5}$, then they reflect on $I_{2}$, and finally on $I_{1}$, so that $K_{t}$ is again parallel to $I_{1}$ at time $t_{6}$. Then $X_{t_{5}}^{1}-Y_{t_{5}}^{1}$ and $X_{t_{6}}^{1}-Y_{t_{6}}^{1}$ will have the same sign and

$$
\left|X_{t_{6}}-Y_{t_{6}}\right|=\cos \alpha_{3} \cos \alpha_{3}\left|X_{t_{5}}-Y_{t_{5}}\right| .
$$

Now suppose that the angles $\alpha_{j}$ are such that for some integers $n$ and $m(n$ odd) we have

$$
\left(\cos \alpha_{2} \cos \alpha_{1} \cos \alpha_{3}\right)^{n}=\left(\cos \alpha_{3}\right)^{2 m} .
$$

The family of triplets $\left(\alpha_{1}, \alpha_{2}, \alpha_{3}\right)$ with this property is dense in the set of all possible acute angles with $\alpha_{1}+\alpha_{2}+\alpha_{3}=\pi$. To show this amounts to considering angles for which the quantity

$$
\frac{1}{2}\left(1+\frac{\cos \left(\alpha_{1}-\alpha_{2}\right)}{\cos \alpha_{3}}\right)^{n} /\left(\cos \alpha_{3}\right)^{2 m-n}
$$

is equal to 1 : the required density statement follows by noting that $x^{n} / y^{2 n-m}$ is dense in $(0, \infty)$ for $x, y \in(0,1)$ with $x / y$ irrational.

If the process $(X, Y)$ repeats $n$ times the motions described in the first scenario, then for some $t_{7}$,

$$
\left|X_{t_{7}}-Y_{t_{7}}\right|=\left(\cos \alpha_{2} \cos \alpha_{1} \cos \alpha_{3}\right)^{n}\left|X_{t_{1}}-Y_{t_{1}}\right|
$$

and $X_{t_{7}}^{1}-Y_{t_{7}}^{1}$ and $X_{t_{1}}^{1}-Y_{t_{1}}^{1}$ will have opposite sign (since $n$ is odd). If the second scenario is repeated $m$ times then for some $t_{8}$,

$$
\left|X_{t_{8}}-Y_{t_{8}}\right|=\left(\cos \alpha_{3}\right)^{2 m}\left|X_{t_{5}}-Y_{t_{5}}\right|=\left(\cos \alpha_{2} \cos \alpha_{1} \cos \alpha_{3}\right)^{n}\left|X_{t_{5}}-Y_{t_{5}}\right|
$$

and $X_{t_{8}}^{1}-Y_{t_{8}}^{1}$ and $X_{t_{5}}^{1}-Y_{t_{5}}^{1}$ will have the same sign. Hence, under the two scenarios the moving particles may become parallel to $I_{1}$ and at the same distance from each other, but with their "order" reversed. Since the densities for the processes $\left(X_{t}, Y_{t}\right)$ and $\left(Y_{t}, X_{t}\right)$ are both continuous, the corresponding distributions must be mutually absolutely continuous on a part of the state space. This is a version of the "transposition property" used in our proofs of inefficiency.

Next we indicate why it is difficult to derive inefficiency of the synchronous coupling directly from the "transposition property." Suppose that the points $x$ and $y$ belong to $D$, the line passing through them is parallel to $I_{1}$, and $X_{0}=x, Y_{0}=y$. Consider the event $A_{u}$ that neither process $X_{t}$ or $Y_{t}$ touches $I_{2}$ or $I_{3}$ before time $u$. Note that given $A_{u}$, we have $\left|X_{u}-Y_{u}\right|=\left|X_{0}-Y_{0}\right|$. Hence,

$$
\mathbb{E}\left(\left|X_{u}-Y_{u}\right| \mid X_{0}=x, Y_{0}=y\right) \geq\left|X_{0}-Y_{0}\right| \mathbb{P}\left(A_{u} \mid X_{0}=x, Y_{0}=y\right)
$$

If the points $x$ and $y$ are very close to each other, the event $A_{u}$ is "almost identical" to the event $A_{u}^{1}$ that $X_{t}$ does not hit $I_{2}$ or $I_{3}$ before time $u$. For large $u$, the probability of $A_{u}^{1}$ is well approximated by $c(x, y) \exp \left(-u \mu^{*}\right)$, where $\mu^{*}$ is the first eigenvalue for the Laplacian in $D$ with Neumann boundary conditions on $I_{1}$ and Dirichlet conditions on $I_{2} \cup I_{3}$. Hence, we have a heuristic estimate $\mathbb{P}\left(A_{u}\right) \approx \mathbb{P}\left(A_{u}^{1}\right) \approx c \exp \left(-u \mu^{*}\right)$. It is conceivable that $\left|X_{0}-Y_{0}\right| \mathbb{P}\left(A_{u}\right)$ is the main contribution to $\mathbb{E}\left(\left|X_{u}-Y_{u}\right|\right)$, since the length of the vector $X_{t}-Y_{t}$ is shortened exponentially fast over long intervals of time owing to reflection of the processes $X_{t}$ and $Y_{t}$ on $\partial D$. If this is the case then we might have $\mathbb{E}\left(\left|X_{u}-Y_{u}\right|\right) \approx$ $c \exp \left(-u \mu^{*}\right)$, for large $u$. So now the question is whether $\mu^{*}$ might be equal to the second Neumann eigenvalue $\mu_{2}$. This is rather doubtful but at present we do not know how to prove that the two eigenvalues are different. It should be noted that for an arbitrary convex planar domain and arbitrary division of the boundary into the "Neumann" and "Dirichlet" parts, the first mixed eigenvalue can be smaller than, equal to or larger than the second Neumann eigenvalue. Hence, there is no general principle that would show that $\mu^{*} \neq \mu_{2}$.

For the remaining part of the section, we switch our attention to the "mirror" coupling for reflected Brownian motion in planar domains. The mirror coupling seems to be the most natural coupling for diffusions in $\mathbb{R}^{d}$ and for reflected Brownian motion in particular. One feels that the mirror coupling is optimal from the point of view of efficiency but we do not have any rigorous results to this effect. See [10, Theorem 5.3] for results on other versions of optimality for couplings of diffusions.

Mirror couplings for reflected processes have been constructed in [29]. We will present a new construction of mirror couplings which is particularly well suited for the study of those of its properties which are important in this paper. We will start with the discussion of the mirror coupling in very simple domains and then (in Section 4) we will progress towards more complicated domains.

First we discuss the mirror coupling for free Brownian motions in $\mathbb{R}^{2}$. Suppose that $x, y \in \mathbb{R}^{2}, x \neq y$, and that $x$ and $y$ are symmetric with respect to a line $M$. Let $X$ be a Brownian motion starting from $x$ and let $\tau$ be the first time $t$ with $X_{t} \in M$. Then we let $Y$ be the mirror image of $X$ with respect to
$M$ for $t \leq \tau$, and we let $Y_{t}=X_{t}$ for $t>\tau$. The process $Y$ is a Brownian motion starting from $y$ and ( $X, Y$ ) is called the mirror coupling for (non-reflecting) Brownian motion.

We start the discussion of the mirror coupling for reflected Brownian motions with the simplest case, that of a half-plane $D$. Suppose $x, y \in D$ and let $M$ be the line of symmetry for $x$ and $y$. The case when $M$ is parallel to $\partial D$ can be easily handled using Skorokhod's lemma [28, Lemma 3.6.14], so we focus on the case when $M$ intersects $\partial D$. By performing rotation and translation, if necessary, we may suppose that $D$ is the upper halfplane. Let $h$ be the point of intersection between the boundary $\partial D$ and the line of symmetry $M$. We write $x=\left(r^{x}, \theta^{x}\right)$ and $y=\left(r^{y}, \theta^{y}\right)$ in polar coordinates based on $h$. The points $x$ and $y$ are at the same distance from $h$ so $r^{x}=r^{y}$. Suppose without loss of generality that $\theta^{x}<\theta^{y}$. We first generate a 2-dimensional Bessel process $R_{t}$ starting from $r^{x}$. Then we generate two coupled one-dimensional processes on the "half-circle" as follows. Let $W$ be a 1-dimensional Brownian motion starting from 0 . We construct $\tilde{\Theta}^{x}$ as the reflected Brownian motion in $[0, \pi]$ started at $\theta^{x}$, solving the Skorokhod equation

$$
\tilde{\Theta}_{t}^{x}=\theta^{x}+W_{t}+L^{o}\left(\tilde{\Theta}^{x}\right)_{t}-L^{\pi}\left(\tilde{\Theta}^{x}\right)_{t}
$$

where $L^{o}\left(\tilde{\Theta}^{x}\right), L^{\pi}\left(\tilde{\Theta}^{x}\right)$ are the local time "pushes" for $\tilde{\Theta}^{x}$ at 0 and $\pi$ (the minimal increasing processes required to keep $\tilde{\Theta}^{x}$ nonnegative and no greater than $\pi)$. We construct $\hat{\Theta}^{y}$ similarly but using a mirror-reflected driving Brownian motion:

$$
\hat{\Theta}_{t}^{y}=\theta^{y}-W_{t}+L^{o}\left(\hat{\Theta}^{y}\right)_{t}-L^{\pi}\left(\hat{\Theta}^{y}\right)_{t} .
$$

These reflecting Brownian motions have to be time-changed in order to serve as the angular parts of reflected Brownian motion in $D$ : fortunately we can use the same time-change in each case, namely

$$
\sigma(t)=\int_{0}^{t} R_{s}^{-2} d s
$$

We set

$$
X_{t}=\left(R_{t}, \tilde{\Theta}_{\sigma(t)}^{x}\right), \quad Y_{t}=\left(R_{t}, \hat{\Theta}_{\sigma(t)}^{y}\right)
$$

This is a generalization of the skew-product representation of planar Brownian motion, as described in [27]. Both $X$ and $Y$ can be viewed as obtained from free Brownian motions using reflection in the boundary of $D$. Indeed the processes $X_{t}$ and $Y_{t}$ behave like free Brownian motions coupled by the mirror coupling as long as they are both strictly inside $D$. The processes will stay together after the first time they meet. The pair $\left(X_{t}, Y_{t}\right)$ will be called the mirror coupling for reflected Brownian motions in a half-plane.

The line of symmetry for $X_{t}$ and $Y_{t}$ will be denoted $M_{t}$ if $X_{t} \neq Y_{t}$. For definiteness, we let $M_{t}$ be the horizontal line passing through $X_{t}$ if $X_{t}=Y_{t}$.

The most important property of the above coupling is that by construction the distances of $X_{t}$ and $Y_{t}$ from $h$ remain equal to each other as time $t$ varies. This property manifests itself in more general domains in the following way.

