Multistage Stochastic Integer Programs: An Introduction

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Abstract. We consider linear multistage stochastic integer programs and study their functional and dynamic programming formulations as well as conditions for optimality and stability of solutions. Furthermore, we study the application of the Rockafellar-Wets dualization approach as well as the structure and algorithmic potential of corresponding dual problems. For discrete underlying probability distributions we discuss possible large scale mixed-integer linear programming formulations and three dual decomposition approaches, namely, scenario, component and nodal decomposition.

1 Introduction

Stochastic programming deals with the optimization of decision making under uncertainty over time. Typical objects of study are random optimization problems where outcomes of random data are unveiled over time, and the decisions to be optimized must not anticipate future outcomes (non-anticipativity). The latter provides a tight link to real-time optimization seen as the need for optimal “here-and-now” decision in an incomplete (or uncertain) data environment. Provided that probabilistic information on the uncertain data is available, operational models suitable for real-time optimization often may be formulated as multi-stage stochastic programs. Basic references for theory, algorithmics, and application of stochastic programming are the textbooks [7,24,34]. The edited volume [50] provides insight into recent research in the field.

Indispensability of integer requirements is a basic modeling experience in practical optimization. Like in other branches of mathematical optimization this has considerable consequences on structural properties and algorithm design in stochastic programming, too. The models best understood so far are (purely) linear stochastic programs. This is mainly due to the fact that the optimal value of a linear minimization problem is a convex function of the right-hand side and a concave function of the objective function vector. This enables application of the machinery of convex analysis in various contexts, such as duality, stability, and subgradient minimization. For an impression on these developments we refer to [12,14,36,49], with accent on theory, and to [5,42], with accent on computation.

With integer requirements, the above convexity/concavity observation is no
longer valid, and the mentioned functions become discontinuous. Thus, comparatively little is known on theory and algorithms for mixed-integer linear stochastic programs. A recent survey is provided in [26]. Impressions on developments in theory can be obtained from [2, 43, 46] and on algorithm design from [9, 11, 19, 27, 30, 33, 47], see also the Ph.D. theses [8, 32, 48]. The present paper aims at a short introduction into some essential theoretical and algorithmic issues in multi-stage stochastic integer programming. Accent is placed on introducing approaches. Proofs are omitted, with references to the original sources instead. The main topics will be modeling, approximation, and algorithms.

2 Multistage Stochastic Integer Programs

2.1 Modeling

We consider a finite horizon sequential decision process under uncertainty, in which a decision made at stage \( t \) is based only on information available at \( t (1 \leq t \leq T) \). We assume that the information is given by a discrete time stochastic process \( \{\xi_t\}_{t=1}^T \) defined on some probability space \((\Omega, \mathcal{F}, P)\) and with \( \xi_t \) taking values in \( \mathbb{R}^n \). The information available at stage \( t \) consists of the random vector \( \xi^t := (\xi_1, \ldots, \xi_t) \), and the stochastic decision \( x_t \) at stage \( t \) varying in \( \mathbb{R}^{m_t} \) is assumed to depend only on \( \xi^t \). The latter property is called nonanticipativity and is equivalent to the measurability of \( x_t \) with respect to the \( \sigma \)-algebra \( \mathcal{F}_t \subseteq \mathcal{F} \) which is generated by \( \xi^t \). Clearly, we have \( \mathcal{F}_1 \subseteq \mathcal{F}_{t+1} \) for \( t = 1, \ldots, T - 1 \) and, with no loss of generality, we may assume that \( \mathcal{F}_1 = \{\emptyset, \Omega\} \), i.e., \( \xi_1 \) and \( x_1 \) are deterministic, and that \( \mathcal{F}_T = \mathcal{F} \).

More precisely, we consider a decision model where the objective is given by expected linear costs and the constraints consist of three groups: the measurability constraints on \( x_t \), a linear constraint describing the relation between decisions at different stages, and constraints characterizing feasibility of the \( t \)-th stage decision \( x_t \). The latter constraints consist of a linear inequality constraint and of the general constraint \( x_t \in X_t \) where the (fixed) set \( X_t \) has the property that its convex hull \( \text{conv}(X_t) \) is polyhedral, allowing for mixed-integer decisions in all stages. Furthermore, the data \( \xi_t \) at stage \( t \) may enter all corresponding cost coefficients, matrices and right-hand sides. This leads to the following stochastic decision model:

\[
\min \{E\left[ \sum_{t=1}^{T} c_t(\xi_t)x_t \right] : x_t \text{ is measurable with respect to } \mathcal{F}_t, \quad (1) \}
\]

\[
x_t \in X_t, \ B_t(\xi_t)x_t \geq d_t(\xi_t), \ P - a.s., \ t = 1, \ldots, T, \quad (2)
\]

\[
\sum_{t=1}^{T} A_{tr}(\xi_t)x_t \geq g_t(\xi_t), \ P - a.s., \ t = 2, \ldots, T \quad (3)
\]

Throughout, the following is imposed: The sets \( X_t \) are nonempty and closed. The matrices \( A_{tr}(\cdot), B_t(\cdot) \) as well as the coefficients \( c_t(\cdot) \) and the right-hand
sides $d_t(\cdot)$, $g_t(\cdot)$ all depend affinely linearly on the corresponding component of $\xi$, for each $\tau = 1, \ldots, t, t = 1, \ldots, T$. In order to have the model (1)-(3) well defined, we need that the scalar products $c_t(\xi_t)x_t$ are integrable. The latter property is implied by the integrability of $\|\xi_t\|\|x_t\|$ and by the conditions $\xi_t \in L_q(\Omega, F_t, P; \mathbb{R}^n)$ and $x_t \in L_p(\Omega, F_t, P; \mathbb{R}^m)$ where $q_t, r_t \in [1, \infty]$ with $\frac{1}{q_t} + \frac{1}{r_t} = 1$. Since it is desirable to impose only weak conditions on the data process $\xi$ and since we assume later on that the set $X_t$ is bounded, we may restrict our attention to decisions $x_t \in L_\infty(\Omega, F_t, P; \mathbb{R}^m)$ and to the first order moment condition $\xi_t \in L_1(\Omega, F_t, P; \mathbb{R}^n)$ on the data at stage $t$ for each $t = 1, \ldots, T$. Then the nonanticipativity constraint (1) may be expressed equivalently as

$$x_t \in L_\infty(\Omega, F_t, P; \mathbb{R}^m) \quad \text{and} \quad x_t = E[x_t | F_{t-1}], \ t = 1, \ldots, T, \quad (4)$$

by using the conditional expectation $E[\cdot | F_{t-1}]$ with respect to the $\sigma$-algebra $F_t$. Condition (4) describes a linear subspace $N_{na}$ of the space $X^T \bigoplus_{t=1}^T L_\infty(\Omega, F_t, P; \mathbb{R}^m)$. This combination of functional and (P-a.s.) pointwise constraints in our model, i.e., the functional condition $x \in N_{na}$ and the P-a.s. constraints (2) and (3), forms the theoretical and algorithmic challenge of multistage stochastic programs. A special role is played by the two-stage case (i.e., $T=2$) where $N_{na}$ takes the specific form $N_{na} = \mathbb{R}^{m_1} \times L_\infty(\Omega, F, P; \mathbb{R}^{m_2})$. An additional complication of the model (1)-(3) is caused by the mixed-integer constraints hidden in the condition $x_t \in X_t, \ t = 1, \ldots, T$.

