

A NEW DERIVATION OF THE FINITE N MASTER LOOP EQUATION FOR LATTICE YANG-MILLS

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ABSTRACT. We give a new derivation of the finite N master loop equation for lattice Yang-Mills theory with structure group $SO(N)$, $U(N)$ or $SU(N)$. The $SO(N)$ case was initially proved by Chatterjee in [Cha19a], and $SU(N)$ was analyzed in a follow-up work by Jafarov [Jaf16]. Our approach is based on the Langevin dynamic, an SDE on the manifold of configurations, and yields a simple proof via Itô's formula.

1. INTRODUCTION

Let G be the Lie group $SO(N)$, $U(N)$, or $SU(N)$. The goal of this paper is to derive the finite N master loop equation for the lattice Yang-Mills theory with gauge group G . This was first obtained in [Cha19a] for $G = SO(N)$ and then in [Jaf16] for $G = SU(N)$.

We recall the setup, closely following the notation in [Cha19a]. Let Λ be a finite subset of \mathbb{Z}^d . We recall that a lattice edge is positively oriented if the beginning point is smaller in lexicographic order than the ending point. Let E^+ (resp. E^-) be the set of positively (resp. negatively) oriented edges, and denote by E_Λ^+ , E_Λ^- the corresponding subsets of edges with both beginning and ending points in Λ . Define $E \stackrel{\text{def}}{=} E^+ \cup E^-$ and let $u(e)$ and $v(e)$ denote the starting point and ending point of an edge $e \in E$, respectively. A path ρ in the lattice \mathbb{Z}^d is defined to be a sequence of edges $e_1 e_2 \cdots e_n$ with $e_i \in E$ and $v(e_i) = u(e_{i+1})$ for $i = 1, 2, \dots, n-1$. The path ρ is called closed if $v(e_n) = u(e_1)$. A plaquette is a closed path of length four which traces out the boundary of a square; more precisely it is non-backtracking in the sense of [Cha19a, Sec. 2]. The set of plaquettes is denoted as \mathcal{P} and \mathcal{P}_Λ is the set of plaquettes whose vertices are all in Λ , and \mathcal{P}_Λ^+ is the subset of plaquettes $p = e_1 e_2 e_3 e_4$ such that the beginning point of e_1 is lexicographically the smallest among all the vertices in p and the ending point of e_1 is the second smallest.

The lattice Yang-Mills theory (or lattice gauge theory) on Λ for the gauge group G , with $\beta \in \mathbb{R}$ the inverse coupling constant, is the probability measure $\mu_{\Lambda, N, \beta}$ on the set of all collections $Q = (Q_e)_{e \in E_\Lambda^+}$ of G -matrices, defined as

$$d\mu_{\Lambda, N, \beta}(Q) := Z_{\Lambda, N, \beta}^{-1} \exp\left(N\beta \operatorname{Re} \sum_{p \in \mathcal{P}_\Lambda^+} \operatorname{Tr}(Q_p)\right) \prod_{e \in E_\Lambda^+} d\sigma_N(Q_e), \quad (1.1)$$

where $Z_{\Lambda, N, \beta}$ is the normalizing constant, $Q_p \stackrel{\text{def}}{=} Q_{e_1} Q_{e_2} Q_{e_3} Q_{e_4}$ for a plaquette $p = e_1 e_2 e_3 e_4$, and σ_N is the Haar measure on G . Note that for $p \in \mathcal{P}_\Lambda^+$ the edges e_3 and e_4 are negatively oriented, so throughout the paper we define $Q_e \stackrel{\text{def}}{=} Q_{e^{-1}}$ for $e \in E^-$, where e^{-1} denotes the edge with orientation reversed. Also, Re is the real part, which can be omitted when $G = SO(N)$. We do not intend to further discuss the background and motivation for the above model (1.1); instead we refer to the review paper [Cha19b].

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For a closed path $\rho = e_1 \cdots e_n$, ρ' is said to be cyclically equivalent to ρ if $\rho' = e_i e_{i+1} \cdots e_n e_1 e_2 \cdots e_{i-1}$ for some $2 \leq i \leq n$. Cyclical equivalence classes are referred to as cycles, and a cycle with no backtracking is called a loop and denoted by l . By [Cha19a, Lemma 2.1], for any cycle l there is a unique loop denoted as $[l]$ by successive backtrack erasures until there are no more backtracks. A loop sequence $s = (l_1, \dots, l_m)$ is a collection of loops; more precisely it is an equivalence class understood up to a insertion and deletion of a null cycle (think of ee^{-1} for instance). The length of a loop l is denoted by $|l|$. For a loop sequence s with minimal representation (l_1, \dots, l_n) , the length is defined as

$$|s| \stackrel{\text{def}}{=} \sum_{i=1}^n |l_i|. \quad (1.2)$$

We refer to [Cha19a, Sec. 2] for precise definitions of loop sequence, minimal representation, location, etc.

Given a loop $l = e_1 e_2 \cdots e_n$, the Wilson loop variable W_l is defined as

$$W_l = \text{Tr}(Q_{e_1} Q_{e_2} \cdots Q_{e_n}).$$

By cyclic invariance of the trace, this definition is independent of the particular representative chosen in the equivalence class l . Write \mathbf{E} for expectation with respect to (1.1). For any non-null loop sequence s with minimal representation (l_1, \dots, l_m) such that each l_i is contained in Λ , define

$$W_s \stackrel{\text{def}}{=} W_{l_1} W_{l_2} \cdots W_{l_m}, \quad \phi(s) \stackrel{\text{def}}{=} \mathbf{E} \frac{W_s}{N^m}. \quad (1.3)$$

The master loop equation is a recursion which expresses $\phi(s)$ in terms of a linear combination of $\phi(s')$, where s' is a loop sequence obtained by performing an operation on s . The operations are called splitting, twisting, merger, deformation, and expansion; each being further divided into a positive or negative type. This leads to a definition of sets $\mathbb{T}^\pm(s)$, $\mathbb{S}^\pm(s)$, $\mathbb{M}^\pm(s)$, $\mathbb{D}^\pm(s)$, $\mathbb{E}^\pm(s)$ of all loop sequences obtained by performing one of these operations on s . The precise definition is given in (O1)-(O5) in Section 4, but we also recommend the graphs and further discussion in [Cha19a, Sec. 2.2] for additional intuition.

Moreover, to state the theorem, we define as in [Jaf16]

$$\ell(s) \stackrel{\text{def}}{=} \sum_{e \in E^+} t(e)^2, \quad t(e) \stackrel{\text{def}}{=} \sum_{i=1}^m t_i(e) \quad (e \in E^+) \quad (1.4)$$

and $t_i(e) \stackrel{\text{def}}{=} \text{the number of occurrences of } e \text{ minus the number of occurrences of } e^{-1} \text{ in the loop } l_i$. The following is the finite N master loop equation for the model (1.1).

Theorem 1.1. *Let s be as above, and suppose that all vertices that are at distance ≤ 1 from any l_i belong to Λ . Then for $G = SO(N)$ ([Cha19a, Theorem 3.6])*

$$\begin{aligned} (N-1)|s|\phi(s) &= N\beta \sum_{s' \in \mathbb{D}^-(s)} \phi(s') - N\beta \sum_{s' \in \mathbb{D}^+(s)} \phi(s') + N \sum_{s' \in \mathbb{S}^-(s)} \phi(s') - N \sum_{s' \in \mathbb{S}^+(s)} \phi(s') \\ &+ \sum_{s' \in \mathbb{T}^-(s)} \phi(s') - \sum_{s' \in \mathbb{T}^+(s)} \phi(s') + \frac{1}{N} \sum_{s' \in \mathbb{M}^-(s)} \phi(s') - \frac{1}{N} \sum_{s' \in \mathbb{M}^+(s)} \phi(s'). \end{aligned} \quad (1.5)$$

For $G = SU(N)$ ([Jaf16, Theorem 4.2])

$$\begin{aligned} \left(N|s| - \frac{\ell(s)}{N}\right)\phi(s) &= \frac{N\beta}{2} \sum_{s' \in \mathbb{D}^-(s)} \phi(s') - \frac{N\beta}{2} \sum_{s' \in \mathbb{D}^+(s)} \phi(s') + N \sum_{s' \in \mathbb{S}^-(s)} \phi(s') - N \sum_{s' \in \mathbb{S}^+(s)} \phi(s') \\ &+ \frac{N\beta}{2} \sum_{s' \in \mathbb{E}^-(s)} \phi(s') - \frac{N\beta}{2} \sum_{s' \in \mathbb{E}^+(s)} \phi(s') + \frac{1}{N} \sum_{s' \in \mathbb{M}_U^-(s)} \phi(s') - \frac{1}{N} \sum_{s' \in \mathbb{M}_U^+(s)} \phi(s'). \end{aligned} \quad (1.6)$$

¹ For $G = U(N)$,

$$\begin{aligned}
 N|s|\phi(s) = & \frac{N\beta}{2} \sum_{s' \in \mathbb{D}^-(s)} \phi(s') - \frac{N\beta}{2} \sum_{s' \in \mathbb{D}^+(s)} \phi(s') + N \sum_{s' \in \mathbb{S}^-(s)} \phi(s') - N \sum_{s' \in \mathbb{S}^+(s)} \phi(s') \\
 & + \frac{1}{N} \sum_{s' \in \mathbb{M}_V^-(s)} \phi(s') - \frac{1}{N} \sum_{s' \in \mathbb{M}_V^+(s)} \phi(s').
 \end{aligned} \tag{1.7}$$

The proof of this theorem in [Cha19a, Jaf16] is based on Stein’s exchangeable pairs and integration by parts (see Sec. 5 - Sec. 8 therein). Here we reprove it using a simple Langevin dynamic and Itô calculus. To illustrate the basic idea, recall that for standard Gaussian X , one has integration by parts (Stein’s lemma) $\mathbf{E}[Xf(X)] = \mathbf{E}[f'(X)]$. This identity can be derived in a simple way from Itô’s formula applied to a Langevin dynamic as follows. Consider the Ornstein-Uhlenbeck process $dX_t = -\frac{1}{2}X_t dt + dB_t$ for a Brownian motion B . Let F be an anti-derivative of f , i.e. $F' = f$. Then $dF(X_t) = F'(X_t)(-\frac{1}{2}X_t dt + dB_t) + \frac{1}{2}F''(X_t)dt$. In stationarity, taking expectation yields Stein’s lemma. Our proof of Theorem 1.1 is based on a generalization of this idea.

