# Curse-of-Complexity Attenuation in the Curse-of-Dimensionality-Free Method for HJB PDEs 

William M. McEneaney, Ameet Deshpande and Stephane Gaubert


#### Abstract

Recently, a curse-of-dimensionality-free method was developed for solution of Hamilton-Jacobi-Bellman partial differential equations (HJB PDEs) for nonlinear control problems, using semiconvex duality and max-plus analysis. The curse-of-dimensionality-free method may be applied to HJB PDEs where the Hamiltonian is given as (or well-approximated by) a pointwise maximum of quadratic forms. Such HJB PDEs also arise in certain switched linear systems. The method constructs the correct solution of an HJB PDE from a maxplus linear combination of quadratics. The method completely avoids the curse-of-dimensionality, and is subject to cubic computational growth as a function of space dimension. However, it is subject to a curse-of-complexity. In particular, the number of quadratics in the approximation grows exponentially with the number of iterations. Efficacy of such a method depends on the pruning of quadratics to keep the complexity growth at a reasonable level. Here we apply a pruning algorithm based on semidefinite programming. Computational speeds are exceptional, with an example HJB PDE in six-dimensional Euclidean space solved to the indicated quality in approximately 30 minutes on a typical desktop machine.


## I. Introduction

Dynamic programming is an extremely robust tool for solving nonlinear optimal control problems. In the case of deterministic optimal control, or in the case of deterministic games where one player's feedback is prespecified, the dynamic programming equation reduces to a Hamilton-JacobiBellman (HJB) PDE. The difficulty is that one must solve the HJB PDE.

Various approaches have been taken to solving the HJB PDE. The most common methods are grid-based methods (c.f., [6], [7]) which still suffer from the curse-ofdimensionality, as the number of grid points and computations grow exponentially with the space dimension. However, in recent years, entirely new classes of numerical methods for HJB PDEs have emerged (c.f., [8], [15], [2], [1], [12]). These methods exploit the max-plus linearity of the associated semigroup.

In the previous work of the first author [12], [14], a new method based on above semigroup linearity was proposed for certain nonlinear HJB PDEs, and this method was free from the curse-of-dimensionality. In fact, the computational

[^0]growth in state-space dimension is cubic. However, there is exponential computational growth in a certain measure of complexity of the Hamiltonian. Under this measure, the minimal complexity Hamiltonian is the linear/quadratic Hamiltonian - corresponding to solution by a Riccati equation. If the Hamiltonian is given as a pointwise maximum of $M$ linear/quadratic Hamiltonians, then one could say the complexity of the Hamiltonian is $M$. Such PDEs can also arise in switched linear systems.
The algorithm constructs the semiconvex dual of the value function as a max-plus sum, i.e., a pointwise maximum, of certain quadratic functions. An infinite time-horizon problem is considered, and as such, the value function is approximated by iterating a finite-horizon semigroup until a large enough propagation horizon is reached. This finite-horizon semigroup itself is approximated as maximum of a finite number of quadratic forms, or as a semigroup for a system switching between $M$ linear-quadratic systems. The dual of the approximate value at each iteration is stored as a set of quadratic functions. Acting on this dual with the above dual semigroup leads to a new approximation, where the number of quadratics grows by a fixed factor at each iteration. This is the curse-of-complexity. To attenuate this computational growth, we develop a pruning method based on semidefinite programming (SDP).

## II. Problem statement and assumptions

The HJB PDEs we consider arise in infinite-horizon nonlinear optimal control problems, and their Hamiltonians are given as (or well-approximated by) pointwise maxima of linear-quadratic functions. Note that pointwise maxima of quadratic forms can approximate, arbitrarily closely, any semiconvex function. More specifically, we consider

$$
\begin{align*}
& 0=-\widetilde{H}(x, \nabla V)=-\max _{m \in\{1,2, \ldots, M\}}\left\{H^{m}(x, \nabla V)\right\}(1) \\
& V(0)=0 \tag{2}
\end{align*}
$$

(i.e., with boundary condition $V=0$ at the origin) where each of the constituent Hamiltonians has the form

$$
\begin{align*}
H^{m}(x, p)= & \frac{1}{2} x^{T} D^{m} x+\frac{1}{2} p^{T} \Sigma^{m} p+\left(A^{m} x\right)^{T} p \\
& +\left(l_{1}^{m}\right)^{T} x+\left(l_{2}^{m}\right)^{T} p+\alpha^{m} \tag{3}
\end{align*}
$$

where $D^{m}, \Sigma^{m}$ are $n \times n$ symmetric matrices, $l_{1}^{m}, l_{2}^{m} \in \mathbb{R}^{n}$ and $\alpha^{m} \in \mathbb{R}$.

Hamiltonian $\widetilde{H}$ is associated with an optimal control problem for switched linear systems. Let $\mathcal{M} \doteq\{1,2, \ldots M\}$.

The corresponding value function is

$$
\begin{align*}
\widetilde{V}(x) & =\sup _{w \in \mathcal{W}} \sup _{\mu \in \mathcal{D}} \sup _{T<\infty} \widetilde{J}(x, T ; w, \mu) \\
& =\sup _{w \in \mathcal{W}} \sup _{\mu \in \mathcal{D}_{\infty}} \sup _{T<\infty} \int_{0}^{T} L^{\mu_{t}}\left(\xi_{t}\right)-\frac{\gamma^{2}}{2}\left|w_{t}\right|^{2} d t \tag{4}
\end{align*}
$$

where

$$
\begin{aligned}
& L^{\mu_{t}}(x)=\frac{1}{2} x^{T} D^{\mu_{t}} x+\left(l_{1}^{\mu_{t}}\right)^{T} x+\alpha^{\mu_{t}} \\
& \mathcal{D}_{\infty}=\{\mu:[0, \infty) \rightarrow \mathcal{M}: \text { measurable }\} \\
& \mathcal{W} \doteq L_{2}^{\text {loc }}\left([0, \infty) ; \mathbb{R}^{k}\right)
\end{aligned}
$$

and the state dynamics are given by

$$
\begin{equation*}
\dot{\xi}=A^{\mu_{t}} \xi+l_{2}^{\mu_{t}}+\sigma^{\mu_{t}} w_{t}, \quad \xi_{0}=x \tag{5}
\end{equation*}
$$

where $\sigma^{m}$ and $\gamma$ are such that $\Sigma^{m}=\frac{1}{\gamma^{2}} \sigma^{m}\left(\sigma^{m}\right)^{T}$. Here $\mu_{t}$ is a switching control which appears in addition to the control $w$.