2.2 Multistage Models, Dynamic Programming and Optimality

We adopt the setting of the previous section and assume that $X_t$ is compact and $\xi_t \in L_1(\Omega, F_t, P; \mathbb{R}^n)$ for $t = 1, \ldots, T$. For each $\omega \in \Omega$ we define the subset $Y(\omega)$ of $X := \bigoplus_{t=1}^T \mathbb{R}^n$ by

$$Y(\omega) := \{y \in X : y_1 \in X_1, B_1(\xi_1(\omega))y_1 \geq d_1(\xi_1(\omega)), t = 1, \ldots, T, \quad \sum_{t=1}^T \sum_{\tau=1}^t A_{1\tau}(\xi_1(\omega))y_\tau \geq g_1(\xi_1(\omega)), \ t = 2, \ldots, T\}$$

and the extended real-valued function $\varphi$

$$\varphi(y_1, \ldots, y_T, \omega) := \begin{cases} \sum_{t=1}^T c_t(\xi_t(\omega))y_t, & (y_1, \ldots, y_T) \in Y(\omega), \\ +\infty, & \text{otherwise} \end{cases}$$

from $X \times \Omega$ to $(-\infty, +\infty]$. With these notations, the model (1)-(3) is equivalent to the optimization problem

$$\min \{E[\varphi(x_1, \ldots, x_T, \omega)] : x_t \text{ is measurable w.r.t. } F_t, \ t = 1, \ldots, T\}. \quad (7)$$

The real-valued function $(y, \omega) \mapsto \sum_{t=1}^T c_t(\xi_t(\omega))y_t$ is continuous in $y$ for each $\omega \in \Omega$ and measurable in $\omega$ for each $y \in X$, and the set-valued mapping $Y$
from \( \mathcal{O} \) to \( \mathcal{X} \) is closed-valued and measurable (cf. Theorem 14.36 in [40]). Hence, the function \( \varphi \) is \( \mathcal{B}(\mathcal{X}) \otimes \mathcal{F} \)-measurable (cf. Example 14.32 in [40]). Furthermore, the following estimate is valid for each \( y \in \chi_{t=1}^T X_t \), and \( \omega \in \Omega \):

\[
|\varphi(y_1, \ldots, y_T, \omega)| \leq \sum_{t=1}^{T} |\phi_t(\xi_t(\omega))| \sup_{y \in X_t} \|y\| \tag{8}
\]

Hence, \( E[\varphi(x_1, \ldots, x_T, \omega)] \) is finite for each decision \( x = (x_1, \ldots, x_T) \) such that \( x(\omega) \in \mathcal{Y}(\omega) \) for \( P \)-almost all \( \omega \in \Omega \). As in [16], we construct recursively two sequences of functions by putting \( \psi_{T+1} := \varphi \) and

\[
\varphi_t(y_1, \ldots, y_t, \omega) := E^{t}[\psi_{t+1}(y_1, \ldots, y_t, \cdot)|\mathcal{F}_t](\omega), \tag{9}
\]

\[
\psi_t(y_1, \ldots, y_{t-1}, \omega) := \inf_{\omega} \varphi_t(y_1, \ldots, y_{t-1}, \cdot, \omega), \tag{10}
\]

for \( t = T, \ldots, 1 \), and for each \( \omega \in \Omega \) and \( y_T \in X_T \), \( \tau = 1, \ldots, T \). Here, \( E^{t}[\cdot|\mathcal{F}_t] \) denotes the regular conditional expectation with respect to \( \mathcal{F}_t \). We recall that the \textit{regular conditional expectation} is a version of the conditional expectation (i.e., \( E^{t}[\cdot|\mathcal{F}_t] = E[\cdot|\mathcal{F}_t] \), \( P \)-a.s.) having the property that the mapping \( (z, \omega) \mapsto \Phi(z, \omega) := E^{t}[\Psi(z, \cdot)|\mathcal{F}_t](\omega) \) from \( Z \times \Omega \) to \( (\mathbb{R}, +, \mathbb{R}) \) is \( \mathcal{B}(Z) \otimes \mathcal{F}_t \)-measurable if \( \Psi \) is \( \mathcal{B}(Z) \otimes \mathcal{F}_t \)-measurable. Here, \( Z \) denotes a closed subset of a Euclidean space. The regular conditional expectation exists if \( \Psi \) is \( \mathcal{B}(Z) \otimes \mathcal{F}_t \)-measurable and uniformly integrable, i.e., there exists a (real) random variable \( \zeta \) with finite first moment such that \( |\Psi(z, \omega)| \leq \zeta(\omega) \) for \( z \in Z_t \) and \( \omega \in \Omega \) (see [15]). Due to condition (8), relation (9) is well defined for \( t = T \) and leads to a \( \mathcal{B}(Z) \otimes \mathcal{F}_t \)-measurable function \( \varphi_T \), where \( Z := \chi_{t=1}^T X_t \). It is shown in [16] that the relations (9) and (10) are well defined for all \( t = T, \ldots, 1 \). Furthermore, the following optimality criterion and existence result for (7) or, equivalently, for (1)--(3) are valid.

**Theorem 1.** Let the general assumptions be satisfied and assume that there exists a feasible solution of (1)--(3). Then \( \hat{x}_1^T \) is a solution of (1)--(3) if and only if

\[
\varphi_t(\hat{x}_t^T(\omega), \omega) = \psi_t(\hat{x}_t^{T-1}(\omega), \omega), \quad P \text{-a.s., } t = 1, \ldots, T. \tag{11}
\]

Moreover, there exists a solution \( \hat{x}_1 \) of the first-stage optimization problem

\[
\min \{ \varphi_1(x_1) = E[\psi_2(x_1, \omega)] : x_1 \in X_1, B_1(z_1) x_1 \geq d_1(z_1) \}. \tag{12}
\]

and, given \( \mathcal{F}_t \)-measurable functions \( \hat{x}_t \) for \( t = 1, \ldots, T-1 \), there exists an \( \mathcal{F}_t \)-measurable function \( \hat{x}_t \) such that \( \varphi_t(\hat{x}_t^T(\omega), \omega) = \psi_t(\hat{x}_t^{T-1}(\omega), \omega), \quad P \text{-a.s.} \)

The theorem is a special case of the more general results (Theorems 1 and 2) in [16], Theorem 1 implies the existence of a solution to (1)--(3) and justifies the solution approach (11) which is usually called \textit{dynamic programming approach}. Due to measurable selection arguments (cf. Chapter 14 in [40]), a feasible solution of (1)--(3) exists if the model (1)--(3) has relatively complete recourse, i.e., if \( \gamma(\omega) \neq \emptyset \).
2.3 Structure and Stability

We adopt the setting of the previous sections, denote by $\mathcal{P}(\mathcal{Z})$ the set of all Borel probability measures on some closed subset $\mathcal{Z}$ of $\mathbb{R}^s$ with $s = \sum_{i=1}^s s_i$, which is chosen such that it contains the support of $\xi$. By $\mu \in \mathcal{P}(\mathcal{Z})$ we denote the probability distribution of $\xi$. We consider the probability space $(\mathcal{Z}, \mathcal{B}(\mathcal{Z}), \mu)$ as the underlying probability space $(\Omega, \mathcal{F}, P)$ in Section 1, and define a function $f$ from $X_1 \times \mathcal{Z}$ to the extended real numbers $\mathbb{R}$ by $f(x_1, \xi) := \psi_2(x_1, \xi)$, where $\psi_2$ is defined by (9) and (10). Then the first-stage optimization problem (12) can be rewritten in the following form:

$$\min \{ \int_{\mathcal{Z}} f(x_1, \xi) \mu(d\xi) : x_1 \in X_1, B_1(\xi_1)x_1 \geq d_1(\xi_1) \} \quad (13)$$

