An Augmented Lagrangian Approach to Constrained MAP Inference

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Abstract

We propose a new algorithm for approximate MAP inference on factor graphs, by combining augmented Lagrangian optimization with the dual decomposition method. Each slave subproblem is given a quadratic penalty, which pushes toward faster consensus than in previous subgradient approaches. Our algorithm is provably convergent, parallelizable, and suitable for fine decompositions of the graph. We show how it can efficiently handle problems with (possibly global) structural constraints via simple sort operations. Experiments on synthetic and real-world data show that our approach compares favorably with the state-of-the-art.

1. Introduction

Graphical models enable compact representations of probability distributions, being widely used in computer vision, natural language processing (NLP), and computational biology (Koller & Friedman, 2009). A prevalent problem is the one of inferring the most probable configuration, the so-called *maximum a posteriori* (MAP). Unfortunately, this problem is intractable, except for a limited class of models. This fact precludes computing the MAP exactly in many important models involving *non-local* features or requiring *structural constraints* to ensure valid predictions.

A significant body of research has thus been placed on *approximate* MAP inference, *e.g.*, via linear programming relaxations (Schlesinger, 1976). Several message-passing algorithms have been proposed that exploit the graph structure in these relaxations (Wainwright et al., 2005; Kolmogorov, 2006; Werner, 2007; Globerson & Jaakkola, 2008; Ravikumar et al., 2010). In the same line, Komodakis et al. (2007) proposed a method based on the classical *dual decomposition* technique (DD; Dantzig & Wolfe 1960; Everett III 1963; Shor 1985), which breaks the original problem into a set of smaller (slave) subproblems, splits the shared variables, and tackles the Lagrange dual with the *subgradient* algorithm. Initially applied in computer vision, DD has also been shown effective in NLP (Koo et al., 2010). The drawback is that the subgradient algorithm is very slow to converge when the number of slaves is large. This led Jojic et al. (2010) to propose an accelerated gradient method by smoothing the objective.

In this paper, we ally the simplicity of DD with the effectiveness of *augmented Lagrangian methods*, which have a long-standing history in optimization (Hestenes, 1969; Powell, 1969; Glowinski & Marroco, 1975; Gabay & Mercier, 1976; Boyd et al., 2011). The result is a novel algorithm for approximate MAP inference: *DD-ADMM* (*Dual Decomposition with the Alternating Direction Method of Multipliers*). Rather than placing all efforts in attempting progress in the dual, DD-ADMM looks for a saddle point of the Lagrangian function, which is augmented with a quadratic term to penalize slave disagreements. Key features of DD-ADMM are:

- It is suitable for heavy parallelization (many slaves);
- it is provably convergent, even when each slave subproblem is only solved approximately;
- consensus among slaves is fast, by virtue of the quadratic penalty term, hence it exhibits faster convergence in the *primal* than competing methods;
- in addition to providing an optimality certificate for the exact MAP, it also provides guarantees that the *LP-relaxed* solution has been found.

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After providing the necessary background (Sect. 2) and introducing and analyzing DD-ADMM (Sect. 3), we turn to the slave subproblems (Sect. 4). Of particular concern to us are problems with *structural constraints*, which arise commonly in NLP, vision, and other structured prediction tasks. We show that, for several important constraints, each slave can be solved *exactly* and *efficiently* via sort operations. Experiments with pairwise MRFs and dependency parsing (Sect. 5) testify to the success of our approach.

2. Background

2.1. Problem Formulation

Let $\boldsymbol{X} \triangleq (X_1, \ldots, X_N) \in \mathcal{X}$ be a vector of discrete random variables, where each $X_i \in \mathcal{X}_i$, with $\mathcal{X}_i$ a finite set. We assume that $\boldsymbol{X}$ has a Gibbs distribution associated with a factor graph $\mathcal{G}$ (Kschischang et al., 2001), composed of a set of variable nodes $\{1, \ldots, N\}$ and a set of factor nodes $\mathcal{A}$, with each $a \in \mathcal{A}$ linked to a subset of variables $\mathcal{N}(a) \subseteq \{1, \ldots, N\}$:

$$P_{\boldsymbol{\theta},\boldsymbol{\phi}}(\boldsymbol{x}) \propto \exp\left(\textstyle\sum_{i=1}^{N} \theta_i(x_i) + \sum_{a\in\mathcal{A}} \phi_a(\boldsymbol{x}_a)\right). \quad (1)$$

Above, $\boldsymbol{x}_a$ stands for the subvector indexed by the elements of $\mathcal{N}(a)$, and $\theta_i(.)$ and $\phi_a(.)$ are, respectively, unary and higher-order log-potential functions. To accommodate hard constraints, we allow these functions to take values in $\mathbb{R} \cup \{-\infty\}$. For simplicity, we write $\boldsymbol{\theta}_i \triangleq (\theta_i(x_i))_{x_i\in\mathcal{X}_i}$ and $\boldsymbol{\phi}_a \triangleq (\phi_a(\boldsymbol{x}_a))_{\boldsymbol{x}_a\in\mathcal{X}_a}$.

We are interested in the task of finding the most probable assignment (the MAP), $\hat{\boldsymbol{x}} \triangleq \arg\max_{\boldsymbol{x}\in\mathcal{X}} P_{\boldsymbol{\theta},\boldsymbol{\phi}}(\boldsymbol{x})$. This (in general NP-hard) combinatorial problem can be transformed into a linear program (LP) by introducing marginal variables $\boldsymbol{\mu} \triangleq (\boldsymbol{\mu}_i)_{i=1}^n$ and $\boldsymbol{\nu} \triangleq (\boldsymbol{\nu}_a)_{a\in\mathcal{A}}$, constrained to the *marginal polytope* of $\mathcal{G}$, *i.e.*, the set of realizable marginals (Wainwright & Jordan, 2008). Denoting this set by $\mathcal{M}(\mathcal{G})$, this yields

$$\text{OPT} \triangleq \max_{(\boldsymbol{\mu},\boldsymbol{\nu})\in\mathcal{M}(\mathcal{G})} \textstyle\sum_i \boldsymbol{\theta}_i^\top \boldsymbol{\mu}_i + \sum_a \boldsymbol{\phi}_a^\top \boldsymbol{\nu}_a, \quad (2)$$

which always admits an integer solution. Unfortunately, $\mathcal{M}(\mathcal{G})$ often lacks a concise representation, which renders (2) intractable. A common workaround is to replace $\mathcal{M}(\mathcal{G})$ by the outer bound $\mathcal{L}(\mathcal{G}) \supseteq \mathcal{M}(\mathcal{G})$—the so-called *local polytope*, defined as

$$\mathcal{L}(\mathcal{G}) = \left\{ (\boldsymbol{\mu},\boldsymbol{\nu}) \;\middle|\; \begin{array}{l} \mathbf{1}^\top \boldsymbol{\mu}_i = 1, \forall i \wedge \\ \mathbf{H}_{ia}\boldsymbol{\nu}_a = \boldsymbol{\mu}_i, \forall a, i \in \mathcal{N}(a) \wedge \\ \boldsymbol{\nu}_a \geq 0, \forall a \end{array} \right\}, \quad (3)$$

where $\mathbf{H}_{ia}(x_i, \boldsymbol{x}_a) = 1$ if $[\boldsymbol{x}_a]_i = x_i$, and 0 otherwise. This yields the following LP relaxation of (2):

$$\text{OPT}' \triangleq \max_{(\boldsymbol{\mu},\boldsymbol{\nu})\in\mathcal{L}(\mathcal{G})} \textstyle\sum_i \boldsymbol{\theta}_i^\top \boldsymbol{\mu}_i + \sum_a \boldsymbol{\phi}_a^\top \boldsymbol{\nu}_a, \quad (4)$$

which will be our main focus throughout. Obviously, $\text{OPT}' \geq \text{OPT}$, since $\mathcal{L}(\mathcal{G}) \supseteq \mathcal{M}(\mathcal{G})$.

2.2. Dual Decomposition

Several message passing algorithms (Wainwright et al., 2005; Kolmogorov, 2006; Globerson & Jaakkola, 2008) are derived via some reformulation of (4) followed by dualization. The DD method (Komodakis et al., 2007) reformulates (4) by adding new variables $\boldsymbol{\nu}_i^a$ (for each factor a and $i \in \mathcal{N}(a)$) that are local "replicas" of the marginals $\boldsymbol{\mu}_i$. Letting $\mathcal{N}(i) \triangleq \{a | i \in \mathcal{N}(a)\}$ and $d_i = |\mathcal{N}(i)|$ (the degree of node i), (4) is rewritten as

$$\begin{array}{ll} \max\limits_{\boldsymbol{\nu},\boldsymbol{\mu}} & \sum_a \left(\sum_{i\in\mathcal{N}(a)} d_i^{-1} \boldsymbol{\theta}_i^\top \boldsymbol{\nu}_i^a + \boldsymbol{\phi}_a^\top \boldsymbol{\nu}_a \right) \quad (5) \\ \text{s.t.} & (\boldsymbol{\nu}^a_{\mathcal{N}(a)}, \boldsymbol{\nu}_a) \in \mathcal{M}(\mathcal{G}_a), \quad \forall a \\ & \boldsymbol{\nu}_i^a = \boldsymbol{\mu}_i, \quad \forall a, i \in \mathcal{N}(a), \end{array}$$

where $\mathcal{G}_a$ is the subgraph of $\mathcal{G}$ comprised only of factor a and the variables in $\mathcal{N}(a)$, $\mathcal{M}(\mathcal{G}_a)$ is the corresponding marginal polytope, and we denote $\boldsymbol{\nu}^a_{\mathcal{N}(a)} \triangleq (\boldsymbol{\nu}_i^a)_{i\in\mathcal{N}(a)}$. (Note that by definition, $\mathcal{L}(\mathcal{G}) = \{(\boldsymbol{\mu},\boldsymbol{\nu}) \mid (\boldsymbol{\mu}_{N_a}, \boldsymbol{\nu}_a) \in \mathcal{M}(\mathcal{G}_a),\ \forall a \in \mathcal{A}\}$.) Problem (5) would be completely separable (over the factors) if it were not the "coupling" constraints $\boldsymbol{\nu}_i^a = \boldsymbol{\mu}_i$. Introducing Lagrange multipliers $\boldsymbol{\lambda}_i^a$ for these constraints, the dual problem (*master*) becomes

$$\begin{array}{ll} \min\limits_{\boldsymbol{\lambda}} & L(\boldsymbol{\lambda}) \triangleq \sum_a s_a \left(\left(d_i^{-1}\boldsymbol{\theta}_i + \boldsymbol{\lambda}_i^a\right)_{i\in\mathcal{N}(a)}, \boldsymbol{\phi}_a \right) \\ \text{s.t.} & \boldsymbol{\lambda} \in \Lambda \triangleq \left\{ \boldsymbol{\lambda} \;\middle|\; \sum_{a\in\mathcal{N}(i)} \boldsymbol{\lambda}_i^a = 0,\ \forall i \right\}, \end{array} \quad (6)$$

where each s_a corresponds to a *slave* subproblem

$$s_a(\boldsymbol{\omega}^a_{\mathcal{N}(a)}, \boldsymbol{\phi}_a) \triangleq \max_{\substack{(\boldsymbol{\nu}^a_{\mathcal{N}(a)}, \boldsymbol{\nu}_a) \\ \in\mathcal{M}(\mathcal{G}_a)}} \sum_{i\in\mathcal{N}(a)} \boldsymbol{\omega}_i^\top \boldsymbol{\nu}_i^a + \boldsymbol{\phi}_a^\top \boldsymbol{\nu}_a. \quad (7)$$