To motivate the assumptions for this rather general problem class, we consider $\widetilde{H}$ as being constructed so as to resemble some given nonlinear control problem which has a (finite) solution. That is, we think of $\stackrel{\rightharpoonup}{H}$ as being chosen to resemble some other Hamiltonian, which may correspond to the originating object of interest. In particular, we suppose that problem

$$
\begin{equation*}
0=-\widetilde{\widetilde{H}}(x, \nabla V), \quad V(0)=0 \tag{6}
\end{equation*}
$$

has finite value, and that we are choosing $\widetilde{H}$ to approximate $\widetilde{\widetilde{H}}$. Let $Q_{K}=\left\{\phi: \mathbb{R}^{n} \rightarrow \mathbb{R} \mid \phi\right.$ is semiconvex, and $0 \leq$ $\left.\phi(x) \leq(K / 2)|x|^{2} \forall x \in \mathbb{R}^{n}\right\}$. We may take $Q_{K}$ as the domain of the semigroup. We make the following block of assumptions.

Assume there exists unique viscosity solution, $\widetilde{\widetilde{V}}$, to (6) in $Q_{K}$ for some $K \in(0, \infty)$.
$\underset{\sim}{\text { Assume that }} \widetilde{H}(x, p)=\max _{m \in \mathcal{M}} H^{m}(x, p) \leq$ $\widetilde{H}(x, p)$ for all $x, p \in \mathbb{R}^{n}$.
Assume $H^{1}(x, p)$ has coefficients satisfying the following: $l_{1}^{1}=l_{2}^{1}=0 ; \alpha^{1}=0$; there exists $c_{A} \in$ $(0, \infty)$ such that $x^{\prime} A^{1} x \leq-c_{A}|x|^{2} \forall x \in \mathbb{R}^{n} ; D^{1}$ is positive definite, symmetric; and $\gamma^{2} /\left|\sigma^{1}\right|^{2}>$ $c_{D} / c_{A}^{2}$, where $c_{D}$ is such that $x^{\prime} D^{1} x \leq c_{D}|x|^{2}$ $\forall x \in \mathbb{R}^{n}$.
Assume that system (5) is controllable in the sense that given $x, y \in \mathbb{R}^{n}$ and $T>0$, there exist processes $w \in \mathcal{W}$ and $\mu$ measurable with range in $\mathcal{M}$, such that $\xi_{T}=y$ when $\xi_{0}=x$ and one applies controls $w, \mu$.

Note that the last of these assumptions, the controllability assumption, is satisfied if there exists at least one $m \in$ $\mathcal{M}$ such that $\sigma^{m}\left(\sigma^{m}\right)^{T}$ (which is $n \times n$ ) has $n$ positive eigenvalues.

Assume there exist $c_{1}, c_{2}<\infty$ such that for any
$\varepsilon$-optimal pair, $\mu^{\varepsilon}, w^{\varepsilon}$ for the $\widetilde{H}$ problem, one has

$$
\begin{equation*}
\left\|w^{\varepsilon}\right\|_{L_{2}[0, T]}^{2} \leq c_{1}+c_{2}|x|^{2} \tag{A.w}
\end{equation*}
$$

for all $\varepsilon \in(0,1]$, all $T<\infty$ and all $x \in \mathbb{R}^{n}$.

Note that the behavior specified in $(A . w)$ is proved in the purely quadratic case (c.f., [12]) under reasonable assumptions on the constituent-Hamiltonian matrices, but in this more general case, we assume it instead. Lastly, we make the following assumption.

Assume there exist $\underline{T}, c_{3} \in(0, \infty)$ such that for all $x \in \mathbb{R}^{n}$, all $\varepsilon_{\sim} \in(0,1]$, and all $\mu^{\varepsilon}, w^{\varepsilon}$ which $\underset{\widetilde{V}}{\text { are }} \varepsilon$-optimal for $\widetilde{V}$ (i.e., such that $\widetilde{J}\left(x, \mu^{\varepsilon}, w^{\varepsilon}\right) \geq$
$\widetilde{V}(x)-\varepsilon)$, one has

$$
\int_{0}^{T} L^{\mu_{t}^{\varepsilon}}\left(\xi_{t}^{\varepsilon}\right) d t \geq c_{3} \int_{0}^{T}\left|\xi_{t}^{\varepsilon}\right|^{2} d t \quad \forall T \geq \underline{T}
$$

where $\dot{\xi}_{t}^{\varepsilon}=A^{\mu_{t}^{\varepsilon}} \xi_{t}^{\varepsilon}+l_{2}^{\mu_{t}^{\varepsilon}}+\sigma^{\mu_{t}^{\varepsilon}} w_{t}^{\varepsilon}, \xi_{0}^{\varepsilon}=x$.
Note that these last two assumptions might be difficult to verify. Easily verifiable assumptions appear in [12], [14], but these generate a significantly smaller class of systems than those for which these methods apply.

Now, define the operator

$$
\widetilde{S}_{T}[\phi]=\sup _{w \in \mathcal{W}} \sup _{\mu \in \mathcal{D}_{T}} \int_{0}^{T} L^{\mu_{t}}\left(\xi_{t}\right)-\frac{\gamma^{2}}{2}\left|w_{t}\right|^{2} d t+\phi\left(\xi_{T}\right)
$$

where $\mathcal{D}_{T}=\{\mu:[0, T) \rightarrow \mathcal{M}:$ measurable $\}$. Under the above assumptions, a viscosity solution, $\widetilde{V}$ of (1),(2) exists, satisfies $0 \leq \widetilde{V} \leq \widetilde{\widetilde{V}}$ and is given by $\widetilde{V}=\lim _{T \rightarrow \infty} \widetilde{S}_{T}\left[V_{0}\right]$ for any $V_{0} \in Q_{K}$ such that $0 \leq V_{0} \leq \widetilde{V}$, [13], [14].