Chatterjee and Jafarov [Cha19a, Jaf16] have used the master equations to study various properties of Wilson loops in the large N limit under a smallness constraint in β . Further investigations were initiated by Basu and Ganguly [BG18], where more structure is deduced on the limit in the two dimensional lattice setting. We also mention that the Langevin dynamic can also be used to prove large N limit, mass gap property and uniqueness of the lattice Yang–Mills model on the whole lattice, see [SZZ22].

Master loop equations for lattice gauge theories were originally stated in Makeenko-Migdal’s paper [MM79], thus are also often called Makeenko-Migdal equations. The first rigorous version was established for two-dimensional Yang-Mills model in continuum in [Lév17], and alternative proofs and extensions were given in [Dah16, Dri19, DGHK17, DHK17] on plane or surfaces. These equations belong to a general class of equations arising in quantum field theory and random matrix theory, known as Integration by Parts or Schwinger–Dyson equations. See for instance [CGMS09, GN15] who derived Schwinger–Dyson equations for orthogonal and unitary multimatrix models, which are to some extent related with the lattice Yang–Mills model.

We remark that the lattice cutoff is a key simplification which allows us to prove Theorem 1.1 in any dimension; it would be much more challenging to prove similar results in the continuum (besides the aforementioned two-dimensional results). It would be natural to conjecture that the Langevin dynamic in Section 3 - which is the main ingredient in our proof - has a scaling limit given by the ones recently constructed in [CCHS22a, CCHS22b] on the two and three dimensional torus. Since our proof is relatively more elementary, we hope that it would make it easier to study related questions or other gauge theory models. We also refer to [Cha19a, Sec. 18] for a more comprehensive list of open questions.

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¹Compared to [Jaf16, Theorem 4.2] we use a different notation \mathbb{M}_V^\pm for the sets of merger terms to distinguish it from $SO(N)$ case.

2. PRELIMINARIES

This section collects some standard facts about Brownian motions on a Lie group G or its Lie algebra \mathfrak{g} , mostly from [Lév17, Sec. 1].

We write the Lie algebras of $SO(N)$, $U(N)$, $SU(N)$ as $\mathfrak{so}(N)$, $\mathfrak{u}(N)$, $\mathfrak{su}(N)$ respectively. Every matrix Q in one of these Lie groups satisfies $QQ^* = I_N$, and every matrix X in one of these Lie algebras satisfies $X + X^* = 0$. Here I_N denotes the identity matrix, and for any matrix M we write M^* for the conjugate transpose of M . Let $M_N(\mathbb{R})$ and $M_N(\mathbb{C})$ be the space of real and complex $N \times N$ matrices.

We endow $M_N(\mathbb{C})$ with the Hilbert-Schmidt inner product

$$\langle X, Y \rangle = \operatorname{Re} \operatorname{Tr}(XY^*) \quad \forall X, Y \in M_N(\mathbb{C}). \quad (2.1)$$

We restrict this inner product to our Lie algebra \mathfrak{g} , which is then invariant under the adjoint action. In particular for $X, Y \in \mathfrak{so}(N), \mathfrak{u}(N)$ or $\mathfrak{su}(N)$ we have $\langle X, Y \rangle = -\operatorname{Tr}(XY)$.

It is well-known that this induces an inner product on the tangent space at every $Q \in G$ via the right multiplication on G . Indeed, given any $X \in \mathfrak{g}$, the curve $t \mapsto e^{tX}Q$ is well approximated near $t = 0$ by $Q + tXQ$ up to an error of order t^2 . Hence, for $X, Y \in \mathfrak{g}$, XQ and YQ are two tangent vectors on the tangent space at $Q \in G$, and their inner product is given by $\operatorname{Tr}((XQ)(YQ)^*) = \operatorname{Tr}(XY^*)$.

Denote by \mathfrak{B} and B the Brownian motions on G and its Lie algebra \mathfrak{g} respectively. B is characterized by

$$\mathbf{E} \left[\langle B(s), X \rangle \langle B(t), Y \rangle \right] = \min(s, t) \langle X, Y \rangle \quad \forall X, Y \in \mathfrak{g}. \quad (2.2)$$

By [Lév17, Sec. 1.4], the Brownian motions \mathfrak{B} and B are related through the following SDE:

$$d\mathfrak{B} = dB \circ \mathfrak{B} = dB \mathfrak{B} + \frac{c_{\mathfrak{g}}}{2} \mathfrak{B} dt, \quad (2.3)$$

where \circ is the Stratonovich product, and $dB \mathfrak{B}$ is in the Itô sense. Here the constant $c_{\mathfrak{g}}$ is determined by $\sum_{\alpha} v_{\alpha}^2 = c_{\mathfrak{g}} I_N$ where $(v_{\alpha})_{\alpha=1}^{\dim_{\mathbb{R}} \mathfrak{g}}$ is an orthonormal basis of \mathfrak{g} . Note that $\dim_{\mathbb{R}} \mathfrak{g}$ indicates that for matrices with complex entries, dimension is counted with respect to \mathbb{R} . Moreover, by [Lév17, Lem. 1.2],

$$c_{\mathfrak{so}(N)} = -\frac{1}{2}(N-1), \quad c_{\mathfrak{u}(N)} = -N, \quad c_{\mathfrak{su}(N)} = -\frac{N^2-1}{N}. \quad (2.4)$$

Note that in [Lév17, Lem. 1.2], the scalar product differs from (2.1) by a scalar multiple depending on N and \mathfrak{g} , so we accounted for this in the expression for $c_{\mathfrak{g}}$ above.

Denote by δ the Kronecker function, i.e. $\delta_{ij} = 1$ if $i = j$ and 0 otherwise. For any matrix M , we write M^{ij} for its (i, j) th entry. The following holds by straightforward calculations:

$$d\langle B^{ij}, B^{k\ell} \rangle = \frac{1}{2}(\delta_{ik}\delta_{j\ell} - \delta_{i\ell}\delta_{jk})dt, \quad \mathfrak{g} = \mathfrak{so}(N); \quad (2.5a)$$

$$d\langle B^{ij}, B^{k\ell} \rangle = -\delta_{i\ell}\delta_{jk} dt, \quad \mathfrak{g} = \mathfrak{u}(N); \quad (2.5b)$$

$$d\langle B^{ij}, B^{k\ell} \rangle = (-\delta_{i\ell}\delta_{jk} + \frac{1}{N}\delta_{ij}\delta_{k\ell}) dt, \quad \mathfrak{g} = \mathfrak{su}(N). \quad (2.5c)$$

Remark 2.1. *The relation (2.5c) can be deduced from (2.5b) as follows. Given a $\mathfrak{u}(N)$ Brownian motion $t \mapsto B_t$, we may define an $\mathfrak{su}(N)$ Brownian motion $t \mapsto \hat{B}_t$ by letting $\hat{B}_t \stackrel{\text{def}}{=} B_t - \frac{1}{N} \operatorname{Tr}(B_t) I_N$. Since any $X \in \mathfrak{su}(N) \subset \mathfrak{u}(N)$ is traceless, the identity (2.2) satisfied by B implies the same identity for \hat{B} in light of the equality $\langle \hat{B}_t, X \rangle = \langle B_t, X \rangle$, which follows from $\langle I_N, X \rangle = \operatorname{Tr}(X) = 0$. To obtain (2.5c) from (2.5b), note that off diagonal entries of \hat{B} are identical to those of B , so (2.5c) holds*

if $i \neq j$ and $k \neq \ell$. Similarly, since on-diagonal entries of B are independent from off-diagonal entries, it holds that $\langle \hat{B}^{ii}, \hat{B}^{k\ell} \rangle = 0$ for $k \neq \ell$, again consistent with (2.5c). For $k = \ell$, note that

$$d\langle \hat{B}^{ii}, \hat{B}^{\ell\ell} \rangle = d\left\langle B^{ii} - \frac{1}{N} \sum_{j=1}^N B^{jj}, B^{\ell\ell} - \frac{1}{N} \sum_{j=1}^N B^{jj} \right\rangle = \left(-\delta_{i\ell} + \frac{1}{N} + \frac{1}{N} - N \cdot \frac{1}{N^2}\right) dt = \left(-\delta_{i\ell} + \frac{1}{N}\right) dt,$$

where we used the independence of diagonal entries of B .