Note that the slaves (7) are MAP problems of the same kind as (2), but *local* to each factor a. Denote by $(\hat{\boldsymbol{\nu}}^a_{\mathcal{N}(a)}, \hat{\boldsymbol{\nu}}_a) = \text{MAP}(\boldsymbol{\omega}^a_{\mathcal{N}(a)}, \boldsymbol{\phi}_a)$ the maximizer of (7). The master problem (6) can be addressed elegantly with a projected subgradient algorithm: note that a subgradient $\nabla_{\boldsymbol{\lambda}_i^a} L(\boldsymbol{\lambda})$ is readily available upon solving the ath slave, via $\nabla_{\boldsymbol{\lambda}_i^a} L(\boldsymbol{\lambda}) = \hat{\boldsymbol{\nu}}_i^a$. These slaves can be handled in parallel and then have their solutions gathered for computing a projection onto Λ, which is simply a centering operation. This results in Alg. 1.

Alg. 1 inherits the properties of subgradient algorithms, hence it converges to the optimal value of OPT' in (4) if the stepsize sequence $(\eta_t)_{t\in T}$ is diminishing and nonsummable (Bertsekas et al., 1999). In practice, convergence can be quite slow if the number of slaves is large. This is because it may be hard to reach a consensus on variables with many replicas.

Algorithm 1 DD-Subgradient

input: factor graph $\mathcal{G}$, parameters $\boldsymbol{\theta}, \boldsymbol{\phi}$, number of iterations T, stepsize sequence $(\eta_t)_{t=1}^T$
Initialize $\boldsymbol{\lambda} = \mathbf{0}$
for $t = 1$ **to** T **do**
  for each factor $a \in \mathcal{A}$ **do**
    Set $\boldsymbol{\omega}_i^a = d_i^{-1}\boldsymbol{\theta}_i + \boldsymbol{\lambda}_i^a$, for $i \in \mathcal{N}(a)$
    Compute $(\hat{\boldsymbol{\nu}}^a_{\mathcal{N}(a)}, \hat{\boldsymbol{\nu}}_a) = \textsc{map}\big(\boldsymbol{\omega}^a_{\mathcal{N}(a)}, \boldsymbol{\phi}_a\big)$
  end for
  Compute average $\boldsymbol{\mu}_i = {d_i}^{-1} \sum_{a:i\in\mathcal{N}(a)} \hat{\boldsymbol{\nu}}_i^a$
  Update $\boldsymbol{\lambda}_i^a \leftarrow \boldsymbol{\lambda}_i^a - \eta_t\,(\hat{\boldsymbol{\nu}}_i^a - \boldsymbol{\mu}_i)$
end for
output: $\boldsymbol{\lambda}$

Algorithm 2 DD-ADMM

input: factor graph $\mathcal{G}$, parameters $\boldsymbol{\theta}, \boldsymbol{\phi}$, number of iterations T, sequence $(\eta_t)_{t=1}^T$, parameter τ
Initialize $\boldsymbol{\mu}$ uniformly, $\boldsymbol{\lambda} = \mathbf{0}$
for $t = 1$ **to** T **do**
  for each factor $a \in \mathcal{A}$ **do**
    Set $\boldsymbol{\omega}_i^a = d_i^{-1}\boldsymbol{\theta}_i + \boldsymbol{\lambda}_i^a + \eta_t\boldsymbol{\mu}_i$, for $i \in \mathcal{N}(a)$
    Update $(\boldsymbol{\nu}^a_{\mathcal{N}(a)}, \boldsymbol{\nu}_a) \leftarrow \textsc{quad}_{\eta_t}\big(\boldsymbol{\omega}^a_{\mathcal{N}(a)}, \boldsymbol{\phi}_a\big)$
  end for
  Update $\boldsymbol{\mu}_i \leftarrow {d_i}^{-1} \sum_{a:i\in\mathcal{N}(a)} \big(\boldsymbol{\nu}_i^a - \eta_t^{-1}\boldsymbol{\lambda}_i^a\big)$
  Update $\boldsymbol{\lambda}_i^a \leftarrow \boldsymbol{\lambda}_i^a - \tau\eta_t\,(\boldsymbol{\nu}_i^a - \boldsymbol{\mu}_i)$
end for
output: $\boldsymbol{\mu}, \boldsymbol{\nu}, \boldsymbol{\lambda}$

3. Augmented Lagrangian Method

In this section we introduce a faster method, DD-ADMM, which replaces the MAP computation at each factor a by a (usually simple) quadratic problem (QP); this method penalizes disagreements among slaves *more aggressively* than the subgradient method.

Given an optimization problem with equality constraints, the augmented Lagrangian (AL) function is the Lagrangian *augmented* with a quadratic constraint violation penalty. For the constraint problem (5), it is

$$A_\eta(\boldsymbol{\mu}, \boldsymbol{\nu}, \boldsymbol{\lambda}) \triangleq \sum_a \Big(\sum_{i\in\mathcal{N}(a)} \big(d_i^{-1}\boldsymbol{\theta}_i + \boldsymbol{\lambda}_i^a\big)^\top \boldsymbol{\nu}_i^a + \boldsymbol{\phi}_a^\top \boldsymbol{\nu}_a \Big) - \sum_a \sum_{i\in\mathcal{N}(a)} \boldsymbol{\lambda}_i^{a\top} \boldsymbol{\mu}_i - \frac{\eta}{2} \sum_a \sum_{i\in\mathcal{N}(a)} \|\boldsymbol{\nu}_i^a - \boldsymbol{\mu}_i\|^2, \quad (8)$$

where η controls the weight of the penalty. Applied to our problem, a traditional AL method (Hestenes, 1969; Powell, 1969) would alternate between the joint maximization of $A_\eta(\boldsymbol{\mu}, \boldsymbol{\nu}, \boldsymbol{\lambda})$ w.r.t. $\boldsymbol{\mu}$ and $\boldsymbol{\nu}$, and an update of the Lagrange multipliers $\boldsymbol{\lambda}$. Unfortunately, the quadratic term in (8) breaks the separability, making the joint maximization w.r.t. $\boldsymbol{\mu}$ and $\boldsymbol{\nu}$ unappealing.

We bypass this problem by using the *alternating direction method of multipliers* (ADMM; Gabay & Mercier 1976; Glowinski & Marroco 1975), in which the joint maximization is replaced by a single Gauss-Seidel step. This yields the following updates:

$$\boldsymbol{\nu}^{(t)} \leftarrow \operatorname*{argmax}_{\boldsymbol{\nu}} A_{\eta_t}(\boldsymbol{\mu}^{(t-1)}, \boldsymbol{\nu}, \boldsymbol{\lambda}^{(t-1)}), \quad (9)$$

$$\boldsymbol{\mu}^{(t)} \leftarrow \operatorname*{argmax}_{\boldsymbol{\mu}} A_{\eta_t}(\boldsymbol{\mu}, \boldsymbol{\nu}^{(t)}, \boldsymbol{\lambda}^{(t-1)}), \quad (10)$$

$$\boldsymbol{\lambda}_i^{a(t)} \leftarrow \boldsymbol{\lambda}_i^{a(t-1)} - \tau\eta_t \left(\boldsymbol{\nu}_i^{a(t)} - \boldsymbol{\mu}_i^{(t)}\right), \forall a, i \in \mathcal{N}(a). \quad (11)$$

Crucially, the maximization w.r.t. $\boldsymbol{\mu}$ (10) has a closed form, while that w.r.t. $\boldsymbol{\nu}$ (9) can be carried out in parallel at each factor, as in Alg. 1. The only difference is that, instead of computing the MAP, each slave now needs to solve a QP of the form

$$\min_{(\boldsymbol{\nu}^a_{\mathcal{N}(a)}, \boldsymbol{\nu}_a)\in\mathcal{M}(\mathcal{G}_a)} \frac{\eta_t}{2} \sum_{i\in\mathcal{N}(a)} \|\boldsymbol{\nu}_i^a - \eta_t^{-1}\boldsymbol{\omega}_i^a\|^2 - \boldsymbol{\phi}_a^\top \boldsymbol{\nu}_a. \quad (12)$$

The resulting algorithm is DD-ADMM (Alg. 2). Let $\textsc{quad}_{\eta_t}(\boldsymbol{\omega}_{\mathcal{N}(a)}, \boldsymbol{\phi}_a)$ denote the solution of (12); as $\eta_t \to 0$, $\textsc{quad}_{\eta_t}(\boldsymbol{\omega}_{\mathcal{N}(a)}, \boldsymbol{\phi}_a)$ approaches $\textsc{map}(\boldsymbol{\omega}_{\mathcal{N}(a)}, \boldsymbol{\phi}_a)$, hence Alg. 2 approaches Alg. 1. However, DD-ADMM converges without the need of decreasing η_t:

Proposition 1 *Let $(\boldsymbol{\mu}^{(t)}, \boldsymbol{\nu}^{(t)}, \boldsymbol{\lambda}^{(t)})_t$ be the sequence of iterates produced by Alg. 2 with a fixed $\eta_t = \eta$ and $0 < \tau \leq (\sqrt{5}+1)/2 \simeq 1.61$. Then the following holds:*

1. *Primal feasibility of (5) is achieved in the limit,* i.e., $\|\boldsymbol{\nu}_i^{a(t)} - \boldsymbol{\mu}_i^{(t)}\| \to 0, \forall a \in \mathcal{A}, i \in \mathcal{N}(a)$;
2. $(\boldsymbol{\mu}^{(t)}, \boldsymbol{\nu}^{(t)})$ *converges to an optimal solution of (4);*
3. $\boldsymbol{\lambda}^{(t)}$ *converges to an optimal solution of (6);*
4. $\boldsymbol{\lambda}^{(t)}$ *is always dual feasible; hence the objective of (6) evaluated at $\boldsymbol{\lambda}^{(t)}$ approaches* OPT' *from above.*