First of all, suppose that $D$ is an arbitrary halfplane, and $x$ and $y$ belong to $D$. Let $M_{t}$ be the line of symmetry for $X_{t}$ and $Y_{t}$, constructed as above as a mirror coupling of reflected Brownian motion begun at $x, y$ respectively. Suppose that $M_{0}$ intersects $\partial D$. Then for every $t$, the distance from $X_{t}$ to $M_{t} \cap \partial D$ is the same as for $Y_{t}$. Note that $M_{t}$ may move, but only in a continuous way, while the point $M_{t} \cap \partial D$ will never move. We will call $M_{t}$ the mirror and the intersection point $h$ of $M_{t}$ and $\partial D$ will be called the hinge. The absolute value of the angle between the mirror and the normal vector to $\partial D$ at $h$ can only decrease; thus if $M_{t}$ is parallel to $\partial D$ then it will stay parallel to $\partial D$ until the coupling time. In this case, $M_{t}$ can move only away from $\partial D$ and only in a continuous fashion.

The mirror coupling of reflected Brownian motions in a convex polygonal domain $D$ can be described as follows. Suppose that $X_{t}$ and $Y_{t}$ start from $x$ and $y$ inside the domain $D$. As soon as one of the particles hits a side $I$ of $\partial D$, the processes will evolve according to the coupling described in the previous paragraph. To be more precise, let $K$ be the straight line containing $I$ where $I$ is the side of $\partial D$ most recently hit by one of the particles. Since the process which hits $I$ does not "feel" the shape of $D$ except for the direction of $I$, it follows that $X_{t}$ and $Y_{t}$ will remain at the same distance from the hinge $\left\{h_{t}\right\}=M_{t} \cap K$, as long as the particles do not hit a side different from $I$. The mirror $M_{t}$ can move but the hinge $h_{t}$ will remain constant as long as $I$ remains the side of $\partial D$ where the reflection takes place. The hinge $h_{t}$ will jump when the reflection location moves from $I$ to another side of $\partial D$. Since $D$ is convex, $h_{t}$ will be always on $\partial D$ or outside $D$.

Remark : We remark in passing a point of methodological interest: this representation was first discovered by accident as we explored the system of mirrorcoupled reflecting Brownian motions using computer algebra, specifically the implementation Itovsn3 of stochastic calculus in the computer algebra package REDUCE. For details (as implemented in the Mathematica version of Itovsn3) see the Mathematica notebook reflect in [30]. Of course with hindsight the properties mentioned above now appear obvious ...

Recall that $p(t, x, y)$ denotes the transition densities for reflecting Brownian motion in the triangle $D$ and recall from (3.2) the following one-term eigenfunction expansion for $p(t, x, y)$,

$$
\begin{equation*}
p(t, x, y)=c_{1}+\varphi_{2}(x) \varphi_{2}(y) e^{-\mu_{2} t}+R(t, x, y) \tag{3.4}
\end{equation*}
$$

where $R(t, x, y)$ converges to 0 faster than $e^{-\mu_{2} t}$ as $t \rightarrow \infty$. The coupling time is denoted by $\tau$. Suppose that one can prove that for some $\mu \geq 0$ and $x, y \in D$,

$$
\begin{equation*}
\mathbb{P}\left(\tau>t \mid X_{0}=x, Y_{0}=y\right) \leq c(x, y) e^{-\mu t}, \quad t \geq 0 \tag{3.5}
\end{equation*}
$$

It is reasonable to expect that $\mu$ is then a lower bound for $\mu_{2}$. However, Example 3.1 applies to the mirror coupling as well and we see that (3.5) may hold for some $\mu>\mu_{2}$ and some $x, y \in D$. The following lemma is entirely analogous to Lemma 3.2.

Lemma 3.5 If $\varphi_{2}(x) \neq \varphi_{2}(y)$ and (3.5) holds, then $\mu \leq \mu_{2}$.

Proof : Consider $x, y \in D$ such that $\varphi_{2}(x) \neq \varphi_{2}(y)$. Let $A=\{v \in D$ : $\left.\varphi_{2}(v)>0\right\}$ and note that $A$ must have positive measure. By (3.4),

$$
\begin{gathered}
\int_{A} p(t, x, z) d z-\int_{A} p(t, y, z) d z=\int_{A}\left[\varphi_{2}(x)-\varphi_{2}(y)\right] \varphi_{2}(z) e^{-\mu_{2} t} d z+R(t, x, y) \\
=c_{1}(x, y) e^{-\mu_{2} t}+R(t, x, y)
\end{gathered}
$$

where $c_{1}(x, y) \neq 0$. On the other hand,

$$
\begin{aligned}
& \left|\int_{A} p(t, x, z) d z-\int_{A} p(t, y, z) d z\right| \\
& \quad=\left|\mathbb{P}\left(X_{t} \in A \mid X_{0}=x\right)-\mathbb{P}\left(Y_{t} \in A \mid Y_{0}=y\right)\right| \\
& \quad=\left|\mathbb{P}\left(X_{t} \in A, t<\tau \mid X_{0}=x\right)-\mathbb{P}\left(Y_{t} \in A, t<\tau \mid Y_{0}=y\right)\right| \\
& \quad \leq \mathbb{P}\left(t<\tau \mid X_{0}=x, Y_{0}=y\right) \leq c_{2}(x, y) e^{-\mu t}
\end{aligned}
$$

Since this inequality and (3.6) hold for arbitrarily large $t$, we see that $\mu \leq \mu_{2}$.

Definition 3.6 A mirror coupling $(X, Y)$ of reflected Brownian motions in $D$ is said to be efficient if the estimate (3.5) holds with $\mu=\mu_{2}$ for some $x$ and $y$ with $\varphi_{2}(x) \neq \varphi_{2}(y)$.

Our main theorem for mirror coupling in triangles identifies the cases of inefficiency and efficiency for this coupling in simple geometric terms.

## Theorem 3.7

(i) If a triangle $D$ has an obtuse angle (which is to say, strictly greater than $\pi / 2)$ then the mirror coupling for the reflected Brownian motion in $D$ is efficient.
(ii) If all angles of the triangle $D$ are distinct from each other and acute (which is to say, strictly less than $\pi / 2$ ) then the mirror coupling for the reflected Brownian motion in $D$ is not efficient.

## Remark :

(i) Note that in Theorem 3.7 (ii) we assume that all angles of $D$ are distinct. This technical assumption is probably unnecessary, but would be tedious to lift.
(ii) Example 2.12 and Theorem 3.7 (ii) naturally lead to the following open question: Are there no efficient Markovian couplings for reflected Brownian motion in generic acute triangles?

Proof of Theorem 3.7 (i): Suppose that an obtuse triangle $D$ is oriented so that its longest side lies on the horizontal axis, its leftmost vertex is at the origin and the triangle is contained in the first quadrant (see Fig. 3). The angle formed by any straight line $K$ with the horizontal axis will be denoted $\angle K$. Let the angles formed by the sides $I_{2}$ and $I_{1}$ of the triangle $D$ with the horizontal axis be $\alpha \in(-\pi / 2,0)$ and $\beta \in(0, \pi / 2)$ (see Fig. 3).
[Figure 3 about here.]
Fix any two points $x=\left(x_{1}, x_{2}\right) \in D$ and $y=\left(y_{1}, y_{2}\right) \in D$, such that $x_{1}<y_{1}$ and $\angle M \in \mathcal{A}=(\pi / 2+\alpha, \pi / 2+\beta)$, where $M$ is the line of symmetry for $x$ and $y$. Consider a mirror coupling $\left(X_{t}, Y_{t}\right)$ with $X_{0}=x, Y_{0}=y$, and recall that $M_{t}$ denotes the mirror, i.e., the line of symmetry for $X_{t}$ and $Y_{t}$. Since $M_{0}=M$ we have $\angle M_{0} \in \mathcal{A}$. We will argue that $\angle M_{t}$ will not leave the interval $\mathcal{A}$ until the coupling time. Let $K_{j}$ denote the straight line containing the side $I_{j}$ of the triangle. Suppose that for some $s$, the angle $\angle M_{s}$ is within this interval and the hinge $h_{t}$ belongs to $K_{3}$ for $t \in(s, u)$. Then $\left|\pi / 2-\angle M_{t}\right|$ will be decreasing on the interval $(s, u)$, and so $\angle M_{t} \in \mathcal{A}$ for all $t \in(s, u)$. Next consider the case when $\angle M_{s} \in \mathcal{A}$ and $h_{t} \in K_{1}$ for $t \in(s, u)$. Then $\left|(\pi / 2+\beta)-M_{t}\right|$ is decreasing for $t \in(s, u)$ and so $\angle M_{t}$ must stay in $\mathcal{A}$ for $t \in(s, u)$. The final case, when the hinge belongs to $K_{2}$, may be treated in the same way. We have shown that $\angle M_{t}$ does not leave $\mathcal{A}$ before the coupling time.

By acuteness of the triangle $D$, the interval $\mathcal{A}$ lies strictly inside $(0, \pi)$, so there exists $k>0$ such that $Y_{t}^{1}-X_{t}^{1} \geq k\left|Y_{t}-X_{t}\right|$ for $t<\tau$ (and so for all $t$ ). We will now analyze the distance between $X_{t}$ and $Y_{t}$. Up to the coupling time, the process $\rho=|Y-X|$ is a one-dimensional Brownian motion with twice the standard variance as long as both $X$ and $Y$ are strictly inside $D$. When one of the processes $X$ or $Y$ is reflecting on $\partial D$, then $\rho$ gets a "push" determined by the local time spent by $X$ or $Y$ on $\partial D$ and by the direction of $M$ relative to the reflecting side of $D$. Since $D$ is convex, the direction of the push for $\rho$ always points towards 0 . This shows that for any $\rho$, the hitting distribution of 0 for the process $\rho$ starting from $\rho_{0}$ is stochastically majorized by the hitting distribution of 0 for the one-dimensional Brownian motion with twice the standard variance and starting from $\rho_{0}$. Hence, we may fix arbitrarily small $p_{0}>0$ and find $\hat{\rho}>0$ such that if $\rho_{t} \leq \hat{\rho}$ (but the positions of $X_{t}$ and $Y_{t}$ are otherwise arbitrary) then $\mathbb{P}\left(\rho_{t+1}>0\right)<p_{0}$. Choose $p_{0}$ such that $\left(2 p_{0}\right)^{j}<e^{-2 \mu_{2} j}$, and find a corresponding $\hat{\rho}>0$ with $\hat{\rho}<|x-y|$.