In the max-plus algebra, addition and multiplication are defined as $a \oplus b=\max \{a, b\}$ and $a \otimes b=a+b$, respectively. It is well known that $\widetilde{S}_{T}$ forms a max-plus linear semigroup.

## III. CURSE-OF-DIMENSIONALITY-FREE ALGORITHM

The key steps in the curse-of-dimensionality-free algorithm developed in [14] are given below. Since we are interested in understanding how the curse-of-complexity arises in this algorithm, we shall sidestep the theoretical foundations which are well covered in [14], [12], and focus on the algorithmic flow.

## A. Approximate propagation

Define the consituent-Hamiltonian semigroup operators as

$$
S_{\tau}^{m}[\phi]=\sup _{w \in \mathcal{W}} \int_{0}^{\tau} L^{m}\left(\xi_{t}\right)-\frac{\gamma^{2}}{2}\left|w_{t}\right|^{2} d t+\phi\left(\xi_{\tau}\right)
$$

Importantly, propagation of a quadratic $\phi$ by an $S_{\tau}^{m}$ operator can be reduced to solution of a differential Riccati equation. Define the time-indexed operators

$$
\bar{S}_{\tau}[\phi](x)=\max _{m \in \mathcal{M}} S_{\tau}^{m}[\phi](x)=\bigoplus_{m \in \mathcal{M}} S_{\tau}^{m}[\phi](x)
$$

Fix any $T<\infty$. Under the above assumptions, we have (c.f., [13])

$$
\lim _{N \rightarrow \infty}\left\{\bar{S}_{T / N}\right\}^{N}[\phi]=\widetilde{S}_{T}[\phi]
$$

where the superscript $N$ represents repeated application of the operator, $N$ times.

## B. Duals of the semigroup operators

This algorithm uses the concept of semiconvex dual (c.f., [12]). For a function $\phi$ which is uniformly semiconvex with constant $\beta$ such that $\beta-c<0$, the semiconvex dual, $a$, is given by

$$
\begin{equation*}
a(z) \quad=-\max _{x \in \mathbb{R}^{N}} \psi(x, z)-\phi(x) \tag{7}
\end{equation*}
$$

and the dual relationship is given by

$$
\begin{equation*}
\phi(x) \quad=\max _{z \in \mathbb{R}^{n}}[\psi(x, z)+a(z)] \tag{8}
\end{equation*}
$$

where $\psi(x, z)=-(c / 2)|x-z|^{2}$.
We may obtain the duals of the $S_{\tau}^{m}$ operators, $\mathcal{B}_{\tau}^{m}$, and these are also max-plus linear semigroup operators. In particular, they are max-plus integral operators with kernels

$$
B_{\tau}^{m}(x, z)=-\max _{y \in \mathbb{R}^{N}}\left\{\psi(y, x)-S_{\tau}^{m}[\psi(\cdot, z)](y)\right\}
$$

Importantly, note that as the $S_{\tau}^{m}[\psi(\cdot, z)](y)$ are quadratic functions, the $B_{\tau}^{m}$ are quadratic functions. Each of these is obtained only once, at the outset of the algorithm.

Rather than computing directly the value function $\tilde{V}$, we shall approximate its semiconvex dual function $\bar{a}$, defined as in Eqn (7), using the iterative scheme described in the next section. Then, the approximated value function will be recovered from $\bar{a}$ by the inverse formula, as in Eqn (8).

## C. Dual space propagation and the curse-of-complexity

Once the kernels $B_{\tau}^{m}$ of the dual semigroup are obtained, one can begin the iteration. One may begin with an initial quadratic (in the dual space), say

$$
\bar{a}^{0}(z)=\widehat{a}^{0}(z)=(1 / 2)(z-\widehat{z})^{T} \widehat{Q}(z-\widehat{z})+\widehat{r}
$$

Given approximation, $\bar{a}^{k}$ at step $k$, one obtains the next iterate from

$$
\begin{aligned}
\bar{a}^{k+1}(z) & =\bigoplus_{m \in \mathcal{M}} \int_{\mathbb{R}^{n}}^{\oplus} B_{\tau}^{m}(z, y) \otimes \bar{a}^{k}(y) d y \\
& =\max _{m \in \mathcal{M}} \max _{y \in \mathbb{R}^{n}}\left[B_{\tau}^{m}(z, y)+\bar{a}^{k}(y)\right]
\end{aligned}
$$

If $\bar{a}^{k}$ has the form $\bar{a}^{k}(z)=\bigoplus_{\left\{m_{i}\right\}_{i=1}^{k} \in \mathcal{M}^{k}} \widehat{a}_{\left\{m_{i}\right\}_{i=1}^{k}}(z)$ where each $\widehat{a}_{\left\{m_{i}\right\}_{i=1}^{k}}$ is a quadratic form, then $\bar{a}^{k+1}$ takes the form

$$
\bar{a}^{k+1}(z)=\bigoplus_{\left\{m_{i}\right\}_{i=1}^{k+1} \in \mathcal{M}^{k+1}} \widehat{a}_{\left\{m_{i}\right\}_{i=1}^{k+1}}^{k+1}(z)
$$

where the $\widehat{a}_{\left\{m_{i}\right\}_{i=1}^{k+1}}^{k+1}$ are also quadratic. Consequently, the computations reduce to obtaining the coefficients of these quadratics at each step, and these computations are analytic (modulo matrix inverses). This is the reason that the computational growth in the space dimension is only cubic. However, note that the number of quadratics comprising the $\bar{a}^{k}$ grows by a factor of $M$ at each iteration - hence the curse-of-complexity. It has been noted that quite typically, most of the quadratics do not contribute to the value (as they never achieve the maximum at any point, $z$ ), and may
be pruned wihtout consequence. We now proceed to discuss the use of semidefinite programming as a means of pruning the constituent quadratics, the $\widehat{a}_{\left\{m_{i}\right\}_{i=1}^{k+1}}^{k}$.