Remark 2.2. Note that the choice of this inner product (2.1) may differ among the literature by a constant multiple. (2.4) will then differ by (the inverse of) the same constant. The r.h.s. of (2.5a) – (2.5c) should then also be multiplied by the suitable constants.

3. YANG MILLS SDE

Recall that for any function $f \in C^\infty(G)$, the right-invariant derivative is given by $\frac{d}{dt}|_{t=0} f(e^{tX}Q)$. Define the configuration space as the Lie group product $\mathcal{Q} = G^{E_\Lambda^+}$, consisting of all maps $Q : e \in E_\Lambda^+ \mapsto Q_e \in G$. Let $\mathfrak{q} = \mathfrak{g}^{E_\Lambda^+}$ be the corresponding direct sum of \mathfrak{g} and note that \mathfrak{q} is the Lie algebra of the Lie group \mathcal{Q} . For any matrix-valued functions A, B on E_Λ^+ , we denote by AB the pointwise product. Given $X \in \mathfrak{q}$, the exponential map $t \mapsto \exp(tX)$ is defined by $\exp(tX)_e \stackrel{\text{def}}{=} e^{tX_e}$.

As above, the tangent space at $Q \in \mathcal{Q}$ consists of the products XQ with $X \in \mathfrak{q}$, and given two such elements XQ and YQ , their inner product is defined by

$$\langle XQ, YQ \rangle = \sum_{e \in E_\Lambda^+} \text{Tr}(X_e Y_e^*).$$

Given any function $f \in C^\infty(\mathcal{Q})$, the right-invariant derivative is given by $\frac{d}{dt}|_{t=0} f(\exp(tX)Q)$. For each $Q \in \mathcal{Q}$, the differential $\nabla f(Q)$ is an element of the tangent space at Q which satisfies for each $X \in \mathfrak{q}$

$$\langle \nabla f(Q), XQ \rangle = \frac{d}{dt} \Big|_{t=0} f(\exp(tX)Q). \quad (3.1)$$

Denote by $\mathfrak{B} = (\mathfrak{B}_e)_{e \in E_\Lambda^+}$ and $B = (B_e)_{e \in E_\Lambda^+}$ the Brownian motions on \mathcal{Q} and \mathfrak{q} respectively. These are independent on distinct edges, and for each edge e , \mathfrak{B}_e and B_e are related through (2.3). Let $\mathcal{S}(Q) \stackrel{\text{def}}{=} N\beta \text{Re} \sum_{p \in \mathcal{P}_\Lambda^+} \text{Tr}(Q_p)$. We consider the Langevin dynamic for the measure (1.1), which is the following SDE on \mathcal{Q}

$$dQ = \frac{1}{2} \nabla \mathcal{S}(Q) dt + d\mathfrak{B}. \quad (3.2)$$

We now derive an explicit expression for $\nabla \mathcal{S}$. For a plaquette $p = e_1 e_2 e_3 e_4 \in \mathcal{P}$, we write $p \succ e_1$ to indicate that p is a plaquette that starts from edge e_1 . Note that for each edge e , there are $2(d-1)$ plaquettes in \mathcal{P} such that $p \succ e$. For any Lie algebra \mathfrak{g} embedded into $M_N(\mathbb{C})$, it forms a closed subspace of M_N , and therefore M_N has an orthogonal decomposition $M_N = \mathfrak{g} \oplus \mathfrak{g}^\perp$. Given $M \in M_N$, we denote by $\mathbf{p}M \in \mathfrak{g}$ the orthogonal projection onto \mathfrak{g} .

Lemma 3.1. *It holds that*

$$\nabla \mathcal{S}(Q)_e = N\beta \sum_{p \in \mathcal{P}_\Lambda, p \succ e} \mathbf{p}Q_p^* \cdot (Q_e^*)^{-1}, \quad (3.3)$$

where \cdot is matrix multiplication.

Proof. To prove the claim, fixing an edge $e \in E_\Lambda^+$, let $X \in \mathfrak{g}$ and with a slight abuse of notation we write $X \in \mathfrak{q}$ to be the function which is X at e and zero elsewhere. Note that for every $\tilde{p} \in \mathcal{P}_\Lambda^+$ that contains both the beginning point and the ending point of e , there is a unique way to obtain a plaquette $p \in \mathcal{P}_\Lambda$ such that $p \succ e$ by a cyclic permutation and possibly a reversal of the four edges.

One then has $\mathrm{Tr}(Q_{\bar{p}}) = \mathrm{Tr}(Q_p)$ (without reversal) or $\mathrm{Tr}(Q_{\bar{p}}) = \mathrm{Tr}(Q_p^{-1}) = \mathrm{Tr}(Q_p^*)$ (with reversal). We have

$$\left. \frac{d}{dt} \right|_{t=0} \mathrm{ReTr}(\exp(tX)Q_p) = \left. \frac{d}{dt} \right|_{t=0} \mathrm{ReTr}((\exp(tX)Q_p)^*) = \mathrm{ReTr}(XQ_p)$$

where we used $(XQ_p)^* = \overline{(XQ_p)^t}$. Therefore the derivative of \mathcal{S} at Q in the tangent direction XQ is equal to $N\beta$ times

$$\sum_{p \in \mathcal{P}_\Lambda, p \succ e} \mathrm{ReTr}(XQ_p) = \sum_{p \in \mathcal{P}_\Lambda, p \succ e} \langle X, Q_p^* \rangle = \sum_{p \in \mathcal{P}_\Lambda, p \succ e} \langle X, \mathbf{p}Q_p^* \rangle = \sum_{p \in \mathcal{P}_\Lambda, p \succ e} \mathrm{ReTr}(X^* \mathbf{p}Q_p^*). \quad (3.4)$$

Furthermore, note that

$$\langle \nabla \mathcal{S}(Q)_e, XQ_e \rangle = \mathrm{ReTr}(\nabla \mathcal{S}(Q)_e Q_e^* X^*) = \mathrm{ReTr}(X^* \nabla \mathcal{S}(Q)_e Q_e^*). \quad (3.5)$$

To ensure that (3.4) and (3.5) agree, we have (3.3). \square

For specific choice of Lie algebras, our SDE system reads

$$dQ_e = \frac{1}{2} \nabla \mathcal{S}(Q)_e dt + \frac{1}{2} c_{\mathfrak{g}} Q_e dt + dB_e Q_e, \quad (e \in E_\Lambda^+) \quad (3.6)$$

where $c_{\mathfrak{g}}$ is as in (2.4) and

$$\frac{1}{2} \nabla \mathcal{S}(Q)_e = \begin{cases} -\frac{1}{4} N\beta \sum_{p \in \mathcal{P}_\Lambda, p \succ e} (Q_p - Q_p^*) Q_e, & \mathfrak{g} \in \{\mathfrak{so}(N), \mathfrak{u}(N)\}, \\ -\frac{1}{4} N\beta \sum_{p \in \mathcal{P}_\Lambda, p \succ e} \left((Q_p - Q_p^*) - \frac{1}{N} \mathrm{Tr}(Q_p - Q_p^*) I_N \right) Q_e, & \mathfrak{g} \in \mathfrak{su}(N). \end{cases} \quad (3.7)$$

While our measure (1.1) and the dynamic are both defined on the configuration space $\mathcal{Q} = G_\Lambda^{E^+}$ (in particular the above SDE system is parametrized by $e \in E_\Lambda^+$), when we apply it to Wilson loops later, we will also need to differentiate $Q_{e^{-1}}$ for $e \in E_\Lambda^+$, which is just Q_e^* . So we give the conjugate transpose of (3.6):

$$dQ_e^* = \frac{1}{2} (\nabla \mathcal{S}(Q)_e)^* dt + \frac{1}{2} c_{\mathfrak{g}} Q_e^* + Q_e^* dB_e^*, \quad (e \in E_\Lambda^+). \quad (3.8)$$

This system is well-posed and has (1.1) as invariant measure, as we show in the next two lemmas.

Lemma 3.2. *For fixed $N \in \mathbb{N}$, $T > 0$ and any initial data $Q(0) = (Q_e(0))_{e \in E_\Lambda^+} \in \mathcal{Q}$, there exists a unique solution $Q = (Q_e)_{e \in E_\Lambda^+} \in C([0, T]; \mathcal{Q})$ to (3.6) a.s..*

Proof. For fixed N and Λ , we can write (3.6) as the system for the entries of the matrices Q_e , which can be viewed as a finite dimensional SDE with locally Lipschitz coefficients. We introduce a stopping time

$$\tau := \inf\{t \geq 0 : \|Q_e(t)\|_\infty > 2, \text{ for at least one } e \in E_\Lambda^+\} \wedge T,$$

and we obtain local solutions $Q = (Q_e)_{e \in E_\Lambda^+}$ with $Q_e \in C([0, \tau]; M_N)$ for $e \in E_\Lambda^+$, which satisfies (3.6) before τ , by fixed point argument (see e.g. [LR15, Chapter 3]).