Recall from the proof of Theorem 3.4 that an eigenfunction must be analytic. In particular, an eigenfunction cannot vanish on a non-empty open set unless it is identically equal to 0 . Fix any $x \in D$ and find $y \in D$ such that $\varphi_{2}(x) \neq \varphi_{2}(y)$ and $\angle M \in(\pi / 2+\alpha, \pi / 2+\beta)$, where $M$ is the line of symmetry for $x$ and $y$. Such a point $y$ must exist because otherwise $\varphi_{2}$ would be constant, and, therefore, it would vanish, on a non-empty open set inside $D$, which is impossible.

Consider $t$ such that

$$
\begin{equation*}
\mathbb{P}\left(\rho_{t}>\hat{\rho} \mid t<\tau, X_{0}=x, Y_{0}=y\right) \geq p_{0} . \tag{3.6}
\end{equation*}
$$

Recall that $Y_{t}^{1}-X_{t}^{1} \geq k\left|Y_{t}-X_{t}\right|=k \rho_{t}$ for $t<\tau$. This and (3.6) show that

$$
\begin{equation*}
\mathbb{P}\left(Y_{t}^{1}-X_{t}^{1}>k \hat{\rho} \mid t<\tau, X_{t}=x, Y_{t}=y\right) \geq p_{0} . \tag{3.7}
\end{equation*}
$$

Since $Y_{t}^{1} \geq X_{t}^{1}$ for all $t$, (3.7) implies that

$$
\mathbb{E}\left(Y_{t}^{1} \mid t<\tau, X_{0}=x, Y_{0}=y\right) \geq \mathbb{E}\left(X_{t}^{1} \mid t<\tau, X_{0}=x, Y_{0}=y\right)+p_{0} k \hat{\rho}
$$

and therefore we have

$$
\begin{gather*}
\mathbb{E}\left(Y_{t}^{1} \mid X_{0}=x, Y_{0}=y\right)-\mathbb{E}\left(X_{t}^{1} \mid X_{0}=x, Y_{0}=y\right) \\
\geq \quad p_{0} k \hat{\rho} \mathbb{P}\left(t<\tau \mid X_{0}=x, Y_{0}=y\right) \tag{3.8}
\end{gather*}
$$

By (3.4),

$$
\begin{aligned}
& \mathbb{E}\left[Y_{t}^{1} \mid X_{0}=x, Y_{0}=y\right]-\mathbb{E}\left[X_{t}^{1} \mid X_{0}=x, Y_{0}=y\right] \\
& =\int_{D} p(t, y, z) z_{1} d z-\int_{D} p(t, x, z) z_{1} d z \\
& =\int_{D}\left[\varphi_{2}(x)-\varphi_{2}(y)\right] z_{1} \varphi_{2}(z) e^{-\mu_{2} t} d z+R(t, x, y)
\end{aligned}
$$

It follows from this and inequality (3.8) that

$$
\begin{equation*}
\mathbb{P}\left(t<\tau \mid X_{0}=x, Y_{0}=y\right) \quad \leq \quad c(x, y) e^{-\mu_{2} t} \tag{3.9}
\end{equation*}
$$

Next we consider a $t$ for which inequality (3.6) fails. Let $s$ be the supremum of times less than $t$ for which (3.6) holds. Let $j_{0}$ and $j_{1}$ be the smallest and largest integers in $(s, t)$. If there are no such $j_{0}, j_{1}$ then

$$
\begin{aligned}
\mathbb{P}\left[t \leq \tau \mid X_{0}=x, Y_{0}=y\right] & \leq \mathbb{P}\left[s \leq \tau \mid X_{0}=0, Y_{0}=y\right] \\
& \leq c(x, y) e^{-\mu_{2} s} \leq c^{\prime}(x, y) e^{-\mu_{2} t}
\end{aligned}
$$

On the other hand if $\ell$ is an integer in $\left[j_{0}, j_{1}-1\right]$, then by the definition of $s$, (3.6) fails for $\ell$. If there is no coupling by time $\ell$ and $\rho_{\ell} \leq \hat{\rho}$ then the probability of no coupling by time $\ell+1$ is less than $p_{0}$. This and the failure of inequality (3.6) at time $\ell$ imply that

$$
\begin{aligned}
\mathbb{P}\left(\ell+1<\tau \mid X_{0}=x, Y_{0}=y\right) & \leq \\
\left(p_{0}+\mathbb{P}\left[\rho_{\ell}>\hat{\rho} \mid \ell<\tau, X_{0}=x, Y_{0}=y\right]\right) & \times \mathbb{P}\left[\ell<\tau \mid X_{0}=x, Y_{0}=0\right] \\
& \leq 2 p_{0} \mathbb{P}\left(\ell<\tau \mid X_{0}=x, Y_{0}=y\right)
\end{aligned}
$$

Thus, applying (3.9) to $s$,

$$
\begin{aligned}
\mathbb{P}\left(t<\tau \mid X_{0}=x, Y_{0}=y\right) & \leq \mathbb{P}\left(j_{1}<\tau \mid X_{0}=x, Y_{0}=y\right) \\
& \leq\left(2 p_{0}\right)^{j_{1}-j_{0}} \mathbb{P}\left(j_{0}<\tau \mid X_{0}=x, Y_{0}=y\right) \\
& \leq\left(2 p_{0}\right)^{t-s-2} \mathbb{P}\left(s<\tau \mid X_{0}=x, Y_{0}=y\right) \\
& \leq e^{-2 \mu_{2}(t-s-2)} c(x, y) e^{-\mu_{2} s} \leq c_{1}(x, y) e^{-\mu_{2} t}
\end{aligned}
$$

We see that (3.9) extends to all $t \geq 0$. Since we have chosen $x$ and $y$ with $\varphi_{2}(x) \neq \varphi_{2}(y)$, we conclude that the mirror coupling is efficient in obtuse triangles.

We defer the proof of Theorem 3.7 (ii) till we have proved several subsidiary lemmas.

Let $\mathcal{D}=\{(x, y) \in \bar{D} \times \bar{D}: x \neq y\}, \mathcal{D}(\varepsilon)=\{(x, y) \in \bar{D} \times \bar{D}:|x-y| \geq \varepsilon\}$ and $\hat{\mathcal{D}}(\varepsilon)=\mathcal{D} \backslash \mathcal{D}(\varepsilon)$. In an abuse of notation we use $T(A)$ to denote the hitting time of $A$ for any process, including for example $X$ and $(X, Y)$. Sometimes the notation will record the process as well, as in $T_{X}(A)$. We work under the hypotheses of Theorem 3.7 (ii); the domain is an acute-angled triangle all of whose angles are different from each other.

Lemma 3.8 For sufficiently small $a \in(0, \operatorname{diam}(D) / 10)$ there exist $s, c_{1}>0$, $\mathcal{D}_{1} \subset \mathcal{D}(a)$, and a (probability) measure $\nu$ on $\mathcal{D}_{1}$ with $\nu\left(\mathcal{D}_{1}\right)>0$, such that for all $(x, y) \in \mathcal{D}(a)$ and for every subset $A$ of $\mathcal{D}_{1}$ we have

$$
\mathbb{P}\left(\left(X_{s}, Y_{s}\right) \in A \mid X_{0}=x, Y_{0}=y\right) \quad>\quad c_{1} \nu(A)
$$

Proof : Most of the proof is concerned with a description of "all possible" trajectories of $(X, Y)$, before the coupling time. Our description will be partly given in terms of possible motions of the mirror process $M$ and will be partly qualitative in nature. We are interested in all trajectory-related events of positive probability, no matter how small that probability might be.

We will say that a positive measure $\nu_{1}$ is a component of a (probability) measure $\nu_{2}$ if $\nu_{1}(A) \leq \nu_{2}(A)$ for all $A$.

Suppose that $X_{0}=x, Y_{0}=y, x, y \in D, x \neq y$. Suppose further that $B_{X}$ and $B_{Y}$ are non-empty open subsets of $D$ which are mirror-images of each other with respect to $M_{0}$ and such that $B_{X}$ lies totally on the same side of $M_{0}$ as $X_{0}$. Then it is easy to see that the coupling process $(X, Y)$ may reach $B_{X} \times B_{Y}$ without touching the boundary of $D$, and moreover this can happen in an arbitrarily short time. In particular $X, Y$ can come arbitrarily close to any one of the points of intersection of the initial position $M_{0}$ of the mirror with $\partial D$, before the mirror $M$ has first moved.

Now fix one of the points in $M_{0} \cap \partial D$. Call this point $h$ and assume that $h$ is not a vertex of the triangle $D$. Let $\theta_{M}(t)$ be the angle between $M_{t}$ and the side $I_{1}$ containing $h$. We will argue that if $X_{0}$ and $Y_{0}$ are close to $h$ (if they are not, they can move close to $h$, by our previous remarks), then the mirror can turn around $h$ in the direction towards the normal (i.e, $\theta_{M}(t)$ will monotonically move towards $\pi / 2$ ), and for each $t>0$, the angle between $M_{t}$ and $I_{1}$ is a random variable which has a non-trivial atom at $\pi / 2$ and a component with a strictly positive continuous density on ( $\theta_{M}(0), \pi / 2$ ) (or $\left(\pi / 2, \theta_{M}(0)\right)$ ). We will show that all this may happen before $X$ and $Y$ leave a small neighborhood of $h$, and so the hinge, i.e., the point of intersection of the mirror $M$ with $\partial D$ around which the mirror is turning, will remain fixed at $h$.

The next part of our argument will be quantitative in nature; note that we actually prove more than is strictly needed in this lemma.