The techniques exploited in [16] and used in the previous section imply that the integrand $f$ is $\mathcal{B}(X_1) \otimes \mathcal{B}(\mathcal{Z})$-measurable. The recursions (9) and (10) together with the Fatou Lemma for (conditional) expectations as well as lower semicontinuity properties of infima in parametric minimization (e.g. Theorem 1.17 of [40]) imply lower semicontinuity of $f$ with respect to $x_1$ and of the objective function $x_1 \mapsto \int_{\mathcal{Z}} f(x_1, \xi) \mu(d\xi)$. If the multistage model (1)–(3) has relatively complete recourse, it holds that $|f(x_1, \xi)| \leq K(\max_{i=1}^n \|\xi\|)$ for each feasible $x_1$, each $\xi \in \mathcal{Z}$ and some constant $K > 0$. Hence, the integrand has a uniform and integrable upper bound, and the objective function is finite at all feasible $x_1$. By Lebesgue’s theorem, the objective function is continuous at some feasible $x_1$ if $\mu(\{\xi \in \mathcal{Z} : f(\cdot, \xi) \text{ is not continuous at } x_1\}) = 0$. Such discontinuity sets of the integrand $f$ have been studied in [43] for the two-stage situation with fixed recourse matrix $A_{22}$ and recourse costs $c_2$.

When developing approximation schemes and algorithmic approaches for solving the model (1)–(3), the behavior of its optimal value $\text{val}(\mu)$ and of the set $\text{Sol}(\mu)$ of first-stage solutions to (1)–(3) is important when perturbing or approximating the underlying distribution $\mu$. We say that the model (1)–(3) is stable if $\text{val}(\mu)$ and $\text{Sol}(\mu)$ satisfy certain continuity properties with respect to some suitable convergence of probability measures. Here, we follow the presentation in [35] and consider the following distance

$$d_f(\mu, \nu) := \sup \{ \int_{\mathcal{Z}} f(x_1, \xi)(\mu - \nu)(d\xi) : x_1 \text{ is feasible} \} \quad (14)$$

of probability measures $\mu$ and $\nu$ belonging to the set $\mathcal{P}_1(\mathcal{Z}) := \{ \nu \in \mathcal{P}(\mathcal{Z}) : \int_{\mathcal{Z}} ||\xi|| \nu(d\xi) < \infty \}$. Then it holds for any perturbation $\nu$ of the original underlying probability distribution $\mu$ that

$$|\text{val}(\mu) - \text{val}(\nu)| \leq d_f(\mu, \nu) \quad (15)$$

$$\emptyset \neq \text{Sol}(\nu) \subseteq \text{Sol}(\mu) + \Psi(d_f(\mu, \nu))B_m, \quad (16)$$

where $B_m$ denotes the closed unit ball in $\mathbb{R}^m$, and $\Psi$ is some monotonically increasing function on $\mathbb{R}$, with $\Psi(0) = 0$, which is related to the growth.
behaviour of the objective function near the set \( \text{Sol}(\mu) \) (see [35]). While (15) represents a Lipschitz type estimate for the optimal value at \( \mu \), the relation (16) says that the sets of first-stage solutions behave upper semicontinuously at \( \mu \) with respect to \( d_f \). In general, the distance \( d_f \) is rather involved and difficult to handle. Hence, it is of considerable interest to derive estimates of \( d_f \) in terms of simpler probability metrics and to expose relations to the classical concept of weak convergence of probability measures. For two-stage models with fixed recourse matrix and costs, such results are obtained in [44] and [35]. We also refer to relevant stability studies in [2,18,49]. Altogether, such stability results justify the approximation of the underlying distribution \( \mu \) by simpler measures and provide techniques for designing approximation schemes. Next, we show that approximations by discrete measures having finitely many atoms or scenarios play a prominent role since they lead to specially structured large-scale mixed-integer linear programs.

### 2.4 Scenario Based Models

We assume throughout this section that \( \Omega \) is finite, i.e., \( \Omega = \{\omega_s\}_{s=1}^S \), \( \mathcal{F} \) is the power set of \( \Omega \) and \( P(\{\omega_s\}) = p_s \), \( s = 1, \ldots, S \). We denote by \( \xi_t^s := \xi_t(\omega_s) \) the value of the data scenario \( s \) at stage \( t \) and by \( x_t^s \) the value of the decision scenario \( s \) at \( t \) for \( s = 1, \ldots, S \), \( t = 1, \ldots, T \). Since \( \Omega \) is finite, there exists a finite subset \( \mathcal{E}_t \) of the \( \sigma \)-algebra \( \mathcal{F}_t \), for each \( t = 1, \ldots, T \), such that \( \mathcal{E}_t \) is a partition of \( \Omega \) and that the smallest \( \sigma \)-algebra containing \( \mathcal{E}_t \) is just \( \mathcal{F}_t \). Then the conditional expectation w.r.t. \( \mathcal{F}_t \) in the nonanticipativity condition (4) takes the form

\[
E[x_t^s | \mathcal{F}_t] = \sum_{C \in \mathcal{E}_t} \frac{1}{P(C)} \int_C x_t(\omega) P(d\omega) \chi_C
\]

where \( \chi_C \) denotes the characteristic function of the set \( C \in \mathcal{E}_t \). Hence, the nonanticipativity condition (4) is equivalent to the following equality constraints

\[
x_t^s = \sum_{C \in \mathcal{E}_t} \left( \sum_{\omega_s \in C} p_s \right)^{-1} \sum_{\omega_s \in C} p_s x_t^s, \quad \sigma = 1, \ldots, S, \ t = 1, \ldots, T. \quad (18)
\]

Clearly, for \( t = 1 \) we have \( \mathcal{E}_1 = \{\Omega\} \) and, hence, condition (18) is equivalent to the equations \( x_1^s = \sum_{\sigma=1}^S p_\sigma x_1^\sigma, \ \sigma = 1, \ldots, S \), i.e., to \( x_1^1 = \ldots = x_1^S \). Hence, the multistage stochastic program (1)-(3) takes the following form
which will be called its scenario formulation:

\[
\min \left\{ \sum_{s=1}^{S} \sum_{t=1}^{T} p_{s,t}(\xi_s^t)x_s^t : x \text{ satisfies the constraints (18)}, \right. \\
\left. x_s^t \in X_t, B_t(\xi_s^t)x_s^t \geq d_t(\xi_s^t), s = 1, \ldots, S, t = 1, \ldots, T, \right. \\
\left. \sum_{s=1}^{S} A_t(\xi_s^t)x_s^t \geq g_t(\xi_s^t), s = 1, \ldots, S, t = 2, \ldots, T \right\}
\]