Proof: 1, 2, and 3 are general properties of ADMM algorithms (Glowinski & Le Tallec, 1989, Thm. 4.2). All necessary conditions are met: problem (5) is convex and the coupling constraint can be written in the form $(\boldsymbol{\nu}^a_{\mathcal{N}(a)})_{a\in\mathcal{A}} = \boldsymbol{M}\boldsymbol{\mu}$ where $\boldsymbol{M}$ has full column rank. To show 4, use induction: $\boldsymbol{\lambda}^{(0)} = \mathbf{0} \in \Lambda$; if $\boldsymbol{\lambda}^{(t-1)} \in \Lambda$, *i.e.*, $\sum_{a\in\mathcal{N}(i)} \boldsymbol{\lambda}_i^{a(t-1)} = 0$, $\forall i$, then, after line 9,

$$\sum_{a\in\mathcal{N}(i)} \boldsymbol{\lambda}_i^{a(t)} = \sum_{a\in\mathcal{N}(i)} \boldsymbol{\lambda}_i^{a(t-1)} - \tau\eta_t \left(\sum_{a\in\mathcal{N}(i)} \boldsymbol{\nu}_i^{a(t)} - d_i\boldsymbol{\mu}_i^{(t)} \right) = (1-\tau) \sum_{a\in\mathcal{N}(i)} \boldsymbol{\lambda}_i^{a(t-1)} = 0 \;\Rightarrow\; \boldsymbol{\lambda}^{(t)} \in \Lambda. \quad \blacksquare$$

Prop. 1 reveals yet another important feature of DD-ADMM: after a sufficient decrease of the residual term,

say $\sum_a \sum_{i\in\mathcal{N}(a)} \|\boldsymbol{\nu}_i^a - \boldsymbol{\mu}_i\|^2 < \epsilon$, we have at hand a ϵ-feasible primal-dual pair. If, in addition, the duality gap (difference between (4) and (6)) is $< \delta$, then we are in possession of an (ϵ, δ)-optimality certificate for the LP relaxation. Such a certificate is not readily available in Alg. 1, unless the relaxation is tight.

The next proposition, based on results of Eckstein & Bertsekas (1992), states that convergence may still hold if (12) is solved approximately.

Proposition 2 *Let $\eta_t = \eta$ be fixed and $\tau = 1$. Consider the sequence of residuals $\boldsymbol{r}^{(t)} \triangleq (r_a^{(t)})_{a\in\mathcal{A}}$, where*

$$r_a^{(t)} \triangleq \|(\boldsymbol{\nu}_{\mathcal{N}(a)}^{a(t)}, \boldsymbol{\nu}_a^{(t)}) - \textsc{quad}_{\eta_t}(\boldsymbol{\omega}_{\mathcal{N}(a)}^{(t)}, \boldsymbol{\phi}_a)\|.$$

Then, Prop. 1 still holds provided $\sum_{t=1}^{\infty} \|\boldsymbol{r}^{(t)}\| < \infty$.

4. Solving the Slave Subproblems

We next show how to efficiently solve (12). Several cases admit closed-form solutions: binary pairwise factors, some factors expressing hard constraints, and factors linked to multi-valued variables, once binarized. In all cases, the asymptotic computational cost is the same as that of MAP computations, up to a logarithmic factor. Detailed proofs of the following facts are included as supplementary material.

4.1. Binary pairwise factors

If factor a is binary and pairwise ($|\mathcal{N}(a)| = 2$), problem (12) can be re-written as the minimization of $\frac{1}{2}(z_1 - c_1)^2 + \frac{1}{2}(z_2 - c_2)^2 - c_{12}z_{12}$, w.r.t. $(z_1, z_2, z_{12}) \in [0,1]^3$, under the constraints $z_{12} \le z_1$, $z_{12} \le z_2$, and $z_{12} \ge z_1 + z_2 - 1$, where c_1, c_2 and c_{12} are functions of $\boldsymbol{\omega}_i^a$ and $\boldsymbol{\phi}_a$. Considering $c_{12} \ge 0$, without loss of generality (if $c_{12} < 0$, we recover this case by redefining $c_1' = c_1 + c_{12}$, $c_2' = 1 - c_2$, $c_{12}' = -c_{12}$, $z_2' = 1 - z_2$, $z_{12}' = z_1 - z_{12}$), the lower bound constraints $z_{12} \ge z_1 + z_2 - 1$ and $z_{12} \ge 0$ are always innactive and can be ignored. By inspecting the KKT conditions we obtain the following closed-form solution: $z_{12}^* = \min\{z_1^*, z_2^*\}$ and

$$(z_1^*, z_2^*) = \begin{cases} ([c_1]_{\mathbb{U}}, [c_2 + c_{12}]_{\mathbb{U}}) & \text{if } c_1 > c_2 + c_{12} \\ ([c_1 + c_{12}]_{\mathbb{U}}, [c_2]_{\mathbb{U}}) & \text{if } c_2 > c_1 + c_{12} \\ ([(c_1 + c_2 + c_{12})/2]_{\mathbb{U}}, & \\ \ [(c_1 + c_2 + c_{12})/2]_{\mathbb{U}}) & \text{otherwise,} \end{cases}$$

where $[x]_{\mathbb{U}} = \min\{\max\{x, 0\}, 1\}$ denotes the projection (clipping) onto the unit interval.

4.2. Hard constraint factors

Many applications, *e.g.*, in error-correcting coding (Richardson & Urbanke, 2008) and NLP (Smith & Eisner, 2008; Martins et al., 2010; Tarlow et al., 2010) involve *hard constraint factors*—these are factors with indicator log-potential functions: $\boldsymbol{\phi}_a(\boldsymbol{x}_a) = 0$, if $\boldsymbol{x}_a \in \mathcal{S}_a$, and $-\infty$ otherwise, where $\mathcal{S}_a$ is an *acceptance set*. For binary variables, these factors impose logical constraints; *e.g.*,

- the one-hot XOR factor, for which $\mathcal{S}_{\text{XOR}} = \{(x_1, \dots, x_n) \in \{0,1\}^n | \sum_{i=1}^n x_i = 1\}$,
- the OR factor, for which $\mathcal{S}_{\text{OR}} = \{(x_1, \dots, x_n) \in \{0,1\}^n | \bigvee_{i=1}^n x_i = 1\}$,
- the OR-WITH-OUTPUT factor, for which $\mathcal{S}_{\text{OR-OUT}} = \{(x_1, \dots, x_n) \in \{0,1\}^n | \bigvee_{i=1}^{n-1} x_i = x_n\}$.

Variants of these factors (*e.g.*, with negated inputs/outputs) allow computing a wide range of other logical functions. It can be shown that the marginal polytope of a hard factor with binary variables and acceptance set $\mathcal{S}_a$ is defined by $\boldsymbol{z} \in \operatorname{conv} \mathcal{S}_a$, where $\boldsymbol{z} \triangleq (\mu_1(1), \dots, \mu_n(1))$ and conv denotes the convex hull. Letting $\boldsymbol{c} \triangleq (\omega_a^i(1) + 1 - \omega_a^i(0))_{i\in\mathcal{N}(a)}$, problem (12) is that of minimizing $\|\boldsymbol{z} - \boldsymbol{c}\|^2$ s.t. $\boldsymbol{z} \in \operatorname{conv} \mathcal{S}_a$, which is a Euclidean projection onto a polyhedron:

- $\operatorname{conv} \mathcal{S}_{\text{XOR}}$ is the probability simplex; the projection is efficiently obtained via a sort (Duchi et al., 2008).
- $\operatorname{conv} \mathcal{S}_{\text{OR}}$ is a "faulty" hypercube with a vertex removed, and $\operatorname{conv} \mathcal{S}_{\text{OR-OUT}}$ is a pyramid whose base is a hypercube with a vertex removed; both projections can be efficiently computed with sort operations.

The algorithms and proofs of correctness are provided as supplementary material. In all cases, complexity is $O(|\mathcal{N}(a)| \log |\mathcal{N}(a)|)$ and can be improved to $O(|\mathcal{N}(a)|)$ using a technique similar to Duchi et al. (2008).

4.3. Larger slaves and multi-valued variables

For general factors, a closed-form solution of problem (12) is not readily available. One possible strategy (exploiting Prop. 2) is to use an *inexact* algorithm that becomes sufficiently accurate as Alg. 2 proceeds; this can be achieved by warm-starting with the solution obtained in the previous iteration. This strategy can be useful for handling coarser decompositions, in which each factor is a subgraph such as a chain or a tree. However, unlike the MAP problem in DD-subgradient, in which dynamic programming can be used to compute an exact solution for these special structures, that does not seem possible in QUAD.

Yet, there is an alternative strategy for handling multi-valued variables, which is to *binarize* the graph and make use of the results established in Sect. 4.2 for hard constraint factors. We illustrate this procedure for pairwise MRFs (but the idea carries over when higher order potentials are used): let $X_1, \dots, X_N$ be the variables of the original graph, and $\mathcal{E} \subseteq \{1, \dots, N\}^2$ be

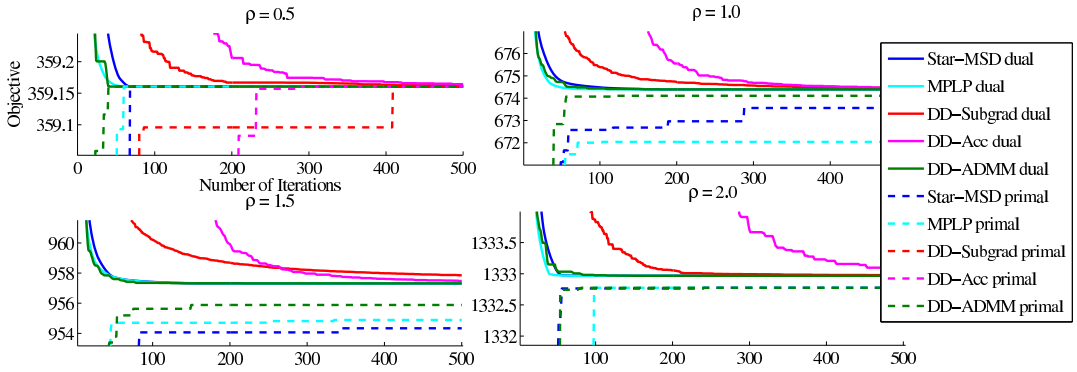

Figure 1. Results for 30×30 random Ising models with several edge couplings. We plot the dual objectives and the best primal feasible solution at each iteration. For DD-subgradient, we set $\eta_t = \eta_0/t$, with η_0 yielding the maximum dual improvement in 10 iterations, with halving steps (those iterations are not plotted). For DD-Acc we plot the most favorable $\epsilon \in \{0.1, 1, 10, 100\}$. For DD-ADMM, we set $\eta = 5.0$ and $\tau = 1.0$. All decompositions are edge-based.

the set of edges. Let $M = |\mathcal{E}|$ be the number of edges and $L = |\mathcal{X}_i|, \forall i$ the number of labels. Then:

- For each node i, define binary variables U_{ik} for each possible $k \in \{1, \dots, L\}$ of X_i. Link these variables to a XOR factor, imposing $\sum_{k=1}^{L} \mu_i(k) = 1$, $\forall i$.
- For each edge $(i, j) \in \mathcal{E}$, define binary variables $U_{ijkk'}$ for each value pair $(k, k') \in \{1, \dots, L\}^2$. Link variables $\{U_{ijkk'}\}_{k'=1}^{L}$ and $\neg U_{ik}$ to a XOR factor, for each $k \in \{1, \dots, L\}$; and link variables $\{U_{ijkk'}\}_{k=1}^{L}$ and $\neg U_{jk'}$ to a XOR factor for each $k' \in \{1, \dots, L\}$. These impose constraints $\mu_i(k) = \sum_{k'=1}^{L} \mu_{ij}(k, k')$, $\forall k$, and $\mu_j(k') = \sum_{k=1}^{l} \mu_{ij}(k, k')$, $\forall k'$.