Suppose that $X_{0}=x$ and $Y_{0}=y,|x-y|=\rho$, and at least one of the points $x$ or $y$ is at distance no more than $\rho$ from $I_{1}$. We moreover assume that the distance from $x$ to each one of the other sides of $\partial D$ is greater than $10 \rho$, likewise for $y$. Consider a polar coordinate system $(r, \theta)$ based on $h$ as origin, such that $I_{1}$ lies on the line $\{(r, \theta)$ : either $\theta=0$ or $\theta=\pi\}$. Without loss of generality assume that the distance from $x$ to $I_{1}$ is not smaller than that for $y$, that $\pi / 2 \geq \theta_{y}>0$, that $\theta_{x}>\theta_{y}$, and that $\pi-\theta_{x}>\theta_{y}$, where $x=\left(r_{x}, \theta_{x}\right)$ and $y=\left(r_{y}, \theta_{y}\right)$. We will write $X_{t}=\left(r_{X}(t), \theta_{X}(t)\right)$, and use corresponding polar coordinates for $Y$. We argue geometrically, considering the ray $\left\{(r, \theta): \theta=\theta_{M}(t)\right\}$ defining the mirror process $M$ for $X$ and $Y$. We will consider two cases. The first case (C1) is when $\theta_{x} \leq 3 \pi / 4$. Let $b_{1}>1$ be defined by

$$
b_{1}^{2}=\inf _{\beta \in(0,3 \pi / 4)} \frac{\sin (7 \beta / 12)}{\sin (\beta / 2)} .
$$

We define a ray segment

$$
Q_{1}=\left\{(r, \theta): r \geq \max \left(r_{x} / b_{1}, r_{x}-\rho\right), \theta=7 \theta_{x} / 6\right\}
$$

and the rest of a 4 -sided curvilinear domain

$$
\begin{aligned}
& Q_{2}=\quad\left\{(r, \theta): r \geq r_{x} / b_{1}, \theta=4 \theta_{x} / 5\right\} \cup \\
& \cup\left\{(r, \theta): r=\max \left(r_{x} / b_{1}, r_{x}-\rho\right) \text { or } r=r_{x}+\rho\right\}
\end{aligned}
$$

By scaling and the effect of the assumption (C1) on possible locations of $x, y$ (and the monotonic effect of reflection in the boundary) it follows that there exists $p_{1}>0$ such that

$$
\mathbb{P}\left(T_{X}\left(Q_{1}\right)<T_{X}\left(Q_{2}\right) \mid X_{0}=x, Y_{0}=y\right) \quad>\quad p_{1}
$$

Elementary trigonometry can be used to show that $\theta_{M}(0)<4 \theta_{x} / 5$, in view of the assumptions that $y$ is closer to $I_{1}$ than $x$, and not further from $I_{1}$ than $\rho$, and that $|x-y|=\rho$. We have

$$
\frac{\rho}{2 r_{x}}=\frac{|x-y| / 2}{r_{x}}=\sin \left(\left(\theta_{x}-\theta_{y}\right) / 2\right) .
$$

We also have $\theta_{M}(t) \leq \max \left(\theta_{M}(0), 7 \theta_{x} / 12\right)<4 \theta_{x} / 5$ as long as $\theta_{X}(t) \leq 7 \theta_{x} / 6$. This gives control over the coupling time $\tau$ : the event $\left\{T_{X}\left(Q_{1}\right)<\bar{T}_{X}\left(Q_{2}\right)\right\}$ implies $\left\{T_{X}\left(Q_{1}\right)<\tau\right\}$. The fact that $\theta_{X}(t)$ reaches the level $7 \theta_{x} / 6$ for the first time when $t=T_{X}\left(Q_{1}\right)$ implies that $\theta_{Y}\left(T_{X}\left(Q_{1}\right)\right) \leq \theta_{y}$. Using the fact that $Y$ is a reflection of $Y$ modified by reflection in $I_{1}$,

$$
\begin{aligned}
\frac{\left|X\left(T_{X}\left(Q_{1}\right)\right)-Y\left(T_{X}\left(Q_{1}\right)\right)\right| / 2}{r_{X}\left(T_{X}\left(Q_{1}\right)\right)} & =\sin \left(\left(\theta_{X}\left(T_{X}\left(Q_{1}\right)\right)-\theta_{Y}\left(T_{X}\left(Q_{1}\right)\right)\right) / 2\right) \\
\geq \sin \left(\left(7 \theta_{x} / 6-\left(\theta_{x}+\theta_{y}\right) / 2\right)\right. & =\sin \left(\theta_{x} / 12+\left(\theta_{x}-\theta_{y}\right) / 2\right) \\
\geq \sin \left(\left(\theta_{x}-\theta_{y}\right) / 12+\left(\theta_{x}-\theta_{y}\right) / 2\right) & =\sin \left(7\left(\theta_{x}-\theta_{y}\right) / 12\right)
\end{aligned}
$$

Using the construction of $Q_{1}$ and recalling the definition of $b_{1}$,

$$
\begin{aligned}
\left|X\left(T_{X}\left(Q_{1}\right)\right)-Y\left(T_{X}\left(Q_{1}\right)\right)\right| & \geq 2 r_{X}\left(T_{X}\left(Q_{1}\right)\right) \sin \left(7\left(\theta_{x}-\theta_{y}\right) / 12\right) \\
\geq 2\left(r_{x} / b_{1}\right) \sin \left(7\left(\theta_{x}-\theta_{y}\right) / 12\right) & \geq 2 r_{x} b_{1} \sin \left(\left(\theta_{x}-\theta_{y}\right) / 2\right) \\
& \geq 2 r_{x} b_{1} \frac{\rho}{2 r_{x}}=b_{1} \rho .
\end{aligned}
$$

We conclude that in case (C1),

$$
\mathbb{P}\left(T\left(\mathcal{D}\left(b_{1} \rho\right)\right)<\tau \mid X_{0}=x, Y_{0}=y\right)>p_{1}
$$

The second case (C2) is when $\theta_{x} \geq 3 \pi / 4$. Let

$$
\begin{aligned}
Q_{3}= & \{(r, \theta): \rho / 4<r<\rho, \theta=\pi\} \\
Q_{4}= & \{(r, \theta): \theta=\pi / 2\} \cup\{(r, \theta): r=3 \rho\} \cup \\
& \cup\{(r, \theta): \theta=\pi, r \leq \rho / 4 \text { or } r \geq \rho\}
\end{aligned}
$$

It is easy to see that there exists $p_{2}>0$ such that for $x$ and $y$ satisfying the assumptions of (C2),

$$
\mathbb{P}\left(T_{X}\left(Q_{3}\right)<T_{X}\left(Q_{4}\right) \mid X_{0}=x, Y_{0}=y\right)>p_{2}
$$

Note that $\left\{T_{X}\left(Q_{3}\right)<T_{X}\left(Q_{4}\right)\right\}$ implies $\left\{T_{X}\left(Q_{3}\right)<\tau\right\}$. If the former event occurs then the mirror for $X$ and $Y$ will be perpendicular to $I_{1}$ after time $T_{X}\left(Q_{3}\right)$ at least as long as $X$ and $Y$ do not leave the ball of radius $7 \rho$ around $h$.

Using the strong Markov property and repeated application of properties proved in cases (C1) and (C2), we see that the mirror may turn around $h$ towards the normal position while all the time $X$ and $Y$ may stay in a very small neighborhood of $h$, without coupling.

Next we will prove that the distribution of $\theta_{M}(t)$ has a component with a continuous density unless $\theta_{M}(0)=\pi / 2$. Recall our current assumptions that $X_{0}=x, Y_{0}=y, x, y \in D, x \neq y$. Suppose without loss of generality that $x$ is not closer to $I_{1}$ than $y$. Let $A$ be a closed disc in $D$, with non-empty interior, not far from $h$, on the same side of the mirror as $x$, and such that $x \notin A$. Assume without loss of generality that $I_{1}$ lies on the horizontal axis, $h=(0,0)$, and $\theta_{M}(0) \in(0, \pi / 2)$. Then apply the complex analytic transformation $z \rightarrow \log z$ to $X$ and to $Y$, viewing $\log$ as a mapping of the upper half-plane to the strip where the imaginary part is between 0 and $\pi$. We can make a single random time-change simultaneously converting each of the processes $\log (X)$ and $\log (Y)$ into reflected Brownian motions $\tilde{X}$ and $\tilde{Y}$ respectively. The same time-change works for both the processes since they are always the same distance from $h$. For the same reason the processes $\tilde{X}$ and $\tilde{Y}$ have the same real parts and they are related by a mirror coupling. Let $\tilde{X}^{2}$ and $\tilde{Y}^{2}$ be the imaginary parts of $\tilde{X}$ and $\tilde{Y}$, and let $\tilde{L}$ measure the local time spent by $\tilde{Y}$ on the real axis. Then

$$
\begin{equation*}
\frac{\tilde{Y}_{t}^{2}+\tilde{X}_{t}^{2}}{2}-\frac{\tilde{Y}_{s}^{2}+\tilde{X}_{s}^{2}}{2}=\frac{\tilde{L}_{t}-\tilde{L}_{s}}{2} . \tag{3.10}
\end{equation*}
$$

Let $\tilde{S}$ be the hitting time of $\log (A)$ by $\tilde{X}$. For a fixed $t$ we have that $\tilde{L}_{\tilde{S}}=\tilde{L}_{t}$, with positive probability. For a fixed $t$ we know that $\tilde{L}_{t}$ has a continuous density, so it follows from equation (3.10) that $\left(\tilde{Y}^{2}(\tilde{S})+\tilde{X}^{2}(\tilde{S})\right) / 2$ has a component with a continuous density. If $S$ is the hitting time of $A$ by $X_{t}$ then $\theta_{M}(S)=$ $\left(\tilde{Y}^{2}(\tilde{S})+\tilde{X}^{2}(\tilde{S})\right) / 2$. This shows that the distribution of $\theta_{M}(S)$ has a component with a continuous density. For a fixed $s$, the event $\left\{L_{S}=L_{s}\right\}$ has a positive probability, so $\theta_{M}(s)$ has a component with a continuous density. Furthermore it is now not hard to see that the density is strictly positive on $\left(\theta_{M}(0), \pi / 2\right)$.

Our argument so far shows that the mirror can turn at either point of intersection with $\partial D$ towards the normal direction, and the angle where it stops before switching the turning point (hinge) is a random variable with an atom at $\pi / 2$ and a component with a continuous positive density on the interval between the starting angle and $\pi / 2$. All this can happen with positive probability before the coupling time, and, moreover, any turning with a finite sum of all turning angles can be done in an arbitrarily small time, with positive probability.

We have explicitly and implicitly assumed in our arguments that $X_{0}$ and $Y_{0}$ belong to $D$. We will briefly discuss what may happen when $X_{0}$ or $Y_{0}$ belong to $\partial D$, including, possibly, one of the vertices. We do not assume that $X_{0}$ and $Y_{0}$ are necessarily close to the mirror $M_{0}$. Since $X$ and $Y$ spend zero time on the boundary of $D$, then with probability 1 , there will be arbitrarily small $s>0$ with $X_{s} \in D$ and $Y_{s} \in D$. Once both processes $X$ and $Y$ are strictly inside $D$, they can move in the way described earlier in the proof.

Next we will use the above results on the possible movements of the mirror to construct a component of the distribution of $\left(X_{1}, Y_{1}\right)$, for any starting points $x \neq y$ for $X$ and $Y$.