Since for $e \in E_\Lambda^+$, $\nabla \mathcal{S}(Q)_e$ belongs to the tangent space of G at Q_e , and the generator associated with $dB_e Q_e + \frac{c_{\mathfrak{g}}}{2} B_e Q_e$ is given by $\frac{1}{2}$ of Laplace-Beltrami operator on G (see (2.3)), exactly the same argument as in [Lév17, Lemma 1.3] implies that $Q_e(t) \in G$, $\forall t \geq 0$, and $\tau = T$ a.s.. \square

Lemma 3.3. (1.1) is invariant under the SDE system (3.6).

Proof. By standard calculation (cf. [Hsu02, Chapter 3]) the generator L for our SDE with $F \in C^\infty(\mathcal{Q})$ is

$$LF = \frac{1}{2} \sum_{e \in E_\Lambda^+} \Delta_{Q_e} F + \sum_{e \in E_\Lambda^+} \frac{1}{2} \langle \nabla \mathcal{S}(Q)_e, \nabla_{Q_e} F \rangle. \quad (3.9)$$

Here Δ_{Q_e} and ∇_{Q_e} are the Laplace-Beltrami operator and the gradient (w.r.t. the variable Q_e) on G endowed with the metric given in Sec 2. Using integration by parts w.r.t. the Haar measure, we have for $F, G \in C^\infty(\mathcal{Q})$

$$\begin{aligned} \int (LF)G d\mu_{\Lambda, N, \beta} &= -\frac{1}{2} \sum_{e \in E_\Lambda^+} \int \langle \nabla_{Q_e} F, \nabla \mathcal{S}(Q)_e \rangle G d\mu_{\Lambda, N, \beta} + \frac{1}{2} \sum_{e \in E_\Lambda^+} \int \langle \nabla_{Q_e} F, \nabla \mathcal{S}(Q)_e \rangle G d\mu_{\Lambda, N, \beta} \\ &\quad - \frac{1}{2} \sum_{e \in E_\Lambda^+} \int \langle \nabla_{Q_e} F, \nabla_{Q_e} G \rangle d\mu_{\Lambda, N, \beta} \\ &= -\frac{1}{2} \sum_{e \in E_\Lambda^+} \int \langle \nabla_{Q_e} F, \nabla_{Q_e} G \rangle d\mu_{\Lambda, N, \beta} = \int (LG)F d\mu_{\Lambda, N, \beta}, \end{aligned} \quad (3.10)$$

where we exchange F and G in the last step. Hence, L is symmetric w.r.t. $\mu_{\Lambda, N, \beta}$ and the result follows by $L1 \equiv 0$. \square

Using Lemma 3.2, we can choose $\mu_{\Lambda, N, \beta}$ as an initial distribution and obtain a stationary solution $Q \in C([0, T]; \mathcal{Q})$. We will fix this stationary solution in the following proof.

Remark 3.4. *Dynamics very similar with ours (3.6) seem to have appeared in physics and are used for Monte Carlo simulations for lattice gauge theory, see e.g. [BKK⁺85, GL83].*

4. PROOF OF THE MAIN THEOREM

The proof of Theorem 1.1 is based on the Itô formula applied to W_s according to the dynamics (3.6). The master theorem for $\phi(s)$ is obtained by normalizing and taking expected value. The nonlinear drift of the SDE yields the deformation and expansion terms, while the Itô correction leads to the splitting, twisting and merger terms, in addition to a multiple of W_s , which combined with $c_{\mathfrak{g}}$ part gives the left-hand side of the master loop equation (1.5).

We recall from [Cha19a, Sec. 2] and [Jaf16, Sec. 2] the following notation and definitions.

- (O1) $\times_{x,y}^1 l$ and $\times_{x,y}^2 l$ denote the pair of loops obtained by positive splitting of l at x and y if l contains the same edge e at x and y , or negative splitting if l contains e at location x and e^{-1} at location y . For $l = aebec$ (where a, b, c are paths and e is an edge), $\times_{x,y}^1 l \stackrel{\text{def}}{=} [aec]$ and $\times_{x,y}^2 l \stackrel{\text{def}}{=} [be]$. For $l = aebe^{-1}c$, $\times_{x,y}^1 l \stackrel{\text{def}}{=} [ac]$ and $\times_{x,y}^2 l \stackrel{\text{def}}{=} [b]$. We say that a loop sequence s' is obtained from splitting s provided that exactly two components of s' arise from splitting a single loop in s . The sets $\mathbb{S}^+(s)$ and $\mathbb{S}^-(s)$ consist of all loop sequences obtained from positive or negative splitting of s , respectively.
- (O2) $\alpha_{x,y} l$ denotes the negative twisting if l contains an edge e at both x and y , or positive twisting if l contains an edge e at location x and e^{-1} at location y . For $l = aebec$, $\alpha_{x,y} l \stackrel{\text{def}}{=} [ab^{-1}c]$. For $l = aebe^{-1}c$, $\alpha_{x,y} l \stackrel{\text{def}}{=} [aeb^{-1}e^{-1}c]$. We say that a loop sequence s' is obtained from twisting s provided that exactly one component of s' arises from twisting one loop in s . The sets $\mathbb{T}^+(s)$ and $\mathbb{T}^-(s)$ consist of all loop sequences obtained from positive or negative twisting of s , respectively.

- (O3) $l \oplus_{x,y} l'$ and $l \ominus_{x,y} l'$ are positive and negative mergers of l and l' at locations x, y . For $l = aeb$ and $l' = ced$ (where a, b, c, d are paths and e is an edge), $l \oplus_{x,y} l' = [aedceb]$, $l \ominus_{x,y} l' = [ac^{-1}d^{-1}b]$. For $l = aeb$ and $l' = ce^{-1}d$, $l \oplus_{x,y} l' = [aec^{-1}d^{-1}eb]$, $l \ominus_{x,y} l' = [adcb]$ (here x and y are the unique location in l and l' , respectively, where e or e^{-1} occurs and e is the edge occurring at location x in l). We say that a loop sequence s' is obtained from merging s provided that exactly one component of s' arises from merging two loops in s . The sets $\mathbb{M}^+(s)$ and $\mathbb{M}^-(s)$ denote all loop sequences obtained from either positive mergers or negative mergers of s . Furthermore, we define two more sets $\mathbb{M}_U^+(s) \subset \mathbb{M}^+(s)$ and $\mathbb{M}_U^-(s) \subset \mathbb{M}^-(s)$; the first consists of positive mergers resulting from an edge e appearing in both of the two merged loops; the second consists of negative mergers where an edge e occurs in one loop and e^{-1} in the other.
- (O4) $l \oplus_x p$ and $l \ominus_x p$ are deformations obtained by merging l and p at locations x and y (here y is the unique location in p where e or e^{-1} occurs and e is the edge occurring at location x in l). We say that a loop sequence s' is obtained from deformations of s provided that exactly one component of s' arises from deformation of one loop in s . The sets $\mathbb{D}^+(s)$ and $\mathbb{D}^-(s)$ consist of all loop sequences obtained from positive or negative deformations of s , respectively.
- (O5) A positive expansion of l at location x by a plaquette p passing through e^{-1} replaces l with the pair of loops (l, p) . A negative expansion of l at location x by a plaquette p passing through e replaces l with the pair of loops (l, p) . The sets $\mathbb{E}^+(s)$ and $\mathbb{E}^-(s)$ consist of all loop sequences obtained from positive or negative expansions of s , respectively.

In preparation for an application of the Itô formula, we recall a convenient matrix analogue of Itô differentials. To treat each of the three groups G in a unified way, we introduce parameters λ , ν , and μ as follows. For $G \in \{SO(N), U(N), SU(N)\}$, we rewrite (2.5) as

$$dB^{ij}dB^{kl} \stackrel{\text{def}}{=} d\langle B^{ij}, B^{kl} \rangle = (\lambda\delta_{i\ell}\delta_{jk} + \nu\delta_{ij}\delta_{k\ell} + \mu\delta_{ik}\delta_{j\ell}) dt, \quad (4.1)$$

where λ, μ, ν depend on G and the choice of the inner product. Since we will apply Itô formula to products of matrices, we will use a matrix variant of the standard $dBdB$ notation for formulating Itô's rule. Given a matrix M , we use the shorthand $dB M$ or $dB M dB$, which should always be understood by writing it in components as a matrix product and in the latter case applying (2.5). This leads us to two useful identities which will be crucial in the proof of Theorem 1.1.