The constraints above define a local polytope which is equivalent to the one of the original graph, hence problem (4) is the same for both graphs. This process increases the number of factors to $N + 2ML$, where each is a XOR of size L or $L + 1$. However, solving QUAD for each of these factors only costs $O(L \log L)$ (see Sect. 4.2), hence the overall cost per iteration of Alg. 2 is $O(ML^2 \log L)$ if the graph is connected. Up to a log factor, this is the same as message-passing algorithms or DD-subgradient when run in the original graph, which have $O(ML^2)$ cost per iteration.

5. Experiments

We compare DD-ADMM (Alg. 2) with four other approximate MAP inference algorithms:

- Star-MSD (Sontag et al., 2011), an acceleration of the max-sum diffusion algorithm (Kovalevsky & Koval, 1975; Werner, 2007) based on star updates;
- Generalized MPLP (Globerson & Jaakkola, 2008);
- DD-subgradient (Komodakis et al. 2007, Alg. 1);
- DD-Acc (accelerated DD, Jojic et al. 2010).

All these algorithms address problem (4) with the same cost per iteration: the first two use message-passing, performing block coordinate descent in the dual; DD-Acc uses a smoothed dual objective by adding an entropic term to each subproblem, and then applies optimal first-order methods (Nesterov, 1983), yielding $O(1/\epsilon)$ complexity. The primal and dual objectives are the same for all algorithms.

5.1. Binary Pairwise MRFs

Fig. 1 shows typical plots for an Ising model (binary pairwise MRF) on a random grid, with single node log-potentials chosen as $\theta_i(1) - \theta_i(0) \sim \mathcal{U}[-1, 1]$ and mixed edge couplings in $\mathcal{U}[-\rho, \rho]$, where $\rho \in \{0.5, 1, 1.5, 2\}$. Decompositions are edge-based for all methods. For MPLP and Star-MSD, primal feasible solutions $(\hat{x}_i)_{i=1}^{N}$ were obtained by decoding the single node messages, as in Globerson & Jaakkola (2008); for the DD methods we use $\hat{x}_i = \text{argmax}_{x_i} \mu_i(x_i)$.

We observe that DD-subgradient is the slowest, taking a long time to find a "good" primal feasible solution, arguably due to the large number of slave subproblems. Surprisingly, DD-Acc is also not competitive in this setting, as it consumes many iterations before it reaches a near-optimal region.[1] MPLP performs slightly better than Star-MSD and both are comparable to DD-ADMM in terms of convergence of the dual objective. However, DD-ADMM outperforms all competitors at obtaining a "good" feasible primal solution in early iterations (it retrieved the true MAP in all cases, in ≤ 200 iterations). We conjecture that this rapid progress in the primal is due to the penalty term in the AL (8), which is very effective at pushing

[1] It is conceivable that the early iterations of DD-Acc could make faster progress by annealing ϵ. Here we have just used the variant described by Jojic et al. (2010).

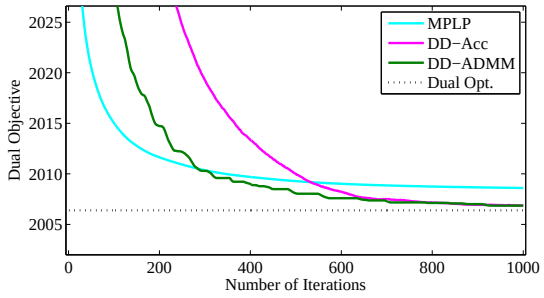

Figure 2. Results for a 20×20 random Potts model with 8-valued nodes and coupling factor $\rho = 10$. We plot the dual objectives, and the value of the true dual optimum. For DD-Acc, $\epsilon = 1.0$; for DD-ADMM, $\eta = 0.5$, $\tau = 1.0$.

for a feasible primal solution of the relaxed LP.

5.2. Multi-valued Pairwise MRFs

To assess the effectiveness of DD-ADMM in the non-binary case, we evaluated it against DD-Acc and MPLP in a Potts model (multi-valued pairwise MRF) with single node log-potentials chosen as $\theta_i(x_i) \sim \mathcal{U}[-1, 1]$ and edge log-potentials as $\theta_{ij}(x_i, x_j) \sim \mathcal{U}[-10, 10]$ if $x_i = x_j$ and 0 otherwise. For DD-Acc and MPLP, we used the same edge decomposition as before, since they can handle multi-valued variables; for DD-ADMM we binarized the graph as described in Sect. 4.3. Fig. 2 shows the best dual solutions obtained at each iteration for the three algorithms. We observe that MPLP decreases the objective very rapidly in the beginning and then slows down. DD-Acc manages to converge faster, but it is relatively slow to take off. DD-ADMM has the best features of both methods.

5.3. Dependency Parsing

A third set of experiments aims at assessing the ability of DD-ADMM for handling problems with heavily constrained outputs. The task is *non-projective dependency parsing* of natural language sentences, to which DD approaches have recently been applied (Koo et al., 2010). Fig. 3 depicts an example of a sentence (the input) and its dependency tree (the output to be predicted). Second-order models are state-of-the-art for this task: they include scores for each possible arc and for certain pairs of arcs (*e.g.*, siblings and grandparents); the goal is to find a directed spanning tree maximizing the overall score.

We experimented with two factor graphs that represent this problem (see Fig. 3). Both models have hard constraints on top of a binary pairwise MRF whose nodes represent arc candidates and whose edges link pairs of arcs which bear a sibling or grandparent relationship. The "tree" model has a TREE hard constraint factor connected to all nodes enforcing the overall assignment to be a directed spanning tree (Smith & Eisner, 2008). Unfortunately, handling this factor poses difficulties for all methods except DD-subgradient.[2]

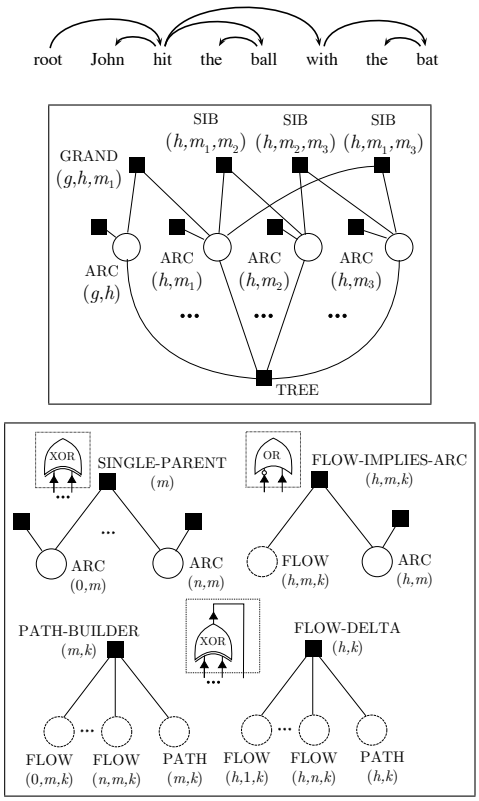

Figure 3. Top: a dependency parse tree, where each arc (h, m) links a *head* word h to a *modifier* word m. Middle: tree-based factor graph corresponding to a second-order dependency parsing model with sibling and grandparent features, including a TREE hard constraint factor. Bottom: the flow-based factor graph is an alternative representation for the same model, in which extra flow and path variables are added, and the TREE factor is replaced by smaller XOR, OR and OR-OUT. See Martins et al. (2010) for details.

[2] Further information about the combinatorial TREE factor can be found in Martins et al. (2010) and references therein. Briefly, solving the MAP subproblem (necessary for DD-subgradient) corresponds to finding a maximum weighted arborescence, which can be done in $O(n^2)$ time, where n is the number of words in the sentence (Tarjan, 1977). Computing all posterior marginals (necessary for DD-Acc) can be done in $O(n^3)$ time invoking the matrix-tree theorem (Smith & Smith, 2007; Koo et al., 2007; McDonald & Satta, 2007). Unfortunately, this procedure suffers from severe numerical problems in the low-temperature setting, which prevents its use in DD-Acc where the temperature must be set as $O(\epsilon/(n \log n))$. Finally, no efficient algorithm is currently known for the simultaneous computation of all max-marginals in the TREE factor (which

The "flow" model imposes the same global constraint by adding extra variables and several XOR, OR and OR-OUT factors (Martins et al., 2010). The two models have the same expressive power, and differ from the one in Koo et al. (2010), which combines a tree constraint with factors that emulate head automata (instead of *all* pairs of arcs) and has fewer slaves. In our case, both factor graphs yield $O(n^3)$ slaves (n is the number of words in the sentence), degrading the performance of standard DD methods.

This is illustrated in Fig. 4, which shows that DD-subgradient converges slowly, even with the more favorable tree-based factor graph (both for synthetic and real-world data). For this problem, MPLP and Star-MSD also have poor performance, while DD-ADMM manages to converge to a near-optimal solution very fast (note the sharp decrease in relative error on the bottom plot, compared with DD-subgradient). This has an impact in final accuracy: setting a maximum of 1000 iterations for all algorithms gave DD-ADMM an advantage of $> 1.5\%$ in unlabeled attachment score (fraction of words with the correct head attached), in the Penn Treebank dataset.