Since \( \mathcal{F}_t \subseteq \mathcal{F}_{t+1} \), every element of \( \mathcal{E}_{t+1} \) can be represented as the union of certain elements of \( \mathcal{E}_t \). Furthermore, formula (17) shows that the number of elements in \( \mathcal{E}_t \) coincides with the number of realizations of \( \xi \) and \( x \) at period \( t \), respectively. Hence, representing the relations between the elements of \( \mathcal{E}_t \) and \( \mathcal{E}_{t+1} \) for \( t = 1, \ldots, T-1 \), leads to a tree having the same structure as the sets of scenarios of \( \xi \) and \( x \), respectively. Therefore, such a tree is called scenario tree. It is based on a finite set \( \mathcal{N} \subseteq \mathbb{N} \) of nodes. Fig. 1 shows an example of a scenario tree where the \( t_k \) denote the branching points of the tree. The root node \( n = 1 \) stands for period \( t = 1 \). Every other node \( n \) has a unique predecessor node \( n_- \) and a transition probability \( \pi_{n/n_-} > 0 \), which is the probability of \( n \) being the successor of \( n_- \). The probability \( \pi_n \) of each node \( n \) is given recursively by \( \pi_1 = 1, \pi_n = \pi_{n/n_-} \pi_{n_-} , n > 1 \). We denote by \( \mathcal{N}_n(n) \) the set of successors to node \( n \), by \( \text{path}(n) \) the path from the root to node \( n \) and by \( t(n) \) its length, i.e., \( t(n) := \text{card}(\text{path}(n)) \). \( \mathcal{N}_n \) denotes the set \( \{ n \in \mathcal{N} : t(n) = t \} \), and it holds \( \sum_{n \in \mathcal{N}_t} \pi_n = 1 \) for each period \( t \). Nodes \( n \) with \( \mathcal{N}_n(n) = \emptyset \) are called leaves; they constitute the terminal set \( \mathcal{N}_T \). A scenario corresponds to a path from the root node to a leaf. Clearly,
it holds that \( \text{card}(\mathcal{N}_T) = S \) and \( \{\pi_n\}_{n \in \mathcal{N}_T} = \{\rho_k\}_{k=1}^S \). Conversely, given these scenario probabilities, the remaining node and transition probabilities are generated recursively by \( \pi_n := \sum_{n \in \mathcal{N}_T(n)} \pi_{n^-} \), \( \pi_{n^-/n} := \pi_{n^-}/\pi_n \) for \( n^- \in \mathcal{N}_T(n) \). We use the following notation for the sequence of predecessors of any node \( n \in \mathcal{N} \): \( n_0 := n \), \( n_{-1} := n_{-} \) if \( n > 1 \), \( n_{-(k+1)} := (n_{-k})_{-} \) if \( t(k) > 1 \). Note that \( t(n_{-k}) = t(n) - k \) for \( k = 1, \ldots, t(n) - 1 \). Furthermore, we denote by \( \{\xi^n\}_{n \in \mathcal{N}_T} \) the realizations of \( \xi_t \) and by \( \{x^n\}_{n \in \mathcal{N}_T} \) the realizations of \( x_t \). After these preparations the scenario tree formulation of the multistage stochastic program reads:

\[
\min \left\{ \sum_{n \in \mathcal{N}_T} \pi_n c_{t(n)}(\xi^n) x^n : x^n \in X_t(n), B_{t(n)}(\xi^n) x^n \geq d_{t(n)}(\xi^n), \right. \\
\left. \vphantom{\min} \sum_{n \in \mathcal{N}_T} \pi_{n^-} A_{t(n), t(n)_{-}}(\xi^n) x^n \geq g_{t(n)}(\xi^n), n \in \mathcal{N}_T \right\}
\]  

Both formulations of the multistage stochastic program will be used for the description of decomposition approaches. We recall that the nonanticipativity condition appears explicitly in the scenario formulation (19), but disappears in the scenario tree formulation (20) because it is incorporated into the tree construction. Since it holds that \( \text{card}(\mathcal{N}) \ll TS \), the dimensions of both model formulations are quite different. More precisely, the model (19) contains \( \sum_{t=1}^T m_t S \) decision variables and \( \sum_{t=1}^T (m_t + k_t + r_t) S \) linear constraints, whereas the model (20) contains \( \sum_{n \in \mathcal{N}_T} m_{t(n)} \) decisions and \( \sum_{t \in \mathcal{N}_T} (k_{t(n)} + r_{t(n)}) \) linear constraints. Here, \( k_t \) and \( r_t \) denote the dimensions of \( d_t(\cdot) \) and \( g_t(\cdot) \), respectively, for \( t = 1, \ldots, T \).

### 2.5 Dualization and the Convex Case

We assume \( \xi \in \times_{t=1}^T L_\infty(\Omega, \mathcal{F}, P; \mathbb{R}^n) \) and consider the multistage stochastic integer program of Section 1 as an abstract (infinite) optimization problem in the Banach space \( \times_{t=1}^T L_\infty(\Omega, \mathcal{F}, P; \mathbb{R}^{m_t}) \), i.e., in the form

\[
\min \{ E \left[ \sum_{t=1}^T c_t(\xi_t) x_t \right] : x \in \times_{t=1}^T L_\infty(\Omega, \mathcal{F}, P; \mathbb{R}^{m_t}), x \in \mathcal{N}_T \},
\]

\[
x_t \in X_t, B_t(\xi_t) x_t \geq d_t(\xi_t), P - a.s., t = 1, \ldots, T,
\]

\[
\sum_{t=1}^T A_t(\xi_t) x_t \geq g_t(\xi_t), P - a.s., t = 2, \ldots, T
\]

Let \( F(\cdot) \) denote the objective function, i.e., \( F(x) := E[\sum_{t=1}^T c_t(\xi_t) x_t] \).

Our aim is to introduce a Lagrangian associated with the essential groups of constraints of problem (21)-(23), namely, the (functional) nonanticipativity constraint \( x \in \mathcal{N}_T \), the \( k_t \) coupling constraints \( B_t(\xi_t) x_t \geq d_t(\xi_t) \) and \( r_t \) dynamic constraints (23). We make use of the concepts and results of
[38] and introduce the following sets $A_1 := \{ \lambda_1 \in \times_{t=1}^T L_1(\Omega, \mathcal{F}, P; \mathbb{R}^{m_t}) : E[\lambda_1 | \mathcal{F}_t] = 0, P - a.s., \ t = 1, \ldots, T \}$, $A_2 := \{ \lambda_2 \in \times_{t=1}^T L_1(\Omega, \mathcal{F}, P; \mathbb{R}^{h_t}) : \lambda_2 \geq 0, P - a.s. \}$ and $A_3 := \{ \lambda_3 \in \times_{t=1}^T L_1(\Omega, \mathcal{F}, P; \mathbb{R}^{t_t}) : \lambda_3 \geq 0, P - a.s. \}$ of Lagrange multipliers. The sets $A_2$ and $A_3$ are convex cones and $A_1$ is a linear space which is complementary to the nonanticipativity subspace $\mathcal{N}_{na}$ with respect to the dual pairing $\langle \cdot, \cdot \rangle$ of $L_1$ and $L_\infty$, i.e., it holds $\langle \lambda_1, x \rangle := E[\sum_{t=1}^T \lambda_1 t x_t] = 0$ for all $\lambda_1 \in A_1$ and $x \in \mathcal{N}_{na}$.