Let the sides of the triangle $D$ be called $I_{1}, I_{2}$ and $I_{3}$. Let us name the vertices of the triangle as follows: $\left\{z_{1}\right\}=I_{1} \cap I_{2},\left\{z_{2}\right\}=I_{2} \cap I_{3}$ and $\left\{z_{3}\right\}=I_{3} \cap I_{1}$. Let $F_{j}=\left(B\left(z_{j}, \rho_{1}\right) \backslash B\left(z_{j}, \rho_{1} / 2\right)\right) \cap \partial D$, where $\rho_{1}>0$ is chosen so that $F_{j}$ 's have the following property. For every $z \in F_{j}$, the line $K$ passing through $z$ and perpendicular to $\partial D$ at this point, crosses only the sides of $\partial D$ which are adjacent to $z_{j}$, and the points in $K \cap \partial D$ are at least $2 \rho_{1}$ units from the other vertices. Such a $\rho_{1}$ exists in view of the fact that all angles of $D$ are acute.

Suppose that $M_{0}$ does not pass through any vertex and let $h_{t}^{1}$ and $h_{t}^{2}$ be the points of intersection of $M_{t}$ with $\partial D$. Note that the hinge $h_{t}$, i.e., the point about which the mirror $M_{t}$ is turning, is sometimes equal to $h_{t}^{1}$ and sometimes to $h_{t}^{2}$. We can and will choose $h_{t}^{1}$ and $h_{t}^{2}$ so that the functions $t \rightarrow h_{t}^{j}$ are continuous up to the time when one of the hinges jumps to the third line segment. Suppose that $h_{0}^{1} \in I_{1}$ and $h_{0}^{2} \in I_{2}$, and both points are close to $z_{1}$; if they are not, the argument requires only minor modifications. It is possible that $h_{t}^{1}$ will not move until time $t_{1}$ when $M_{t_{1}}$ is perpendicular to $\partial D$ at $h_{t_{1}}^{1}$. Then $h_{t}^{1}$ will start moving while $h_{t}^{2}$ will remain fixed at the position $h_{t_{1}}^{2}$ until time $t_{2}$ when $M_{t_{2}}$ is perpendicular to $\partial D$ at $h_{t_{1}}^{2}=h_{t_{2}}^{2}$. The effect of these motions is that $h_{t_{2}}^{1}$ is at a greater distance from $z_{1}$ than $h_{0}^{1}$. If the same type of motions are repeated again, then $h_{t}^{1}$ will move away from $z_{1}$ by even greater distance. Since this process cannot be continued indefinitely, either $h_{t}^{1}$ or $h_{t}^{2}$ must hit one of the
vertices $z_{2}$ or $z_{3}$. Before this happens, either $h_{t}^{1}$ or $h_{t}^{2}$ must reach $F_{2}$ or $F_{3}$. Suppose, for example, that $h_{t}^{1}$ hits $F_{3}$ first. Then it follows from the definition of $F_{3}$ that $h_{t}^{2}$ may slide along $I_{2}$ and reach $F_{2}$ while $h_{t}^{1}$ remains fixed at a point of $F_{3}$. Hence we may have $h_{t}^{1} \in F_{3}$ and $h_{t}^{2} \in F_{2}$, for some $t$, with positive probability.

Next we may suppose that $h_{t}^{1}$ does not move until the mirror is perpendicular to $\partial D$ at this point. The other point, $h_{t}^{2}$, will then move to the side $I_{3}$. By repeating the process discussed in the previous paragraph, the points $h_{t}^{1}$ and $h_{t}^{2}$ may move so that $h_{t}^{1} \in F_{1}$ and $h_{t}^{2} \in F_{2}$. At some future times we may have, in succession, $h_{t}^{1} \in F_{1}$ and $h_{t}^{2} \in F_{3}$, and then, $h_{t}^{1} \in F_{2}$ and $h_{t}^{2} \in F_{3}$.

We now discuss the case when $M_{0}$ passes through a vertex. For example, assume that $h_{0}^{1}=z_{1}$. If the mirror is not perpendicular to the opposite side $I_{3}$, then it can turn around the point $h_{0}^{2}$ and it will no longer pass through a vertex. Suppose that the mirror is perpendicular to $I_{3}$. We have assumed in Theorem 3.7 (ii) that all angles of $D$ are different, so the angles formed by $M_{0}$ with $I_{1}$ and $I_{2}$ are not equal. It follows that when $X_{t}$ or $Y_{t}$ reflect on $I_{1}$ or $I_{2}$, the mirror will turn around $h_{0}^{1}=z_{1}$ and it will no longer be perpendicular to $I_{3}$.

Consider any point $z \in F_{1} \cap I_{1}$ and an orthonormal coordinate system $C S(z)$ in which $I_{1}$ lies on the horizontal axis, $z$ is the origin, and $D$ lies in the upper half-plane. Then choose $\rho_{3}>0$ such that for all $z \in F_{1} \cap I_{1}$ we have (viewed in the $C S(z)$ coordinate system) $\hat{B}(z)=B\left(\left(3 \rho_{3}, 3 \rho_{3}\right), 2 \rho_{3}\right) \subset D$ and $B\left(\left(-3 \rho_{3}, 3 \rho_{3}\right), 2 \rho_{3}\right) \subset D$.

Let $\hat{\theta}_{M}(t)$ be the angle formed by $M_{t}$ and the line containing $I_{1}$; let $\hat{h}_{t}$ be the intersection point of $M_{t}$ and $I_{1}$, with the convention that $\hat{h}_{t}=z_{1}$ if $M_{t} \cap I_{1}=\emptyset$; and let $\hat{X}_{t}$ be the position of $X_{t}$ expressed in terms of the $C S\left(\hat{h}_{t}\right)$ coordinate system.

Our argument has shown that for any starting points $x$ and $y$ for $X$ and $Y$ there is positive probability that at time $t=1$ the mirror $M_{1}$ passes through $F_{1} \cap I_{1}$ and is perpendicular to $I_{1}$, and $\hat{X}_{1} \in \hat{B}\left(\hat{h}_{1}\right)$. Hence, the event $A_{1}=$ $\left\{\hat{\theta}_{M}(1)=\pi / 2\right\}$ has a positive probability. Moreover, given the event $A_{1}, \hat{h}_{1}$ has a strictly positive density on $F_{1} \cap I_{1}$. Given $A_{1}$ and $\hat{h}_{1}$, the density of $\hat{X}_{1}$ is strictly positive on $\hat{B}\left(\hat{h}_{1}\right)$.

We sketch a proof, using compactness, that there exist lower and upper bounds for the densities of $\hat{h}_{1}$ and $\hat{X}_{1}$, uniform in $x, y \in \mathcal{D}(a)$ for any fixed small $a>0$. Let $\psi^{x, y}(v)$ be the density of $\hat{h}_{1}$ restricted to $F_{1} \cap I_{1}$ (the proof for the density of $\hat{X}_{1}$ is analogous and so it is omitted). Suppose that there exist $v_{0} \in F_{1} \cap I_{1}$ and a sequence $\left(x_{k}, y_{k}\right) \in \mathcal{D}(a)$ such that $\psi^{x_{k}, y_{k}}\left(v_{0}\right) \rightarrow 0$ as $k \rightarrow \infty$. By compactness, we may suppose that $x_{k} \rightarrow x_{\infty}$ and $y_{k} \rightarrow y_{\infty}$. Note that we must have $\left(x_{\infty}, y_{\infty}\right) \in \mathcal{D}(a)$. Now going back to our argument, it is not hard to see that the infimum of $\psi^{x, y}\left(v_{0}\right)$ taken over $(x, y)$ in a neighborhood of $\left(x_{\infty}, y_{\infty}\right)$ must be strictly positive. The crucial observation here is that the distance between $x_{\infty}$ and $y_{\infty}$ is strictly positive. This gives us the desired contradiction.

We fix some small $a>0$ and $x_{0}, y_{0} \in \mathcal{D}(a)$ and take $\nu$ to be the restric-
tion of the distribution of $\left(X_{1}, Y_{1}\right)$ to the event $A_{1} \cap\left\{\hat{h}_{1} \in F_{1} \cap I_{1}\right\}$, given $\left(X_{0}, Y_{0}\right)=\left(x_{0}, y_{0}\right)$. Chosen in this way, $\nu$ satisfies the condition in the lemma.

## Lemma 3.9

(i) There exist $a, \alpha, c_{1}>0$ such that for any $\varepsilon \in(0, a)$ and $(x, y) \in \mathcal{D}(\varepsilon)$, we have

$$
\mathbb{P}\left[T(\mathcal{D}(a))<\tau \mid X_{0}=x, Y_{0}=y\right] \geq c_{1} \varepsilon^{\alpha} .
$$

(ii) There exists $c_{2}<\infty$ such that for all $\varepsilon>0$ and $(x, y) \in \hat{\mathcal{D}}(\varepsilon)$,

$$
\mathbb{E}\left[T\left(\hat{\mathcal{D}}^{c}(\varepsilon)\right) \mid X_{0}=x, Y_{0}=y\right] \leq c_{2} \varepsilon^{2}
$$

## Proof :

(i) Let $\rho_{0}$ be so small that any disc of radius $100 \rho_{0}$ can intersect at most 2 sides of the triangle $D$. Let $A(\rho)$ be the event that the process $(X, Y)$ will hit $\mathcal{D}(2 \rho)$ before exiting $\mathcal{D}$ and, moreover, this will happen before $X$ or $Y$ move more than $4 \rho$ away from their starting points. For a fixed $\rho \leq \rho_{0}$, let $p=p(\rho)$ be the infimum of $\mathbb{P}\left(A(\rho) \mid X_{0}=x, Y_{0}=y\right)$, evaluated over all $x$ and $y$ with $|x-y|=\rho$. By the arguments presented in the proof of Lemma 3.8, we know that $p(\rho)>0$. We will argue that $p(\rho)=p\left(\rho_{1}\right)$ for some $\rho_{1}>0$ and all $\rho<\rho_{1}$.
Consider a vertex $z_{1}$ of the triangle $D$. Then

$$
p_{1}(\rho)=\inf _{\substack{|x-y|=\rho \\\left|x-z_{1}\right| \leq 10 \rho}} \mathbb{P}\left(A(\rho) \mid X_{0}=x, Y_{0}=y\right)
$$

depends only on the angle at the vertex $z_{1}$, because neither process $X$ nor $Y$ can hit the side of $D$ opposite to $z_{1}$ before moving more than $4 \rho$ units away from its starting point. By scaling, we obtain $p_{1}(\rho)=p\left(\rho_{1}\right)$ for some $\rho_{1}>0$ and all $\rho<\rho_{1}$. The same argument applies to the neighborhoods of the other two vertices, and to the points of $D$ more than $9 \rho$ units away from any vertex.
By repeatedly applying the strong Markov property at the hitting times of $\mathcal{D}\left(2^{j} \rho\right)$, we see that

$$
\mathbb{P}\left(T(\mathcal{D}(a))<\tau \mid X_{0}=x, Y_{0}=y\right) \geq c_{3} p^{k}
$$

for $(x, y) \in \mathcal{D}\left(a / 2^{k}\right)$. This can be easily rewritten as the estimate in part (i) of the lemma.
(ii) Recall the process $\rho=|X-Y|$ from the proof of Theorem 3.7 (i). It is the sum of Brownian motion (with variance twice the standard one) and a non-increasing process. Let $T_{\rho}(0)$ be the hitting time of 0 for $\rho_{t}$. By
comparing $\rho$ with the Brownian motion with diffusion rate 2 , we obtain the following estimate. There exists $p>0$ such that for all $\varepsilon>0$ and $(x, y) \in \hat{\mathcal{D}}(\varepsilon)$,

$$
\mathbb{P}\left(T_{\rho}(0)<\varepsilon^{2} \mid X_{t}=x, Y_{t}=y\right)>p
$$

which clearly implies

$$
\mathbb{P}\left(T\left(\hat{\mathcal{D}}^{c}(\varepsilon)\right) \geq \varepsilon^{2} \mid X_{t}=x, Y_{t}=y\right) \leq 1-p
$$