Given adapted matrix-valued processes M, N , we have the following two identities

$$dB M dB = (\lambda \text{Tr} M + \nu M + \mu M^t) dt. \quad (4.2)$$

$$\text{Tr}(dB M) \text{Tr}(dB N) = \left(\lambda \text{Tr}(MN) + \nu \text{Tr}(M) \text{Tr}(N) + \mu \text{Tr}(MN^t) \right) dt. \quad (4.3)$$

The first follows from (4.1) by fixing components i, ℓ and writing

$$\begin{aligned} (dB M dB)^{i\ell} &= \sum_{j,k} dB^{ij} M^{jk} dB^{k\ell} = \sum_{j,k} M^{jk} (\lambda\delta_{i\ell}\delta_{jk} + \nu\delta_{ij}\delta_{k\ell} + \mu\delta_{ik}\delta_{j\ell}) dt, \\ &= (\lambda\delta_{i\ell} \text{Tr} M + \nu M^{i\ell} + \mu M^{\ell i}) dt. \end{aligned}$$

The second follows in a similar way, since

$$\begin{aligned} \text{Tr}(dB M) \text{Tr}(dB N) &= \sum_{i,j,k,\ell} dB^{ij} M^{ji} dB^{k\ell} N^{\ell k} = \sum_{i,j,k,\ell} M^{ji} N^{\ell k} (\lambda\delta_{i\ell}\delta_{jk} + \nu\delta_{ij}\delta_{k\ell} + \mu\delta_{ik}\delta_{j\ell}) dt \\ &= \left(\lambda \text{Tr}(MN) + \nu \text{Tr}(M) \text{Tr}(N) + \mu \text{Tr}(MN^t) \right) dt. \end{aligned}$$

We remark that the identities (4.2)-(4.3) are not new, and are sometimes called the “magic formulas”, see e.g. [DL22, Lemma 7.1] for more background and literature, or [Sen08, Lemma 4.1].

We now turn to the proof of the main theorem.

Proof of Theorem 1.1. Let s be a string with minimal representation (l_1, \dots, l_m) , then we may view the quantity $\phi(s)$ as the mean of the stationary Itô process $N^{-m} \prod_{i=1}^m W_{l_i}$. For each constituent loop $l \in \{l_1, \dots, l_m\}$, there exist edges (depending on l) labelled $e_1, \dots, e_n \in E$ such that $W_l = \text{Tr}(Q_{e_1} \cdots Q_{e_n})$. Applying Itô's product rule with (3.6) and (3.8) yields

$$dW_l = \left(\frac{c_{\mathfrak{g}}}{2} |l| W_l + \mathcal{D}_l + \mathcal{I}_l \right) dt + dM_l, \quad (4.4)$$

where we define

$$\mathcal{D}_l \stackrel{\text{def}}{=} \sum_{x=1}^n \frac{1}{2} \text{Tr} \left(\prod_{i=1}^{x-1} Q_{e_i} \left[\nabla \mathcal{S}(Q)_{e_x} \mathbf{1}_{e_x \in E^+} + (\nabla \mathcal{S}(Q)_{e_x^{-1}})^* \mathbf{1}_{e_x^{-1} \in E^+} \right] \prod_{i=x+1}^n Q_{e_i} \right), \quad (4.5)$$

$$dM_l \stackrel{\text{def}}{=} \sum_{x=1}^n \text{Tr} \left(\prod_{i=1}^{x-1} Q_{e_i} \left[dB_{e_x} Q_{e_x} \mathbf{1}_{e_x \in E^+} + Q_{e_x} dB_{e_x^{-1}}^* \mathbf{1}_{e_x^{-1} \in E^+} \right] \prod_{i=x+1}^n Q_{e_i} \right),$$

$$\mathcal{I}_l dt \stackrel{\text{def}}{=} \sum_{x < y} \text{Tr} (Q_a dQ_{e_x} Q_b dQ_{e_y} Q_c), \quad (4.6)$$

and recall that $e_x \in E^- \Leftrightarrow e_x^{-1} \in E^+$ (with the usual convention used above and below that an empty product of matrices is I_N). In the definition of \mathcal{I}_l we use the shorthand notation

$$Q_a \stackrel{\text{def}}{=} \prod_{i=1}^{x-1} Q_{e_i}, \quad Q_b \stackrel{\text{def}}{=} \prod_{i=x+1}^{y-1} Q_{e_i}, \quad Q_c \stackrel{\text{def}}{=} \prod_{i=y+1}^n Q_{e_i}$$

and omit the dependence of these quantities on x and y .

Below, in Step 1 we write \mathcal{I}_l in terms of the splitting and twisting operations, leading to the identity (4.9). At this stage we have the dynamic for each fixed loop in s and now want to analyze W_s , which is itself a product of Itô processes, so we apply Itô's product rule again and obtain

$$\begin{aligned} dW_s &= d\left(\prod_{i=1}^m W_{l_i}\right) = \sum_{i=1}^m dW_{l_i} \prod_{j \neq i} W_{l_j} + \mathcal{I}_s dt \\ &= \frac{c_{\mathfrak{g}}}{2} |s| W_s dt + \sum_{i=1}^m (\mathcal{I}_{l_i} + \mathcal{D}_{l_i}) \prod_{j \neq i} W_{l_j} dt + \mathcal{I}_s dt + dM_s, \end{aligned} \quad (4.7)$$

where M_s is a martingale and \mathcal{I}_s denotes the Itô correction defined by

$$\mathcal{I}_s dt \stackrel{\text{def}}{=} \sum_{i < j} dW_{l_i} dW_{l_j} \prod_{k \neq i, j} W_{l_k},$$

which is calculated in Step 2 in terms of the merger operation, leading to the identity (4.12). In Step 3, we express \mathcal{D}_l in terms of the deformation and expansion operations using the expression (3.7) for the drift $\nabla \mathcal{S}$, leading to the identity (4.14). Next, we normalize W_s by dividing (4.7) by N^m , taking expectation on both sides, and using stationarity to obtain

$$-\frac{c_{\mathfrak{g}}}{2} |s| \phi(s) = \frac{1}{N^m} \mathbf{E} \left[\sum_{i=1}^m (\mathcal{I}_{l_i} + \mathcal{D}_{l_i}) \prod_{j \neq i} W_{l_j} + \mathcal{I}_s \right]. \quad (4.8)$$

In the final step, we consider each particular Lie group, specifying $c_{\mathfrak{g}}$ according to (2.4), the parameters λ, μ, ν as in (2.5), and use the output of Steps 1-3 to show that the RHS of (4.8) can be closed in terms of ϕ , leading to the master equations (1.5)-(1.7).

STEP 1. In this step we analyze an individual loop W_l and argue that

$$\mathcal{I}_l = -\frac{\nu}{2} (|l| - \ell(l)) W_l - \frac{\lambda}{2} \sum_{s' \in \mathbb{S}^-(l)} W_{s'} + \frac{\lambda}{2} \sum_{s' \in \mathbb{S}^+(l)} W_{s'} + \frac{\mu}{2} \sum_{l' \in \mathbb{T}^-(l)} W_{l'} - \frac{\mu}{2} \sum_{l' \in \mathbb{T}^+(l)} W_{l'} \quad (4.9)$$

where $\ell(l)$ is as in (1.4) (with $m = 1$ there). To prove the claim, we apply (3.6)+(3.8) to (4.6). Since $(B_e)_{e \in E_\Lambda^+}$ are independent, the contribution to (4.6) is restricted to x, y with the property that $e_y = e_x$ or $e_y = e_x^{-1}$. Since the dynamics in (4.4) also depends on the orientation of the edge, we sub-divide each case into two further cases $e_x \in E^+$ and $e_x^{-1} \in E^+$ to obtain the identity

$$\mathcal{I}_l dt = \sum_{x < y} \left(J_l^{(1)} \mathbf{1}_{e_y=e_x \in E^+} + J_l^{(2)} \mathbf{1}_{e_y^{-1}=e_x^{-1} \in E^+} + J_l^{(3)} \mathbf{1}_{e_y^{-1}=e_x \in E^+} + J_l^{(4)} \mathbf{1}_{e_y=e_x^{-1} \in E^+} \right), \quad (4.10)$$

where $J_l^{(i)}$, $i = 1, 2, 3, 4$, are defined by ²

$$\begin{aligned} J_l^{(1)} dt &\stackrel{\text{def}}{=} \text{Tr} \left(Q_a (dB_{e_x} Q_{e_x}) Q_b (dB_{e_y} Q_{e_y}) Q_c \right), \\ J_l^{(2)} dt &\stackrel{\text{def}}{=} \text{Tr} \left(Q_a (Q_{e_x} dB_{e_x^{-1}}^*) Q_b (Q_{e_y} dB_{e_y^{-1}}^*) Q_c \right), \\ J_l^{(3)} dt &\stackrel{\text{def}}{=} \text{Tr} \left(Q_a (dB_{e_x} Q_{e_x}) Q_b (Q_{e_y} dB_{e_y^{-1}}^*) Q_c \right), \\ J_l^{(4)} dt &\stackrel{\text{def}}{=} \text{Tr} \left(Q_a (Q_{e_x} dB_{e_x^{-1}}^*) Q_b (dB_{e_y} Q_{e_y}) Q_c \right). \end{aligned} \quad (4.11)$$

Recall that $B_e^* = -B_e$. To analyze each term, we will apply (4.2) with a suitable choice of M and B , while taking into account the relation between e_x and e_y imposed by the indicator function according to (4.10). For $J_l^{(1)}$ and $J_l^{(2)}$, the role of M is played by $Q_{e_x} Q_b$ and $Q_b Q_{e_y}$ respectively, while the role of B is played by B_{e_x} and $B_{e_x^{-1}}^* = -B_{e_x^{-1}}$ leading to

$$J_l^{(1)} = J_l^{(2)} = \lambda \text{Tr}(Q_{e_x} Q_b) \text{Tr}(Q_a Q_{e_x} Q_c) + \nu \text{Tr}(Q_a Q_{e_x} Q_b Q_{e_x} Q_c) + \mu \text{Tr}(Q_a Q_{b^{-1}} Q_c),$$

where we note that $\mu \neq 0$ only if $G = SO(N)$ in which case $Q^* = Q^t$. In a similar way, we obtain

$$J_l^{(3)} = J_l^{(4)} = -\lambda \text{Tr}(Q_b) \text{Tr}(Q_a Q_c) - \nu \text{Tr}(Q_a Q_{e_x} Q_b Q_{e_x^{-1}} Q_c) - \mu \text{Tr}(Q_a Q_{e_x} Q_{b^{-1}} Q_{e_x^{-1}} Q_c),$$

where we used cyclic invariance of the trace and $Q_{e_x} Q_{e_y} = Q_{e_x} Q_{e_x}^* = I_N$.