The Lagrangian is defined to be the function

$$L(x, \lambda) := E\left[ \sum_{t=1}^T \left\{ c_t(\xi_t) x_t - \lambda_1 t x_t + \lambda_2 t (d_t(\xi_t) - B_t(\xi_t) x_t) \right\} \right]$$

$$+ \sum_{t=2}^T \lambda_3 g_t(\xi_t) - \sum_{t=1}^T A_{t-1}(\xi_t) x_t \right]$$

from $\times_{t=1}^T L_\infty(\Omega, \mathcal{F}, P; \mathbb{R}^{m_t}) \times A$ to $\mathbb{R}$, where $A := \times_{t=1}^T A_i$. The dual function $D$ from $A$ to $\mathbb{R}$ is defined by

$$D(\lambda) := \inf \left\{ L(x, \lambda) : x \in \times_{t=1}^T L_\infty(\Omega, \mathcal{F}, P; \mathbb{R}^{m_t}), \ x_t \in X_t, P - a.s., t = 1, \ldots, T \right\},$$

and the dual problem associated with (21)–(23) is

$$\max D(\lambda) : \lambda \in A.$$  

(26)

We assume again that the sets $X_t$, $t = 1, \ldots, T$, are compact. Then the Lagrangian $L$ and the dual function $D$ are well-defined, $D$ is concave and the weak duality estimate

$$D(\lambda) \leq F(x)$$

for all $\lambda \in A$ and all $x$ satisfying (21)–(23). (27)

is valid. In the following, we say that the model (21)–(23) is strictly feasible if there exist $x \in \mathcal{N}_{na}$ and $\varepsilon > 0$ such that

$$\tilde{x}_t + \varepsilon B_m \subseteq \text{conv}(X_t), B_1(\xi_t) \tilde{x}_t \geq d_t(\xi_t) + \varepsilon, P - a.s., t = 1, \ldots, T,$$

$$\sum_{t=1}^T A_{t-1}(\xi_t) \tilde{x}_t \geq g_t(\xi_t) + \varepsilon, P - a.s., t = 2, \ldots, T,$$

where $B_m$ denotes the closed unit ball in $\mathbb{R}^m$. Then we conclude from Theorem 1 and from Theorem 3 of [39] that the following holds.

**Theorem 2.** Assume that the general assumptions are satisfied, that the sets $X_t$, $t = 1, \ldots, T$, are convex compact and that the model (21)–(23) has relatively complete recourse and is strictly feasible. Then there exist optimal solutions $\bar{x}$ to (26) and $\bar{\tilde{x}}$ to (21)–(23), and it holds $D(\bar{\lambda}) = F(\bar{\tilde{x}})$.
Since the sets $X_t$, $t = 1, \ldots, T$, fail to be convex, such a duality result is not available in our setting and, due to (27), we are faced with a duality gap

$$DG := F(\tilde{x}) - \sup_{\lambda \in A} D(\lambda) \geq 0. \quad (28)$$

This inequality is strict, in general. On the other hand, in case of a discrete underlying probability distribution, the theory of Lagrangian relaxation in mixed-integer linear programming (cf. e.g. Chapter II.3.6 of [31]) implies that the optimal value of (26) is greater than or equal to the optimal value of the linear programming relaxation to (19) or (20). In other words, the lower bound obtained by dualizing constraints is never worse the bound obtained by relaxing the integer requirements.

So far we have associated Lagrange multipliers with nonanticipativity, coupling as well as dynamic constraints. Of course, it is also of interest to consider restricted Lagrangians and restricted duals by associating multipliers with one or with two of these three groups of constraints, only. For such restricted dualization schemes, duality results for the convex case that are similar to Theorem 2 may be derived as well (see [37] for dualizing the nonanticipativity constraints and [38] for other inequality constraints). It is worth recalling that the duality gap increases when dualizing additional constraints (see Section 3.1 in [29]). Since small duality gaps are of algorithmic interest, we take a closer look at dualization schemes where either nonanticipativity or coupling or dynamic constraints are associated with Lagrange multipliers. We denote the corresponding dual functions from $A_i$ to $\mathbb{R}$ by $D_i$ for $i = 1, 2, 3$ and start with dualizing nonanticipativity constraints, i.e.,

$$D_i(\lambda_i) := \inf \{ E[\sum_{t=1}^T (c_t(\xi_t|x_t) - \lambda_{1t} x_t) ] : x_t \in X_t, B_t(\xi_t|x_t) \geq d_t(\xi_t), P - \text{a.s., } t = 1, \ldots, T, \sum_{t=1}^T A_{i\tau}(\xi_t) x_{\tau} \geq g_t(\xi_t), P - \text{a.s., } t = 2, \ldots, T \}$$

$$= E[\inf \{ \sum_{t=1}^T (c_t(\xi_t|x_t) - \lambda_{1t} x_t) : x_t \in X_t, B_t(\xi_t|x_t) \geq d_t(\xi_t), t = 1, \ldots, T, \sum_{t=1}^T A_{i\tau}(\xi_t) x_{\tau} \geq g_t(\xi_t), t = 2, \ldots, T \}].$$

where the infimum and expectation may be interchanged since the minimization problem only contains $P$-a.s. pointwise constraints (see e.g. Theorem 14.60 of [40]). Hence, the multistage stochastic program defining $D_i$ decomposes into pathwise minimization problems. This effect becomes more transparent if the underlying probability distribution of $\xi$ is discrete, i.e., if $\Omega = \{\omega_1, \ldots, \omega_s\}$. Adopting the notation of Section 2.4, the dual function
takes the form
\[
D_1(\lambda_1) := \sum_{s=1}^{S} p_s \inf \{ \sum_{t=1}^{T} [\ell_t(\xi^x_t)x^x_t - \lambda^x_t x^x_t] : x^x_t \in X_t, \quad (29) \}
\]
\[
B_t(\xi^x_t)x^x_t \geq d_t(\xi^x_t), \ t = 1, \ldots, T,
\]
\[
\sum_{t=1}^{T} A_{tT}(\xi^x_t)x^x_t \geq g_t(\xi^x_t), \ t = 2, \ldots, T),
\]
where \( \lambda_1 \in A_1 \) has the scenarios \( \{ \lambda^x_1 \}^{T}_{t=1} \) with probabilities \( p_s \) for \( s = 1, \ldots, S \) and \( A_1 \) is given by the linear subspace
\[
A_1 = \{ \lambda_1 : E[\lambda_1|F_T] = \sum_{C \in \mathcal{E}_t} (\sum_{s=1}^{S} p_s)^{-1} \sum_{s=1}^{S} p_s \lambda^x_t \chi_C = 0, t = 1, \ldots, T \}
\]
\[
= \{ \lambda_1 : \sum_{s=1}^{S} p_s \lambda^x_t = 0, C \in \mathcal{E}_t, t = 1, \ldots, T \} \quad (30)
\]
of the Euclidean space of dimension \( (\sum_{t=1}^{T} k_t + \sum_{t=2}^{T} r_t) S \). Since \( \mathcal{E}_1 = \{ \Omega \} \) and \( \mathcal{E}_T = \{ \{ \omega_1 \}, \ldots, \{ \omega_S \} \} \), the conditions for \( t = 1 \) and \( t = T \) in (30) are equivalent to \( \sum_{s=1}^{S} p_s \lambda^x_1 = 0 \) and \( \lambda^x_T = 0, s = 1, \ldots, S \), respectively. We note that the constraint \( \lambda_1 \in A_1 \) means that each \( \lambda_1 \neq 0 \) is anticipative, i.e., \( \lambda_1 \) is not \( \mathcal{F}_t \)-measurable for some \( t \) (see also the example in [22]). Sometimes, one might find it more convenient that the dual function is defined and maximized on the whole space, i.e., without regard to the subspace constraint \( \lambda_1 \in A_1 \). This can be done by replacing \( \lambda_1 \) in the right-hand side of (29) by \( \lambda_1 - E[\lambda_1|F_T] \) for \( t = 1, \ldots, T \). Then the subspace constraint for the multiplier is automatically satisfied and the dual maximization problem is unconstrained.