By the Markov property,

$$
\mathbb{P}\left(T\left(\hat{\mathcal{D}}^{c}(\varepsilon)\right) \geq k \varepsilon^{2} \mid X_{t}=x, Y_{t}=y\right) \leq(1-p)^{k}
$$

and so

$$
\mathbb{E}\left(T\left(\hat{\mathcal{D}}^{c}(\varepsilon)\right) \mid X_{t}=x, Y_{t}=y\right) \leq \sum_{k \geq 1} k \varepsilon^{2}(1-p)^{k}=c_{2} \varepsilon^{2}
$$

The following lemma is almost the same as [4, Lemma 5.1]. We reproduce that result here as many details of the original proof have to be changed. The intuitive meaning of the lemma is that if we condition $X$ and $Y$ on not coupling before time $s$, then the processes are likely to move apart for a considerable distance at time $s$. Hence, Lemma 3.10 below is a version of the parabolic boundary Harnack principle for the process $(X, Y)$.

Recall the mirror coupling $(X, Y)$. It will be convenient to write $Z=(X, Y)$ as the separate components of $Z$ will play no role in Lemma 3.10. The state space of $Z$ is $\bar{D} \times \bar{D}$. Recall that $\mathcal{D}=\{(x, y) \in \bar{D} \times \bar{D}: x \neq y\}, \mathcal{D}(\varepsilon)=\{(x, y) \in$ $\bar{D} \times \bar{D}:|x-y| \geq \varepsilon\}$ and $\hat{\mathcal{D}}(\varepsilon)=\mathcal{D} \backslash \mathcal{D}(\varepsilon)$.

In the following lemma, $\mathbb{P}^{z}$ and $\mathbb{E}^{z}$ will denote the distribution of $Z$ starting from $z$ and the corresponding expectation. Conditioning by a harmonic function $h$ will be reflected in the notation by writing $\mathbb{P}_{h}^{z}$ and $\mathbb{E}_{h}^{z}$. See [20] for the discussion of conditioned Brownian motion and [44] for conditioning of general Markov processes.

We will denote the space-time counterpart of $Z$ by $V$. More precisely, if $Z$ has law $\mathbb{P}^{z}$, then the law of the space-time process $\mathrm{V}=\left\{V_{t}=\left(Z_{t}, s-t\right), t \geq 0\right\}$ will be denoted $\mathbb{P}^{z, s}$. The distribution of space-time process conditioned by a parabolic function $g$ will be denoted $P_{g}^{z, s}$. By abuse of notation, $T(A)$ will denote the first hitting time of $A$ for $V$ as well as for $Z$.

Lemma 3.10 There exist $\varepsilon, c, u>0$ such that for all $z \in \mathcal{D}$,

$$
\mathbb{P}^{z}\left(Z_{u} \in \mathcal{D}(\varepsilon), \tau>u\right) \geq c \mathbb{P}^{z}(\tau>u) .
$$

Proof : Fix some small $\varepsilon_{0}>0$ such that $M=\mathcal{D}\left(\varepsilon_{0}\right)$ contains a non-empty open ball and let $\mathcal{D}_{1}=\mathcal{D} \backslash M$. Let $h(z)=\mathbb{P}^{z}(T(M)<\tau)$ and $U_{k}=\left\{x \in D_{1}\right.$ : $\left.h(x) \in\left[2^{k-1}, 2^{k}\right]\right\}$ for integer $k$.

By Lemma 3.9 (i), $U_{k} \subset \hat{\mathcal{D}}\left(c_{1} 2^{k / \alpha}\right)$, for some $c_{1}, \alpha>0$. Then Lemma 3.9 (ii) shows that $\sup _{z \in U_{k}} \mathbb{E}^{z}\left(T\left(U_{k}^{c}\right)\right) \leq c_{2} 2^{2 k / \alpha}$. It follows that

$$
\sum_{k=0}^{\infty} \sup _{z \in U_{-k}} \mathbb{E}^{z}\left(T\left(U_{-k}^{c}\right)\right)<\infty
$$

An argument of Chung [13] (see also [4]) shows that for suitable $c_{3}$,

$$
c_{3} \sum_{k=0}^{\infty} \sup _{z \in U_{-k}} \mathbb{E}^{z}\left(T\left(U_{-k}^{c}\right)\right)
$$

is an upper bound for $\mathbb{E}_{h}^{z}\left(T\left(\mathcal{D}_{1}^{c}\right)\right)$. It follows that for a suitable $u>0$ and every $z \in \mathcal{D}$,

$$
\begin{equation*}
\left.\mathbb{P}_{h}^{z}\left(T\left(\mathcal{D}_{1}^{c}\right)\right)<u / 4\right)>1 / 2 . \tag{3.11}
\end{equation*}
$$

Recall the discussion of space-time processes before the statement of the lemma. The function

$$
(z, t) \mapsto g(z, t)=\mathbb{P}^{z}(\tau>t)
$$

is parabolic in $\mathcal{D} \times[0, \infty)$ with boundary values 1 on $\mathcal{D} \times\{0\}$ and 0 otherwise.
Let $g_{1}$ be a parabolic function in $\mathcal{D} \times[0, \infty)$ which has the same boundary values as $g$ except that $g_{1}(z, 0)$ is changed from 1 to $\delta$ for $z \in \mathcal{D}_{1}$, where $\delta \in(0,1)$ will be chosen later. Now we will estimate $g_{1}$ on $\mathcal{D} \times[u / 2, u]$.

It is easy to see that $g_{1}(z, s)>c_{4}$ for all $z \in M$ and $s \in[u / 4, u]$. We obviously have $h(y) \leq 1$ for all $y$. Let $h(x, s)=h(x)$. For $x \in \mathcal{D}_{1}$ and $s \geq u / 2$ we have, by (3.11),

$$
\begin{aligned}
g_{1}(x, s) & \geq \int_{\substack{t \in[u / 4, u] \\
y \in \mathcal{D} 1}} g_{1}(y, t) \mathbb{P}^{x, s}\left(T\left(\mathcal{D}_{1}^{c}\right) \in d t, X\left(T\left(\mathcal{D}_{1}^{c}\right)\right) \in d y\right) \\
& =\int_{\substack{t \in[u / 4, u] \\
y \in \mathcal{D} 1}} \frac{h(x, s)}{h(y, t)} \frac{h(y, t)}{h(x, s)} g_{1}(y, t) \mathbb{P}^{x, s}\left(T\left(\mathcal{D}_{1}^{c}\right) \in d t, X\left(T\left(\mathcal{D}_{1}^{c}\right)\right) \in d y\right) \\
& =\int_{\substack{t \in[u / 4, u] \\
y \in \mathcal{D} \mathcal{D}_{1}}} \frac{h(x, s)}{h(y, t)} g_{1}(y, t) \mathbb{P}_{h}^{x, s}\left(T\left(\mathcal{D}_{1}^{c}\right) \in d t, X\left(T\left(\mathcal{D}_{1}^{c}\right)\right) \in d y\right) \\
& \geq \int_{\substack{t \in[u / 4, u] \\
y \in \mathcal{D} 1}} h(x, s) c_{4} \mathbb{P}_{h}^{x, s}\left(T\left(\mathcal{D}_{1}^{c}\right) \in d t, X\left(T\left(\mathcal{D}_{1}^{c}\right)\right) \in d y\right) \\
& =h(x, s) c_{4} \mathbb{P}_{h}^{x, s}\left(T\left(\mathcal{D}_{1}^{c}\right) \in[u / 4, s]\right) \\
& \geq h(x, s) c_{4} / 2=c_{5} h(x, s)=c_{5} h(x) .
\end{aligned}
$$

Let

$$
\begin{gathered}
W_{k}=\left\{(z, s): g_{1}(z, s) \in\left[2^{k}, 2^{k+1}\right], s \in[u / 2, u]\right\} \\
W=\bigcup_{k=-\infty}^{k_{1}} W_{k}
\end{gathered}
$$

where $k_{1}<0$ will be chosen later. If $2^{-m}<c_{5}$ then $W_{k} \subset U_{k+m} \times[u / 2, u]$. Using the estimate of Chung [13] we obtain for small $k_{1}$ and all $z \in \mathcal{D}$,

$$
\begin{aligned}
\mathbb{E}_{g_{1}}^{z, u}\left(T\left(W^{c}\right)\right) & \leq c_{6} \sum_{k=-\infty}^{k_{1}} \sup _{(y, s) \in W_{k}} \mathbb{E}^{y, s} T\left(W_{k}^{c}\right) \\
& \leq c_{6} \sum_{k=-\infty}^{k_{1}} \sup _{(y, s) \in U_{k+m}} \mathbb{E}^{y, s} T\left(U_{k+m}^{c}\right)<\infty
\end{aligned}
$$

Choose $k_{1}$ so small that for any $z \in \mathcal{D}$,

$$
\begin{equation*}
\mathbb{E}_{g_{1}}^{z, u} T\left(W^{c}\right)<u / 8 \tag{3.12}
\end{equation*}
$$

Let

$$
Q=\left\{(x, s): g_{1}(x, s) \geq 2^{k_{1}}, s \in[u / 2, u]\right\}
$$

Since the $g_{1}$-process cannot exit $\mathcal{D} \times[0, \infty)$ through $\partial \mathcal{D} \times[0, \infty)$, (3.12) implies

$$
\begin{equation*}
P_{g_{1}}^{z, u}(T(Q)>u / 4)<1 / 2 \tag{3.13}
\end{equation*}
$$