From the definition of splitting terms in (O1), the terms above with a coefficient λ contribute the splitting terms in (4.9). Indeed, recall that after choosing a closed path in l and writing the loop as $l = e_1 \cdots e_n$, the sets $\mathbb{S}^+(l)$ and $\mathbb{S}^-(l)$ consist of all loop sequences $s' = (\times_{x,y}^1 l, \times_{x,y}^2 l)$ where the locations $x \neq y$ have the property that $e_y = e_x$ or $e_y = e_x^{-1}$ respectively, according to whether the splitting is positive or negative. Note that in the calculation above $x < y$, however if we let $s'' := (\times_{y,x}^1 l, \times_{y,x}^2 l)$, then $W_{s'} = W_{s''}$. Here, we recall from the discussion in [Cha19a, Sec 2.2] (below definitions of $\mathbb{S}^-, \mathbb{S}^+$ therein) that if a loop l has a splitting at x and y , then it also has a splitting at y and x , and they are reverse of each other, and should be counted as two distinct splittings of l . This gives the coefficient $\lambda/2$ before the splitting terms. The same applies to twisting and we will keep this in mind below.

The sets $\mathbb{T}^+(l)$ and $\mathbb{T}^-(l)$ consist of all loops of the form $\alpha_{x,y} l$ where the locations $x \neq y$ have the property that $e_y = e_x$ or $e_y = e_x^{-1}$ respectively, according to whether the twisting is positive or negative. Hence, from the definition of the twisting terms (O2), the terms with a coefficient μ lead to the final two terms in (4.9).

Finally, we turn to the terms with coefficient ν , and begin by introducing the following notation: for any edge $e \in E^+$ we let $A(e)$ be the set of locations in l where e occurs and $B(e)$ be the set of locations in l where e^{-1} occurs. The terms with coefficient ν are given by

$$\nu W_l \sum_{x < y} \left(\mathbf{1}_{e_x=e_y} - \mathbf{1}_{e_x=e_y^{-1}} \right)$$

²When $e \in E^-$, (3.8) yields $dQ_e = dQ_{e^{-1}}^* = (\cdots) + Q_{e^{-1}}^* dB_{e^{-1}}^* = (\cdots) + Q_e dB_{e^{-1}}^*$ where (\cdots) is the drift.

$$\begin{aligned}
 &= \nu W_l \sum_{x < y} \sum_{e \in E^+} \left(\mathbf{1}_{e_x=e} \mathbf{1}_{e_y=e} + \mathbf{1}_{e_x=e^{-1}} \mathbf{1}_{e_y=e^{-1}} - \mathbf{1}_{e_x=e} \mathbf{1}_{e_y=e^{-1}} - \mathbf{1}_{e_x=e^{-1}} \mathbf{1}_{e_y=e} \right) \\
 &= \nu W_l \sum_{e \in E^+} \left(\frac{1}{2} |A(e)| (|A(e)| - 1) + \frac{1}{2} |B(e)| (|B(e)| - 1) - |A(e)| |B(e)| \right) \\
 &= -\frac{\nu}{2} W_l \sum_{e \in E^+} \left(|A(e)| + |B(e)| - (|A(e)| - |B(e)|)^2 \right) = -\frac{\nu}{2} W_l (|l| - \ell(l)),
 \end{aligned}$$

with $\ell(l) = \sum_{e \in E^+} (|A(e)| - |B(e)|)^2$, which is precisely the first term in (4.9), completing the proof of the claim.

STEP 2. In this step we consider the m constituent Wilson loops W_{l_i} in W_s and argue that

$$\begin{aligned}
 \mathcal{I}_s = & \frac{\lambda}{2} \sum_{s' \in \mathbb{M}_U^+(s)} W_{s'} - \frac{\mu}{2} \sum_{s' \in \mathbb{M}^+(s) \setminus \mathbb{M}_U^+(s)} W_{s'} - \frac{\lambda}{2} \sum_{s' \in \mathbb{M}_U^-(s)} W_{s'} + \frac{\mu}{2} \sum_{s' \in \mathbb{M}^-(s) \setminus \mathbb{M}_U^-(s)} W_{s'} \\
 & + \nu \sum_{i < j} \sum_{e \in E^+} t_i(e) t_j(e) W_s.
 \end{aligned} \tag{4.12}$$

To analyze \mathcal{I}_s , we start by fixing two loops l_i and l_j with $i \neq j$, and analyze $dW_{l_i} dW_{l_j}$. First we choose a path to represent each loop and write them as

$$l_i = \prod_{k=1}^{|l_i|} e_k^i = a e_x^i b, \quad l_j = \prod_{k=1}^{|l_j|} e_k^j = c e_y^j d,$$

where we use a shorthand notation

$$a \stackrel{\text{def}}{=} \prod_{k=1}^{x-1} e_k^i, \quad b \stackrel{\text{def}}{=} \prod_{k=x+1}^{|l_i|} e_k^i, \quad c \stackrel{\text{def}}{=} \prod_{k=1}^{y-1} e_k^j, \quad d \stackrel{\text{def}}{=} \prod_{k=y+1}^{|l_j|} e_k^j.$$

Using again the independence of edges, taking into account (4.4) we obtain

$$dW_{l_i} dW_{l_j} = \sum_{x=1}^{|l_i|} \sum_{y=1}^{|l_j|} (\mathbf{1}_{e_y^j=e_x^i} + \mathbf{1}_{e_y^j=(e_x^i)^{-1}}) \text{Tr}(Q_a dQ_{e_x^i} Q_b) \text{Tr}(Q_c dQ_{e_y^j} Q_d).$$

To ease the notation below, we will drop the superscript from the edge and simply write $e_x = e_x^i$ and $e_y = e_y^j$. Apply (3.6)+(3.8) to the r.h.s. of the above equation. Using $B_e^* = -B_e$ and cyclic invariance of the trace, we may re-write the r.h.s. above as

$$\sum_{x=1}^{|l_i|} \sum_{y=1}^{|l_j|} \left(\mathbf{1}_{e_y=e_x \in E^+} J_s^{(1)} + \mathbf{1}_{e_y^{-1}=e_x^{-1} \in E^+} J_s^{(2)} + \mathbf{1}_{e_y^{-1}=e_x \in E^+} J_s^{(3)} + \mathbf{1}_{e_y=e_x^{-1} \in E^+} J_s^{(4)} \right) \tag{4.13}$$

where $J_s^{(i)}$ are defined by

$$\begin{aligned}
 J_s^{(1)} dt &\stackrel{\text{def}}{=} \text{Tr}(dB_{e_x} Q_{e_x} Q_b Q_a) \text{Tr}(dB_{e_y} Q_{e_y} Q_d Q_c), \\
 J_s^{(2)} dt &\stackrel{\text{def}}{=} \text{Tr}(dB_{e_x^{-1}} Q_b Q_a Q_{e_x}) \text{Tr}(dB_{e_y^{-1}} Q_d Q_c Q_{e_y}), \\
 J_s^{(3)} dt &\stackrel{\text{def}}{=} -\text{Tr}(dB_{e_x} Q_{e_x} Q_b Q_a) \text{Tr}(dB_{e_y^{-1}} Q_d Q_c Q_{e_y}), \\
 J_s^{(4)} dt &\stackrel{\text{def}}{=} -\text{Tr}(dB_{e_x^{-1}} Q_b Q_a Q_{e_x}) \text{Tr}(dB_{e_y} Q_{e_y} Q_d Q_c).
 \end{aligned}$$

We calculate these terms similarly as in Step 1 using cyclic invariance of the trace, $Q_{e_x} Q_{e_x}^* = I_N$, the fact that $\mu \neq 0$ only if $G = SO(N)$ in which case $Q_{e_x}^* = Q_{e_x}^t$, and taking into account the relation between e_x and e_y imposed by the indicator functions.