Lastly, note that once one has propagated sufficiently far (say $k=K$ steps), the value function approximation is recovered from $\bar{a}^{K}$ via (8).

## IV. Pruning Algorithms

In the above curse-of-dimensionality-free algorithm, at step $k, \bar{a}^{k}$ is represented as a max-plus sum of quadratics. Let us index the elements of this sum by integers $i \in \mathcal{I}_{k}$ (rather than by the sequences $\left\{m_{i}\right\}_{i=1}^{k}$ ). That is, we have

$$
\bar{a}^{k}(z)=\bigoplus_{i \in \mathcal{I}^{k}} \widehat{a}_{i}^{k}(z)
$$

where we let each $\widehat{a}_{i}^{k}$ be given in the form

$$
\widehat{a}_{i}^{k}(z)=\widehat{a}_{i}(z)=z^{T} A_{i} z+2 b_{i}^{T} z+c_{i}
$$

where we delete the superscript $k$ for simplicity of notation here and in the sequel.

Recall that we are reducing computational cost by pruning quadratics ( $\widehat{a}_{i}$ ) which do not contribute to the solution approximation (not achieving the maximum at any $z \in$ $\mathbb{R}^{n}$ ). Consequently, we want to determine whether the $p^{t h}$ quadratic contributes to the pointwise maximum. i.e. whether there is a region where it is greater than all other quadratics. Fix $p \in \mathcal{I}^{k}$. Thus we want to ensure feasibility of

$$
\begin{equation*}
\widehat{a}_{p}(z) \geq \widehat{a}_{i}(z) \quad \forall i \neq p \tag{9}
\end{equation*}
$$

Alternatively, we consider the problem:

$$
\begin{align*}
& \text { Maximize } G(z, \nu) \doteq \nu \quad \text { subject to }  \tag{10}\\
& \widehat{a}_{p}(z)-\widehat{a}_{i}(z) \geq \nu \quad \forall i \neq p
\end{align*}
$$

Then, the maximum value of $\nu, \bar{\nu}$, is the maximum amount by which the $p^{t h}$ quadratic can be lowered before it submerges below the max-plus sum of the rest. If $\bar{\nu} \leq 0$, then $p^{t h}$ quadratic does not contribute to the max-plus sum, and hence it can be pruned without consequence. If $\bar{\nu}>0$, the quadratic contributes to the max-plus sum, and $\bar{\nu}$ can serve as some measure of contribution of the $p^{t h}$ quadratic to the value function, enabling us to rank the quadratics. This is useful in over-pruning.

Note that (10) implies that the importance metric,

$$
\begin{equation*}
\bar{\nu}=\max _{z \in \Re^{n}} \min _{i \neq p} \widehat{a}_{p}(z)-\widehat{a}_{i}(z) \tag{11}
\end{equation*}
$$

This ranking scheme is independent of the location of the quadratic. Since in the curse-of-dimensionality-free method, the solution is grown from the origin, and the region of interest is often near the origin, we would like to have ranking scheme which reflects this bias. Hence we consider following importance metric

$$
\begin{equation*}
\nu^{0}=\max _{z \in \Re^{n}} \frac{1}{1+|z|^{2}} \min _{i \neq p}\left(\widehat{a}_{p}(z)-\widehat{a}_{i}(z)\right) \tag{12}
\end{equation*}
$$

which discounts the quadratics which contribute to the value function far away from origin. The extent of this
bias between location and contribution, can be tweaked by multiplier on the term $|z|^{2}$. Similar to (10), the above metric can be reformulated as

$$
\begin{align*}
& \text { Maximize } G(z, \nu) \doteq \nu \quad \text { subject to }  \tag{13}\\
& \widehat{a}_{p}(z)-\widehat{a}_{i}(z) \geq \nu\left(1+|z|^{2}\right) \quad \forall i \neq p
\end{align*}
$$

and $\widehat{a}_{p}(z)$ can be pruned if and only if the maximum value $\nu^{0} \leq 0$.

## A. Pairwise pruning

Before undertaking the pruning using semidefinite programming, pairwise pruning is used, which checks between all pairs of quadratic basis functions, and prunes those which are completely dominated by another. Let $A=A_{i}-A_{j}, b=$ $b_{i}-b_{j}, c=c_{i}-c_{j}$, and define $q(z)=z^{T} A z+2 b^{T} z+c$. Then $q$ is nonnegative everywhere if and only if the homogeneous quadratic form, $z^{T} A z+2 t b^{T} z+c t^{2}$ is nonnegative for all $z \in \mathbb{R}^{n}$ and all $t \in \mathbb{R}$ (easily proved using $q\left(t^{-1} z\right) \geq 0$ when $t \neq 0$ ). Latter statement is true if and only if

$$
\left[\begin{array}{cc}
c & b^{T}  \tag{14}\\
b & A
\end{array}\right] \succeq 0
$$

If we define for any index $p, Q_{p}=\left[\begin{array}{cc}c_{p} & b_{p}^{T} \\ b_{p} & A_{p}\end{array}\right]$. Using (14), in the pairwise comparison between $i$ 'th and $j^{\prime}$ 'th quadratics, later can be pruned if $Q_{i}-Q_{j} \succeq 0$. Pairwise pruning reduces the computational effort of the semidefinite pruning by getting rid of obviously dominated quadratics.