Next we consider dualizations of certain inequality constraints by some multiplier, but leave the nonanticipativity constraint untouched. In contrast to the anticipativity of multipliers in the previous case, the multipliers may now be chosen nonanticipative, i.e., as elements of \( \times_{t=1}^{T} L_1(\Omega, \mathcal{F}_t, P) \). This is due to the linear separability properties of (21)–(23) (Theorem 7 of [38]). In particular, when dualizing the coupling constraints, the restricted dual function
\[
D_2(\lambda_2) := \inf \{ E[\sum_{t=1}^{T} (c_t(\xi_t)x_t + \lambda_2(d_t(\xi_t) - B_t(\xi_t)x_t))] : x \in \mathcal{A}_{ad}, \quad (31) \}
\]
\[
x_t \in X_t, \sum_{t=1}^{T} A_{tT}(\xi_t)x_T \geq g_t(\xi_t), \ P - a.s., \ t = 2, \ldots, T,
\]
has to be maximized on the convex cone \( \mathcal{A}_2 := \{ \lambda_2 \in \times_{t=1}^{T} L_1(\Omega, \mathcal{F}_t, P; \mathbb{R}^k) : \lambda_2 \geq 0, P - a.s., \ t = 1, \ldots, T \} \). Dualizing the dynamic constraints leads to
maximizing the restricted dual

\[
D_3(\lambda_3) := \inf \left\{ \sum_{t=1}^{T} \alpha_t(x_t) x_t + \sum_{t=2}^{T} \lambda_3 t (g_t(x_t) - \sum_{\tau=1}^{t} A_{\tau}(x_{\tau} x_{\tau}) : \right. \\
\left. x \in \mathcal{N}_{\text{no}}, x_t \in X_t, B_t(x_t) x_t \geq d_t(x_t), P - a.s., t = 1, \ldots, T \right\}
\]

subject to the convex cone \( A_3 := \{ \lambda_3 \in \times_{t=1}^{T} L_1(\Omega, \mathcal{F}_t, P; \mathbb{R}^m) : \lambda_3 \geq 0, P - a.s., t = 1, \ldots, T \} \). Clearly, both optimization problems on the right-hand sides of (31) and (32), respectively, are stochastic integer programs. While the program in (31) exhibits the typical multistage structure, the specific feature of the program in (32) is the lack of a dynamic constraint. In Sections 3.2 and 3.3 we gain further information on these programs in case of a discrete underlying probability distribution, i.e., when the data, decisions and multipliers form scenario trees.

3 Decomposition Methods

Due to the enormous size of scenario based models in multi-stage stochastic programming, decomposition is the method of choice when it comes to numerical solution. This is further enhanced by special structures met, both in the scenario formulation (19) and in the scenario tree formulation (20) of multi-stage stochastic programs. If integer requirements are missing in (21) - (23), powerful convexity and duality results (cf. Theorem 2) are the basis of efficient decomposition methods. These methods can be subdivided into primal and dual ones.

Primal decomposition methods employ the scenario tree formulation (20). Starting from the root node, primal proposals are passed down the tree where they are used to compute so called feasibility and optimality cuts that are passed upward to be included into convex optimization problems whose solutions lead to updated primal proposals that are again passed down the tree, and so on. This procedure (nested decomposition) is enhanced by regularization and cut deletion. Its mathematical backbone is convexity, in particular ideas from the area of bundle-trust and proximal point methods.

Dual decomposition circles around duality results such as Theorem 2. The approaches discussed in Section 2.5 then all benefit from a zero duality gap. Particular attention has been paid to dualizing nonanticipativity in the framework of augmented Lagrangians and related proximal point algorithms (progressive hedging, cf. [39]). The survey papers [5,42] provide further insights into both primal and dual decomposition of multi-stage stochastic linear programs.

With integer requirements in (21) - (23) the mentioned powerful convexity and duality results are lost. Approaches to decomposition, that have proven efficient for purely linear models, have to be rethought from their very beginnings.
The impact of integrality on *primal decomposition* is twofold: Feasibility and optimality cuts can no longer be obtained as linear functionals but as merely subadditive functionals instead. Primal proposals can no longer be obtained via convex programs but via merely lower semicontinuous (discontinuous) nonconvex programs instead. For algorithmic realization this leads to obstacles impossible to overcome with existing methods [8,11]. Two-stage models have been tackled with limited success by solving the mentioned lower semicontinuous programs via enumeration [45] or branch-and-bound [1] and exploiting problem similarities in the second stage.

The impact of integrality on *dual decomposition* has already been mentioned in Section 2.5: Theorem 2 is no longer valid, and we face a non-zero duality gap (28). Although progressive hedging then is no longer formally justified, quite satisfactory results have been observed empirically for specific applications [30,47].

In what follows, we will return to the dualization schemes introduced in Section 2.5 in case that the underlying probability distribution is discrete. We will discuss the solution of the corresponding dual maximization problems

\[
\max \{ D_i(\lambda_i) : \lambda_i \in A_i \} \quad (i = 1, 2, 3)
\]

by subgradient type methods and examine the decoupling potential of the different dualizations. Under the conditions imposed in Section 2.5 the dual functions \( D_i \) are finite, concave and polyhedral. They have the form

\[
D_i(\lambda_i) = \inf_x \{ F(x) + \langle \lambda_i, G_i(x) \rangle \},
\]

where \( F \) is the objective function, \( G_i \) is some affine linear function from \( L_\infty \) to \( L_\infty \), and \( \langle \cdot, \cdot \rangle \) denotes the dual pairing of \( L_1 \) and \( L_\infty \). Hence, \( G_i(x_i(\lambda_i)) \) is a subgradient of \( D_i \) if \( x_i(\lambda_i) \) is a solution to the minimization problem (33) defining \( D_i \). Furthermore, the solution sets of the dual problems are nonempty since their objectives are polyhedral and their supremum finite. Therefore, subgradient bundle methods may be used for solving the duals [23,28,25]. Let us consider the *proximal bundle method* [17,23,25] in more detail. Starting from an arbitrary point \( \lambda_i^1 = \lambda_i^k \in A_i \), this method generates a sequence \( \{ \lambda_i^k \}_{k \in \mathbb{N}} \) in \( A_i \) converging to some solution of the dual problem, and trial points \( \lambda_i^k \) for evaluating the solutions \( x_i^k = x_i(\lambda_i^k) \) of (33), the subgradients \( G_i(x_i^k) \) of \( D_i \) and its linearizations

\[
D_i^k(\cdot) := D_i(\lambda_i^k) + \langle \cdot, -\lambda_i^k, G_i(x_i^k) \rangle \geq D_i(\cdot).
\]

Iteration \( k \) uses the polyhedral model \( D_{ik}(\cdot) := \min_{i \in A_i} D_i^k(\cdot) \) with \( k \in \mathcal{N}_k \subset \{1, \ldots, k\} \) for finding the next trial point \( \lambda_i^{k+1} \) as a solution of the quadratic subproblem

\[
\max \{ D_{ik}(\lambda) - \frac{1}{2} u_k |\lambda - \lambda_i^k|^2 : \lambda \in A_i \},
\]

(34)
where the proximity weight \( u_k > 0 \) and the penalty term \( J(x) := \langle s, x \rangle \) should keep \( \lambda_{k+1} \) close to the prox-center \( \lambda_k \). An ascent step to \( \lambda_{k+1} = \lambda_k + \delta_k \) occurs if

\[
D_k(\lambda_{k+1}) \geq D_k(\lambda_k) + \kappa \delta_k,
\]

where \( \kappa \in (0, 1) \) is a fixed Armijo-like parameter and \( \delta_k := D_k(\lambda_{k+1}) - D_k(\lambda_k) \geq 0 \) is the predicted ascent (if \( \delta_k = 0 \) then \( \lambda_{k+1} \) is a solution and the method may stop). Otherwise, a null step \( \lambda_{k+1} = \lambda_k \) improves the next model \( D_{k+1} \) with the new linearization \( D_k(\lambda_k) \). The choices of the weights \( u_k \) and of the index set \( N^{k+1} \) are discussed in [17, 25] (see also Section 3.4 of [19]). The quadratic subproblem (34) is essentially influenced by the dual pairing \( \langle s, x \rangle \). The latter reads \( \langle s, x \rangle = \sum_{n=1}^N \rho_n \sum_{t=1}^T \lambda_{nt} y_t^s \) and \( \langle s, y \rangle = \sum_{n \in N} \pi_n \lambda_n y^n \) for the scenario and the node formulations, respectively.