Now let $\delta=2^{k_{1}-1}$. Since $0 \leq g_{1} \leq 1$, the process $g_{1}\left(V_{t}\right)$ is a martingale under $\mathbb{P}^{z, s}$, and $g_{1}(z, s) \geq 2^{k_{1}}$ for $(z, s) \in Q$, we see that there is at least $2^{k_{1}-1} / 2$ chance that $V$ under $\mathbb{P}^{z, s}$ will hit $M \times\{0\}$ before hitting any other part of $\partial(\mathcal{D} \times[0, \infty))$. Thus we have for $(z, s) \in Q$,

$$
\begin{aligned}
& \mathbb{P}_{g_{1}}^{z, s}\left[V_{s} \in M \times\{0\}\right]= \\
& \int_{M}\left(g_{1}(y, 0) / g_{1}(z, s)\right) \mathbb{P}^{z, s}\left[V_{s} \in d y, T(\partial(\mathcal{D} \times[0, \infty)))=s\right] \\
& \geq \int_{M} \mathbb{P}^{z, s}\left[V_{s} \in d y, T(\partial(\mathcal{D} \times[0, \infty)))=s\right] \\
& \geq 2^{k_{1}-1} / 2 .
\end{aligned}
$$

This and (3.13) yield, by the strong Markov property, for all $z \in \mathcal{D}$,

$$
\mathbb{P}_{g_{1}}^{z, u}\left[V_{u} \in M \times\{0\}\right] \geq c_{7}>0
$$

The ratio of $g$ and $g_{1}$ is bounded away from 0 and $\infty$ on the accessible boundary of $\mathcal{D} \times[0, \infty)$, so

$$
P_{g}^{z, u}\left[V_{u} \in M \times\{0\}\right] \geq c_{8}>0
$$

for all $z \in \mathcal{D}$. This is equivalent to the statement in the lemma.

Proof of Theorem 3.7 (ii): We start by constructing a coupling of couplings. More precisely, we will construct processes $(X, Y)$ and $(\tilde{X}, \tilde{Y})$ such that each one of them is a mirror coupling of reflected Brownian motions in $D$. Hence, each of these processes is Markov. The two processes will also form a coupling, but the combined process $((X, Y),(\tilde{X}, \tilde{Y}))$ will not be Markov since the coupling will fail to have the co-adapted property.

Let $\tau$ and $\tilde{\tau}$ denote the coupling times for $(X, Y)$ and $(\tilde{X}, \tilde{Y})$, respectively.
Fix some $a_{1}, c_{1}, u_{1}>0$ which satisfy Lemma 3.10 in place of $\varepsilon, c$ and $u$. Find $a_{2} \in\left(0, a_{1}\right), c_{2}>0$, a set $\mathcal{D}_{1} \subset \mathcal{D}\left(a_{2}\right)$ and a measure $\nu$ supported by $\mathcal{D}_{1}$ which satisfy Lemma 3.8. By Lemmas 3.8 and 3.10 , for $A \subset \mathcal{D}_{1}$ and any $x, y \in \bar{D}$,

$$
\begin{aligned}
& \mathbb{P}\left(\left(X_{u_{1}+1}, Y_{u_{1}+1}\right) \in A \mid X_{0}=x, Y_{0}=y\right) \\
& =\int_{\mathcal{D}} \mathbb{P}\left(\left(X_{1}, Y_{1}\right) \in A \mid X_{0}=x^{\prime}, Y_{0}=y^{\prime}\right) \times \\
& \quad \times \mathbb{P}\left(\left(X_{u_{1}}, Y_{u_{1}}\right) \in\left(d x^{\prime}, d y^{\prime}\right) \mid X_{0}=x, Y_{0}=y\right) \\
& \geq \int_{\mathcal{D}_{1}} \mathbb{P}\left(\left(X_{1}, Y_{1}\right) \in A \mid X_{0}=x^{\prime}, Y_{0}=y^{\prime}\right) \times \\
& \quad \times \mathbb{P}\left(\left(X_{u_{1}}, Y_{u_{1}}\right) \in\left(d x^{\prime}, d y^{\prime}\right) \mid X_{0}=x, Y_{0}=y\right) \\
& \geq c_{2} \int_{\mathcal{D}_{1}} \nu(A) \mathbb{P}\left(\left(X_{u_{1}}, Y_{u_{1}}\right) \in\left(d x^{\prime}, d y^{\prime}\right) \mid X_{0}=x, Y_{0}=y\right) \\
& \geq c_{1} c_{2} \nu(A) \mathbb{P}\left(\tau>u_{1} \mid X_{0}=x, Y_{0}=y\right) \\
& \geq c_{1} c_{2} \nu(A) \mathbb{P}\left(\tau>u_{1}+1 \mid X_{0}=x, Y_{0}=y\right) .
\end{aligned}
$$

Let $u_{2}=u_{1}+1$. The last formula implies that

$$
\begin{equation*}
\mathbb{P}\left(\left(X_{u_{2}}, Y_{u_{2}}\right) \in A \mid \tau>u_{2}, X_{0}=x, Y_{0}=y\right) \geq \nu_{1}(A)=c_{1} c_{2} \nu(A) . \tag{3.14}
\end{equation*}
$$

Consider any $x, y, \tilde{x}$ and $\tilde{y}$ in $\bar{D}$. On a single probability space, we will construct processes $\left(X_{t}, Y_{t}\right)$ and $\left(\tilde{X}_{t}, \tilde{Y}_{t}\right)$ starting from $(x, y)$ and $(\tilde{x}, \tilde{y})$ respectively. In view of (3.14) we may construct random vectors $\mathcal{V}=\left(\mathcal{V}_{1}, \mathcal{V}_{2}\right)$ and $\tilde{\mathcal{V}}=\left(\tilde{\mathcal{V}}_{1}, \tilde{\mathcal{V}}_{2}\right)$ such that for $A \subset \mathbb{R}^{4}$,

$$
\mathbb{P}(\mathcal{V} \in A)=\mathbb{P}\left(\left(X_{u_{2}}, Y_{u_{2}}\right) \in A \mid X_{0}=x, Y_{0}=y\right)
$$

and

$$
\mathbb{P}(\tilde{\mathcal{V}} \in A)=\mathbb{P}\left(\left(\tilde{X}_{u_{2}}, \tilde{Y}_{u_{2}}\right) \in A \mid \tilde{X}_{0}=\tilde{x}, \tilde{Y}_{0}=\tilde{y}\right) ;
$$

moreover, for $A \subset \mathcal{D}_{1}$,

$$
\begin{equation*}
\mathbb{P}\left(\mathcal{V}=\tilde{\mathcal{V}} \in A \mid \mathcal{V}_{1} \neq \mathcal{V}_{2}, \tilde{\mathcal{V}}_{1} \neq \tilde{\mathcal{V}}_{2}\right) \geq \nu_{1}(A) . \tag{3.15}
\end{equation*}
$$

Now we set

$$
\begin{aligned}
\left(X_{0}, Y_{0}\right)=(x, y), & \left(\tilde{X}_{0}, \tilde{Y}_{0}\right)=(\tilde{x}, \tilde{y}), \\
\left(X_{u_{2}}, Y_{u_{2}}\right) & =\mathcal{V}, \\
\left(\tilde{X}_{u_{2}}, \tilde{Y}_{u_{2}}\right) & =\tilde{\mathcal{V}}
\end{aligned}
$$

Next we construct $\left\{\left(X_{t}, Y_{t}\right), t \in\left[0, u_{2}\right]\right\}$ and $\left\{\left(\tilde{X}_{t}, \tilde{Y}_{t}\right), t \in\left[0, u_{2}\right]\right\}$ by adding bridges between the endpoints of the trajectories in such a way that each of these processes is a mirror coupling of reflected Brownian motions in $D$. Let $Q_{u_{2}}^{x, y, \tilde{x}, \tilde{y}}$ denote the distribution of $\left\{\mathcal{Z}_{t}=\left(\left(X_{t}, Y_{t}\right),\left(\tilde{X}_{t}, \tilde{Y}_{t}\right)\right), t \in\left[0, u_{2}\right]\right\}$. We inductively define $Q_{(k+1) u_{2}}^{x, y, \tilde{x}, \tilde{y}}$ for $k=1,2,3, \ldots$, by the following "Markov-like property" formula (the Markov property does not extend to other times besides $k u_{2}$ ),

$$
\begin{gather*}
Q_{(k+1) u_{2}}^{x, y, \tilde{x}, \tilde{y}} \quad\left(\left\{\mathcal{Z}_{t}, t \in\left[0, k u_{2}\right]\right\} \in A_{1},\left\{\mathcal{Z}_{t}, t \in\left[k u_{2},(k+1) u_{2}\right]\right\} \in A_{2}\right) \\
=\int Q_{k u_{2}}^{x, y, \tilde{x}, \tilde{y}}\left(\left\{\mathcal{Z}_{t}, t \in\left[0, k u_{2}\right]\right\} \in A_{1}, \mathcal{Z}_{k u_{2}} \in(d v, d z, d \tilde{v}, d \tilde{z})\right) \\
\times Q_{u_{2}}^{v, z, \tilde{v}, \tilde{z}}\left(\left\{\mathcal{Z}_{t}, t \in\left[0, u_{2}\right]\right\} \in A_{2}^{k u_{2}}\right), \tag{3.16}
\end{gather*}
$$

for all $A_{1} \subset C\left(\left[0, k u_{2}\right], \mathbb{R}^{8}\right)$ and $A_{2} \subset C\left(\left[k u_{2},(k+1) u_{2}\right], \mathbb{R}^{8}\right)$. Here $A_{2}^{k u_{2}}$ is the family of functions in $A_{2}$, shifted to the left by $k u_{2}$ units. Note that the measure $Q_{(k+1) u_{2}}^{x, y, \tilde{x}, \tilde{y}}$ is uniquely defined if we specify its value on cylinders of the form $A_{1} \times A_{2}$, as in (3.16). We define $Q_{\infty}^{x, y, \tilde{x}, \tilde{y}}$, i.e., the distribution of $\left\{\left(\left(X_{t}, Y_{t}\right),\left(\tilde{X}_{t}, \tilde{Y}_{t}\right)\right), t \in[0, \infty)\right\}$ using the measures $Q_{k u_{2}}^{x, y, \tilde{x}, \tilde{y}}$ and consistency, in the obvious way.