## B. Shor's semidefinite relaxation based pruning

The problem of evaluating an individual quadratic $\widehat{a}_{p}(z)$ for pruning, (13), can be rephrased as below. Let $q_{i}(z)=$ $\widehat{a}_{p}(z)-\widehat{a}_{i}(z)$ for all $i \neq p$. Then, $\widehat{a}_{p}$ can be pruned if and only if

$$
\begin{equation*}
\nu^{0} \doteq \max _{z, \nu}\left\{\nu: q_{i}(z)-\nu\left(1+|z|^{2}\right) \geq 0 \quad \forall i \neq p\right\} \leq 0 \tag{15}
\end{equation*}
$$

Lemma 4.1: With $\lambda \in \Re^{\# I^{k}-1}$ such that $\lambda_{i} \geq 0$ and $\lambda \neq 0, \nu \in \Re$ is an upper bound on $\nu^{0}$ if following condition is satisfied.

$$
\begin{equation*}
\sum_{i \neq p} \lambda_{i}\left(q_{i}(z)-\nu(1+|z|)^{2}\right) \leq 0 \forall z \tag{16}
\end{equation*}
$$

Proof: From (15), $\exists z$ such that, $q_{i}(z)-v^{0}\left(1+|z|^{2}\right) \geq$ $0, \forall i \neq p$, therefore,

$$
\begin{equation*}
\sum_{i \neq p} \lambda_{i}\left(q_{i}(z)-v^{0}\left(1+|z|^{2}\right)\right) \geq 0 \tag{17}
\end{equation*}
$$

Subtracting (16) from (17)

$$
\left(\nu-\nu^{0}\right)\left(1+|z|^{2}\right) \sum_{i} \lambda_{i} \geq 0
$$

with assumptions on $\lambda, \sum_{i \neq p} \lambda_{i}>0$, so that we can divide by it, to get, $\nu \geq \nu^{0}$. Hence proved.

Now we will seek to minimize this upper bound $\nu$ by varying $\lambda$ and $\nu$ subject to constraint (16). Also note that, if $(\lambda, \nu)$ are feasible, so is $(k \lambda, \nu)$ for $k>0$. using this we
can normalize $\lambda$ by dividing by $\sum \lambda_{i}$. Which implies, $\lambda$ lies within a simplex $\mathcal{S}, \lambda_{i} \geq 0, \sum \lambda_{i}=1$.

$$
\begin{equation*}
\nu_{0}=\min _{\lambda \in \mathcal{S}, \nu \in \Re}\left\{\nu: \sum_{i \neq p} \lambda_{i} q_{i}(z) \leq \nu(1+|z|)^{2} \forall z\right\} \tag{18}
\end{equation*}
$$

Since $q_{i}(z)=\widehat{a}_{p}-\widehat{a}_{i}$, using linear superposition and result (14), (18) can be reposed with following semidefinite program:

$$
\begin{equation*}
\nu_{0}=\min _{\lambda \in \mathcal{S}, \nu}\left\{\nu: \sum_{i \neq p} \lambda_{i}\left(Q_{p}-Q_{i}\right) \preceq \nu I\right\} \tag{19}
\end{equation*}
$$

Note that if such minimal $\nu_{0}<0$ value, the by Lemma 4.1, $\nu^{0} \leq \nu_{0} \leq 0$. Hence as per (15), p'th quadratic, $\widehat{a}_{p}(z)$ can be pruned. Since this gives sufficient condition for pruning, it leads to conservative pruning. If $\nu_{0}>0$, the prunability is not conclusive. Nevertheless, it does give us a working indication of the importance of the quadratic. Since (19) can be restated as,

$$
\begin{equation*}
\nu_{0}=\min _{\lambda \in \mathcal{S}, \nu}\left\{\nu: Q_{p} \preceq \nu I+\sum_{i \neq p} \lambda_{i} Q_{i}\right\} \tag{20}
\end{equation*}
$$

if $\nu_{0}>0$, it indicates that the $p^{t h}$ quadratic has to come down by at least $\nu_{0}\left(1+|z|^{2}\right)$, before it is dominated by the convex hull of the remaining quadratics.

An additional way to develop intuition for result (20) is as follows. The above test evaluates the $p^{t h}$ quadratic $\widehat{a}_{p}$, which can be pruned if $\nu=0$ satisfies the inequality in (20). Thus $\widehat{a}_{p}$ can be pruned if, $Q_{p} \preceq \sum_{i \neq p} \lambda_{i} Q_{i}$. Thus, if the convex hull of remaining quadratics intersects the semidefinite cone of quadratics greater than $\widehat{a}_{p}$ for all $z \in \Re^{n}$, then $\widehat{a}_{p}$ can be pruned.

## C. Dual of Shor's relaxation based pruning

If $q_{i}(z)=\widehat{a}_{p}(z)-\widehat{a}_{i}(z)=z^{T} \bar{A}_{i} z+2 \bar{b}_{i}^{T} z+\bar{c}_{i}$, and if we define,

$$
\bar{Q}_{i}=\left[\begin{array}{cc}
\bar{c}_{i} & \bar{b}_{i}^{T} \\
\bar{b}_{i} & \bar{A}_{i}
\end{array}\right]=Q_{p}-Q_{i} \text { and } Z(z)=\left[\begin{array}{cc}
1 & z^{T} \\
z & z z^{T}
\end{array}\right]
$$

Then value of a quadratic form $q_{i}(z)$ can be written as Frobenius inner product of $\bar{Q}_{i}$ and $Z(z)$.

$$
\begin{aligned}
q_{i}(z) & =z^{T} \bar{A}_{i} z+2 \bar{b}_{i}^{T} z+\bar{c}_{i} \\
& =\left[\begin{array}{l}
1 \\
z
\end{array}\right]^{T} \bar{Q}_{i}\left[\begin{array}{l}
1 \\
z
\end{array}\right]=\operatorname{Tr}\left(\bar{Q}_{i}\left[\begin{array}{l}
1 \\
z
\end{array}\right]\left[\begin{array}{l}
1 \\
z
\end{array}\right]^{T}\right) \\
& =\left\langle\bar{Q}_{i}, Z(z)\right\rangle
\end{aligned}
$$