### 3.1 Scenario Decomposition

Scenario decomposition rests on the dualization of nonanticipativity constraints if the probability distribution of \( \xi \) is discrete. This leads to the dual maximization problem

\[
\max \{ D_1(\lambda_1) : \lambda_1 \in A_1 \}
\]

where \( D_1 \) and \( A_1 \) are defined as in (29), (30) of Section 2.5. Since the computation of \( D_1 \) decomposes into solving pathwise minimization problems, function values and subgradients of \( D_1 \) are obtained by solving the single-scenario problems

\[
\begin{align*}
\min \sum_{t=1}^T \left\{ c_t(\xi_t^s)x_t^s - \lambda_{1t}^s x_t^s \right\} : & x_t^s \in X_t, \\
B_t(\xi_t^s)x_t^s & \geq d_t(\xi_t^s), \quad t = 1, \ldots, T, \\
\sum_{t=1}^T A_{tr}(\xi_t^s)x_t^s & \geq g_{r}(\xi_t^s), \quad t = 2, \ldots, T
\end{align*}
\]

for all \( s = 1, \ldots, S \).

Indeed, if \( \bar{x}^s, s = 1, \ldots, S \), denote optimal solutions to these problems, then

\[
D_1(\lambda_1) = \sum_{s=1}^S \rho_s \sum_{t=1}^T \left\{ c_t(\xi_t^s)\bar{x}_t^s - \lambda_{1t}^s \bar{x}_t^s \right\}
\]

and \( \tilde{G}_1(\bar{x}) = \bar{x} \) is a subgradient of \( D_1 \) at \( \lambda_1 \), where \( \bar{x} \) has the scenarios \( \bar{x}^s, s = 1, \ldots, S \). Compared with the scenario formulation (19) of the multi-stage stochastic program (1)-(3), which is a mixed-integer linear program in dimension \( S \cdot \sum_{t=1}^T m_t \), the above single-scenario problems are \( S \) mixed-integer linear programs each of dimension \( \sum_{t=1}^T m_t \), only.

In view of (28), solving (35) provides a lower bound to the optimal value of the multi-stage stochastic integer program (19). If the single-scenario solutions
$x_1^*, \ldots, x_T^*$ for the optimal $\lambda_1$ in (35) fulfilled the nonanticipativity constraints then $\bar{x}$ would be optimal to (19). In general, however, one faces a non-zero duality gap (28). Therefore the lower bounding has to be accompanied by upper bounding procedures resting on the generation of “promising” feasible solutions. This can be accomplished by primal heuristics starting from the results of the dual optimization, i.e., from single-scenario solutions $x_1^*, \ldots, x_T^*$ corresponding to optimal or nearly optimal $\lambda_1$.

An algorithmic realization of scenario decomposition for the case $T = 2$, i.e., for two-stage stochastic integer programs, has been proposed in [8–10]. The nonanticipativity constraints then read $x_1^* = \sum_{\sigma=1}^{S} p_\sigma x_1^\sigma$, $\sigma = 1, \ldots, S$. In [8–10], the equivalent representation $x_1^1 = \ldots = x_1^S$ is employed, and the scenario formulation (19) is set up with (18) replaced by $x_1^1 = \ldots = x_1^S$. Then, the usual Lagrangian relaxation of mixed-integer linear programming is performed with respect to the constraints $x_1^1 = \ldots = x_1^S$. In particular, this leads to a non-probabilistic Lagrangian, in contrast to the probabilistic Lagrangian (24) introduced in Section 2.5. As a consequence, the Lagrangian dual of [8–10] is unconstrained and lives in dimension $(S-1) \cdot m_1$. In the setting of Section 2.5, cf. (30), we obtain a dual in dimension $S \cdot m_1$ constrained by $\sum_{\sigma=1}^{S} p_\sigma \lambda_1^\sigma = 0$, i.e., essentially an unconstrained program in dimension $(S-1) \cdot m_1$ as well.

In [8–10], the scheme of lower and upper bounding outlined above is further enhanced by embedding into a branch-and-bound algorithm in the spirit of global optimization. As stated in (13), the stochastic program can be rewritten as a nonconvex global optimization problem. In the branching part of the algorithm, the feasible region of (13) is subdivided. On each member of the subdivision, the bounding part employs dualization of nonanticipativity for the lower and a primal heuristic for the upper bounds. For further details on scenario decomposition for two-stage stochastic integer programs we refer to [21].

Only little is known about algorithmic realizations of scenario decomposition for multi-stage stochastic integer programs with $T > 2$. First experiences on extending the approach of [8–10] will be reported in [4].

### 3.2 Component Decomposition

Dualization of component coupling constraints results in the dual maximization problem

$$
\max \{ D_2(\lambda_2) : \lambda_2 \in A_2 \},
$$

where $D_2$ and $A_2$ are defined in Section 2.5. We assume that the underlying probability distribution of the data process $\xi$ is discrete and, hence, given in form of a scenario tree $\{\xi^n\}_{n \in N}$, where $N$ denotes the finite set of nodes. The notation of Section 2.4 is used, and we denote by $x = \{x^n\}_{n \in N}$ the decision scenario tree and by $\lambda_2 = \{\lambda_2^n\}_{n \in N}$ the multiplier scenario tree. Then the
dual function (31) may be rewritten in the following form (see also (20)):

\[
D_2(\lambda_2) := \inf \left\{ \sum_{n \in \mathcal{N}} \pi_n \left[ c_{t(n)}(\xi^n) x^n + \lambda_2^n (d_{t(n)}(\xi^n) - B_{t(n)}(\xi^n) x^n) \right] \right\} \quad (36)
\]

where \( \lambda_2 \in \Lambda_2 = \{ (\lambda_2^n)_{n \in \mathcal{N}} : \lambda_2^n \geq 0, n \in \mathcal{N} \} \). In order to demonstrate the component decoupling potential hidden in \( D_2 \), we assume that \( X_t \) has the specific structure \( X_t = \times_{i=1}^m X_{ti} \), where the \( X_{ti} \) are closed subsets of \( \mathbb{R} \), that \( m_t = m \), \( k_t = k \) and \( r_t = m_r \) for \( t = 1, \ldots, T \) and some \( r \in \mathbb{N} \), and that the matrices \( A_{t\tau}(\cdot) \) are block-diagonal with \( m \) blocks \( a_{t\tau}^i(\cdot) \in \mathbb{R}^r \) for \( i = 1, \ldots, m \). In particular, this condition means that the constraints in (36) are expressible as componentwise constraints. We denote by \( c_{t\tau}^i(\cdot) \) the \( i \)-th component of \( c_{t\tau}(\cdot) \), by \( g_{t\tau}^i(\cdot) \in \mathbb{R}^r \) the \( i \)-th component vector of \( g_{t\tau}(\cdot) \), and by \( b_{t\tau}^i(\cdot) \) the \( i \)-th column of the matrix \( B_{t\tau}(\cdot) \). With \( x^n_i \) denoting the \( i \)-th component of \( x^n \), we obtain by exchanging summation w.r.t. \( n \) and \( i \)

\[
D_2(\lambda_2) = \inf \left\{ \sum_{n \in \mathcal{N}} \pi_n \left[ \sum_{i=1}^m c_{t(n)}^i(\xi^n) x^n_i + \lambda_2^n b_{t(n)}^i(\xi^n) x^n_i + \lambda_2^n d_{t(n)}(\xi^n) \right] \right\} \quad (37)
\]

where the functions \( D_{2i}, i = 1, \ldots, m \), from \( \Lambda_2 \) to \( \mathbb{R} \) are defined by

\[
D_{2i}(\lambda_2) = \inf \left\{ \sum_{n \in \mathcal{N}} \pi_n \left[ c_{t(n)}^i(\xi^n) x^n_i + \lambda_2^n b_{t(n)}^i(\xi^n) x^n_i \right] : x^n_i \in X_{ti}^i \right\} \quad (38)
\]