Applying (4.3) with $M = Q_{e_x} Q_b Q_a$ and $N = Q_{e_y} Q_d Q_c$ for $J_s^{(1)}$ and with $M = Q_b Q_a Q_{e_x}$ and $N = Q_d Q_c Q_{e_y}$ for $J_s^{(2)}$, under the assumption $e_x = e_y$ we obtain

$$J_s^{(1)} = J_s^{(2)} = \lambda \text{Tr}(Q_a Q_{e_x} Q_d Q_c Q_{e_x} Q_b) + \nu \text{Tr}(Q_{l_i}) \text{Tr}(Q_{l_j}) + \mu \text{Tr}(Q_a Q_{c^{-1}} Q_{d^{-1}} Q_b).$$

Similarly, applying (4.3) with $M = Q_{e_x} Q_b Q_a$, $N = Q_d Q_c Q_{e_y}$ and $M = Q_b Q_a Q_{e_x}$, $N = Q_{e_y} Q_d Q_c$, under the assumption $e_x = e_y^{-1}$ we obtain

$$J_s^{(3)} = J_s^{(4)} = -\lambda \text{Tr}(Q_a Q_d Q_c Q_b) - \nu \text{Tr}(Q_{l_i}) \text{Tr}(Q_{l_j}) - \mu \text{Tr}(Q_a Q_{e_x} Q_{c^{-1}} Q_{d^{-1}} Q_{e_x} Q_b).$$

We first note that for $s = (l_1, \dots, l_m)$, the set $\mathbb{M}^+(s)$ is the collection of s' which can be obtained from s by merging some l_i and l_j with $i \neq j$. If l_i and l_j can be merged at locations x and y respectively, then s', s'' defined by $s' = (l_1, \dots, l_{i-1}, l_i \oplus_{x,y} l_j, l_i, \dots, l_m)$ and $s'' = (l_1, \dots, l_{j-1}, l_i \oplus_{y,x} l_j, l_j, \dots, l_m)$ both belong to $\mathbb{M}^+(s)$. Since $i < j$ we need to account for both contributions, leading to the factors of $1/2$ in (4.12). An analogous characterization holds for $\mathbb{M}_U^+(s), \mathbb{M}^-(s)$, and $\mathbb{M}_U^-(s)$.

To find the contribution of the terms with coefficient μ to \mathcal{I}_s , we recall the definition of mergers terms in (O3) and multiply by $\prod_{k \neq i,j} W_{l_k}$ and sum over $i < j$ to obtain

$$\begin{aligned} & \mu \sum_{i < j} \sum_{x=1}^{|l_i|} \sum_{y=1}^{|l_j|} \left(\mathbf{1}_{e_y=e_x} W_{l_i \ominus_{x,y} l_j} - \mathbf{1}_{e_y=e_x^{-1}} W_{l_i \oplus_{x,y} l_j} \right) \prod_{k \neq i,j} W_{l_k} \\ &= \frac{\mu}{2} \sum_{s' \in \mathbb{M}^-(s) \setminus \mathbb{M}_U^-(s)} W_{s'} - \frac{\mu}{2} \sum_{s' \in \mathbb{M}^+(s) \setminus \mathbb{M}_U^+(s)} W_{s'}. \end{aligned}$$

Similarly, collecting the terms with coefficient λ , multiplying by $\prod_{k \neq i,j} W_{l_k}$ and summing over $i < j$ yields

$$\lambda \sum_{i < j} \sum_{x=1}^{|l_i|} \sum_{y=1}^{|l_j|} \left(\mathbf{1}_{e_y=e_x} W_{l_i \oplus_{x,y} l_j} - \mathbf{1}_{e_y=e_x^{-1}} W_{l_i \ominus_{x,y} l_j} \right) \prod_{k \neq i,j} W_{l_k} = \frac{\lambda}{2} \sum_{s' \in \mathbb{M}_U^+(s)} W_{s'} - \frac{\lambda}{2} \sum_{s' \in \mathbb{M}_U^-(s)} W_{s'}.$$

Finally, we turn to the terms with coefficient ν , and begin by introducing the following notation: for any edge $e \in E^+$ we let $A_r(e)$ be the set of locations in l_r where e occurs and $B_r(e)$ be the set of locations in l_r where e^{-1} occurs. The terms with coefficient ν are given by $\sum_{x,y}$ of

$$\begin{aligned} & \nu W_{l_i} W_{l_j} \left(\mathbf{1}_{e_x=e_y} - \mathbf{1}_{e_x=e_y^{-1}} \right) \\ &= \nu W_{l_i} W_{l_j} \sum_{e \in E^+} \left(\mathbf{1}_{e_x=e} \mathbf{1}_{e_y=e} + \mathbf{1}_{e_x=e^{-1}} \mathbf{1}_{e_y=e^{-1}} - \mathbf{1}_{e_x=e} \mathbf{1}_{e_y=e^{-1}} - \mathbf{1}_{e_x=e^{-1}} \mathbf{1}_{e_y=e} \right). \end{aligned}$$

Summing over x, y we obtain

$$\begin{aligned} & \nu W_{l_i} W_{l_j} \sum_{e \in E^+} \left(|A_i(e)| |A_j(e)| + |B_i(e)| |B_j(e)| - |A_i(e)| |B_j(e)| - |B_i(e)| |A_j(e)| \right) \\ &= \nu W_{l_i} W_{l_j} \sum_{e \in E^+} \left((|A_i(e)| - |B_i(e)|) (|A_j(e)| - |B_j(e)|) \right) = \nu W_{l_i} W_{l_j} \sum_{e \in E^+} t_i(e) t_j(e), \end{aligned}$$

with $t_i(e) = |A_i(e)| - |B_i(e)|$. Substituting the ν part into \mathcal{I}_s , we get the last term in (4.12).

STEP 3. In this step, we analyze the gradient terms \mathcal{D}_l and claim that

$$\mathcal{D}_l = \begin{cases} -\frac{1}{4} N \beta \sum_{s' \in \mathbb{D}^+(l)} W_{s'} & + \frac{1}{4} N \beta \sum_{s' \in \mathbb{D}^-(l)} W_{s'} & \text{for } G \in \{SO(N), U(N)\} \\ -\frac{1}{4} N \beta \sum_{s' \in \mathbb{D}^+(l) \cup \mathbb{E}^+(l)} W_{s'} & + \frac{1}{4} N \beta \sum_{s' \in \mathbb{D}^-(l) \cup \mathbb{E}^-(l)} W_{s'} & \text{for } G = SU(N) \end{cases} \quad (4.14)$$

Recalling $\frac{1}{2}\nabla S(Q)_e$ given in (3.7), we first claim that for each of our Lie groups, and for every $e \in E^-$,

$$\frac{1}{2}(\nabla S(Q))_{e^{-1}}^* = \frac{1}{2}\nabla S(Q)_e \quad (4.15)$$

where the r.h.s. is given by the formula (3.7). Indeed, to calculate the l.h.s. of (4.15), note that

$$\sum_{p \in \mathcal{P}_\Lambda, p \succ e^{-1}} \left((Q_p - Q_p^*) Q_{e^{-1}} \right)^* = \sum_{p \in \mathcal{P}_\Lambda, p \succ e^{-1}} Q_e (Q_p^* - Q_p) = \sum_{\bar{p} \in \mathcal{P}_\Lambda, \bar{p} \succ e} (Q_{\bar{p}} - Q_{\bar{p}}^*) Q_e$$

where in the last step we made a change of variable $p = e^{-1}e_2e_3e_4 \mapsto \bar{p} = e e_4^{-1}e_3^{-1}e_2^{-1}$. This establishes (4.15) in case where $G \in \{SO(N), U(N)\}$. For the $SU(N)$ case we need to analyze the additional trace term, so noting that $\text{Tr}(Q_p - Q_p^*)$ is purely imaginary,

$$\begin{aligned} \sum_{p \in \mathcal{P}_\Lambda, p \succ e^{-1}} \left(\text{Tr}(Q_p - Q_p^*) Q_{e^{-1}} \right)^* &= - \sum_{p \in \mathcal{P}_\Lambda, p \succ e^{-1}} \text{Tr}(Q_p - Q_p^*) Q_e \\ &= - \sum_{\bar{p} \in \mathcal{P}_\Lambda, \bar{p} \succ e} \text{Tr}(Q_{\bar{p}}^{-1} - Q_{\bar{p}}) Q_e = \sum_{\bar{p} \in \mathcal{P}_\Lambda, \bar{p} \succ e} \text{Tr}(Q_{\bar{p}} - Q_{\bar{p}}^*) Q_e \end{aligned}$$

so (4.15) holds. In light of (4.15), the constraint in (4.5) on the orientation of the edge may be removed, and the expression for \mathcal{D}_l simplifies to

$$\mathcal{D}_l = \frac{1}{2} \sum_{x=1}^n \text{Tr} \left(\prod_{i=1}^{x-1} Q_{e_i} \nabla S(Q)_{e_x} \prod_{i=x+1}^n Q_{e_i} \right).$$

In light of (3.7), in the case $G \in \{SO(N), U(N)\}$, we find that \mathcal{D}_l is given by

$$-\frac{1}{4} N \beta \sum_{x=1}^n \sum_{p \in \mathcal{P}_\Lambda, p \succ e_x} \text{Tr} \left(\prod_{i=1}^{x-1} Q_{e_i} (Q_p - Q_p^{-1}) Q_{e_x} \prod_{i=x+1}^n Q_{e_i} \right) = -\frac{1}{4} N \beta \sum_{x=1}^n \sum_{p \succ e_x} \left(W_{l \oplus_x p} - W_{l \ominus_x p} \right).$$

which yields the first case in (4.14) taking into account that each $s' \in \mathbb{D}^+(l)$ is of the form $l \oplus_x p$ for a plaquette p which contains an edge e_x and analogously for $s' \in \mathbb{D}^-(l)$. Note that here and below, in the cycle p we choose the path where e_x is the first edge.