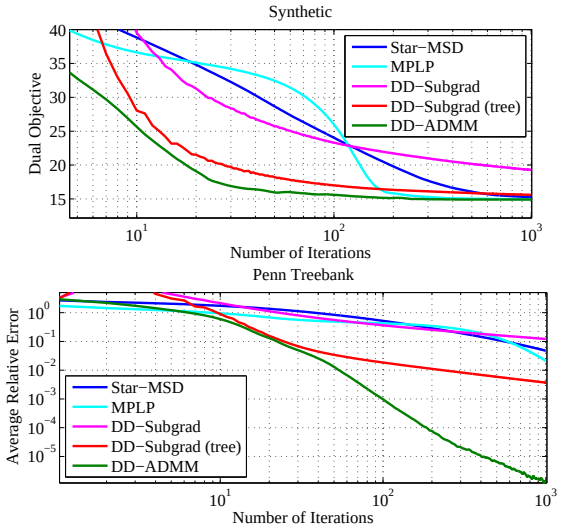

Figure 4. Dependency parsing with 2nd-order models. For DD-subgradient, we consider both tree-based and flow-based factor graphs, and η_0 as in Fig. 1. DD-ADMM ($\eta = 0.05$, $\tau = 1.5$), MPLP, and Star-MSD ran only on the flow-based factor graph (see footnote 2). DD-Acc is not shown due to numerical problems when computing some low-temperature marginals. Top: synthetic 10-word sentences; we randomly generated (unary) arc log-potentials from $\mathcal{N}(0, 1)$ and (pairwise) grandparent and sibling log-potentials from $\mathcal{N}(0, 0.1)$. Bottom: §23 of the Penn Treebank; the plot shows relative errors per iteration w.r.t. the dual optimum, averaged over the 2,399 test sentences.

6. Related Work and Final Remarks

DD-subgradient was first proposed for image segmentation using pairwise (Komodakis et al., 2007) and higher order factor graphs (Komodakis & Paragios, 2009). It was recently adopted for NLP (Koo et al., 2010), with only a few slave subproblems handled with dynamic programming or combinatorial algorithms.

Johnson et al. (2007) proposed smoothing the objective by adding an *entropic* term to each subproblem, with the goal of enabling gradient-based optimization; each subproblem becomes that of computing marginals at a particular temperature. Jojic et al. (2010) combined this with optimal first order methods to accelerate DD to $O(1/\epsilon)$ complexity. This rate, however, relies on the ability to compute low-temperature marginals with arbitrary precision, which poses problems for some of the hard constraint factors considered in this paper. In contrast, DD-ADMM uses *exact* solutions of the corresponding slave subproblems, efficiently computed using sort operations (see Sect. 4.2).

We point out that replacing the quadratic penalty of the AL (8) by an entropic one would *not* lead to the same subproblems as in Jojic et al. (2010): it would lead to the problem of minimizing non-strictly convex free energies with different counting numbers. Although extensions of ADMM to Bregman penalties have been considered in the literature, convergence has been shown only for quadratic penalties.

Quadratic problems were also recently considered in a sequential algorithm (Ravikumar et al., 2010); however, that algorithm tackles the *primal* formulation and only pairwise models are considered. A similar cyclic projection can be adopted in DD-ADMM to approximately solve QUAD for larger slaves.

DD-ADMM is dual decomposable, hence the slaves can be solved in parallel, making it suitable for multi-core architectures with obvious speed-ups. A significant amount of computation can be saved by *caching* and *warm-starting* the subproblems, which tend to become more and more similar across later iterations.

In the future, we plan to experiment with larger slaves, by using approximate ADMM steps as enabled by Prop. 2. The encouraging results of DD-ADMM for solving LP relaxations suggest that it can also be useful for *tightening* these relaxations towards the true MAP, as the MPLP algorithm in Sontag et al. (2008).

would be necessary for MPLP and Star-MSD); or for computing an Euclidean projection onto the arborescence polytope (which would be necessary for DD-ADMM).

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An Augmented Lagrangian Aproach to Constrained MAP Inference
Supplementary Material

A. Derivation of QUAD for Binary Pairwise Factors

In this section, we derive in detail the closed form solution of problem (12) for binary pairwise factors (Sect. 4.1). Recall that the marginal polytope $\mathcal{M}(\mathcal{G}_a)$ is given by:

$$\mathcal{M}(\mathcal{G}_a) = \left\{ (\boldsymbol{\nu}^a_{\mathcal{N}(a)}, \boldsymbol{\nu}_a) \;\middle|\; \begin{array}{ll} \sum_{x_i \in \mathcal{X}_i} \nu^a_i(x_i) = 1, & \forall i \in \mathcal{N}(a) \\ \nu^a_i(x_i) = \sum_{\boldsymbol{x}_a \sim x_i} \nu_a(\boldsymbol{x}_a), & \forall i \in \mathcal{N}(a), x_i \in \mathcal{X}_i \\ \nu_a(\boldsymbol{x}_a) \geq 0, & \forall \boldsymbol{x}_a \in \mathcal{X}_a \end{array} \right\}. \tag{13}$$

If factor a is binary and pairwise ($|\mathcal{N}(a)| = 2$), we may reparameterize our problem by introducing new variables $z_1 \triangleq \nu^a_1(1)$, $z_2 \triangleq \nu^a_2(1)$, and $z_{12} \triangleq \nu_a(1,1)$. Noting that $\boldsymbol{\nu}^a_1 = (1 - z_1, z_1)$, $\boldsymbol{\nu}^a_2 = (1 - z_2, z_2)$, and $\boldsymbol{\nu}_a = (1 - z_1 - z_2 + z_{12}, z_1 - z_{12}, z_2 - z_{12}, z_{12})$, problem (12) becomes

$$\begin{array}{lll} \min\limits_{z_1,z_2,z_{12}} & \frac{\eta_t}{2}[(1 - z_1 - \eta_t^{-1}\omega^a_1(0))^2 + (z_1 - \eta_t^{-1}\omega^a_1(1))^2 + (1 - z_2 - \eta_t^{-1}\omega^a_2(0))^2 + (z_2 - \eta_t^{-1}\omega^a_2(1))^2] & \\ & -\phi_a(00)(1 - z_1 - z_2 + z_{12}) - \phi_a(10)(z_1 - z_{12}) - \phi_a(01)(z_2 - z_{12}) - \phi_a(11) z_{12} & \\ \text{s.t.} & z_{12} \leq z_1, \quad z_{12} \leq z_2, \quad z_{12} \geq z_1 + z_2 - 1, \quad (z_1, z_2, z_{12}) \in [0,1]^3 & \end{array} \tag{14}$$

or, multiplying the objective by the constant $1/(2\eta_t)$:

$$\begin{array}{lll} \min\limits_{z_1,z_2,z_{12}} & \frac{1}{2}(z_1 - c_1)^2 + \frac{1}{2}(z_2 - c_2)^2 - c_{12} z_{12} & \\ \text{s.t.} & z_{12} \leq z_1, \quad z_{12} \leq z_2, \quad z_{12} \geq z_1 + z_2 - 1, \quad (z_1, z_2, z_{12}) \in [0,1]^3, & \end{array} \tag{15}$$

where we have substituted

$$c_1 = (\eta_t^{-1}\omega^a_1(1) + 1 - \eta_t^{-1}\omega^a_1(0) + \eta_t^{-1}\phi_a(00) - \eta_t^{-1}\phi_a(10))/2 \tag{16}$$

$$c_2 = (\eta_t^{-1}\omega^a_2(1) + 1 - \eta_t^{-1}\omega^a_2(0) + \eta_t^{-1}\phi_a(00) - \eta_t^{-1}\phi_a(01))/2 \tag{17}$$

$$c_{12} = (\eta_t^{-1}\phi_a(00) - \eta_t^{-1}\phi_a(10) - \eta_t^{-1}\phi_a(00) + \eta_t^{-1}\phi_a(11))/2. \tag{18}$$

Now, notice that in (15) we can assume $c_{12} \geq 0$ without loss of generality—indeed, if $c_{12} < 0$, we recover this case by redefining $c'_1 = c_1 + c_{12}$, $c'_2 = 1 - c_2$, $c'_{12} = -c_{12}$, $z'_2 = 1 - z_2$, $z'_{12} = z_1 - z_{12}$. Thus, assuming that $c_{12} \geq 0$, the lower bound constraints $z_{12} \geq z_1 + z_2 - 1$ and $z_{12} \geq 0$ are always innactive and can be ignored. Hence, (15) can be simplified to:

$$\begin{array}{lll} \min\limits_{z_1,z_2,z_{12}} & \frac{1}{2}(z_1 - c_1)^2 + \frac{1}{2}(z_2 - c_2)^2 - c_{12} z_{12} & \\ \text{s.t.} & z_{12} \leq z_1, \quad z_{12} \leq z_2, \quad z_1 \in [0,1], \quad z_2 \in [0,1]. & \end{array} \tag{19}$$

Second, if $c_{12} = 0$, the problem becomes separable, and the solution is

$$z^*_1 = [c_1]_{\mathbb{U}}, \quad z^*_2 = [c_2]_{\mathbb{U}}, \quad z^*_{12} = \min\{z^*_1, z^*_2\}, \tag{20}$$

where $[x]_{\mathbb{U}} = \min\{\max\{x, 0\}, 1\}$ is the projection (clipping) onto the unit interval. We next analyze the case where $c_{12} > 0$. The Lagrangian function of (19) is:

$$\begin{array}{rcl} L(\boldsymbol{z}, \boldsymbol{\mu}, \boldsymbol{\lambda}, \boldsymbol{\nu}) & = & \frac{1}{2}(z_1 - c_1)^2 + \frac{1}{2}(z_2 - c_2)^2 - c_{12} z_{12} + \mu_1(z_{12} - z_1) + \mu_2(z_{12} - z_2) \\ & & -\lambda_1 z_1 - \lambda_2 z_2 + \nu_1(z_1 - 1) + \nu_2(z_2 - 1). \end{array} \tag{21}$$