Similarly, $q_{i}(z)-\nu\left(1+z^{2}\right)=\left\langle\bar{Q}_{i}-\nu I, Z(z)\right\rangle$ Thus from (15), $\widehat{a}_{p}$ can be pruned if and only if

$$
\begin{equation*}
\nu^{0}=\max _{z \in \Re^{n}, \nu \in \Re}\left\{\nu:\left\langle\bar{Q}_{i}-\nu I, Z(z)\right\rangle \geq 0 \quad \forall i \neq p\right\} \leq 0 \tag{21}
\end{equation*}
$$

Let $\mathcal{Z}$ be the nonlinear manifold of set of all symmetric dyadic matrices $Z(z), z \in \Re^{n}$. All matrices in $\mathcal{Z}$ are positive semidefinite with northwestern entry 1 . Let $\overline{\mathcal{Z}}$ be set of all
such matrices. Replacing $\mathcal{Z}$ with $\overline{\mathcal{Z}} \supset \mathcal{Z}$, , we get relaxation and an upper bound for the maximum in (21).

$$
\nu^{0} \leq \bar{\nu}^{0} \doteq \max _{Z \in \overline{\mathcal{Z}}, \nu \in \Re}\left\{\nu:\left\langle\bar{Q}_{i}-\nu I, Z\right\rangle \geq 0 \quad \forall i \neq p\right\}
$$

If $\bar{\nu}^{0} \leq 0$, then $\nu^{0} \leq 0$, implying prunability. Thus $\widehat{a}_{p}$ can be pruned if

$$
\max _{Z, \nu \in \Re}\left\{\begin{array}{c|c}
\left\langle\bar{Q}_{i}-\nu I, Z\right\rangle \geq 0 \quad \forall i \neq p  \tag{22}\\
Z \succeq 0, Z_{11}=1
\end{array}\right\} \leq 0
$$

Since $Z \succeq 0$ and $Z_{11}=1, \operatorname{Tr}(Z)>0$. Also note that $\langle I, Z\rangle=\operatorname{Tr}(Z)$. If we define $Y=Z / \operatorname{Tr}(Z)$, then $Y \succeq 0$, $\operatorname{Tr}(Y)=1$. Also

$$
\left\langle\bar{Q}_{i}-\nu I, Z\right\rangle=\left\langle\bar{Q}_{i}, Z\right\rangle-\nu \operatorname{Tr}(Z)=\operatorname{Tr}(Z)\left\{\left\langle\bar{Q}_{i}, Y\right\rangle-\nu\right\}
$$

constraint set in (22) can be simplified to following semidefinite program. Thus $\widehat{a}_{p}(z)$ can be pruned if

$$
\max _{Y, \nu \in \Re}\left\{\begin{array}{c|c}
\nu & \left\langle\bar{Q}_{i}, Y\right\rangle \geq \nu \quad \forall i \neq p  \tag{23}\\
Y \succeq 0, \operatorname{Tr}(Y)=1
\end{array}\right\} \leq 0
$$

This program is a dual of (19).
Intuitively, $L(Q)=\langle Q, Y\rangle$ can be thought as a linear functional over space of $n \times n$ symmetric matrices, taking constant values over hyperplanes normal (in the sense of above inner product) to $Y$. So we are searching in the space of hyperplane normals $Y \succeq 0, \operatorname{Tr}(Y)=1$, which is a slice of the cone of semidefinite matrices. If we can find a hyperplane, separating $Q_{p}$ from convex hull of rest of $Q_{i}, \forall i \neq p$, then prunability of $p$ th quadratic is not conclusive. Hence it need not be pruned.

Since above is the sufficient condition for pruning, it leads to conservative pruning. It is not a necessary condition, due to the nonconvexity of $\mathcal{Z}$, manifold of dyadic matrices. To achive better pruning, higher order pruning techniques can be used.

## V. COMPUTATIONAL COMPLEXITY

Since our aim is to reduce the curse-of-complexity without loosing the freedom from the curse-of-dimensionality, it is worthwhile to discuss the computational overhead involved in these pruning methods. They are polynomial in both dimensionality and the number of quadratic functions. In particular, they are free from the curse-of-dimensionality.

A generic semi-definite program $\mathcal{P}$ is given by

$$
P_{0}=\min _{\eta \in \mathbb{R}^{N}}\left\{c^{T} \eta: A_{0}+\sum_{j=1}^{N} \eta_{j} A_{j} \geq 0,\|\eta\|_{2} \leq R\right\}
$$

where the $A_{j}$ are symmetric matrices with $\widetilde{M}$ diagonal blocks of size $k_{i} \times k_{i}, i=1, \ldots \widetilde{M}$. We say that $\eta^{\varepsilon}$ is an $\varepsilon$-optimal solution if

$$
\left\|\eta^{\varepsilon}\right\|_{2} \leq R, A_{0}+\sum_{j=1}^{N} \eta_{j}^{\varepsilon} A_{j} \geq-\epsilon I, c^{T} \eta^{\varepsilon} \leq P_{0}+\epsilon
$$

In [4], the authors derive the computational complexity of obtaining such an $\eta^{\varepsilon}$.

$$
\begin{align*}
\mathcal{C}(\mathcal{P}, \epsilon)= & \mathcal{O}(1)\left(1+\sum_{1}^{\widetilde{M}} k_{i}\right)^{1 / 2}  \tag{24}\\
& \cdot N\left(N^{2}+N \sum_{1}^{\widetilde{M}} k_{i}^{2}+\sum_{1}^{\widetilde{M}} k_{i}^{3}\right) \mathcal{D}(\mathcal{P}, \epsilon)
\end{align*}
$$

where $\mathcal{D}(\mathcal{P}, \epsilon)$ depends on the specific problem data (as indicated by the $\mathcal{P}$ in the argument as well as $\epsilon$ ). Using this expression, we can obtain an upper bound for the arithmetic complexity of the pruning algorithms. Assuming the worst case scenario where no quadratic gets pruned, we find the complexity of testing one quadratic for pruning as follows.