For the $SU(N)$ case, the quantity \mathcal{D}_l is given by adding to the quantity above the term

$$\frac{\beta}{4} \sum_{x=1}^n \sum_{p \in \mathcal{P}_\Lambda, p \succ e_x} \text{Tr} \left(\prod_{i=1}^n Q_{e_i} \right) \text{Tr} \left(Q_p - Q_p^{-1} \right) = \frac{\beta}{4} \sum_{x=1}^n \sum_{p \in \mathcal{P}_\Lambda, p \succ e_x} \left(W_l W_p - W_l W_{p^{-1}} \right), \quad (4.16)$$

which yields the second case in (4.14) taking into account that each $s' \in \mathbb{E}^-(l)$ is of the form $s' = (l, p)$ for a plaquette p which contains an edge e_x and analogously for $s' \in \mathbb{E}^+(l)$. Here we use $W_{p^{-1}} = W_{\bar{p}}$ with $\bar{p} \succ e_x^{-1}$.

STEP 4. In this final step, we consider each group $G \in \{SO(N), U(N), SU(N)\}$ and use Steps 1-3 to conclude the proof of the master equations (1.5)-(1.7). To this end, we use the identity (4.9) for each constituent loop l_i and insert it, together with (4.12) and (4.14) into (4.8). We note in advance that

$$N^{-m} \sum_{i=1}^m \sum_{s' \in \mathcal{O}(l_i)} W_{s'} \prod_{j \neq i} W_{l_j} = N^{-m} \sum_{s' \in \mathcal{O}(s)} W_{s'} \quad (4.17)$$

for each of the operations $\mathcal{O} \in \{\mathbb{S}^\pm, \mathbb{T}^\pm, \mathbb{D}^\pm, \mathbb{E}^\pm\}$. Depending on the operation \mathcal{O} , the quantity $W_{s'}$ on the r.h.s. of (4.17) should be normalized in one of three possible ways according to the following:

$$s' \in \{\mathbb{M}^+(s), \mathbb{M}^-(s)\} \rightarrow \frac{\mathbf{E} W_{s'}}{N^{m-1}} = \phi(s'). \quad (4.18a)$$

$$s' \in \{\mathbb{D}^+(s), \mathbb{D}^-(s), \mathbb{T}^+(s), \mathbb{T}^-(s)\} \rightarrow \frac{\mathbf{E}W_{s'}}{N^m} = \phi(s'). \quad (4.18b)$$

$$s' \in \{\mathbb{S}^+(s), \mathbb{S}^-(s), \mathbb{E}^+(s), \mathbb{E}^-(s)\} \rightarrow \frac{\mathbf{E}W_{s'}}{N^{m+1}} = \phi(s'). \quad (4.18c)$$

This follows since we have $m - 1$ loops in $W_{s'}$ for $s' \in \{\mathbb{M}^+(s), \mathbb{M}^-(s)\}$, m loops in $W_{s'}$ for $s' \in \{\mathbb{D}^+(s), \mathbb{D}^-(s), \mathbb{T}^+(s), \mathbb{T}^-(s)\}$, and $m + 1$ loops in $W_{s'}$ for $s' \in \{\mathbb{S}^+(s), \mathbb{S}^-(s)\}$. Also we have an extra W_p for the expansion term, which requires an extra $\frac{1}{N}$ in $W_{s'}$ for $s' \in \{\mathbb{E}^+(s), \mathbb{E}^-(s)\}$.

We now turn to each of the groups and simplify the r.h.s. of (4.8) in accordance with the observations above.

Let $G = SO(N)$. By (2.5a), we have $\lambda = -\frac{1}{2}$, $\mu = \frac{1}{2}$, $\nu = 0$. Since $c_{\mathfrak{so}(N)} = -\frac{1}{2}(N - 1)$ from (2.4), the identity (4.8) takes the form

$$(N - 1)|s|\phi(s) = \frac{4}{N^m} \mathbf{E} \left[\sum_{i=1}^m (\mathcal{D}_{l_i} + \mathcal{I}_{l_i}) \prod_{j \neq i} W_{l_j} + \mathcal{I}_s \right].$$

Substituting (4.14) for \mathcal{D}_{l_i} , (4.9) for \mathcal{I}_{l_i} and (4.12) for \mathcal{I}_s , the r.h.s. is equal to

$$\begin{aligned} & -N\beta \sum_{s' \in \mathbb{D}^+(s)} \frac{\mathbf{E}[W_{s'}]}{N^m} + N\beta \sum_{s' \in \mathbb{D}^-(s)} \frac{\mathbf{E}[W_{s'}]}{N^m} - N \sum_{s' \in \mathbb{S}^+(s)} \frac{\mathbf{E}[W_{s'}]}{N^{m+1}} + N \sum_{s' \in \mathbb{S}^-(s)} \frac{\mathbf{E}[W_{s'}]}{N^{m+1}} \\ & - \sum_{s' \in \mathbb{T}^+(s)} \frac{\mathbf{E}[W_{s'}]}{N^m} + \sum_{s' \in \mathbb{T}^-(s)} \frac{\mathbf{E}[W_{s'}]}{N^m} - \frac{1}{N} \sum_{s' \in \mathbb{M}^+(s)} \frac{\mathbf{E}[W_{s'}]}{N^{m-1}} + \frac{1}{N} \sum_{s' \in \mathbb{M}^-(s)} \frac{\mathbf{E}[W_{s'}]}{N^{m-1}} \end{aligned}$$

where we used (4.17). Identifying the summands with ϕ by (4.18) completes the proof of (1.5).

Let $G = U(N)$. By (2.5b), we have $\lambda = -1$, $\mu = \nu = 0$. Since $c_{\mathfrak{u}(N)} = -N$, the identity (4.8) takes the form

$$N|s|\phi(s) = \frac{2}{N^m} \mathbf{E} \left[\sum_{i=1}^m (\mathcal{D}_{l_i} + \mathcal{I}_{l_i}) \prod_{j \neq i} W_{l_j} + \mathcal{I}_s \right].$$

Again substituting (4.14)+(4.9)+(4.12) and using (4.17), the r.h.s. is equal to

$$\begin{aligned} & -\frac{N\beta}{2} \sum_{s' \in \mathbb{D}^+(s)} \frac{\mathbf{E}[W_{s'}]}{N^m} + \frac{N\beta}{2} \sum_{s' \in \mathbb{D}^-(s)} \frac{\mathbf{E}[W_{s'}]}{N^m} \\ & - \frac{1}{N} \sum_{s' \in \mathbb{M}_U^+(s)} \frac{\mathbf{E}[W_{s'}]}{N^{m-1}} + \frac{1}{N} \sum_{s' \in \mathbb{M}_U^-(s)} \frac{\mathbf{E}[W_{s'}]}{N^{m-1}} - N \sum_{s' \in \mathbb{S}^+(s)} \frac{\mathbf{E}[W_{s'}]}{N^{m+1}} + N \sum_{s' \in \mathbb{S}^-(s)} \frac{\mathbf{E}[W_{s'}]}{N^{m+1}}. \end{aligned} \quad (4.19)$$

Identifying the summands with ϕ by (4.18) completes the proof of (1.7).

Let $G = SU(N)$. By (2.5c), we have $\lambda = -1$, $\mu = 0$, $\nu = \frac{1}{N}$. Since $c_{\mathfrak{su}(N)} = -N + \frac{1}{N}$ by (2.4), the identity (4.8) takes the form

$$(N - \frac{1}{N})|s|\phi(s) = \frac{2}{N^m} \mathbf{E} \left[\sum_{i=1}^m (\mathcal{D}_{l_i} + \mathcal{I}_{l_i}) \prod_{j \neq i} W_{l_j} + \mathcal{I}_s \right].$$

We again apply (4.14)+(4.9)+(4.12) using (4.17), and note that the only differences from the $U(N)$ case are the ν -terms and the expansion terms. The r.h.s. is then equal to

$$\begin{aligned} (4.19) & -\frac{N\beta}{2} \sum_{s' \in \mathbb{E}^+(s)} \frac{\mathbf{E}[W_{s'}]}{N^{m+1}} + \frac{N\beta}{2} \sum_{s' \in \mathbb{E}^-(s)} \frac{\mathbf{E}[W_{s'}]}{N^{m+1}} \\ & - \frac{1}{N} |s| \frac{\mathbf{E}[W_s]}{N^m} + \frac{1}{N} \sum_{e \in E^+} \left(\sum_{i=1}^m t_i(e)^2 + 2 \sum_{i < j} t_i(e) t_j(e) \right) \frac{\mathbf{E}[W_s]}{N^m} \end{aligned} \quad (4.20)$$

where we used $\sum_{i=1}^m |l_i| = |s|$ and $\ell(l_i) = \sum_{e \in E^+} t_i(e)^2$. By (1.4), the second line of (4.20) is equal to $-\frac{1}{N}(|s| - \ell(s))\phi(s)$. This implies (1.6) and completes the proof. \square

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