At optimality, the following KKT conditions need to be satisfied:

$$\nabla_{z_1} L(\boldsymbol{z}^*, \boldsymbol{\mu}^*, \boldsymbol{\lambda}^*, \boldsymbol{\nu}^*) = 0 \quad \Rightarrow \quad z_1^* = c_1 + \mu_1^* + \lambda_1^* - \nu_1^* \tag{22}$$
$$\nabla_{z_2} L(\boldsymbol{z}^*, \boldsymbol{\mu}^*, \boldsymbol{\lambda}^*, \boldsymbol{\nu}^*) = 0 \quad \Rightarrow \quad z_2^* = c_2 + \mu_2^* + \lambda_2^* - \nu_2^* \tag{23}$$
$$\nabla_{z_{12}} L(\boldsymbol{z}^*, \boldsymbol{\mu}^*, \boldsymbol{\lambda}^*, \boldsymbol{\nu}^*) = 0 \quad \Rightarrow \quad c_{12} = \mu_1^* + \mu_2^* \tag{24}$$
$$\lambda_1^* z_1^* = 0 \tag{25}$$
$$\lambda_2^* z_2^* = 0 \tag{26}$$
$$\mu_1^*(z_{12}^* - z_1^*) = 0 \tag{27}$$
$$\mu_2^*(z_{12}^* - z_2^*) = 0 \tag{28}$$
$$\nu_1^*(z_1^* - 1) = 0 \tag{29}$$
$$\nu_2^*(z_2^* - 1) = 0 \tag{30}$$
$$\boldsymbol{\mu}^*, \boldsymbol{\lambda}^*, \boldsymbol{\nu}^* \geq 0 \tag{31}$$
$$z_{12}^* \leq z_1^*, \quad z_{12}^* \leq z_2^*, \quad z_1^* \in [0,1], \quad z_2^* \in [0,1] \tag{32}$$

We are going to consider three cases separately:

1. $\boxed{z_1^* > z_2^*}$

   From the primal feasibility conditions (32), this implies $z_1^* > 0$, $z_2^* < 1$, and $z_{12}^* < z_1^*$. Complementary slackness (25,30,27) implies in turn $\lambda_1^* = 0$, $\nu_2^* = 0$, and $\mu_1^* = 0$. From (24) we have $\mu_2^* = c_{12}$. Since we are assuming $c_{12} > 0$, we then have $\mu_2^* > 0$, and complementary slackness (28) implies $z_{12}^* = z_2^*$.

   Plugging the above into (22)–(23) we obtain

   $$z_1^* = c_1 - \nu_1^* \leq c_1, \qquad z_2^* = c_2 + \lambda_2^* + c_{12} \geq c_2 + c_{12}. \tag{33}$$

   Now we have the following:

   - Either $z_1^* = 1$ or $z_1^* < 1$. In the latter case, $\nu_1^* = 0$ by complementary slackness (29), hence $z_1^* = c_1$. Since in any case we must have $z_1^* \leq c_1$, we conclude that $z_1^* = \min\{c_1, 1\}$.
   - Either $z_2^* = 0$ or $z_2^* > 0$. In the latter case, $\lambda_2^* = 0$ by complementary slackness (26), hence $z_2^* = c_2 + c_{12}$. Since in any case we must have $z_2^* \geq \lambda_2$, we conclude that $z_2^* = \max\{0, c_2 + c_{12}\}$.

   In sum:

   $$z_1^* = \min\{c_1, 1\}, \quad z_{12}^* = z_2^* = \max\{0, c_2 + c_{12}\}, \tag{34}$$

   and our assumption $z_1^* > z_2^*$ can only be valid if $c_1 > c_2 + c_{12}$.

2. $\boxed{z_1^* < z_2^*}$

   By symmetry, we have

   $$z_2^* = \min\{c_2, 1\}, \quad z_{12}^* = z_1^* = \max\{0, c_1 + c_{12}\}, \tag{35}$$

   and our assumption $z_1^* < z_2^*$ can only be valid if $c_2 > c_1 + c_{12}$.

3. $\boxed{z_1^* = z_2^*}$

   In this case, it is easy to verify that we must have $z_{12}^* = z_1^* = z_2^*$, and we can rewrite our optimization problem in terms of one variable only (call it z). The problem becomes that of minimizing $\frac{1}{2}(z - c_1)^2 + \frac{1}{2}(z - c_2)^2 - c_{12} z$, which equals a constant plus $(z - \frac{c_1 + c_2 + c_{12}}{2})^2$, subject to $z \in \mathbb{U} \triangleq [0, 1]$. Hence:

   $$z_{12}^* = z_1^* = z_2^* = \left[\frac{c_1 + c_2 + c_{12}}{2}\right]_{\mathbb{U}}. \tag{36}$$

Putting all the pieces together, we have the following solution assuming $c_{12} \geq 0$:

$$z_{12}^* = \min\{z_1^*, z_2^*\}, \qquad (z_1^*, z_2^*) = \begin{cases} ([c_1]_{\mathbb{U}}, & [c_2 + c_{12}]_{\mathbb{U}}) & \text{if } c_1 > c_2 + c_{12} \\ ([c_1 + c_{12}]_{\mathbb{U}}, & [c_2]_{\mathbb{U}}) & \text{if } c_2 > c_1 + c_{12} \\ ([(c_1 + c_2 + c_{12})/2]_{\mathbb{U}}, & [(c_1 + c_2 + c_{12})/2]_{\mathbb{U}}) & \text{otherwise.} \end{cases} \tag{37}$$

B. Derivation of QUAD for Several Hard Constraint Factors

In this section, we consider hard constraint factors with binary variables (Sect. 4.2). These are factors whose log-potentials are indicator functions, *i.e.*, they can be written as $\phi_a : \{0,1\}^m \to \bar{\mathbb{R}}$ with

$$\phi_a(\boldsymbol{x}_a) = \begin{cases} 0 & \text{if } \boldsymbol{x}_a \in \mathcal{S}_a \\ -\infty & \text{otherwise,} \end{cases} \tag{38}$$

where $\mathcal{S}_a \subseteq \{0,1\}^m$ is the acceptance set. Since any probability distribution over $\mathcal{S}_a$ has to assign zero mass to points not in $\mathcal{S}_a$, this choice of potential will always lead to $\nu_a(\boldsymbol{x}_a) = 0, \forall \boldsymbol{x}_a \notin \mathcal{S}_a$. Also, because the variables are binary, we always have $\nu_i^a(0) + \nu_i^a(1) = 1$. In fact, if we introduce the set

$$\mathcal{Z}_a \triangleq \{(\nu_1^a(1), \dots, \nu_m^a(1)) \mid (\boldsymbol{\nu}_{N(a)}^a, \boldsymbol{\nu}_a) \in \mathcal{M}(\mathcal{G}_a) \text{ for some } \boldsymbol{\nu}_a \text{ s.t. } \nu_a(\boldsymbol{x}_a) = 0, \forall \boldsymbol{x}_a \notin \mathcal{S}_a\} \tag{39}$$

we have that the two following optimization problems are equivalent for any function f:

$$\min_{\substack{(\boldsymbol{\nu}_{N(a)}^a, \boldsymbol{\nu}_a) \\ \in \mathcal{M}(\mathcal{G}_a)}} f(\boldsymbol{\nu}_{N(a)}^a) + \boldsymbol{\phi}_a^\top \boldsymbol{\nu}_a \quad = \quad \min_{\boldsymbol{z} \in \mathcal{Z}_a} \tilde{f}(\boldsymbol{z}), \tag{40}$$

where $\tilde{f}(z_1, \dots, z_m) \triangleq f(1 - z_1, z_1, \dots, 1 - z_m, z_m)$. Hence the set $\mathcal{Z}_a$ can be used as a "replacement" of the marginal polytope $\mathcal{M}(\mathcal{G}_a)$. By abuse of language, we will sometimes refer to $\mathcal{Z}_a$ (which is also a polytope) as "the marginal polytope of $\mathcal{G}_a$." As a particularization of (40), we have that the quadratic problem (12) becomes that of computing a projection onto $\mathcal{Z}_a$. Of particular interest is the following result.

Proposition 3 *We have* $\mathcal{Z}_a = \operatorname{conv} \mathcal{S}_a$.

Proof: From the definition of $\mathcal{M}(\mathcal{G}_a)$ and the fact that we are constraining $\nu_a(\boldsymbol{x}_a) = 0, \forall \boldsymbol{x}_a \notin \mathcal{S}_a$, it follows:

$$\begin{aligned} \mathcal{Z}_a &= \left\{ \boldsymbol{z} \geq 0 \;\middle|\; \exists \boldsymbol{\nu}_a \geq 0 \text{ s.t. } \forall i \in N(a), z_i = \sum_{\substack{\boldsymbol{x}_a \in \mathcal{S}_a \\ \boldsymbol{x}_i = 1}} \nu_a(\boldsymbol{x}_a) = 1 - \sum_{\substack{\boldsymbol{x}_a \in \mathcal{S}_a \\ \boldsymbol{x}_i = 0}} \nu_a(\boldsymbol{x}_a) \right\} \\ &= \left\{ \boldsymbol{z} \geq 0 \;\middle|\; \exists \boldsymbol{\nu}_a \geq 0, \sum_{\boldsymbol{x}_a \in \mathcal{S}_a} \nu_a(\boldsymbol{x}_a) = 1 \text{ s.t. } \boldsymbol{z} = \sum_{\boldsymbol{x}_a \in \mathcal{S}_a} \nu_a(\boldsymbol{x}_a) \boldsymbol{x}_a \right\} = \operatorname{conv} \mathcal{S}_a. \end{aligned} \tag{41}$$

■

Note also that $\|\boldsymbol{\nu}_i^a - \eta_t^{-1} \boldsymbol{\omega}_i^a\|^2 = (\nu_i^a(1) - \eta_t^{-1}\omega_i^a(1))^2 + (1 - \nu_i^a(1) - \eta_t^{-1}\omega_i^a(0))^2$ equals a constant plus $2(\nu_i^a(1) - \eta_t^{-1}(\omega_i^a(1) + 1 - \omega_i^a(0))/2)^2$. Hence, (12) reduces to computing the following projection:

$$\operatorname{proj}_a(\boldsymbol{z}_0) \triangleq \operatorname*{argmin}_{\boldsymbol{z} \in \operatorname{conv} \mathcal{S}_a} \frac{1}{2} \|\boldsymbol{z} - \boldsymbol{z}_0\|^2, \qquad \text{where } \boldsymbol{z}_{0i} \triangleq (\omega_i^a(1) + 1 - \omega_i^a(0))/2, \; \forall i. \tag{42}$$

Another important fact has to do with *negated inputs*. Let factor a' be constructed from a by negating one of the inputs (without loss of generality, the first one, x_1)—*i.e.*, $\mathcal{S}_{a'} = \{\boldsymbol{x}_{a'} \mid (1 - x_1, x_2, \dots, x_m) \in \mathcal{S}_a\}$. Then, if we have a procedure for evaluating the operator proj_a, we can use it for evaluating $\operatorname{proj}_{a'}$ through the change of variable $z_1' \triangleq 1 - z_1$, which turns the objective function into $(1 - z_1' - z_{01})^2 = (z_1' - (1 - z_{01}))^2$. Naturally, the same idea holds when there is more than one negated input. The overall procedure computes $\boldsymbol{z} = \operatorname{proj}_{a'}(\boldsymbol{z}_0)$:

1. For each input i, set $z_{0i}' = z_{0i}$ if it is not negated and $z_{0i}' = 1 - z_{0i}$ otherwise.
2. Obtain $\boldsymbol{z}'$ as the solution of $\operatorname{proj}_a(\boldsymbol{z}_0')$.
3. For each input i, set $z_i = z_i'$ if it is not negated and $z_i = 1 - z_i'$ otherwise.