For the Shor's relaxation pruning of (19), we have $\eta=$ $\left[\begin{array}{ll}\nu & \lambda^{T}\end{array}\right]^{T}$. Hence $N=\# \mathcal{I}_{k}$. In addition to main LMI, we have $\lambda_{i} \geq 0$ and an equality $\sum_{\lambda_{i}}=1$, which can be split up into two inequalities, giving number of LMI blocks $\widetilde{M}=$ $\# \mathcal{I}_{k}+2$, and $k_{1}=n+1, k_{2}=k_{3}=\ldots=k_{\widetilde{M}}=1$. Substituting these into (24), complexity is

$$
\begin{align*}
\mathcal{C}(\mathcal{P}, \epsilon)= & \mathcal{O}(1)\left(\# \mathcal{I}_{k}+n+3\right)^{1 / 2} \# \mathcal{I}_{k}^{2}  \tag{25}\\
& \cdot\left[3 \# \mathcal{I}_{k}+2+(n+1)^{2}(n+2)\right] \mathcal{D}(\mathcal{P}, \epsilon)
\end{align*}
$$

The complexity grows as $n^{3.5}$ and $\# \mathcal{I}_{k}^{3.5}$.
For the dual of Shor's relaxation pruning of (23), independent variables are $Y$, which has $(n+1)(n+2) / 2$ free variables (corresponding to the upper triangle), and $\nu$. Thus $N=\frac{1}{2}\left(n^{2}+3 n+4\right)$. Number of block LMIs is $\widetilde{M}=\# \mathcal{I}_{k}+2$, and $k_{1}=\ldots=k_{\widetilde{M}-3}=1, k_{\widetilde{M}-2}=n+1$, $k_{\widetilde{M}-1}=k_{\widetilde{M}_{5}}=1$. Substituting these into (24), complexity grows as $n^{6.5}$ and $\# \mathcal{I}_{k}^{1.5}$.

Note that these computational cost growth rates as a function of space dimension, $n$, are poorer than the cubic growth of the basic algorithm. However, they are still tremendously helpful for real-world problems, and the curse-of-dimensionality is still very far off.

## VI. IMPORTANCE-BASED OVER-PRUNING

Both the Shor's relaxation and its dual pruning schemes generate importance metric for unpruned quadratics. Overpruning keeps a limited number of more important quadratics, due to memory/computational constraints. This is useful in containing the complexity growth, while ensuring an accurate solution in a region of interest. At end of $k$ 'th iteration, the algorithm sorts $\widehat{a}_{i}^{k}$ according to importance metric, and keeps at most $L(k)$ quadratics, pruning the rest.

However, there does not yet exist a theory which allows us to map the importance measure of a quadratic at step $k$ to an error bound in the approximation at the terminal step. Regardless, application of these methods has proven extremely fruitful, as can be seen in the example below.

## VII. Six-dimensional, SIX-Hamiltonian example

We applied the curse-of-dimensionality-free method to the following problem over $\mathbb{R}^{6}$ with six constituent quadratic

Hamiltonians. Note that the problem was tweaked to exhibit sufficiently complex and interesting behavior, such that there is interaction amongst dimensions, and each operator is important somewhere in the domain. Hence the following data yields a reasonably rich problem.

We shall specify the matrices in terms of the following building blocks for the dynamics:

$$
\begin{array}{cc}
A_{a}=\left[\begin{array}{cc}
-1 & .5 \\
.1 & -1
\end{array}\right], & A_{b}=A_{a}, \\
A_{c}=A_{a}, & A_{d}=\left[\begin{array}{cc}
-1 & .5 \\
.3 & -1
\end{array}\right] \\
A_{e}=A_{a}, & A_{f}=\left[\begin{array}{cc}
-1 & .5 \\
.1 & -1
\end{array}\right] \\
\Sigma_{a}=0.4 \times\left[\begin{array}{cc}
0.27 & -.01 \\
-.01 & 0.27
\end{array}\right], & \Sigma_{b}=0.4 \Sigma_{a} \\
\Sigma_{c}=\Sigma_{a}, & \Sigma_{d}=0.4 \Sigma_{a} \\
\Sigma_{e}=\Sigma_{a}, & \Sigma_{f}=0.4 \Sigma_{a}
\end{array}
$$

and the following building blocks for the payoff functions:

$$
\begin{array}{cc}
D_{a}=\left[\begin{array}{cc}
1.5 & .2 \\
.2 & 1.5
\end{array}\right], & D_{b}=1.4 \times D_{a}, \\
D_{c}=1.4 \times\left[\begin{array}{ll}
0.2 & 1.5 \\
1.5 & 0.2
\end{array}\right], & D_{d}=1.2 \times\left[\begin{array}{cc}
1.6 & 0 \\
0 & 0
\end{array}\right], \\
D_{e}=1.1 \times\left[\begin{array}{ll}
0.3 & 1.5 \\
1.5 & 0.3
\end{array}\right], & D_{f}=1.3 \times\left[\begin{array}{cc}
0 & 0 \\
0 & 1.6
\end{array}\right] .
\end{array}
$$

We will use a parameter to adjust the interaction in the dynamics across the dimensions, and this will be $\gamma=-0.1$. Now we are ready to define each of the Hamiltonians. We need to specify parameters for the dynamics $\left(A, \Sigma, l_{2}\right)$ and the payoff $\left(D, l_{1}, \alpha\right)$. For the example below, $l_{1}=0$ and $l_{2}=0$ for all the Hamiltonians. The remaining parameters are as follows.