By specifying (33) we obtain that \( G_2(\bar{x}) = \{ d_{t(n)}(\xi^n) - \sum_{i=1}^m b_{t(n)}^i(\xi^n) x^n_i \}_{n \in \mathcal{N}} \) is a subgradient of \( D_2 \) at \( \lambda_2 \), where \( \bar{x} = \{ x^n_i \}_{n \in \mathcal{N}} \) is a solution of (37). The dual function (36), which is defined by a multistage stochastic integer program of dimension \( m|\mathcal{N}| \), decomposes into \( m \) functions each given by a multistage stochastic integer program of dimension \( |\mathcal{N}| \). Since the dimension of the dual problem is \( k|\mathcal{N}| \), the computational potential of this dualization approach takes effect in situations, where the number \( k \) of coupling constraints to be dualized is much smaller than the decision dimension \( m \) (i.e., \( k < m \)) and where the \( m \) subproblems (37) of dimension \( |\mathcal{N}| \) can be solved.
much faster than the original multistage model of dimension $m|\mathcal{N}|$. The latter could appear, for example, if complex mixed-integer models decompose into pure integer and pure linear programs.

Component decomposition has been applied successfully under the label Lagrangian relaxation of coupling constraints to solving hydro-thermal power management models under data uncertainty. Lagrangian relaxation has a long tradition for solving (deterministic) unit commitment problems of power systems operation. Recently, this technique has been extended to stochastic power management models, where the stochasticity enters the model, for example, via the electric load, streamflows to hydro units, and electricity prices.

When letting the production decisions of individual power units play the role of components, the above dualization scheme leads to a decomposition into single (thermal or hydro) power unit models. Such approaches for determining lower bounds have been proposed and implemented in [3, 13, 19, 33, 41]. In [19, 20, 32] encouraging numerical results and computing times have been reported for both solving the dual and determining a nearly optimal primal solution by a Lagrangian based heuristic.

### 3.3 Nodal Decomposition

Finally, we return to the dualization of the dynamic constraints of (21) - (23) in case of a discrete underlying probability distribution and show that the dual function exhibits a nodewise decoupling structure. We let $D_{3}$ and $A_{3}$ be defined as in Section 2.5 and consider the corresponding dual problem

\[
\max\{D_{3}(\lambda_{3}) : \lambda_{3} \in A_{3}\}.
\]

Let \(\{\xi^{n}\}_{n \in \mathcal{N}}\) be the scenario tree representing the data process \(\xi, \mathcal{N}\) the finite set of nodes, \(\{\pi_{n}\}_{n \in \mathcal{N}}\) the node probabilities, and \(\{x^{n}\}_{n \in \mathcal{N}}\) and \(\{\lambda_{3}^{n}\}_{n \in \mathcal{N}}\) the corresponding scenario trees of the decision and of the multiplier process, respectively. Using the notation of Section 2.4, the dual function $D_{3}$ takes the following scenario tree representation

\[
D_{3}(\lambda_{3}) = \inf \{ \mathbf{c}^{T} x^{1} + \sum_{n \in \mathcal{N} \setminus \{1\}} \pi_{n} \left[ \phi_{\tau(n)}(\xi^{n}) x^{n} + \lambda_{3}^{n} (g_{\tau(n)}(\xi^{n})) \right] \\
- \sum_{n=0}^{t(n) - 1} A_{\tau(n), \tau(n) - n}(\xi^{n}) x^{n_{n-1}} \right] : \\
x^{n} \in X_{\tau(n)}, B_{\tau(n)}(\xi^{n}) x^{n} \geq d_{\tau(n)}(\xi^{n}), n \in \mathcal{N},
\]

where \(\lambda_{3} \in A_{3} = \{\{\lambda_{3}^{n}\}_{n \in \mathcal{N}} : \lambda_{3}^{n} \geq 0, n \in \mathcal{N}\}\). Since the minimization problem in (38) contains only node constraints for the decision tree, we rearrange
its objective function with respect to the decision nodes and obtain

\[
D_3(\lambda_3) = \inf \left\{ \sum_{n \in \mathcal{N}} \pi_n(c_{t(n)}(\xi^n)) - \sum_{\ell \in \text{Tr}(n)} \pi_{\ell} \lambda_{\ell,3} A_{t(\ell),t(n)}(\xi^{\ell}) x^n \right. \\
\left. + \sum_{n \in \mathcal{N} \setminus \{1\}} \pi_n \lambda_{3}^{n} g_{t(n)}(\xi^n) : \right. \\
x^n \in X_{t(n)}, B_{t(n)}(\xi^n) x^n \geq d_{t(n)}(\xi^n), n \in \mathcal{N} \right\},
\]

where \( \text{Tr}(1) := \mathcal{N} \setminus \{1\} \), and \( \text{Tr}(n) \) for \( n > 1 \) denotes the set of all nodes belonging to the subtree with root node \( n \), i.e., \( \text{Tr}(n) := \bigcup_{n_T \in \mathcal{N} \setminus \{1\}} \{ \text{path}(n_T) : n_T \in \text{path}(n_T) \} \setminus \text{path}(n_\text{root}) \). Now, we may interchange summation and minimization and arrive at the node decomposed formulation

\[
D_3(\lambda_3) = \sum_{n \in \mathcal{N}} D_{3,n}(\lambda_3) + \sum_{n \in \mathcal{N} \setminus \{1\}} \pi_n \lambda_{3}^{n} g_{t(n)}(\xi^n) 
\]  

(39)

of \( D_3 \), where the functions \( D_{3,n}, n \in \mathcal{N}, \) are defined on \( \lambda_3 \) and given by

\[
D_{3,n}(\lambda_3) := \inf \left\{ (\pi_n c_{t(n)}(\xi^n)) - \sum_{\ell \in \text{Tr}(n)} \pi_{\ell} \lambda_{\ell,3} A_{t(\ell),t(n)}(\xi^{\ell}) x^n : \right. \\
\left. x^n \in X_{t(n)}, B_{t(n)}(\xi^n) x^n \geq d_{t(n)}(\xi^n) \right\}.
\]

Hence, the representation (39) of \( D_3 \) provides a decomposition of the original mixed-integer program of dimension \( \sum_{n \in \mathcal{N}} m_{t(n)} \) into \( |\mathcal{N}| \) subproblems (40) of dimension \( m_{t(n)} \) for \( n \in \mathcal{N} \). Formulas for computing subgradients of \( D_3 \) may be derived similarly to the previous section. Computational experience of such nodal decomposition schemes for determining lower bounds of multi-stage stochastic integer programs is not available yet.

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References

44. Schultz, R.: Rates of convergence in stochastic programs with complete integer recourse. SIAM J. Optim. 6 (1996), 1138–1152