Below, we will also use the following

Algorithm 3 Projection onto simplex (Duchi et al., 2008)

Input: $\boldsymbol{z}_0$
Sort $\mathbf{z}_0$ into $\mathbf{y}_0$: $y_1 \geq \ldots \geq y_m$
Find $\rho = \max\left\{j \in [m] \mid y_{0j} - \frac{1}{j}\left(\sum_{r=1}^{j} y_{0r} - 1\right) > 0\right\}$
Define $\tau = \frac{1}{\rho}\left(\sum_{r=1}^{\rho} y_{0r} - 1\right)$
Output: $\boldsymbol{z}$ s.t. $z_i = \max\{z_{0i} - \tau, 0\}$.

Lemma 4 *Consider a problem of the form*

$$P: \qquad \min_{\boldsymbol{x} \in \mathcal{X}} f(\boldsymbol{x}) \quad \text{s.t.} \quad g(\boldsymbol{x}) \leq 0, \tag{43}$$

where $\mathcal{X}$ is nonempty convex subset of $\mathbb{R}^d$ and $f : \mathcal{X} \to \mathbb{R}$ and $g : \mathcal{X} \to \mathbb{R}$ are convex functions. Suppose that the problem (43) is feasible and bounded below, and let $\mathcal{A}$ be the set of solutions of the relaxed problem $\min_{\boldsymbol{x} \in \mathcal{X}} f(\boldsymbol{x})$, *i.e. $\mathcal{A} = \operatorname{Argmin}_{\boldsymbol{x} \in \mathcal{X}} f(\boldsymbol{x})$. Then:*

1. *if for some $\tilde{\boldsymbol{x}} \in \mathcal{A}$ we have $g(\tilde{\boldsymbol{x}}) \leq 0$, then $\tilde{\boldsymbol{x}}$ is also a solution of the original problem P;*
2. *otherwise (if for all $\tilde{\boldsymbol{x}} \in \mathcal{A}$ we have $g(\tilde{\boldsymbol{x}}) > 0$), the inequality constraint is necessarily active in P, i.e., problem P is equivalent to* $\min_{\boldsymbol{x} \in \mathcal{X}} f(\boldsymbol{x})$ *s.t.* $g(\boldsymbol{x}) = 0$.

Proof: Let f^* be the optimal value of P. The first statement is obvious: since $\tilde{\boldsymbol{x}}$ is a solution of a relaxed problem we have $f(\tilde{\boldsymbol{x}}) \leq f^*$; hence if $\tilde{\boldsymbol{x}}$ is feasible this becomes an equality. For the second statement, assume that $\exists \boldsymbol{x} \in \mathcal{X}$ s.t. $g(\boldsymbol{x}) < 0$ (otherwise, the statement holds trivially). The nonlinear Farkas' lemma (Bertsekas et al., 2003, Prop. 3.5.4, p. 204) implies that there exists some $\lambda^* \geq 0$ s.t. $f(\boldsymbol{x}) - f^* + \lambda^* g(\boldsymbol{x}) \geq 0$ holds for all $\boldsymbol{x} \in \mathcal{X}$. In particular, this also holds for an optimal $\boldsymbol{x}^*$ (i.e., such that $f^* = f(\boldsymbol{x}^*)$), which implies that $\lambda^* g(\boldsymbol{x}^*) \geq 0$. However, since $\lambda^* \geq 0$ and $g(\boldsymbol{x}^*) \leq 0$ (since $\boldsymbol{x}^*$ has to be feasible), we also have $\lambda^* g(\boldsymbol{x}^*) \leq 0$, *i.e.*, $\lambda^* g(\boldsymbol{x}^*) = 0$. Now suppose that $\lambda^* = 0$. Then we have $f(\boldsymbol{x}) - f^* \geq 0$, $\forall \boldsymbol{x} \in \mathcal{X}$, which implies that $\boldsymbol{x}^* \in \mathcal{A}$ and contradicts the assumption that $g(\tilde{\boldsymbol{x}}) > 0, \forall \tilde{\boldsymbol{x}} \in \mathcal{A}$. Hence we must have $g(\boldsymbol{x}^*) = 0$. ∎

B.1. Derivation of QUAD for the XOR Factor

The XOR factor is defined as:

$$\phi_{\text{XOR}}(x_1, \ldots, x_m) = \begin{cases} 0 & \text{if } \sum_{i=1}^{m} x_i = 1 \\ -\infty & \text{otherwise,} \end{cases} \tag{44}$$

where each $x_i \in \{0, 1\}$. When $m = 2$, $\exp(\phi_{\text{XOR}}(x_1, x_2)) = x_1 \oplus x_2$ (the Boolean XOR function), hence the name. When $m > 2$, it is a "one-hot" generalization of $\oplus$.[3] For this case, we have that $\operatorname{conv} \mathcal{S}_{\text{XOR}} = \{\boldsymbol{z} \geq 0 \mid \mathbf{1}^\top \boldsymbol{z} = 1\}$ is the probability simplex. Hence the quadratic problem (12) reduces to that of projecting onto the simplex, which can be done efficiently by Alg. 3. This algorithm requires a sort operation and its cost is $O(m \log m)$.[4]

B.2. Derivation of QUAD for the OR Factor

The OR factor is defined as:

$$\phi_{\text{OR}}(x_1, \ldots, x_m) = \begin{cases} 0 & \text{if } \sum_{i=1}^{m} x_i \geq 1 \\ -\infty & \text{otherwise,} \end{cases} \tag{45}$$

where each $x_i \in \{0, 1\}$. This factor indicates whether any of its input variables is 1, hence the name. For this case, we have that $\operatorname{conv} \mathcal{S}_{\text{OR}} = \{\boldsymbol{z} \in [0, 1]^m \mid \mathbf{1}^\top \boldsymbol{z} \geq 1\}$ is a "faulty" unit hypercube where one of the vertices (the origin) is missing. From Lemma 4, we have that the following procedure solves (42) for $\mathcal{S}_{\text{OR}}$:

[3] There is another generalization of $\oplus$ which returns the parity of $x_1, \ldots, x_m$; ours should not be confused with that one.

[4] In the high-dimensional case, a red-black tree can be used to reduce this cost to $O(m)$ (Duchi et al., 2008). In later iterations of DD-ADMM, great speed ups can be achieved in practice since this procedure is repeatedly invoked with only small changes on the coefficients.

1. Set $\tilde{\boldsymbol{z}}$ as the projection of $\boldsymbol{z}_0$ onto the unit cube. This can be done by clipping each coordinate to the unit interval $\mathbb{U} = [0, 1]$, i.e., by setting $\tilde{z}_i = [z_{0i}]_{\mathbb{U}} = \min\{1, \max\{0, z_{0i}\}\}$. If $\mathbf{1}^\top \tilde{\boldsymbol{z}} \geq 1$, then return $\tilde{\boldsymbol{z}}$. Else go to step 2.
2. Return the projection of $\boldsymbol{z}_0$ onto the simplex (use Alg. 3).

The validity of the second step stems from the fact that, if the relaxed problem in the first step does not return a feasible point, then the constraint $\mathbf{1}^\top \boldsymbol{z} \geq 1$ has to be active, i.e., we must have $\mathbf{1}^\top \boldsymbol{z} = 1$. This, in turn, implies that $\boldsymbol{z} \leq 1$ hence the problem becomes equivalent to the XOR case.

B.3. Derivation of QUAD for the OR-WITH-OUTPUT Factor

The OR-WITH-OUTPUT factor is defined as:

$$\phi_{\text{OR-OUT}}(x_1, \ldots, x_m) = \begin{cases} 0 & \text{if } x_m = \bigvee_{i=1}^{m-1} x_i \\ -\infty & \text{otherwise,} \end{cases} \tag{46}$$

where each $x_i \in \{0, 1\}$. In other words, this factor indicates (via the "output" variable x_m) if any of variables x_1 to x_{m-1} (the "input" variables) is active. Solving the quadratic problem for this factor is slightly more complicated than in the previous two cases; however, we next see that it can also be addressed in $O(m \log m)$ with a sort operation. For this case, we have that

$$\operatorname{conv} \mathcal{S}_{\text{OR-OUT}} = \left\{ \boldsymbol{z} \in [0, 1]^m \;\middle|\; z_m \leq \sum_{i=1}^{m-1} z_i \text{ and } z_m \geq z_i, i = 1, \ldots, m-1 \right\}.$$

It is instructive to write this polytope as the intersection of the three sets $[0, 1]^m$, $\mathcal{A}_1 \triangleq \{\boldsymbol{z} \mid z_m \geq z_i, i = 1, \ldots, m-1\}$, and $\mathcal{A}_2 \triangleq \{\boldsymbol{z} \in [0, 1]^m \mid z_m \leq \mathbf{1}^\top \boldsymbol{z}_1^{m-1}\}$. We further define $\mathcal{A}_0 \triangleq [0, 1]^m \cap \mathcal{A}_1$. From Lemma 4, we have that the following procedure is correct:

1. Set $\tilde{\boldsymbol{z}}$ as the projection of $\boldsymbol{z}_0$ onto the unit cube. If $\tilde{\boldsymbol{z}} \in \mathcal{A}_1 \cap \mathcal{A}_2$, then we are done: just return $\tilde{\boldsymbol{z}}$. Else, if $\tilde{\boldsymbol{z}} \in \mathcal{A}_1$ but $\tilde{\boldsymbol{z}} \notin \mathcal{A}_2$, go to step 3. Otherwise, go to step 2.
2. Set $\tilde{\boldsymbol{z}}$ as the projection of $\boldsymbol{z}_0$ onto $\mathcal{A}_0$ (we will describe how to do this later). If $\tilde{\boldsymbol{z}} \in \mathcal{A}_2$, return $\tilde{\boldsymbol{z}}$. Otherwise, go to step 3.
3. Return the projection of $\boldsymbol{z}_0$ onto the set $\{\boldsymbol{z} \in [0, 1]^m \mid z_m = \mathbf{1}^\top z_1^{m-1}\}$. This corresponds to the QUAD problem of a XOR factor with the mth input negated (we call such factor XOR-WITH-OUTPUT because it is analogous to OR-WITH-FACTOR but with the role of OR replaced by that of XOR). As explained above, it can be solved by projecting onto the simplex (use Alg. 3).