For the first Hamiltonian, $H^{1}$, we let

$$
\begin{aligned}
A^{1} & =\left[\begin{array}{ccc}
A_{a} & \gamma I & \gamma I \\
\gamma I & A_{a} & 0 \\
\gamma I & 0 & A_{a}
\end{array}\right],
\end{aligned} \quad \Sigma^{1}=\left[\begin{array}{ccc}
\Sigma_{a} & 0 & 0 \\
0 & \Sigma_{a} & 0 \\
0 & 0 & \Sigma_{a}
\end{array}\right], ~\left[\begin{array}{ccc}
D_{a} & 0 & 0 \\
0 & D_{a} & 0 \\
0 & 0 & D_{a}
\end{array}\right], \quad \alpha_{1}=0 . ~ l l
$$

For the second Hamiltonian, $H^{2}$, we let

$$
\begin{aligned}
A^{2}=\left[\begin{array}{ccc}
A_{b} & \gamma I & \gamma I \\
\gamma I & A_{b} & 0 \\
\gamma I & 0 & A_{b}
\end{array}\right], & \Sigma^{2}=\left[\begin{array}{ccc}
\Sigma_{b} & 0 & 0 \\
0 & \Sigma_{b} & 0 \\
0 & 0 & \Sigma_{b}
\end{array}\right] \\
D^{2} & =\left[\begin{array}{ccc}
D_{b} & 0 & 0 \\
0 & D_{b} & 0 \\
0 & 0 & D_{b}
\end{array}\right],
\end{aligned}
$$

For the third Hamiltonian, $H^{3}$, we let

$$
\begin{aligned}
A^{3}=\left[\begin{array}{ccc}
A_{c} & \gamma I & \gamma I \\
\gamma I & A_{c} & 0 \\
\gamma I & 0 & A_{c}
\end{array}\right], & \Sigma^{3}=\left[\begin{array}{ccc}
\Sigma_{c} & 0 & 0 \\
0 & \Sigma_{c} & 0 \\
0 & 0 & \Sigma_{c}
\end{array}\right] \\
D^{3} & =\left[\begin{array}{ccc}
D_{c} & 0 & 0 \\
0 & D_{c} & 0 \\
0 & 0 & D_{c}
\end{array}\right],
\end{aligned}
$$

For the fourth Hamiltonian, $H^{4}$, we let

$$
\begin{aligned}
A^{4}=\left[\begin{array}{ccc}
A_{d} & \gamma I & \gamma I \\
\gamma I & A_{d} & 0 \\
\gamma I & 0 & A_{d}
\end{array}\right], & \Sigma^{4}=\left[\begin{array}{ccc}
\Sigma_{d} & 0 & 0 \\
0 & \Sigma_{d} & 0 \\
0 & 0 & \Sigma_{d}
\end{array}\right] \\
D^{4} & =\left[\begin{array}{ccc}
D_{d} & 0 & 0 \\
0 & D_{d} & 0 \\
0 & 0 & D_{d}
\end{array}\right],
\end{aligned}
$$

For the fifth Hamiltonian, $H^{5}$, we let

$$
\begin{aligned}
A^{5} & =\left[\begin{array}{ccc}
A_{e} & \gamma I & \gamma I \\
\gamma I & A_{e} & 0 \\
\gamma I & 0 & A_{e}
\end{array}\right],
\end{aligned} \quad \Sigma^{5}=\left[\begin{array}{ccc}
\Sigma_{e} & 0 & 0 \\
0 & \Sigma_{e} & 0 \\
0 & 0 & \Sigma_{e}
\end{array}\right]
$$

For the sixth Hamiltonian, $H^{6}$, we let

$$
\begin{aligned}
A^{6}=\left[\begin{array}{ccc}
A_{f} & \gamma I & \gamma I \\
\gamma I & A_{f} & 0 \\
\gamma I & 0 & A_{f}
\end{array}\right], & \Sigma^{6}=\left[\begin{array}{ccc}
\Sigma_{f} & 0 & 0 \\
0 & \Sigma_{f} & 0 \\
0 & 0 & \Sigma_{f}
\end{array}\right] \\
D^{6} & =\left[\begin{array}{ccc}
D_{f} & 0 & 0 \\
0 & D_{f} & 0 \\
0 & 0 & D_{f}
\end{array}\right],
\end{aligned}
$$

For this example, we let the time-discretization step-size be $\tau=0.2$, and propagation was carried out with the Shor's semidefinite relaxation based pruning. The overpruning threshold was set heuristically to $L(k)=20+6 k$. That is, a maximum of $L(k)$ quadratics, $\widehat{a}_{i}^{k}$, were retained at the $k^{t h}$ step. In this test, 25 iteration steps were carried out in 30 minutes, yielding a rather accurate solution in a compact domain in all six dimensions. This computation-time is for an Apple mac desktop, from roughly 2005. Slices of statistics for this value function along the 1-2 axes are shown in the accompanying figures. The backsubstitution error depends on the propagation as well as the time-discretization. The theoretical error bounds in [13] are of the form $\varepsilon\left(1+|x|^{2}\right)$ (over the entire space) where $\varepsilon \downarrow 0$ as the number of propagation steps goes to infinity and time-discretization go to zero, with the convergence rates derived in above reference.

## VIII. Conclusions

Thus in this paper, two semi-definite programming schemes for pruning the quadratics were proposed for containing the curse-of-complexity in the curse-ofdimensionality free method. Computational complexity for both is polynomial in space dimension. Both give us an


Fig. 1. $x_{1}$ partial on the $x_{1}, x_{2}$ plane


Fig. 2. $\quad x_{2}$ partial on the $x_{1}, x_{2}$ plane


Fig. 3. Optimal switching policy on the $x_{1}, x_{2}$ plane


Fig. 4. Backsubstitution error on the $x_{1}, x_{2}$ plane
importance metric to rank the quadratics according the importance, which is a function of contribution to the pointwise maximum and its location. This is useful in case we need to over-prune. These methods have been applied to solve a sample 6 dimensional, 6 hamiltonian problem in reasonable amount of time.

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    William McEneaney is with the Dept. of Mech. and Aero. Eng., University of California San Diego, San Diego, CA 92093-0411, USA. wmcenean@ucsd.edu

    Ameet deshpande is with the Dept. of Mech. and Aero. Eng., University of California San Diego, San Diego, CA 92093-0411, USA. adeshpan@ucsd.edu

    Stephane Gaubert is with INRIA, Domaine de Voluceau 78153 LeChesnay Cédex, France. Stephane. Gaubert@inria.fr