Let $\tau^{*}=\inf \left\{t:\left(X_{t}, Y_{t}\right)=\left(\tilde{X}_{t}, \tilde{Y}_{t}\right)\right\}$. It follows from (3.15) that

$$
Q_{\infty}^{x, y, \tilde{x}, \tilde{y}}\left(\tau^{*}>u_{2} \mid \tau>u_{2}, \tilde{\tau}>u_{2}\right) \leq p_{1}=1-\nu_{1}\left(\mathcal{D}_{1}\right)<1
$$

for all $x, y, \tilde{x}, \tilde{y} \in \bar{D}$. By the "Markov-like property" (3.16),

$$
Q_{\infty}^{x, y, \tilde{x}, \tilde{y}}\left(\tau^{*}>k u_{2} \mid \tau>k u_{2}, \tilde{\tau}>k u_{2}\right) \leq p_{1}^{k}
$$

Hence, for some $c_{3}>0$, all $t \geq u_{2}$ and all $x, y, \tilde{x}, \tilde{y} \in \bar{D}$,

$$
\begin{equation*}
Q_{\infty}^{x, y, \tilde{x}, \tilde{y}}\left(\tau^{*}>t \mid \tau>t, \tilde{\tau}>t\right) \leq e^{-c_{3} t} \tag{3.17}
\end{equation*}
$$

The rest of the proof is very similar to the end of the proof of Theorem 2.6 (i). Consider $x, y \in \bar{D}$ with $\varphi_{2}(x) \neq \varphi_{2}(y)$. Recall (3.5);

$$
\mathbb{P}\left(\tau>t \mid X_{0}=x, Y_{0}=y\right) \quad \leq \quad c(x, y) e^{-\mu t}, \quad t \geq 0
$$

for some $\mu>0$. In the following we will use the generic notation $\mathbb{P}$ for probability but we will assume that the process $\left\{\left(\left(X_{t}, Y_{t}\right),\left(\tilde{X}_{t}, \tilde{Y}_{t}\right)\right), t \in[0, \infty)\right\}$ has the distribution $Q_{\infty}^{x, y, y, x}$. Hence, $X_{0}=\tilde{Y}_{0}=x$ and $Y_{0}=\tilde{X}_{0}=y$. Then we can use (3.17) to show, for $A \subset D$, and large $t$,

$$
\begin{aligned}
& \left|\int_{A} p(t, x, z) d z-\int_{A} p(t, y, z) d z\right| \\
& =\left|\mathbb{P}\left(X_{t} \in A \mid X_{0}=x\right)-\mathbb{P}\left(Y_{t} \in A \mid Y_{0}=y\right)\right| \\
& =\left|\mathbb{P}\left(X_{t} \in A, t<\tau \mid X_{0}=x, Y_{0}=y\right)-\mathbb{P}\left(Y_{t} \in A, t<\tau \mid X_{0}=x, Y_{0}=y\right)\right| \\
& =\mid \mathbb{P}\left(X_{t} \in A \mid t<\tau, X_{0}=x, Y_{0}=y\right) \mathbb{P}\left(t<\tau \mid X_{0}=x, Y_{0}=y\right)
\end{aligned}
$$

$$
\begin{aligned}
& \quad-\mathbb{P}\left(Y_{t} \in A \mid t<\tau, X_{0}=x, Y_{0}=y\right) \mathbb{P}\left(t<\tau \mid X_{0}=x, Y_{0}=y\right) \mid \\
& \leq\left|\mathbb{P}\left(X_{t} \in A \mid t<\tau, X_{0}=x, Y_{0}=y\right)-\mathbb{P}\left(Y_{t} \in A \mid t<\tau, X_{0}=x, Y_{0}=y\right)\right| \times \\
& \quad \times c(x, y) e^{-\mu t} \\
& =\left|\mathbb{P}\left(X_{t} \in A \mid t<\tau, X_{0}=x, Y_{0}=y\right)-\mathbb{P}\left(\tilde{X}_{t} \in A \mid t<\tilde{\tau}, \tilde{X}_{0}=y, \tilde{Y}_{0}=x\right)\right| \\
& \quad \times c e^{-\mu t} \\
& \leq \mid \mathbb{P}\left(X_{t} \in A, t<\tau^{*} \mid t<\tilde{\tau}, t<\tau, X_{0}=x, Y_{0}=y\right) \\
& \quad-\mathbb{P}\left(\tilde{X}_{t} \in A, t<\tau^{*} \mid t<\tilde{\tau}, t<\tau, \tilde{X}_{0}=y, \tilde{Y}_{0}=x\right) \mid c e^{-\mu t} \\
& \leq \mathbb{P}\left(t<\tau^{*} \mid t<\tilde{\tau}, t<\tau, X_{0}=x, \tilde{X}_{0}=y\right) \mid c e^{-\mu t} \quad \leq e^{-c_{3} t} c e^{-\mu t}
\end{aligned}
$$

This and (3.6) show that $\mu_{2} \geq \mu+c_{3}>\mu$. Thus the mirror coupling for reflected Brownian motion in a triangle with distinct acute angles is not efficient.

## 4 Some further examples

This section contains a selection of rather informal examples, giving some indication of how far the results on mirror couplings in triangles, presented in the last section, can be generalized to mirror couplings in other planar sets. The motivation for our efforts comes from two different sources. First, mirror couplings have been used to estimate the spectral gap for some diffusions [48] so it is a natural question how sharp those estimates are. This we cannot say but our examples may be a source of inspiration for future research in this direction. Second, the mirror coupling techniques developed for this project have already been applied $[3,6,5]$ to a problem on "hot spots." Hence, the techniques seem to have some interest beyond the efficiency of Markovian couplings.

The reader might have noticed that the assumption that $D$ is a triangle, adopted in Section 3, does not play a major role in the arguments. The following example makes this point explicit.

Example 4.1 Fig. 4 shows a convex polygonal domain $D_{1}$ whose boundary is naturally divided into "upper" and "lower" parts. The angles between line segments in the upper part of $\partial D_{1}$ and those in the lower part are less than $\pi / 2$. The arguments presented in Section 3 carry over to this case and it is easy to see that both synchronous and mirror couplings for reflected Brownian motion in $D_{1}$ are efficient.
[Figure \& about here.]
On the other hand, Theorem 3.7 (ii) can be also generalized to some other domains besides triangles. The domain $D_{2}$ illustrated in Fig. 5 is an acute triangle whose corners have been cut. The mirror coupling can be proved to be inefficient in $D_{2}$ just as in the case of a triangle with acute angles.
[Figure 5 about here.]

What about non-convex domains? Convexity is used in Theorem 3.7 (i) to show that the distance between the processes $X$ and $Y$ is a Brownian motion plus a process which always pushes $X$ and $Y$ towards each other. This is true only in convex domains. However it can be circumvented, at a price of obscure conditions and tedious details.

Example 4.2 It is easy to check that the mirror will never turn more than the angle $\pi$ in some non-convex domains, for example in the domain $D_{3}$ in Fig. 6. The angles between any two line segments in $\partial D_{3}$ are less than $\pi / 2$. We believe that the mirror coupling is efficient in $D_{3}$ but the proof of Theorem 3.7 (i) does not completely apply in this case because of the lack of convexity as indicated above. However we believe that one can circumvent the need for convexity for this kind of example, albeit with the need for more involved arguments.
[Figure 6 about here.]
Our next example addresses the question of what happens if the domain has smooth boundary, rather than polygonal boundary. The construction of a mirror coupling in such a domain has to proceed along different lines than that presented in Section 3, which works only for polygonal domains. However, the construction does not present major problems; an example in a similar context is to be found in [29].

Example 4.3 The domain $D_{4}$ in Fig. 7 has a piecewise smooth boundary. If we consider two tangent lines to $\partial D_{4}$, one to the upper part of $\partial D_{4}$ and the other tangent to the lower part then they form an angle less than $\pi / 2$. In a domain $D_{4}$, the hinge will be the point of the intersection of the mirror and the tangent line to $\partial D_{4}$ at the point where one of the processes is reflecting from the boundary. Hence, the hinge will move not only by jumps but also in a continuous fashion. The general qualitative behavior of the mirror movement does not change in a fundamental way, however, from the polygonal domain case. Hence, the mirror cannot turn more than $\pi$ in $D_{4}$. This is all we have to know to prove that the mirror coupling is efficient in $D_{4}$.
[Figure 7 about here.]
It is also possible to make some progress if a symmetry is present:
Example 4.4 In [3] it is proved that the mirror cannot turn more than $\pi$ in a convex domain $D_{5}$ if we assume in addition that
(1) $D_{5}$ has a line of symmetry $S$ which intersects $\partial D_{5}$ at $x$ and $y$,
(2) $|x-z| \vee|y-z|<|x-y|$ for all $z \in \bar{D}_{5} \backslash\{x, y\}$, and
(3) for all $r>0$ the sets $\partial B(x, r) \cup D$ and $\partial B(y, r) \cup D$ are connected.

Just as in Example 4.2, we believe that the mirror coupling is efficient in domains $D_{5}$ satisfying these assumptions but the proof given in Section 3 would have to be modified. In the present case, we cannot claim that $Y_{t}^{1}-X_{t}^{1}>$ $\alpha\left|Y_{t}^{1}-X_{t}^{1}\right|$ for some $\alpha>0$. This property holds if we impose some more assumptions on the slope of $\partial D_{5}$.

It should be noted that our methods cannot decide whether the mirror coupling is efficient for reflected Brownian motion in the domains considered in [3, Theorem 1.3 (A1)]. Those domains are assumed to have two perpendicular axes of symmetry. The proof of that result is based on the behavior of the mirror coupling for the reflected Brownian motion with absorption on one of the axes of symmetry. Hence, the technique does not directly apply to reflected Brownian motion in the whole domain.

Finally we will provide some details of the mirror coupling behavior in the case when the domain $D$ is a disc. This highly symmetric case makes the analysis especially easy and complete.

Example 4.5 Recall from Example 4.3 that the hinge $h_{t}$ lies at the intersection of the mirror and the line tangent to the circle $\partial D$ where one of the processes $X$ or $Y$ is reflecting. A quick look at Fig. 8 should convince the reader that the mirror $M$ must move towards the center of the disc (i.e., its distance from the center can only decrease). Moreover, the points of intersection of the mirror with $\partial D$ can only move upwards in Fig. 8. These remarks follow from the fact that the effect of reflection is the counterclockwise motion of the mirror $M$ around the (instantaneous) hinge position $h$. If the mirror passes through the center at some time s, it will never change its position after time s, because after that time, the processes $X$ and $Y$ will reflect at the boundary of $D$ at the same time, until their coupling time. Hence, the mirror can never intersect $\partial D$ inside the (smaller) arc between $a$ and $b$ and likewise not between $a_{1}$ and $b_{1}$ (the antipodal points to $a$ and $b$ ). The process $Y$ can not start reflecting on $\partial D$ before the time when the mirror passes through the center. Hence, the mirror must hit the center of the disc before or at the same time when $Y$ hits the smaller arc between $a_{1}$ and $b_{1}$. Given these properties of the mirror coupling in a disc, it is not hard to prove that it is efficient.
[Figure 8 about here.]

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