Note that the first step above can be ommited; however, it avoids performing step 2 (which requires a sort) unless it is really necessary. To completely specify the algorithm, we only need to explain how to compute the projection onto $\mathcal{A}_0$ (step 2). Recall that $\mathcal{A}_0 = [0, 1]^m \cap \mathcal{A}_1$. Fortunatelly, it turns out that the following sequential projection is correct:

Procedure 5 *To project onto $\mathcal{A}_0 = [0, 1]^m \cap \mathcal{A}_1$:*

2a. Set $\tilde{\tilde{\boldsymbol{z}}}$ as the projection of $\boldsymbol{z}_0$ onto $\mathcal{A}_1$. Alg. 4 shows how to do this.

2b. Set $\tilde{\boldsymbol{z}}$ as the projection of $\tilde{\tilde{\boldsymbol{z}}}$ onto the unit cube (with the usual clipping procedure).

The proof that the composition of these two projections yields the desired projection onto $\mathcal{A}_0$ is a bit involved, so we defer it to Prop. 6.[5] We only need to describe how to project onto $\mathcal{A}_1$ (step 2a), which is written as the following problem:

$$\min_{\boldsymbol{z}} \frac{1}{2} \|\boldsymbol{z} - \boldsymbol{z}_0\|^2 \quad \text{s.t.} \quad z_m \geq z_i, \; i = 1, \ldots, m-1. \tag{47}$$

[5] Note that in general, the composition of individual projections is not equivalent to projecting onto the intersection. In particular, commuting steps 2a and 2b would make our procedure incorrect.

Algorithm 4 Projection onto $\mathcal{A}_1$

Input: $\boldsymbol{z}_0$
Sort $z_{01}, \ldots, z_{0(m-1)}$ into $y_1 \geq \ldots \geq y_{m-1}$
Find $\rho = \min \left\{ j \in [m] \mid \frac{1}{j} \left(z_{0m} + \sum_{r=1}^{j-1} y_r \right) > y_j \right\}$
Define $\tau = \frac{1}{\rho} \left(z_{0m} + \sum_{r=1}^{\rho-1} y_r \right)$
Output: $\boldsymbol{z}$ s.t. $z_m = \tau$ and $z_i = \min\{z_{0i}, \tau\}$, $i = 1, \ldots, m-1$.

Algorithm 5 Dijkstra's algorithm for projecting onto $\bigcap_{j=1}^{J} \mathcal{C}_j$

Input: Point $\boldsymbol{x}_0 \in \mathbb{R}^d$, convex sets $\mathcal{C}_1, \ldots, \mathcal{C}_J$
Initialize $\boldsymbol{x}^{(0)} = \boldsymbol{x}_0$, $\boldsymbol{u}_j^{(0)} = \boldsymbol{0}$ for all $j = 1, \ldots, J$
for $t = 1, 2, \ldots$ **do**
  for $j = 1$ **to** J **do**
    Set $s = j + (t-1)J$
    Set $\tilde{\boldsymbol{x}}_0 = \boldsymbol{x}^{(s-1)} - \boldsymbol{u}_j^{(t-1)}$
    Set $\boldsymbol{x}^{(s)} = \text{proj}_{\mathcal{C}_j}(\tilde{\boldsymbol{x}}_0)$, and $\boldsymbol{u}_j^{(t)} = \boldsymbol{x}^{(s)} - \tilde{\boldsymbol{x}}_0$
  end for
end for
Output: $\boldsymbol{x}$

It can be successively rewritten as:

$$\begin{aligned} & \min_{z_m} \frac{1}{2}(z_m - z_{0m})^2 + \sum_{i=1}^{m-1} \min_{z_i \leq z_m} \frac{1}{2}(z_i - z_{0i})^2 \\ = \; & \min_{z_m} \frac{1}{2}(z_m - z_{0m})^2 + \sum_{i=1}^{m-1} \frac{1}{2}(\min\{z_m, z_{0i}\} - z_{0i})^2 \\ = \; & \min_{z_m} \frac{1}{2}(z_m - z_{0m})^2 + \frac{1}{2} \sum_{i \in \mathcal{I}(z_m)} (z_m - z_{0i})^2. \end{aligned} \tag{48}$$

where $\mathcal{I}(z_m) \triangleq \{i : z_{0i} \geq z_m\}$. Assuming that the set $\mathcal{I}(z_m)$ is given, the previous is a sum-of-squares problem whose solution is $z_m^* = \frac{z_{0m} + \sum_{i \in \mathcal{I}(z_m)} z_{0i}}{1 + |\mathcal{I}(z_m)|}$. The set $\mathcal{I}(z_m)$ can be determined by inspection after sorting $z_{01}, \ldots, z_{0(m-1)}$. The procedure is shown in Alg. 4.

Proposition 6 *Procedure 5 is correct.*

Proof: The proof is divided into the following parts:

1. We show that Procedure 5 corresponds to the first iteration of Dijkstra's projection algorithm (Boyle & Dykstra, 1986) applied to sets $\mathcal{A}_1$ and $[0,1]^m$;
2. We show that Dijkstra's converges in one iteration if a specific condition is met;
3. We show that with the two sets above that condition is met.

The first part is rather trivial. Dijkstra's algorithm is shown as Alg. 5; when $J = 2$, $\mathcal{C}_1 = \mathcal{A}_1$ and $\mathcal{C}_2 = [0,1]^m$, and noting that $\boldsymbol{u}_1^{(1)} = \boldsymbol{u}_2^{(1)} = 0$, we have that the first iteration becomes Procedure 5.

We turn to the second part, to show that, when $J = 2$, the fact that $\boldsymbol{x}^{(3)} = \boldsymbol{x}^{(2)}$ implies that $\boldsymbol{x}^{(s)} = \boldsymbol{x}^{(2)}$, $\forall s > 3$. In words, if at the second iteration t of Dijkstra's, the value of $\boldsymbol{x}$ does not change after computing the first projection, then it will never change, so the algorithm has converged and $\boldsymbol{x}$ is the desired projection. To see that,

consider the moment in Alg. 5 when $t = 2$ and $j = 1$. After the projection, we update $\boldsymbol{u}_1^{(2)} = \boldsymbol{x}^{(3)} - (\boldsymbol{x}^{(2)} - \boldsymbol{u}_1^{(1)})$, which when $\boldsymbol{x}^{(3)} = \boldsymbol{x}^{(2)}$ equals $\boldsymbol{u}_1^{(1)}$, *i.e.*, $\boldsymbol{u}_1$ keeps unchanged. Then, when $t = 2$ and $j = 2$, one first computes $\tilde{\boldsymbol{x}}_0 = \boldsymbol{x}^{(3)} - \boldsymbol{u}_2^{(1)} = \boldsymbol{x}^{(3)} - (\boldsymbol{x}^{(2)} - \boldsymbol{x}_0) = \boldsymbol{x}_0$, *i.e.*, the projection is the same as the one already computed at $t = 1$, $j = 2$. Hence the result is the same, *i.e.*, $\boldsymbol{x}^{(4)} = \boldsymbol{x}^{(2)}$, and similarly $\boldsymbol{u}_2^{(2)} = \boldsymbol{u}_2^{(1)}$. Since neither $\boldsymbol{x}$, $\boldsymbol{u}_1$ and $\boldsymbol{u}_2$ changed in the second iteration, and subsequent iterations only depend on these values, we have that $\boldsymbol{x}$ will never change afterwards.

Finally, we are going to see that, regardless of the choice of $\boldsymbol{z}_0$ in Procedure 5 ($\boldsymbol{x}_0$ in Alg. 5) we will always have $\boldsymbol{x}^{(3)} = \boldsymbol{x}^{(2)}$. Looking at Alg. 4, we see that after $t = 1$:

$$x_i^{(1)} = \begin{cases} \tau, & \text{if } i = m \text{ or } x_{0i} \geq \tau \\ x_{0i}, & \text{otherwise,} \end{cases} \qquad u_{1i}^{(1)} = \begin{cases} \tau - x_{0i}, & \text{if } i = m \text{ or } x_{0i} \geq \tau \\ 0, & \text{otherwise,} \end{cases}$$
$$x_i^{(2)} = [x_i^{(1)}]_{\mathbb{U}} = \begin{cases} [\tau]_{\mathbb{U}}, & \text{if } i = m \text{ or } x_{0i} \geq \tau \\ [x_{0i}]_{\mathbb{U}}\,, & \text{otherwise.} \end{cases} \tag{49}$$

Hence in the beginning of the second iteration ($t = 2$, $j = 1$), we have

$$\tilde{x}_{0i} = x_i^{(2)} - u_{1i}^{(1)} = \begin{cases} [\tau]_{\mathbb{U}} - \tau + x_{0i}, & \text{if } i = m \text{ or } x_{0i} \geq \tau \\ [x_{0i}]_{\mathbb{U}}\,, & \text{otherwise.} \end{cases} \tag{50}$$

Now two things should be noted about Alg. 4:

- If a constant is added to all entries in $\boldsymbol{z}_0$, the set $\mathcal{I}(z_m)$ remains the same, and τ and $\boldsymbol{z}$ are affected by the same constant;
- Let $\boldsymbol{z}_0'$ be such that $z_{0i}' = z_{0i}$ if $i = m$ or $z_{0i} \geq \tau$, and $z_{0i}' \leq \tau$ otherwise. Let $\boldsymbol{z}'$ be the projected point when such $\boldsymbol{z}_0'$ is given as input. Then $\mathcal{I}(z_m') = \mathcal{I}(z_m)$, $\tau' = \tau$, $z_i' = z_i$ if $i = m$ or $z_{0i} \geq \tau$, and $z_i' = z_{0i}'$ otherwise.

The two facts above allow to relate the projection of $\tilde{\boldsymbol{x}}_0$ (in the second iteration) with that of $\boldsymbol{x}_0$ (in the first iteration). Using $[\tau]_{\mathbb{U}} - \tau$ as the constant, and noting that, for $i \neq m$ and $x_{0i} < \tau$, we have $[x_{0i}]_{\mathbb{U}} - [\tau]_{\mathbb{U}} + \tau \geq \tau$ if $x_{0i} < \tau$, the two facts imply that:

$$x_i^{(3)} = \begin{cases} x_i^{(1)} + [\tau]_{\mathbb{U}} - \tau = [\tau]_{\mathbb{U}}, & \text{if } i = m \text{ or } x_{0i} \geq \tau \\ [x_{0i}]_{\mathbb{U}}\,, & \text{otherwise;} \end{cases} \tag{51}$$

hence $\boldsymbol{x}^{(3)} = \boldsymbol{x}^{(2)}$, which concludes the proof. ■